CALCULUS

Gilbert Strang
Massachusetts Institute of Technology

WELLESLEY-CAMBRIDGE PRESS
Box 82-279 Wellesley MA 02181
Contents

CHAPTER 1  Introduction to Calculus
1.1 Velocity and Distance  1
1.2 Calculus Without Limits  8
1.3 The Velocity at an Instant  16
1.4 Circular Motion  22
1.5 A Review of Trigonometry  29
1.6 A Thousand Points of Light  34
1.7 Computing in Calculus  36

CHAPTER 2  Derivatives
2.1 The Derivative of a Function  44
2.2 Powers and Polynomials  50
2.3 The Slope and the Tangent Line  58
2.4 Derivative of the Sine and Cosine  64
2.5 The Product and Quotient and Power Rules  71
2.6 Limits  78
2.7 Continuous Functions  85

CHAPTER 3  Applications of the Derivative
3.1 Linear Approximation  91
3.2 Maximum and Minimum Problems  96
3.3 Second Derivatives: Minimum vs. Maximum  105
3.4 Graphs  112
3.5 Ellipses, Parabolas, and Hyperbolas  121
3.6 Iterations $x_{n+1} = F(x_n)$  130
3.7 Newton's Method and Chaos  137
3.8 The Mean Value Theorem and l'Hôpital's Rule  146
CHAPTER 4 The Chain Rule
4.1 Derivatives by the Chain Rule 154
4.2 Implicit Differentiation and Related Rates 160
4.3 Inverse Functions and Their Derivatives 164
4.4 Inverses of Trigonometric Functions 171

CHAPTER 5 Integrals
5.1 The Idea of the Integral 177
5.2 Antiderivatives 182
5.3 Summation vs. Integration 187
5.4 Indefinite Integrals and Substitutions 195
5.5 The Definite Integral 201
5.6 Properties of the Integral and the Average Value 206
5.7 The Fundamental Theorem and Its Consequences 213
5.8 Numerical Integration 220

CHAPTER 6 Exponentials and Logarithms
6.1 An Overview 228
6.2 The Exponential $e^x$ 236
6.3 Growth and Decay in Science and Economics 242
6.4 Logarithms 252
6.5 Separable Equations Including the Logistic Equation 259
6.6 Powers Instead of Exponentials 267
6.7 Hyperbolic Functions 277

CHAPTER 7 Techniques of Integration
7.1 Integration by Parts 283
7.2 Trigonometric Integrals 288
7.3 Trigonometric Substitutions 294
7.4 Partial Fractions 300
7.5 Improper Integrals 305

CHAPTER 8 Applications of the Integral
8.1 Areas and Volumes by Slices 311
8.2 Length of a Plane Curve 320
8.3 Area of a Surface of Revolution 325
8.4 Probability and Calculus 328
8.5 Masses and Moments 336
8.6 Force, Work, and Energy 342
<table>
<thead>
<tr>
<th>CHAPTER 9</th>
<th><strong>Polar Coordinates and Complex Numbers</strong></th>
</tr>
</thead>
<tbody>
<tr>
<td>9.1 Polar Coordinates</td>
<td>348</td>
</tr>
<tr>
<td>9.2 Polar Equations and Graphs</td>
<td>351</td>
</tr>
<tr>
<td>9.3 Slope, Length, and Area for Polar Curves</td>
<td>356</td>
</tr>
<tr>
<td>9.4 Complex Numbers</td>
<td>360</td>
</tr>
</tbody>
</table>

<table>
<thead>
<tr>
<th>CHAPTER 10</th>
<th><strong>Infinite Series</strong></th>
</tr>
</thead>
<tbody>
<tr>
<td>10.1 The Geometric Series</td>
<td>368</td>
</tr>
<tr>
<td>10.2 Convergence Tests: Positive Series</td>
<td>374</td>
</tr>
<tr>
<td>10.3 Convergence Tests: All Series</td>
<td>381</td>
</tr>
<tr>
<td>10.4 The Taylor Series for $e^x, \sin x,$ and $\cos x$</td>
<td>385</td>
</tr>
<tr>
<td>10.5 Power Series</td>
<td>391</td>
</tr>
</tbody>
</table>

<table>
<thead>
<tr>
<th>CHAPTER 11</th>
<th><strong>Vectors and Matrices</strong></th>
</tr>
</thead>
<tbody>
<tr>
<td>11.1 Vectors and Dot Products</td>
<td>398</td>
</tr>
<tr>
<td>11.2 Planes and Projections</td>
<td>407</td>
</tr>
<tr>
<td>11.3 Cross Products and Determinants</td>
<td>416</td>
</tr>
<tr>
<td>11.4 Matrices and Linear Equations</td>
<td>425</td>
</tr>
<tr>
<td>11.5 Linear Algebra in Three Dimensions</td>
<td>435</td>
</tr>
</tbody>
</table>

<table>
<thead>
<tr>
<th>CHAPTER 12</th>
<th><strong>Motion along a Curve</strong></th>
</tr>
</thead>
<tbody>
<tr>
<td>12.1 The Position Vector</td>
<td>446</td>
</tr>
<tr>
<td>12.2 Plane Motion: Projectiles and Cycloids</td>
<td>453</td>
</tr>
<tr>
<td>12.3 Tangent Vector and Normal Vector</td>
<td>459</td>
</tr>
<tr>
<td>12.4 Polar Coordinates and Planetary Motion</td>
<td>464</td>
</tr>
</tbody>
</table>

<table>
<thead>
<tr>
<th>CHAPTER 13</th>
<th><strong>Partial Derivatives</strong></th>
</tr>
</thead>
<tbody>
<tr>
<td>13.1 Surfaces and Level Curves</td>
<td>472</td>
</tr>
<tr>
<td>13.2 Partial Derivatives</td>
<td>475</td>
</tr>
<tr>
<td>13.3 Tangent Planes and Linear Approximations</td>
<td>480</td>
</tr>
<tr>
<td>13.4 Directional Derivatives and Gradients</td>
<td>490</td>
</tr>
<tr>
<td>13.5 The Chain Rule</td>
<td>497</td>
</tr>
<tr>
<td>13.6 Maxima, Minima, and Saddle Points</td>
<td>504</td>
</tr>
<tr>
<td>13.7 Constraints and Lagrange Multipliers</td>
<td>514</td>
</tr>
</tbody>
</table>
## Contents

### CHAPTER 14  Multiple Integrals
- 14.1 Double Integrals ........................................ 521
- 14.2 Changing to Better Coordinates ......................... 527
- 14.3 Triple Integrals ........................................... 536
- 14.4 Cylindrical and Spherical Coordinates ................. 541

### CHAPTER 15  Vector Calculus
- 15.1 Vector Fields ............................................ 549
- 15.2 Line Integrals ............................................ 555
- 15.3 Green’s Theorem .......................................... 563
- 15.4 Surface Integrals ......................................... 573
- 15.5 The Divergence Theorem ................................ 582
- 15.6 Stokes’ Theorem and the Curl of \( \mathbf{F} \) ................. 589

### CHAPTER 16  Mathematics after Calculus
- 16.1 Linear Algebra ........................................... 599
- 16.2 Differential Equations .................................... 603
- 16.3 Discrete Mathematics ..................................... 611

  Study Guide For Chapter 1 .............................. G-0

  Answers to Odd-Numbered Problems ..................... A-0

  Index .............................................................. I-0

  Table of Integrals ............................................ T-0
Contents

CHAPTER 1  Introduction to Calculus
1.1 Velocity and Distance 1
1.2 Calculus Without Limits 8
1.3 The Velocity at an Instant 16
1.4 Circular Motion 22
1.5 A Review of Trigonometry 29
1.6 A Thousand Points of Light 34
1.7 Computing in Calculus 36

CHAPTER 2  Derivatives
2.1 The Derivative of a Function 44
2.2 Powers and Polynomials 50
2.3 The Slope and the Tangent Line 58
2.4 Derivative of the Sine and Cosine 64
2.5 The Product and Quotient and Power Rules 71
2.6 Limits 78
2.7 Continuous Functions 85

CHAPTER 3  Applications of the Derivative
3.1 Linear Approximation 91
3.2 Maximum and Minimum Problems 96
3.3 Second Derivatives: Minimum vs. Maximum 105
3.4 Graphs 112
3.5 Ellipses, Parabolas, and Hyperbolas 121
3.6 Iterations \( x_{n+1} = F(x_n) \) 130
3.7 Newton’s Method and Chaos 137
3.8 The Mean Value Theorem and l’Hôpital’s Rule 146
The right way to begin a calculus book is with calculus. This chapter will jump directly into the two problems that the subject was invented to solve. You will see what the questions are, and you will see an important part of the answer. There are plenty of good things left for the other chapters, so why not get started?

The book begins with an example that is familiar to everybody who drives a car. It is calculus in action—the driver sees it happening. The example is the relation between the speedometer and the odometer. One measures the speed (or velocity); the other measures the distance traveled. We will write $v$ for the velocity, and $f$ for how far the car has gone. The two instruments sit together on the dashboard:

![Fig. 1.1 Velocity $v$ and total distance $f$ (at one instant of time).](image)

Notice that the units of measurement are different for $v$ and $f$. The distance $f$ is measured in kilometers or miles (it is easier to say miles). The velocity $v$ is measured in km/hr or miles per hour. A unit of time enters the velocity but not the distance. Every formula to compute $v$ from $f$ will have $f$ divided by time.

The central question of calculus is the relation between $v$ and $f$. 

The right way to begin a calculus book is with calculus. This chapter will jump directly into the two problems that the subject was invented to solve. You will see what the questions are, and you will see an important part of the answer. There are plenty of good things left for the other chapters, so why not get started?
Can you find $v$ if you know $f$, and vice versa, and how? If we know the velocity over the whole history of the car, we should be able to compute the total distance traveled. In other words, if the speedometer record is complete but the odometer is missing, its information could be recovered. One way to do it (without calculus) is to put in a new odometer and drive the car all over again at the right speeds. That seems like a hard way; calculus may be easier. But the point is that the information is there. If we know everything about $v$, there must be a method to find $f$.

What happens in the opposite direction, when $f$ is known? If you have a complete record of distance, could you recover the complete velocity? In principle you could drive the car, repeat the history, and read off the speed. Again there must be a better way.

The whole subject of calculus is built on the relation between $v$ and $f$. The question we are raising here is not some kind of joke, after which the book will get serious and the mathematics will get started. On the contrary, I am serious now—and the mathematics has already started. We need to know how to find the velocity from a record of the distance. (That is called differentiation, and it is the central idea of differential calculus.) We also want to compute the distance from a history of the velocity. (That is integration, and it is the goal of integral calculus.)

Differentiation goes from $f$ to $v$; integration goes from $v$ to $f$. We look first at examples in which these pairs can be computed and understood.

**CONSTANT VELOCITY**

Suppose the velocity is fixed at $v = 60$ (miles per hour). Then $f$ increases at this constant rate. After two hours the distance is $f = 120$ (miles). After four hours $f = 240$ and after $t$ hours $f = 60t$. We say that $f$ increases linearly with time—its graph is a straight line.

![Constant velocity](image)

Notice that this example starts the car at full velocity. No time is spent picking up speed. (The velocity is a "step function.") Notice also that the distance starts at zero; the car is new. Those decisions make the graphs of $v$ and $f$ as neat as possible. One is the horizontal line $v = 60$. The other is the sloping line $f = 60t$. This $v, f, t$ relation needs algebra but not calculus:

*if $v$ is constant and $f$ starts at zero then $f = vt$.*

The opposite is also true. When $f$ increases linearly, $v$ is constant. The division by time gives the slope. The distance is $f_1 = 120$ miles when the time is $t_1 = 2$ hours. Later $f_2 = 240$ at $t_2 = 4$. At both points, the ratio $f/t$ is 60 miles/hour. Geometrically, the velocity is the slope of the distance graph:

\[
\text{slope} = \frac{\text{change in distance}}{\text{change in time}} = \frac{vt}{t} = v.
\]
1.1 Velocity and Distance

The slope of the $f$-graph gives the $v$-graph. Figure 1.3 shows two more possibilities:

1. The distance starts at 20 instead of 0. The distance formula changes from $60t$ to $20 + 60t$. The number 20 cancels when we compute change in distance—so the slope is still 60.

2. When $v$ is negative, the graph of $f$ goes downward. The car goes backward and the slope of $f = -30t$ is $v = -30$.

I don't think speedometers go below zero. But driving backwards, it's not that safe to watch. If you go fast enough, Toyota says they measure "absolute values"—the speedometer reads $+30$ when the velocity is $-30$. For the odometer, as far as I know it just stops. It should go backward.?+

VELOCITY vs. DISTANCE: SLOPE vs. AREA

How do you compute $f$ from $v$? The point of the question is to see $f = vt$ on the graphs. We want to start with the graph of $v$ and discover the graph of $f$. Amazingly, the opposite of slope is area.

The distance $f$ is the area under the $v$-graph. When $v$ is constant, the region under the graph is a rectangle. Its height is $v$, its width is $t$, and its area is $v$ times $t$. This is integration, to go from $v$ to $f$ by computing the area. We are glimpsing two of the central facts of calculus.

| 4A | The slope of the $f$-graph gives the velocity $v$. The area under the $v$-graph gives the distance $f$. |

That is certainly not obvious, and I hesitated a long time before I wrote it down in this first section. The best way to understand it is to look first at more examples. The whole point of calculus is to deal with velocities that are not constant, and from now on $v$ has several values.

EXAMPLE (Forward and back) There is a motion that you will understand right away. The car goes forward with velocity $V$, and comes back at the same speed. To say it more correctly, the velocity in the second part is $-V$. If the forward part lasts until $t = 3$, and the backward part continues to $t = 6$, the car will come back where it started. The total distance after both parts will be $f = 0$.

?This actually happened in Ferris Bueller's Day Off, when the hero borrowed his father's sports car and ran up the mileage. At home he raised the car and drove in reverse. I forget if it worked.
1 Introduction to Calculus

\[ v(t) = \text{slope of } f(t) \]

\[ f(t) = 3V \]

\[ \text{velocity } V \]

\[ \text{velocity } -V \]

\[ t \]

\[ 3 \]

\[ 6 \]

The \( v \)-graph shows velocities \( V \) and \( -V \). The distance starts up with slope \( V \) and reaches \( f = 3V \). Then the car starts backward. The distance goes down with slope \( -V \) and returns to \( f = 0 \) at \( t = 6 \).

Notice what that means. The total area “under” the \( v \)-graph is zero! A negative velocity makes the distance graph go downward (negative slope). The car is moving backward. \textit{Area below the axis in the } v \text{-graph is counted as negative.}

FUNCTIONS

This forward-back example gives practice with a crucially important idea—the concept of a “function.” We seize this golden opportunity to explain functions:

\textit{The number } v(t) \textit{is the value of the function } v \textit{at the time } t.\]

The time \( t \) is the input to the function. The velocity \( v(t) \) at that time is the output. Most people say “\( v \) of \( t \)” when they read \( v(t) \). The number “\( v \) of \( 2 \)” is the velocity when \( t = 2 \). The forward-back example has \( v(2) = V \) and \( v(4) = -V \). The function contains the whole history, like a memory bank that has a record of \( v \) at each \( t \).

It is simple to convert forward-back motion into a formula. Here is \( v(t) \):

\[ v(t) = \begin{cases} +V & \text{if } 0 < t < 3 \\ ? & \text{if } t = 3 \\ -V & \text{if } 3 < t < 6 \end{cases} \]

The right side contains the instructions for finding \( v(t) \). The input \( t \) is converted into the output \( +V \) or \( -V \). The velocity \( v(t) \) depends on \( t \). In this case the function is “discontinuous,” because the needle jumps at \( t = 3 \). \textit{The velocity is not defined at that instant.} There is no \( v(3) \). (You might argue that \( v \) is zero at the jump, but that leads to trouble.) The graph of \( v \) has a corner, and we can’t give its slope.

The problem also involves a second function, namely the distance. The principle behind \( f(t) \) is the same: \( f(t) \) \textit{is the distance at time } t. It is the net distance forward, and again the instructions change at \( t = 3 \). In the forward motion, \( f(t) \) equals \( Vt \) as before. In the backward half, a calculation is built into the formula for \( f(t) \):

\[ f(t) = \begin{cases} Vt & \text{if } 0 \leq t \leq 3 \\ V(6 - t) & \text{if } 3 \leq t \leq 6 \end{cases} \]

At the switching time the right side gives two instructions (one on each line). This would be bad except that they agree: \( f(3) = 3V \).† The distance function is “con-
There is no jump in \( f \), even when there is a jump in \( v \). After \( t = 3 \) the distance decreases because of \(- V\). At \( t = 6 \) the second instruction correctly gives \( f(6) = 0 \).

Notice something more. The functions were given by graphs before they were given by formulas. The graphs tell you \( f \) and \( v \) at every time \( t \)—sometimes more clearly than the formulas. The values \( f(t) \) and \( v(t) \) can also be given by tables or equations or a set of instructions. (In some way all functions are instructions—the function tells how to find \( f \) at time \( t \).) Part of knowing \( f \) is knowing all its inputs and outputs—its domain and range:

The domain of a function is the set of inputs. The range is the set of outputs.

The domain of \( f \) consists of all times \( 0 \leq t \leq 6 \). The range consists of all distances \( 0 \leq f(t) \leq 3V \). (The range of \( v \) contains only the two velocities \( +V \) and \(-V \).) We mention now, and repeat later, that every “linear” function has a formula \( f(t) = vt + C \). Its graph is a line and \( v \) is the slope. The constant \( C \) moves the line up and down. It adjusts the line to go through any desired starting point.

SUMMARY: MORE ABOUT FUNCTIONS

May I collect together the ideas brought out by this example? We had two functions \( v \) and \( f \). One was velocity, the other was distance. Each function had a domain, and a range, and most important a graph. For the \( f \)-graph we studied the slope (which agreed with \( v \)). For the \( v \)-graph we studied the area (which agreed with \( f \)).

Calculus produces functions in pairs, and the best thing a book can do early is to show you more of them.

\[
\begin{align*}
\text{input} & \rightarrow \text{function} & \text{output} \\
\text{domain} & \rightarrow f(t) & \text{range}
\end{align*}
\]

Note about the definition of a function. The idea behind the symbol \( f(t) \) is absolutely crucial to mathematics. Words don’t do it justice! By definition, a function is a “rule” that assigns one member of the range to each member of the domain. Or, a function is a set of pairs \((t, f(t))\) with no \( t \) appearing twice. (These are “ordered pairs” because we write \( t \) before \( f(t) \).) Both of those definitions are correct—but somehow they are too passive.

In practice what matters is the active part. The number \( f(t) \) is produced from the number \( t \). We read a graph, plug into a formula, solve an equation, run a computer program. The input \( t \) is “mapped” to the output \( f(t) \), which changes as \( t \) changes. Calculus is about the rate of change. This rate is our other function \( v \).

![Fig. 1.5 Subtracting 2 from \( f \) affects the range. Subtracting 2 from \( t \) affects the domain.](image)
1 Introduction to Calculus

It is quite hard at the beginning, and not automatic, to see the difference between \( f(t) - 2 \) and \( f(t^2) \). Those are both new functions, created out of the original \( f(t) \). In \( f(t) - 2 \), we subtract 2 from all the distances. That moves the whole graph down. In \( f(t^2) \), we subtract 2 from the time. That moves the graph over to the right. Figure 1.5 shows both movements, starting from \( f(t) = 2t + 1 \). The formula to find \( f(t^2) \) is \( 2(t - 2) + 1 \), which is \( 2t - 3 \).

A graphing calculator also moves the graph, when you change the viewing window. You can pick any rectangle \( A < t < B \), \( C < f(t) < D \). The screen shows that part of the graph. But on the calculator, the function \( f(t) \) remains the same. It is the axes that get renumbered. In our figures the axes stay the same and the function is changed.

There are two more basic ways to change a function. (We are always creating new functions—that is what mathematics is all about.) Instead of subtracting or adding, we can multiply the distance by 2. Figure 1.6 shows \( 2f(t) \). And instead of shifting the time, we can speed it up. The function becomes \( f(2t) \). Everything happens twice as fast (and takes half as long). On the calculator those changes correspond to a "zoom"—on the \( f \) axis or the \( t \) axis. We soon come back to zooms.

```
<table>
<thead>
<tr>
<th>domain 1</th>
<th>1</th>
<th>1</th>
<th>2</th>
</tr>
</thead>
<tbody>
<tr>
<td>range</td>
<td>3</td>
<td>1</td>
<td>2</td>
</tr>
</tbody>
</table>

\( f(t) = 2t - 1 \)

slop 2

<table>
<thead>
<tr>
<th>domain 1/2</th>
<th>0</th>
<th>1</th>
</tr>
</thead>
<tbody>
<tr>
<td>range</td>
<td>3</td>
<td>1</td>
</tr>
</tbody>
</table>

\( 2f(t) = 4t + 2 \)
slop 4

<table>
<thead>
<tr>
<th>domain 1</th>
<th>0</th>
<th>1</th>
</tr>
</thead>
<tbody>
<tr>
<td>range</td>
<td>3</td>
<td>1</td>
</tr>
</tbody>
</table>

\( f(2t) = 4t + 1 \)
slop 4

<table>
<thead>
<tr>
<th>domain 1/2</th>
<th>0</th>
<th>1/2</th>
</tr>
</thead>
<tbody>
<tr>
<td>range</td>
<td>3</td>
<td>1</td>
</tr>
</tbody>
</table>
```

Fig. 1.6 Doubling the distance or speeding up the time doubles the slope.

### 1.1 Exercises

Each section of the book contains read-through questions. They allow you to outline the section yourself—more actively than reading a summary. This is probably the best way to remember the important ideas.

Starting from \( f(0) = 0 \) at constant velocity \( v \), the distance function is \( f(t) = \_a\_ \). When \( f(t) = 5s \) the velocity is \( v = \_b\_. \) When \( f(t) = 5s + 1000 \) the velocity is still \( c \) and the starting value is \( f(0) = \_d\_. \) In each case \( v \) is the \( \_e\_ \) of the graph of \( f \). When \( t \) is negative, the graph of \( f \) goes downward. In that case area in the \( t \)-graph counts as \( \_f\_. \)

Forward motion from \( f(0) = 0 \) to \( f(2) = 10 \) has \( v = \_g\_. \) Then backward motion to \( f(4) = 0 \) has \( v = \_h\_. \) The distance function is \( f(t) = 5t \) for \( 0 < t < 2 \) and then \( f(t) = \_k\_ \) (not \(-5t\)). The slopes are \( \_l\_ \) and \( \_m\_. \) The distance \( f(3) = \_n\_. \) The area under the \( t \)-graph up to time 1.5 is \( \_o\_. \) The domain of \( f \) is the time interval \( \_p\_ \) and the range is the distance interval \( \_q\_. \) The range of \( v(t) \) is only \( \_r\_. \)

The value of \( f(t) = 3t + 1 \) at \( t = 2 \) is \( f(2) = \_s\_. \) The value \( 19 \) equals \( f(\_t\_ \_u\_. \) The difference \( f(4) - f(1) = \_v\_. \) That is the change in distance, when \( 4 - 1 \) is the change in \( \_w\_. \) The ratio of those changes equals \( \_x\_ \), which is the \( \_y\_ \) of the graph. The formula for \( f(t) + 2 \) is \( 3t + 3 \) whereas \( f(t + 2) \) equals \( \_z\_. \) Those functions have the same \( x \) as \( f \): the graph of \( f(t) + 2 \) is shifted \( \_A\_ \) and \( f(t + 2) \) is shifted \( \_B\_ \). The formula for \( 5f(t) \) is \( \_C\_ \). The formula for \( 5f(t) \) is \( \_D\_ \). The slope has jumped from 3 to \( \_F\_. \)
### 1.1 Velocity and Distance

The set of inputs to a function is its \( F \). The set of outputs is its \( G \). The functions \( f(t) = 7 + 3(t - 2) \) and \( f(t) = vt + C \) are \( H \). Their graphs are \( I \) with slopes equal to \( J \) and \( K \). They are the same function, if \( v = L \) and \( C = M \).

Draw the velocity graph that goes with each distance graph.

#### Problem 3
Write down three-part formulas for the velocities \( u(t) \) in Problem 2, starting from \( v(t) = 2 \) for \( 0 < t < 1 \).

#### Problem 4
The distance in \( lb \) starts with \( f(t) = 10 - 10t \) for \( 0 < t < 1 \). Give a formula for the second part.

#### Problem 5
In the middle of graph 2a find \( f(15) \) and \( f(12) \) and \( f(t) \).

#### Problem 6
In graph 2b find \( f(1.4T) \). If \( T = 3 \) what is \( f(4) \)?

#### Problem 7
Find the average speed between \( t = 0 \) and \( t = 5 \) in graph 1a. What is the speed at \( t = 5 \)?

#### Problem 8
What is the average speed between \( t = 0 \) and \( t = 2 \) in graph 1b? The average speed is zero between \( t = \frac{1}{2} \) and \( t = \) ________.

#### Problem 9 (recommended)
A car goes at speed \( v = 20 \) into a brick wall at distance \( f = 4 \). Give two-part formulas for \( v(t) \) and \( f(t) \) (before and after), and draw the graphs.

#### Problem 10
Draw any reasonable graphs of \( v(t) \) and \( f(t) \) when
(a) the driver backs up, stops to shift gear, then goes fast;
(b) the driver slows to 55 for a police car;
(c) in a rough gear change, the car accelerates in jumps;
(d) the driver waits for a light that turns green.

#### Problem 11
Your bank account earns simple interest on the opening balance \( f(0) \). What are the interest rates per year?

#### Problem 12
The earth's population is growing at \( v = 100 \) million a year, starting from \( f = 5.2 \) billion in 1990. Graph \( f(t) \) and find \( f(2000) \).

#### Problem 13
Draw the distance graph that goes with each velocity graph. Start from \( f = 0 \) at \( t = 0 \) and mark the distance.

#### Problem 14
Write down formulas for \( v(t) \) in Problem 14, starting with \( v = -40 \) for \( 0 < t < 1 \). Find the average velocities to \( t = 2.5 \) and \( t = 3T \).

#### Problem 15
Give 3-part formulas for the areas \( f(t) \) under \( v(t) \) in 13.

#### Problem 16
The distance in 14a starts with \( f(t) = -40t \) for \( 0 < t < 1 \). Find \( f(t) \) in the other part, which passes through \( f = 0 \) at \( t = 2 \).

#### Problem 17
Draw rough graphs of \( y = \sqrt{x} \) and \( y = \sqrt{x - 4} \) and \( y = \sqrt{x} - 4 \). They are "half-parabolas" with infinite slope at the start.

#### Problem 18
What is the break-even point if \( x \) yearbooks cost \$1200 + 30x to produce and the income is 40x? The slope of the cost line is ________ (cost per additional book). If it goes above ________, you can't break even.

#### Problem 19
What are the domains and ranges of the distance functions in 14a and 14b—all values of \( t \) and \( f(t) \) if \( f(0) = 0 \)?

#### Problem 20
What is the range of \( u(t) \) in 14b? Why is \( t = 1 \) not in the domain of \( v(t) \) in 14a?

Problems 23–28 involve linear functions \( f(t) = vt + C \). Find the constants \( v \) and \( C \).

#### Problem 23
What linear function has \( f(0) = 3 \) and \( f(2) = -11 \)?

#### Problem 24
Find two linear functions whose domain is \( 0 \leq t \leq 2 \) and whose range is \( 1 \leq f(t) \leq 9 \).

#### Problem 25
Find the linear function with \( f(1) = 4 \) and slope 6.

#### Problem 26
What functions have \( f(t + 1) = f(t) + 2 \)?

#### Problem 27
Find the linear function with \( f(t + 2) = f(t) + 6 \) and \( f(1) = 10 \).

#### Problem 28
Find the only \( f = vt \) that has \( f(2t) = 4f(t) \). Show that every \( f = \frac{1}{2}vt^2 \) has this property. To go ________ times as far in twice the time, you must accelerate.
Sketch the graph of \( f(t) = |3 - 2t| \) (absolute value) for \(|t| \leq 2\) and find its slopes and range.

Sketch the graph of \( f(t) = 4 - t - |4 - t| \) for \( 2 \leq t \leq 5 \) and find its slope and range.

Suppose \( v = 8 \) up to time \( T \), and after that \( v = -2 \). Starting from zero, when does \( f \) return to zero? Give formulas for \( v(t) \) and \( f(t) \).

What function \( F(C) \) converts Celsius temperature \( C \) to Fahrenheit temperature \( F \)? The slope is \( \frac{9}{5} \), which is the number of Fahrenheit degrees equivalent to 1°C.

What function \( C(F) \) converts Fahrenheit to Celsius (or Centigrade), and what is its slope?

What function converts the weight \( w \) in grams to the weight \( f(w) \) in kilograms? Interpret the slope of \( f(w) \).

(Food of March 1989) Ten hours after the accident the alcohol reading was .061. Blood alcohol is eliminated at .015 per hour. What was the reading at the time of the accident? Which points between 39 and 43 can be in the domain of \( f(t) \)? With this domain find the range in 37–42.

(a) Draw the graph of \( f(t) = \frac{1}{2}t + 3 \) with domain \( 0 \leq t \leq 2 \). Then give a formula and graph for

\( f(t) = \frac{1}{2}t + 3 \)
\( f(t) = \frac{1}{2}t + 1 \)
\( f(t) = f(t+1) \)
\( f(t+1) \)
\( f(4t) \)
\( f(4t) \)

(a) Draw the graph of \( U(t) \) = step function = 0 for \( t < 0 \), 1 for \( t \geq 0 \). Then draw

\( U(t) + 2 \)
\( U(t+2) \)
\( 3U(t) \)
\( U(3t) \)

(a) Draw the graph of \( f(t) = t + 1 \) for \(-1 \leq t \leq 1\). Find the domain, range, slope, and formula for

\( 2f(t) \)
\( f(t-3) \)
\( -f(t) \)
\( f(-t) \)

If \( f(t) = t - 1 \) what are \( 2f(3t) \) and \( f(1-t) \) and \( f(t-1) \)?

In the forward-back example find \( f(\frac{1}{2}T) \) and \( f(\frac{1}{3}T) \). Verify that these agree with the areas "under" the v-graph in Figure 1.4.

Find formulas for the outputs \( f_1(t) \) and \( f_2(t) \) which come from the input \( t \):

(1) inside = input + 3
output = inside + 3
(2) inside ← input + 6
output ← inside + 3

Note BASIC and FORTRAN (and calculus itself) use = instead of ←. But the symbol ← or := is in some ways better. The instruction \( t ← t + 6 \) produces a new \( t \) equal to the old \( t \) plus six. The equation \( t = t + 6 \) is not intended.

Your computer can add and multiply. Starting with the number 1 and the input called \( t \), give a list of instructions to lead to these outputs:

\( f_1(t) = t^2 + t \)
\( f_2(t) = f_1(f_1(t)) \)
\( f_3(t) = f_1(t+1) \)

In fifty words or less explain what a function is.

The last questions are challenging but possible.

If \( f(t) = 3t - 1 \) for \( 0 \leq t \leq 2 \) give formulas (with domain) and find the slopes of these six functions:

\( f(t+2) \)
\( f(t) + 2 \)
\( 2f(t) \)
\( f(2t) \)
\( f(-t) \)
\( f(f(t)) \)

For \( f(t) = vt + C \) find the formulas and slopes of

\( f(3t+1) \)
\( f(3t+1) \)
\( 2f(4t) \)
\( f(-t) \)
\( f(t)-f(0) \)
\( f(f(t)) \)

(hardest) The forward-back function is \( f(t) = 2t \) for \( 0 \leq t \leq 3 \), \( f(t) = 12 - 2t \) for \( 3 \leq t \leq 6 \). Graph \( f(f(t)) \) and find its four-part formula. First try \( t = 1.5 \) and 3.

Why is the letter \( X \) not the graph of a function?

Which capital letters are the graphs of functions?

Draw graphs of their slopes.

### 1.2 Calculus Without Limits

The next page is going to reveal one of the key ideas behind calculus. The discussion is just about numbers—functions and slopes can wait. The numbers are not even special, they can be any numbers. The crucial point is to look at their differences:

Suppose the numbers are \( f = 0 \ 2 \ 6 \ 7 \ 4 \ 9 \)
Their differences are \( \Delta f = 2 \ 4 \ 1 \ -3 \ 5 \)

The differences are printed in between, to show \( 2 - 0 = 2 \) and \( 6 - 2 = 4 \) and \( 7 - 6 = 1 \).
1.2 Calculus Without Limits

Notice how $4 - 7$ gives a negative answer $-3$. The numbers in $f$ can go up or down, the differences in $v$ can be positive or negative. The idea behind calculus comes when you add up those differences:

$$2 + 4 + 1 - 3 + 5 = 9$$

The sum of differences is 9. This is the last number on the top line (in $f$). Is this an accident, or is this always true? If we stop earlier, after $2 + 4 + 1$, we get the 7 in $f$. Test any prediction on a second example:

Suppose the numbers are $f = 1 \ 3 \ 7 \ 8 \ 5 \ 10$
Their differences are $v = 2 \ 4 \ 1 \ -3 \ 5$

The $f$'s are increased by 1. The differences are exactly the same—no change. The sum of differences is still 9. But the last $f$ is now 10. That prediction is not right, we don't always get the last $f$.

The first $f$ is now 1. The answer 9 (the sum of differences) is $10 - 1$, the last $f$ minus the first $f$. What happens when we change the $f$'s in the middle?

Suppose the numbers are $f = 1 \ 5 \ 12 \ 7 \ 10$
Their differences are $v = 4 \ 7 \ -5 \ 3$

The differences add to $4 + 7 - 5 + 3 = 9$. This is still $10 - 1$. No matter what $f$'s we choose or how many, the sum of differences is controlled by the first $f$ and last $f$.

If this is always true, there must be a clear reason why the middle $f$'s cancel out.

The sum of differences is $(5 - 1) + (12 - 5) + (7 - 12) + (10 - 7) = 10 - 1$.
The 5's cancel, the 12's cancel, and the 7's cancel. It is only $10 - 1$ that doesn't cancel. This is the key to calculus!

1B The differences of the $f$'s add up to $(f_{\text{last}} - f_{\text{first}})$.

**EXAMPLE 1** The numbers grow linearly: $f = 2 \ 3 \ 4 \ 5 \ 6 \ 7$
Their differences are constant: $v = 1 \ 1 \ 1 \ 1 \ 1$

The sum of differences is certainly 5. This agrees with $7 - 2 = f_{\text{last}} - f_{\text{first}}$. The numbers in $v$ remind us of constant velocity. The numbers in $f$ remind us of a straight line $f = vt + C$. This example has $v = 1$ and the $f$'s start at 2. The straight line would come from $f = t + 2$.

**EXAMPLE 2** The numbers are squares: $f = 0 \ 1 \ 4 \ 9 \ 16$
Their differences grow linearly: $v = 1 \ 3 \ 5 \ 7$

$1 + 3 + 5 + 7$ agrees with $4^2 = 16$. It is a beautiful fact that the first $j$ odd numbers always add up to $j^2$. The $v$'s are the odd numbers, the $f$'s are perfect squares.

**Note** The letter $j$ is sometimes useful to tell which number in $f$ we are looking at. For this example the zeroth number is $f_0 = 0$ and the $j$th number is $f_j = j^2$. This is a part of algebra, to give a formula for the $f$'s instead of a list of numbers. We can also use $j$ to tell which difference we are looking at. The first $v$ is the first odd number $v_1 = 1$. The $j$th difference is the $j$th odd number $v_j = 2j - 1$. (Thus $v_4$ is $8 - 1 = 7$.) It is better to start the differences with $j = 1$, since there is no zeroth odd number $v_0$.

**With this notation the $j$th difference is** $v_j = f_j - f_{j-1}$. Sooner or later you will get comfortable with subscripts like $j$ and $j - 1$, but it can be later. The important point is that the sum of the $v$'s equals $f_{\text{last}} - f_{\text{first}}$. We now connect the $v$'s to slopes and the $f$'s to areas.
1 Introduction to Calculus

Figure 1.7 shows a natural way to graph Example 2, with the odd numbers in \( v \) and the squares in \( f \). Notice an important difference between the \( v \)-graph and the \( f \)-graph. The graph of \( f \) is "piecewise linear." We plotted the numbers in \( f \) and connected them by straight lines. The graph of \( v \) is "piecewise constant." We plotted the differences as constant over each piece. This reminds us of the distance-velocity graphs, when the distance \( f(t) \) is a straight line and the velocity \( v(t) \) is a horizontal line.

Now make the connection to slopes:

The slope of the \( f \)-graph is \( \frac{\text{distance up}}{\text{distance across}} = \frac{\text{change in } f}{\text{change in } t} = v. \)

Over each piece, the change in \( t \) (across) is 1. The change in \( f \) (upward) is the difference that we are calling \( v \). The ratio is the slope \( v/1 \) or just \( v \). The slope makes a sudden change at the breakpoints \( t = 1, 2, 3, \ldots \). At those special points the slope of the \( f \)-graph is not defined—we connected the \( v \)'s by vertical lines but this is very debatable. The main idea is that between the breakpoints, the slope of \( f(t) \) is \( v(t) \).

Now make the connection to areas:

The total area under the \( v \)-graph is \( f_{\text{last}} - f_{\text{first}} \).

This area, underneath the staircase in Figure 1.7, is composed of rectangles. The base of every rectangle is 1. The heights of the rectangles are the \( v \)'s. So the areas also equal the \( v \)'s, and the total area is the sum of the \( v \)'s. This area is \( f_{\text{last}} - f_{\text{first}} \).

Even more is true. We could start at any time and end at any later time—not necessarily at the special times \( t = 0, 1, 2, 3, 4 \). Suppose we stop at \( t = 3.5. \) Only half of the last rectangular area (under \( v = 7 \)) will be counted. The total area is \( 1 + 3 + 5 + \frac{1}{2}(7) = 12.5 \). This still agrees with \( f_{\text{last}} - f_{\text{first}} = 12.5 - 0 \). At this new ending time \( t = 3.5 \), we are only halfway up the last step in the \( f \)-graph. Halfway between 9 and 16 is 12.5.

4C The \( v \)'s are slopes of \( f(t) \). The area under the \( v \)-graph is \( f(t_{\text{end}}) - f(t_{\text{start}}) \).

This is nothing less than the Fundamental Theorem of Calculus. But we have only used algebra (no curved graphs and no calculations involving limits). For now the Theorem is restricted to piecewise linear \( f(t) \) and piecewise constant \( v(t) \). In Chapter 5 that restriction will be overcome.

Notice that a proof of \( 1 + 3 + 5 + 7 = 4^2 \) is suggested by Figure 1.7a. The triangle under the dotted line has the same area as the four rectangles under the staircase. The area of the triangle is \( \frac{1}{2} \cdot \text{base} \cdot \text{height} = \frac{1}{2} \cdot 4 \cdot 8 \), which is the perfect square \( 4^2 \). When there are \( j \) rectangles instead of 4, we get \( \frac{1}{2} \cdot j \cdot 2j = j^2 \) for the area.
The next examples show other patterns, where \( f \) and \( v \) increase exponentially or oscillate around zero. I hope you like them but I don’t think you have to learn them. They are like the special functions \( 2^t \) and \( \sin t \) and \( \cos t \)—except they go in steps. You get a first look at the important functions of calculus, but you only need algebra. \textit{Calculus is needed for a steadily changing velocity, when the graph of \( f \) is curved.}

The last example will be \textit{income tax}—which really does go, in steps. Then Section 1.3 will introduce the slope of a curve. The crucial step for curves is working with \textit{limits}. That will take us from algebra to calculus.

**EXPONENTIAL VELOCITY AND DISTANCE**

Start with the numbers \( f = 1, 2, 4, 8, 16 \). These are “powers of 2.” They start with the zeroth power, which is \( 2^0 = 1 \). The exponential starts at 1 and not 0. After \( j \) steps there are \( j \) factors of 2, and \( f_j \) equals \( 2^j \). \textit{Please recognize the difference between} \( 2^j \) \textit{and} \( j^2 \) \textit{and} \( 2^j \). The numbers \( 2^j \) grow linearly, the numbers \( j^2 \) grow quadratically, the numbers \( 2^j \) grow exponentially. At \( j = 10 \) these are 20 and 100 and 1024. The exponential \( 2^j \) quickly becomes much larger than the others.

The differences of \( f = 1, 2, 4, 8, 16 \) are exactly \( v = 1, 2, 4, 8 \). We get the same beautiful numbers. \textit{When the} \( f \) \textit{’s are powers of 2, so are the} \( v \) \textit{’s.} The formula \( v_j = 2^{j-1} \) is slightly different from \( f_j = 2^j \), because the first \( v \) is numbered \( v_1 \). (Then \( v_1 = 2^0 = 1 \). The zeroth power of every number is 1, except that 0\(^0\) is meaningless.) The two graphs in Figure 1.8 use the same numbers but they look different, because \( f \) is piecewise linear and \( v \) is piecewise constant.

![Fig. 1.6](image)

Where will calculus come in? It works with the smooth curve \( f(t) = 2^t \). This exponential growth is critically important for population and money in a bank and the national debt. You can spot it by the following test: \( v(t) \) is \textit{proportional to} \( f(t) \).

\textbf{Remark} The function \( 2^t \) is trickier than \( t^2 \). For \( f = t^2 \) the slope is \( v = 2t \). It is proportional to \( t \) and not \( t^2 \). For \( f = 2^t \) the slope is \( v = c2^t \), and we won’t find the constant \( c = .693 \ldots \) until Chapter 6. (The number \( c \) is the natural logarithm of 2.) Problem 37 estimates \( c \) with a calculator—the important thing is that it’s constant.

**OSCILLATING VELOCITY AND DISTANCE**

We have seen a forward-back motion, velocity \( V \) followed by \(-V\). That is oscillation of the simplest kind. The graph of \( f \) goes linearly up and linearly down. Figure 1.9 shows another oscillation that returns to zero, but the path is more interesting.

The numbers in \( f \) are now 0, 1, 1, 0, \(-1\), \(-1\), 0. Since \( f_6 = 0 \) the motion brings us back to the start. The whole oscillation can be repeated.
The differences in $v$ are 1, 0, -1, -1, 0, 1. They add up to zero, which agrees with $f_{\text{last}} - f_{\text{first}}$. It is the same oscillation as in $f$ (and also repeatable), but shifted in time.

The $f$-graph resembles (roughly) a sine curve. The $v$-graph resembles (even more roughly) a cosine curve. The waveforms in nature are smooth curves, while these are "digitized"—the way a digital watch goes forward in jumps. You recognize that the change from analog to digital brought the computer revolution. The same revolution is coming in CD players. Digital signals (off or on, 0 or 1) seem to win every time.

The piecewise $v$ and $f$ start again at $t = 6$. The ordinary sine and cosine repeat at $t = 2\pi$. A repeating motion is periodic—here the "period" is 6 or $2\pi$. (With $t$ in degrees the period is 360—a full circle. The period becomes $2\pi$ when angles are measured in radians. We virtually always use radians—which are degrees times $2\pi/360$.) A watch has a period of 12 hours. If the dial shows AM and PM, the period is ______.

![Fig. 1.9 Piecwise constant "cosine" and piecewise linear "sine." They both repeat.](image)

**A SHORT BURST OF SPEED**

The next example is a car that is driven fast for a short time. The speed is $V$ until the distance reaches $f = 1$, when the car suddenly stops. The graph of $f$ goes up linearly with slope $V$, and then across with slope zero:

$$v(t) = \begin{cases} V & \text{up to } t = T \\ 0 & \text{after } t = T \end{cases} \quad f(t) = \begin{cases} Vt & \text{up to } t = T \\ 1 & \text{after } t = T \end{cases}$$

This is another example of "function notation." Notice the general time $t$ and the particular stopping time $T$. The distance is $f(t)$. The domain of $f$ (the inputs) includes all times $t \geq 0$. The range of $f$ (the outputs) includes all distances $0 \leq f \leq 1$.

Figure 1.10 allows us to compare three cars—a Jeep and a Corvette and a Maserati. They have different speeds but they all reach $f = 1$. So the areas under the $v$-graphs are all 1. The rectangles have height $V$ and base $T = 1/V$.

![Fig. 1.10 Bursts of speed with $V_M T_M = V_C T_C = V_J T_J = 1$. Step function has infinite slope.](image)

**Optional remark** It is natural to think about faster and faster speeds, which means steeper slopes. The $f$-graph reaches 1 in shorter times. The extreme case is a step function, when the graph of $f$ goes straight up. This is the unit step $U(t)$, which is zero up to $t = 0$ and jumps immediately to $U = 1$ for $t > 0$. 
What is the slope of the step function? It is zero except at the jump. At that moment, which is \( t = 0 \), the slope is infinite. We don’t have an ordinary velocity \( v(t) \) — instead we have an impulse that makes the car jump. The graph is a spike over the single point \( t = 0 \), and it is often denoted by \( \delta \) — so the slope of the step function is called a "delta function." The area under the infinite spike is 1.

You are absolutely not responsible for the theory of delta functions! Calculus is about curves, not jumps.

Our last example is a real-world application of slopes and rates—to explain "how taxes work." Note especially the difference between tax rates and tax brackets and total tax. The rates are \( r \), the brackets are on \( x \), the total tax is \( f \).

**EXAMPLE 3** *Income tax is piecewise linear. The slopes are the tax rates .15, .28, .31.*

Suppose you are single with taxable income of \( x \) dollars (Form 1040, line 37—after all deductions). These are the 1991 instructions from the Internal Revenue Service:

If \( x \) is not over $20,350, the tax is 15% of \( x \).

If \( $20,350 < x < $49,300 \), the tax is $3052.50 + 28% of the amount over $20,350.

If \( x \) is over $49,300, the tax is $11,158.50 + 31% of the amount over $49,300.

The first bracket is \( 0 < x < $20,350 \). (The IRS never uses this symbol \( < \), but I think it is OK here. We know what it means.) The second bracket is \( $20,350 < x < $49,300 \). The top bracket \( x > $49,300 \) pays tax at the top rate of 31%. But only the income in that bracket is taxed at that rate.

Figure 1.11 shows the rates and the brackets and the tax due. Those are not average rates, they are marginal rates. Total tax divided by total income would be the average rate. The marginal rate of .28 or .31 gives the tax on each additional dollar of income—it is the slope at the point \( x \). Tax is like area or distance—it adds up. Tax rate is like slope or velocity—it depends where you are. This is often unclear in the news media.

**Question** What is the equation for the straight line in the top bracket?

**Answer** The bracket begins at \( x = $49,300 \) when the tax is \( f(x) = $11,158.50 \). The slope of the line is the tax rate .31. When we know a point on the line and the slope, we know the equation. This is important enough to be highlighted.

For \( x \) in the top bracket the tax is \( f(x) = $11,158.50 + .31(x - $49,300) \). This is the tax on $49,300 plus the extra tax on extra income.

Section 2.3 presents this "point-slope equation" for any straight line. Here you see it for one specific example. Where does the number $11,158.50 come from? It is the tax at the end of the middle bracket, so it is the tax at the start of the top bracket.
1 Introduction to Calculus

Figure 1.1 also shows a distance-velocity example. The distance at \( t = 2 \) is \( f(2) = 40 \) miles. After that time the velocity is 60 miles per hour. So the line with slope 60 on the \( f \)-graph has the equation

\[
f(t) = \text{starting distance} + \text{extra distance} = 40 + 60(t - 2).
\]

The starting point is \((2, 40)\). The new speed 60 multiplies the extra time \( t - 2 \). The point-slope equation makes sense. We now review this section, with comments.

Central idea Start with any numbers in \( f \). Their differences go in \( v \). Then the sum of those differences is \( f_{\text{last}} - f_{\text{first}} \).

Subscript notation The numbers are \( f_0, f_1, \ldots \) and the first difference is \( v_1 = f_1 - f_0 \). A typical number is \( f_j \) and the \( j \)th difference is \( v_j = f_j - f_{j-1} \). When those differences are added, all \( f \)'s in the middle (like \( f_1 \)) cancel out:

\[
v_1 + v_2 + \cdots + v_j = (f_1 - f_0) + (f_2 - f_1) + \cdots + (f_j - f_{j-1}) = f_j - f_0.
\]

Examples \( f_j = j \) or \( j^2 \) or \( 2^j \). Then \( v_j = 1 \) (constant) or \( 2j - 1 \) (odd numbers) or \( 2^{j-1} \).

Functions Connect the \( f \)'s to be piecewise linear. Then the slope \( v \) is piecewise constant. The area under the \( v \)-graph from any \( t_{\text{start}} \) to any \( t_{\text{end}} \) equals \( f(t_{\text{end}}) - f(t_{\text{start}}) \).

Units Distance in miles and velocity in miles per hour. Tax in dollars and tax rate in \((\text{dollars paid})/(\text{dollars earned})\). Tax rate is a percentage like .28, with no units.

4.2 Exercises

Read-through questions

Start with the numbers \( f = 1, 6, 2, 5 \). Their differences are \( v = \underline{\ldots} \). The sum of those differences is \( \underline{\ldots} \). This is equal to \( f_{\text{last}} \) minus \( f_{\text{first}} \). The numbers 6 and 2 have no effect on this answer, because in \((6-1) + (2-6) + (5-2)\) the numbers 6 and 2 \( \underline{\ldots} \). The slope of the line between \( f(0) = 1 \) and \( f(1) = 6 \) is \( \underline{\ldots} \). The equation of that line is \( f(t) = \underline{\ldots} \).

With distances 1, 5, 25 at unit times, the velocities are \( \underline{\ldots} \). These are the \( \underline{\ldots} \) of the \( f \)-graph. The slope of the tax graph is the tax \( \underline{\ldots} \). If \( f(t) \) is the postage cost for \( t \) ounces or \( t \) grams, the slope is \( \underline{\ldots} \) per \( \underline{\ldots} \). For distances 0, 1, 4, 9 the velocities are \( \underline{\ldots} \). The sum of the first \( j \) odd numbers is \( f_j = \underline{\ldots} \). Then \( f_{10} = \underline{\ldots} \) and the velocity \( v_{10} \) is \( \underline{\ldots} \).

The piecewise linear sine has slopes \( \underline{\ldots} \). Those form a piecewise \( \underline{\ldots} \) cosine. Both functions have \( \underline{\ldots} \) equal to 6, which means that \( f(t + 6) = \underline{\ldots} \) for every \( t \). The velocities \( v = 1, 2, 4, 8, \ldots \) have \( v_j = \underline{\ldots} \). Then \( f_{10} = \underline{\ldots} \). The sum of 1, 2, 4, 8, 16 is \( \underline{\ldots} \). The difference \( 2^j - 2^{j-1} \) equals \( \underline{\ldots} \). After a burst of speed \( V \) to time \( T \), the distance is \( \underline{\ldots} \). If \( f(T) = 1 \) and \( V \) increases, the burst lasts only to \( T = \underline{\ldots} \). When \( V \) approaches infinity, \( f(t) \) approaches a \( \underline{\ldots} \) function. The velocities approach a \( \underline{\ldots} \) function, which is concentrated at \( t = 0 \) but has area \( \underline{\ldots} \) under its graph. The slope of a step function is \( \underline{\ldots} \).

Problems 1–4 are about numbers \( f \) and differences \( v \).

1. From the numbers \( f = 0, 2, 7, 10 \) find the differences \( v \) and the sum of the three \( v \)'s. Write down another \( f \) that leads to the same \( v \)'s. For \( f = 0, 3, 12, 10 \) the sum of the \( v \)'s is still \( \underline{\ldots} \).

2. Starting from \( f = 1, 3, 2, 4 \) draw the \( f \)-graph (linear pieces) and the \( v \)-graph. What are the areas “under” the \( v \)-graph that add to \( 4 \)? If the next number in \( f \) is 11, what is the area under the next \( v \)?

3. From \( v = 1, 2, 1, 0, -1 \) find the \( f \)'s starting at \( f_0 = 3 \). Graph \( v \) and \( f \). The maximum value of \( f \) occurs when \( v = \underline{\ldots} \). Where is the maximum \( f \) when \( v = 1, 2, 1, -1 \)?

4. For \( f = 1, b, c, 7 \) find the differences \( v_1, v_2, v_3 \) and add them up. Do the same for \( f = a, b, c, d \).

Problems 5–11 are about linear functions and constant slopes.

5. Write down the slopes of these linear functions:

(a) \( f(t) = 1.1t \) (b) \( f(t) = 1 - 2t \) (c) \( f(t) = 4 + 5(t - 6) \).

Compute \( f(6) \) and \( f(7) \) for each function and confirm that \( f(7) - f(6) \) equals the slope.

6. If \( f(t) = 5 + 3(t - 1) \) and \( g(t) = 1.5 + 2.5(t - 1) \) what is \( h(t) = f(t) - g(t) \)? Find the slopes of \( f, g, \) and \( h \).
7 Suppose \( v(t) = 2 \) for \( t < 5 \) and \( v(t) = 3 \) for \( t > 5 \).
   (a) If \( f(0) = 0 \) find a two-part formula for \( f(t) \).
   (b) Check that \( f(10) \) equals the area under the graph of \( v(t) \) (two rectangles) up to \( t = 10 \).
8 Suppose \( v(t) = 10 \) for \( t < 1/10 \), \( v(t) = 0 \) for \( t > 1/10 \). Starting from \( f(0) = 1 \) find \( f(t) \) in two pieces.
9 Suppose \( g(t) = 2t + 1 \) and \( f(t) = 4t \). Find \( g(3) \) and \( f(g(3)) \) and \( f(g(t)) \). How is the slope of \( f(g(t)) \) related to the slopes of \( f \) and \( g \)?
10 For the same functions, what are \( f(3) \) and \( g(f(3)) \) and \( g(f(t)) \)? When \( t \) is changed to \( 4t \), distance increases ______ times as fast and the velocity is multiplied by ________.
11 Compute \( f(6) \) and \( f(8) \) for the functions in Problem 5. Confirm that the slopes \( a \) agree with
   \[
   \text{slope} = \frac{f(8) - f(6)}{8 - 6} = \text{change in } f \over \text{change in } t.
   \]

Problems 12–18 are based on Example 3 about income taxes.
12 What are the income taxes on \( x = $10,000 \) and \( x = $30,000 \) and \( x = $50,000 \)?
13 What is the equation for income tax \( f(x) \) in the second bracket \( $20,350 \leq x \leq $49,300 \)? How is the number \( 11,158.50 \) connected with the other numbers in the tax instructions?
14 Write the tax function \( F(x) \) for a married couple if the IRS treats them as two single taxpayers each with taxable income \( x/2 \). (This is not done.)
15 In the 15% bracket, with 5% state tax as a deduction, the combined rate is not 20% but _______. Think about the tax on an extra $100.
16 A piecewise linear function is continuous when \( f(t) \) at the end of each interval equals \( f(t) \) at the start of the following interval. If \( f(t) = 5t \) up to \( t = 1 \) and \( v(t) = 2 \) for \( t > 1 \), define \( f \) beyond \( t = 1 \) so it is (a) continuous (b) discontinuous.
   (c) Define a tax function \( f(x) \) with rates .15 and .28 so you would lose by earning an extra dollar beyond the breakpoint.
17 The difference between a tax credit and a deduction from income is the difference between \( f(x) - c \) and \( f(x - d) \). Which is more desirable, a credit of \( c = $1000 \) or a deduction of \( d = $1000 \), and why? Sketch the tax graphs when \( f(x) = .15x \).
18 The average tax rate on the taxable income \( x \) is \( a(x) = f(x)/x \). This is the slope between \( (0, 0) \) and the point \( (x, f(x)) \).
   Draw a rough graph of \( a(x) \). The average rate \( a \) is below the marginal rate \( v \) because ________.

Problems 19–30 involve numbers \( f_0, f_1, f_2, \ldots \) and their differences \( v_j = f_j - f_{j-1} \). They give practice with subscripts 0, ..., \( j \).
19 Find the velocities \( v_1, v_2, v_3 \) and formulas for \( v_j \) and \( f_j \):
   (a) \( f = 1, 3, 5, 7 \ldots \) (b) \( f = 0, 1, 0, 1, \ldots \) (c) \( f = 0, \frac{1}{2}, \frac{1}{2}, \frac{1}{4}, \sin t \).
20 Find \( f_1, f_2, f_3 \) and a formula for \( f_j \) with \( f_0 = 0 \):
   (a) \( v = 1, 2, 4, 8, \ldots \) (b) \( v = -1, 1, -1, 1, \ldots \)
21 The areas of these nested squares are \( 1^2, 2^2, 3^2, \ldots \). What are the areas of the L-shaped bands (the differences between squares)? How does the figure show that \( 1 + 3 + 5 + 7 = 4^2 \)?
22 From the area under the staircase (by rectangles and then by triangles) show that the first \( j \) whole numbers 1 to \( j \) add up to \( j^2 + j/2 \). Find \( 1 + 2 + \ldots + 100 \).
23 If \( v = 1, 3, 5, \ldots \) then \( f_j = j^2 \). If \( v = 1, 1, 1, \ldots \) then \( f_j = _____ \). Add those to find the sum of 2, 4, 6, ..., 2\( j \). Divide by 2 to find the sum of 1, 2, 3, ..., \( j \). (Compare Problem 22.)
24 True (with reason) or false (with example).
   (a) When the \( f \)'s are increasing so are the \( v \)'s.
   (b) When the \( v \)'s are increasing so are the \( f \)'s.
   (c) When the \( f \)'s are periodic so are the \( v \)'s.
   (d) When the \( v \)'s are periodic so are the \( f \)'s.
25 If \( f(t) = t^2 \), compute \( f(99) \) and \( f(101) \). Between those times, what is the increase in \( f \) divided by the increase in \( t \)?
26 If \( f(t) = t^2 + t \), compute \( f(99) \) and \( f(101) \). Between those times, what is the increase in \( f \) divided by the increase in \( t \)?
27 If \( f_j = j^2 + j + 1 \) find a formula for \( v_j \).
28 Suppose the \( v \)'s increase by 4 at every step. Show by example and then by algebra that the "second difference" \( f_{j+1} - 2f_j + f_{j-1} \) equals 4.
29 Suppose \( f_0 = 0 \) and the \( v \)'s are \( 1, \frac{1}{2}, \frac{1}{2}, \frac{1}{4}, \ldots \). For which \( j \) does \( f_j = 3^2 \)?
30 Show that \( a_j = f_{j+1} - 2f_j + f_{j-1} \) always equals \( v_{j+1} - v_j \). If \( v \) is velocity then \( a \) stands for ________.

Problems 31–34 involve periodic \( f \)'s and \( v \)'s (like \( \sin t \) and \( \cos t \)).
31 For the discrete sine \( f_0, 0, 1, 0, -1, -1, 0 \) find the second differences \( a_j = f_{j+1} - 2f_j + f_{j-1} \) and \( a_3 = f_3 - 2f_2 + f_1 \) and \( a_5 \). Compare \( a_3 \) with \( f_3 \).
32 If the sequence \( v_1, v_2, \ldots \) has period 6 and \( w_1, w_2, \ldots \) has period 10, what is the period of \( v_1 + w_1, v_2 + w_2, \ldots \)?
33 Draw the graph of \( f(t) \) starting from \( f_0 = 0 \) when \( v = 1, -1, -1, 1, \ldots \). If \( v \) has period 4 find \( f(12), f(13) \), \( f(100, 1) \).
1. Introduction to Calculus

Problems 35–42 are about exponential v's and f's.

35 Find the v's for f = 1, 3, 9, 27. Predict v₄ and v₅. Algebra gives \(3^j - 3^{j-1} = (3 - 1)3^{j-1}\).

36 Find \(1 + 2 + 4 + \cdots + 32\) and also \(1 + \frac{1}{2} + \frac{1}{4} + \cdots + \frac{1}{32}\).

37 Estimate the slope of \(f(t) = 2^t\) at \(t = 0\). Use a calculator to compute \((\text{increase in } f)/(\text{increase in } t)\) when \(t\) is small:

\[
\frac{f(t) - f(0)}{t} = \frac{2 - 1}{1} + \frac{2^1 - 1}{.1} + \frac{2^2 - 1}{.01} + \frac{2^3 - 1}{.001}.
\]

38 Suppose \(f_0 = 1\) and \(v_j = 2f_{j-1}\). Find \(f_4\).

39 (a) From \(f = 1, \frac{1}{2}, \frac{1}{4}, \frac{1}{8}\) find \(v_1, v_2, v_3\) and predict \(v_4\).

(b) Check \(f_2 - f_0 = v_1 + v_2 + v_3\) and \(f_4 - f_2 = v_3\). (Then this is the correct \(f_j = 1 + r + \cdots + r^j = \text{sum of a geometric series}\).)

40 Suppose \(v_j = r^j\). Show that \(f_j = (r^{j+1} - 1)/(r - 1)\) starts from \(f_0 = 1\) and has \(f_j - f_{j-1} = v_j\). (This is the correct \(f_j = 1 + r + \cdots + r^j = \text{sum of a geometric series}\).)

41 From \(f_j = (-1)^j\) compute \(v_j\). What is \(v_1 + v_2 + \cdots + v_4\)?

42 Estimate the slope of \(f(t) = e^t\) at \(t = 0\). Use a calculator that knows \(e\) (or else take \(e = 2.78\)) to compute

\[
\frac{f(t) - f(0)}{t} = \frac{e - 1}{1} + \frac{e^1 - 1}{.1} + \frac{e^2 - 1}{.01} + \frac{e^3 - 1}{.001}.
\]

Problems 43–47 are about \(U(t) = \text{step from } 0 \text{ to } 1 \text{ at } t = 0\).

43 Graph the four functions \(U(t-1)\) and \(U(t-2)\) and \(U(3t)\) and \(4U(t)\). Then graph \(f(t) = 4U(3t-1) - 2\).

44 Graph the square wave \(U(t) - U(t-1)\). If this is the velocity \(v(t)\), graph the distance \(f(t)\). If this is the distance \(f(t)\), graph the velocity.

45 Two bursts of speed lead to the same distance \(f = 10\):

\[v = \ldots \text{to } t = 0.001 \quad v = \ldots \text{to } t = \ldots\]

As \(v \to \infty\) the limit of the \(f(t)\)'s is ________.

46 Draw the staircase function \(U(t) + U(t-1) + U(t-2)\). Its slope is a sum of three _____ functions.

47 Which capital letters like \(L\) are the graphs of functions when steps are allowed? The slope of \(L\) is minus a delta function. Graph the slopes of the others.

48 Write a subroutine \(\text{FINDV}\) whose input is a sequence \(f_0, f_1, \ldots, f_N\) and whose output is \(v_1, v_2, \ldots, v_N\). Include graphical output if possible. Test on \(f_j = 2j\) and \(j^2\) and \(2^j\).

49 Write a subroutine \(\text{FINDF}\) whose input is \(v_1, \ldots, v_N\) and \(f_0\), and whose output is \(f_0, f_1, \ldots, f_N\). The default value of \(f_0\) is zero. Include graphical output if possible. Test \(v_j = j\).

50 If \(\text{FINDV}\) is applied to the output of \(\text{FINDF}\), what sequence is returned? If \(\text{FINDF}\) is applied to the output of \(\text{FINDV}\), what sequence is returned? Watch \(f_0\).

51 Arrange \(2j\) and \(j^2\) and \(\sqrt{j}\) in increasing order

(a) when \(j\) is large: \(j = 9\) (b) when \(j\) is small: \(j = \frac{1}{2}\).

52 The average age of your family since 1970 is a piecewise linear function \(A(t)\). Is it continuous or does it jump? What is its slope? Graph it the best you can.

1.3 The Velocity at an Instant

We have arrived at the central problems that calculus was invented to solve. There are two questions, in opposite directions, and I hope you could see them coming.

1. If the velocity is changing, how can you compute the distance traveled?
2. If the graph of \(f(t)\) is not a straight line, what is its slope?

Find the distance from the velocity, find the distance from the function. Our goal is to do both—but not in one section. Calculus may be a good course, but it is not magic. The first step is to let the velocity change in the steadiest possible way.

Question 1 Suppose the velocity at each time \(t\) is \(v(t) = 2t\). Find \(f(t)\).

With \(v = 2t\), a physicist would say that the acceleration is constant (it equals \(2\)). The driver steps on the gas, the car accelerates, and the speedometer goes steadily up. The distance goes up too—faster and faster. If we measure \(t\) in seconds and \(v\) in feet per second, the distance \(f\) comes out in feet. After 10 seconds the speed is 20 feet per second. After 44 seconds the speed is 88 feet/second (which is 60 miles/hour). The acceleration is clear, but how far has the car gone?
**Question 2** The distance traveled by time $t$ is $f(t) = t^2$. Find the velocity $v(t)$.

The graph of $f(t) = t^2$ is on the right of Figure 1.12. It is a parabola. The curve starts at zero, when the car is new. At $t = 5$ the distance is $f = 25$. By $t = 10$, $f$ reaches 100.

Velocity is distance divided by time, but what happens when the speed is changing? Dividing $f = 100$ by $t = 10$ gives $v = 10$—the average velocity over the first ten seconds. Dividing $f = 121$ by $t = 11$ gives the average speed over 11 seconds. But how do we find the instantaneous velocity—the reading on the speedometer at the exact instant when $t = 10$?

![Graphs showing the velocity $v = 2t$, distance $f = t^2$, change in distance $(t + h)^2 - t^2$, and slope $2t + h$ approaches $v = 2t$.](image)

Fig. 1.12 The velocity $v = 2t$ is linear. The distance $f = t^2$ is quadratic.

I hope you see the problem. As the car goes faster, the graph of $t^2$ gets steeper—because more distance is covered in each second. The average velocity between $t = 10$ and $t = 11$ is a good approximation—but only an approximation—to the speed at the moment $t = 10$. Averages are easy to find:

- distance at $t = 10$ is $f(10) = 10^2 = 100$
- distance at $t = 11$ is $f(11) = 11^2 = 121$

$$\text{average velocity} = \frac{f(11) - f(10)}{11 - 10} = \frac{121 - 100}{1} = 21.$$

The car covered 21 feet in that 1 second. Its average speed was 21 feet/second. Since it was gaining speed, the velocity at the beginning of that second was below 21.

Geometrically, what is the average? It is a slope, but not the slope of the curve. **The average velocity is the slope of a straight line.** The line goes between two points on the curve in Figure 1.12. When we compute an average, we pretend the velocity is constant—so we go back to the easiest case. It only requires a division of distance by time:

$$\text{average velocity} = \frac{\text{change in } f}{\text{change in } t}. \tag{1}$$

**Calculus and the Law** You enter a highway at 1:00. If you exit 150 miles away at 3:00, your average speed is 75 miles per hour. I'm not sure if the police can give you a ticket. You could say to the judge, “When was I doing 75?” The police would have
1 Introduction to Calculus

to admit that they have no idea—but they would have a definite feeling that you must have been doing 75 sometime.†

We return to the central problem—computing \( v(10) \) at the instant \( t = 10 \). The average velocity over the next second is 21. We can also find the average over the half-second between \( t = 10.0 \) and \( t = 10.5 \). Divide the change in distance by the change in time:

\[
\frac{f(10.5) - f(10.0)}{10.5 - 10.0} = \frac{(10.5)^2 - (10.0)^2}{.5} = \frac{110.25 - 100}{.5} = 20.5.
\]

That average of 20.5 is closer to the speed at \( t = 10 \). It is still not exact.

The way to find \( v(10) \) is to keep reducing the time interval. This is the basis for Chapter 2, and the key to differential calculus. Find the slope between points that are closer and closer on the curve. The “limit” is the slope at a single point.

Algebra gives the average velocity between \( t = 10 \) and any later time \( t = 10 + h \). The distance increases from \( 10^2 \) to \( (10 + h)^2 \). The change in time is \( h \). So divide:

\[
v_{\text{average}} = \frac{(10 + h)^2 - 10^2}{h} = \frac{100 + 20h + h^2 - 100}{h} = 20 + h.
\] (2)

This formula fits our previous calculations. The interval from \( t = 10 \) to \( t = 11 \) had \( h = 1 \), and the average was \( 20 + h = 21 \). When the time step was \( h = \frac{1}{2} \), the average was \( 20 + \frac{1}{2} = 20.5 \). Over a millionth of a second the average will be 20 plus \( \frac{1}{1,000,000} \)—which is very near 20.

Conclusion: The velocity at \( t = 10 \) is \( v = 20 \). That is the slope of the curve. It agrees with the \( v \)-graph on the left side of Figure 1.12, which also has \( v(10) = 20 \).

We now show that the two graphs match at all times. If \( f(t) = t^2 \) then \( v(t) = 2t \). You are seeing the key computation of calculus, and we can put it into words before equations. Compute the distance at time \( t + h \), subtract the distance at time \( t \), and divide by \( h \). That gives the average velocity:

\[
v_{\text{ave}} = \frac{f(t + h) - f(t)}{h} = \frac{(t + h)^2 - t^2}{h} = \frac{t^2 + 2th + h^2 - t^2}{h} = 2t + h.
\] (3)

This fits the previous calculation, where \( t \) was 10. The average was \( 20 + h \). Now the average is \( 2t + h \). It depends on the time step \( h \), because the velocity is changing. But we can see what happens as \( h \) approaches zero. The average is closer and closer to the speedometer reading of \( 2t \), at the exact moment when the clock shows time \( t \):

\[\text{1E As } h \text{ approaches zero, the average velocity } 2t + h \text{ approaches } v(t) = 2t.\]

Note The computation (3) shows how calculus needs algebra. If we want the whole \( v \)-graph, we have to let time be a “variable.” It is represented by the letter \( t \). Numbers are enough at the specific time \( t = 10 \) and the specific step \( h = 1 \)—but algebra gets beyond that. The average between any \( t \) and any \( t + h \) is \( 2t + h \). Please don't hesitate to put back numbers for the letters—that checks the algebra.

†This is our first encounter with the much despised “Mean Value Theorem.” If the judge can prove the theorem, you are dead. A few \( v \)-graphs and \( f \)-graphs will confuse the situation (possibly also a delta function).
1.3 The Velocity at an Instant

There is also a step beyond algebra! Calculus requires the limit of the average. As \( h \) shrinks to zero, the points on the graph come closer. "Average over an interval" becomes "velocity at an instant." The general theory of limits is not particularly simple, but here we don't need it. (It isn't particularly hard either.) In this example the limiting value is easy to identify. The average \( 2t + h \) approaches \( 2t \), as \( h \to 0 \).

What remains to do in this section? We answered Question 2—to find velocity from distance. We have not answered Question 1. If \( v(t) = 2t \) increases linearly with time, what is the distance? This goes in the opposite direction (it is integration).

The Fundamental Theorem of Calculus says that no new work is necessary. If the slope of \( f(t) \) leads to \( v(t) \), then the area under that \( v \)-graph leads back to the \( f \)-graph. The odometer readings \( f = t^2 \) produced speedometer readings \( v = 2t \). By the Fundamental Theorem, the area under \( 2t \) should be \( t^2 \). But we have certainly not proved any fundamental theorems, so it is better to be safe—by actually computing the area.

Fortunately, it is the area of a triangle. The base of the triangle is \( t \) and the height is \( v = 2t \). The area agrees with \( f(t) \):

\[
\text{area} = \frac{1}{2} \text{(base)} \times \text{(height)} = \frac{1}{2} t \times 2t = t^2. \tag{4}
\]

EXAMPLE 1 The graphs are shifted in time. The car doesn't start until \( t = 1 \). Therefore \( v = 0 \) and \( f = 0 \) up to that time. After the car starts we have \( v = 2(t - 1) \) and \( f = (t - 1)^2 \). You see how the time delay of 1 enters the formulas. Figure 1.13 shows how it affects the graphs.

\[ \begin{align*}
\text{Fig. 1.13} & \quad \text{Delayed velocity and distance. The pairs } v = at + b \text{ and } f = \frac{1}{2} at^2 + bt. \\
\end{align*} \]

EXAMPLE 2 The acceleration changes from 2 to another constant \( a \). The velocity changes from \( v = 2t \) to \( v = at \). The acceleration is the slope of the velocity curve! The distance is also proportional to \( a \), but notice the factor \( \frac{1}{2} \):

\[
\text{acceleration } a \leftrightarrow \text{velocity } v = at \leftrightarrow \text{distance } f = \frac{1}{2} at^2.
\]

If \( a \) equals 1, then \( v = t \) and \( f = \frac{1}{2} t^2 \). That is one of the most famous pairs in calculus. If \( a \) equals the gravitational constant \( g \), then \( v = gt \) is the velocity of a falling body. The speed doesn't depend on the mass (tested by Galileo at the Leaning Tower of Pisa). Maybe he saw the distance \( f = \frac{1}{2} gt^2 \) more easily than the speed \( v = gt \). Anyway, this is the most famous pair in physics.
EXAMPLE 3  Suppose \( f(t) = 3t + t^2 \). The average velocity from \( t \) to \( t + h \) is
\[
\frac{f(t + h) - f(t)}{h} = \frac{3(t + h) + (t + h)^2 - 3t - t^2}{h}.
\]
The change in distance has an extra \( 3h \) (coming from \( 3(t + h) \) minus \( 3t \)). The velocity contains an additional \( 3 \) (coming from \( 3h \) divided by \( h \)). When \( 3t \) is added to the distance, \( 3 \) is added to the velocity. If Galileo had thrown a weight instead of dropping it, the starting velocity \( v_0 \) would have added \( v_0 t \) to the distance.

FUNCTIONS ACROSS TIME

The idea of slope is not difficult—for one straight line. Divide the change in \( f \) by the change in \( t \). In Chapter 2, divide the change in \( y \) by the change in \( x \). Experience shows that the hard part is to see what happens to the slope as the line moves.

Figure 1.14a shows the line between points \( A \) and \( B \) on the curve. This is a "secant line." Its slope is an average velocity. What calculus does is to bring that point \( B \) down the curve toward \( A \).

**Question 1**  What happens to the "change in \( f \)"—the height of \( B \) above \( A \)?
**Answer**  The change in \( f \) decreases to zero. So does the change in \( t \).

**Question 2**  As \( B \) approaches \( A \), does the slope of the line increase or decrease?
**Answer**  I am not going to answer that question. It is too important. Draw another secant line with \( B \) closer to \( A \). Compare the slopes.

This question was created by Steve Monk at the University of Washington—where 57% of the class gave the right answer. Probably 97% would have found the right slope from a formula. Figure 1.14b shows the opposite problem. We know the velocity, not the distance. But calculus answers questions about both functions.

**Question 3**  Which car is going faster at time \( t = 3/4 \)?
**Answer**  Car \( C \) has higher speed. Car \( D \) has greater acceleration.

**Question 4**  If the cars start together, is \( D \) catching up to \( C \) at the end? Between \( t = \frac{1}{2} \) and \( t = 1 \), do the cars get closer or further apart?
**Answer**  This time more than half the class got it wrong. You won't but you can see why they did. You have to look at the speed graph and imagine the distance graph. When car \( C \) is going faster, the distance between them _______.

| Fig. 1.14 | Slope of line, slope of curve. Two velocity graphs. Which is which? |
1.3 The Velocity at an Instant

To repeat: The cars start together, but they don't finish together. They reach the same speed at $t = 1$, not the same distance. Car C went faster. You really should draw their distance graphs, to see how they bend.

These problems help to emphasize one more point. Finding the speed (or slope) is entirely different from finding the distance (or area):

1. To find the slope of the $f$-graph at a particular time $t$, you don't have to know the whole history.
2. To find the area under the $v$-graph up to a particular time $t$, you do have to know the whole history.

A short record of distance is enough to recover $v(t)$.

On the graph of $f(t)$, the average velocity between $A$ and $B$ is the slope of the line through $A$ and $B$. The velocity at $A$ is found by $\frac{1}{2}$.

When the velocity is increasing, the distance is positive.

When the velocity is positive, the distance is increasing.

11 If $f(t) = 6t^2$ find the slope of the $f$-graph and also the $v$-graph. The slope of the $v$-graph is the ________.

13 Find $f(t)$ from $v(t) = 20t$ iff $f(0) = 12$. Also if $f(1) = 12$.

For the same $f(t) = t^2 + t$, find the average speed between

(a) $t = 0$ and 1 (b) $t = 0$ and $\frac{1}{2}$ (c) $t = 0$ and $h$.

5 In the answer to 3(c), find the limit as $h \to 0$. What does that limit tell us?

6 Set $h = 0$ in your answer to 4(c). Draw the graph of $f(t) = t^2 + t$ and show its slope at $t = 0$.

7 Draw the graph of $v(t) = 1 + 2t$. From geometry find the area under it from 0 to $t$. Find the slope of that area function $f(t)$.

8 Draw the graphs of $v(t) = 3 - 2t$ and the area $f(t)$.

9 True or false

(a) If the distance $f(t)$ is positive, so is $v(t)$.
(b) If the distance $f(t)$ is increasing, so is $v(t)$.
(c) If $f(t)$ is positive, $v(t)$ is increasing.
(d) If $v(t)$ is positive, $f(t)$ is increasing.

10 If $f(t) = 6t^2$ find the slope of the $f$-graph and also the $v$-graph. The slope of the $v$-graph is the ________.

11 If $f(t) = t^2$ what is the average velocity between $t = 9$ and $t = 11$? What is the average between $t - h$ and $t + h$?

12 (a) Show that for $f(t) = \frac{1}{2}at^2$ the average velocity between $t - h$ and $t + h$ is exactly the velocity at $t$.
(b) The area under $v(t) = at$ from $t - h$ to $t + h$ is exactly the base $2h$ times ________.

13 Find $f(t)$ from $v(t) = 20t$ if $f(0) = 12$. Also if $f(1) = 12$.

14 True or false, for any distance curves.

(a) The slope of the line from $A$ to $B$ is the average velocity between those points.
(b) Secant lines have smaller slopes than the curve.
(c) If $f(t)$ and $F(t)$ start together and finish together, the average velocities are equal.
(d) If $v(t)$ and $V(t)$ start together and finish together, the increases in distance are equal.

15 When you jump up and fall back your height is $y = 2t - t^2$ in the right units.
(a) Graph this parabola and its slope.
(b) Find the time in the air and maximum height.
(c) Prove: Half the time you are above $y = \frac{1}{2}$.

Basketball players "hang" in the air partly because of (c).

16 Graph $f(t) = t^2$ and $g(t) = f(t) - 2$ and $h(t) = f(2t)$, all from $t = 0$ to $t = 1$. Find the velocities.

17 (Recommended) An up and down velocity is $v(t) = 2t$ for $t \leq 3$, $v(t) = 12 - 2t$ for $t \geq 3$. Draw the piecewise parabola $f(t)$. Check that $f(6)$ = area under the graph of $v(t)$.

18 Suppose $v(t) = 1$ for $t \leq 2$ and $v(t) = 2$ for $t \geq 2$. Draw the graph of $f(t)$ out to $t = 3$.

19 Draw $f(t)$ up to $t = 4$ when $u(t)$ increases linearly from
(a) $0$ to $2$
(b) $-1$ to $1$
(c) $-2$ to $0$.

20 (Recommended) Suppose $v(t)$ is the piecewise linear sine function of Section 1.2. (In Figure 1.8 it was the distance.) Find the area under $v(t)$ between $t = 0$ and $t = 1, 2, 3, 4, 5, 6$. Plot those points $f(1), \ldots, f(6)$ and draw the complete piecewise parabola $f(t)$.

21 Draw the graph of $f(t) = \lfloor 1 - t^2 \rfloor$ for $0 \leq t \leq 2$. Find a three-part formula for $v(t)$.

22 Draw the graphs of $f(t)$ for these velocities (to $t = 2$):
(a) $v(t) = 1 - t$
(b) $v(t) = |1 - t|$
(c) $v(t) = (1 - t) + |1 - t|$.

23 When does $f(t) = t^2 - 3t$ reach $0$? Find the average velocity up to that time and the instantaneous velocity at that time.

24 If $f(t) = \frac{1}{2}at^2 + bt + c$, what is $v(t)$? What is the slope of $v(t)$? When does $f(t)$ equal $41$, if $a = b = c = \frac{1}{2}$?

25 If $f(t) = t^2$ then $v(t) = 2t$. Does the speeded-up function $f(4t)$ have velocity $v(4t)$ or $4v(t)$ or $4v(4t)$?

26 If $f(t) = t^2$ find $v(t)$ and $f(3t)$. Does the slope of $f(3t)$ equal $v(3t)$ or $3v(t)$ or $3v(3t)$?

27 For $f(t) = t^2$ find $v_{ave}(t)$ between $0$ and $t$. Graph $v_{ave}(t)$ and $v(t)$.

28 If you know the average velocity $v_{ave}(t)$, how can you find the distance $f(t)$? Start from $f(0) = 0$.

### 1.4 Circular Motion

This section introduces completely new distances and velocities—*the sines and cosines from trigonometry*. As I write that last word, I ask myself how much trigonometry it is essential to know. There will be the basic picture of a right triangle, with sides $\cos t$ and $\sin t$ and $1$. There will also be the crucial equation $(\cos t)^2 + (\sin t)^2 = 1$, which is Pythagoras' law $a^2 + b^2 = c^2$. The squares of two sides add to the square of the hypotenuse (and the 1 is really $1^2$). Nothing else is needed immediately. If you don't know trigonometry, don't stop—an important part can be learned now.

You will recognize the wavy graphs of the sine and cosine. We *intend to find the slopes of those* graphs. That can be done without using the formulas for $\sin(x + y)$ and $\cos(x + y)$—which later give the same slopes in a more algebraic way. Here it is only basic things that are needed.† And anyway, how complicated can a triangle be?

**Remark** You might think trigonometry is only for surveyors and navigators (people with triangles). Not at all! By far the biggest applications are to *rotation and vibration* and *oscillation*. It is fantastic that sines and cosines are so perfect for "repeating motion"—around a circle or up and down.

†Sines and cosines are so important that I added a review of trigonometry in Section 1.5. But the concepts in this section can be more valuable than formulas.
1.4 Circular Motion

Our underlying goal is to offer one more example in which the velocity can be computed by common sense. Calculus is mainly an extension of common sense, but here that extension is not needed. We will find the slope of the sine curve. The straight line $f = vt$ was easy and the parabola $f = \frac{1}{2}at^2$ was harder. The new example also involves realistic motion, seen every day. We start with circular motion, in which the position is given and the velocity will be found.

**A ball goes around a circle of radius one.** The center is at $x = 0$, $y = 0$ (the origin). The $x$ and $y$ coordinates satisfy $x^2 + y^2 = 1^2$, to keep the ball on the circle. We specify its position in Figure 1.16a by giving its angle with the horizontal. And we make the ball travel with constant speed, by requiring that the angle is equal to the time $t$. The ball goes counterclockwise. At time 1 it reaches the point where the angle equals 1. The angle is measured in radians rather than degrees, so a full circle is completed at $t = 2\pi$ instead of $t = 360$.

The ball starts on the $x$ axis, where the angle is zero. Now find it at time $t$:

**The ball is at the point where $x = \cos t$ and $y = \sin t$.**

This is where trigonometry is useful. The cosine oscillates between 1 and $-1$, as the ball goes from far right to far left and back again. The sine also oscillates between 1 and $-1$, starting from $\sin 0 = 0$. At time $\pi/2$ the sine (the height) increases to one. The cosine is zero and the ball reaches the top point $x = 0$, $y = 1$. At time $\pi$ the cosine is $-1$ and the sine is back to zero—the coordinates are $(1, 0)$. At $t = 2\pi$ the circle is complete (the angle is also $2\pi$), and $x = \cos 2\pi = 1$, $y = \sin 2\pi = 0$.
1 Introduction to Calculus

Important point: The distance around the circle (its circumference) is \(2\pi r = 2\pi\), because the radius is 1. The ball travels a distance \(2\pi\) in a time \(2\pi\). The speed equals 1. It remains to find the velocity, which involves not only speed but direction.

Degrees vs. radians A full circle is 360 degrees and \(2\pi\) radians. Therefore

\[
1 \text{ radian} = \frac{360}{2\pi} \text{ degrees} \approx 57.3 \text{ degrees}
\]

\[
1 \text{ degree} = \frac{2\pi}{360} \text{ radians} \approx .01745 \text{ radians}
\]

Radians were invented to avoid those numbers! The speed is exactly 1, reaching \(t\) radians at time \(t\). The speed would be .01745, if the ball only reached \(t\) degrees. The ball would complete the circle at time \(T = 360\). We cannot accept the division of the circle into 360 pieces (by whom?), which produces these numbers.

To check degree mode vs. radian mode, verify that \(\sin(\theta) \approx .017\) and \(\sin(\pi) \approx .84\).

VELOCITY OF THE BALL

At time \(t\), which direction is the ball going? Calculus watches the motion between \(t\) and \(t + h\). For a ball on a string, we don't need calculus—just let go. The direction of motion is tangent to the circle. With no force to keep it on the circle, the ball goes off on a tangent. If the ball is the moon, the force is gravity. If it is a hammer swinging around on a chain, the force is from the center. When the thrower lets go, the hammer takes off—and it is an art to pick the right moment. (I once saw a friend hit by a hammer at MIT. He survived, but the thrower quit track.) Calculus will find that same tangent direction, when the points at \(t\) and \(t + h\) come close.

The "velocity triangle" is in Figure 1.16b. It is the same as the position triangle, but rotated through 90°. The hypotenuse is tangent to the circle, in the direction the ball is moving. Its length equals 1 (the speed). The angle \(t\) still appears, but now it is the angle with the vertical. The upward component of velocity is \(\cos t\), when the upward component of position is \(\sin t\). That is our common sense calculation, based on a figure rather than a formula. The rest of this section depends on it—and we check \(v = \cos t\) at special points.

At the starting time \(t = 0\), the movement is all upward. The height is \(\sin 0 = 0\) and the upward velocity is \(\cos 0 = 1\). At time \(\pi/2\), the ball reaches the top. The height is \(\sin \pi/2 = 1\) and the upward velocity is \(\cos \pi/2 = 0\). At that instant the ball is not moving up or down.

The horizontal velocity contains a minus sign. At first the ball travels to the left. The value of \(x\) is \(\cos t\), but the speed in the \(x\) direction is \(-\sin t\). Half of trigonometry is in that figure (the good half), and you see how \(\sin^2 t + \cos^2 t = 1\) is so basic. That equation applies to position and velocity, at every time.

Application of plane geometry: The right triangles in Figure 1.16 are the same size and shape. They look congruent and they are—the angle \(t\) above the ball equals the angle \(t\) at the center. That is because the three angles at the ball add to 180°.

OSCILLATION: UP AND DOWN MOTION

We now use circular motion to study straight-line motion. That line will be the \(y\) axis. Instead of a ball going around a circle, a mass will move up and down. It oscillates between \(y = 1\) and \(y = -1\). The mass is the "shadow of the ball," as we explain in a moment.
There is a jumpy oscillation that we do not want, with \( v = 1 \) and \( v = -1 \). That "bang-bang" velocity is like a billiard ball, bouncing between two walls without slowing down. If the distance between the walls is 2, then at \( t = 4 \) the ball is back to the start. The distance graph is a zigzag (or sawtooth) from Section 1.2.

We prefer a smoother motion. Instead of velocities that jump between +1 and −1, a real oscillation slows down to zero and gradually builds up speed again. The mass is on a spring, which pulls it back. The velocity drops to zero as the spring is fully stretched. Then \( v \) is negative, as the mass goes the same distance in the opposite direction. Simple harmonic motion is the most important back and forth motion, while \( f = vt \) and \( f = \frac{1}{2}at^2 \) are the most important one-way motions.

![Circular motion of the ball and harmonic motion of the mass](image)

How do we describe this oscillation? The best way is to match it with the ball on the circle. The height of the ball will be the height of the mass. The "shadow of the ball" goes up and down, level with the ball. As the ball passes the top of the circle, the mass stops at the top and starts down. As the ball goes around the bottom, the mass stops and turns back up the \( y \) axis. Halfway up (or down), the speed is 1. Figure 1.17a shows the mass at a typical time \( t \). The height is \( y = f(t) = \sin t \), level with the ball. This height oscillates between \( f = 1 \) and \( f = -1 \). But the mass does not move with constant speed. The speed of the mass is changing although the speed of the ball is always 1. The time for a full cycle is still \( 2\pi \), but within that cycle the mass speeds up and slows down. The problem is to find the changing velocity \( v \). Since the distance is \( f = \sin t \), the velocity will be the slope of the sine curve.

**THE SLOPE OF THE SINE CURVE**

At the top and bottom (\( t = \pi/2 \) and \( t = 3\pi/2 \)) the ball changes direction and \( v = 0 \). The slope at the top and bottom of the sine curve is zero.† At time zero, when the ball is going straight up, the slope of the sine curve is \( v = 1 \). At \( t = \pi \), when the ball and mass and \( f \)-graph are going down, the velocity is \( v = -1 \). The mass goes fastest at the center. The mass goes slowest (in fact it stops) when the height reaches a maximum or minimum. The velocity triangle yields \( v \) at every time \( t \).

To find the upward velocity of the mass, look at the upward velocity of the ball. Those velocities are the same! The mass and ball stay level, and we know \( v \) from circular motion: The upward velocity is \( v = \cos t \).

†That looks easy but you will see later that it is extremely important. At a maximum or minimum the slope is zero. The curve levels off.
Figure 1.18 shows the result we want. On the right, $f = \sin t$ gives the height. On the left is the velocity $v = \cos t$. That velocity is the slope of the $f$-curve. The height and velocity (red lines) are oscillating together, but they are out of phase—just as the position triangle and velocity triangle were at right angles. This is absolutely fantastic, that in calculus the two most famous functions of trigonometry form a pair: The slope of the sine curve is given by the cosine curve.

When the distance is $f(t) = \sin t$, the velocity is $v(t) = \cos t$.

Admission of guilt: The slope of $\sin t$ was not computed in the standard way. Previously we compared $(t + h)^2$ with $t^2$, and divided that distance by $h$. This average velocity approached the slope $2t$ as $h$ became small. For $\sin t$ we could have done the same:

$$\text{average velocity} = \frac{\text{change in } \sin t}{\text{change in } t} = \frac{\sin (t + h) - \sin t}{h}.$$  \hspace{1cm} (1)

This is where we need the formula for $\sin (t + h)$, coming soon. Somehow the ratio in (1) should approach $\cos t$ as $h \to 0$. (It does.) The sine and cosine fit the same pattern as $t^2$ and $2t$—our shortcut was to watch the shadow of motion around a circle.

![Graphs showing $v = \cos t$ and $f = \sin t$](image)

**Fig. 1.18** $v = \cos t$ when $f = \sin t$ (red); $v = -\sin t$ when $f = \cos t$ (black).

**Question 1** What if the ball goes twice as fast, to reach angle $2t$ at time $t$?

**Answer** The speed is now 2. The time for a full circle is only $\pi$. The ball's position is $x = \cos 2t$ and $y = \sin 2t$. The velocity is still tangent to the circle—but the tangent is at angle $2t$ where the ball is. Therefore $\cos 2t$ enters the upward velocity and $-\sin 2t$ enters the horizontal velocity. The difference is that the velocity triangle is twice as big. The upward velocity is not $\cos 2t$ but $2 \cos 2t$. The horizontal velocity is $-2 \sin 2t$. Notice these 2's!

**Question 2** What is the area under the cosine curve from $t = 0$ to $t = \pi/2$?

You can answer that, if you accept the Fundamental Theorem of Calculus—computing areas is the opposite of computing slopes. The slope of $\sin t$ is $\cos t$, so the area under $\cos t$ is the increase in $\sin t$. No reason to believe that yet, but we use it anyway.

From $\sin 0 = 0$ to $\sin \pi/2 = 1$, the increase is 1. Please realize the power of calculus. No other method could compute the area under a cosine curve so fast.
1.4 Circular Motion

THE SLOPE OF THE COSINE CURVE

I cannot resist uncovering another distance and velocity (another $f$-$v$ pair) with no extra work. This time $f$ is the cosine. The time clock starts at the top of the circle. The old time $t = \pi/2$ is now $t = 0$. The dotted lines in Figure 1.18 show the new start. But the shadow has exactly the same motion—the ball keeps going around the circle, and the mass follows it up and down. The $f$-graph and $v$-graph are still correct, both with a time shift of $\pi/2$.

The new $f$-graph is the cosecant. The new $v$-graph is minus the sine. The slope of the cosine curve follows the negative of the sine curve. That is another famous pair, twins of the first:

When the distance is $f(t) = \cos t$, the velocity is $v(t) = -\sin t$.

You could see that coming, by watching the ball go left and right (instead of up and down). Its distance across is $f = \cos t$. Its velocity across is $v = -\sin t$. That twin pair completes the calculus in Chapter 1 (trigonometry to come). We review the ideas:

$v$ is the velocity
the slope of the distance curve
the limit of average velocity over a short time
the derivative of $f$.

$f$ is the distance
the area under the velocity curve
the limit of total distance over many short times
the integral of $v$.

Differential calculus: Compute $v$ from $f$. Integral calculus: Compute $f$ from $v$.

With constant velocity, $f$ equals $vt$. With constant acceleration, $v = at$ and $f = \frac{1}{2}at^2$. In harmonic motion, $v = \cos t$ and $f = \sin t$. One part of our goal is to extend that list—for which we need the tools of calculus. Another and more important part is to put these ideas to use.

Before the chapter ends, may I add a note about the book and the course? The book is more personal than usual, and I hope readers will approve. What I write is very close to what I would say, if you were in this room. The sentences are spoken before they are written.† Calculus is alive and moving forward—it needs to be taught that way.

One new part of the subject has come with the computer. It works with a finite step $h$, not an "infinitesimal" limit. What it can do, it does quickly—even if it cannot find exact slopes or areas. The result is an overwhelming growth in the range of problems that can be solved. We landed on the moon because $f$ and $v$ were so accurate. (The moon's orbit has sines and cosines, the spacecraft starts with $v = at$ and $f = \frac{1}{2}at^2$. Only the computer can account for the atmosphere and the sun's gravity and the changing mass of the spacecraft.) Modern mathematics is a combination of exact formulas and approximate computations. Neither part can be ignored, and I hope you will see numerically what we derive algebraically. The exercises are to help you master both parts.

†On television you know immediately when the words are live. The same with writing.
1 Introduction to Calculus

The course has made a quick start—not with an abstract discussion of sets or functions or limits, but with the concrete questions that led to those ideas. You have seen a distance function \( f \) and a limit \( v \) of average velocities. We will meet more functions and more limits (and their definitions!) but it is crucial to study important examples early. There is a lot to do, but the course has definitely begun.

1.4 Exercises

Read-through questions

A ball at angle \( t \) on the unit circle has coordinates \( x = \ldots \) and \( y = \ldots \). It completes a full circle at \( t = \ldots \). Its speed is \( \ldots \). Its velocity points in the direction of \( \ldots \), which is \( \ldots \) to the radius coming out from the center. The upward velocity is \( \ldots \) and the horizontal velocity is \( \ldots \).

A mass going up and down level with the ball has height \( f(t) = \ldots \). This is called simple \( \ldots \) motion. The velocity is \( v(t) = \ldots \). When \( t = \pi/2 \) the height \( f = \ldots \) and the velocity is \( v = \ldots \). If a speeded-up mass reaches \( f = \sin 2t \) at time \( t \), its velocity is \( v = \ldots \). A shadow traveling \( \ldots \) under the ball has \( f = \cos t \) and \( v = \ldots \). When \( f \) is distance \( \ldots \) area = integral, \( v \) is \( \ldots \).

1. For a ball going around a unit circle with speed 1,
   (a) how long does it take for 5 revolutions?
   (b) at time \( t = 3\pi/2 \) where is the ball?
   (c) at \( t = 22 \) where is the ball (approximately)?

2. For the same motion find the exact \( x \) and \( y \) coordinates at \( t = 2\pi/3 \). At what time would the ball hit the \( x \) axis, if it goes off on the tangent at \( t = 2\pi/3 \)?

3. A ball goes around a circle of radius 4. At time \( t \) (when it reaches angle \( t \)) find
   (a) its \( x \) and \( y \) coordinates
   (b) the speed and the distance traveled
   (c) the vertical and horizontal velocity.

4. On a circle of radius \( R \) find the \( x \) and \( y \) coordinates at time \( t \) (and angle \( t \)). Draw the velocity triangle and find the \( x \) and \( y \) velocities.

5. A ball travels around a unit circle (radius 1) with speed 3, starting from angle zero. At time \( t \),
   (a) what angle does it reach?
   (b) what are its \( x \) and \( y \) coordinates?
   (c) what are its \( x \) and \( y \) velocities? This part is harder.

6. If another ball stays \( \pi/2 \) radians ahead of the ball with speed 3, find its angle, its \( x \) and \( y \) coordinates, and its vertical velocity at time \( t \).

7. A mass moves on the \( x \) axis under or over the original ball (on the unit circle with speed 1). What is the position \( x = f(t) \)? Find \( x \) and \( v \) at \( t = \pi/4 \). Plot \( x \) and \( v \) up to \( t = \pi \).

8. Does the new mass (under or over the ball) meet the old mass (level with the ball)? What is the distance between the masses at time \( t \)?

9. Draw graphs of \( f(t) = \cos 3t \) and \( \cos 2\pi t \) and \( 2\pi \cos t \), marking the time axes. How long until each \( f \) repeats?

10. Draw graphs of \( f = \sin (t + \pi) \) and \( v = \cos (t + \pi) \). This oscillation stays level with what ball?

11. Draw graphs of \( f = \sin (\pi/2 - t) \) and \( v = -\cos (\pi/2 - t) \). This oscillation stays level with a ball going which way starting where?

12. Draw a graph of \( f(t) = \sin t + \cos t \). Estimate its greatest height (maximum \( f \)) and the time it reaches that height. By computing \( f^2 \) check your estimate.

13. How fast should you run across the circle to meet the ball again? It travels at speed 1.

14. A mass falls from the top of the unit circle when the ball of speed 1 passes by. What acceleration \( a \) is necessary to meet the ball at the bottom?

Find the area under \( v = \cos t \) from the change in \( f = \sin t \):

15. from \( t = 0 \) to \( t = \pi \)
16. from \( t = 0 \) to \( t = \pi/6 \)
17. from \( t = 0 \) to \( t = 2\pi \)
18. from \( t = \pi/2 \) to \( t = 3\pi/2 \).

19. The distance curve \( f = \sin 4t \) yields the velocity curve \( v = 4 \cos 4t \). Explain both 4's.

20. The distance curve \( f = 2 \cos 3t \) yields the velocity curve \( v = -6 \sin 3t \). Explain the -6.

21. The velocity curve \( v = \cos 4t \) yields the distance curve \( f = \frac{1}{4} \sin 4t \). Explain the \( \frac{1}{4} \).

22. The velocity \( v = 5 \sin 5t \) yields what distance?
23 Find the slope of the sine curve at \( t = \pi/3 \) from \( v = \cos t \). Then find an average slope by dividing \( \sin \pi/2 - \sin \pi/3 \) by the time difference \( \pi/2 - \pi/3 \).

24 The slope of \( f = \sin t \) at \( t = 0 \) is \( \cos 0 = 1 \). Compute average slopes \( (\sin t)/t \) for \( t = 1, .1, .01, .001 \).

The ball at \( x = \cos t, y = \sin t \) circles (1) counterclockwise (2) with radius 1 (3) starting from \( x = 1, y = 0 \) (4) at speed 1. Find (1)(2)(3)(4) for the motions 25–30.

25 \( x = \cos 3t, y = -\sin 3t \)
26 \( x = 3 \cos 4t, y = 3 \sin 4t \)
27 \( x = 5 \sin 2t, y = 5 \cos 2t \)
28 \( x = 1 + \cos t, y = \sin t \)
29 \( x = \cos (t + 1), y = \sin (t + 1) \)
30 \( x = \cos (-t), y = \sin (-t) \)

The oscillation \( x = 0, y = \sin t \) goes (1) up and down (2) between \(-1\) and \(1\) (3) starting from \( x = 0, y = 0 \) (4) at velocity \( v = \cos t \). Find (1)(2)(3)(4) for the oscillations 31–36.

31 \( x = \cos t, y = 0 \)
32 \( x = 0, y = \sin 5t \)
33 \( x = 0, y = 2 \sin (t + \theta) \)
34 \( x = \cos t, y = \cos t \)
35 \( x = 0, y = -2 \cos \frac{t}{2} \)
36 \( x = \cos^2 t, y = \sin^2 t \)

37 If the ball on the unit circle reaches \( t \) degrees at time \( t \), find its position and speed and upward velocity.

38 Choose the number \( k \) so that \( x = \cos kt, y = \sin kt \) completes a rotation at \( t = 1 \). Find the speed and upward velocity.

39 If a pitcher doesn't pause before starting to throw, a balk is called. The American League decided mathematically that there is always a stop between backward and forward motion, even if the time is too short to see it. (Therefore no balk.) Is that true?

---

### 1.5 A Review of Trigonometry

Trigonometry begins with a right triangle. The size of the triangle is not as important as the angles. We focus on one particular angle—call it \( \theta \)—and on the ratios between the three sides \( x, y, r \). The ratios don't change if the triangle is scaled to another size. Three sides give six ratios, which are the basic functions of trigonometry:

\[
\begin{align*}
\cos \theta &= \frac{x}{r} = \frac{\text{near side}}{\text{hypotenuse}} \\
\sin \theta &= \frac{y}{r} = \frac{\text{opposite side}}{\text{hypotenuse}} \\
\tan \theta &= \frac{y}{x} = \frac{\text{opposite side}}{\text{near side}} \\
\sec \theta &= \frac{r}{x} = \frac{1}{\cos \theta} \\
\csc \theta &= \frac{r}{y} = \frac{1}{\sin \theta} \\
\cot \theta &= \frac{x}{y} = \frac{1}{\tan \theta}
\end{align*}
\]

Of course those six ratios are not independent. The three on the right come directly from the three on the left. And the tangent is the sine divided by the cosine:

\[
\tan \theta = \frac{\sin \theta}{\cos \theta} = \frac{y/r}{x/r} = \frac{y}{x}
\]

Note that "tangent of an angle" and "tangent to a circle" and "tangent line to a graph" are different uses of the same word. As the cosine of \( \theta \) goes to zero, the tangent of \( \theta \) goes to infinity. The side \( x \) becomes zero, \( \theta \) approaches 90°, and the triangle is infinitely steep. The sine of 90° is \( y/r = 1 \).

Triangles have a serious limitation. They are excellent for angles up to 90°, and they are OK up to 180°, but after that they fail. We cannot put a 240° angle into a triangle. Therefore we change now to a circle.
1 Introduction to Calculus

Fig. 1.20 Trigonometry on a circle. Compare \(2 \sin \theta\) with \(2 \theta\) and \(\tan \theta\) (periods \(2\pi, \pi, \pi\)).

Angles are measured from the positive \(x\) axis (counterclockwise). Thus \(90^\circ\) is straight up, \(180^\circ\) is to the left, and \(360^\circ\) is in the same direction as \(0^\circ\). (Then \(450^\circ\) is the same as \(90^\circ\).) Each angle yields a point on the circle of radius \(r\). The coordinates \(x\) and \(y\) of that point can be negative (but never \(r\)). As the point goes around the circle, the six ratios \(\cos \theta, \sin \theta, \tan \theta, \ldots\) trace out six graphs. The cosine waveform is the same as the sine waveform—just shifted by \(90^\circ\).

One more change comes with the move to a circle. Degrees are out. Radians are in. The distance around the whole circle is \(2\pi r\). The distance around to other points is \(\theta r\). We measure the angle by that multiple \(\theta\). For a half-circle the distance is \(\pi r\), so the angle is \(\pi\) radians—which is \(180^\circ\). A quarter-circle is \(\pi/2\) radians or \(90^\circ\).

The distance around to angle \(\theta\) is \(r\) times \(\theta\).

When \(r = 1\) this is the ultimate in simplicity: The distance is \(\theta\). A \(45^\circ\) angle is \(\frac{\pi}{4}\) of a circle and \(2\pi/8\) radians—and the length of the circular arc is \(2\pi/8\). Similarly for \(1^\circ\):

\[
360^\circ = 2\pi \text{ radians} \quad 1^\circ = \frac{2\pi}{360} \text{ radians} \quad 1 \text{ radian} = \frac{360}{2\pi} \text{ degrees}.
\]

An angle going clockwise is negative. The angle \(-\pi/3\) is \(-60^\circ\) and takes us \(\frac{1}{2}\) of the wrong way around the circle. What is the effect on the six functions?

Certainly the radius \(r\) is not changed when we go to \(-\theta\). Also \(x\) is not changed (see Figure 1.20a). But \(y\) reverses sign, because \(-\theta\) is below the axis when \(+\theta\) is above. This change in \(y\) affects \(y/r\) and \(y/x\) but not \(x/r\):

\[
\cos(-\theta) = \cos \theta \quad \sin(-\theta) = -\sin \theta \quad \tan(-\theta) = -\tan \theta.
\]

The cosine is even (no change). The sine and tangent are odd (change sign).

The same point is \(\frac{1}{2}\) of the right way around. Therefore \(\frac{\pi}{2}\) of \(2\pi\) radians (or \(300^\circ\)) gives the same direction as \(-\pi/3\) radians or \(-60^\circ\). A difference of \(2\pi\) makes no difference to \(x, y, r\). Thus \(\sin \theta\) and \(\cos \theta\) and the other four functions have period \(2\pi\).

We can go five times or a hundred times around the circle, adding \(10\pi\) or \(200\pi\) to the angle, and the six functions repeat themselves.

**EXAMPLE** Evaluate the six trigonometric functions at \(\theta = 2\pi/3\) (or \(\theta = -4\pi/3\)).

This angle is shown in Figure 1.20a (where \(r = 1\)). The ratios are

\[
\cos \theta = \frac{x}{r} = -1/2 \quad \sin \theta = \frac{y}{r} = \frac{\sqrt{3}}{2} \quad \tan \theta = \frac{y}{x} = -\sqrt{3}
\]

\[
\sec \theta = -2 \quad \csc \theta = 2/\sqrt{3} \quad \cot \theta = -1/\sqrt{3}
\]

Those numbers illustrate basic facts about the sizes of four functions:

\[
|\cos \theta| \leq 1 \quad |\sin \theta| \leq 1 \quad |\sec \theta| \geq 1 \quad |\csc \theta| \geq 1.
\]

The tangent and cotangent can fall anywhere, as long as \(\cot \theta = 1/\tan \theta\).
The numbers reveal more. The tangent $-\sqrt{3}$ is the ratio of sine to cosine. The secant $-2$ is $1/\cos \theta$. Their squares are $3$ and $4$ (differing by $1$). That may not seem remarkable, but it is. There are three relationships in the squares of those six numbers, and they are the key identities of trigonometry:

$$\cos^2 \theta + \sin^2 \theta = 1 \quad 1 + \tan^2 \theta = \sec^2 \theta \quad \cot^2 \theta + 1 = \csc^2 \theta$$

**Everything flows from the Pythagoras formula** $x^2 + y^2 = r^2$. Dividing by $r^2$ gives $(x/r)^2 + (y/r)^2 = 1$. That is $\cos^2 \theta + \sin^2 \theta = 1$. Dividing by $x^2$ gives the second identity, which is $1 + (y/x)^2 = (r/x)^2$. Dividing by $y^2$ gives the third. All three will be needed throughout the book—and the first one has to be unforgettable.

**DISTANCES AND ADDITION FORMULAS**

To compute the distance between points we stay with Pythagoras. The points are in Figure 1.21a. They are known by their $x$ and $y$ coordinates, and $d$ is the distance between them. The third point completes a right triangle.

For the $x$ distance along the bottom we don’t need help. It is $x_2 - x_1$ (or $|x_2 - x_1|$ since distances can’t be negative). The distance up the side is $|y_2 - y_1|$. Pythagoras immediately gives the distance $d$:

$$\text{distance between points} = d = \sqrt{(x_2 - x_1)^2 + (y_2 - y_1)^2}. \quad (1)$$

![Fig. 1.21 (a) Distance between points and equal distances in two circles.](image)

By applying this distance formula in two identical circles, we discover the cosine of $s - t$. (Subtracting angles is important.) In Figure 1.21b, the distance squared is

$$d^2 = (\text{change in } x)^2 + (\text{change in } y)^2$$

$$= (\cos s - \cos t)^2 + (\sin s - \sin t)^2. \quad (2)$$

Figure 1.21c shows the same circle and triangle (but rotated). The same distance squared is

$$d^2 = (\cos(s - t) - 1)^2 + (\sin(s - t))^2. \quad (3)$$

Now multiply out the squares in equations (2) and (3). Whenever $(\cos s)^2 + (\sin s)^2$ appears, replace it by $1$. The distances are the same, so (2) = (3):

$$1 + 1 - 2 \cos s \cos t - 2 \sin s \sin t$$

$$1 + 1 - 2 \cos(s - t)$$
1 Introduction to Calculus

After canceling $1 + 1$ and then $-2$, we have the "addition formula" for $\cos(s - t)$:

\[ \cos s \cos t + \sin s \sin t. \]  

(4)

The cosine of $s + t$ equals $\cos s \cos t - \sin s \sin t$.  

(5)

The easiest is $t = 0$. Then $\cos t = 1$ and $\sin t = 0$. The equations reduce to $\cos s = \cos s$.

To go from (4) to (5) in all cases, replace $t$ by $-t$. No change in $\cos t$, but a "minus" appears with the sine. In the special case $s = t$, we have $\cos(t + t) = (\cos t)(\cos t) - (\sin t)(\sin t)$. This is a much-used formula for $\cos 2t$:

\[ \cos 2t = \cos^2 t - \sin^2 t = 2 \cos^2 t - 1 = 1 - 2 \sin^2 t. \]  

(6)

I am constantly using $\cos^2 t + \sin^2 t = 1$, to switch between sines and cosines.

We also need addition formulas and double-angle formulas for the sine of $s - t$ and $s + t$ and $2t$. For that we connect sine to cosine, rather than $(\sin t)^2$ to $(\cos t)^2$.

The connection goes back to the ratio $y/r$ in our original triangle. This is the sine of the angle $\theta$ and also the cosine of the complementary angle $\pi/2 - \theta$:

\[ \sin \theta = \cos(\pi/2 - \theta) \quad \text{and} \quad \cos \theta = \sin(\pi/2 - \theta). \]  

(7)

The complementary angle is $\pi/2 - \theta$ because the two angles add to $\pi/2$ (a right angle). By making this connection in Problem 19, formulas (4–5–6) move from cosines to sines:

\[ \sin(s - t) = \sin s \cos t - \cos s \sin t \]  

(8)

\[ \sin(s + t) = \sin s \cos t + \cos s \sin t \]  

(9)

\[ \sin 2t = \sin(t + t) = 2 \sin t \cos t \]  

(10)

I want to stop with these ten formulas, even if more are possible. Trigonometry is full of identities that connect its six functions—basically because all those functions come from a single right triangle. The $x, y, r$ ratios and the equation $x^2 + y^2 = r^2$ can be rewritten in many ways. But you have now seen the formulas that are needed by calculus.\textsuperscript{†} They give derivatives in Chapter 2 and integrals in Chapter 5. And it is typical of our subject to add something of its own—a limit in which an angle approaches zero. The essence of calculus is in that limit.

\textbf{Review of the ten formulas}  

Figure 1.22 shows $d^2 = (0 - \frac{1}{2})^2 + (1 - \sqrt{3}/2)^2$.

\[ \cos \frac{\pi}{6} = \cos \frac{\pi}{2} \cos \frac{\pi}{3} + \sin \frac{\pi}{2} \sin \frac{\pi}{3} \quad (s - t) \]  

\[ \sin \frac{\pi}{6} = \sin \frac{\pi}{2} \cos \frac{\pi}{3} - \cos \frac{\pi}{2} \sin \frac{\pi}{3} \]

\[ \cos \frac{5\pi}{6} = \cos \frac{\pi}{2} \cos \frac{\pi}{3} - \sin \frac{\pi}{2} \sin \frac{\pi}{3} \quad (s + t) \]  

\[ \sin \frac{5\pi}{6} = \sin \frac{\pi}{2} \cos \frac{\pi}{3} + \cos \frac{\pi}{2} \sin \frac{\pi}{3} \]

\[ \cos \frac{2\pi}{3} = \cos^2 \frac{\pi}{3} - \sin^2 \frac{\pi}{3} \quad (2t) \]  

\[ \sin \frac{2\pi}{3} = 2 \sin \frac{\pi}{3} \cos \frac{\pi}{3} \]

\[ \cos \frac{\pi}{6} = \sin \frac{\pi}{3} = \sqrt{3}/2 \]  

\[ \left( \frac{\pi}{2} - \theta \right) \]  

\[ \sin \frac{\pi}{6} = \cos \frac{\pi}{3} = 1/2 \]

\[ \]  

\textsuperscript{†}Calculus turns (6) around to $\cos^2 t = \frac{1}{2}(1 + \cos 2t)$ and $\sin^2 t = \frac{1}{2}(1 - \cos 2t)$.
Read-through questions

Starting with a triangle, the six basic functions are the **a** of the sides. Two ratios (the cosine x/r and the **b**) are below 1. Two ratios (the secant r/x and the **c**) are above 1. Two ratios (the **d** and the **e**) can take any value. The six functions are defined for all angles **θ**, by changing from a triangle to a **f**.

The angle **θ** is measured in **g**. A full circle is **8 = i**, when the distance around is 2πr. The distance to angle **θ** is **L**. All six functions have period **k**. Going clockwise changes the sign of **θ** and **L** and **m**. Since cos (-θ) = cos **θ**, the cosine is **n**.

Coming from \(x^2 + y^2 = r^2\) are the three identities \(\sin^2 \theta + \cos^2 \theta = 1\) and **o** and **p**. (Divide by \(r^2\) and **q** and **r**.) The distance from (2, 5) to (3, 4) is **s**. The distance from (1, 0) to (cos (s-t), sin (s-t)) leads to the addition formula cos (s-t) = **t**. Changing the sign of t gives cos (s+t) = **u**. Choosing s=t gives cos 2t = **v** or **w**. Therefore \(\frac{1}{2}(1 + \cos 2t) = **x**\), a formula needed in calculus.

1 In a 60-60-60 triangle show why sin 30° = \(\frac{1}{2}\).

2 Convert \(π, 3\pi, -\pi/4\) to degrees and 60°, 90°, 270° to radians. What angles between 0 and 2π correspond to \(90° = \frac{π}{2}\) and \(θ = -1^o\)?

3 Draw graphs of tan θ and cot θ from 0 to 2π. What is their (shortest) period?

4 Show that cos 2θ and cos³ θ have period π and draw them on the same graph.

5 At \(θ = 3\pi/2\) compute the six basic functions and check cos² θ + sin² θ, sec² θ - tan² θ, csc² θ - cot² θ.

6 Prepare a table showing the values of the six basic functions at \(θ = 0, \pi/4, \pi/3, \pi/2, \pi\).

7 The area of a circle is \(πr^2\). What is the area of the sector that has angle θ? It is a fraction _____ of the whole area.

8 Find the distance from (1, 0) to (0, 1) along (a) a straight line (b) a quarter-circle (c) a semicircle centered at \(\left(\frac{1}{2}, \frac{1}{2}\right)\).

9 Find the distance d from (1, 0) to \(\left(\frac{1}{2}, \sqrt{3}/2\right)\) and show on a circle why 6d is less than 2π.

10 In Figure 1.22 compute \(d^2\) and (with calculator) 12d. Why is 12d close to and below 2π?

11 Decide whether these equations are true or false:

   (a) \(\frac{\sin \theta}{1 - \cos \theta} = \frac{1 + \cos \theta}{\sin \theta}\)

   (b) \(\sec \theta + \csc \theta = \frac{\tan \theta + \cot \theta}{\sin \theta + \cos \theta}\)

   (c) \(\cos \theta - \sec \theta = \sin \theta \tan \theta\)

   (d) \(\sin (2\pi - \theta) = \sin \theta\)

12 Simplify sin (π - θ), cos (π - θ), sin (π/2 + θ), cos (π/2 + θ).

13 From the formula for cos (2t + t) find cos 3t in terms of cos t.

14 From the formula for sin (2t + t) find sin 3t in terms of sin t.

15 By averaging cos (s-t) and cos (s+t) in (4-5) find a formula for cos s cos t. Find a similar formula for sin s sin t.

16 Show that \((\cos t + i \sin t)^2 = \cos 2t + i \sin 2t\), if \(i^2 = -1\).

17 Draw cos θ and sec θ on the same graph. Find all points where cos θ = sec θ.

18 Find all angles s and t between 0 and 2π where sin (s+t) = sin s + sin t.

19 Complementary angles have sin θ = cos (π/2 - θ). Write sin (s+t) as cos (π/2 - s - t) and apply formula (4) with \(π/2 - t\) instead of s. In this way derive the addition formula (9).

20 If formula (9) is true, how do you prove (8)?

21 Check the addition formulas (4-5) and (8-9) for \(s = t = \pi/4\).

22 Use (5) and (9) to find a formula for tan (s + t).
In 23–28 find every $\theta$ that satisfies the equation.

23 $\sin \theta = -1$
24 $\sec \theta = -2$
25 $\sin \theta = \cos \theta$
26 $\sin \theta = \theta$
27 $\sec^2 \theta + \csc^2 \theta = 1$
28 $\tan \theta = 0$

29 Rewrite $\cos \theta + \sin \theta = \sqrt{2} \sin(\theta + \phi)$ by choosing the correct "phase angle" $\phi$. (Make the equation correct at $\theta = 0$. Square both sides to check.)

30 Match $a \sin x + b \cos x$ with $A \sin(x + \phi)$. From equation (9) show that $a = A \cos \phi$ and $b = A \sin \phi$. Square and add to find $A = \ldots$. Divide to find $\tan \phi = b/a$.

31 Draw the base of a triangle from the origin $O = (0, 0)$ to $P = (a, 0)$. The third corner is at $Q = (b \cos \theta, b \sin \theta)$. What are the side lengths $OP$ and $OQ$? From the distance formula (1) show that the side $PQ$ has length

$$d^2 = a^2 + b^2 - 2ab \cos \theta \quad \text{(law of cosines).}$$

32 Extend the same triangle to a parallelogram with its fourth corner at $R = (a + b \cos \theta, b \sin \theta)$. Find the length squared of the other diagonal $OR$.

33 Draw graphs for equations 33–36, and mark three points.

33 $y = \sin 2x$
34 $y = 2 \sin \pi x$
35 $y = \frac{1}{2} \cos 2\pi x$
36 $y = \sin x + \cos x$

37 Which of the six trigonometric functions are infinite at what angles?

38 Draw rough graphs or computer graphs of $t \sin t$ and $\sin 4t \sin t$ from 0 to $2\pi$.

---

### 1.6 A Thousand Points of Light

The graphs on the back cover of the book show $y = \sin \pi$. This is very different from $y = \sin x$. The graph of $\sin x$ is one continuous curve. By the time it reaches $x = 10,000$, the curve has gone up and down $10,000/2\pi$ times. Those 1591 oscillations would be so crowded that you couldn't see anything. The graph of $\sin \pi$ has picked 10,000 points from the curve—and for some reason those points seem to lie on more than 40 separate sine curves.

The second graph shows the first 1000 points. They don't seem to lie on sine curves. Most people see hexagons. But they are the same thousand points! It is hard to believe that the graphs are the same, but I have learned what to do. **Tilt the second graph and look from the side at a narrow angle.** Now the first graph appears. You see "diamonds." The narrow angle compresses the $x$ axis—back to the scale of the first graph.

The effect of scale is something we don't think of. We understand it for maps. Computers can zoom in or zoom out—those are changes of scale. What our eyes see
depends on what is "close." We think we see sine curves in the 10,000 point graph, and they raise several questions:

1. Which points are near (0, 0)?
2. How many sine curves are there?
3. Where does the middle curve, going upward from (0, 0), come back to zero?

A point near (0, 0) really means that \( \sin n \) is close to zero. That is certainly not true of \( \sin 1 \) (1 is one radian!). In fact \( \sin 1 \) is up the axis at .84, at the start of the seventh sine curve. Similarly \( \sin 2 \) is .91 and \( \sin 3 \) is .14. (The numbers 3 and .14 make us think of \( \pi \). The sine of 3 equals the sine of \( \pi - 3 \). Then \( \sin .14 \) is near .14.) Similarly \( \sin 4, \sin 5, \ldots, \sin 21 \) are not especially close to zero.

**The first point to come close is** \( \sin 22 \). This is because \( 22/7 \) is near \( \pi \). Then 22 is close to \( 7\pi \), whose sine is zero:

\[
\sin 22 = \sin (7\pi - 22) \approx \sin (-.01) \approx -.01.
\]

That is the first point to the right of (0, 0) and slightly below. You can see it on graph 1, and more clearly on graph 2. It begins a curve downward.

The next point to come close is \( \sin 44 \). This is because 44 is just past \( 14\pi \). 44 is close to \( 14\pi + .02 \) so \( \sin 44 \approx \sin .02 \approx .02 \).

**This point** \( (44, \sin 44) \) **starts the middle sine curve**. Next is \( (88, \sin 88) \).

Now we know something. **There are 44 curves**. They begin near the heights \( \sin 0, \sin 1, \ldots, \sin 43 \). Of these 44 curves, 22 start upward and 22 start downward. I was confused at first, because I could only find 42 curves. The reason is that \( \sin 11 = -0.99999 \) and \( \sin 33 \) equals .9999. Those are so close to the bottom and top that you can’t see their curves. The sine of 11 is near \(-1\) because \( \sin 22 \) is near zero. It is almost impossible to follow a single curve past the top—coming back down it is not the curve you think it is.

The points on the middle curve are at \( n = 0 \) and 44 and 88 and every number \( 44N \). Where does that curve come back to zero? In other words, when does \( 44N \) come very close to a multiple of \( \pi \)? We know that 44 is \( 14\pi + .02 \). More exactly 44 is \( 14\pi + .0177 \). So we multiply .0177 until we reach \( \pi \):

\[
\text{if } N = \pi/.0177 \text{ then } 44N = (14\pi + .0177)N = 14\pi N + \pi.
\]

This gives \( N = 177.5 \). At that point \( 44N = 7810 \). This is half the period of the sine curve. The sine of 7810 is very near zero.

If you follow the middle sine curve, you will see it come back to zero above 7810. The actual points on that curve have \( n = 44 + 177 \) and \( n = 44 + 178 \), with sines just above and below zero. Halfway between is \( n = 7810 \). **The equation for the middle sine curve is** \( y = \sin (\pi x/7810) \). Its period is 15,620—beyond our graph.

**Question** The fourth point on that middle curve looks the same as the fourth point coming down from \( \sin 3 \). What is this “double point?”

**Answer** 4 times 44 is 176. On the curve going up, the point is \( (176, \sin 176) \). On the curve coming down it is \( (179, \sin 179) \). **The sines of 176 and 179 differ only by .00003.**

The second graph spreads out this double point. Look above 176 and 179, at the center of a hexagon. You can follow the sine curve all the way across graph 2.

Only a little question remains. Why does graph 2 have hexagons? \( \text{i don’t know.} \)

The problem is with your eyes. To understand the hexagons, Doug Hardin plotted points on straight lines as well as sine curves. Graph 3 shows \( y = \text{fractional part of } n/2\pi \). Then he made a second copy, turned it over, and placed it on top. That produced graph 4—with hexagons. Graphs 3 and 4 are on the next page.
This is called a Moiré pattern. If you can get a transparent copy of graph 3, and turn it slowly over the original, you will see fantastic hexagons. They come from interference between periodic patterns—in our case $44/7$ and $25/4$ and $19/3$ are near $2\pi$. This interference is an enemy of printers, when color screens don't line up. It can cause vertical lines on a TV. Also in making cloth, operators get dizzy from seeing Moiré patterns move. There are good applications in engineering and optics—but we have to get back to calculus.

**1.7 Computing in Calculus**

Software is available for calculus courses—a lot of it. The packages keep getting better. Which program to use (if any) depends on cost and convenience and purpose. How to use it is a much harder question. These pages identify some of the goals, and also particular packages and calculators. Then we make a beginning (this is still Chapter 1) on the connection of computing to calculus.

The discussion will be informal. It makes no sense to copy the manual. Our aim is to support, with examples and information, the effort to use computing to help learning.

For calculus, the greatest advantage of the computer is to offer graphics. You see the function, not just the formula. As you watch, $f(x)$ reaches a maximum or a minimum or zero. A separate graph shows its derivative. Those statements are not 100% true, as everybody learns right away—as soon as a few functions are typed in. But the power to see this subject is enormous, because it is adjustable. If we don't like the picture we change to a new viewing window.

This is computer-based graphics. It combines numerical computation with graphical computation. You get pictures as well as numbers—a powerful combination. The computer offers the experience of actually working with a function. The domain and range are not just abstract ideas. You choose them. May I give a few examples.

**EXAMPLE 1** Certainly $x^3$ equals $3^x$ when $x = 3$. Do those graphs ever meet again? At this point we don't know the full meaning of $3^x$, except when $x$ is a nice number. (Neither does the computer.) Checking at $x = 2$ and 4, the function $x^3$ is smaller both times: $2^3$ is below $3^2$ and $4^3 = 64$ is below $3^4 = 81$. If $x^3$ is always less than $3^x$ we ought to know—these are among the basic functions of mathematics.
The computer will answer numerically or graphically. At our command, it solves \( x^3 = 3^x \). At another command, it plots both functions—this shows more. The screen proves a point of logic (or mathematics) that escaped us. If the graphs cross once, they must cross again—because \( 3^x \) is higher at 2 and 4. A crossing point near 2.5 is seen by zooming in. I am less interested in the exact number than its position—it comes before \( x = 3 \) rather than after.

A few conclusions from such a basic example:

1. A supercomputer is not necessary.
2. High-level programming is not necessary.
3. We can do mathematics without completely understanding it.

The third point doesn’t sound so good. Write it differently: We can learn mathematics while doing it. The hardest part of teaching calculus is to turn it from a spectator sport into a workout. The computer makes that possible.

**EXAMPLE 2** (mental computer) Compare \( x^2 \) with \( 2^x \). The functions meet at \( x = 2 \). Where do they meet again? Is it before or after 2?

That is mental computing because the answer happens to be a whole number (4). Now we are on a different track. Does an accident like \( 2^4 = 4^2 \) ever happen again? Can the machine tell us about integers? Perhaps it can plot the solutions of \( x^b = b^x \). I asked Mathematica for a formula, hoping to discover \( x \) as a function of \( b \)—but the program just gave back the equation. For once the machine typed HELP instead of the user.

Well, mathematics is not helpless. I am proud of calculus. There is a new exercise at the end of Section 6.4, to show that we never see whole numbers again.

**EXAMPLE 3** Find the number \( b \) for which \( x^b = b^x \) has only one solution (at \( x = b \)).

When \( b \) is 3, the second solution is below 3. When \( b \) is 2, the second solution (4) is above 2. If we move \( b \) from 2 to 3, there must be a special “double point”—where the curves barely touch but don’t cross. For that particular \( b \)—and only for that one value—the curve \( x^b \) never goes above \( b^x \).

This special point \( b \) can be found with computer-based graphics. In many ways it is the “center point of calculus.” Since the curves touch but don’t cross, they are tangent. They have the same slope at the double point. Calculus was created to work with slopes, and we already know the slope of \( x^2 \). Soon comes \( x^6 \). Eventually we discover the slope of \( b^x \), and identify the most important number in calculus.

The point is that this number can be discovered first by experiment.

**EXAMPLE 4** Graph \( y(x) = e^x - x^x \). Locate its minimum.

The next example was proposed by Don Small. Solve \( x^4 - 11x^3 + 5x - 2 = 0 \). The first tool is algebra—try to factor the polynomial. That succeeds for quadratics, and then gets extremely hard. Even if the computer can do algebra better than we can, factoring is seldom the way to go. In reality we have two good choices:

1. (Mathematics) Use the derivative. Solve by Newton’s method.
2. (Graphics) Plot the function and zoom in.

Both will be done by the computer. Both have potential problems! Newton’s method is fast, but that means it can fail fast. (It is usually terrific.) Plotting the graph is also fast—but solutions can be outside the viewing window. This particular function is
zero only once, in the standard window from $-10$ to $10$. The graph seems to be leaving zero, but mathematics again predicts a second crossing point. So we zoom out before we zoom in.

The use of the zoom is the best part of graphing. Not only do we choose the domain and range, we change them. The viewing window is controlled by four numbers. They can be the limits $A < x < B$ and $C < y < D$. They can be the coordinates of two opposite corners: $(A, C)$ and $(B, D)$. They can be the center position $(a, b)$ and the scale factors $c$ and $d$. Clicking on opposite corners of the zoom box is the fastest way, unless the center is unchanged and we only need to give scale factors. (Even faster: Use the default factors.) Section 3.4 discusses the centering transform and zoom transform—a change of picture on the screen and a change of variable within the function.

**Example 5** Find all real solutions to $x^4 - 11x^3 + 5x - 2 = 0$.

**Example 6** Zoom out and in on the graphs of $y = \cos 40x$ and $y = x \sin (1/x)$. Describe what you see.

**Example 7** What does $y = (\tan x - \sin x)/x^3$ become at $x = 0$? For small $x$ the machine eventually can’t separate $\tan x$ from $\sin x$. It may give $y = 0$. Can you get close enough to see the limit of $y$?

For these examples, and for most computer exercises in this book, a menu-driven system is entirely adequate. There is a list of commands to choose from. The user provides a formula for $y(x)$, and many functions are built in. A calculus supplement can be very useful—MicroCalc or True BASIC or Exploring Calculus or MPP (in the public domain). Specific to graphics are Surface Plotter and Master Grapher and Gyrographics (animated). The best software for linear algebra is MATLAB.

Powerful packages are increasing in convenience and decreasing in cost. They are capable of symbolic computation—which opens up a third avenue of computing in calculus.

**Symbolic Computation**

In symbolic computation, answers can be formulas as well as numbers and graphs. The derivative of $y = x^2$ is seen as “$2x$.” The derivative of $\sin t$ is “$\cos t$.” The slope of $b^x$ is known to the program. The computer does more than substitute numbers into formulas—it operates directly on the formulas. We need to think where this fits with learning calculus.

In a way, symbolic computing is close to what we ourselves do. Maybe too close—there is some danger that symbolic manipulation is all we do. With a higher-level language and enough power, a computer can print the derivative of $\sin(x^2)$. So why learn the chain rule? Because mathematics goes deeper than “algebra with formulas.” We deal with ideas.

I want to say clearly: Mathematics is not formulas, or computations, or even proofs, but ideas. The symbols and pictures are the language. The book and the professor and the computer can join in teaching it. The computer should be non-threatening (like this book and your professor)—you can work at your own pace. Your part is to learn by doing.

**Example 8** A computer algebra system quickly finds $100! = (100)(99)(98)\ldots(1)$. The number has 158 digits (not written out here). The last 24
digits are zeros. For $10! = 3628800$ there are seven digits and two zeros. Between 10 and 100, and beyond, are simple questions that need ideas:

1. How many digits (approximately) are in the number $N!$?
2. How many zeros (exactly) are at the end of $N!$?

For Question 1, the computer shows more than $N$ digits when $N = 100$. It will never show more than $N^2$ digits, because none of the $N$ terms can have more than $N$ digits. A much tighter bound would be $2N$, but is it true? Does $N!$ always have fewer than $2N$ digits?

For Question 2, the zeros in $10!$ can be explained. One comes from 10, the other from 5 times 2. (10 is also 5 times 2.) Can you explain the 24 zeros in $100!$? An idea from the card game blackjack applies here too: Count the fives.

Hard question: How many zeros at the end of $200!$?

The outstanding package for full-scale symbolic computation is Mathematica. It was used to draw graphs for this book, including $y = \sin n$ on the back cover. The complete command was `ListPlot[Table[\{Sin[n], \{n, 10000\}\]]`. This system has rewards and also drawbacks, including the price. Its original purpose, like MathCAD and MACSYMA and REDUCE, was not to teach calculus—but it can. The computer algebra system MAPLE is good.

As I write in 1990, DERIVE is becoming well established for the PC. For the Macintosh, Calculus TIL is a “sleeper” that deserves to be widely known. It builds on MAPLE and is much more accessible for calculus. An important alternative is Theorist. These are menu-driven (therefore easier at the start) and not expensive.

I strongly recommend that students share terminals and work together. Two at a terminal and 3–5 in a working group seems to be optimal. Mathematics can be learned by talking and writing—it is a human activity. Our goal is not to test but to teach and learn.

Writing in Calculus

May I emphasize the importance of writing? We totally miss it, when the answer is just a number. A one-page report is harder on instructors as well as students—but much more valuable. A word processor keeps it neat. You can’t write sentences without being forced to organize ideas—and part of yourself goes into it.

I will propose a writing exercise with options. If you have computer-based graphing, follow through on Examples 1–4 above and report. Without a computer, pick a paragraph from this book that should be clearer and make it clearer. Rewrite it with examples. Identify the key idea at the start, explain it, and come back to express it differently at the end. Ideas are like surfaces—they can be seen many ways.

Every reader will understand that in software there is no last word. New packages keep coming (Analyzer and EPIC among them). The biggest challenges at this moment are three-dimensional graphics and calculus workbooks. In 3D, the problem is the position of the eye—since the screen is only 2D. In workbooks, the problem is to get past symbol manipulation and reach ideas. Every teacher, including this one, knows how hard that is and hopes to help.

GRAPHING CALCULATORS

The most valuable feature for calculus—computer-based graphing—is available on hand calculators. With trace and zoom their graphs are quite readable. By creating the graphs you subconsciously learn about functions. These are genuinely personal computers, and the following pages aim to support and encourage their use.
Programs for a hand-held machine tend to be simple and short. We don’t count the zeros in 100 factorial (probably we could). A calculator finds crossing points and maximum points to good accuracy. Most of all it allows you to explore calculus by yourself. You set the viewing window and define the function. Then you see it.

There is a choice of calculators—*which one to buy?* For this book there was also a choice—*which one to describe?* To provide you with listings for useful programs, we had to choose. Fortunately the logic is so clear that you can translate the instructions into any language—for a computer as well as a calculator. The programs given here are the "greatest common denominator" of computing in calculus.

The range of choices starts with the Casio *fx 7000G*—the first and simplest, with very limited memory but a good screen. The Casio 7500, 8000, and 8500 have increasing memory and extra features. The Sharp *E1-5200* (or 9000 in Canada and Europe) is comparable to the Casio 8000. These machines have *algebraic entry*—the normal order as in $y = x + 3$. They are inexpensive and good. More expensive and much more powerful are the Hewlett-Packard calculators—the *HP-28S* and *HP-48SX*. They have large memories and extensive menus (and symbolic algebra). They use *reverse Polish notation*—numbers first in the stack, then commands. They require extra time and effort, and other books do justice to their amazing capabilities. It is estimated that those calculators could get 95 on a typical calculus exam.

While this book was being written, Texas Instruments produced a new graphing calculator: the *TI-81*. It is closer to the Casio and Sharp (emphasis on graphing, easy to learn, no symbolic algebra, moderate price). With earlier machines as a starting point, many improvements were added. There is some risk in a choice that is available only at before this textbook is published, and we hope that the experts we asked are right. Anyway, *our programs are for the TI-81*. It is impressive.

These few pages are no substitute for the manual that comes with a calculator. A valuable supplement is a guide directed especially at calculus—my absolute favorites are *Calculus Activities for Graphic Calculators* by Dennis Pence (PWS-Kent, 1990 for the Casio and Sharp and * HP-28S*, 1991 for the *TI-81*). A series of *Calculator Enhancements*, using *HP’s*, is being published by Harcourt Brace Jovanovich. What follows is an introduction to one part of a calculus laboratory. Later in the book, we supply *TI-81* programs close to the mathematics and the exercises that they are prepared for.

A few words to start: To select from a menu, press the item number and ENTER. Edit a command line using *DEL*(ete) and *INS*(ert). Every line ends with ENTER. For calculus select radians on the MODE screen. For powers use *x*¹. For special powers choose $x^{2}$, $x^{-1}$, $\sqrt{x}$. Multiplication has priority, so $(−)3 + 2 \times 2$ produces 1. Use keys for *SIN*, *IF*, *IS*, ... When you press letters, I multiplies S. If a program says $3 \rightarrow I$, type $3 \text{STO} \rightarrow I$. Storage locations are A to Z or Greek θ.

**Functions** A graphing calculator helps you (forces you?) to understand the concept of a function. It also helps you to understand specific functions—especially when changing the viewing window.

To evaluate $y = x^2 - 2x$ *just once*, use the home screen. To define $y(x)$ for repeated use, move to the function edit screen. Press MODE, choose FUNCTION, and press Y=. Then type in the formula. *Important tip:* for $x$ on the *TI-81*, the key $X \uparrow T$ is faster than two steps $\text{ALPHA} \rightarrow X$. The $Y=$ edit screen is the same place where the formula is needed for graphing.
Example \( Y_1 = X^2 - 2X \) ENTER on the \( Y = \) screen. 4 STO X ENTER on the home screen. \( Y_1 \) ENTER on the \( Y-VARS \) screen. The screen shows 8, which is \( Y(4) \). The formula remains when the calculator is off.

**Graphing** You specify the \( X \) range and \( Y \) range. (We should say \( X \) domain but we don’t.) The screen is a grid of \( 96 \times 64 \) little rectangles called “pixels.” The first column of pixels represents \( X_{\text{min}} \) and the last column is \( X_{\text{max}} \). Press RANGE to reset.

With \( X_{\text{res}} = 1 \) the function is evaluated 96 times as it is graphed. \( X_{\text{scL}} \) and \( Y_{\text{scL}} \) give the spaces between ticks on the axes.

The \textit{ZOOM} menu is a fast way to set ranges. \textit{ZOOM \textbf{S}tandard} gives the default \(-10 < x < 10, -10 < y < 10\). \textit{ZOOM \textbf{T}rig} gives \(-2\pi < x < 2\pi, -3 < y < 3\).

The keystroke \textit{GRAPH} shows the graphing screen with the current functions.

Example Set the ranges \((-)2 < X < 3 \) and \((-)150 < Y < 50\). Press \( Y= \) and store \( Y_1 = X \) (in \textit{MATH}) \(- 28X^2 + 15X + 36 \) ENTER. Press \textit{GRAPH}. You won’t see much of the graph! Press \textit{RANGE} and reset \((-)10 < X < 30,(--)4000 < Y < 2300\). Press \textit{GRAPH}. See a cubic polynomial.

“Smart Graph” recalls the graph instantly without redrawing it, if no settings have changed. The \textit{DRAW} menu is for points, lines, and shaded regions. This is perfect for our piecewise linear functions—just connect the breakpoints with lines. In Section 3.6 the lines show an iteration by its “cobweb.”

**Programming** This book contains programs that you can type in once and save. We chose \textit{Autoscaling}, \textit{Newton’s Method}, \textit{Secant Method}, \textit{Cobweb Iteration}, and \textit{Numerical Integration}. You will create others—to do calculations or to add features that are not available as single keystrokes. The calculator is like a computer, with a fairly small set of instructions. \textit{One difference}: Memory is too precious to store comments with the code. You have to see the logic by rereading the program.

To enter the world of programming, press \textit{PRGM}. Each \textit{PRGM} submenu lists all programs by name—a digit, a letter, or \( \theta \) (37 names). The program \textit{title} has up to eight characters. Select the \textit{EDIT} submenu and press \( G \) for the edit screen. Type the title \textit{GRAPHS} and press \( \text{ENTER} \). Practice on this one:

\[
: \text{"x}^2+x\text{" STO (Y-VARS) Y1 ENTER} \\
: \text{"x-1" STO (Y-VARS) Y2 ENTER} \\
: \text{(PRGM)(I/O) DispGraph}
\]

The menus to call are in parentheses. Leave the edit screen with \textit{QUIT} (not \textit{CLEAR}—that erases the line with the cursor). Set the default window by \textit{ZOOM \textbf{S}tandard}.

To execute, press \textit{PRGM} \( \text{(EXEC)} \text{ G ENTER} \). The program draws the graphs. It leaves \( Y_1 \) and \( Y_2 \) on the \( Y= \) screen. To erase the program from the home screen, press \textit{(PRGM)}(\text{ERASE})\text{ G}. Practice again by creating \textit{Program2:FUNC}. Type \( : \sqrt{X} \) STO Y and : \textit{(PRGM)(I/O)DispY}. Move to the home screen, store \( X \) by \( 4 \) STO X ENTER, and execute by \textit{(PRGM)(EXEC)2 ENTER}. Also try \( X = -1 \). When it fails to imagine \( i \), select \( 1: \text{Goto Error} \).

**Piecewise functions** and \textit{Input} (to a running program). The definition of a piecewise function includes the \textit{domain of each piece}. Logical tests like "\( \text{IF } X \geq 7 \)" determine which domain the input value \( X \) falls into. An \textit{IF} statement only affects the following line—which is executed when \textit{TEST = 1} (meaning \textit{true}) and skipped when \textit{TEST = 0} (meaning \textit{false}). \textit{IF} commands are in the \textit{PRGM} \( \text{(CTL)} \) submenu; \textit{TEST} calls the menu of inequalities.
1 Introduction to Calculus

An input value $X = 4$ need not be stored in advance. Program $P$ stops while running to request input. Execute with $P$ ENTER after selecting the PRGM (EXEC) menu. Answer $?$ with $4$ and ENTER. After completion, rerun by pressing ENTER again. The function is $y = 14 - x$ if $x < 7$, $y = x$ if $x \geq 7$.

**PrgmP: PIECES**

:Disp "X=" PGRM (I/O) Ask for input
:Input X PGRM (I/O) Screen ? ENTER X
:14-X -> Y First formula for all $X$
:If 7 < X PRGM (CTL) TEST
:Y -> X Overwrite if TEST = 1
:Disp Y Display $Y(X)$

Overwriting is faster than checking both ends $A \leq X \leq B$ for each piece. Even faster: a whole formula $(14 - X)(X < 7) + (X)(7 \leq X)$ can go on a single line using 1 and 0 from the tests. Compute-store-display $Y(X)$ as above, or define $Y 1$ on the edit screen.

**Exercise** Define a third piece $Y = 8 + X$ if $X < 3$. Rewrite $P$ using $Y 1 =$.

**TRACE and ZOOM** The best feature is graphing. But a whole graph can be like a whole book—too much at once. You want to focus on one part. A computer or calculator will trace along the graph, stop at a point, and zoom in.

There is also ZOOM OUT, to widen the ranges and see more. Our eyes work the same way—they put together information on different scales. Looking around the room uses an amazingly large part of the human brain. With a big enough computer we can try to imitate the eyes—this is a key problem in artificial intelligence. With a small computer and a zoom feature, we can use our eyes to understand functions.

Press TRACE to locate a point on the graph. A blinking cursor appears. Move left or right—the cursor stays on the graph. Its coordinates appear at the bottom of the screen. When $x$ changes by a pixel, the calculator evaluates $y(x)$. To solve $y(x) = 0$, read off $x$ at the point when $y$ is nearest to zero. To minimize or maximize $y(x)$, read off the smallest and largest $y$. In all these problems, zoom in for more accuracy.

To blow up a figure we can choose new ranges. The fast way is to use a ZOOM command. For a preset range, use ZOOM Standard or ZOOM Trig. To shrink or stretch by $X \text{Fact}$ or $Y \text{Fact}$ (default values 4), use ZOOM In or ZOOM Out. Choose the center point and press ENTER. The new graph appears. Change those scaling factors with ZOOM Set Factors. Best of all, create your own viewing window. Press ZOOM Box.

To draw the box, move the cursor to one corner. Press ENTER and this point is a small square. The same keys move a second (blinking) square to the opposite corner—the box grows as you move. Press ENTER, and the box is the new viewing window. The graphs show the same function with a change of scale. Section 3.4 will discuss the mathematics—here we concentrate on the graphics.

**EXAMPLE 9** Place $Y 1 = X \sin (1/X)$ in the $Y =$ edit screen. Press ZOOM Trig for a first graph. Set $X \text{Fact} = 1$ and $Y \text{Fact} = 2.5$. Press ZOOM In with center at $(0,0)$. To see a larger picture, use $X \text{Fact} = 10$ and $Y \text{Fact} = 1$. Then ZOOM Out again. As $X$ gets large, the function $X \sin (1/X)$ approaches _____.

Now return to ZOOM Trig. ZOOM In with the factors set to 4 (default). Zoom again by pressing ENTER. With the center and the factors fixed, this is faster than drawing a zoom box.
1.7 Computing in Calculus

EXAMPLE 10  Repeat for the more erratic function $Y = \sin(1/X)$. After ZOOM Trig, create a box to see this function near $X = .01$. The $Y$ range is now ________.

Scaling is crucial. For a new function it can be tedious. A formula for $y(x)$ does not easily reveal the range of $y$'s, when $A \leq x \leq B$ is given. The following program is often more convenient than zooms. It samples the function $L = 19$ times across the $x$-range (every 5 pixels). The inputs $X_{\text{min}}, X_{\text{max}}, Y_1$ are previously stored on other screens. After sampling, the program sets the $y$-range from $C = Y_{\text{min}}$ to $D = Y_{\text{max}}$ and draws the graph.

Notice the loop with counter $K$. The loop ends with the command $IS > (K, L)$, which increases $K$ by 1 and skips a line if the new $K$ exceeds $L$. Otherwise the command Goto 1 restarts the loop. The screen shows the short form on the left.

Example:  $Y_1 = x^3 + 10x^2 - 7x + 42$ with range $X_{\text{min}} = -12$ and $X_{\text{max}} = 10$. Set tick spacing $X_{\text{scl}} = 4$ and $Y_{\text{scl}} = 250$. Execute with PRGM (EXEC) A ENTER. For this program we also list menu locations and comments.

```
ProgMA :AUTOSCL  Menu (Submenu) Comment
:All-Off             Y-VARS (OFF) Turn off functions
:Xmin→ A             VARS (RNG) Store Xmin using STO
:19 → L              Store number of evaluations (19)
:((Xmax-A)/L → H     Spacing between evaluations
:A → X               Start at $x = A$
:Y1 → C              Y-VARS (Y) Evaluate the function
:C → D               Start $C$ and $D$ with this value
:1 → K               Initialize counter $K = 1$
:Lbl 1               PRGM (CTL) Mark loop start
:A+KH → X            Calculate next $x$
:Y1 → Y              Evaluate function at $x$
:IF Y < C            PRGM (CTL) New minimum?
:Y → C               Update $C$
:IF D < Y            PRGM (CTL) New maximum?
:Y → D               Update $D$
:IS>(K,L)           PRGM (CTL) Add 1 to $K$, skip Goto if $>L$
:Goto 1              PRGM (CTL) Loop return to Lbl 1
:Y1-0n               Y-VARS (ON) Turn on $Y_1$
:C → Ymin            STO VARS (RNG) Set $Y_{\text{min}} = C$
:D → Ymax            STO VARS (RNG) Set $Y_{\text{max}} = D$
:DispGraph           PRGM (I/O) Generate graph
```
# Contents

## CHAPTER 1  
**Introduction to Calculus**

1.1 Velocity and Distance ........................................... 1  
1.2 Calculus Without Limits ........................................... 8  
1.3 The Velocity at an Instant ......................................... 16  
1.4 Circular Motion .................................................... 22  
1.5 A Review of Trigonometry ......................................... 29  
1.6 A Thousand Points of Light ........................................ 34  
1.7 Computing in Calculus ............................................. 36  

## CHAPTER 2  
**Derivatives**

2.1 The Derivative of a Function ....................................... 44  
2.2 Powers and Polynomials .......................................... 50  
2.3 The Slope and the Tangent Line ................................... 58  
2.4 Derivative of the Sine and Cosine ................................. 64  
2.5 The Product and Quotient and Power Rules ..................... 71  
2.6 Limits .............................................................. 78  
2.7 Continuous Functions ............................................. 85  

## CHAPTER 3  
**Applications of the Derivative**

3.1 Linear Approximation .............................................. 91  
3.2 Maximum and Minimum Problems ................................ 96  
3.3 Second Derivatives: Minimum vs. Maximum .................... 105  
3.4 Graphs ............................................................ 112  
3.5 Ellipses, Parabolas, and Hyperbolas ............................. 121  
3.6 Iterations $x_{n+1} = F(x_n)$ ...................................... 130  
3.7 Newton’s Method and Chaos ...................................... 137  
3.8 The Mean Value Theorem and l'Hôpital’s Rule .................. 146
CHAPTER 2

Derivatives

2.1 The Derivative of a Function

This chapter begins with the definition of the derivative. Two examples were in Chapter 1. When the distance is \( t^2 \), the velocity is \( 2t \). When \( f(t) = \sin t \) we found \( v(t) = \cos t \). The velocity is now called the derivative of \( f(t) \). As we move to a more formal definition and new examples, we use new symbols \( f' \) and \( df/dt \) for the derivative.

At time \( t \), the derivative \( f'(t) \) or \( df/dt \) or \( v(t) \) is

\[
f'(t) = \lim_{\Delta t \to 0} \frac{f(t + \Delta t) - f(t)}{\Delta t}.
\]  

(1)

The ratio on the right is the average velocity over a short time \( \Delta t \). The derivative, on the left side, is its limit as the step \( \Delta t \) (delta \( t \)) approaches zero.

Go slowly and look at each piece. The distance at time \( t + \Delta t \) is \( f(t + \Delta t) \). The distance at time \( t \) is \( f(t) \). Subtraction gives the change in distance, between those times. We often write \( \Delta f \) for this difference: \( \Delta f = f(t + \Delta t) - f(t) \). The average velocity is the ratio \( \Delta f/\Delta t \)—change in distance divided by change in time.

The limit of the average velocity is the derivative, if this limit exists:

\[
\frac{df}{dt} = \lim_{\Delta t \to 0} \frac{\Delta f}{\Delta t}.
\]

(2)

This is the neat notation that Leibniz invented: \( \Delta f/\Delta t \) approaches \( df/dt \). Behind the innocent word “limit” is a process that this course will help you understand.

Note that \( \Delta f \) is not \( \Delta \) times \( f \)! It is the change in \( f \). Similarly \( \Delta t \) is not \( \Delta \) times \( t \). It is the time step, positive or negative and eventually small. To have a one-letter symbol we replace \( \Delta t \) by \( h \).

The right sides of (1) and (2) contain average speeds. On the graph of \( f(t) \), the distance up is divided by the distance across. That gives the average slope \( \Delta f/\Delta t \).

The left sides of (1) and (2) are instantaneous speeds \( df/dt \). They give the slope at the instant \( t \). This is the derivative \( df/dt \) (when \( \Delta t \) and \( \Delta f \) shrink to zero). Look again
at the calculation for $f(t) = t^2$:

$$\frac{\Delta f}{\Delta t} = \frac{f(t + \Delta t) - f(t)}{\Delta t} = \frac{t^2 + 2t \Delta t + (\Delta t)^2 - t^2}{\Delta t} = 2t + \Delta t. \quad (3)$$

Important point: Those steps are taken before $\Delta t$ goes to zero. **If we set $\Delta t = 0$ too soon, we learn nothing.** The ratio $\Delta f/\Delta t$ becomes $0/0$ (which is meaningless). The numbers $\Delta f$ and $\Delta t$ must approach zero together, not separately. Here their ratio is $2t + \Delta t$, the average speed.

To repeat: Success came by writing out $(t + \Delta t)^2$ and subtracting $t^2$ and dividing by $\Delta t$. Then and only then can we approach $\Delta t = 0$. The limit is the derivative $2t$.

There are several new things in formulas (1) and (2). Some are easy but important, others are more profound. The idea of a function we will come back to, and the definition of a limit. But the notations can be discussed right away. They are used constantly and you also need to know how to read them aloud:

- $f(t) =$ "$f$ of $t$" = the value of the function $f$ at time $t$
- $\Delta t =$ "delta $t$" = the time step forward or backward from $t$
- $f(t + \Delta t) =$ "$f$ of $t$ plus delta $t$" = the value of $f$ at time $t + \Delta t$
- $\Delta f =$ "delta $f$" = the change $f(t + \Delta t) - f(t)$
- $\Delta f / \Delta t =$ "delta $f$ over delta $t$" = the average velocity
- $f'(t) =$ "$f$ prime of $t$" = the value of the derivative at time $t$
- $df/dt =$ "$d f \; d t$" = the same as $f'$ (the instantaneous velocity)
- $\lim_{\Delta t \to 0} =$ "limit as delta $t$ goes to zero" = the process that starts with numbers $\Delta f/\Delta t$ and produces the number $df/dt$.

From those last words you see what lies behind the notation $df/dt$. The symbol $\Delta t$ indicates a nonzero (usually short) length of time. The symbol $dt$ indicates an infinitesimal (even shorter) length of time. Some mathematicians work separately with $df$ and $dt$, and $df/dt$ is their ratio. For us $df/dt$ is a single notation (don't cancel $d$ and don't cancel $\Delta$). The derivative $df/dt$ is the limit of $\Delta f/\Delta t$. **When that notation $df/dt$ is awkward, use $f'$ or $v$.**

**Remark** The notation hides one thing we should mention. The time step can be negative just as easily as positive. We can compute the average $\Delta f/\Delta t$ over a time interval before the time $t$, instead of after. This ratio also approaches $df/dt$.

The notation also hides another thing: **The derivative might not exist.** The averages $\Delta f/\Delta t$ might not approach a limit (it has to be the same limit going forward and backward from time $t$). In that case $f'(t)$ is not defined. At that instant there is no clear reading on the speedometer. This will happen in Example 2.

**EXAMPLE 1** (Constant velocity $V = 2$) The distance $f$ is $V$ times $t$. The distance at time $t + \Delta t$ is $V$ times $t + \Delta t$. The difference $\Delta f$ is $V$ times $\Delta t$:

$$\frac{\Delta f}{\Delta t} = \frac{V \Delta t}{\Delta t} = V \quad \text{so the limit is} \quad \frac{df}{dt} = V.$$  

The derivative of $Vt$ is $V$. The derivative of $2t$ is 2. The averages $\Delta f/\Delta t$ are always $V = 2$, in this exceptional case of a constant velocity.
EXAMPLE 2  Constant velocity $2$ up to time $t = 3$, then stop.

For small times we still have $f(t) = 2t$. But after the stopping time, the distance is fixed at $f(t) = 6$. The graph is flat beyond time 3. Then $f(t + \Delta t) = f(t)$ and $\Delta f = 0$ and the derivative of a constant function is zero:

$$t > 3: \quad f'(t) = \lim_{\Delta t \to 0} \frac{f(t + \Delta t) - f(t)}{\Delta t} = \lim_{\Delta t \to 0} \frac{0}{\Delta t} = 0.$$  (4)

In this example the derivative is not defined at the instant when $t = 3$. The velocity falls suddenly from 2 to zero. The ratio $\Delta f / \Delta t$ depends, at that special moment, on whether $\Delta t$ is positive or negative. The average velocity after time $t = 3$ is zero. The average velocity before that time is 2. When the graph of $f$ has a corner, the graph of $v$ has a jump. It is a step function.

One new part of that example is the notation $(df/dt$ or $f'$ instead of $v)$. Please look also at the third figure. It shows how the function takes $t$ (on the left) to $f(t)$. Especially it shows $\Delta t$ and $\Delta f$. At the start, $\Delta f / \Delta t$ is 2. After the stop at $t = 3$, all $t$'s go to the same $f(t) = 6$. So $\Delta f = 0$ and $df/dt = 0$.

---

**Fig. 2.1** The derivative is 2 then 0. It does not exist at $t = 3$.

---

**THE DERIVATIVE OF $1/t$**

Here is a completely different slope, for the “demand function” $f(t) = 1/t$. The demand is $1/t$ when the price is $t$. A high price $t$ means a low demand $1/t$. Increasing the price reduces the demand. The calculus question is: How quickly does $1/t$ change when $t$ changes? The “marginal demand” is the slope of the demand curve.

The big thing is to find the derivative of $1/t$ once and for all. It is $-1/t^2$.

EXAMPLE 3  $f(t) = \frac{1}{t}$ has $\Delta f = \frac{1}{t + \Delta t} - \frac{1}{t}$. This equals $\frac{t - (t + \Delta t)}{t(t + \Delta t)} = \frac{-\Delta t}{t(t + \Delta t)}$.

Divide by $\Delta t$ and let $\Delta t \to 0$: $\frac{\Delta f}{\Delta t} = \frac{-1}{t(t + \Delta t)}$ approaches $\frac{df}{dt} = \frac{-1}{t^2}$.

Line 1 is algebra, line 2 is calculus. The first step in line 1 subtracts $f(t)$ from $f(t + \Delta t)$. The difference is $1/(t + \Delta t)$ minus $1/t$. The common denominator is $t$ times $t + \Delta t$—this makes the algebra possible. We can’t set $\Delta t = 0$ in line 2, until we have divided by $\Delta t$.

The average is $\Delta f / \Delta t = -1/(t + \Delta t)$. Now set $\Delta t = 0$. The derivative is $-1/t^2$.

Section 2.4 will discuss the first of many cases when substituting $\Delta t = 0$ is not possible, and the idea of a limit has to be made clearer.
2.1 The Derivative of a Function

Check the algebra at \( t = 2 \) and \( t + \Delta t = 3 \). The demand \( 1/t \) drops from 1/2 to 1/3. The difference is \( \Delta f = -1/6 \), which agrees with \(-1/(2)(3)\) in line 1. As the steps \( \Delta f \) and \( \Delta t \) get smaller, their ratio approaches \(-1/(2)(2) = -1/4\).

**This derivative is negative.** The function \( 1/t \) is decreasing, and \( \Delta f \) is below zero. The graph is going downward in Figure 2.2, and its slope is negative:

> **An increasing** \( f(t) \) **has positive slope.** **A decreasing** \( f(t) \) **has negative slope.**

The slope \(-1/t^2\) is very negative for small \( t \). A price increase severely cuts demand.

The next figure makes a small but important point. There is nothing sacred about \( t \). Other letters can be used—especially \( x \). A quantity can depend on **position instead of time**. The height changes as we go west. The area of a square changes as the side changes. Those are not affected by the passage of time, and there is no reason to use \( t \). You will often see \( y = f(x) \), with \( x \) across and \( y \) up—connected by a function \( f \).

Similarly, \( f \) is not the only possibility. Not every function is named \( f \)! That letter is useful because it stands for the word function—but we are perfectly entitled to write \( y(x) \) or \( y(t) \) instead of \( f(x) \) or \( f(t) \). The distance up is a function of the distance across. This relationship "\( y \) of \( x \)" is all-important to mathematics.

The slope is also a function. Calculus is about two functions, \( y(x) \) and \( dy/dx \).

**Question** If we add 1 to \( y(x) \), what happens to the slope? **Answer** Nothing.

**Question** If we add 1 to the slope, what happens to the height? **Answer** ________.

The symbols \( t \) and \( x \) represent **independent variables**—they take any value they want to (in the domain). Once they are set, \( f(t) \) and \( y(x) \) are determined. Thus \( f \) and \( y \) represent **dependent variables**—they depend on \( t \) and \( x \). A change \( \Delta t \) produces a
change $\Delta f$. A change $\Delta x$ produces $\Delta y$. The independent variable goes inside the parentheses in $f(t)$ and $y(x)$. It is not the letter that matters, it is the idea:

- independent variable $t$ or $x$
- dependent variable $f$ or $g$ or $y$ or $z$ or $u$
- derivative $df/dt$ or $df/dx$ or $dy/dx$ or $\ldots$

The derivative $dy/dx$ comes from $[\text{change in } y]$ divided by $[\text{change in } x]$. The time step becomes a space step, forward or backward. The slope is the rate at which $y$ changes with $x$. **The derivative of a function is its “rate of change.”**

I mention that physics books use $x(t)$ for distance. Darn it.

To emphasize the definition of a derivative, here it is again with $y$ and $x$:

$$\frac{\Delta y}{\Delta x} = \frac{y(x + \Delta x) - y(x)}{\Delta x} = \frac{\text{distance up}}{\text{distance across}} \quad dy = \lim_{\Delta x \to 0} \frac{\Delta y}{\Delta x} = y'(x).$$

The notation $y'(x)$ pins down the point $x$ where the slope is computed. In $dy/dx$ that extra precision is omitted. This book will try for a reasonable compromise between logical perfection and ordinary simplicity. The notation $dy/dx(x)$ is not good; $y'(x)$ is better; when $x$ is understood it need not be written in parentheses.

You are allowed to say that the function is $y = x^2$ and the derivative is $y' = 2x$—even if the strict notation requires $y(x) = x^2$ and $y'(x) = 2x$. You can even say that the function is $x^2$ and its derivative is $2x$ and its **second derivative** is $2$—provided everybody knows what you mean.

Here is an example. It is a little early and optional but terrific. You get excellent practice with letters and symbols, and out come new derivatives.

**EXAMPLE 4** If $u(x)$ has slope $du/dx$, what is the slope of $f(x) = (u(x))^2$?

From the derivative of $x^2$ this will give the derivative of $x^4$. In that case $u = x^2$ and $f = x^4$. First point: **The derivative of $u^2$ is not $(du/dx)^2$**. We do not square the derivative $2x$. To find the “square rule” we start as we have to—with $\Delta f = f(x + \Delta x) - f(x)$:

$$\Delta f = (u(x + \Delta x))^2 - (u(x))^2 = [u(x + \Delta x) + u(x)] [u(x + \Delta x) - u(x)].$$

This algebra puts $\Delta f$ in a convenient form. We factored $a^2 - b^2$ into $(a + b)$ times $(a - b)$. Notice that we don’t have $(\Delta u)^2$. We have $\Delta f$, the change in $u^2$. Now divide by $\Delta x$ and take the limit:

$$\frac{\Delta f}{\Delta x} = \left[ u(x + \Delta x) + u(x) \right] \frac{u(x + \Delta x) - u(x)}{\Delta x} \text{ approaches } 2u(x) \frac{du}{dx}. \quad (5)$$

This is the square rule: **The derivative of $(u(x))^2$ is $2u(x)$ times $du/dx$.** From the derivatives of $x^2$ and $1/x$ and $\sin x$ (all known) the examples give new derivatives.

**EXAMPLE 5** ($u = x^2$) The derivative of $x^4$ is $2u \, du/dx = 2(x^2)(2x) = 4x^3$.

**EXAMPLE 6** ($u = 1/x$) The derivative of $1/x^2$ is $2u \, du/dx = (2/x)(-1/x^2) = -2/x^3$.

**EXAMPLE 7** ($u = \sin x$, $du/dx = \cos x$) The derivative of $u^2 = \sin^2 x$ is $2 \sin x \cos x$.

Mathematics is really about ideas. The notation is created to express those ideas. Newton and Leibniz invented calculus independently, and Newton’s friends spent a lot of time proving that he was first. He was, but it was Leibniz who thought of
writing $dy/dx$—which caught on. It is the perfect way to suggest the limit of $\Delta y/\Delta x$. Newton was one of the great scientists of all time, and calculus was one of the great inventions of all time—but the notation must help. You now can write and speak about the derivative. What is needed is a longer list of functions and derivatives.

### 2.1 Exercises

#### Read-through questions

The derivative is the \_ \_ \_ \_ of $\Delta f/\Delta t$ as $\Delta t$ approaches \_ \_ \_ \_. Here $\Delta f$ equals \_ \_ \_ \_. The step $\Delta t$ can be positive or \_ \_ \_. The derivative is written \_ \_ \_ \_ or \_ \_ \_ \_. If $f(t) = 2x + 3$ and $\Delta x = 4$ then $\Delta f = \_ \_ \_$. If $\Delta x = -1$ then $\Delta f = \_ \_ \_$. If $\Delta x = 0$ then $\Delta f = \_ \_ \_$. The slope is not 0/0 but $df/dx = \_ \_ \_ \_$. The derivative does not exist where $f(t)$ has a \_ \_ \_ \_.

1. Which of the following numbers (as is) gives $df/dt$ at time $t$? If in doubt test on $f(t) = t^2$.
   
   (a) $f(t + \Delta t) - f(t)$
   
   (b) $\lim_{h \to 0} \frac{f(t + 2h) - f(t)}{2h}$
   
   (c) $\lim_{\Delta t \to 0} \frac{f(t + \Delta t) - f(t)}{-\Delta t}$
   
   (d) $\lim_{h \to 0} \frac{f(t + \Delta t) - f(t)}{\Delta t}$

2. Suppose $f(x) = x^2$. Compute each ratio and set $h = 0$:
   
   (a) $\frac{f(x + h) - f(x)}{h}$
   
   (b) $\frac{f(x + 5h) - f(x)}{5h}$
   
   (c) $\frac{f(x + h) - f(x - h)}{2h}$
   
   (d) $\frac{f(x + 1) - f(x)}{h}$

3. For $f(x) = 3x$ and $g(x) = x + 3x$, find $f(4 + h)$ and $g(4 + h)$ and $f'(4)$ and $g'(4)$. Sketch the graphs of $f$ and $g$—why do they have the same slope?

4. Find three functions with the same slope as $f(x) = x^2$.

5. For $f(x) = 1/x$, sketch the graphs of $f(x) + 1$ and $f(x + 1)$. Which one has the derivative $-1/x^2$?

6. Choose $c$ so that the line $y = x$ is tangent to the parabola $y = x^2 + c$. They have the same slope where they touch.

7. Sketch the curve $y(x) = 1 - x^2$ and compute its slope at $x = 3$.

8. If $f(t) = 1/t$, what is the average velocity between $t = \frac{1}{2}$ and $t = 2$? What is the average between $t = \frac{1}{2}$ and $t = 1$? What is the average (to one decimal place) between $t = \frac{1}{2}$ and $t = 101/200$?

9. Find $\Delta y/\Delta x$ for $y(x) = x + x^2$. Then find $dy/dx$.

10. Find $\Delta y/\Delta x$ and $dy/dx$ for $y(x) = 1 + 2x + 3x^2$.

11. When $f(t) = 4/t$, simplify the difference $f(t + \Delta t) - f(t)$, divide by $\Delta t$, and set $\Delta t = 0$. The result is $f'(t)$.

12. Find the derivative of $1/t^2$ from $\Delta f(t) = 1/(t + \Delta t)^2 - 1/t^2$. Write $\Delta f$ as a fraction with the denominator $t^2(t + \Delta t)^2$. Divide the numerator by $\Delta t$ to find $\Delta f/\Delta t$. Set $\Delta t = 0$.

13. Suppose $f(t) = 7t$ to $t = 1$. Afterwards $f(t) = 7 + 9(t - 1)$.
   
   (a) Find $df/dt$ at $t = \frac{1}{2}$ and $t = \frac{3}{2}$.
   
   (b) Why doesn’t $f(t)$ have a derivative at $t = 1$?

14. Find the derivative of the derivative (the second derivative) of $y = 3x^2$. What is the third derivative?

15. Find numbers $A$ and $B$ so that the straight line $y = x$ fits smoothly with the curve $Y = A + Bx + x^2$ at $x = 1$. Smoothly means that $y = Y$ and $dy/dx = dY/dx$ at $x = 1$.

16. Find numbers $A$ and $B$ so that the horizontal line $y = 4$ fits smoothly with the curve $y = A + Bx + x^2$ at the point $x = 2$.

17. True (with reason) or false (with example):
   
   (a) If $f(t) < 0$ then $df/dt < 0$.
   
   (b) The derivative of $(f(t))^2$ is $2df/dt$.
   
   (c) The derivative of $2f(t)$ is $2df/dt$.
   
   (d) The derivative is the limit of $\Delta f$ divided by the limit of $\Delta t$.

18. For $f(x) = 1/x$ the centered difference $f(x + h) - f(x - h)$ is $1/(x + h) - 1/(x - h)$. Subtract by using the common denominator $(x + h)(x - h)$. Then divide by $2h$ and set $h = 0$. Why divide by $2h$ to obtain the correct derivative?

19. Suppose $y = mx + b$ for negative $x$ and $y = Mx + B$ for $x > 0$. The graphs meet if ______. The two slopes are ______. The slope at $x = 0$ is ______ (what is possible?)

20. The slope of $y = 1/x$ at $x = 1/4$ is $y' = -1/x^2 = -16$. At $h = 1/12$, which of these ratios is closest to $-16$?
   
   - $\frac{y(x + h) - y(x)}{h}$
   
   - $\frac{y(x) - y(x - h)}{h}$
   
   - $\frac{y(x + h) - y(x - h)}{2h}$

21. Find the average slope of $y = x^2$ between $x = x_1$ and $x = x_2$. What does this average approach as $x_2$ approaches $x_1$?

22. Redraw Figure 2.1 when $f(t) = 3 - 2t$ for $t \leq 2$ and $f(t) = -1$ for $t > 2$. Include $df/dt$.
23 Redraw Figure 2.3 for the function \( y(x) = 1 - (1/x) \). Include \( dy/dx \).

24 The limit of \( O/\Delta t \) as \( \Delta t \to 0 \) is not \( 0/0 \). Explain.

25 Guess the limits by an informal working rule. Set \( \Delta t = 0.1 \) and \(-0.1\) and imagine \( \Delta t \) becoming smaller:

(a) \( \frac{1 + \Delta t}{2 + \Delta t} \)  
(b) \( \frac{|\Delta t|}{\Delta t} \)  
(c) \( \frac{\Delta t + (\Delta t)^2}{\Delta t - (\Delta t)^2} \)  
(d) \( \frac{t + \Delta t}{t - \Delta t} \)

*26 Suppose \( f(x)/x \to 7 \) as \( x \to 0 \). Deduce that \( f(0) = 0 \) and \( f'(0) = 7 \). Give an example other than \( f(x) = 7x \).

27 What is \( \lim_{x \to 0} \frac{f(3 + x) - f(3)}{x} \) if it exists? What if \( x \to 1 \)?

Problems 28–31 use the square rule: \( d(u^2)/dx = 2u (du/dx) \).

28 Take \( u = x \) and find the derivative of \( x^2 \) (a new way).

29 Take \( u = x^4 \) and find the derivative of \( x^8 \) (using \( du/dx = 4x^3 \)).

30 If \( u = 1 \) then \( u^2 = 1 \). Then \( d1/dx \) is 2 times \( d1/dx \). How is this possible?

31 Take \( u = \sqrt{x} \). The derivative of \( u^2 = x \) is \( 1 = 2u (du/dx) \). So what is \( du/dx \), the derivative of \( \sqrt{x} \)?

32 The left figure shows \( f(t) = t^2 \). Indicate distances \( f(t + \Delta t) \) and \( \Delta t \) and \( \Delta f \). Draw lines that have slope \( \Delta f/\Delta t \) and \( f'(t) \).

33 The right figure shows \( f(x) \) and \( \Delta x \). Find \( \Delta f/\Delta x \) and \( f'(2) \).

34 Draw \( f(x) \) and \( \Delta x \) so that \( \Delta f/\Delta x = 0 \) but \( f'(x) \neq 0 \).

35 If \( f = u^2 \) then \( df/dx = 2u du/dx \). If \( g = f^2 \) then \( dg/dx = 2f df/dx \). Together those give \( g = u^4 \) and \( dg/dx = \quad \).

36 True or false, assuming \( f(0) = 0 \):

(a) If \( f(x) \leq x \) for all \( x \), then \( df/dx \leq 1 \).

(b) If \( df/dx \leq 1 \) for all \( x \), then \( f(x) \leq x \).

37 The graphs show \( \Delta f \) and \( \Delta f/\Delta x \) for \( f(x) = x^2 \). Why is \( 2x + h \) the equation for \( \Delta f/\Delta x \)? If \( h \) is cut in half, draw in the new graphs.

38 Draw the corresponding graphs for \( f(x) = \frac{1}{2}x \).

39 Draw \( 1/x \) and \( 1/(x + h) \) and \( \Delta f/\Delta x \)—either by hand with \( h = \frac{1}{2} \) or by computer to show \( h \to 0 \).

40 For \( y = e^x \), show on computer graphs that \( dy/dx = y \).

41 Explain the derivative in your own words.

---

2.2 Powers and Polynomials

This section has two main goals. One is to find the derivatives of \( f(x) = x^3 \) and \( x^4 \) and \( x^5 \) (and more generally \( f(x) = x^n \)). The power or exponent \( n \) is at first a positive integer. Later we allow \( x^n \) and \( x^{2.2} \) and every \( x^n \).

The other goal is different. While computing these derivatives, we look ahead to their applications. In using calculus, we meet equations with derivatives in them—"differential equations." It is too early to solve those equations. But it is not too early to see the purpose of what we are doing. Our examples come from economics and biology.
With \( n = 2 \), the derivative of \( x^2 \) is \( 2x \). With \( n = -1 \), the slope of \( x^{-1} \) is \(-1x^{-2}\). Those are two pieces in a beautiful pattern, which it will be a pleasure to discover. We begin with \( x^3 \) and its derivative \( 3x^2 \), before jumping to \( x^n \).

**EXAMPLE 1** If \( f(x) = x^3 \) then \( \Delta f = (x + h)^3 - x^3 = (x^3 + 3x^2 h + 3xh^2 + h^3) - x^3 \).

**Step 1:** Cancel \( x^3 \).

**Step 2:** Divide by \( h \).

**Step 3:** \( h \) goes to zero.

\[
\frac{\Delta f}{h} = 3x^2 + 3xh + h^2 \quad \text{approaches} \quad \frac{df}{dx} = 3x^2.
\]

That is straightforward, and you see the crucial step. The power \( (x + h)^3 \) yields four separate terms \( x^3 + 3x^2 h + 3xh^2 + h^3 \). (Notice 1, 3, 3, 1.) After \( x^3 \) is subtracted, we can divide by \( h \). At the limit \( (h = 0) \) we have \( 3x^2 \).

For \( f(x) = x^n \) the plan is the same. A step of size \( h \) leads to \( f(x + h) = (x + h)^n \). One reason for algebra is to calculate powers like \( (x + h)^n \), and if you have forgotten the binomial formula we can recapture its main point. Start with \( n = 4 \):

\[
(x + h)(x + h)(x + h)(x + h) = x^4 + ??\quad + h^4.
\]

Multiplying the four \( x \)'s gives \( x^4 \). Multiplying the four \( h \)'s gives \( h^4 \). These are the easy terms, but not the crucial ones. The subtraction \( (x + h)^4 - x^4 \) will remove \( x^4 \), and the limiting step \( h \to 0 \) will wipe out \( h^4 \) (even after division by \( h \)). *The products that matter are those with exactly one \( h \).* In Example 1 with \( (x + h)^3 \), this key term was \( 3x^2 h \). Division by \( h \) left \( 3x^2 \).

*With only one \( h \), there are \( n \) places it can come from.* Equation (1) has four \( h \)'s in parentheses, and four ways to produce \( x^3 h \). Therefore the key term is \( 4x^3 h \). (Division by \( h \) leaves \( 4x^3 \).) In general there are \( n \) parenthesees and \( n \) ways to produce \( x^{n-1} \), so the *binomial formula* contains \( nx^{n-1} h \):

\[
(x + h)^n = x^n + nx^{n-1}h + \cdots + h^n.
\]

2B For \( n = 1, 2, 3, 4, \ldots \), the derivative of \( x^n \) is \( nx^{n-1} \).

Subtract \( x^n \) from (2). Divide by \( h \). The key term is \( nx^{n-1} \). The rest disappears as \( h \to 0 \):

\[
\frac{\Delta f}{\Delta x} = \frac{(x + h)^n - x^n}{h} = \frac{nx^{n-1}h + \cdots + h^n}{h} \quad \text{so} \quad \frac{df}{dx} = nx^{n-1}.
\]

The terms replaced by the dots involve \( h^2 \) and \( h^3 \) and higher powers. After dividing by \( h \), they still have at least one factor \( h \). All those terms vanish as \( h \) approaches zero.

**EXAMPLE 2** \( (x + h)^4 = x^4 + 4x^3 h + 6x^2 h^2 + 4xh^3 + h^4 \). This is \( n = 4 \) in detail.

Subtract \( x^4 \), divide by \( h \), let \( h \to 0 \). The derivative is \( 4x^3 \). The coefficients 1, 4, 6, 4, 1 are in Pascal’s triangle below. For \( (x + h)^5 \) the next row is 1, 5, 10, ?.

Remark The missing terms in the binomial formula (replaced by the dots) contain all the products \( x^{n-j} h^j \). An \( x \) or an \( h \) comes from each parenthesis. The binomial coefficient "\( n \) choose \( j \)" is the number of ways to choose \( j \) \( h \)'s out of \( n \) parentheses. It involves \( n \) factorial, which is \( n(n - 1) \cdots (1) \). Thus \( 5! = 5 \cdot 4 \cdot 3 \cdot 2 \cdot 1 = 120 \).
2 Derivatives

These are numbers that gamblers know and love:

\[
\binom{n}{j} = \frac{n!}{j!(n-j)!}
\]

In the last row, the coefficient of \(x^3h\) is \(4!/1!3! = 4 \cdot 3 \cdot 2 \cdot 1 / 1 \cdot 3 \cdot 2 \cdot 1 = 4\). For the \(x^2h^2\) term, with \(j = 2\), there are \(4 \cdot 3 \cdot 2 \cdot 1/2 \cdot 1 \cdot 2 \cdot 1 = 6\) ways to choose two \(h\)'s. Notice that \(1 + 4 + 6 + 4 + 1\) equals 16, which is \(2^4\). Each row of Pascal's triangle adds to a power of 2.

Choosing 6 numbers out of 49 in a lottery, the odds are \(49 \cdot 48 \cdot 47 \cdot 46 \cdot 45 \cdot 44 / 6!\) to 1. That number is \(\binom{49}{6} = 13,983,816\). It is the coefficient of \(h^6\) in \((x + h)^49\). If \(\lambda\) times \(N\) tickets are bought, the expected number of winners is \(\lambda\). The chance of no winner is \(e^{-\lambda}\). The chance of one winner is \(\lambda e^{-\lambda}\). See Section 8.4.

Florida's lottery in September 1990 (these rules) had six winners out of 109,163,978 tickets.

DERIVATIVES OF POLYNOMIALS

Now we have an infinite list of functions and their derivatives:

\[
x, x^2, x^3, x^4, x^5, \ldots
\]

\[
1, 2x, 3x^2, 4x^3, 5x^4, \ldots
\]

The derivative of \(x^n\) is \(n\) times the next lower power \(x^{n-1}\). That rule extends beyond these integers 1, 2, 3, 4, 5 to all powers:

\[
f = 1/x \quad \text{has} \quad f' = -1/x^2: \quad \text{Example 3 of Section 2.1} \quad (n = -1)
\]

\[
f = 1/x^2 \quad \text{has} \quad f' = -2/x^3: \quad \text{Example 6 of Section 2.1} \quad (n = -2)
\]

\[
f = \sqrt{x} \quad \text{has} \quad f' = \frac{1}{2}x^{-1/2}: \quad \text{true but not yet checked} \quad (n = \frac{1}{2})
\]

Remember that \(x^{-2}\) means \(1/x^2\) and \(x^{-1/2}\) means \(1/\sqrt{x}\). Negative powers lead to decreasing functions, approaching zero as \(x\) gets large. Their slopes have minus signs.

**Question** What are the derivatives of \(x^{10}\) and \(x^{2.2}\) and \(x^{-1.2}\)?

**Answer** \(10x^9\) and \(2.2x^{1.2}\) and \(-\frac{1}{5}x^{-3.2}\). Maybe \((x + h)^{2.2}\) is a little unusual. Pascal's triangle can't deal with this fractional power, but the formula stays firm: After \(x^{2.2}\) comes \(2.2x^{1.2}h\). The complete binomial formula is in Section 10.5.

That list is a good start, but plenty of functions are left. What comes next is really simple. A tremendous number of new functions are “linear combinations” like

\[
f(x) = 6x^3 \quad \text{or} \quad 6x^3 + \frac{1}{2}x^2 \quad \text{or} \quad 6x^3 - \frac{1}{2}x^2.
\]

What are their derivatives? The answers are known for \(x^3\) and \(x^2\), and we want to multiply by 6 or divide by 2 or add or subtract. Do the same to the derivatives:

\[
f'(x) = 18x^2 \quad \text{or} \quad 18x^2 + x \quad \text{or} \quad 18x^2 - x.
\]

**2C** The derivative of \(c\) times \(f(x)\) is \(c\) times \(f'(x)\).

**2D** The derivative of \(f(x) + g(x)\) is \(f'(x) + g'(x)\).

The number \(c\) can be any constant. We can add (or subtract) any functions. The rules allow any combination of \(f\) and \(g\): The derivative of \(9f(x) - 7g(x)\) is \(9f'(x) - 7g'(x)\).
The reasoning is direct. When \( f(x) \) is multiplied by \( c \), so is \( f(x + h) \). The difference \( \Delta f \) is also multiplied by \( c \). All averages \( \Delta f / h \) contain \( c \), so their limit is \( cf' \). The only incomplete step is the last one (the limit). *We still have to say what "limit" means.*

Rule 2D is similar. Adding \( f + g \) means adding \( \Delta f + \Delta g \). Now divide by \( h \). In the limit as \( h \to 0 \) we reach \( f' + g' \)—because a limit of sums is a sum of limits. Any example is easy and so is the proof—it is the definition of limit that needs care (Section 2.6).

**You can now find the derivative of every polynomial.** A "polynomial" is a combination of \( 1, x, x^2, \ldots, x^n \)—for example \( 9 + 2x - x^5 \). That particular polynomial has slope \( 2 - 5x^4 \). Note that the derivative of 9 is zero! A constant just raises or lowers the graph, without changing its slope. It alters the mileage before starting the car.

The disappearance of constants is one of the nice things in differential calculus. The reappearance of those constants is one of the headaches in integral calculus. When you find \( v \) from \( f \), the starting mileage doesn’t matter. The constant in \( f \) has no effect on \( v \). (\( \Delta f \) is measured by a trip meter; \( \Delta t \) comes from a stopwatch.) To find distance from velocity, you need to know the mileage at the start.

### A LOOK AT DIFFERENTIAL EQUATIONS (FIND \( y \) FROM \( dy/dx \))

We know that \( y = x^3 \) has the derivative \( dy/dx = 3x^2 \). Starting with the function, we found its slope. Now reverse that process. **Start with the slope and find the function.** This is what science does all the time—and it seems only reasonable to say so.

Begin with \( dy/dx = 3x^2 \). The slope is given, the function \( y \) is not given.

**Question** Can you go backward to reach \( y = x^3 \)?

**Answer** Almost but not quite. You are only entitled to say that \( y = x^3 + C \). The constant \( C \) is the starting value of \( y \) (when \( x = 0 \)). Then the **differential equation** \( dy/dx = 3x^2 \) is solved.

Every time you find a derivative, you can go backward to solve a differential equation. The function \( y = x^2 + x \) has the slope \( dy/dx = 2x + 1 \). In reverse, the slope \( 2x + 1 \) produces \( x^2 + x \)—and all the other functions \( x^2 + x + C \), shifted up and down. After going from distance \( f \) to velocity \( v \), we return to \( f + C \). But there is a lot more to differential equations. Here are two crucial points:

1. **We reach \( dy/dx \) by way of \( \Delta y/\Delta x \), but we have no system to go backward.** With \( dy/dx = (\sin x)/x \) we are lost. What function has this derivative?
2. **Many equations have the same solution \( y = x^3 \).** Economics has \( dy/dx = 3y/x \). Geometry has \( dy/dx = 3y^{2/3} \). These equations involve \( y \) as well as \( dy/dx \). Function and slope are mixed together! This is typical of differential equations.

**To summarize:** Chapters 2–4 compute and use derivatives. Chapter 5 goes in reverse. Integral calculus discovers the function from its slope. Given \( dy/dx \) we find \( y(x) \). Then Chapter 6 solves the differential equation \( dy/dt = y \), function mixed with slope. Calculus moves from derivatives to integrals to differential equations.

This discussion of the purpose of calculus should mention a specific example. Differential equations are applied to an epidemic (like AIDS). In most epidemics the number of cases grows exponentially. The peak is quickly reached by \( e^t \), and the epidemic dies down. Amazingly, exponential growth is not happening with AIDS—the best fit to the data through 1988 is a **cubic polynomial** (*Los Alamos Science, 1989*):

\[
y = 174.6(t - 1981.2)^3 + 340.
\]
This is dramatically different from other epidemics. Instead of \( \frac{dy}{dt} = y \) we have \( \frac{dy}{dt} = 3y/t \). Before this book is printed, we may know what has been preventing \( e' \) (fortunately). Eventually the curve will turn away from a cubic—I hope that mathematical models will lead to knowledge that saves lives.

Added in proof: In 1989 the curve for the U.S. dropped from \( t^3 \) to \( t^2 \).

**MARGINAL COST AND ELASTICITY IN ECONOMICS**

First point about economics: The **marginal** cost and **marginal** income are crucially important. The average cost of making automobiles may be $10,000. But it is the $8000 cost of the next car that decides whether Ford makes it. "The average describes the past, the marginal predicts the future." For bank deposits or work hours or wheat, which come in smaller units, the amounts are continuous variables. Then the word "marginal" says one thing: **Take the derivative.**

The average pay over all the hours we ever worked may be low. We wouldn't work another hour for that! This average is rising, but the pay for each additional hour rises faster—possibly it jumps. When $10/hour increases to $15/hour after a 40-hour week, a 50-hour week pays $550. The average income is $11/hour. The marginal income is $15/hour—the overtime rate.

Concentrate next on cost. Let \( y(x) \) be the cost of producing \( x \) tons of steel. The cost of \( x + \Delta x \) tons is \( y(x + \Delta x) \). The extra cost is the difference \( \Delta y \). Divide by \( \Delta x \), the number of extra tons. The ratio \( \Delta y/\Delta x \) is the **average cost per extra ton**. When \( \Delta x \) is an ounce instead of a ton, we are near the marginal cost \( dy/dx \).

Example: When the cost is \( x^2 \), the average cost is \( x^2/x = x \). The marginal cost is \( 2x \). Figure 2.4 has increasing slope—an example of "diminishing returns to scale."

![Fig. 2.4 Marginal exceeds average. Constant elasticity \( E = \pm 1 \). Perfectly elastic to perfectly inelastic (Γ curve).](image)

This raises another point about economics. The units are arbitrary. In yen per kilogram the numbers look different. The way to correct for arbitrary units is to work with **percentage change** or **relative change**. An increase of \( \Delta x \) tons is a relative increase of \( \Delta x/x \). A cost increase \( \Delta y \) is a relative increase of \( \Delta y/y \). Those are **dimensionless**, the same in tons/tons or dollars/dollars or yen/yen.

A third example is the demand \( y \) at price \( x \). Now \( dy/dx \) is negative. But again the units are arbitrary. The demand is in liters or gallons, the price is in dollars or pesos.

†These paragraphs show how calculus applies to economics. You do **not** have to be an economist to understand them. Certainly the author is not, probably the instructor is not, possibly the student is not. We can all use \( dy/dx \).
Relative changes are better. When the price goes up by 10%, the demand may drop by 5%. If that ratio stays the same for small increases, the elasticity of demand is $\frac{1}{2}$.

Actually this number should be $-\frac{1}{2}$. The price rose, the demand dropped. In our definition, the elasticity will be $-\frac{1}{2}$. In conversation between economists the minus sign is left out (I hope not forgotten).

**Definition**
The elasticity of the demand function $y(x)$ is

$$E(x) = \lim_{\Delta x \to 0} \frac{\Delta y}{y} \frac{y}{\Delta x} = \frac{dy/dx}{y/x}. \quad (3)$$

*Elasticity is “marginal” divided by “average.”* $E(x)$ is also relative change in $y$ divided by relative change in $x$. Sometimes $E(x)$ is the same at all prices—this important case is discussed below.

**Example 1**
Suppose the demand is $y = c/x$ when the price is $x$. The derivative $dy/dx = -c/x^2$ comes from calculus. The division $y/x = c/x^2$ is only algebra. *The ratio is $E = -1$.*

*For the demand $y = c/x$, the elasticity is $(-c/x^2)/(c/x^2) = -1$.*

All demand curves are compared with this one. The demand is inelastic when $|E| < 1$. It is elastic when $|E| > 1$. The demand $20/\sqrt{x}$ is inelastic ($E = -\frac{1}{2}$), while $x^{-3}$ is elastic ($E = -3$). *The power $y = cx^n$, whose derivative we know, is the function with constant elasticity $n$:*

if $y = cx^n$ then $dy/dx = cnx^{n-1}$ and $E = cnx^{-1}/(cx^n/x) = n$.

It is because $y = cx^n$ sets the standard that we could come so early to economics.

In the special case when $y = c/x$, consumers spend the same at all prices. Price $x$ times quantity $y$ remains constant at $xy = c$.

**Example 2**
The supply curve has $E > 0$—supply increases with price. Now the baseline case is $y = cx$. The slope is $c$ and the average is $y/x = c$. *The elasticity is $E = c/c = 1$.***

Compare $E = 1$ with $E = 0$ and $E = \infty$. A constant supply is “perfectly inelastic.”

The power $n$ is zero and the slope is zero: $y = c$. No more is available when the harvest is over. Whatever the price, the farmer cannot suddenly grow more wheat. Lack of elasticity makes farm economics difficult.

The other extreme $E = \infty$ is “perfectly elastic.” The supply is unlimited at a fixed price $x$. Once this seemed true of water and timber. In reality the steep curve $x = constant$ is leveling off to a flat curve $y = constant$. Fixed price is changing to fixed supply, $E = \infty$ is becoming $E = 0$, and the supply of water follows a “gamma curve” shaped like $\Gamma$.

**Example 3**
Demand is an increasing function of *income*—more income, more demand. The *income elasticity* is $E(I) = (dy/dI)/(y/I)$. A luxury has $E > 1$ (elastic). Doubling your income more than doubles the demand for caviar. A necessity has $E < 1$ (inelastic). The demand for bread does not double. Please recognize how the central ideas of calculus provide a language for the central ideas of economics.

*Important note on supply = demand* This is the basic equation of microeconomics. Where the supply curve meets the demand curve, the economy finds the equilibrium price. *Supply = demand assumes perfect competition.* With many suppliers, no one can raise the price. If someone tries, the customers go elsewhere.
2 Derivatives

The opposite case is a monopoly—no competition. Instead of many small producers of wheat, there is one producer of electricity. An airport is a monopolist (and maybe the National Football League). If the price is raised, some demand remains.

Price fixing occurs when several producers act like a monopoly—which antitrust laws try to prevent. The price is not set by supply = demand. The calculus problem is different—to maximize profit. Section 3.2 locates the maximum where the marginal profit (the slope!) is zero.

Question on income elasticity From an income of $10,000 you save $500. The income elasticity of savings is $E = 2$. Out of the next dollar what fraction do you save?

Answer The savings is $y = cx^2$ because $E = 2$. The number $c$ must give $500 = \frac{1}{10,000}$, so $c$ is $5 \cdot 10^{-6}$. Then the slope $dy/dx$ is $2cx = 10 \cdot 10^{-6} \cdot 10^4 = \frac{1}{16}$. This is the marginal savings, ten cents on the dollar. Average savings is $5\%$, marginal savings is $10\%$, and $E = 2$.

2.2 EXERCISES

Read-through questions

The derivative of $f = x^4$ is $f' = \underline{a}$. That comes from expanding $(x + h)^4$ into the five terms $\underline{b}$. Subtracting $x^4$ and dividing by $h$ leaves the four terms $\underline{c}$. This is $\Delta f / h$, and its limit is $\underline{d}$.

The derivative of $f = x^n$ is $f' = \underline{e}$. Now $(x + h)^n$ comes from the $\underline{t}$ theorem. The terms to look for are are $x^{n-1}h$, containing only one $\underline{g}$. There are $\underline{h}$ of those terms, so $(x + h)^n = x^n + \underline{i} + \cdots$. After subtracting $\underline{j}$ and dividing by $h$, the limit of $\Delta f / h$ is $\underline{k}$. The coefficient of $x^{n-1}h$, not needed here, is "$n$ choose $\underline{m}$", where $n!$ means $\underline{n}$.

The derivative of $x^{-2}$ is $\underline{n}$. The derivative of $x^{1/2}$ is $\underline{o}$. The derivative of $3x + (1/x)$ is $\underline{p}$, which uses the following rules: The derivative of $3f(x)$ is $\underline{q}$ and the derivative of $f(x) + g(x)$ is $\underline{r}$. Integral calculus recovers $\underline{s}$ from $dy/dx$. If $dy/dx = x^2$ then $y(x) = \underline{t}$.

1 Starting with $f = x^6$, write down $f'$ and then $f''$. (This is "$f$ double prime," the derivative of $f'$.) After $\underline{u}$ derivatives of $x^6$ you reach a constant. What constant?

2 Find a function that has $x^6$ as its derivative.

Find the derivatives of the functions in 3–10. Even if $n$ is negative or a fraction, the derivative of $x^n$ is $nx^{n-1}$.

3 $x^2 + 7x + 5$ 4 $1 + (7/x) + (5/x^2)$
5 $1 + x^2 + x^3 + x^4$ 6 $(x^2 + 1)^2$
7 $x^n + x^{-n}$ 8 $x^n/n!$
9 $1 + x + \frac{1}{2}x^2 + \frac{1}{6}x^3 + \frac{1}{24}x^4$ 10 $\frac{2}{3}x^{3/2} + \frac{2}{5}x^{5/2}$

11 Name two functions with $df/dx = 1/x^2$.

12 Find the mistake: $x^2$ is $x + x + \cdots + x$ (with $x$ terms). Its derivative is $1 + 1 + \cdots + 1$ (also $x$ terms). So the derivative of $x^2$ seems to be $x$.

13 What are the derivatives of $3x^{1/3}$ and $-3x^{-1/3}$ and $(3x^{1/3})^{-1}$?

14 The slope of $x + (1/x)$ is zero when $x = \underline{i}$. What does the graph do at that point?

15 Draw a graph of $y = x^3 - x$. Where is the slope zero?

16 If $df/dx$ is negative, is $f(x)$ always negative? Is $f(x)$ negative for large $x$? If you think otherwise, give examples.

17 A rock thrown upward with velocity 16 ft/sec reaches height $f = 16t - 16t^2$ at time $t$.

(a) Find its average speed $\Delta f / \Delta t$ from $t = 0$ to $t = \frac{1}{2}$.
(b) Find its average speed $\Delta f / \Delta t$ from $t = \frac{1}{2}$ to $t = 1$.
(c) What is $df/dt$ at $t = \frac{1}{2}$?

18 When $f$ is in feet and $t$ is in seconds, what are the units of $f'$ and its derivative $f''$? In $f = 16t - 16t^2$, the first 16 is ft/sec but the second 16 is $\underline{u}$.

19 Graph $y = x^3 + x^2 - x$ from $x = -2$ to $x = 2$ and estimate where it is decreasing. Check the transition points by solving $dy/dx = 0$.

20 At a point where $dy/dx = 0$, what is special about the graph of $y(x)$? Test case: $y = x^2$.

21 Find the slope of $y = \sqrt{x}$ by algebra (then $h \to 0$):

$$\frac{\Delta y}{h} = \frac{\sqrt{x+h} - \sqrt{x}}{h} = \frac{\sqrt{x+h} - \sqrt{x}}{h} \cdot \frac{\sqrt{x+h} + \sqrt{x}}{\sqrt{x+h} + \sqrt{x}}.$$

22 Imitate Problem 21 to find the slope of $y = 1/\sqrt{x}$. 

2.2 Powers and Polynomials

23 Complete Pascal's triangle for \( n = 5 \) and \( n = 6 \). Why do the numbers across each row add to \( 2^n \)?

24 Complete \( (x + h)^2 = x^2 + \_ \_ \_ \). What are the binomial coefficients \( \binom{5}{1} \) and \( \binom{5}{2} \) and \( \binom{5}{3} \)?

25 Compute \( (x + h)^3 - (x - h)^3 \), divide by \( 2h \), and set \( h = 0 \). Why divide by \( 2h \) to find this slope?

26 Solve the differential equation \( y'' = x \) to find \( y(x) \).

27 For \( f(x) = x^2 + x^3 \), write out \( f(x + Ax) \) and \( df/\Delta x \). What is the limit at \( \Delta x = 0 \) and what rule about sums is confirmed?

28 The derivative of \( (u(x))^2 \) is \( \_ \_ \_ \_ \) from Section 2.1. Test this rule on \( u = x^3 \).

29 What are the derivatives of \( x^7 + 1 \) and \( (x + 1)^7 \)? Shift the graph of \( x^7 \).

30 If \( df/dx \) is \( v(x) \), what functions have these derivatives?
   (a) \( 4e(x) \)
   (b) \( v(x) + 1 \)
   (c) \( v(x + 1) \)
   (d) \( v(x) + v'(x) \).

31 What function \( f(x) \) has fourth derivative equal to 1?

32 What function \( f(x) \) has \( n \)th derivative equal to 1?

33 Suppose \( df/dx = 1 + x + x^2 + x^3 \). Find \( f(x) \).

34 Suppose \( df/dx = x^{-2} - x^{-3} \). Find \( f(x) \).

35 \( f(x) \) can be its own derivative. In the infinite polynomial \( f = 1 + x + x^2 + x^3 + \_ \_ \_ \_ \_ \_ \), what numbers multiply \( x^4 \) and \( x^5 \) if \( df/dx \) equals \( f \)?

36 Write down a differential equation \( dy/dx = \_ \_ \_ \) that is solved by \( y = x^2 \). Make the right side involve \( y \) (not just 2x).

37 True or false: (a) The derivative of \( x^n \) is \( nx^n \).
   (b) The derivative of \( ax^n/bx^n \) is \( a/b \).
   (c) If \( df/dx = x^4 \) and \( dx/dx = x^4 \) then \( f(x) = g(x) \).
   (d) \( (f(x) - f(a))/(x - a) \) approaches \( f'(a) \) as \( x \to a \).
   (e) The slope of \( y = (x - 1)^2 \) is \( y' = 3(x - 1)^2 \).

Problems 38–44 are about calculus in economics.

38 When the cost is \( y = y_0 + cx \), find \( E(x) = (dy/dx)/(y/x) \). It approaches \_ \_ \_ \_ \_ \_ \_ \_ for large \( x \).

39 From an income of \( x = 10,000 \), you spend \( y = 1200 \) on your car. If \( E = 1/2 \), what fraction of your next dollar will be spent on the car? Compare \( dy/dx \) (marginal) with \( y/x \) (average).

40 Name a product whose price elasticity is
   (a) high
   (b) low
   (c) negative (?)

41 The demand \( y = c/x \) has \( dy/dx = -y/x \). Show that \( \Delta y/\Delta x \) is \( \not{= -y/x} \). (Use numbers or algebra.) Finite steps miss the special feature of infinitesimal steps.

42 The demand \( y = x^n \) has \( E = \_ \_ \_ \_ \). The revenue \( xy \) (price times demand) has elasticity \( E = \_ \_ \_ \).

43 \( y = 2x + 3 \) grows with marginal cost 2 from the fixed cost 3. Draw the graph of \( E(x) \).

44 From an income \( I \) we save \( S(I) \). The marginal propensity to save is \_ \_ \_ \_ \_. Elasticity is not needed because \( S \) and \( I \) have the same \_ \_ \_ \_. Applied to the whole economy this is \_ \_ \_.

45 \( 2^t \) is doubled when \( t \) increases by \_ \_ \_. \( t^3 \) is doubled when \( t \) increases to \_ \_ \_ \_ \_. The doubling time for AIDS is proportional to \( t \).

46 Biology also leads to \( dy/y = n \ dx/x \), for the relative growth of the head \( (dy/y) \) and the body \( (dx/x) \). Is \( n > 1 \) or \( n < 1 \) for a child?

47 What functions have \( df/dx = x^9 \) and \( df/dx = x^n \)? Why does \( n = -1 \) give trouble?

48 The slope of \( y = x^3 \) comes from this identity:
   \[
   \frac{(x + h)^3 - x^3}{h} = (x + h)^2 + (x + h)x + x^2. 
   \]
   (a) Check the algebra. Find \( dy/dx \) as \( h \to 0 \).
   (b) Write a similar identity for \( y = x^4 \).

49 (Computer graphing) Find all the points where \( y = x^4 + 2x^3 - 7x^2 + 3 = 0 \) and where \( dy/dx = 0 \).

50 The graphs of \( y_1(x) = x^4 + x^2 \) and \( y_2(x) = 7x - 5 \) touch at the point where \( y_3(x) = \_ \_ \_ \_ \_ \_ \_ = 0 \). Plot \( y_3(x) \) to see what is special. What does the graph of \( y(x) \) do at a point where \( y = y' = 0 \)?

51 In the Massachusetts lottery you choose 6 numbers out of 36. What is your chance to win?

52 In what circumstances would it pay to buy a lottery ticket for every possible combination, so one of the tickets would win?
Chapter 1 started with straight line graphs. The velocity was constant (at least piecewise). The distance function was linear. Now we are facing polynomials like $x^3 - 2$ or $x^4 - x^2 + 3$, with other functions to come soon. Their graphs are definitely curved. Most functions are not close to linear—except if you focus all your attention near a single point. That is what we will do.

*Over a very short range a curve looks straight.* Look through a microscope, or zoom in with a computer, and there is no doubt. The graph of distance versus time becomes nearly linear. Its slope is the velocity at that moment. We want to find the line that the graph stays closest to—the "tangent line"—before it curves away.

The tangent line is easy to describe. We are at a particular point on the graph of $y = f(x)$. At that point $x$ equals $a$ and $y$ equals $f(a)$ and the slope equals $f'(a)$. The tangent line goes through that point $x = a$, $y = f(a)$ with that slope $m = f'(a)$.

Figure 2.5 shows the line more clearly than any equation, but we have to turn the geometry into algebra. We need the equation of the line.

**EXAMPLE 1** Suppose $y = x^4 - x^2 + 3$. At the point $x = a = 1$, the height is $y = f(a) = 3$. The slope is $dy/dx = 4x^3 - 2x$. At $x = 1$ the slope is $4 - 2 = 2$. That is $f'(a)$:

The numbers $x = 1$, $y = 3$, $dy/dx = 2$ determine the tangent line.
The equation of the tangent line is $y - 3 = 2(x - 1)$, and this section explains why.

**THE EQUATION OF A LINE**

A straight line is determined by two conditions. We know the line if we know two of its points. (We still have to write down the equation.) Also, if we know one point and the slope, the line is set. That is the situation for the tangent line, which has a known slope at a known point:

1. The equation of a line has the form $y = mx + b$
2. The number $m$ is the slope of the line, because $dy/dx = m$
3. The number $b$ adjusts the line to go through the required point.

I will take those one at a time—first $y = mx + b$, then $m$, then $b$.

1. The graph of $y = mx + b$ is not curved. How do we know? For the specific example $y = 2x + 1$, take two points whose coordinates $x$, $y$ satisfy the equation:

   $x = 0$, $y = 1$ and $x = 4$, $y = 9$ both satisfy $y = 2x + 1$. 

Fig. 2.5 The tangent line has the same slope 2 as the curve (especially after zoom).
Those points (0, 1) and (4, 9) lie on the graph. The point halfway between has \(x = 2\) and \(y = 5\). That point also satisfies \(y = 2x + 1\). The halfway point is on the graph. If we subdivide again, the midpoint between (0, 1) and (2, 5) is (1, 3). This also has \(y = 2x + 1\). The graph contains all halfway points and must be straight.

2. What is the correct slope \(m\) for the tangent line? In our example it is \(m = f'(a) = 2\). The curve and its tangent line have the same slope at the crucial point: \(dy/dx = 2\).

Allow me to say in another way why the line \(y = mx + b\) has slope \(m\). At \(x = 0\) its height is \(y = b\). At \(x = 1\) its height is \(y = m + b\). The graph has gone one unit across (0 to 1) and \(m\) units up (\(b\) to \(m + b\)). The whole idea is

\[
slope = \frac{\text{distance up}}{\text{distance across}} = \frac{m}{1}. \tag{1}
\]

Each unit across means \(m\) units up, to \(2m + b\) or \(3m + b\). A straight line keeps a constant slope, whereas the slope of \(y = x^4 - x^2 + 3\) equals 2 only at \(x = 1\).

3. Finally we decide on \(b\). The tangent line \(y = 2x + b\) must go through \(x = 1, y = 3\). Therefore \(b = 1\). With letters instead of numbers, \(y = mx + b\) leads to \(f(a) = ma + b\). So we know \(b\):

\[
2E \quad \text{The equation of the tangent line has } b = f(a) - ma:
\]

\[
y = mx + f(a) - ma \quad \text{or} \quad y - f(a) = m(x - a). \tag{2}
\]

That last form is the best. You see immediately what happens at \(x = a\). The factor \(x - a\) is zero. Therefore \(y = f(a)\) as required. This is the point-slope form of the equation, and we use it constantly:

\[
y - 3 = 2(x - 1) \quad \text{or} \quad \frac{y - 3}{x - 1} = \frac{\text{distance up}}{\text{distance across}} = \text{slope } 2.
\]

**EXAMPLE 2** The curve \(y = x^3 - 2\) goes through \(y = 6\) when \(x = 2\). At that point \(dy/dx = 3x^2 = 12\). The point-slope equation of the tangent line uses 2 and 6 and 12:

\[
y - 6 = 12(x - 2), \quad \text{which is also} \quad y = 12x - 18.
\]

There is another important line. It is perpendicular to the tangent line and perpendicular to the curve. This is the normal line in Figure 2.6. Its new feature is its slope. When the tangent line has slope \(m\), the normal line has slope \(-1/m\). (Rule: Slopes of

![Fig. 2.6](image_url) Tangent line \(y - y_0 = m(x - x_0)\). Normal line \(y - y_0 = -\frac{1}{m} (x - x_0)\). Leaving a roller-coaster and catching up to a car.
2 Derivatives

perpendicular lines multiply to give \(-1\).) Example 2 has \(m = 12\), so the normal line has slope \(-\frac{1}{12}\):

\[
tangent \ line: \ y - 6 = 12(x - 2) \quad normal \ line: \ y - 6 = -\frac{1}{12}(x - 2).
\]

Light rays travel in the normal direction. So do brush fires—they move perpendicular to the fire line. Use the point-slope form! The tangent is \(y = 12x - 18\), the normal is not \(y = -\frac{1}{12}x - 18\).

**EXAMPLE 3** You are on a roller-coaster whose track follows \(y = x^2 + 4\). You see a friend at \((0, 0)\) and want to get there quickly. Where do you step off?

**Solution** Your path will be the tangent line (at high speed). The problem is to choose \(x = a\) so the tangent line passes through \(x = 0, y = 0\). When you step off at \(x = a\),

- the height is \(y = a^2 + 4\) and the slope is \(2a\)
- the equation of the tangent line is \(y - (a^2 + 4) = 2a(x - a)\)
- this line goes through \((0, 0)\) if \(- (a^2 + 4) = -2a^2\) or \(a = \pm 2\).

The same problem is solved by spacecraft controllers and baseball pitchers. Releasing a ball at the right time to hit a target 60 feet away is an amazing display of calculus. Quarterbacks with a moving target should read Chapter 4 on related rates.

Here is a better example than a roller-coaster. Stopping at a red light wastes gas. It is smarter to slow down early, and then accelerate. When a car is waiting in front of you, the timing needs calculus:

**EXAMPLE 4** How much must you slow down when a red light is 72 meters away? In 4 seconds it will be green. The waiting car will accelerate at 3 meters/sec². You cannot pass the car.

**Strategy** Slow down immediately to the speed \(V\) at which you will just catch that car. (If you wait and brake later, your speed will have to go below \(V\).) At the catch-up time \(T\), the cars have the same speed and same distance. Two conditions, so the distance functions in Figure 2.6d are tangent.

**Solution** At time \(T\), the other car’s speed is \(3(T - 4)\). That shows the delay of 4 seconds. Speeds are equal when \(3(T - 4) = V\) or \(T = \frac{1}{3}V + 4\). Now require equal distances. Your distance is \(V\) times \(T\). The other car’s distance is \(72 + \frac{1}{2}at^2\):

\[
72 + \frac{1}{2} \cdot 3(T - 4)^2 = VT \quad becomes \quad 72 + \frac{1}{2} \cdot \frac{1}{3}V^2 = V(\frac{1}{3}V + 4).
\]

The solution is \(V = 12\) meters/second. This is 43 km/hr or 27 miles per hour.

Without the other car, you only slow down to \(V = 72/4 = 18\) meters/second. As the light turns green, you go through at 65 km/hr or 40 miles per hour. Try it.

**THE SECANT LINE CONNECTING TWO POINTS ON A CURVE**

Instead of the tangent line through one point, consider the secant line through two points. For the tangent line the points came together. Now spread them apart. The point-slope form of a linear equation is replaced by the two-point form.

The equation of the curve is still \(y = f(x)\). The first point remains at \(x = a, y = f(a)\). The other point is at \(x = c, y = f(c)\). The secant line goes between them, and we want its equation. This time we don’t start with the slope—but \(m\) is easy to find.
2.3 The Slope and the Tangent Line

EXAMPLE 5 The curve \( y = x^3 - 2 \) goes through \( x = 2, y = 6 \). It also goes through \( x = 3, y = 25 \). The slope between those points is

\[
m = \frac{\text{change in } y}{\text{change in } x} = \frac{25 - 6}{3 - 2} = 19.
\]

The point-slope form (at the first point) is \( y - 6 = 19(x - 2) \). This line automatically goes through the second point \( (3, 25) \). Check: \( 25 - 6 \) equals \( 19(3 - 2) \). The secant has the right slope 19 to reach the second point. It is the average slope \( \Delta y/\Delta x \).

A look ahead The second point is going to approach the first point. The secant slope \( \Delta y/\Delta x \) will approach the tangent slope \( dy/dx \). We discover the derivative (in the limit). That is the main point now—but not forever.

Soon you will be fast at derivatives. The exact \( dy/dx \) will be much easier than \( \Delta y/\Delta x \). The situation is turned around as soon as you know that \( x^9 \) has slope \( 9x^8 \). Near \( x = 1 \), the distance up is about 9 times the distance across. To find \( \Delta y = 1.001^9 - 1^9 \), just multiply \( \Delta x = .001 \) by 9. The quick approximation is .009, the calculator gives \( \Delta y = .009036 \). It is easier to follow the tangent line than the curve.

Come back to the secant line, and change numbers to letters. What line connects \( x = a, y = f(a) \) to \( x = c, y = f(c) \)? A mathematician puts formulas ahead of numbers, and reasoning ahead of formulas, and ideas ahead of reasoning:

(1) The slope is \( m = \frac{\text{distance up}}{\text{distance across}} = \frac{f(c) - f(a)}{c - a} \)

(2) The height is \( y = f(a) \) at \( x = a \)

(3) The height is \( y = f(c) \) at \( x = c \) (automatic with correct slope).

The two-point form uses the slope between the points:

\[
\text{secant line: } y - f(a) = \left( \frac{f(c) - f(a)}{c - a} \right) (x - a).
\]

At \( x = a \) the right side is zero. So \( y = f(a) \) on the left side. At \( x = c \) the right side has two factors \( c - a \). They cancel to leave \( y = f(c) \). With equation (2) for the tangent line and equation (3) for the secant line, we are ready for the moment of truth.

THE SECANT LINE APPROACHES THE TANGENT LINE

What comes now is pretty basic. It matches what we did with velocities:

\[
\text{average velocity} = \frac{\text{\Delta distance}}{\text{\Delta time}} = \frac{f(t + \Delta t) - f(t)}{\Delta t}.
\]

The limit is \( df/dt \). We now do exactly the same thing with slopes. The secant line turns into the tangent line as \( c \) approaches \( a \):

slope of secant line: \( \frac{\Delta f}{\Delta x} = \frac{f(c) - f(a)}{c - a} \)

slope of tangent line: \( \frac{df}{dx} = \text{limit of } \frac{\Delta f}{\Delta x} \).
2 Derivatives

There stands the fundamental idea of differential calculus! You have to imagine more secant lines than I can draw in Figure 2.7, as $c$ comes close to $a$. Everybody recognizes $c - a$ as $\Delta x$. Do you recognize $f(c) - f(a)$ as $f(x + \Delta x) - f(x)$? It is $\Delta f$, the change in height. All lines go through $x = a, y = f(a)$. Their limit is the tangent line.

$$
\text{secant } y - f(a) = \frac{f(c) - f(a)}{c - a} (x - a)
$$

$$
\text{tangent } y - f(a) = f'(a)(x - a)
$$

Fig. 2.7 Secants approach tangent as their slopes $\Delta f/\Delta x$ approach $df/dx$.

Intuitively, the limit is pretty clear. The two points come together, and the tangent line touches the curve at one point. (It could touch again at faraway points.) Mathematically this limit can be tricky—it takes us from algebra to calculus. Algebra stays away from 0/0, but calculus gets as close as it can.

The new limit for $df/dx$ looks different, but it is the same as before:

$$f'(a) = \lim_{c \to a} \frac{f(c) - f(a)}{c - a}.
$$

EXAMPLE 6 Find the secant lines and tangent line for $y = f(x) = \sin x$ at $x = 0$.

The starting point is $x = 0, y = \sin 0$. This is the origin (0, 0). The ratio of distance up to distance across is $(\sin c)/c$:

$$\text{secant equation } y = \frac{\sin c}{c} x \quad \text{tangent equation } y = 1x.$$

As $c$ approaches zero, the secant line becomes the tangent line. The limit of $(\sin c)/c$ is not 0/0, which is meaningless, but 1, which is $dy/dx$.

EXAMPLE 7 The gold you own will be worth $\sqrt{t}$ million dollars in $t$ years. When does the rate of increase drop to 10% of the current value, so you should sell the gold and buy a bond? At $t = 25$, how far does that put you ahead of $\& = 5$?

Solution The rate of increase is the derivative of $\sqrt{t}$, which is $1/2\sqrt{t}$. That is 10% of the current value $\sqrt{t}$ when $1/2\sqrt{t} = \sqrt{t}/10$. Therefore $2t = 10$ or $t = 5$. At that time you sell the gold, leave the curve, and go onto the tangent line:

$$y - \sqrt{5} = \frac{\sqrt{5}}{10} (t - 5) \quad \text{becomes } \quad y - \sqrt{5} = 2\sqrt{5} \quad \text{at } \quad t = 25.$$

With straight interest on the bond, not compounded, you have reached $y = 3\sqrt{5} = 6.7$ million dollars. The gold is worth a measly five million.

2.3 Exercises

Read-through questions

A straight line is determined by one point and the ___b___. The slope of the tangent line equals the slope of the ____c____. The point-slope form of the tangent equation is $y - f(a) = ____d____$. The tangent line to $y = x^3 + x$ at $x = 1$ has slope ____e____. Its
equation is \( y - 2 = \frac{1}{m} (x - 1) \). It crosses the \( y \) axis at \( \frac{2}{m} \) and the \( x \) axis at \( \frac{1}{m} \). The normal line at this point \((1, 2)\) has slope \( -m \). Its equation is \( y - 2 = -m (x - 1) \). The secant line from \((1, 2)\) to \((2, k)\) has slope \( \frac{1}{m} \). Its equation is \( y - 2 = \frac{1}{m} (x - 1) \).

The point \((c, f(c))\) is on the line \( y - f(a) = m(x - a) \) provided \( m = \frac{f(c) - f(a)}{c - a} \). As \( c \) approaches \( a \), the slope \( m \) approaches \( f'(a) \). The secant line approaches the \( \text{p} \) line.

1. (a) Find the slope of \( y = 12/x \).
   (b) Find the equation of the tangent line at \((2, 6)\).
   (c) Find the equation of the normal line at \((2, 6)\).
   (d) Find the equation of the secant line to \((4, 3)\).

2. For \( y = x^2 + x \) find equations for
   (a) the tangent line and normal line at \((1, 2)\);
   (b) the secant line to \( x = 1 + h, y = (1 + h)^2 + (1 + h) \).

3. A line goes through \((1, -1)\) and \((4, 8)\). Write its equation in point-slope form. Then write it as \( y = mx + b \).

4. The tangent line to \( y = x^3 + 6x \) at the origin is \( y = \) ________. Does it cross the curve again?

5. The tangent line to \( y = x^3 - 3x^2 + x \) at the origin is \( y = \) ________. It is also the secant line to the point ________.

6. Find the tangent line to \( x = y^2 \) at \( x = 4, y = 2 \).

7. For \( y = x^2 \) the secant line from \((a, a^2)\) to \((c, c^2)\) has the equation ________. Do the division by \( c - a \) to find the tangent line as \( c \) approaches \( a \).

8. Construct a function that has the same slope at \( x = 1 \) and \( x = 2 \). Then find two points where \( y = x^4 - 2x^2 \) has the same tangent line (draw the graph).

9. Find a curve that is tangent to \( y = 2x - 3 \) at \( x = 5 \). Find the normal line to that curve at \((5, 7)\).

10. For \( y = 1/x \) the secant line from \((a, 1/a)\) to \((c, 1/c)\) has the equation ________. Simplify its slope and find the limit as \( c \) approaches \( a \).

11. What are the equations of the tangent line and normal line to \( y = \sin x \) at \( x = \pi/2 \)?

12. If \( c \) and \( a \) both approach an in-between value \( x = b \), then the secant slope \( (f(c) - f(a))/(c - a) \) approaches ________.

13. At \( x = a \) on the graph of \( y = 1/x \), compute
    (a) the equation of the tangent line
    (b) the points where that line crosses the axes.

The triangle between the tangent line and the axes always has area ________.

14. Suppose \( g(x) = f(x) + 7 \). The tangent lines to \( f \) and \( g \) at \( x = 4 \) are ________. True or false: The distance between those lines is 7.

15. Choose \( c \) so that \( y = 4x \) is tangent to \( y = x^2 + c \). Match heights as well as slopes.

16. Choose \( c \) so that \( y = 5x - 7 \) is tangent to \( y = x^2 + cx \).

17. For \( y = x^3 + 4x^2 - 3x + 1 \), find all points where the tangent is horizontal.

18. \( y = 4x \) can’t be tangent to \( y = cx^2 \). Try to match heights and slopes, or draw the curves.

19. Determine \( c \) so that the straight line joining \((0, 3)\) and \((5, -2)\) is tangent to the curve \( y = e/(x + 1) \).

20. Choose \( b, c, d \) so that the two parabolas \( y = x^2 + bx + c \) and \( y = dx - x^2 \) are tangent to each other at \( x = 1, y = 0 \).

21. The graph of \( f(x) = x^3 \) goes through \((1, 1)\).
    (a) Another point is \( x = c = 1 + h, y = f(c) = \) ________.
    (b) The change in \( f \) is \( \Delta f = \) ________.
    (c) The slope of the secant is \( m = \) ________.
    (d) As \( h \) goes to zero, \( m \) approaches ________.

22. Construct a function \( y = f(x) \) whose tangent line at \( x = 1 \) is the same as the secant that meets the curve again at \( x = 3 \).

23. Draw two curves bending away from each other. Mark the points \( P \) and \( Q \) where the curves are closest. At those points, the tangent lines are ________ and the normal lines are ________.

*24. If the parabolas \( y = x^2 + 1 \) and \( y = x - x^2 \) come closest at \( (a, a^2 + 1) \) and \( (c, c - c^2) \), set up two equations for \( a \) and \( c \).

25. A light ray comes down the line \( x = a \). It hits the parabolic reflector \( y = x^2 \) at \( P = (a, a^2) \).
    (a) Find the tangent line at \( P \). Locate the point \( Q \) where that line crosses the \( y \) axis.
    (b) Check that \( P \) and \( Q \) are the same distance from the focus at \( F = (0, \frac{1}{4}) \).
    (c) Show from (b) that the figure has equal angles.
    (d) What law of physics makes every ray reflect off the parabola to the focus at \( F \)?

26. In a bad reflector \( y = 2/x \), a ray down one special line \( x = a \) is reflected horizontally. What is \( a \)?
27 For the parabola $4py = x^2$, where is the slope equal to 1? At that point a vertical ray will reflect horizontally. So the focus is at $(0, \ldots)$. 

28 Why are these statements wrong? Make them right. 
   (a) If $y = 2x$ is the tangent line at $(1, 2)$ then $y = -\frac{1}{x}$ is the normal line. 
   (b) As $c$ approaches $a$, the secant slope $(f(c) - f(a))/(c - a)$ approaches $(f(a) - f(a))/(a - a)$. 
   (c) The line through $(2, 3)$ with slope 4 is $y - 2 = 4(x - 3)$. 

29 A ball goes around a circle: $x = \cos t$, $y = \sin t$. At $t = 344$ the ball flies off on the tangent line. Find the equation of that line and the point where the ball hits the ground ($y = 0$). 

30 If the tangent line to $y = f(x)$ at $x = a$ is the same as the tangent line to $y = g(x)$ at $x = b$, find two equations that must be satisfied by $a$ and $b$. 

31 Draw a circle of radius 1 resting in the parabola $y = x^2$. At the touching point $(a, a^2)$, the equation of the normal line is $\ldots$. That line has $x = 0$ when $y = \ldots$. The distance to $(a, a^2)$ equals the radius 1 when $a = \ldots$. This locates the touching point. 

32 Follow Problem 31 for the flatter parabola $y = \frac{1}{2}x^2$ and explain where the circle rests. 

33 You are applying for a $1000 scholarship and your time is worth $10 per hour. If the chance of success is $1 - (1/x)$ from $x$ hours of writing, when should you stop? 

34 Suppose $|f(c) - f(a)| \leq |c - a|$ for every pair of points $a$ and $c$. Prove that $|df/dx| \leq 1$. 

35 From which point $x = a$ does the tangent line to $y = 1/x^2$ hit the $x$ axis at $x = 3$? 

36 If $u(x)/v(x) = 7$ find $u'(x)/v'(x)$. Also find $(u(x)/v(x))'$. 

37 Find $f(c) = 1.001^{10}$ in two ways—by calculator and by $f(c) - f(a) \approx f'(a)(c - a)$. Choose $a = 1$ and $f(x) = x^{10}$. 

38 At a distance $\Delta x$ from $x = 1$, how far is the curve $y = 1/x$ above its tangent line? 

39 At a distance $\Delta x$ from $x = 2$, how far is the curve $y = x^3$ above its tangent line? 

40 Based on Problem 38 or 39, the distance between curve and tangent line grows like what power $(\Delta x)^p$? 

41 The tangent line to $f(x) = x^2 - 1$ at $x_0 = 2$ crosses the $x$ axis at $x_1 = \ldots$. The tangent line at $x_1$ crosses the $x$ axis at $x_2 = \ldots$. Draw the curve and the two lines, which are the beginning of Newton's method to solve $f(x) = 0$. 

42 (Puzzle) The equation $y = mx + b$ requires two numbers, the point-slope form $y - f(a) = f'(a)(x - a)$ requires three, and the two-point form requires four: $a$, $f(a)$, $c$, $f(c)$. How can this be? 

43 Find the time $T$ at the tangent point in Example 4, when you catch the car in front. 

44 If the waiting car only accelerates at 2 meters/sec$^2$, what speed $V$ must you slow down to? 

45 A thief 40 meters away runs toward you at 8 meters per second. What is the smallest acceleration so that $v = at$ keeps you in front? 

46 With 8 meters to go in a relay race, you slow down badly ($f = -8 + 6t - \frac{1}{2}t^2$). How fast should the next runner start (choose $v$ in $f = vt$) so you can just pass the baton? 

2.4 The Derivative of the Sine and Cosine

This section does two things. One is to compute the derivatives of $\sin x$ and $\cos x$. The other is to explain why these functions are so important. They describe oscillation, which will be expressed in words and equations. You will see a “differential equation.” It involves the derivative of an unknown function $y(x)$. 

The differential equation will say that the second derivative—the derivative of the derivative—is equal and opposite to $y$. In symbols this is $y'' = -y$. Distance in one direction leads to acceleration in the other direction. That makes $y$ and $y'$ and $y''$ all oscillate. The solutions to $y'' = -y$ are $\sin x$ and $\cos x$ and all their combinations. 

We begin with the slope. The derivative of $y = \sin x$ is $y' = \cos x$. There is no reason for that to be a mystery, but I still find it beautiful. Chapter 1 followed a ball around a circle; the shadow went up and down. Its height was $\sin t$ and its velocity was $\cos t$. 
2.4 The Derivative of the Sine and Cosine

We now find that derivative by the standard method of limits, when \( y(x) = \sin x \):

\[
\frac{dy}{dx} = \lim_{h \to 0} \frac{\Delta y}{\Delta x} = \lim_{h \to 0} \frac{\sin(x + h) - \sin x}{h}.
\]

(1)

The sine is harder to work with than \( x^2 \) or \( x^3 \). Where we had \((x + h)^2\) or \((x + h)^3\), we now have \(\sin(x + h)\). This calls for one of the basic “addition formulas” from trigonometry, reviewed in Section 1.5:

\[
\sin(x + h) = \sin x \cos h + \cos x \sin h \quad \text{(2)}
\]

\[
\cos(x + h) = \cos x \cos h - \sin x \sin h. \quad \text{(3)}
\]

Equation (2) puts \( \Delta y = \sin(x + h) - \sin x \) in a new form:

\[
\frac{\Delta y}{\Delta x} = \frac{\sin x \cos h + \cos x \sin h - \sin x}{h} = \sin x \left( \frac{\cos h - 1}{h} \right) + \cos x \left( \frac{\sin h}{h} \right). \quad \text{(4)}
\]

The ratio splits into two simpler pieces on the right. Algebra and trigonometry got us this far, and now comes the calculus problem. What happens as \( h \to 0 \)? It is no longer easy to divide by \( h \). (I will not even mention the unspeakable crime of writing \( \frac{\sin h}{h} = \sin \).) There are two critically important limits—the first is zero and the second is one:

\[
\lim_{h \to 0} \frac{\cos h - 1}{h} = 0 \quad \text{and} \quad \lim_{h \to 0} \frac{\sin h}{h} = 1. \quad \text{(5)}
\]

The careful reader will object that limits have not been defined! You may further object to computing these limits separately, before combining them into equation (4). Nevertheless—following the principle of ideas now, rigor later—I would like to proceed. It is entirely true that the limit of (4) comes from the two limits in (5):

\[
\frac{dy}{dx} = (\sin x)(\text{first limit}) + (\cos x)(\text{second limit}) = 0 + \cos x. \quad \text{(6)}
\]

The secant slope \( \Delta y/\Delta x \) has approached the tangent slope \( dy/dx \).

\[
2.4 \quad \text{The derivative of } \quad y = \sin x \quad \text{is} \quad \frac{dy}{dx} = \cos x.
\]

We cannot pass over the crucial step—the two limits in (5). They contain the real ideas. Both ratios become 0/0 if we just substitute \( h = 0 \). Remember that the cosine of a zero angle is 1, and the sine of a zero angle is 0. Figure 2.8a shows a small angle \( h \) (as near to zero as we could reasonably draw). The edge of length \( \sin h \) is close to zero, and the edge of length \( \cos h \) is near 1. Figure 2.8b shows how the ratio of \( \sin h \) to \( h \) (both headed for zero) gives the slope of the sine curve at the start.

When two functions approach zero, their ratio might do anything. We might have

\[
\frac{h^2}{h} \to 0 \quad \text{or} \quad \frac{h}{h} \to 1 \quad \text{or} \quad \frac{\sqrt{h}}{h} \to \infty.
\]

No clue comes from 0/0. What matters is whether the top or bottom goes to zero more quickly. Roughly speaking, we want to show that \((\cos h - 1)/h\) is like \( h^2/h \) and \((\sin h)/h\) is like \( h/h \).

\textbf{Time out} The graph of \( \sin x \) is in Figure 2.9 (in black). The graph of \( \sin(x + \Delta x) \) sits just beside it (in red). The height difference is \( \Delta f \) when the shift distance is \( \Delta x \).
2 Derivatives

**Fig. 2.9** \( \sin(x + h) \) with \( h = 10^\circ = \pi/18 \) radians. \( \Delta f/\Delta x \) is close to \( \cos x \).

Now divide by that small number \( \Delta x \) (or \( h \)). The second figure shows \( \Delta f/\Delta x \). It is close to \( \cos x \). (Look how it starts—it is not quite \( \cos x \).) Mathematics will prove that the limit is \( \cos x \) exactly, when \( \Delta x \to 0 \). Curiously, the reasoning concentrates on only one point (\( x = 0 \)). The slope at that point is \( \cos 0 = 1 \).

We now prove this: \( \sin \Delta x \) divided by \( \Delta x \) goes to 1. The sine curve starts with slope 1. By the addition formula for \( \sin(x + h) \), this answer at one point will lead to the slope \( \cos x \) at all points.

**Question** Why does the graph of \( f(x + \Delta x) \) shift left from \( f(x) \) when \( \Delta x > 0 \)?

**Answer** When \( x = 0 \), the shifted graph is already showing \( f(\Delta x) \). In Figure 2.9a, the red graph is shifted left from the black graph. The red graph shows \( \sin h \) when the black graph shows \( \sin 0 \).

**THE LIMIT OF \( (\sin h)/h \) IS 1**

There are several ways to find this limit. The direct approach is to let a computer draw a graph. Figure 2.10a is very convincing. The function \( (\sin h)/h \) approaches 1 at the key point \( h = 0 \). So does \( (\tan h)/h \). In practice, the only danger is that you might get a message like “undefined function” and no graph. (The machine may refuse to divide by zero at \( h = 0 \). Probably you can get around that.) Because of the importance of this limit, I want to give a mathematical proof that it equals 1.

**Fig. 2.10** \( (\sin h)/h \) squeezed between \( \cos x \) and 1; \( (\tan h)/h \) decreases to 1.

Figure 2.10b indicates, but still only graphically, that \( \sin h \) stays below \( h \). (The first graph shows that too; \( (\sin h)/h \) is below 1.) We also see that \( \tan h \) stays above \( h \). Remember that the tangent is the ratio of sine to cosine. Dividing by the cosine is enough to push the tangent above \( h \). The crucial inequalities (to be proved when \( h \) is small and positive) are

\[
\sin h < h \quad \text{and} \quad \tan h > h. \tag{7}
\]
Since $\tan h = (\sin h)/(\cos h)$, those are the same as
\[
\frac{\sin h}{h} < 1 \quad \text{and} \quad \frac{\sin h}{h} > \cos h.
\]
(8)

What happens as $h$ goes to zero? The ratio $(\sin h)/h$ is squeezed between $\cos h$ and 1. But $\cos h$ is approaching 1! The squeeze as $h \to 0$ leaves only one possibility for $(\sin h)/h$, which is caught in between: The ratio $(\sin h)/h$ approaches 1.

Figure 2.10 shows that “squeeze play.” If two functions approach the same limit, so does any function caught in between. This is proved at the end of Section 2.6.

For negative values of $h$, which are absolutely allowed, the result is the same. To the left of zero, $h$ reverses sign and $\sin h$ reverses sign. The ratio $(\sin h)/h$ is unchanged. (The sine is an odd function: $\sin(-h) = -\sin h$) The ratio is an even function, symmetric around zero and approaching 1 from both sides.

The proof depends on $\sin h < h < \tan h$, which is displayed by the graph but not explained. We go back to right triangles.

Figure 2.11 shows why $\sin h < h$. The straight line $PQ$ has length $2 \sin h$. The circular arc must be longer, because the shortest distance between two points is a straight line. The arc $PQ$ has length $2h$. (Important: When the radius is 1, the arc length equals the angle. The full circumference is $2\pi$ and the full angle is also $2\pi$.)

The straight distance $2 \sin h$ is less than the circular distance $2h$, so $\sin h < h$.

Figure 2.11b shows why $h < \tan h$. This time we look at areas. The triangular area is $\frac{1}{2}(\text{base})(\text{height}) = \frac{1}{2}(1)(\tan h)$. Inside that triangle is the shaded sector of the circle. Its area is $h/2\pi$ times the area of the whole circle (because the angle is that fraction of the whole angle). The circle has area $\pi r^2 = \pi$, so multiplication by $h/2\pi$ gives $\frac{1}{2}h$ for the area of the sector. Comparing with the triangle around it, $\frac{1}{2} \tan h > \frac{1}{2}h$.

The inequalities $\sin h < h < \tan h$ are now proved. The squeeze in equation (8) produces $(\sin h)/h \to 1$. Q.E.D. Problem 13 shows how to prove $\sin h < h$ from areas.

Note All angles $x$ and $h$ are being measured in radians. In degrees, $\cos x$ is not the derivative of $\sin x$. A degree is much less than a radian, and $dy/dx$ is reduced by the factor $2\pi/360$.

THE LIMIT OF $(\cos h - 1)/h$ IS 0

This second limit is different. We will show that $1 - \cos h$ shrinks to zero more quickly than $h$. Cosines are connected to sines by $(\sin h)^2 + (\cos h)^2 = 1$. We start from the

†If we try to prove that, we will be here all night. Accept it as true.
known fact $\sin h < h$ and work it into a form involving cosines:

$$ (1 - \cos h)(1 + \cos h) = 1 - (\cos h)^2 = (\sin h)^2 < h^2. \quad (9) $$

Note that everything is positive. Divide through by $h$ and also by $1 + \cos h$:

$$ 0 < \frac{1 - \cos h}{h} < \frac{h}{1 + \cos h}. \quad (10) $$

Our ratio is caught in the middle. The right side goes to zero because $h \to 0$. This is another “squeeze”—there is no escape. Our ratio goes to zero.

For $\cos h - 1$ or for negative $h$, the signs change but minus zero is still zero. This confirms equation (6). The slope of $\sin x$ is $\cos x$.

**Remark** Equation (10) also shows that $1 - \cos h$ is approximately $\frac{1}{2}h^2$. The 2 comes from $1 + \cos h$. This is a basic purpose of calculus—to find simple approximations like $\frac{1}{2}h^2$. A “tangent parabola” $1 - \frac{1}{2}h^2$ is close to the top of the cosine curve.

### THE DERIVATIVE OF THE COSINE

This will be easy. The quick way to differentiate $\cos x$ is to shift the sine curve by $\pi/2$. That yields the cosine curve (solid line in Figure 2.12b). The derivative also shifts by $\pi/2$ (dotted line). **The derivative of $\cos x$ is $-\sin x$.**

Notice how the dotted line (the slope) goes below zero when the solid line turns downward. The slope equals zero when the solid line is level. **Increasing functions have positive slopes. Decreasing functions have negative slopes.** That is important, and we return to it.

There is more information in $dy/dx$ than “function rising” or “function falling.” The slope tells how quickly the function goes up or down. It gives the rate of change. The slope of $y = \cos x$ can be computed in the normal way, as the limit of $\Delta y/\Delta x$:

$$ \frac{\Delta y}{\Delta x} = \cos(x + h) - \cos x = \cos x \left( \frac{\cos h - 1}{h} \right) - \sin x \left( \frac{\sin h}{h} \right) $$

$$ \frac{dy}{dx} = (\cos x)(0) - (\sin x)(1) = -\sin x. \quad (11) $$

The first line came from formula (3) for $\cos(x + h)$. The second line took limits, reaching 0 and 1 as before. This confirms the graphical proof that the slope of $\cos x$ is $-\sin x$.

#### Fig. 2.12

$y(x)$ increases where $y'$ is positive. $y(x)$ bends up where $y''$ is positive.

### THE SECOND DERIVATIVES OF THE SINE AND COSINE

We now introduce the derivative of the derivative. That is the second derivative of the original function. It tells how fast the slope is changing, not how fast $y$ itself is...
2.4 The Derivative of the Sine and Cosine

changing. The second derivative is the “rate of change of the velocity.” A straight line has constant slope (constant velocity), so its second derivative is zero:

$$f(t) = 5t \text{ has } df/dt = 5 \text{ and } d^2f/dt^2 = 0.$$  

The parabola $y = x^2$ has slope $2x$ (linear) which has slope 2 (constant). Similarly

$$f(t) = \frac{1}{2}at^2 \text{ has } df/dt = at \text{ and } d^2f/dt^2 = a.$$  

There stands the notation $d^2f/dt^2$ (or $d^2y/dx^2$) for the second derivative. A short form is $f''$ or $y''$. (This is pronounced $f$ double prime or $y$ double prime). Example: The second derivative of $y = x^3$ is $y'' = 6x$.

In the distance-velocity problem, $f''$ is acceleration. It tells how fast $v$ is changing, while $v$ tells how fast $f$ is changing. Where $df/dt$ was distance/time, the second derivative is distance/(time)$^2$. The acceleration due to gravity is about 32 ft/sec$^2$ or 9.8 m/sec$^2$, which means that $v$ increases by 32 ft/sec in one second. It does not mean that the distance increases by 32 feet!

The graph of $y = \sin t$ increases at the start. Its derivative $\cos t$ is positive. However the second derivative is $-\sin t$. The curve is bending down while going up. The arch is “concave down” because $y'' = -\sin t$ is negative.

At $t = \pi$ the curve reaches zero and goes negative. The second derivative becomes positive. Now the curve bends upward. The lower arch is “concave up.”

$y'' > 0$ means that $y'$ increases so $y$ bends upward (concave up)  

$y'' < 0$ means that $y'$ decreases so $y$ bends down (concave down).

Chapter 3 studies these things properly—here we get an advance look for $\sin t$.

The remarkable fact about the sine and cosine is that $y'' = -y$. That is unusual and special: acceleration = distance. The greater the distance, the greater the force pulling back:

$$y = \sin t \text{ has } dy/dt = +\cos t \text{ and } d^2y/dt^2 = -\sin t = -y.$$  

$$y = \cos t \text{ has } dy/dt = -\sin t \text{ and } d^2y/dt^2 = -\cos t = -y.$$  

**Question** Does $d^2y/dt^2 < 0$ mean that the distance $y(t)$ is decreasing?  

**Answer** No. Absolutely not! It means that $dy/dt$ is decreasing, not necessarily $y$. At the start of the sine curve, $y$ is still increasing but $y'' < 0$.

Sines and cosines give simple harmonic motion—up and down, forward and back, out and in, tension and compression. Stretch a spring, and the restoring force pulls it back. Push a swing up, and gravity brings it down. These motions are controlled by a differential equation:

$$\frac{d^2y}{dt^2} = -y. \quad (12)$$

All solutions are combinations of the sine and cosine: $y = A \sin t + B \cos t$.

This is not a course on differential equations. But you have to see the purpose of calculus. It models events by equations. It models oscillation by equation (12). Your heart fills and empties. Balls bounce. Current alternates. The economy goes up and down:

$$\text{high prices} \rightarrow \text{high production} \rightarrow \text{low prices} \rightarrow \cdots$$

We can’t live without oscillations (or differential equations).
### 2.4 EXERCISES

#### 2 Derivatives

**Read-through questions**

The derivative of \( y = \sin x \) is \( y' = \text{______} \). The second derivative (the _____ of the derivative) is \( y'' = \text{______} \). The fourth derivative is \( y'''' = \text{______} \). Thus \( y = \sin x \) satisfies the differential equations \( y'' = \text{______} \) and \( y''' = \text{______} \). So does \( y = \cos x \), whose second derivative is \( \text{______} \).

All these derivatives come from one basic limit: \( (\sin h)/h \) approaches \( \text{______} \). The sine of .01 radians is very close to \( \text{______} \). So is the \( \text{______} \) of .01. The cosine of .01 is not .99, because \( 1 - \cos h \) is much \( \text{______} \) than \( h \). The ratio \( (1 - \cos h)/h^2 \) approaches \( \text{______} \). Therefore \( \cos h \) is close to \( 1 - \frac{1}{2} h^2 \) and \( \cos .01 \approx \text{______} \). We can replace \( h \) by \( x \).

The differential equation \( y'' = -y \) leads to \( \text{______} \). When \( y \) is positive, \( y'' \) is \( \text{______} \). Therefore \( y' \) is \( \text{______} \). Eventually \( y \) goes below zero and \( y' \) becomes \( \text{______} \). Then \( y' \) is \( \text{______} \). Examples of oscillation in real life are \( \text{______} \) and \( \text{______} \).

1. Which of these ratios approach 1 as \( h \to 0 \)?
   (a) \( \frac{h}{\sin h} \) (b) \( \frac{\sin^2 h}{h^2} \) (c) \( \frac{\sin h}{2h} \) (d) \( \frac{\sin(-h)}{h} \)

2. (Calculator) Find \( (\sin h)/h \) at \( h = 0.5 \) and 0.1 and .01. Where does \( (\sin h)/h \) go above .99?

3. Find the limits as \( h \to 0 \) of
   (a) \( \frac{\sin^2 h}{h} \) (b) \( \frac{\sin 3h}{5h} \) (c) \( \frac{\sin 5h}{h} \) (d) \( \frac{\sin h}{5h} \)

4. Where does \( \tan h = 1.01h \)? Where does \( \tan h = h \)?

5. \( y = \sin x \) has period \( 2\pi \), which means that \( \sin x = \text{______} \). The limit of \( (\sin(2\pi + h) - \sin 2\pi)/h \) is 1 because \( \text{______} \). This gives \( \frac{dy}{dx} \) at \( x = \text{______} \).

6. Draw \( \cos (x + \Delta x) \) next to \( \cos x \). Mark the height difference \( \Delta y \). Then draw \( \frac{\Delta y}{\Delta x} \) as in Figure 2.9.

7. The key to trigonometry is \( \cos^2 \theta = 1 - \sin^2 \theta \). Set \( \sin \theta \approx \theta \) to find \( \cos^2 \theta \approx 1 - \theta^2 \). The square root is \( \cos \theta \approx 1 - \frac{1}{2} \theta^2 \). Reason: Squaring gives \( \cos^2 \theta \approx \text{______} \) and the correction term \( \text{______} \) is very small near \( \theta = 0 \).

8. (Calculator) Compare \( \cos \theta \) with \( 1 - \frac{1}{2} \theta^2 \) for
   (a) \( \theta = 0.1 \) (b) \( \theta = 0.5 \) (c) \( \theta = 30^\circ \) (d) \( \theta = 3^\circ \).

9. Trigonometry gives \( \cos \theta = 1 - 2 \sin^2 \frac{1}{2} \theta \). The approximation \( \sin \frac{1}{2} \theta \approx \text{______} \) leads directly to \( \cos \theta \approx 1 - \frac{1}{2} \theta^2 \).

10. Find the limits as \( h \to 0 \):
    (a) \( \frac{1 - \cos h}{h^2} \) (b) \( \frac{1 - \cos^2 h}{h^2} \) (c) \( \frac{1 - \cos^2 h}{\sin^2 h} \) (d) \( \frac{1 - \cos 2h}{h} \)

11. Find by calculator or calculus:
    (a) \( \lim_{h \to 0} \frac{\sin 3h}{\sin 2h} \) (b) \( \lim_{h \to 0} \frac{1 - \cos 2h}{1 - \cos h} \)

12. Compute the slope at \( x = 0 \) directly from limits:
    (a) \( y = \tan x \) (b) \( y = \sin(-x) \)

13. The unmarked points in Figure 2.11 are \( P \) and \( S \). Find the height \( PS \) and the area of triangle \( OPR \). Prove by areas that \( \sin h < h \).

14. The slopes of \( \cos x \) and \( 1 - \frac{1}{2} x^2 \) are \( -\sin x \) and \( \text{______} \).

15. Chapter 10 gives an infinite series for \( \sin x \):
    \[ \sin x = x - \frac{x^3}{1 \cdot 3 \cdot 2 \cdot 1} + \frac{x^5}{5 \cdot 4 \cdot 3 \cdot 2 \cdot 1} - \cdots. \]

   From the derivative find the series for \( \cos x \). Then take its derivative to get back to \( -\sin x \).

16. A centered difference for \( f(x) = \sin x \) is
    \[ \frac{f(x + h) - f(x - h)}{2h} = \frac{\sin(x + h) - \sin(x - h)}{2h} = ? \]

   Use the addition formula (2). Then let \( h \to 0 \).

17. Repeat Problem 16 to find the slope of \( \cos x \). Use formula (3) to simplify \( \cos(x + h) - \cos(x - h) \).

18. Find the tangent line to \( y = \sin x \) at
    (a) \( x = 0 \) (b) \( x = \pi \) (c) \( x = \pi/4 \)

19. Where does \( y = \sin x + \cos x \) have zero slope?

20. Find the derivative of \( \sin(x + 1) \) in two ways:
    (a) Expand to \( \sin x \cos 1 + \cos x \sin 1 \). Compute \( dy/dx \).
    (b) Divide \( \Delta y = \sin(x + 1 + \Delta x) - \sin(x + 1) \) by \( \Delta x \). Write \( X \) instead of \( x + 1 \). Let \( \Delta x \) go to zero.

21. Show that \( (\tan h)/h \) is squeezed between 1 and \( 1/\cos h \). As \( h \to 0 \) the limit is \( \text{______} \).

22. For \( y = \sin 2x \), the ratio \( \Delta y/h \) is
    \[ \sin 2(x + h) - \sin 2x = \frac{\sin 2(x + h - 1) + \cos 2x \sin 2h}{h} \]

   Explain why the limit \( dy/dx \) is \( 2 \cos 2x \).

23. Draw the graph of \( y = \sin \frac{1}{2} x \). State its slope at \( x = 0, \pi/2, \) \( \pi, \) and \( 2\pi \). Does \( \frac{1}{2} \) \( \sin x \) have the same slopes?

24. Draw the graph of \( y = \sin x + \sqrt{3} \cos x \). Its maximum value is \( y = \text{______} \) at \( x = \text{______} \). The slope at that point is \( \text{______} \).

25. By combining \( \sin x \) and \( \cos x \), find a combination that starts at \( x = 0 \) from \( y = 2 \) with slope 1. This combination also solves \( y'' = \text{______} \).
26 True or false, with reason:
(a) The derivative of $\sin^2 x$ is $\cos^2 x$
(b) The derivative of $\cos(-x)$ is $\sin x$
(c) A positive function has a negative second derivative.
(d) If $y'$ is increasing then $y''$ is positive.

27 Find solutions to $dy/dx = \sin 3x$ and $dy/dx = \cos 3x$.

28 If $y = \sin 5x$ then $y' = 5 \cos 5x$ and $y'' = -25 \sin 5x$. So this function satisfies the differential equation $y'' = \text{______}$.

29 If $h$ is measured in degrees, find $\lim_{h \to 0} (\sin h)/h$. You could set your calculator in degree mode.

30 Write down a ratio that approaches $dy/dx$ at $x = \pi$. For $y = \sin x$ and $\Delta x = 0.01$ compute that ratio.

31 By the square rule, the derivative of $(u(x))^2$ is $2u \frac{du}{dx}$. Take the derivative of each term in $\sin^2 x + \cos^2 x = 1$.

32 Give an example of oscillation that does not come from physics. Is it simple harmonic motion (one frequency only)?

33 Explain the second derivative in your own words.

---

### 2.5 The Product and Quotient and Power Rules

What are the derivatives of $x + \sin x$ and $x \sin x$ and $1/\sin x$ and $x/\sin x$ and $\sin^n x$? Those are made up from the familiar pieces $x$ and $\sin x$, but we need new rules. Fortunately they are rules that apply to every function, so they can be established once and for all. If we know the separate derivatives of two functions $u$ and $v$, then the derivatives of $u + v$ and $uv$ and $1/v$ and $u/v$ and $u^n$ are immediately available.

This is a straightforward section, with those five rules to learn. It is also an important section, containing most of the working tools of differential calculus. But I am afraid that five rules and thirteen examples (which we need—the eyes glaze over with formulas alone) make a long list. At least the easiest rule comes first. **When we add functions, we add their derivatives.**

**Sum Rule**

The derivative of the sum $u(x) + v(x)$ is $\frac{d}{dx}(u + v) = \frac{du}{dx} + \frac{dv}{dx}$. (1)

**EXAMPLE 1** The derivative of $x + \sin x$ is $1 + \cos x$. That is tremendously simple, but it is fundamental. The interpretation for distances may be more confusing (and more interesting) than the rule itself:

Suppose a train moves with velocity 1. The distance at time $t$ is $t$. On the train a professor paces back and forth (in simple harmonic motion). His distance from his seat is $\sin t$. Then the total distance from his starting point is $t + \sin t$, and his velocity (train speed plus walking speed) is $1 + \cos t$.

If you add distances, you add velocities. Actually that example is ridiculous, because the professor's maximum speed equals the train speed ($= 1$). He is running like mad, not pacing. Occasionally he is standing still with respect to the ground.

The sum rule is a special case of a bigger rule called "linearity." It applies when we add or subtract functions and multiply them by constants—as in $3x - 4 \sin x$. By linearity the derivative is $3 - 4 \cos x$. The rule works for all functions $u(x)$ and $v(x)$. A **linear combination** is $y(x) = au(x) + bv(x)$, where $a$ and $b$ are any real numbers. Then $\Delta y/\Delta x$ is

$$\frac{au(x + \Delta x) + bv(x + \Delta x) - au(x) - bv(x)}{\Delta x} = a \frac{u(x + \Delta x) - u(x)}{\Delta x} + b \frac{v(x + \Delta x) - v(x)}{\Delta x}.$$
2 Derivatives

The limit on the left is \( \frac{dy}{dx} \). The limit on the right is \( a \frac{du}{dx} + b \frac{dv}{dx} \). We are allowed to take limits separately and add. The result is what we hope for:

**Rule of Linearity**

The derivative of \( au(x) + bv(x) \) is

\[
\frac{d}{dx} (au + bv) = a \frac{du}{dx} + b \frac{dv}{dx}.
\]

(2)

The **product rule** comes next. It can't be so simple—products are not linear. The sum rule is what you would have done anyway, but products give something new. The **derivative of \( u \) times \( v \) is not \( du/dx \) times \( dv/dx \)**. Example: The derivative of \( x^5 \) is \( 5x^4 \). Don't multiply the derivatives of \( x^3 \) and \( x^2 \). (\( 3x^2 \) times \( 2x \) is not \( 5x^4 \).) 

*For a product of two functions, the derivative has two terms.*

**Product Rule** (the key to this section)

The derivative of \( u(x)v(x) \) is

\[
\frac{d}{dx} (uv) = u \frac{dv}{dx} + v \frac{du}{dx}.
\]

(3)

**EXAMPLE 2** \( u = x^3 \) times \( v = x^2 \) is \( uv = x^5 \). The product rule leads to \( 5x^4 \):

\[
x^3 \frac{dv}{dx} + x^2 \frac{du}{dx} = x^3(2x) + x^2(3x^2) = 2x^4 + 3x^4 = 5x^4.
\]

**EXAMPLE 3** In the slope of \( x \sin x \), I don't write \( dx/dx = 1 \) but it's there:

\[
\frac{d}{dx} (x \sin x) = x \cos x + \sin x.
\]

**EXAMPLE 4** If \( u = \sin x \) and \( v = \sin x \) then \( uv = \sin^2 x \). We get two equal terms:

\[
\sin x \frac{d}{dx} (\sin x) + \sin x \frac{d}{dx} (\sin x) = 2 \sin x \cos x.
\]

This confirms the "square rule" \( 2u \frac{du}{dx} \), when \( u \) is the same as \( v \). Similarly the slope of \( \cos^2 x \) is \( -2 \cos x \sin x \) (minus sign from the slope of the cosine).

**Question** Those answers for \( \sin^2 x \) and \( \cos^2 x \) have opposite signs, so the derivative of \( \sin^2 x + \cos^2 x \) is zero (sum rule). How do you see that more quickly?

**EXAMPLE 5** The derivative of \( uvw \) is \( uvw' + uv'w + u'vw \)—one derivative at a time. The derivative of \( xxx \) is \( xx + xx + xx \).

**Fig. 2.13** Change in length = \( \Delta u + \Delta v \). Change in area = \( u \Delta v + v \Delta u + \Delta u \Delta v \).
After those examples we prove the product rule. Figure 2.13 explains it best. The area of the big rectangle is $uv$. The \textbf{important changes in area are the two strips $u \Delta v$ and $v \Delta u$.} The corner area $\Delta u \Delta v$ is much smaller. When we divide by $\Delta x$, the strips give $u \Delta v/\Delta x$ and $v \Delta u/\Delta x$. The corner gives $\Delta u \Delta v/\Delta x$, which approaches zero.

Notice how the sum rule is in one dimension and the product rule is in two dimensions. The rule for $uvw$ would be in three dimensions.

The extra area comes from the whole top strip plus the side strip. By algebra,

$$u(x + h)v(x + h) - u(x)v(x) = u(x + h)[v(x + h) - v(x)] + v(x)[u(x + h) - u(x)].$$

This increase is $u(x + h)\Delta v + v(x)\Delta u$—top plus side. \textit{Now divide by $h$ (or $\Delta x$) and let $h \to 0$.} The left side of equation (4) becomes the derivative of $u(x)v(x)$. The right side becomes $u(x)$ times $dv/dx$—we can multiply the two limits—plus $v(x)$ times $du/dx$.

That proves the product rule—definitely useful.

We could go immediately to the quotient rule for $u(x)/v(x)$. But start with $u = 1$. The derivative of $1/x$ is $-1/x^2$ (known). What is the derivative of $1/v(x)$?

\textbf{Reciprocal Rule}

The derivative of \( \frac{1}{v(x)} \) is \( -\frac{dv}{dx} \frac{1}{v^2} \).

The proof starts with $(v)(1/v) = 1$. The derivative of 1 is 0. Apply the product rule:

$$\frac{d}{dx} \left( \frac{1}{v} \right) + \frac{1}{v} \frac{dv}{dx} = 0 \quad \text{so that} \quad \frac{d}{dx} \left( \frac{1}{v} \right) = -\frac{dv}{dx} \frac{1}{v^2}. \quad (6)$$

It is worth checking the units—in the reciprocal rule and others. A test of dimensions is automatic in science and engineering, and a good idea in mathematics. The test ignores constants and plus or minus signs, but it prevents bad errors. If $v$ is in dollars and $x$ is in hours, $dv/dx$ is in dollars/hour. Then dimensions agree:

$$\frac{d}{dx} \left( \frac{1}{v} \right) \approx \frac{(1/\text{dollars})}{\text{hour}} \quad \text{and also} \quad -\frac{dv}{dx} \frac{1}{v^2} \approx \frac{\text{dollars/hour}}{(\text{dollars})^2}. \quad (5)$$

From this test, the derivative of $1/v$ cannot be $1/(dv/dx)$. A similar test shows that Einstein's formula $e = mc^2$ is dimensionally possible. The theory of relativity might be correct! Both sides have the dimension of $(\text{mass})(\text{distance})^2/(\text{time})^2$, when mass is converted to energy.†

\textbf{Example 6} \ The derivatives of $x^{-1}$, $x^{-2}$, $x^{-n}$ are $-1x^{-2}$, $-2x^{-3}$, $-nx^{-n-1}$.

Those come from the reciprocal rule with $v = x$ and $x^2$ and any $x^n$:

$$\frac{d}{dx} (x^{-n}) = \frac{d}{dx} \left( \frac{1}{x^n} \right) = -\frac{nx^{n-1}}{(x^n)^2} = -nx^{-n-1}. \quad (7)$$

\textit{For negative and positive exponents the derivative of $x^n$ is $nx^{n-1}$.}
EXAMPLE 7 The derivatives of $\frac{1}{\cos x}$ and $\frac{1}{\sin x}$ are $\frac{+\sin x}{\cos^2 x}$ and $\frac{-\cos x}{\sin^2 x}$.

Those come directly from the reciprocal rule. In trigonometry, $1/\cos x$ is the secant of the angle $x$, and $1/\sin x$ is the cosecant of $x$. Now we have their derivatives:

$$\frac{d}{dx}(\sec x) = \frac{\sin x}{\cos x \cos x} = \frac{1}{\cos x} \frac{\sin x}{\cos x} = \sec x \tan x.$$

$$\frac{d}{dx}(\csc x) = -\frac{\cos x}{\sin x \sin x} = -\frac{1}{\sin x} \frac{\cos x}{\sin x} = -\csc x \cot x.$$

Those formulas are often seen in calculus. If you have a good memory they are worth storing. Like most mathematicians, I have to check them every time before using them (maybe once a year). It is really the rules that are basic, not the formulas.

The next rule applies to the quotient $u(x)/v(x)$. That is $u$ times $1/v$. Combining the product rule and reciprocal rule gives something new and important:

**Quotient Rule**

The derivative of $\frac{u(x)}{v(x)}$ is $\frac{1}{v} \frac{du}{dx} - u \frac{dv}{dx} = \frac{v \frac{du}{dx} - u \frac{dv}{dx}}{v^2}.$

You must memorize that last formula. The $v^2$ is familiar. The rest is new, but not very new. If $v = 1$ the result is $du/dx$ (of course). For $u = 1$ we have the reciprocal rule. Figure 2.14b shows the difference $(u + \Delta u)/(v + \Delta v) - (u/v)$. The denominator $v(v + \Delta v)$ is responsible for $v^2$.

EXAMPLE 8 (only practice) If $u/v = x^5/x^3$ (which is $x^2$) the quotient rule gives $2x$:

$$\frac{d}{dx} \left( \frac{x^5}{x^3} \right) = \frac{x^3(5x^4) - x^5(3x^2)}{x^6} = \frac{5x^7 - 3x^7}{x^6} = 2x.$$

EXAMPLE 9 (important) For $u = \sin x$ and $v = \cos x$, the quotient is $\sin x/\cos x = \tan x$. The derivative of $\tan x$ is $\sec^2 x$. Use the quotient rule and $\cos^2 x + \sin^2 x = 1$:

$$\frac{d}{dx} \left( \frac{\sin x}{\cos x} \right) = \frac{\cos x(\cos x) - \sin x(-\sin x)}{\cos^2 x} = \frac{1}{\cos^2 x} = \sec^2 x.$$

Again to memorize: $(\tan x)' = \sec^2 x$. At $x = 0$, this slope is 1. The graphs of $\sin x$ and $x$ and $\tan x$ all start with this slope (then they separate). At $x = \pi/2$ the sine curve is flat ($\cos x = 0$) and the tangent curve is vertical ($\sec^2 x = \infty$).

The slope generally blows up faster than the function. We divide by $\cos x$, once for the tangent and twice for its slope. The slope of $1/x$ is $-1/x^2$. The slope is more sensitive than the function, because of the square in the denominator.

EXAMPLE 10

$$\frac{d}{dx} \left( \frac{\sin x}{x} \right) = \frac{x \cos x - \sin x}{x^2}.$$
That one I hesitate to touch at $x = 0$. Formally it becomes $0/0$. In reality it is more like $0^3/0^2$, and the true derivative is zero. Figure 2.10 showed graphically that $(\sin x)/x$ is flat at the center point. The function is even (symmetric across the $y$ axis) so its derivative can only be zero.

This section is full of rules, and I hope you will allow one more. It goes beyond $x^n$ to $(u(x))^n$. A power of $x$ changes to a power of $u(x)$—as in $(\sin x)^6$ or $(\tan x)^7$ or $(x^2 + 1)^6$. The derivative contains $nu^{n-1}$ (copying $nx^{n-1}$), but there is an extra factor $du/dx$. Watch that factor in $6(\sin x)^5\cos x$ and $7(\tan x)^6\sec^2 x$ and $8(x^2 + 1)^7(2x)$:

**Power Rule**

\[
\frac{d}{dx} \left[ u(x)^n \right] = n \left[ u(x)^{n-1} \right] \frac{du}{dx}.
\]

For $n = 1$ this reduces to $du/dx = du/dx$. For $n = 2$ we get the square rule $2u \frac{du}{dx}$. Next comes $u^3$. The best approach is to use mathematical induction, which goes from each $n$ to the next power $n + 1$ by the product rule:

\[
\frac{d}{dx} \left( u^{n+1} \right) = \frac{d}{dx} (u^n u) = u^n \frac{du}{dx} + u \left( nu^{n-1} \frac{du}{dx} \right) = (n + 1)u^n \frac{du}{dx}.
\]

That is exactly equation (12) for the power $n + 1$. We get all positive powers this way, going up from $n = 1$—then the negative powers come from the reciprocal rule.

Figure 2.15 shows the power rule for $n = 1, 2, 3$. The cube makes the point best. The three thin slabs are $u$ by $u$ by $\Delta u$. The change in volume is essentially $3u^2\Delta u$. From multiplying out $(u + \Delta u)^3$, the exact change in volume is $3u^2\Delta u + 3u(\Delta u)^2 + (\Delta u)^3$—which also accounts for three narrow boxes and a midget cube in the corner. This is the binomial formula in a picture.

---

**Example 11**

\[
\frac{d}{dx} (\sin x)^n = n(\sin x)^{n-1} \cos x. \quad \text{The extra factor } \cos x \text{ is } du/dx.
\]

Our last step finally escapes from a very undesirable restriction—that $n$ must be a whole number. We want to allow fractional powers $n = p/q$, and keep the same formula. The derivative of $x^n$ is still $nx^{n-1}$.

To deal with square roots I can write $\left(\sqrt{x}\right)^2 = x$. Its derivative is $2\sqrt{x}\left(\sqrt{x}\right)' = 1$. Therefore $(\sqrt{x})'$ is $1/2\sqrt{x}$, which fits the formula when $n = 1/2$. Now try $n = p/q$:  

---

**Fig. 2.15** Length change $= \Delta u$; area change $\approx 2u \Delta u$; volume change $\approx 3u^2 \Delta u$. 

---

---

---

---
2 Derivatives

**Fractional powers** Write \( u = x^{p/q} \) as \( u^q = x^p \). Take derivatives, assuming they exist:

\[
\frac{du}{dx} q u^{q-1} \frac{du}{dx} = p x^{p-1} \quad \text{(power rule on both sides)}
\]

\[
\frac{du}{dx} = \frac{p x^{p-1}}{q u^{q-1}} \quad \text{(cancel } x^p \text{ with } u^q)\]

\[
\frac{du}{dx} = n x^{n-1} \quad \text{(replace } p/q \text{ by } n \text{ and } u \text{ by } x^n)\]

**Example 12** The slope of \( x^{1/3} \) is \( \frac{1}{3} x^{-2/3} \). The slope is infinite at \( x = 0 \) and zero at \( x = \infty \). But the curve in Figure 2.16 keeps climbing. It doesn't stay below an "asymptote."

![Graph of y = x^{1/3} and y = x^{4/3}](image)

**Example 13** The slope of \( x^{4/3} \) is \( \frac{4}{3} x^{1/3} \). The slope is zero at \( x = 0 \) and infinite at \( x = \infty \). The graph climbs faster than a line and slower than a parabola (\( \frac{4}{3} \) is between 1 and 2). Its slope follows the cube root curve (times \( \frac{4}{3} \)).

**WE STOP NOW!** I am sorry there were so many rules. A computer can memorize them all, but it doesn't know what they mean and you do. Together with the chain rule that dominates Chapter 4, they achieve virtually all the derivatives ever computed by mankind. We list them in one place for convenience.

<table>
<thead>
<tr>
<th>Rule of Linearity</th>
<th>((au + bv)' = au' + bv')</th>
</tr>
</thead>
<tbody>
<tr>
<td>Product Rule</td>
<td>((uv)' = uv' + vu')</td>
</tr>
<tr>
<td>Reciprocal Rule</td>
<td>((1/v)' = -v'/v^2)</td>
</tr>
<tr>
<td>Quotient Rule</td>
<td>((u/v)' = (uv' - uv')/v^2)</td>
</tr>
<tr>
<td>Power Rule</td>
<td>((u^n)' = nu^{n-1}u')</td>
</tr>
</tbody>
</table>

The power rule applies when \( n \) is negative, or a fraction, or any real number. The derivative of \( x^n \) is \( nx^{n-1} \), according to Chapter 6. The derivative of \( (\sin x)^n \) is ________.

And the derivatives of all six trigonometric functions are now established:

\[
\sin x' = \cos x \quad \tan x' = \sec^2 x \quad \sec x' = \sec x \tan x \quad \\
\cos x' = -\sin x \quad \cot x' = -\csc^2 x \quad \csc x' = -\csc x \cot x.
\]
2.5 The Product and Quotient and Power Rules

Read-through questions

The derivatives of \( \sin x \cos x \) and \( \frac{1}{\cos x} \) and \( \sin x \cos x \) and \( \tan^3 x \) come from the \( \boxed{\text{a}} \) rule, \( \boxed{\text{b}} \) rule, \( \boxed{\text{c}} \) rule, and \( \boxed{\text{d}} \) rule. The product of \( \sin x \times \cos x \) has \((uv)' = uv' + e = f\). The derivative of \( \frac{1}{v} \) is \( \boxed{\text{g}} \), so the slope of \( \sec x \) is \( \boxed{\text{h}} \). The derivative of \( u/v \) is \( \boxed{\text{i}} \), so the slope of \( \tan x \) is \( \boxed{\text{j}} \). The derivative of \( \tan^3 x \) is \( \boxed{\text{k}} \). The slope of \( x^n \) is \( \boxed{\text{l}} \) and the slope of \( (u(x))^n \) is \( \boxed{\text{m}} \). With \( n = -1 \) the derivative of \( (\cos x)^{-1} \) is \( \boxed{\text{n}} \), which agrees with the rule for \( \sec x \).

Even simpler is the rule of \( \boxed{\text{o}} \), which applies to \( au(x) + bv(x) \). The derivative is \( \boxed{\text{p}} \). The slope of \( 3 \sin x + 4 \cos x \) is \( \boxed{\text{q}} \). The derivative of \( (3 \sin x + 4 \cos x)^2 \) is \( \boxed{\text{r}} \). The derivative of \( \boxed{\text{s}} \) is \( 4 \sin^3 x \cos x \).

Find the derivatives of the functions in 1–26.

1. \( (x + 1)(x - 1) \)
2. \( (x^2 + 1)(x^2 - 1) \)
3. \( \frac{x^3}{1 + x} + \frac{1}{1 + \sin x} \)
4. \( \frac{1}{1 + x^2} + \frac{1}{1 - \sin x} \)
5. \( (x - 1)(x - 2)(x - 3) \)
6. \( (x - 1)^2(x - 2)^2 \)
7. \( x^2 \cos x + 2x \sin x \)
8. \( x^{1/2}(x + \sin x) \)
9. \( \frac{x^3 + 1}{x + 1} + \frac{\cos x}{\sin x} \)
10. \( \frac{x^2 + 1}{x^2 - 1} + \frac{\sin x}{\cos x} \)
11. \( x^{1/2} \sin^2 x + (\sin x)^{1/2} \)
12. \( x^{3/2} \sin^3 x + (\sin x)^{3/2} \)
13. \( x^4 \cos x + x \cos^4 x \)
14. \( \sqrt{x}(\sqrt{x} + 1)(\sqrt{x} + 2) \)
15. \( \frac{1}{2} x^2 \sin x - x \cos x + \sin x \)
16. \( (x - 6)^{10} + \sin^{10} x \)
17. \( \sec^2 x - x^2 \)
18. \( \csc^2 x - \cot^2 x \)
19. \( \frac{4}{(x - 5)^{3/2}} + \frac{4}{(5 - x)^{3/2}} \)
20. \( \frac{\sin x - \cos x}{\sin x + \cos x} \)
21. \( (\sin x \cos x)^3 + \sin 2x \)
22. \( x \cos x \csc x \)
23. \( u(x)\nu(x)w(x)z(x) \)
24. \( [u(x)]^2 [\nu(x)]^2 \)
25. \( \frac{1}{\tan x} - \frac{1}{\cot x} \)

27. A growing box has length \( t \), width \( \frac{1}{1 + t} \), and height \( \cos t \).

(a) What is the rate of change of the volume?
(b) What is the rate of change of the surface area?

28. With two applications of the product rule show that the derivative of \( uvw \) is \( uv'w + w'u'v \). When a box with sides \( u, v, w \) grows by \( \Delta u, \Delta v, \Delta w \), three slabs are added with volume \( u \Delta w \) and \boxed{\text{a}} and \boxed{\text{b}}.

29. Find the velocity if the distance is \( f(t) = \frac{5t^2}{2} \) for \( t < 10 \), \( 500 + 100\sqrt{t - 100} \) for \( t \geq 10 \).

30. A cylinder has radius \( r = \frac{t^{3/2}}{1 + t^{3/2}} \) and height \( h = \frac{t}{1 + t} \).

(a) What is the rate of change of its volume?
(b) What is the rate of change of its surface area (including top and base)?

31. The height of a model rocket is \( f(t) = \frac{t^3}{1 + t} \).

(a) What is the velocity \( v(t) \)?
(b) What is the acceleration \( dv/dt \)?

32. Apply the product rule to \( u(x)u^2(x) \) to find the power rule for \( u^3(x) \).

33. Find the second derivative of the product \( u(x)v(x) \). Find the third derivative. Test your formulas on \( u = v = x \).

34. Find functions \( y(x) \) whose derivatives are

(a) \( x^3 \)
(b) \( 1/x^3 \)
(c) \( (1 - x)^{3/2} \)
(d) \( \cos^2 x \sin x \)

35. Find the distances \( f(t) \), starting from \( f(0) = 0 \), to match these velocities:

(a) \( v(t) = \cos t \sin t \)
(b) \( v(t) = \tan t \sec^2 t \)
(c) \( v(t) = \sqrt{1 + t} \)

36. Apply the quotient rule to \( (u(x))^3/(u(x))^2 \) and \(-v'/v^2 \).

37. Draw a figure like 2.13 to explain the square rule.

38. Give an example where \( u(x)/v(x) \) is increasing but \( du/dx = dv/dx = 1 \).

39. True or false, with a good reason:

(a) The derivative of \( x^{2n} \) is \( 2nx^{2n-1} \).
(b) By linearity the derivative of \( a(x)u(x) + b(x)v(x) \) is \( a(x)du/dx + b(x)dv/dx \).
(c) The derivative of \( |x|^3 \) is \( 3|x|^2 \).
(d) \( \sec^2 x \) and \( \cot^2 x \) have the same derivative.
(e) \( (uv)' = u'v' \) is true when \( u(x) = 1 \).

40. The cost of \( u \) shares of stock at \( v \) dollars per share is \( uv \) dollars. Check dimensions of \( d(uv)/dt \) and \( u dv/dt \) and \( u du/dt \).

41. If \( u(x)/v(x) \) is a ratio of polynomials of degree \( n \), what are the degrees for its derivative?

42. For \( y = 5x + 3 \), is \( dy/dx \) the same as \( d^2y/dx^2 \)?

43. If you change from \( f(t) = t \cos t \) to its tangent line at \( t = \pi/2 \), find the two-part function \( df/dt \).

44. Explain in your own words why the derivative of \( u(x)v(x) \) has two terms.

45. A plane starts its descent from height \( y = h \) at \( x = -L \) to land at \((0, 0) \). Choose \( a, b, c, d \) so its landing path \( y = ax^3 + bx^2 + cx + d \) is smooth. With \( dx/dt = V = \) constant, find \( dy/dt \) and \( d^2y/dt^2 \) at \( x = 0 \) and \( x = -L \). (To keep \( d^2y/dt^2 \) small, a coast-to-coast plane starts down \( L > 100 \) miles from the airport.)
You have seen enough limits to be ready for a definition. It is true that we have survived this far without one, and we could continue. But this seems a reasonable time to define limits more carefully. The goal is to achieve rigor without rigor mortis.

First you should know that limits of \( \Delta y/\Delta x \) are by no means the only limits in mathematics. Here are five completely different examples. They involve \( n \to \infty \), not \( \Delta x \to 0 \):

1. \( a_n = (n - 3)/(n + 3) \) (for large \( n \), ignore the 3's and find \( a_n \to 1 \))
2. \( a_n = \frac{1}{2} a_{n-1} + 4 \) (start with any \( a_1 \) and always \( a_n \to 8 \))
3. \( a_n = \) probability of living to year \( n \) (unfortunately \( a_n \to 0 \))
4. \( a_n = \) fraction of zeros among the first \( n \) digits of \( n \) (an \( a_n + h? \))
5. \( a_1 = .4, a_2 = .49, a_3 = .493, \ldots \) No matter what the remaining decimals are, the \( a \)'s converge to a limit. Possibly \( a_n \to .493000 \ldots \), but not likely.

The problem is to say what the limit symbol \( \to \) really means.

A good starting point is to ask about convergence to zero. When does a sequence of positive numbers approach zero? What does it mean to write \( a, \to 0 \)? The numbers \( a_1, a_2, a_3, \ldots \), must become “small,” but that is too vague. We will propose four definitions of convergence to zero, and I hope the right one will be clear.

1. All the numbers \( a_n \) are below \( 10^{-10} \). That may be enough for practical purposes, but it certainly doesn’t make the \( a_n \) approach zero.
2. The sequence is getting closer to zero—each \( a_{n+1} \) is smaller than the preceding \( a_n \). This test is met by 1, 1.1, 1.01, 1.001, ... which converges to 1 instead of 0.
3. For any small number you think of, at least one of the \( a_n \)'s is smaller. That pushes something toward zero, but not necessarily the whole sequence. The condition would be satisfied by 1, \( \frac{1}{2} \), 1, \( \frac{1}{3} \), 1, \( \frac{1}{4} \), ..., which does not approach zero.
4. For any small number you think of, the \( a_n \)'s eventually go below that number and stay below. This is the correct definition.

I want to repeat that. To test for convergence to zero, start with a small number—say \( 10^{-10} \). The \( a_n \)'s must go below that number. They may come back up and go below again—the first million terms make absolutely no difference. Neither do the next billion, but eventually all terms must go below \( 10^{-10} \). After waiting longer (possibly a lot longer), all terms drop below \( 10^{-20} \). The tail end of the sequence decides everything.

**Question 1** Does the sequence \( 10^{-3}, 10^{-2}, 10^{-6}, 10^{-5}, 10^{-9}, 10^{-8}, \ldots \) approach 0?

**Answer** Yes. These up and down numbers eventually stay below any \( \varepsilon \).

**Fig. 2.17** Convergence means: Only a finite number of \( a \)'s are outside any strip around \( L \).
2.6 Limits

**Question 2** Does $10^{-4}$, $10^{-6}$, $10^{-4}$, $10^{-8}$, $10^{-4}$, $10^{-10}$, … approach zero?

**Answer** No. This sequence goes below $10^{-4}$ but does not stay below.

There is a recognized symbol for "an arbitrarily small positive number." By worldwide agreement, it is the Greek letter \( \varepsilon \) (epsilon). Convergence to zero means that the sequence eventually goes below \( \varepsilon \) and stays there. The smaller the \( \varepsilon \), the tougher the test and the longer we wait. Think of \( \varepsilon \) as the tolerance, and keep reducing it.

To emphasize that \( \varepsilon \) comes from outside, Socrates can choose it. Whatever \( \varepsilon \) he proposes, the \( a \)'s must eventually be smaller. After some \( a_n \), all the \( a \)'s are below the tolerance \( \varepsilon \). Here is the exact statement:

\[
\text{for any } \varepsilon \text{ there is an } N \text{ such that } a_n < \varepsilon \text{ if } n > N.
\]

Once you see that idea, the rest is easy. Figure 2.17 has \( N = 3 \) and then \( N = 6 \).

**Example 1** The sequence \( \frac{1}{2}, \frac{1}{4}, \frac{1}{8}, \ldots \) starts upward but goes to zero. Notice that \( 1, 4, 9, \ldots, 100, \ldots \) are squares, and \( 2, 4, 8, \ldots, 1024, \ldots \) are powers of 2. Eventually \( 2^n \) grows faster than \( n^2 \), as in \( a_{10} = 100/1024 \). The ratio goes below any \( \varepsilon \).

**Example 2** \( 1, 0, \frac{1}{2}, 0, \frac{1}{4}, 0, \ldots \) approaches zero. These \( a \)'s do not decrease steadily (the mathematical word for steadily is "monotonically") but still their limit is zero. The choice \( \varepsilon = 1/1000 \) produces the right response: Beyond \( a_{2001} \) all terms are below \( 1/1000 \). So \( N = 2001 \) for that \( \varepsilon \).

The sequence \( 1, \frac{1}{2}, \frac{1}{3}, \frac{1}{4}, \frac{1}{5}, \ldots \) is much slower—but it also converges to zero.

Next we allow the numbers \( a_n \) to be negative as well as positive. They can converge upward toward zero, or they can come in from both sides. The test still requires the \( a_n \) to go inside any strip near zero (and stay there). But now the strip starts at \( -\varepsilon \).

The distance from zero is the absolute value \( |a_n| \). Therefore \( a_n \to 0 \) means \( |a_n| \to 0 \). The previous test can be applied to \( |a_n| \):

\[
\text{for any } \varepsilon \text{ there is an } N \text{ such that } |a_n| < \varepsilon \text{ if } n > N.
\]

**Example 3** \( 1, -\frac{1}{2}, \frac{1}{3}, -\frac{1}{4}, \ldots \) converges to zero because \( 1, \frac{1}{2}, \frac{1}{3}, \frac{1}{4}, \ldots \) converges to zero.

It is a short step to limits other than zero. **The limit is \( L \) if the numbers \( a_n - L \) converge to zero.** Our final test applies to the absolute value \( |a_n - L| \):

\[
\text{for any } \varepsilon \text{ there is an } N \text{ such that } |a_n - L| < \varepsilon \text{ if } n > N.
\]

This is the definition of convergence! Only a finite number of \( a \)'s are outside any strip around \( L \) (Figure 2.18). We write \( a_n \to L \) or \( \lim_{n \to \infty} a_n = L \).

![Fig. 2.18](image_url)
EXAMPLE 4  The numbers $\frac{3}{2}$, $\frac{5}{3}$, $\frac{7}{5}$, ... converge to $L = 1$. After subtracting 1 the differences $\frac{1}{2}$, $\frac{2}{3}$, $\frac{1}{5}$, ... converge to zero. Those difference are $|a_n - L|$.

EXAMPLE 5  The sequence $1 + \frac{1}{2}$, $1 + \frac{1}{2} + \frac{1}{3}$, $1 + \frac{1}{2} + \frac{1}{3} + \frac{1}{4}$, ... fails to converge.

The distance between terms is getting smaller. But those numbers $a_1$, $a_2$, $a_3$, $a_4$, ... go past any proposed limit $L$. The second term is $1\frac{1}{2}$. The fourth term adds on $\frac{1}{3} + \frac{1}{4}$, so $a_4$ goes past 2. The eighth term has four new fractions $\frac{1}{5} + \frac{1}{6} + \frac{1}{7} + \frac{1}{8}$, totaling more than $\frac{1}{6} + \frac{1}{7} + \frac{1}{8} = \frac{1}{2}$. Therefore $a_8$ exceeds $2\frac{1}{2}$. Eight more terms will add more than 8 times $\frac{1}{16}$, so $a_{16}$ is beyond 3. The lines in Figure 2.18c are infinitely long, not stopping at any $L$.

In the language of Chapter 10, the harmonic series $1 + \frac{1}{2} + \frac{1}{3} + \ldots$ does not converge. The sum is infinite, because the “partial sums” $a_n$ go beyond every limit $L$ ($a_{5000}$ is past $L = 9$). We will come back to infinite series, but this example makes a subtle point: The steps between the $a_n$ can go to zero while still $a_n \to \infty$.

Thus the condition $a_{n+1} - a_n \to 0$ is not sufficient for convergence. However this condition is necessary. If we do have convergence, then $a_{n+1} - a_n \to 0$. That is a good exercise in the logic of convergence, emphasizing the difference between “sufficient” and “necessary.” We discuss this logic below, after proving that [statement $A$] implies [statement $B$]:

If $[a_n$ converges to $L]$ then $[a_{n+1} - a_n$ converges to zero].

Proof  Because the $a_n$ converge, there is a number $N$ beyond which $|a_n - L| < \varepsilon$ and also $|a_{n+1} - L| < \varepsilon$. Since $a_{n+1} - a_n$ is the sum of $a_{n+1} - L$ and $L - a_n$, its absolute value cannot exceed $\varepsilon + \varepsilon = 2\varepsilon$. Therefore $a_{n+1} - a_n$ approaches zero.

Objection by Socrates: We only got below $2\varepsilon$ and he asked for $\varepsilon$. Our reply: If he particularly wants $|a_{n+1} - a_n| < 1/10$, we start with $\varepsilon = 1/20$. Then $2\varepsilon = 1/10$. But this juggling is not necessary. To stay below $2\varepsilon$ is just as convincing as to stay below $\varepsilon$.

THE LOGIC OF "IF" AND "ONLY IF"

The following page is inserted to help with the language of mathematics. In ordinary language we might say “I will come if you call.” Or we might say “I will come only if you call.” That is different! A mathematician might even say “I will come if and only if you call.” Our goal is to think through the logic, because it is important and not so familiar.†

Statement $A$ above implies statement $B$. Statement $A$ is $a_n \to L$; statement $B$ is $a_{n+1} - a_n \to 0$. Mathematics has at least five ways of writing down $A \Rightarrow B$, and I though you might like to see them together. It seems excessive to have so many expressions for the same idea, but authors get desperate for a little variety. Here are the five ways that come to mind:

$A \Rightarrow B$

$A$ implies $B$

if $A$ then $B$

$A$ is a sufficient condition for $B$

$B$ is true if $A$ is true

†Logical thinking is much more important than $\varepsilon$ and $\delta$. 
EXAMPLES  If [positive numbers are decreasing] then [they converge to a limit].
If [sequences \(a_n\) and \(b_n\) converge] then [the sequence \(a_n + b_n\) converges].
If [\(f(x)\) is the integral of \(v(x)\)] then [\(v(x)\) is the derivative of \(f(x)\)].

Those are all true, but not proved. \(A\) is the hypothesis, \(B\) is the conclusion.

Now we go in the other direction. (It is called the “converse,” not the inverse.) We exchange \(A\) and \(B\). Of course stating the converse does not make it true! \(B\) might imply \(A\), or it might not. In the first two examples the converse was false—the \(a_n\) can converge without decreasing, and \(a_n + b_n\) can converge when the separate sequences do not. The converse of the third statement is true—and there are five more ways to state it:

\[
A \equiv B
\]

\(A\) is implied by \(B\)

if \(B\) then \(A\)

\(A\) is a necessary condition for \(B\)

\(B\) is true only if \(A\) is true

Those words “necessary” and “sufficient” are not always easy to master. The same is true of the deceptively short phrase “if and only if.” The two statements \(A \Rightarrow B\) and \(A \equiv B\) are completely different and they both require proof. That means two separate proofs. But they can be stated together for convenience (when both are true):

\[
A \equiv B
\]

\(A\) implies \(B\) and \(B\) implies \(A\)

\(A\) is equivalent to \(B\)

\(A\) is a necessary and sufficient condition for \(B\)

\(A\) is true if and only if \(B\) is true

EXAMPLES  \([a_n \to L] \leftrightarrow [2a_n \to 2L] \leftrightarrow [a_n + 1 \to L + 1] \leftrightarrow [a_n - L \to 0].\)

RULES FOR LIMITS

Calculus needs a definition of limits, to define \(dy/dx\). That derivative contains two limits: \(\Delta x \to 0\) and \(\Delta y/\Delta x \to dy/dx\). Calculus also needs rules for limits, to prove the sum rule and product rule for derivatives. We started on the definition, and now we start on the rules.

Given two convergent sequences, \(a_n \to L\) and \(b_n \to M\), other sequences also converge:

Addition: \(a_n + b_n \to L + M\)  Subtraction: \(a_n - b_n \to L - M\)

Multiplication: \(a_nb_n \to LM\)  Division: \(a_n/b_n \to L/M\) (provided \(M \neq 0\))

We check the multiplication rule, which uses a convenient identity:

\[
a_nb_n - LM = (a_n - L)(b_n - M) + M(a_n - L) + L(b_n - M).
\]  \hspace{1cm} \text{(2)}

Suppose \(|a_n - L| < \varepsilon\) beyond some point \(N\), and \(|b_n - M| < \varepsilon\) beyond some other point \(N'\). Then beyond the larger of \(N\) and \(N'\), the right side of (2) is small. It is less than \(\varepsilon \cdot \varepsilon + M\varepsilon + L\varepsilon\). This proves that (2) gives \(a_nb_n \to LM\).

An important special case is \(ca_n \to cL\). (The sequence of \(b's\) is \(c, c, c, c, \ldots\)) Thus a constant can be brought “outside” the limit, to give \(\lim ca_n = c \lim a_n\).
THE LIMIT OF \( f(x) \) AS \( x \to a \)

The final step is to replace sequences by functions. Instead of \( a_1, a_2, \ldots \) there is a continuum of values \( f(x) \). The limit is taken as \( x \) approaches a specified point \( a \) (instead of \( n \to \infty \)). Example: As \( x \) approaches \( a = 0 \), the function \( f(x) = 4 - x^2 \) approaches \( L = 4 \). As \( x \) approaches \( a = 2 \), the function \( 5x \) approaches \( L = 10 \). Those statements are fairly obvious, but we have to say what they mean. Somehow it must be this:

if \( x \) is close to \( a \) then \( f(x) \) is close to \( L \).

If \( x - a \) is small, then \( f(x) - L \) should be small. As before, the word small does not say everything. We really mean “arbitrarily small,” or “below any \( \varepsilon \).” The difference \( f(x) - L \) must become as small as anyone wants, when \( x \) gets near \( a \). In that case \( \lim_{x \to a} f(x) = L \). Or we write \( f(x) \to L \) as \( x \to a \).

The statement is awkward because it involves two limits. The limit \( x \to a \) is forcing \( f(x) \to L \). (Previously \( n \to \infty \) forced \( a_n \to a \).) But it is wrong to expect the same \( \varepsilon \) in both limits. We do not and cannot require that \( |x - a| < \varepsilon \) produces \( |f(x) - L| < \varepsilon \). It may be necessary to push \( x \) extremely close to \( a \) (closer than \( \varepsilon \)). We must guarantee that if \( x \) is close enough to \( a \), then \( |f(x) - L| < \varepsilon \).

We have come to the “epsilon-delta definition” of limits. First, Socrates chooses \( \varepsilon \). He has to be shown that \( f(x) \) is within \( \varepsilon \) of \( L \), for every \( x \) near \( a \). Then somebody else (maybe Plato) replies with a number \( \delta \). That gives the meaning of “near \( a \).” Plato’s goal is to get \( f(x) \) within \( \varepsilon \) of \( L \), by keeping \( x \) within \( \delta \) of \( a \):

if \( 0 < |x - a| < \delta \) then \( |f(x) - L| < \varepsilon \). (3)

The input tolerance is \( \delta \) (delta), the output tolerance is \( \varepsilon \). When Plato can find a \( \delta \) for every \( \varepsilon \), Socrates concedes that the limit is \( L \).

EXAMPLE  Prove that \( \lim_{x \to 2} 5x = 10 \). In this case \( a = 2 \) and \( L = 10 \).

Socrates asks for \( |5x - 10| < \varepsilon \). Plato responds by requiring \( |x - 2| < \delta \). What \( \delta \) should he choose? In this case \( |5x - 10| \) is exactly 5 times \( |x - 2| \). So Plato picks \( \delta \) below \( \varepsilon/5 \) (a smaller \( \delta \) is always OK). Whenever \( |x - 2| < \varepsilon/5 \), multiplication by 5 shows that \( |5x - 10| < \varepsilon \).

Remark 1  In Figure 2.19, Socrates chooses the height of the box. It extends above and below \( L \), by the small number \( \varepsilon \). Second, Plato chooses the width. He must make the box narrow enough for the graph to go out the sides. Then \( |f(x) - L| < \varepsilon \).

**Fig. 2.19**  S chooses height \( 2\varepsilon \), then P chooses width \( 2\delta \). Graph must go out the sides.
When $f(x)$ has a jump, the box can't hold it. A step function has no limit as $x$ approaches the jump, because the graph goes through the top or bottom—no matter how thin the box.

**Remark 2** The second figure has $f(x) \to L$, because in taking limits we ignore the final point $x = a$. The value $f(a)$ can be anything, with no effect on $L$. The first figure has more: $f(a) = L$. Then a special name applies—$f$ is continuous. The left figure shows a continuous function, the other figures do not.

We soon come back to continuous functions.

**Remark 3** In the example with $f = 5x$ and $\delta = \varepsilon/5$, the number 5 was the slope. That choice barely kept the graph in the box—it goes out the corners. A little narrower, say $\delta = \varepsilon/10$, and the graph goes safely out the sides. *A reasonable choice is to divide $\varepsilon$ by $2|f'(a)|$. (We double the slope for safety.) I want to say why this $\delta$ works—even if the $\varepsilon$-$\delta$ test is seldom used in practice.*

The ratio of $f(x) - L$ to $x - a$ is distance up over distance across. This is $\Delta f/\Delta x$, close to the slope $f'(a)$. When the distance across is $\delta$, the distance up or down is near $\delta |f'(a)|$. That equals $\varepsilon/2$ for our “reasonable choice” of $\delta$—so we are safely below $\varepsilon$. This choice solves most exercises. But Example 7 shows that a limit might exist even when the slope is infinite.

**EXAMPLE 7** \[
\lim_{x \to 1^+} \sqrt{x - 1} = 0 \quad (a \text{ one-sided limit}).
\]

Notice the plus sign in the symbol $x \to 1^+$. The number $x$ approaches $a = 1$ only from above. An ordinary limit $x \to 1$ requires us to accept $x$ on both sides of 1 (the exact value $x = 1$ is not considered). Since negative numbers are not allowed by the square root, we have a one-sided limit. It is $L = 0$.

Suppose $\varepsilon$ is 1/10. Then the response could be $\delta = 1/100$. A number below 1/100 has a square root below 1/10. In this case the box must be made extremely narrow, $\delta$ much smaller than $\varepsilon$, because the square root starts with infinite slope.

Those examples show the point of the $\varepsilon$-$\delta$ definition. (Given $\varepsilon$, look for $\delta$. This came from Cauchy in France, not Socrates in Greece.) We also see its bad feature: The test is not convenient. Mathematicians do not go around proposing $\varepsilon$'s and replying with $\delta$'s. We may live a strange life, but not that strange.

It is easier to establish once and for all that $5x$ approaches its obvious limit 5a. The same is true for other familiar functions: $x^n \to a^n$ and $\sin x \to \sin a$ and $(1 - x)^{-1} \to (1 - a)^{-1}$—except at $a = 1$. *The correct limit L comes by substituting $x = a$ into the function.* This is exactly the property of a “continuous function.” Before the section on continuous functions, we prove the Squeeze Theorem using $\varepsilon$ and $\delta$.

**Squeeze Theorem** Suppose $f(x) \leq g(x) \leq h(x)$ for $x$ near $a$. If $f(x) \to L$ and $h(x) \to L$ as $x \to a$, then the limit of $g(x)$ is also $L$.

**Proof** $g(x)$ is squeezed between $f(x)$ and $h(x)$. After subtracting $L$, $g(x) - L$ is between $f(x) - L$ and $h(x) - L$. Therefore

$$|g(x) - L| < \varepsilon \quad \text{if} \quad |f(x) - L| < \varepsilon \quad \text{and} \quad |h(x) - L| < \varepsilon.$$  

For any $\varepsilon$, the last two inequalities hold in some region $0 < |x - a| < \delta$. So the first one also holds. This proves that $g(x) \to L$. Values at $x = a$ are not involved—until we get to continuous functions.
2 Derivatives

2.6 EXERCISES

Read-through questions

The limit of \( a_n = (\sin n)/n \) is \( \_a\_ \). The limit of \( a_n = n^n/2^n \) is \( \_b\_ \). The limit of \( a_n = (-1)^n \) is \( \_c\_ \). The meaning of \( a_n \to 0 \) is: Only \( \_d\_ \) of the numbers \( |a_n| \) can be \( \_e\_ \). The meaning of \( a_n \to L \) is: For every \( \_f\_ \) there is an \( \_g\_ \) such that \( \_h\_ \) if \( n > \_i\_ \). The sequence \( 1, 1 + \frac{1}{2}, 1 + \frac{1}{2} + \frac{1}{3}, \ldots \) is not \( \_j\_ \) because eventually those sums go past \( \_k\_ \).

The limit of \( f(x) = \sin x \) as \( x \to a \) is \( \_l\_ \). The limit of \( f(x) = x/|x| \) as \( x \to -2 \) is \( \_m\_ \), but the limit as \( x \to 0 \) does not \( \_n\_ \). This function only has \( \_o\_ \)-sided limits. The meaning of \( \lim_{x \to a} f(x) = L \) is: For every \( \varepsilon \) there is a \( \delta \) such that \( |f(x) - L| < \varepsilon \) whenever \( \_p\_ \).

Two rules for limits, when \( a_n \to L \) and \( b_n \to M \), are \( a_n + b_n \to \_q\_ \) and \( a_n b_n \to \_r\_ \). The corresponding rules for functions, when \( f(x) \to L \) and \( g(x) \to M \) as \( x \to a \), are \( \_s\_ \) and \( \_t\_ \). In all limits, \( |a_n - L| \) or \( |f(x) - L| \) must eventually go below and \( \_u\_ \) any positive \( \_v\_ \).

\( A \Rightarrow B \) means that \( A \) is a \( \_w\_ \) condition for \( B \). Then \( B \) is true \( \_x\_ \) \( A \) is true. \( A \Leftrightarrow B \) means that \( A \) is a \( \_y\_ \) condition for \( B \). Then \( B \) is true \( \_z\_ \) \( A \) is true.

1. What is \( a_n \) and what is the limit \( L \)? After which \( N \) is \( |a_n - L| < 1/10 \)? (Calculator allowed)

(a) \(-1, 1/2, -1/2, \ldots \)  
(b) \(1/2, 1 + 1/2, 1 + 1/2 + 1/2, \ldots \)  
(c) \(1, 2, 3, \ldots \) \( a_n = n/2^n \)  
(d) \(1, 2, 3, 4, 5, \ldots \) \( a_n = n/2^n \)  
(e) \(a_n = \sqrt{n} \)  
(f) \(a_n = \sqrt{n^2 + n - n} \)  
(g) \(1, 1 + 1/2, 1 + 1/2 + 1/3, \ldots \)  

2. Show by example that these statements are false:

(a) If \( a_n \to L \) and \( b_n \to L \), then \( a_n/b_n \to 1 \)  
(b) \( a_n \to L \) if and only if \( a_n^2 \to L^2 \)  
(c) If \( a_n < 0 \) and \( a_n \to L \), then \( L < 0 \)  
(d) If infinitely many \( a_n \)'s are inside every strip around zero then \( a_n \to 0 \)  

3. Which of these statements are equivalent to \( B \Rightarrow A \)?

(a) \( A \) is true so is \( B \)  
(b) \( A \) is true if and only if \( B \) is true  
(c) \( B \) is a sufficient condition for \( A \)  
(d) \( A \) is a necessary condition for \( B \)  

4. Decide whether \( A \Rightarrow B \) or \( B \Rightarrow A \) or neither or both:

(a) \( A = [a_n \to 1] \)  
(b) \( A = [a_n \to 0] \)  
(c) \( A = [a_n \leq n] \)  
(d) \( A = [a_n \to 0] \)  
(e) \( A = [a_n \to 0] \)  
(f) \( A = [a_n < n] \)  

5. If the sequence \( a_1, a_2, a_3, \ldots \) approaches zero, prove that we can put those numbers in any order and the new sequence still approaches zero.

6. Suppose \( f(x) \to L \) and \( g(x) \to M \) as \( x \to a \). Prove from the definitions that \( f(x) + g(x) \to L + M \) as \( x \to a \).  

Find the limits 7–24 if they exist. An \( \varepsilon-\delta \) test is not required.

7. \( \lim_{t \to 2} \frac{t + 3}{t^2 - 2} \)  
8. \( \lim_{t \to 2} \frac{t^2 + 3}{t - 2} \)  
9. \( \lim_{x \to 0} \frac{f(x + h) - f(x)}{h} \) (careful)  
10. \( \lim_{h \to 0} \frac{f(1 + h) - f(1)}{h} \)  
11. \( \lim_{h \to 0} \frac{\sin^2 h \cos^2 h}{h^2} \)  
12. \( \lim_{x \to 0} \frac{2 \tan x}{\sin x} \)  
13. \( \lim_{x \to 0^+} \frac{|x|}{x} \) (one-sided)  
14. \( \lim_{x \to 0^-} \frac{|x|}{x} \) (one-sided)  
15. \( \lim_{x \to 1} \frac{\sin x}{x} \)  
16. \( \lim_{x \to a} \frac{f(c) - f(a)}{c - a} \)  
17. \( \lim_{x \to \frac{5}{2}} \frac{x^2 + 25}{x - 5} \)  
18. \( \lim_{x \to \frac{5}{2}} \frac{x^2 - 25}{x - 5} \)  
19. \( \lim_{x \to 0} \frac{\sqrt{1 + x - 1}}{x} \) (test \( x = .01 \))  
20. \( \lim_{x \to 0} \frac{\sqrt{4 - x}}{x^2 + 6 + x} \)  
21. \( \lim_{x \to 0} \frac{[f(x) - f(a)]}{a} \)  
22. \( \lim_{x \to 0} \frac{[\cos x - \tan x]}{x^2} \)  
23. \( \lim_{x \to 0} \frac{\sin x}{\sin x/2} \)  
24. \( \lim_{x \to 1} \frac{\sin(x - 1)}{x^2 - 1} \)

25. Choose \( \delta \) so that \( |f(x)| < 1/100 \) if \( 0 < x < \delta \).  
26. \( f(x) = 10x \)  
27. \( f(x) = \sqrt{x} \)  
28. \( f(x) = \sin 2x \)  
29. \( f(x) = x \sin x \)  
30. Which of the definitions of a limit require?

(a) \( |f(x) - L| < \varepsilon \Rightarrow 0 < |x - a| < \delta \)  
(b) \( |f(x) - L| < \varepsilon \Rightarrow 0 < |x - a| < \delta \)  
(c) \( |f(x) - L| < \varepsilon \Rightarrow 0 < |x - a| < \delta \)  

28. Give a correct definition of \( \varepsilon-\delta \) for \( f(x) \to 0 \) as \( x \to -\infty \).

29. The limit of \( f(x) = (\sin x)/x \) as \( x \to \infty \) is \( \_\) for \( \varepsilon = .01 \) find a point \( X \) beyond which \( |f(x)| < \varepsilon \).

30. The limit of \( f(x) = 2x/(1 + x) \) as \( x \to \infty \) is \( L = 2 \). For \( \varepsilon = .01 \) find a point \( X \) beyond which \( |f(x) - 2| < \varepsilon \).

31. The limit of \( f(x) = \sin x \) as \( x \to \infty \) does not exist. Explain why not.
Continuous Functions

2.7 Continuous Functions

This will be a brief section. It was originally included with limits, but the combination was too long. We are still concerned with the limit of \( f(x) \) as \( x \rightarrow a \), but a new number is involved. That number is \( f(a) \), the value of \( f \) at \( x = a \). For a “limit,” \( x \) approached \( a \) but never reached it—so \( f(a) \) was ignored. For a “continuous function,” this final number \( f(a) \) must be right.

May I summarize the usual (good) situation as \( x \) approaches \( a \)?

1. The number \( f(a) \) exists (\( f \) is defined at \( a \))
2. The limit of \( f(x) \) exists (it was called \( L \))
3. The limit \( L \) equals \( f(a) \) (\( f(a) \) is the right value)

In such a case, \( f(x) \) is continuous at \( x = a \). These requirements are often written in a single line: \( f(x) \rightarrow f(a) \) as \( x \rightarrow a \). By way of contrast, start with four functions that are not continuous at \( x = 0 \).

32 (Calculator) Estimate the limit of \( \left( 1 + \frac{1}{x} \right)^x \) as \( x \rightarrow \infty \).

33 For the polynomial \( f(x) = 2x - 5x^2 + 7x^3 \) find
   (a) \( \lim_{x \to 1} f(x) \)
   (b) \( \lim_{x \to \infty} f(x) \)
   (c) \( \lim_{x \to \infty} \frac{f(x)}{x^3} \)
   (d) \( \lim_{x \to -\infty} \frac{f(x)}{x^3} \)

34 For \( f(x) = 6x^3 + 1000x \) find
   (a) \( \lim_{x \to \infty} \frac{f(x)}{x} \)
   (b) \( \lim_{x \to -\infty} \frac{f(x)}{x^2} \)
   (c) \( \lim_{x \to \infty} \frac{f(x)}{x^4} \)
   (d) \( \lim_{x \to -\infty} \frac{f(x)}{x^3 + 1} \)

Important rule: As \( x \rightarrow \infty \) the ratio of polynomials \( f(x)/g(x) \) has the same limit as the ratio of their leading terms. \( f(x) = x^3 - x + 2 \) has leading term \( x^3 \) and \( g(x) = 5x^6 + x + 1 \) has leading term \( 5x^6 \). Therefore \( f(x)/g(x) \) behaves like \( x^3/5x^6 \to 0 \), \( g(x)/f(x) \) behaves like \( 5x^6/x^3 \to \infty \), \( f(x)^2/g(x) \) behaves like \( x^6/5x^6 \to 1/5 \).

35 Find the limit as \( x \to \infty \) if it exists:
   \[
   \frac{3x^2 + 2x + 1}{3 + 2x + x^2} \quad \frac{x^4}{x^3 + x^2} \quad \frac{x^2 + 1000}{x^3 - 1000} \quad \frac{x \sin \frac{1}{x}}{x}.
   \]

36 If a particular \( \delta \) achieves \( |f(x) - L| < \varepsilon \), why is it OK to choose a smaller \( \delta \)?

37 The sum of \( 1 + r + r^2 + \ldots + r^{n-1} \) is \( a_n = (1 - r^n)/(1 - r) \). What is the limit of \( a_n \) as \( n \to \infty \)? For which \( r \) does the limit exist?

38 If \( a_n \to L \) prove that there is a number \( N \) with this property: If \( n > N \) and \( m > N \) then \( |a_n - a_m| < 2\varepsilon \). This is Cauchy’s test for convergence.

39 No matter what decimals come later, \( a_1 = .4, a_2 = .49, a_3 = .493, \ldots \) approaches a limit \( L \). How do we know (when we can’t know \( L \))? Cauchy’s test is passed: the \( a’s \) get closer to each other.

   (a) From \( a_4 \) onwards we have \( |a_n - a_m| < \ldots \).
   (b) After which \( a_n \) is \( |a_n - a_m| < 10^{-7} \)?

40 Choose decimals in Problem 39 so the limit is \( L = .494 \). Choose decimals so that your professor can’t find \( L \).

41 If every decimal in \( .abcde\ldots \) is picked at random from \( 0, 1, \ldots, 9 \), what is the “average” limit \( L \)?

42 If every decimal is 0 or 1 (at random), what is the average limit \( L \)?

43 Suppose \( a_n = \frac{1}{2} a_{n-1} + 4 \) and start from \( a_1 = 10 \). Find \( a_2 \) and \( a_3 \) and a connection between \( a_{n-8} \) and \( a_{n-1} \). Deduce that \( a_n \to 8 \).

44 “For every \( \delta \) there is an \( \varepsilon \) such that \( |f(x)| < \varepsilon \) if \( |x| < \delta \).” That test is twisted around. Find \( \varepsilon \) when \( f(x) = \cos x \), which does not converge to zero.

45 Prove the Squeeze Theorem for sequences, using \( \varepsilon \): If \( a_n \to L \) and \( c_n \to L \) and \( a_n \leq b_n \leq c_n \) for \( n > N \), then \( b_n \to L \).

46 Explain in 110 words the difference between “we will get there if you hurry” and “we will get there only if you hurry” and “we will get there if and only if you hurry.”
In Figure 2.20, the first function would be continuous if it had \( f(0) = 0 \). But it has \( f(0) = 1 \). After changing \( f(0) \) to the right value, the problem is gone. The discontinuity is removable. Examples 2, 3, 4 are more important and more serious. There is no "correct" value for \( f(0) \):

2. \( f(x) = \text{step function} \) (jump from 0 to 1 at \( x = 0 \))
3. \( f(x) = 1/x^2 \) (infinite limit as \( x \to 0 \))
4. \( f(x) = \sin(1/x) \) (infinite oscillation as \( x \to 0 \)).

The graphs show how the limit fails to exist. The step function has a jump discontinuity. It has one-sided limits, from the left and right. It does not have an ordinary (two-sided) limit. The limit from the left (\( x \to 0^- \)) is 0. The limit from the right (\( x \to 0^+ \)) is 1. Another step function is \( x/|x| \), which jumps from \(-1\) to 1.

In the graph of \( 1/x^2 \), the only reasonable limit is \( L = +\infty \). I cannot go on record as saying that this limit exists. Officially, it doesn't—but we often write it anyway: \( 1/x^2 \to \infty \) as \( x \to 0 \). This means that \( 1/x^2 \) goes (and stays) above every \( L \) as \( x \to 0 \).

In the same unofficial way we write one-sided limits for \( f(x) = 1/x \):

\[
\lim_{x \to 0^-} \frac{1}{x} = -\infty , \quad \lim_{x \to 0^+} \frac{1}{x} = +\infty .
\]

**Remark** 1/x has a "pole" at \( x = 0 \). So has \( 1/x^2 \) (a double pole). The function \( 1/(x^2 - x) \) has poles at \( x = 0 \) and \( x = 1 \). In each case the denominator goes to zero and the function goes to \( +\infty \) or \( -\infty \). Similarly \( 1/\sin x \) has a pole at every multiple of \( \pi \) (where \( \sin x \) is zero). Except for \( 1/x^2 \) these poles are "simple"—the functions are completely smooth at \( x = 0 \) when we multiply them by \( x \):

\[
(x)\left(\frac{1}{x}\right) = 1 \quad \text{and} \quad (x)\left(\frac{1}{x^2-x}\right) = \frac{1}{x-1} \quad \text{and} \quad (x)\left(\frac{1}{\sin x}\right)
\]

are continuous at \( x = 0 \).

\( 1/x^2 \) has a double pole, since it needs multiplication by \( x^2 \) (not just \( x \)). A ratio of polynomials \( P(x)/Q(x) \) has poles where \( Q = 0 \), provided any common factors like \( (x + 1)/(x + 1) \) are removed first.

Jumps and poles are the most basic discontinuities, but others can occur. The fourth graph shows that \( \sin(1/x) \) has no limit as \( x \to 0 \). This function does not blow up; the sine never exceeds 1. At \( x = 1/3 \) and \( 1/1000 \) it equals \( \sin 3 \) and \( \sin 4 \) and \( \sin 1000 \). Those numbers are positive and negative and (?). As \( x \) gets small and \( 1/x \) gets large, the sine oscillates faster and faster. Its graph won't stay in a small box of height \( \epsilon \), no matter how narrow the box.

**CONTINUOUS FUNCTIONS**

**Definition** \( f \) is "continuous at \( x = a \)" if \( f(a) \) is defined and \( f(x) \to f(a) \) as \( x \to a \). If \( f \) is continuous at every point where it is defined, it is a *continuous function*. 

---

**Fig. 2.20** Four types of discontinuity (others are possible) at \( x = 0 \).
2.7 Continuous Functions

Objection  The definition makes \( f(x) = \frac{1}{x} \) a continuous function! It is not defined at \( x = 0 \), so its continuity can't fail. The logic requires us to accept this, but we don't have to like it. Certainly there is no \( f(0) \) that would make \( \frac{1}{x} \) continuous at \( x = 0 \).

It is amazing but true that the definition of "continuous function" is still debated (Mathematics Teacher, May 1989). You see the reason—we speak about a discontinuity of \( \frac{1}{x} \), and at the same time call it a continuous function. The definition misses the difference between \( \frac{1}{x} \) and \( \frac{\sin x}{x} \). The function \( f(x) = \frac{\sin x}{x} \) can be made continuous at all \( x \). Just set \( f(0) = 1 \).

We call a function "continuable" if its definition can be extended to all \( x \) in a way that makes it continuous. Thus \( \frac{\sin x}{x} \) and \( \sqrt{x} \) are continuable. The functions \( \frac{1}{x} \) and \( \tan x \) are not continuable. This suggestion may not end the debate, but I hope it is helpful.

**EXAMPLE 1**  \( \sin x \) and \( \cos x \) and all polynomials \( P(x) \) are continuous functions.

**EXAMPLE 2**  The absolute value \( |x| \) is continuous. Its slope jumps (not continuable).

**EXAMPLE 3**  Any rational function \( \frac{P(x)}{Q(x)} \) is continuous except where \( Q = 0 \).

**EXAMPLE 4**  The function that jumps between 1 at fractions and 0 at non-fractions is discontinuous everywhere. There is a fraction between every pair of non-fractions and vice versa. (Somehow there are many more non-fractions.)

**EXAMPLE 5**  The function \( 0^x \) is zero for every \( x \), except that \( 0^0 \) is not defined. So define it as zero and this function is continuous. But see the next paragraph where \( 0^0 \) has to be 1.

We could fill the book with proofs of continuity, but usually the situation is clear. "A function is continuous if you can draw its graph without lifting up your pen."

At a jump, or an infinite limit, or an infinite oscillation, there is no way across the discontinuity except to start again on the other side. The function \( x^n \) is continuous for \( n > 0 \). It is not continuable for \( n < 0 \). The function \( x^0 \) equals 1 for every \( x \), except that \( 0^0 \) is not defined. This time continuity requires \( 0^0 = 1 \).

The interesting examples are the close ones—we have seen two of them:

**EXAMPLE 6**  \( \frac{\sin x}{x} \) and \( \frac{1 - \cos x}{x} \) are both continuable at \( x = 0 \).

Those were crucial for the slope of \( \sin x \). The first approaches 1 and the second approaches 0. Strictly speaking we must give these functions the correct values (1 and 0) at the limiting point \( x = 0 \)—which of course we do.

It is important to know what happens when the denominators change to \( x^2 \).

**EXAMPLE 7**  \( \frac{\sin x}{x^2} \) blows up but \( \frac{1 - \cos x}{x^2} \) has the limit \( \frac{1}{2} \) at \( x = 0 \).

Since \( \frac{\sin x}{x} \) approaches 1, dividing by another \( x \) gives a function like \( \frac{1}{x} \). There is a simple pole. It is an example of \( 0/0 \), in which the zero from \( x^2 \) is reached more quickly than the zero from \( \sin x \). The "race to zero" produces almost all interesting problems about limits.
For $1 - \cos x$ and $x^2$ the race is almost even. Their ratio is $1$ to $2$:

$$\frac{1 - \cos x}{x^2} = \frac{1 - \cos^2 x}{x^2(1 + \cos x)} = \frac{\sin^2 x}{x^2} \cdot \frac{1}{1 + \cos x} \to \frac{1}{1 + 1}$$

as $x \to 0$.

This answer $\frac{1}{2}$ will be found again (more easily) by "l'Hôpital's rule." Here I emphasize not the answer but the problem. A central question of differential calculus is to know how fast the limit is approached. The speed of approach is exactly the information in the derivative.

These three examples are all continuous at $x = 0$. The race is controlled by the slope—because $f(x) - f(0)$ is nearly $f'(0)$ times $x$:

- derivative of $\sin x$ is $1$ $\iff$ $\sin x$ decreases like $x$
- derivative of $\sin^2 x$ is $0$ $\iff$ $\sin^2 x$ decreases faster than $x$
- derivative of $x^{1/3}$ is $\infty$ $\iff$ $x^{1/3}$ decreases more slowly than $x$.

**DIFFERENTIABLE FUNCTIONS**

The absolute value $|x|$ is continuous at $x = 0$ but has no derivative. The same is true for $x^{1/3}$. Asking for a derivative is more than asking for continuity. The reason is fundamental, and carries us back to the key definitions:

- **Continuous** at $x$: $f(x + \Delta x) - f(x) \to 0$ as $\Delta x \to 0$
- **Derivative** at $x$: $\frac{f(x + \Delta x) - f(x)}{\Delta x} \to f'(x)$ as $\Delta x \to 0$.

In the first case, $\Delta f$ goes to zero (maybe slowly). In the second case, $\Delta f$ goes to zero as fast as $\Delta x$ (because $\Delta f/\Delta x$ has a limit). That requirement is stronger:

21 At a point where $f(x)$ has a derivative, the function must be continuous. But $f(x)$ can be continuous with no derivative.

**Proof** The limit of $\Delta f = (\Delta x)(\Delta f/\Delta x)$ is $0(0)(df/dx) = 0$. So $f(x + \Delta x) - f(x) \to 0$.

The continuous function $x^{1/3}$ has no derivative at $x = 0$, because $\frac{1}{3} x^{-2/3}$ blows up. The absolute value $|x|$ has no derivative because its slope jumps. The remarkable function $\frac{1}{3} \cos 3x + \frac{1}{9} \cos 9x + \cdots$ is continuous at all points and has a derivative at no points. You can draw its graph without lifting your pen (but not easily—it turns at every point). To most people, it belongs with space-filling curves and unmeasurable areas—in a box of curiosities. Fractals used to go into the same box! They are beautiful shapes, with boundaries that have no tangents. The theory of fractals is very alive, for good mathematical reasons, and we touch on it in Section 3.7.

I hope you have a clear idea of these basic definitions of calculus:

1 Limit (n $\to \infty$ or $x \to a$) 2 Continuity (at $x = a$) 3 Derivative (at $x = a$).

Those go back to $\varepsilon$ and $\delta$, but it is seldom necessary to follow them so far. In the same way that economics describes many transactions, or history describes many events, a function comes from many values $f(x)$. A few points may be special, like market crashes or wars or discontinuities. At other points $df/dx$ is the best guide to the function.
Continuous Functions

This chapter ends with two essential facts about a continuous function on a closed interval. The interval is $a \leq x \leq b$, written simply as $[a, b]$. At the endpoints $a$ and $b$ we require $f(x)$ to approach $f(a)$ and $f(b)$.

**Extreme Value Property** A continuous function on the finite interval $[a, b]$ has a maximum value $M$ and a minimum value $m$. There are points $x_{\text{max}}$ and $x_{\text{min}}$ in $[a, b]$ where it reaches those values:

$$f(x_{\text{max}}) = M \geq f(x) \geq f(x_{\text{min}}) = m \text{ for all } x \in [a, b].$$

**Intermediate Value Property** If the number $F$ is between $f(a)$ and $f(b)$, there is a point $c$ between $a$ and $b$ where $f(c) = F$. Thus if $F$ is between the minimum $m$ and the maximum $M$, there is a point $c$ between $x_{\text{min}}$ and $x_{\text{max}}$ where $f(c) = F$.

Examples show why we require closed intervals and continuous functions. For $0 < x < 1$ the function $f(x) = x$ never reaches its minimum (zero). If we close the interval by defining $f(0) = 3$ (discontinuous) the minimum is still not reached. Because of the jump, the intermediate value $F = 2$ is also not reached. The idea of continuity was inescapable, after Cauchy defined the idea of a limit.

### 2.7 Exercises

**Read-through questions**

Continuity requires the ___ of $f(x)$ to exist as $x \to a$ and to agree with ___ at $a$. The reason that $x/|x|$ is not continuous at $x = 0$ is ___. This function does have ___ limits. The reason that $1/\cos x$ is discontinuous at ___ is ___. The reason that $\cos(1/x)$ is discontinuous at $x = 0$ is ___. The function $f(x) = ___$ has a simple pole at $x = 3$, where $f^2$ has a ___ pole.

The power $x^n$ is continuous at all $x$ provided $n$ is ___. It has no derivative at $x = 0$ when $n$ is ___. $f(x) = \sin(-x)/x$ approaches ___ as $x \to 0$, so this is a ___ function provided we define $f(0) = ___$. A “continuous function” must be continuous at all ___. A “continuable function” can be extended to every point $x$ so that ___.

If $f$ has a derivative at $x = a$ then $f$ is necessarily ___ at $x = a$. The derivative controls the speed at which $f(x)$ approaches ___. On a closed interval $[a, b]$, a continuous $f$ has the ___ value property and the ___ value property. It reaches its ___ $M$ and its ___ $m$, and it takes on every value ___.

In Problems 1–20, find the numbers $c$ that make $f(x)$ into (A) a continuous function and (B) a differentiable function. In one case $f(x) \to f(a)$ at every point, in the other case $f(x) \to f(a)$ at every point.

$$1 \ f(x) = \begin{cases} \sin x & x < 1 \\ c & x \geq 1 \end{cases} \quad 2 \ f(x) = \begin{cases} \cos^3 x & x \neq \pi \\ c & x = \pi \end{cases} \quad 3 \ f(x) = \begin{cases} cx & x < 0 \\ 2cx & x \geq 0 \end{cases} \quad 4 \ f(x) = \begin{cases} cx & x < 1 \\ 2cx & x \geq 1 \end{cases} \quad 5 \ f(x) = \begin{cases} c + x & x < 0 \\ c^2 + x^2 & x \geq 0 \end{cases} \quad 6 \ f(x) = \begin{cases} x^3 & x \neq c \\ -8 & x = c \end{cases} \quad 7 \ f(x) = \begin{cases} 2x & x < c \\ x + 1 & x \geq c \end{cases} \quad 8 \ f(x) = \begin{cases} x^c & x \neq 0 \\ 0 & x = 0 \end{cases} \quad 9 \ f(x) = \begin{cases} (\sin x)/x^2 & x \neq 0 \\ c & x = 0 \end{cases} \quad 10 \ f(x) = \begin{cases} x + c & x \leq c \\ 1 & x > c \end{cases} \quad 11 \ f(x) = \begin{cases} c & x \neq 4 \\ 1/x^3 & x = 4 \end{cases} \quad 12 \ f(x) = \begin{cases} c & x \leq 0 \\ \sec x & x \geq 0 \end{cases} \quad 13 \ f(x) = \begin{cases} x^2 + c & x \neq 1 \\ 2 & x = 1 \end{cases} \quad 14 \ f(x) = \begin{cases} x^2 - 1 & x \neq c \\ 2c & x = c \end{cases} \quad 15 \ f(x) = \begin{cases} (\tan x)/x & x \neq 0 \\ c & x = 0 \end{cases} \quad 16 \ f(x) = \begin{cases} x^2 & x \leq c \\ 2x & x > c \end{cases} \quad 17 \ f(x) = \begin{cases} (c + \cos x)/x & x \neq 0 \\ 0 & x = 0 \end{cases} \quad 18 \ f(x) = |x + c|$$

†The interval $[a, b]$ is closed (endpoints included). The interval $(a, b)$ is open ($a$ and $b$ left out).

The infinite interval $[0, \infty)$ contains all $x \geq 0$. 
\[ f(x) = \begin{cases} \frac{(\sin x - x)}{x^2} & x \neq 0 \\ 0 & x = 0 \end{cases} \]

20 \( f(x) = |x^2 + c^2| \)

Construct your own \( f(x) \) with these discontinuities at \( x = 1 \).

21 Removable discontinuity

22 Infinite oscillation

23 Limit for \( x \to 1^+ \), no limit for \( x \to 1^- \)

24 A double pole

25 \[ \lim_{x \to 1^-} f(x) = 4 + \lim_{x \to 1^+} f(x) \]

26 \[ \lim_{x \to 1} f(x) = \infty \text{ but } \lim_{x \to 1} (x - 1)f(x) = 0 \]

27 \[ \lim_{x \to 1} (x - 1)f(x) = 5 \]

28 The statement \( "3x \to 7\) as \( x \to 1" \) is false. Choose an \( \varepsilon \) for which no \( \delta \) can be found. The statement \( "3x \to 3\) as \( x \to 1" \) is true. For \( \varepsilon = \frac{1}{2} \) choose a suitable \( \delta \).

29 How many derivatives \( f', f'', \ldots \) are continuable functions?
   (a) \( f = x^{3/2} \)
   (b) \( f = x^{3/2} \sin x \)
   (c) \( f = (\sin x)^{3/2} \)

30 Find one-sided limits at points where there is no two-sided limit. Give a 3-part formula for function (c).
   \[ (a) \frac{|x|}{7x} \quad (b) \sin |x| \quad (c) \frac{d}{dx} |x^2 - 1| \]

31 Let \( f(1) = 1 \) and \( f(-1) = 1 \) and \( f(x) = (x^2 - x)/(x^2 - 1) \) otherwise. Decide whether \( f \) is continuous at
   (a) \( x = 1 \)
   (b) \( x = 0 \)
   (c) \( x = -1 \).

32 Let \( f(x) = x^2 \sin 1/x \) for \( x \neq 0 \) and \( f(0) = 0 \). If the limits exist, find
   (a) \( \lim_{x \to 0} f(x) \)
   (b) \( \frac{df}{dx} \) at \( x = 0 \)
   (c) \( \lim_{x \to 0} f'(x) \).

33 If \( f(0) = 0 \) and \( f'(0) = 3 \), rank these functions from smallest to largest as \( x \) decreases to zero:
   \( f(x), \ x, \ xf(x), \ f(x) + 2x, \ 2(f(x) - x), \ (f(x))^2 \).

34 Create a discontinuous function \( f(x) \) for which \( f^2(x) \) is continuous.

35 True or false, with an example to illustrate:
   (a) If \( f(x) \) is continuous at all \( x \), it has a maximum value \( M \).

(b) If \( f(x) \leq 7 \) for all \( x \), then \( f \) reaches its maximum.

(c) If \( f(1) = 1 \) and \( f(2) = -2 \), then somewhere \( f(x) = 0 \).

(d) If \( f(1) = 1 \) and \( f(2) = -2 \) and \( f \) is continuous on \([1, 2]\), then somewhere on that interval \( f(x) = 0 \).

36 The functions \( \cos x \) and \( 2x \) are continuous. Show from the ________ property that \( \cos x = 2x \) at some point between \( 0 \) and \( 1 \).

37 Show by example that these statements are false:
   (a) If a function reaches its maximum and minimum then the function is continuous.
   (b) If \( f(x) \) reaches its maximum and minimum and all values between \( f(0) \) and \( f(1) \), it is continuous at \( x = 0 \).
   (c) (mostly for instructors) If \( f(x) \) has the intermediate value property between all points \( a \) and \( b \), it must be continuous.

38 Explain with words and a graph why \( f(x) = x \sin (1/x) \) is continuous but has no derivative at \( x = 0 \). Set \( f(0) = 0 \).

39 Which of these functions are continuable, and why?
   \[ f_1(x) = \begin{cases} \sin x & x < 0 \\ \cos x & x > 1 \end{cases} \]
   \[ f_2(x) = \begin{cases} \sin 1/x & x < 0 \\ \cos 1/x & x > 1 \end{cases} \]
   \[ f_3(x) = \begin{cases} x & \text{when } \sin x \neq 0 \\ \frac{x}{\sin x} \end{cases} \]
   \[ f_4(x) = x^0 + 0^2 \]

30 Explain the difference between a continuous function and a continuable function. Are continuous functions always continuable?

41 \( f(x) \) is any continuous function with \( f(0) = f(1) \).
   (a) Draw a typical \( f(x) \). Mark where \( f(x) = f(x + \frac{1}{2}) \).
   (b) Explain why \( g(x) = f(x + \frac{1}{2}) - f(x) \) has \( g(\frac{1}{2}) = -g(0) \).
   (c) Deducing from (b) that (a) is always possible: There must be a point where \( g(x) = 0 \) and \( f(x) = f(x + \frac{1}{2}) \).

42 Create an \( f(x) \) that is continuous only at \( x = 0 \).

43 If \( f(x) \) is continuous and \( 0 \leq f(x) \leq 1 \) for all \( x \), then there is a point where \( f(x^*) = x^* \). Explain with a graph and prove with the intermediate value theorem.

44 In the \( \varepsilon-\delta \) definition of a limit, change \( 0 < |x - a| < \delta \) to \( |x - a| < \delta \). Why is \( f(x) \) now continuous at \( x = a \)?

45 A function has a ________ at \( x = 0 \) if and only if \( (f(x) - f(0))/x \) is ________ at \( x = 0 \).
Contents

CHAPTER 1  Introduction to Calculus
1.1 Velocity and Distance 1
1.2 Calculus Without Limits 8
1.3 The Velocity at an Instant 16
1.4 Circular Motion 22
1.5 A Review of Trigonometry 29
1.6 A Thousand Points of Light 34
1.7 Computing in Calculus 36

CHAPTER 2  Derivatives
2.1 The Derivative of a Function 44
2.2 Powers and Polynomials 50
2.3 The Slope and the Tangent Line 58
2.4 Derivative of the Sine and Cosine 64
2.5 The Product and Quotient and Power Rules 71
2.6 Limits 78
2.7 Continuous Functions 85

CHAPTER 3  Applications of the Derivative
3.1 Linear Approximation 91
3.2 Maximum and Minimum Problems 96
3.3 Second Derivatives: Minimum vs. Maximum 105
3.4 Graphs 112
3.5 Ellipses, Parabolas, and Hyperbolas 121
3.6 Iterations $x_{n+1} = F(x_n)$ 130
3.7 Newton's Method and Chaos 137
3.8 The Mean Value Theorem and l'Hôpital's Rule 146
Chapter 2 concentrated on computing derivatives. This chapter concentrates on using them. Our computations produced \( \frac{dy}{dx} \) for functions built from \( x^n \) and \( \sin x \) and \( \cos x \). Knowing the slope, and if necessary also the second derivative, we can answer the questions about \( y = f(x) \) that this subject was created for:

1. How does \( y \) change when \( x \) changes?
2. What is the maximum value of \( y \)? Or the minimum?
3. How can you tell a maximum from a minimum, using derivatives?

The information in \( \frac{dy}{dx} \) is entirely local. It tells what is happening close to the point and nowhere else. In Chapter 2, \( \Delta x \) and \( \Delta y \) went to zero. Now we want to get them back. The local information explains the larger picture, because \( \Delta y \) is approximately \( \frac{dy}{dx} \times \Delta x \).

The problem is to connect the finite to the infinitesimal—the average slope to the instantaneous slope. Those slopes are close, and occasionally they are equal. Points of equality are assured by the Mean Value Theorem—which is the local-global connection at the center of differential calculus. But we cannot predict where \( \frac{dy}{dx} \) equals \( \Delta y/\Delta x \). Therefore we now find other ways to recover a function from its derivatives—or to estimate distance from velocity and acceleration.

It may seem surprising that we learn about \( y \) from \( \frac{dy}{dx} \). All our work has been going the other way! We struggled with \( y \) to squeeze out \( \frac{dy}{dx} \). Now we use \( \frac{dy}{dx} \) to study \( y \). That's life. Perhaps it really is life, to understand one generation from later generations.

### 3.1 Linear Approximation

The book started with a straight line \( f = vt \). The distance is linear when the velocity is constant. As soon as \( v \) begins to change, \( f = vt \) falls apart. Which velocity do we choose, when \( v(t) \) is not constant? The solution is to take very short time intervals,
in which \( v \) is nearly constant:

\[
\begin{align*}
\Delta f &= v \Delta t & \text{is nearly true,} \\
df &= vdt & \text{is exactly true.}
\end{align*}
\]

For a brief moment the function \( f(t) \) is linear—and stays near its tangent line.

In Section 2.3 we found the tangent line to \( y = f(x) \). At \( x = a \), the slope of the curve and the slope of the line are \( f'(a) \). For points on the line, start at \( y = f(a) \). Add the slope times the "increment" \( x - a \): 

\[
Y = f(a) + f'(a)(x - a). \tag{1}
\]

We write a capital \( Y \) for the line and a small \( y \) for the curve. The whole point of tangents is that they are close (provided we don't move too far from \( a \)):

\[
y \approx Y \quad \text{or} \quad f(x) \approx f(a) + f'(a)(x - a). \tag{2}
\]

That is the all-purpose linear approximation. Figure 3.1 shows the square root function \( y = \sqrt{x} \) and its tangent line at \( x = a = 100 \). At the point \( y = \sqrt{100} = 10 \), the slope is \( 1/2 \sqrt{x} = 1/20 \). The table beside the figure compares \( y(x) \) with \( Y(x) \).

<table>
<thead>
<tr>
<th>( x )</th>
<th>( Y )</th>
<th>( y = \sqrt{x} )</th>
</tr>
</thead>
<tbody>
<tr>
<td>100</td>
<td>10.0</td>
<td>10.0</td>
</tr>
<tr>
<td>102</td>
<td>10.1</td>
<td>10.0995</td>
</tr>
<tr>
<td>110</td>
<td>10.5</td>
<td>10.49</td>
</tr>
<tr>
<td>200</td>
<td>15.0</td>
<td>14.1</td>
</tr>
<tr>
<td>400</td>
<td>25.0</td>
<td>20.0</td>
</tr>
</tbody>
</table>

Fig. 3.1 \( Y(x) \) is the linear approximation to \( \sqrt{x} \) near \( x = a = 100 \).

The accuracy gets worse as \( x \) departs from 100. The tangent line leaves the curve. The arrow points to a good approximation at 102, and at 101 it would be even better. In this example \( Y \) is larger than \( y \)—the straight line is above the curve. The slope of the line stays constant, and the slope of the curve is decreasing. Such a curve will soon be called "concave downward," and its tangent lines are above it.

Look again at \( x = 102 \), where the approximation is good. In Chapter 2, when we were approaching \( dy/dx \), we started with \( \Delta y/\Delta x \):

\[
\text{slope} \approx \frac{\sqrt{102} - \sqrt{100}}{102 - 100}. \tag{3}
\]

Now that is turned around! The slope is 1/20. What we don't know is \( \sqrt{102} \):

\[
\sqrt{102} \approx \sqrt{100 + \text{slope}(102 - 100)}. \tag{4}
\]

You work with what you have. Earlier we didn't know \( dy/dx \), so we used (3). Now we are experts at \( dy/dx \), and we use (4). After computing \( y' = 1/20 \) once and for
all, the tangent line stays near \( \sqrt{x} \) for every number near 100. When that nearby number is \( 100 + \Delta x \), notice the error as the approximation is squared:

\[
\left( \sqrt{100 + \frac{1}{20} \Delta x} \right)^2 = 100 + \Delta x + \frac{1}{400} (\Delta x)^2.
\]

The desired answer is \( 100 + \Delta x \), and we are off by the last term involving \((\Delta x)^2\). The whole point of linear approximation is to ignore every term after \( \Delta x \).

There is nothing magic about \( x = 100 \), except that it has a nice square root. Other points and other functions allow \( y \approx Y \). I would like to express this same idea in different symbols. Instead of starting from \( a \) and going to \( x \), we start from \( x \) and go a distance \( \Delta x \) to \( x + \Delta x \). The letters are different but the mathematics is identical.

### 3A
At any point \( x \), and for any smooth function \( y = f(x) \),

\[
\text{slope at } x \approx \frac{f(x + \Delta x) - f(x)}{\Delta x}.
\]

For the approximation to \( f(x + \Delta x) \), multiply both sides by \( \Delta x \) and add \( f(x) \):

\[
f(x + \Delta x) \approx f(x) + (\text{slope at } x)(\Delta x).
\]

#### EXAMPLE 1
**An important linear approximation**: \((1 + x)^n \approx 1 + nx\) for \( x \) near zero.

#### EXAMPLE 2
**A second important approximation**: \(1/(1 + x)^n \approx 1 - nx\) for \( x \) near zero.

**Discussion**  Those are really the same. By changing \( n \) to \(-n\) in Example 1, it becomes Example 2. These are linear approximations using the slopes \( n \) and \(-n\) at \( x = 0 \):

\[
(1 + x)^n \approx 1 + (\text{slope at zero}) \times (x - 0) = 1 + nx.
\]

Here is the same thing with \( f(x) = x^n \). The basepoint in equation (6) is now 1 or \( x \):

\[
(1 + \Delta x)^n \approx 1 + n \Delta x \quad (x + \Delta x)^n \approx x^n + nx^{n-1} \Delta x.
\]

Better than that, here are numbers. For \( n = 3 \) and \(-1 \) and 100, take \( \Delta x = .01 \):

\[
(1.01)^3 \approx 1.03 \quad \frac{1}{1.01} \approx .99 \quad \left(1 + \frac{1}{100}\right)^{100} \approx 2
\]

Actually that last number is no good. The 100th power is too much. Linear approximation gives \( 1 + 100 \Delta x = 2 \), but a calculator gives \((1.01)^{100} = 2.7 \ldots\). This is close to \( e \), the all-important number in Chapter 6. The binomial formula shows why the approximation failed:

\[
(1 + \Delta x)^{100} = 1 + 100 \Delta x + \frac{(100)(99)}{(2)(1)} (\Delta x)^2 + \ldots.
\]

Linear approximation forgets the \((\Delta x)^2\) term. For \( \Delta x = 1/100 \) that error is nearly \( \frac{1}{2} \). It is too big to overlook. The exact error is \( \frac{1}{2} (\Delta x)^2 f''(c) \), where the Mean Value Theorem in Section 3.8 places \( c \) between \( x \) and \( x + \Delta x \). You already see the point:

\( y - Y \) is of order \((\Delta x)^2\). **Linear approximation, quadratic error.**

### DIFFERENTIALS

There is one more notation for this linear approximation. It has to be presented, because it is often used. The notation is suggestive and confusing at the same time—
it keeps the same symbols \(dx\) and \(dy\) that appear in the derivative. Earlier we took great pains to emphasize that \(dy/dx\) is not an ordinary fraction. Until this paragraph, \(dx\) and \(dy\) have had no independent meaning. Now they become separate variables, like \(x\) and \(y\) but with their own names. These quantities \(dx\) and \(dy\) are called differentials.

The symbols \(dx\) and \(dy\) measure changes along the tangent line. They do for the approximation \(Y(x)\) exactly what \(\Delta x\) and \(\Delta y\) did for \(y(x)\). Thus \(dx\) and \(\Delta x\) both measure distance across.

Figure 3.2 has \(\Delta x = dx\). But the change in \(y\) does not equal the change in \(Y\). One is \(\Delta y\) (exact for the function). The other is \(dy\) (exact for the tangent line). The differential \(dy\) is \(\text{equal to } \Delta Y\), the change along the tangent line. Where \(\Delta y\) is the true change, \(dy\) is its linear approximation \((dy/dx)dx\).

You often see \(dy\) written as \(f'(x)dx\).

\[
\begin{align*}
\Delta y &= \text{change in } y \text{ (along curve)} \\
\Delta y &= \text{change in } Y \text{ (along tangent)} \\
\text{Fig. 3.2 } &\text{The linear approximation to } \Delta y \text{ is } dy = f'(x)dx.
\end{align*}
\]

**EXAMPLE 3** \(y = x^2\) has \(dy/dx = 2x\) so \(dy = 2x \, dx\). The table has basepoint \(x = 2\). The prediction \(dy\) differs from the true \(\Delta y\) by exactly \((\Delta x)^2 = .01\) and \(.04\) and \(.09\).

\[
\begin{array}{cccc}
dx & dy & \Delta x & \Delta y \\
y = x^2 & .1 & 0.4 & .1 & 0.41 & \Delta y = (2 + \Delta x)^2 - 2^2 \\
dy = 4dx & .2 & 0.8 & .2 & 0.84 & \Delta y = 4\Delta x + (\Delta x)^2 \\
& .3 & 1.2 & .3 & 1.29 & \\
\end{array}
\]

The differential \(dy = f'(x)dx\) is consistent with the derivative \(dy/dx = f'(x)\). We finally have \(dy = (dy/dx)dx\), but this is not as obvious as it seems! It looks like cancellation—it is really a definition. Entirely new symbols could be used, but \(dx\) and \(dy\) have two advantages: They suggest small steps and they satisfy \(dy = f'(x)dx\). Here are three examples and three rules:

\[
\begin{align*}
d(x^n) &= nx^{n-1}dx \\
d(sin x) &= cos x \, dx \\
d(tan x) &= sec^2 x \, dx \\
d(f + g) &= df + dg \\
d(cf) &= c \, df \\
d(fg) &= f \, dg + g \, df \\
\end{align*}
\]

Science and engineering and virtually all applications of mathematics depend on linear approximation. The true function is "linearized," using its slope \(v\):

Increasing the time by \(\Delta t\) increases the distance by \(\approx v\Delta t\)

Increasing the force by \(\Delta f\) increases the deflection by \(\approx v\Delta f\)

Increasing the production by \(\Delta p\) increases its value by \(\approx v\Delta p\).

\footnote{Fraction or not, it is absolutely forbidden to cancel the \(d\)’s.}
The goal of dynamics or statics or economics is to predict this multiplier \( v \) - the derivative that equals the slope of the tangent line. The multiplier gives a local prediction of the change in the function. The exact law is nonlinear—but Ohm’s law and Hooke’s law and Newton’s law are linear approximations.

### Absolute Change, Relative Change, Percentage Change

The change \( \Delta y \) or \( \Delta f \) can be measured in three ways. So can \( \Delta x \):

- **Absolute change**
  \[
  \Delta f = \frac{\Delta y}{\Delta x} 
  \]

- **Relative change**
  \[
  \frac{\Delta f}{f(x)} = \frac{\Delta x}{x} 
  \]

- **Percentage change**
  \[
  \frac{\Delta f}{f(x)} \times 100 = \frac{\Delta x}{x} \times 100 
  \]

Relative change is often more realistic than absolute change. If we know the distance to the moon within three miles, that is more impressive than knowing our own height within one inch. Absolutely, one inch is closer than three miles. Relatively, three miles is much closer:

\[
\frac{3 \text{ miles}}{300,000 \text{ miles}} < \frac{1 \text{ inch}}{70 \text{ inches}} \quad \text{or} \quad 0.001\% < 1.4\%.
\]

**Example 4** The radius of the Earth is within 80 miles of \( r = 4000 \) miles.

(a) Find the variation \( dV \) in the volume \( V = \frac{4}{3} \pi r^3 \), using linear approximation.

(b) Compute the relative variations \( dr/r \) and \( dV/V \) and \( AV/V \).

**Solution**

The job of calculus is to produce the derivative. After \( dV/dr = 4\pi r^2 \), its work is done. The variation in volume is \( dV = 4\pi(4000)^2(80) \) cubic miles. A 2% relative variation in \( r \) gives a 6% relative variation in \( V \):

\[
\frac{dr}{r} = \frac{80}{4000} = 2\% \quad \frac{dV}{V} = \frac{4\pi(4000)^2(80)}{4\pi(4000)^3/3} = 6\%.
\]

Without calculus we need the exact volume at \( r = 4000 + 80 \) (also at \( r = 3920 \)):

\[
\Delta V = \frac{4\pi(4080)^3/3 - 4\pi(4000)^3/3}{4\pi(4000)^3/3} \approx 6.1\% 
\]

One comment on \( dV = 4\pi r^2 dr \). This is \((\text{area of sphere}) \times \text{change in radius}\). It is the volume of a thin shell around the sphere. The shell is added when the radius grows by \( dr \). The exact \( AV/V \) is 3917312/640000%, but calculus just calls it 6%.

## Exercises

**Read-through questions**

On the graph, a linear approximation is given by the __\( a \)___ line. At \( x = a \), the equation for that line is \( Y = f(a) + _b\). Near \( x = a = 10 \), the linear approximation to \( y = x^3 \) is \( Y = \quad c\). At \( x = 11 \) the exact value is \( 11^3 = d \). The approximation is \( Y = e\). In this case \( \Delta y = f \) and \( dy = g\). If we know \( \sin x \), then to estimate \( \sin(x + \Delta x) \) we add ___ _h___.

In terms of \( x \) and \( \Delta x \), linear approximation is \( f(x + \Delta x) \approx f(x) \quad i\). The error is of order \( (\Delta x)^p \) or \( (x-a)^p \) with \( p = j \). The differential \( dy \) equals __ k__ times the differential __ l___. Those movements are along the __ m__ line, where \( \Delta y \) is along the __ n__.

Find the linear approximation \( Y \) to \( y = f(x) \) near \( x = a \):

1. \( f(x) = x + x^4, \quad a = 0 \)
2. \( f(x) = 1/x, \quad a = -2 \)
Compute 7-12 within 0.01 by deciding on \( f(x) \), choosing the basepoint \( a \), and evaluating \( f(a) + f'(a)(x - a) \). A calculator shows the error.

Compute 7.2001)²

\[ 8 \cdot \cos(0.02) \]

\[ 9 \cdot \cos(0.03) \]

\[ 10 \cdot (15.99)^{1.4} \]

\[ 11 \cdot 0.98 \]

\[ 12 \cdot \sin(3.14) \]

Calculate the numerical error in these linear approximations and compare with \( \frac{1}{2}(\Delta x)^2 f''(x) \):

\[ 13 (1.01)^3 \approx 1 + 3(0.01) \]

\[ 14 \cos(0.01) \approx 1 + 0(0.01) \]

\[ 15 (\sin 0.01)^2 \approx 0 + 0(0.01) \]

\[ 16 (1.01)^{-3} \approx 1 - 3(0.01) \]

\[ 17 (1 + \frac{3}{10})^{1/6} \approx 2 \]

\[ 18 \sqrt{8.99} \approx 3 + \frac{1}{6}(-0.01) \]

Confirm the approximations 19-21 by computing \( f''(0) \):

\[ 19 \sqrt{1 - \frac{1}{4}x} \approx 1 - \frac{1}{4}x \]

\[ 20 1/sqrt{1 + x^2} \approx 1 + \frac{1}{2}x^2 \text{ (use } f = 1/sqrt{1 - u}, \text{ then put } u = x^2) \]

\[ 21 \sqrt{c^2 + x^2} \approx c + \frac{1}{2}x^2 \text{ (use } f(u) = \sqrt{c^2 + u}, \text{ then put } u = x^2) \]

\[ 22 \text{ Write down the differentials } df \text{ for } f(x) = \cos x \text{ and } (x + 1)(x - 1) \text{ and } (x^2 + 1)^2. \]

In 23-27 find the linear change \( dV \) in the volume or \( dA \) in the surface area.

23 \( dV \) if the sides of a cube change from 10 to 10.1.

24 \( dA \) if the sides of a cube change from \( x \) to \( x + dx \).

25 \( dA \) if the radius of a sphere changes by \( dr \).

26 \( dV \) if a circular cylinder with \( r = 2 \) changes height from 3 to 3.05 (recall \( V = \pi r^2 h \)).

27 \( dV \) if a cylinder of height 3 changes from \( r = 2 \) to \( r = 1.9 \).

Extra credit: What is \( dV \) if \( r \) and \( h \) both change \( (dr \text{ and } dh) \)?

28 In relativity the mass is \( m_0/\sqrt{1 - (v/c)^2} \) at velocity \( v \). By Problem 20 this is near \( m_0 + _______ \) for small \( v \). Show that the kinetic energy \( \frac{1}{2}mv^2 \) and the change in mass satisfy Einstein's equation \( e = (\Delta m)c^2 \).

29 Enter 1.1 on your calculator. Press the square root key 5 times (slowly). What happens each time to the number after the decimal point? This is because \( \sqrt{1 + x} \approx _______ \).

30 In Problem 29 the numbers you see are less than 1.05, 1.025, .... The second derivative of \( \sqrt{1 + x} \) is _______ so the linear approximation is higher than the curve.

31 Enter 0.9 on your calculator and press the square root key 4 times. Predict what will appear the fifth time and press again. You now have the _______ root of 0.9. How many decimals agree with \( 1 - \frac{1}{2}(0.1) \) ?

3.2 Maximum and Minimum Problems

Our goal is to learn about \( f(x) \) from \( df/dx \). We begin with two quick questions. If \( df/dx \) is positive, what does that say about \( f \)? If the slope is negative, how is that reflected in the function? Then the third question is the critical one:

How do you identify a maximum or minimum? Normal answer: The slope is zero.

This may be the most important application of calculus, to reach \( df/dx = 0 \).

Take the easy questions first. Suppose \( df/dx \) is positive for every \( x \) between \( a \) and \( b \). All tangent lines slope upward. The function \( f(x) \) is increasing as \( x \) goes from \( a \) to \( b \).

If \( df/dx > 0 \) then \( f(x) \) is increasing. If \( df/dx < 0 \) then \( f(x) \) is decreasing.

To define increasing and decreasing, look at any two points \( x < X \). "Increasing" requires \( f(x) < f(X) \). "Decreasing" requires \( f(x) > f(X) \). A positive slope does not mean a positive function. The function itself can be positive or negative.

EXAMPLE 1 \( f(x) = x^2 - 2x \) has slope \( 2x - 2 \). This slope is positive when \( x > 1 \) and negative when \( x < 1 \). The function increases after \( x = 1 \) and decreases before \( x = 1 \).
Fig. 3.3  Slopes are \(- + \). Slope is \(+ - + - +\) so \(f\) is up-down-up-down-up.

We say that without computing \(f(x)\) at any point! The parabola in Figure 3.3 goes down to its minimum at \(x = 1\) and up again.

**EXAMPLE 2** \(x^2 - 2x + 5\) has the same slope. Its graph is shifted up by 5, a number that disappears in \(df/dx\). All functions with slope \(2x - 2\) are parabolas \(x^2 - 2x + C\), shifted up or down according to \(C\). Some parabolas cross the \(x\) axis (those crossings are solutions to \(f(x) = 0\)). Other parabolas stay above the axis. The solutions to \(x^2 - 2x + 5 = 0\) are complex numbers and we don't see them. The special parabola \(x^2 - 2x + 1 = (x - 1)^2\) grazes the axis at \(x = 1\). It has a "double zero," where \(f(x) = df/dx = 0\).

**EXAMPLE 3** Suppose \(df/dx = (x - 1)(x - 2)(x - 3)(x - 4)\). This slope is positive beyond \(x = 4\) and up to \(x = 1\) \((df/dx = 24\) at \(x = 0)\). And \(df/dx\) is positive again between 2 and 3. At \(x = 1, 2, 3, 4\), this slope is zero and \(f(x)\) changes direction.

Here \(f(x)\) is a fifth-degree polynomial, because \(f'(x)\) is fourth-degree. The graph of \(f\) goes up-down-up-down-up. It might cross the \(x\) axis five times. It must cross at least once (like this one). When complex numbers are allowed, every fifth-degree polynomial has five roots.

You may feel that "positive slope implies increasing function" is obvious—perhaps it is. But there is still something delicate. Starting from \(df/dx > 0\) at every single point, we have to deduce \(f(X) > f(x)\) at pairs of points. That is a "local to global" question, to be handled by the Mean Value Theorem. It could also wait for the Fundamental Theorem of Calculus: The difference \(f(X) - f(x)\) equals the area under the graph of \(df/dx\). That area is positive, so \(f(X)\) exceeds \(f(x)\).

**MAXIMA AND MINIMA**

Which \(x\) makes \(f(x)\) as large as possible? Where is the smallest \(f(x)\)? Without calculus we are reduced to computing values of \(f(x)\) and comparing. With calculus, the information is in \(df/dx\).

Suppose the maximum or minimum is at a particular point \(x\). It is possible that the graph has a corner—and no derivative. But if \(df/dx\) exists, it must be zero. The tangent line is level. The parabolas in Figure 3.3 change from decreasing to increasing. The slope changes from negative to positive. At this crucial point the slope is zero.
**3C Local Maximum or Minimum** Suppose the maximum or minimum occurs at a point \( x \) inside an interval where \( f(x) \) and \( df/dx \) are defined. Then \( f'(x) = 0 \).

The word "local" allows the possibility that in other intervals, \( f(x) \) goes higher or lower. We only look near \( x \), and we use the definition of \( df/dx \).

Start with \( f(x + \Delta x) - f(x) \). If \( f(x) \) is the maximum, this difference is negative or zero. The step \( \Delta x \) can be forward or backward:

- If \( \Delta x > 0 \):
  \[
  \frac{f(x + \Delta x) - f(x)}{\Delta x} = \text{negative} \leq 0 \quad \text{and in the limit } \frac{df}{dx} \leq 0.
  \]

- If \( \Delta x < 0 \):
  \[
  \frac{f(x + \Delta x) - f(x)}{\Delta x} = \text{negative} \geq 0 \quad \text{and in the limit } \frac{df}{dx} \geq 0.
  \]

Both arguments apply. Both conclusions \( df/dx \leq 0 \) and \( df/dx \geq 0 \) are correct. Thus \( df/dx = 0 \).

Maybe Richard Feynman said it best. He showed his friends a plastic curve that was made in a special way — "no matter how you turn it, the tangent at the lowest point is horizontal." They checked it out. It was true.

Surely You're Joking, Mr. Feynman! is a good book (but rough on mathematicians).

**EXAMPLE 3** (continued) Look back at Figure 3.3b. The points that stand out are not the "ups" or "downs" but the "turns." Those are stationary points, where \( df/dx = 0 \). We see two maxima and two minima. None of them are absolute maxima or minima, because \( f(x) \) starts at \(-\infty\) and ends at \(+\infty\).

**EXAMPLE 4** \( f(x) = 4x^3 - 3x^4 \) has slope \( 12x^2 - 12x^3 \). That derivative is zero when \( x^2 \) equals \( x^3 \), at the two points \( x = 0 \) and \( x = 1 \). To decide between minimum and maximum (local or absolute), the first step is to evaluate \( f(x) \) at these stationary points. We find \( f(0) = 0 \) and \( f(1) = 1 \).

Now look at large \( x \). The function goes down to \(-\infty\) in both directions. (You can mentally substitute \( x = 1000 \) and \( x = -1000 \)). For large \( x \), \(-3x^4\) dominates \( 4x^3 \).

**Conclusion** \( f = 1 \) is an absolute maximum. \( f = 0 \) is not a maximum or minimum (local or absolute). We have to recognize this exceptional possibility, that a curve (or a car) can pause for an instant \( (f' = 0) \) and continue in the same direction. The reason is the "double zero" in \( 12x^2 - 12x^3 \), from its double factor \( x^2 \).

---

**Fig. 3.4** The graphs of \( 4x^3 - 3x^4 \) and \( x - x^{-1} \). Check rough points and endpoints.
3.2 Maximum and Minimum Problems

**EXAMPLE 5** Define \( f(x) = x + x^{-1} \) for \( x > 0 \). Its derivative \( 1 - 1/x^2 \) is zero at \( x = 1 \). At that point \( f(1) = 2 \) is the minimum value. Every combination like \( \frac{1}{3} + 3 \) or \( \frac{1}{3} + \frac{1}{2} \) is larger than \( f_{\text{min}} = 2 \). Figure 3.4 shows that the maximum of \( x + x^{-1} \) is \(+\infty\).†

**Important** The maximum always occurs at a **stationary point** (where \( df/dx = 0 \)) or a **rough point** (no derivative) or an **endpoint** of the domain. These are the three types of **critical points**. All maxima and minima occur at critical points! At every other point \( df/dx > 0 \) or \( df/dx < 0 \). Here is the procedure:

1. Solve \( df/dx = 0 \) to find the stationary points \( f(x) \).
2. Compute \( f(x) \) at every critical point—**stationary point**, **rough point**, **endpoint**.
3. Take the maximum and minimum of those critical values of \( f(x) \).

**EXAMPLE 6** (Absolute value \( f(x) = |x| \)) The minimum is zero at a rough point. The maximum is at an endpoint. There are no stationary points.

The derivative of \( y = |x| \) is never zero. Figure 3.4 shows the maximum and minimum on the interval \([-3, 2]\). This is typical of piecewise linear functions.

**Question** Could the minimum be zero when the function never reaches \( f(x) = 0 \)?

**Answer** Yes, \( f(x) = 1/(1 + x)^2 \) approaches but never reaches zero as \( x \to \infty \).

**Remark 1** \( x \to \pm \infty \) and \( f(x) \to \pm \infty \) are avoided when \( f \) is **continuous on a closed interval** \( a \leq x \leq b \). Then \( f(x) \) reaches its maximum and its minimum (Extreme Value Theorem). But \( x \to \infty \) and \( f(x) \to \infty \) are too important to rule out. You test \( x \to \infty \) by considering large \( x \). You recognize \( f(x) \to \infty \) by going above every finite value.

**Remark 2** Note the difference between critical points (specified by \( x \)) and critical values (specified by \( f(x) \)). The example \( x + x^{-1} \) had the minimum point \( x = 1 \) and the minimum value \( f(1) = 2 \).

**MAXIMUM AND MINIMUM IN APPLICATIONS**

To find a maximum or minimum, solve \( f'(x) = 0 \). The slope is zero at the top and bottom of the graph. The idea is clear—and then check rough points and endpoints. But to be honest, that is not where the problem starts.

In a real application, the first step (often the hardest) is to choose the unknown and find the function. It is we ourselves who decide on \( x \) and \( f(x) \). The equation \( df/dx = 0 \) comes in the middle of the problem, not at the beginning. I will start on a new example, with a question instead of a function.

**EXAMPLE 7** Where should you get onto an expressway for minimum driving time, if the expressway speed is 60 mph and ordinary driving speed is 30 mph?

I know this problem well—it comes up every morning. The Mass Pike goes to MIT and I have to join it somewhere. There is an entrance near Route 128 and another entrance further in. I used to take the second one, now I take the first. Mathematics should decide which is faster—some mornings I think they are maxima.

Most models are simplified, to focus on the key idea. We will allow the expressway to be entered at any point \( x \) (Figure 3.5). Instead of two entrances (a discrete problem)

---

†A good word is *approach* when \( f(x) \to \infty \). Infinity is not reached. But I still say “the maximum is \( \infty \)”.
we have a continuous choice (a calculus problem). The trip has two parts, at speeds 30 and 60:

- a distance \( \sqrt{a^2 + x^2} \) up to the expressway, in \( \sqrt{a^2 + x^2}/30 \) hours
- a distance \( b - x \) on the expressway, in \( (b - x)/60 \) hours

**Problem** Minimize \( f(x) = \text{total time} = \frac{1}{30} \sqrt{a^2 + x^2} + \frac{1}{60} (b - x) \).

We have the function \( f(x) \). Now comes calculus. The first term uses the power rule: \( \text{The derivative of } u^{1/2} \text{ is } \frac{1}{2} u^{-1/2} du/dx \). Here \( u = a^2 + x^2 \) has \( du/dx = 2x \):

\[
\frac{df}{dx} = \frac{1}{30} \left( a^2 + x^2 \right)^{-1/2} (2x) - \frac{1}{60}.
\]

To solve \( \frac{df}{dx} = 0 \), multiply by 60 and square both sides:

\[
(a^2 + x^2)^{-1/2}(2x) = 1 \quad \text{gives} \quad 2x = (a^2 + x^2)^{1/2} \quad \text{and} \quad 4x^2 = a^2 + x^2. \]

Thus \( 3x^2 = a^2 \). This yields two candidates, \( x = a/\sqrt{3} \) and \( x = -a/\sqrt{3} \). But a negative \( x \) would mean useless driving on the expressway. In fact \( f' \) is not zero at \( x = -a/\sqrt{3} \). That false root entered when we squared \( 2x \).

![Fig. 3.5](image)

Fig. 3.5 Join the freeway at \( x \) minimize the driving time \( f(x) \).

I notice something surprising. The stationary point \( x = a/\sqrt{3} \) does not depend on \( b \). The total time includes the constant \( b/60 \), which disappeared in \( df/dx \). Somehow \( b \) must enter the answer, and this is a warning to go carefully. The minimum might occur at a rough point or an endpoint. Those are the other critical points of \( f \) and our drawing may not be realistic. Certainly we expect \( x \leq b \), or we are entering the expressway beyond MIT.

Continue with calculus. Compute the driving time \( f(x) \) for an entrance at \( x^* = a/\sqrt{3} \):

\[
f(x^*) = \frac{1}{30} \sqrt{a^2 + \left(\frac{a^2}{3}\right)} + \frac{1}{60} \left( b - \frac{a}{\sqrt{3}} \right) = \sqrt{\frac{3a}{60} + \frac{b}{60}} = f^*.
\]

The square root of \( 4a^2/3 \) is \( 2a/\sqrt{3} \). We combined \( 2/30 - 1/60 = 3/60 \) and divided by \( \sqrt{3} \). **Is this stationary value \( f^* \) a minimum?** You must look also at endpoints:

- enter at \( x = 0 \): travel time is \( a \cdot 30 + b \cdot 60 = f^{**} \)
- enter at \( x = b \): travel time is \( \sqrt{a^2 + b^2} \cdot 30 = f^{***} \).
3.2 Maximum and Minimum Problems

The comparison \( f^* < f^{**} \) should be automatic. Entering at \( x = 0 \) was a candidate and calculus didn't choose it. The derivative is not zero at \( x = 0 \). It is not smart to go perpendicular to the expressway.

The second comparison has \( x = b \). We drive directly to MIT at speed 30. This option has to be taken seriously. In fact it is optimal when \( b \) is small or \( a \) is large.

This choice \( x = b \) can arise mathematically in two ways. If all entrances are between 0 and \( b \), then \( b \) is an endpoint. If we can enter beyond MIT, then \( b \) is a rough point. The graph in Figure 3.5c has a corner at \( x = b \), where the derivative jumps. The reason is that distance on the expressway is the absolute value \( |b - x| \) — never negative.

Either way \( x = b \) is a critical point. The optimal \( x \) is the smaller of \( a/\sqrt{3} \) and \( b \).

- if \( a/\sqrt{3} \leq b \): stationary point wins, enter at \( x = a/\sqrt{3} \), total time \( f^* \)
- if \( a/\sqrt{3} > b \): no stationary point, drive directly to MIT, time \( f^{**} \)

The heart of this subject is in “word problems.” All the calculus is in a few lines, computing \( f' \) and solving \( f'(x) = 0 \). The formulation took longer. Step 1 usually does:

1. Express the quantity to be minimized or maximized as a function \( f(x) \).
2. Compute \( f'(x) \), solve \( f'(x) = 0 \), check critical points for \( f_{\text{min}} \) and \( f_{\text{max}} \).

A picture of the problem (and the graph of \( f(x) \)) makes all the difference.

**EXAMPLE 7** (continued) Choose \( x \) as an angle instead of a distance. Figure 3.6 shows the triangle with angle \( x \) and side \( a \). The driving distance to the expressway is \( a \sec x \). The distance on the expressway is \( b - a \tan x \). Dividing by the speeds 30 and 60, the driving time has a nice form:

\[
f(x) = \text{total time} = \frac{a \sec x}{30} - \frac{b - a \tan x}{60}.
\]

(3)

The derivatives of \( \sec x \) and \( \tan x \) go into \( df/dx \):

\[
\frac{df}{dx} = \frac{a}{30} \sec x \tan x - \frac{a}{60} \sec^2 x.
\]

(4)

Now set \( df/dx = 0 \), divide by \( a \), and multiply by \( 30 \cos^2 x \):

\[
\sin x = \frac{1}{3},
\]

(5)

This answer is beautiful. The angle \( x \) is 30°! That optimal angle \( (\pi/6 \text{ radians}) \) has \( \sin x = \frac{1}{3} \). The triangle with side \( a \) and hypotenuse \( a/\sqrt{3} \) is a 30 60 90 right triangle.

I don't know whether you prefer \( \sqrt{a^2 + x^2} \) or trigonometry. The minimum is exactly as before — either at 30° or going directly to MIT.

![Fig. 3.6](a) Driving at angle \( x \). (b) Energies of spring and mass. (c) Profit = income — cost.)
EXAMPLE 8  In mechanics, nature chooses minimum energy. A spring is pulled down by a mass, the energy is \( f(x) \), and \( df/dx = 0 \) gives equilibrium. It is a philosophical question why so many laws of physics involve minimum energy or minimum time—which makes the mathematics easy.

The energy has two terms—for the spring and the mass. The spring energy is \( \frac{1}{2}kx^2 \)—positive in stretching \((x > 0)\) is downward) and also positive in compression \((x < 0)\). The potential energy of the mass is taken as \(-mx\)—decreasing as the mass goes down. The balance is at the minimum of \( f(x) = \frac{1}{2}kx^2 - mx \).

I apologize for giving you such a small problem, but it makes a crucial point. *When \( f(x) \) is quadratic, the equilibrium equation \( df/dx = 0 \) is linear.*

\[
\frac{df}{dx} = kx - m = 0.
\]

Graphically, \( x = m/k \) is at the bottom of the parabola. Physically, \( kx = m \) is a balance of forces—the spring force against the weight. *Hooke’s law* for the spring force is elastic constant \( k \) times displacement \( x \).

EXAMPLE 9  *Derivative of cost = marginal cost* (our first management example).

The paper to print \( x \) copies of this book might cost \( C = 1000 + 3x \) dollars. The derivative is \( dC/dx = 3 \). This is the *marginal cost* of paper for each additional book. If \( x \) increases by one book, the cost \( C \) increases by \$3. The marginal cost is like the velocity and the total cost is like the distance.

*Marginal cost is in dollars per book. Total cost is in dollars.* On the plus side, the income is \( I(x) \) and the marginal income is \( dI/dx \). To apply calculus, we overlook the restriction to whole numbers.

Suppose the number of books increases by \( dx \). The cost goes up by \((dC/dx)dx \). The income goes up by \((dI/dx)dx \). If we skip all other costs, then profit \( P(x) = \text{income} I(x) - \text{cost} C(x) \). In most cases \( P \) increases to a maximum and falls back.

At the high point on the profit curve, the *marginal profit is zero*:

\[
dP/dx = 0 \quad \text{or} \quad dI/dx = dC/dx.
\]  

Profit is maximized when marginal income \( I' \) equals marginal cost \( C' \).

This basic rule of economics comes directly from calculus, and we give an example:

\[
C(x) = \text{cost of } x \text{ advertisements} = 900 + 400x - x^2
\]

\[
\text{setup cost 900, print cost 400x, volume savings } x^2
\]

\[
I(x) = \text{income due to } x \text{ advertisements} = 600x - 6x^2
\]

\[
\text{sales 600 per advertisement, subtract } 6x^2 \text{ for diminishing returns}
\]

*optimal decision* \( dC/dx = dI/dx \) or \( 400 - 2x = 600 - 12x \) or \( x = 20 \)

\[
\text{profit} = \text{income} - \text{cost} = 9600 - 8500 = 1100.
\]

The next section shows how to verify that this profit is a maximum not a minimum.

The first exercises ask you to solve \( df/dx = 0 \). Later exercises also look for \( f(x) \).

\[\text{Maybe } dx \text{ is a differential calculus book. I apologize for that.}\]
3.2 Maximum and Minimum Problems

3.2 EXERCISES

Read-through questions

If \( df/dx > 0 \) in an interval then \( f(x) \) is ____ a _____. If a maximum or minimum occurs at \( x \) then \( f'(x) = b _____. \) Points where \( f'(x) = 0 \) are called __ c ____ points. The function \( f(x) = 3x^2 - x \) has a (minimum) (maximum) at \( x = d _____. \) A stationary point that is not a maximum or minimum occurs for \( f(x) = e _____. \)

Extreme values can also occur where ___ f ____ is not defined or at the ___ g ____ of the domain. The minima of \( |x| \) and \( 5x \) for \(-2 \leq x \leq 2\) are at \( x = h ____ \) and \( x = i ____ \), even though \( df/dx \) is not zero. \( x^* \) is an absolute ___ j ____ when \( f(x^*) \geq f(x) \) for all \( x \). A ___ k ____ minimum occurs when \( f(x^*) < f(x) \) for all \( x \) near \( x^* \).

Find the stationary points and rough points and endpoints. Decide whether each point is a local or absolute minimum or maximum.

1. \( f(x) = x^2 + 4x + 5 \), \(-\infty < x < \infty\)
2. \( f(x) = x^3 - 12x \), \(-\infty < x < \infty\)
3. \( f(x) = x^2 + 3 \), \(-1 < x < 4\)
4. \( f(x) \) = \( x^2 + 1 \) for \( x \leq 1 \), \( x^2 - 4x + 5 \) for \( x > 1 \).
5. \( f(x) = 202x^2 - 9x \), \(-\infty < x < \infty\)
6. \( f(x) = x \) for \( 0 < x < 10 \)
7. \( f(x) = x^2 + \sqrt{x} \), \( 0 < x < 2\)
8. \( f(x) = x^2 - 2x + 3 \), \(-1 < x < 1\)
9. \( f(x) = x + \sin x \), \( 0 < x < 2\)
10. \( f(x) = 4 \) for \( x \leq 1 \), \( x^2 - 4x + 5 \) for \( x > 1 \).
11. \( f(x) = x \) for \( 0 < x < \pi \)
12. \( f(x) = x^2 + x \) for \( -1 < x < 2\)
13. \( f(x) \) = distance from \( x > 0 \) to nearest whole number
14. \( f(x) \) = distance from \( x > 0 \) to nearest prime number
15. \( f(x) = |x + 1| + |x - 1| \), \(-3 \leq x \leq 2\)
16. \( f(x) = x^2 + x \) for \( 0 \leq x \leq 1 \)
17. \( f(x) = x^2/2 - x^3/2 \) for \( 0 \leq x \leq 4 \)
18. \( f(x) = \sin x + \cos x \) for \( 0 \leq x \leq 2\pi \)
19. \( f(x) = x + \sin x \) for \( 0 \leq x \leq 2\pi \)
20. \( f(x) = \cos^2 \theta \sin \theta \) for \( -\pi < \theta < \pi \)
21. \( f(x) = 4 \) sin \( \theta - 3 \) cos \( \theta \) for \( 0 \leq \theta \leq 2\pi \)
22. \( f(x) = x^2 + 1 \) for \( x \leq 1 \), \( x^2 - 4x + 5 \) for \( x > 1 \).

In applied problems, choose metric units if you prefer.

23. The airlines accept a box if length + width + height = \( l + w + h \leq 62" \) or 158 cm. If \( h \) is fixed show that the maximum volume \( (62 - w - h)wh \) is \( V = h(31 - \frac{h}{2})^2 \). Choose \( h \) to maximize \( V \). The box with greatest volume is a _____.

24. If a patient's pulse measures 70, then 80, then 120, what least squares value minimizes \( (x - 70)^2 + (x - 80)^2 + (x - 120)^2 \)? If the patient got nervous, assign 120 a lower weight and minimize \( (x - 70)^2 + (x - 80)^2 + \frac{1}{2}(x - 120)^2 \).

25. At speed \( v \), a truck uses \( av + \frac{(b + w)/a}{a} \) gallons of fuel per mile. How many miles per gallon at speed \( v \)? Minimize the fuel consumption. Maximize the number of miles per gallon.

26. A limousine gets \( 120(2x)/5 \) miles per gallon. The chauffeur costs $10/hour, the gas costs $1/gallon.

(a) Find the cost per mile at speed \( v \).
(b) Find the cheapest driving speed.

27. You should shoot a basketball at the angle \( \theta \) requiring minimum speed. Avoid line drives and rainbows. Shooting from \((0,0)\) with the basket at \((a, b)\), minimize \( J(\theta) = 1/(a \sin \theta \cos \theta - b \cos^2 \theta) \).

(a) If \( b = 0 \) you are level with the basket. Show that \( \theta = 45° \) is best (Jabbar sky hook).
(b) Reduce \( df/d\theta = 0 \) to \( \tan 2\theta = -a/b \). Solve when \( a = b \).
(c) Estimate the best angle for a free throw.

The same angle allows the largest margin of error (Sports Science by Peter Brancazio). Section 12.2 gives the flight path.

28. On the longest and shortest days, in June and December, why does the length of day change the least?

29. Find the shortest \( Y \) connecting \( P, Q, \) and \( B \) in the figure. Originally \( B \) was a birdfeeder. The length of \( Y \) is \( L(x) = \sqrt{x^2 + a^2} \).

(a) Choose \( x \) to minimize \( L \) (not allowing \( x > b \)).
(b) Show that the center of the \( Y \) has 120° angles.
(c) The best \( Y \) becomes a \( V \) when \( a/b = _____. \)

30. If the distance function is \( f(t) = (1 + 3t)/(1 + 3t^2) \), when does the forward motion end? How far have you traveled? Extra credit: Graph \( f(t) \) and \( df/dt \).
In 31–34, we make and sell \( x \) pizzas. The income is \( R(x) = ax + bx^2 \) and the cost is \( C(x) = c + dx + ex^2 \).

31 The profit is \( \Pi(x) = \frac{R(x)}{x} \). The average profit per pizza is \( \frac{\Pi(x)}{x} \). The marginal profit per additional pizza is \( \frac{d\Pi(x)}{dx} \). We should maximize the (profit) (average profit) (marginal profit).

32 We receive \( R(x) = ax + bx^2 \) when the price per pizza is \( p(x) = \frac{ax + bx^2}{x} \). In reverse: When the price is \( p \) we sell \( x = \frac{R(p)}{p} \) pizzas (a function of \( p \)). We expect \( b < 0 \) because

33 Find \( x \) to maximize the profit \( \Pi(x) \). At that \( x \) the marginal profit is \( \frac{d\Pi(x)}{dx} \). We should maximize the

34 Figure B shows \( R(x) = 3x - x^2 \) and \( C_1(x) = 1 + x^2 \) and \( C_2(x) = 2 + x^2 \). With cost \( C_2 \), which sales \( x \) makes a profit? Which \( x \) makes the most profit? With higher fixed cost in \( C_2 \), the best plan is

The cookie box and popcorn box were created by Kay Dundas from a 12" x 12" square. A box with no top is a calculus classic.

35 Choose \( x \) to find the maximum volume of the cookie box.
36 Choose \( x \) to maximize the volume of the popcorn box.
37 A high-class chocolate box adds a strip of width \( x \) down across the front of the cookie box. Find the new volume \( V(x) \) and the \( x \) that maximizes it. Extra credit: Show that \( V_{\text{max}} \) is reduced by more than 20%.
38 For a box with no top, cut four squares of side \( x \) from the corners of the 12" square. Fold up the sides so the height is \( x \). Maximize the volume.

Geometry provides many problems, more applied than they seem.
39 A wire four feet long is cut in two pieces. One piece forms a circle of radius \( r \), the other forms a square of side \( x \). Choose \( r \) to minimize the sum of their areas. Then choose \( r \) to maximize.

40 A fixed wall makes one side of a rectangle. We have 200 feet of fence for the other three sides. Maximize the area \( A \) in 4 steps:

1. Draw a picture of the situation.
2. Select one unknown quantity as \( x \) (but not \( A \)).
3. Find all other quantities in terms of \( x \).
4. Solve \( \frac{dA}{dx} = 0 \) and check endpoints.

In 53–56, \( x \) and \( y \) are nonnegative numbers with \( x + y = 10 \).
Maximize and minimize:

53 \( xy \), 54 \( x^2 + y^2 \), 55 \( y - (1/x) \), 56 \( \sin x \sin y \)
57 Find the total distance \( f(x) \) from \( A \) to \( X \) to \( C \). Show that \( df/dx = 0 \) leads to \( \sin a = \sin c \). Light reflects at an equal angle to minimize travel time.
Fermat’s principle says that light travels from A to B on the quickest path. Its velocity above the x axis is v and below the x axis is w.

(a) Find the time $T(x)$ from A to X to B. On AX, time $= \frac{\text{distance}}{\text{velocity}} = \frac{\sqrt{r^2 + x^2}}{v}$.

(b) Find the equation for the minimizing x.

(c) Deduce Snell’s law $(\sin a)/v = (\sin b)/w$.

“Closest point problems” are models for many applications.

59 Where is the parabola $y = x^2$ closest to $x = 0$, $y = 2$?

60 Where is the line $y = 5 - 2x$ closest to $(0, 0)$?

61 What point on $y = -x^2$ is closest to what point on $y = 5 - 2x$? At the nearest points, the graphs have the same slope. Sketch the graphs.

62 Where is $y = x^2$ closest to $(0, \frac{1}{2})$? Minimizing $x^2 + (y - \frac{1}{2})^2 + y = (y - \frac{1}{2})^2$ gives $y < 0$. What went wrong?

63 Draw the line $y = mx$ passing near $(2, 3), (1, 1)$, and $(-1, 1)$. For a least squares fit, minimize $(3 - 2m)^2 + (1 - m)^2 + (1 + m)^2$.

Problems 66–72 go back to the theory.

66 To find where the graph of $y(x)$ has greatest slope, solve _______. For $y = 1/(1 + x^2)$ this point is _______.

67 When the difference between $f(x)$ and $g(x)$ is smallest, their slopes are _______. Show this point on the graphs of $f = 2 + x^2$ and $g = 2x - x^2$.

68 Suppose $y$ is fixed. The minimum of $x^2 + xy - y^2$ (a function of $x$) is $m(y) = \______$. Find the maximum of $m(y)$.

Now $x$ is fixed. The maximum of $x^2 + xy - y^2$ (a function of $y$) is $M(x) = \______$. Find the minimum of $M(x)$.

69 For each $m$ the minimum value of $f(x) - mx$ occurs at $x = m$. What is $f(x)$?

70 $y = x + 2x^2 \sin(1/x)$ has slope 1 at $x = 0$. But show that $y$ is not increasing on an interval around $x = 0$, by finding points where $dy/dx = 1 - 2 \cos(1/x) + 4x \sin(1/x)$ is negative.

71 True or false, with a reason: Between two local minima of a smooth function $f(x)$ there is a local maximum.

72 Create a function $y(x)$ that has its maximum at a rough point and its minimum at an endpoint.

73 Draw a circular pool with a lifeguard on one side and a drowner on the opposite side. The lifeguard swims with velocity $v$ and runs around the rest of the pool with velocity $w = 10v$. If the swim direction is at angle $\theta$ with the direct line, choose $\theta$ to minimize and maximize the arrival time.

3.3 Second Derivatives: Bending and Acceleration

When $f'(x)$ is positive, $f(x)$ is increasing. When $dy/dx$ is negative, $y(x)$ is decreasing. That is clear, but what about the second derivative? From looking at the curve, can you decide the sign of $f''(x)$ or $d^2y/dx^2$? The answer is yes and the key is in the bending.

A straight line doesn’t bend. The slope of $y = mx + b$ is $m$ (a constant). The second derivative is zero. We have to go to curves, to see a changing slope. Changes in the derivative show up in $f''(x)$:

- $f = x^2$ has $f' = 2x$ and $f'' = 2$ (this parabola bends up)
- $y = \sin x$ has $dy/dx = \cos x$ and $d^2y/dx^2 = -\sin x$ (the sine bends down)
The slope $2x$ gets larger even when the parabola is falling. The sign of $f'$ or $f''$ is not revealed by $f''$. The second derivative tells about change in slope.

A function with $f''(x) > 0$ is concave up. It bends upward as the slope increases. It is also called convex. A function with decreasing slope—this means $f''(x) < 0$—is concave down. Note how $\cos x$ and $1 + \cos x$ and even $1 + \frac{1}{2}x + \cos x$ change from concave down to concave up at $x = \pi/2$. At that point $f'' = -\cos x$ changes from negative to positive. The extra $1 + \frac{1}{2}x$ tilts the graph but the bending is the same.

Here is another way to see the sign of $f''$. Watch the tangent lines. When the curve is concave up, the tangent stays below it. A linear approximation is too low. This section computes a quadratic approximation—which includes the term with $f'' > 0$. When the curve bends down ($f'' < 0$), the opposite happens—the tangent lines are above the curve. The linear approximation is too high, and $f''$ lowers it.

In physical motion, $f''(t)$ is the acceleration—in units of distance/(time)$^2$. Acceleration is rate of change of velocity. The oscillation $\sin 2t$ has $v = 2 \cos 2t$ (maximum speed 2) and $a = -4 \sin 2t$ (maximum acceleration 4).

An increasing population means $f' > 0$. An increasing growth rate means $f'' > 0$. Those are different. The rate can slow down while the growth continues.

**MAXIMUM VS. MINIMUM**

Remember that $f'(x) = 0$ locates a stationary point. That may be a minimum or a maximum. The second derivative decides! Instead of computing $f(x)$ at many points, we compute $f''(x)$ at one point—the stationary point. It is a minimum if $f''(x) > 0$.

<table>
<thead>
<tr>
<th>3D</th>
<th>When $f'(x) = 0$ and $f''(x) &gt; 0$, there is a local minimum at $x$.</th>
</tr>
</thead>
<tbody>
<tr>
<td></td>
<td>When $f'(x) = 0$ and $f''(x) &lt; 0$, there is a local maximum at $x$.</td>
</tr>
</tbody>
</table>

To the left of a minimum, the curve is falling. After the minimum, the curve rises. The slope has changed from negative to positive. The graph bends upward and $f''(x) > 0$.

At a maximum the slope drops from positive to negative. In the exceptional case, when $f'(x) = 0$ and also $f''(x) = 0$, anything can happen. An example is $x^3$, which pauses at $x = 0$ and continues up (its slope is $3x^2 \geq 0$). However $x^4$ pauses and goes down (with a very flat graph).

We emphasize that the information from $f'(x)$ and $f''(x)$ is only "local." To be certain of an absolute minimum or maximum, we need information over the whole domain.
3.3 Second Derivatives: Bending and Acceleration

EXAMPLE 1 \( f(x) = x^3 - x^2 \) has \( f'(x) = 3x^2 - 2x \) and \( f''(x) = 6x - 2 \).

To find the maximum and/or minimum, solve \( 3x^2 - 2x = 0 \). The stationary points are \( x = 0 \) and \( x = \frac{2}{3} \). At those points we need the second derivative. It is \( f''(0) = -2 \) (local maximum) and \( f''(\frac{2}{3}) = +2 \) (local minimum).

Between the maximum and minimum is the inflection point. That is where \( f''(x) = 0 \). The curve changes from concave down to concave up. This example has \( f''(x) = 6x - 2 \), so the inflection point is at \( x = \frac{1}{3} \).

INFLECTION POINTS

In mathematics it is a special event when a function passes through zero. When the function is \( f \), its graph crosses the axis. When the function is \( f' \), the tangent line is horizontal. When \( f'' \) goes through zero, we have an inflection point.

The direction of bending changes at an inflection point. Your eye picks that out in a graph. For an instant the graph is straight (straight lines have \( f'' = 0 \)). It is easy to see crossing points and stationary points and inflection points. Very few people can recognize where \( f'' = 0 \) or \( f''' = 0 \). I am not sure if those points have names.

There is a genuine maximum or minimum when \( f'(x) \) changes sign. Similarly, there is a genuine inflection point when \( f''(x) \) changes sign. The graph is concave down on one side of an inflection point and concave up on the other side. The tangents are above the curve on one side and below it on the other side. At an inflection point, the tangent line crosses the curve (Figure 3.7b).

Notice that a parabola \( y = ax^2 + bx + c \) has no inflection points: \( y'' \) is constant. A cubic curve has one inflection point, because \( f'' \) is linear. A fourth-degree curve might or might not have inflection points—the quadratic \( f''(x) \) might or might not cross the axis.

EXAMPLE 2 \( x^4 - 2x^2 \) is W-shaped, \( 4x^3 - 4x \) has two bumps, \( 12x^2 - 4 \) is U-shaped.

The table shows the signs at the important values of \( x \):

<table>
<thead>
<tr>
<th>( x )</th>
<th>(-\sqrt{2})</th>
<th>(-1)</th>
<th>(-1/\sqrt{3})</th>
<th>0</th>
<th>(1/\sqrt{3})</th>
<th>1</th>
<th>(\sqrt{2})</th>
</tr>
</thead>
<tbody>
<tr>
<td>( f(x) )</td>
<td>0</td>
<td>-</td>
<td>-</td>
<td>0,0</td>
<td>-</td>
<td>-</td>
<td>0</td>
</tr>
<tr>
<td>( f'(x) )</td>
<td>0</td>
<td>+</td>
<td>0</td>
<td>-</td>
<td>0</td>
<td></td>
<td></td>
</tr>
<tr>
<td>( f''(x) )</td>
<td>0</td>
<td>-</td>
<td>0</td>
<td></td>
<td></td>
<td></td>
<td></td>
</tr>
</tbody>
</table>

Between zeros of \( f(x) \) come zeros of \( f'(x) \) (stationary points). Between zeros of \( f'(x) \) come zeros of \( f''(x) \) (inflection points). In this example \( f(x) \) has a double zero at the origin, so a single zero of \( f'' \) is caught there. It is a local maximum, since \( f''(0) < 0 \).

Inflection points are important—not just for mathematics. We know the world population will keep rising. We don’t know if the rate of growth will slow down. Remember: The rate of growth stops growing at the inflection point. Here is the 1990 report of the UN Population Fund.

The next ten years will decide whether the world population trebles or merely doubles before it finally stops growing. This may decide the future of the earth as a habitation for humans. The population, now 5.3 billion, is increasing by a quarter of a million every day. Between 90 and 100 million people will be added every year.

\[ \text{†That rules out } f(x) = x^4, \text{ which has } f'' = 12x^2 > 0 \text{ on both sides of zero. Its tangent line is the } x \text{ axis. The line stays below the graph—so no inflection point.} \]
3 Applications of the Derivative

during the 1990s; a billion people—a whole China—over the decade. The fastest
growth will come in the poorest countries.

A few years ago it seemed as if the rate of population growth was slowing everywhere except in Africa and parts of South Asia. The world’s population seemed set to stabilize around 10.2 billion towards the end of the next century.

Today, the situation looks less promising. The world has overshot the marker points of the 1984 “most likely” medium projection. It is now on course for an eventual total that will be closer to 11 billion than to 10 billion.

If fertility reductions continue to be slower than projected, the mark could be missed again. In that case the world could be headed towards a total of up to 14 billion people.

Starting with a census, the UN follows each age group in each country. They estimate the death rate and fertility rate—the medium estimates are published. This report is saying that we are not on track with the estimate.

Section 6.5 will come back to population, with an equation that predicts 10 billion. It assumes we are now at the inflection point. But China’s second census just started on July 1, 1990. When it’s finished we will know if the inflection point is still ahead.

You now understand the meaning of \( f''(x) \). Its sign gives the direction of bending—the change in the slope. The rest of this section computes how much the curve bends—using the size of \( f'' \) and not just its sign. We find quadratic approximations based on \( f''(x) \). In some courses they are optional—the main points are highlighted.

**CENTERED DIFFERENCES AND SECOND DIFFERENCES**

Calculus begins with average velocities, computed on either side of \( x \):

\[
\frac{f(x + \Delta x) - f(x)}{\Delta x} \quad \text{and} \quad \frac{f(x) - f(x - \Delta x)}{\Delta x} \quad \text{are close to} \quad f'(x). \quad (1)
\]

We never mentioned it, but a better approximation to \( f'(x) \) comes from averaging those two averages. This produces a **centered difference**, which is based on \( x + \Delta x \) and \( x - \Delta x \). It divides by \( 2\Delta x \):

\[
f'(x) \approx \frac{1}{2} \left[ \frac{f(x + \Delta x) - f(x)}{\Delta x} + \frac{f(x) - f(x - \Delta x)}{\Delta x} \right] = \frac{f(x + \Delta x) - f(x - \Delta x)}{2\Delta x}. \quad (2)
\]

We claim this is better. The test is to try it on powers of \( x \).

For \( f(x) = x \) these ratios all give \( f' = 1 \) (exactly). For \( f(x) = x^2 \), only the centered difference correctly gives \( f' = 2x \). The one-sided ratio gave \( 2x + \Delta x \) (in Chapter 1, it was \( 2t + h \)). It is only “first-order accurate.” But centering leaves no error. We are averaging \( 2x - \Delta x \) with \( 2x + \Delta x \). Thus the centered difference is “second-order accurate.”

I ask now: What ratio converges to the second derivative? One answer is to take differences of the first derivative. Certainly \( \Delta f'/\Delta x \) approaches \( f'' \). But we want a ratio involving \( f \) itself. A natural idea is to take differences of differences, which brings us to “second differences”:

\[
\frac{f(x + \Delta x) - f(x)}{\Delta x} \quad \frac{f(x) - f(x - \Delta x)}{\Delta x} = \frac{f(x + \Delta x) - 2f(x) + f(x - \Delta x)}{(\Delta x)^2} \to \frac{d^2f}{dx^2}. \quad (3)
\]

†The United Nations watches the second derivative!
On the top, the difference of the difference is $A^2f$. It corresponds to $d^2f$. On the bottom, $(Ax)^2$ corresponds to $dx^2$. This explains the way we place the 2's in $d^2f/dx^2$. To say it differently: $dx$ is squared, $df$ is not squared—as in distance/(time)$^2$.

Note that $(Ax)^2$ becomes much smaller than $Ax$. If we divide $df$ by $(Ax)^2$, the ratio blows up. It is the extra cancellation in the second difference $A^2f$ that allows the limit to exist. That limit is $f''(x)$.

**Application**  

The great majority of differential equations can't be solved exactly. A typical case is $f''(x) = -\sin f(x)$ (the pendulum equation). To compute a solution, I would replace $f''(x)$ by the second difference in equation (3). Approximations at points spaced by $Ax$ are a very large part of scientific computing.

To test the accuracy of these differences, here is an experiment on $f(x) = \sin x + \cos x$. The table shows the errors at $x = 0$ from formulas (1), (2), (3):

<table>
<thead>
<tr>
<th>step length $\Delta x$</th>
<th>one-sided errors</th>
<th>centered errors</th>
<th>second difference errors</th>
</tr>
</thead>
<tbody>
<tr>
<td>1/4</td>
<td>.1347</td>
<td>.0104</td>
<td>-.0052</td>
</tr>
<tr>
<td>1/8</td>
<td>.0650</td>
<td>.0026</td>
<td>-.0013</td>
</tr>
<tr>
<td>1/16</td>
<td>.0319</td>
<td>.0007</td>
<td>-.0003</td>
</tr>
<tr>
<td>1/32</td>
<td>.0158</td>
<td>.0002</td>
<td>-.0001</td>
</tr>
</tbody>
</table>

The one-sided errors are cut in half when $\Delta x$ is cut in half. The other columns decrease like $(\Delta x)^2$. Each reduction divides those errors by 4. **The errors from one-sided differences are $O(\Delta x)$ and the errors from centered differences are $O(\Delta x)^2$.**

The "big $O$" notation  

When the errors are of order $\Delta x$, we write $E = O(\Delta x)$. This means that $E \leq C\Delta x$ for some constant $C$. We don't compute $C$—in fact we don't want to deal with it. The statement "one-sided errors are $Ox$ of delta $x$" captures what is important. The main point of the other columns is $E = O(\Delta x)^2$.

**LINEAR APPROXIMATION VS. QUADRATIC APPROXIMATION**

The second derivative gives a tremendous improvement over linear approximation $f(a) + f'(a)(x - a)$. A tangent line starts out close to the curve, but the line has no way to bend. After a while it overshoots or undershoots the true function (see Figure 3.8). That is especially clear for the model $f(x) = x^2$, when the tangent is the $x$ axis and the parabola curves upward.

You can almost guess the term with bending. It should involve $f''$, and also $(\Delta x)^2$. It might be exactly $f''(x)$ times $(\Delta x)^2$ but it is not. The model function $x^2$ has $f'' = 2$. There must be a factor $\frac{1}{2}$ to cancel that 2:

$$f(x) \approx f(a) + f'(a)(x - a) + \frac{1}{2} f''(a)(x - a)^2.$$  \hspace{1cm} (4)

At the basepoint this is $f(a) = f(a)$. The derivatives also agree at $x = a$. Furthermore the second derivatives agree. On both sides of (4), the second derivative at $x = a$ is $f''(a)$.

The quadratic approximation bends with the function. It is not the absolutely final word, because there is a cubic term $\frac{1}{6} f'''(a)(x - a)^3$ and a fourth-degree term $\frac{1}{24} f''''(a)(x - a)^4$ and so on. The whole infinite sum is a "Taylor series." Equation (4) carries that series through the quadratic term—which for practical purposes gives a terrific approximation. You will see that in numerical experiments.
Applications of the Derivative

Two things to mention. First, equation (4) shows why \( f'' > 0 \) brings the curve above the tangent line. The linear part gives the line, while the quadratic part is positive and bends upward. Second, equation (4) comes from (2) and (3). Where one-sided differences give \( f(x + Ax) \approx f(x) + f'(x) Ax + \frac{1}{2} f''(x)(Ax)^2 \), centered differences give the quadratic:

\[
\begin{align*}
\text{from (2):} & \quad f(x + Ax) \approx f(x - Ax) + 2f'(x) Ax + f''(x)(Ax)^2, \\
\text{from (3):} & \quad f(x + Ax) = f(x) + f'(x) Ax + \frac{1}{2} f''(x)(Ax)^2.
\end{align*}
\]

Add and divide by 2. The result is \( f(x + Ax) \approx f(x) + f'(x) Ax + \frac{1}{2} f''(x)(Ax)^2 \). This is correct through \((Ax)^2\) and misses by \((Ax)^3\), as examples show:

**EXAMPLE 3** \((x + Ax)^3 \approx (x^3) + (3x^2)(Ax) + \frac{1}{2}(6x)(Ax)^2 + \text{error } (Ax)^3\).

**EXAMPLE 4** \((1 + x)^n \approx 1 + nx + \frac{1}{2}n(n - 1)x^2\).

The first derivative at \( x = 0 \) is \( n \). The second derivative is \( n(n - 1) \). The cubic term would be \( \frac{1}{6}n(n - 1)(n - 2)x^3 \). We are just producing the binomial expansion!

**EXAMPLE 5** \(\frac{1}{1 - x} \approx 1 + x + x^2 = \text{start of a geometric series.}\)

\(1/(1 - x)\) has derivative \(1/(1 - x^2)\). Its second derivative is \(2/(1 - x)^2\). At \( x = 0 \) those equal 1,1,2. The factor \(\frac{1}{2}\) cancels the 2, which leaves 1,1,1. This explains \(1 + x + x^2\).

The next terms are \(x^3\) and \(x^4\). The whole series is \(1/(1 - x) = 1 + x + x^2 + x^3 + \cdots\).

**Numerical experiment** \(1/\sqrt{1 + x} \approx 1 - \frac{1}{2}x + \frac{5}{8}x^2\) is tested for accuracy. Dividing \(x\) by 2 almost divides the error by 8. If we only keep the linear part \(1 - \frac{1}{2}x\), the error is only divided by 4. Here are the errors at \(x = \frac{1}{2}, \frac{1}{4}, \frac{1}{8}\): 

- Linear approximation \(\text{error } \approx \frac{3}{8}x^2\): \(0.194 \quad 0.0053 \quad 0.0014\)
- Quadratic approximation \(\text{error } \approx \frac{5}{16}x^3\): \(-0.00401 \quad -0.00055 \quad -0.00007\)

### 3.3 Exercises

**Read-through questions**

The direction of bending is given by the sign of \(a\). If the second derivative is \(b\) in an interval, the function is concave up (or convex). The graph bends \(c\). The tangent lines are \(d\) the graph. If \(f''(x) < 0\) then the graph is concave \(e\), and the slope is \(f\). At a point where \(f'(x) = 0\) and \(f''(x) > 0\), the function has a \(g\). At a point where \(h\), the function has a maximum. A point where \(i\) is an \(j\) point, provided \(k\) changes sign. The tangent line \(l\) the graph.

The centered approximation to \(f'(x)\) is \([m]2Ax\). The 3-point approximation to \(f'(x)\) is \([n]2\Delta x^2\). The second-order approximation to \(f(x + Ax)\) is \(f(x) + f'(x)Ax + [\text{m}]\). Without that extra term this is just the \(\text{n}\) approximation. With that term the error is \(O(\text{o})\).

1. A graph that is concave upward is inaccurately said to "hold water." Sketch a graph with \(f''(x) > 0\) that would not hold water.
2. Find a function that is concave down for \(x < 0\) and concave up for \(0 < x < 1\) and concave down for \(x > 1\).
3. Can a function be always concave down and never cross zero? Can it be always concave down and positive? Explain.
4. Find a function with \(f''(2) = 0\) and no other inflection point.

**True or false,** when \(f(x)\) is a 9th degree polynomial with \(f'(1) = 0\) and \(f'(3) = 0\). Give (or draw) a reason.

- \(j\) \(k\) somewhere between \(x = 1\) and \(x = 3\).
- \(f''(x) = 0\) somewhere between \(x = 1\) and \(x = 3\).
3.3 Second Derivatives: Bending and Acceleration

There is no absolute maximum at \( x = 3 \).

There are seven points of inflection.

If \( f(x) \) has nine zeros, it has seven inflection points.

If \( f(x) \) has seven inflection points, it has nine zeros.

In 11–16 decide which stationary points are maxima or minima.

\[ f(x) = x^2 - 6x \]

\[ f(x) = x^3 - 6x^2 \]

\[ f(x) = (x - 2)^2(x - 4)^2 \]

\[ f(x) = x + x^2 \]

Locate the inflection points and the regions where \( f(x) \) is concave up or down.

\[ f(x) = x + x^4 - x^3 \]

\[ f(x) = \sin x + \tan x \]

\[ f(x) = x + x^3 \]

\[ f(x) = \sin x \]

\[ f(x) = \sin x - \cos x \]

\[ f(x) = x + \sin 2x \]

Write down the quadratic \( f(0) + f'(0)x + \frac{1}{2}f''(0)x^2 \) in 23–26.

\[ f(x) = \cos x + \sin x \]

\[ f(x) = \tan x \]

\[ f(x) = (\sin x)/x \]

\[ f(x) = 1 + x + x^2 \]

In 26, find \( f(1) + f'(1)(x - 1) + \frac{1}{2}f''(1)(x - 1)^2 \) around \( a = 1 \).

Find \( A \) and \( B \) in \( \sqrt{1 - x} \approx 1 + Ax + Bx^2 \).

Find \( A \) and \( B \) in \( 1/(1 - x)^2 \approx 1 + Ax + Bx^2 \).

Substitute the quadratic approximation into \( [f(x + \Delta x) - f(x)]/\Delta x \), to estimate the error in this one-sided approximation to \( f'(x) \).

What is the quadratic approximation at \( x = 0 \) to \( f(-\Delta x) \)?

Substitute for \( f(x + \Delta x) \) and \( f(x - \Delta x) \) in the centered approximation \([f(x + \Delta x) - f(x - \Delta x)]/2\Delta x \), to get \( f'(x) + \text{error} \). Find the \( \Delta x \) and \( (\Delta x)^2 \) terms in this error. Test on \( f(x) = x^3 \) at \( x = 0 \).

Guess a third-order approximation \( f(\Delta x) \approx f(0) + f'(0)\Delta x + \frac{1}{2}f''(0)(\Delta x)^2 + \ldots \). Test it on \( f(x) = x^3 \).

Construct a table as in the text, showing the actual errors at \( x = 0 \) in one-sided differences, centered differences, second differences, and quadratic approximations. By hand take two values of \( \Delta x \), by calculator take three, by computer take four.

\[ f(x) = x^2 + x^4 \]

\[ f(x) = 1/(1 - x) \]

\[ f(x) = x^2 + \sin x \]

Example 5 was \( 1/(1 - x) \approx 1 + x + x^2 \). What is the error at \( x = 0.1 \)? What is the error at \( x = 2? \)

Substitute \( x = 0.1 \) and \( x = -0.1 \) in the geometric series \( 1/(1 - x) = 1 + x + x^2 + \cdots \) to find \( 1/1.01 \) and \( 1/1.01 - 1 \) first to four decimals and then to all decimals.

Compute \( \cos x^2 \) by equation (4) with \( a = 0 \). OK to check on a calculator. Also compute \( \cos 1 \). Why so far off?

Why is \( \sin x \approx x \) not only a linear approximation but also a quadratic approximation? \( x = 0 \) is an \text{_____} point.

If \( f(x) \) is an even function, find its quadratic approximation at \( x = 0 \). What is the equation of the tangent line?

For \( f(x) \) is an even function, what is the center difference \([f(3) - f(1)]/2 \), and what is the true slope \( f''(2) \)?

For \( f(x) = x + x^2 + x^3 \), what is the second difference \([f(3) - 2f(2) + f(1)]/2 \), and what is the exact \( f''(2) \)?

The error in \( f(a) + f'(a)(x - a) \) is approximately \( \frac{1}{2}f''(a)(x - a)^2 \). This error is positive when the function is \text{_____}. Then the tangent line is \text{_____} the curve.

Draw a piecewise linear \( y(x) \) that is concave up. Define "concave up" without using the test \( d^2y/dx^2 \). If derivatives don't exist, a new definition is needed.

What do these sentences say about \( f \) or \( f' \) or \( f'' \) or \( f''' \)?

1. The population is growing more slowly.
2. The plane is landing smoothly.
3. The economy is picking up speed.
4. The tax rate is constant.
5. A bike accelerates faster but a car goes faster.
6. Stock prices have peaked.
7. The rate of acceleration is slowing down.
8. This course is going downhill.

(Recommended) Draw a curve that goes up-down-up. Below it draw its derivative. Then draw its second derivative. Mark the same points on all curves—the maximum, minimum, and inflection points of the first curve.

Repeat Problem 46 on a printout showing \( y(x) = x^3 - 4x^2 + x + 2 \) and \( dy/dx \) and \( d^2y/dx^2 \) on the same graph.
Reading a graph is like appreciating a painting. Everything is there, but you have to know what to look for. One way to learn is by sketching graphs yourself, and in the past that was almost the only way. Now it is obsolete to spend weeks drawing curves—a computer or graphing calculator does it faster and better. That doesn't remove the need to appreciate a graph (or a painting), since a curve displays a tremendous amount of information.

This section combines two approaches. One is to study actual machine-produced graphs (especially electrocardiograms). The other is to understand the mathematics of graphs—slope, concavity, asymptotes, shifts, and scaling. We introduce the centering transform and zoom transform. These two approaches are like the rest of calculus, where special derivatives and integrals are done by hand and day-to-day applications are by computer. Both are essential—the machine can do experiments that we could never do. But without the mathematics our instructions miss the point. To create good graphs you have to know a few of them personally.

READING AN ELECTROCARDIOGRAM (ECG or EKG)

The graphs of an ECG show the electrical potential during a heartbeat. There are twelve graphs—six from leads attached to the chest, and six from leads to the arms and left leg. (It doesn't hurt, but everybody is nervous. You have to lie still, because contraction of other muscles will mask the reading from the heart.) The graphs record electrical impulses, as the cells depolarize and the heart contracts.

What can I explain in two pages? The graph shows the fundamental pattern of the ECG. Note the P wave, the QRS complex, and the T wave. Those patterns, seen differently in the twelve graphs, tell whether the heart is normal or out of rhythm—or suffering an infarction (a heart attack).

First of all the graphs show the heart rate. The dark vertical lines are by convention 1/2 second apart. The light lines are 1/3 second apart. If the heart beats every 1/2 second (one dark line) the rate is 5 beats per second or 300 per minute. That is extreme tachycardia—not compatible with life. The normal rate is between three dark lines per beat (1/3 second, or 100 beats per minute) and five dark lines (one second between beats, 60 per minute). A baby has a faster rate, over 100 per minute. In this figure the rate is _____. A rate below 60 is bradycardia, not in itself dangerous. For a resting athlete that is normal.

Doctors memorize the six rates 300, 150, 100, 75, 60, 50. Those correspond to 1, 2, 3, 4, 5, 6 dark lines between heartbeats. The distance is easiest to measure between spikes (the peaks of the R wave). Many doctors put a printed scale next to the chart. One textbook emphasizes that “Where the next wave falls determines the rate. No mathematical computation is necessary.” But you see where those numbers come from.
3.4 Graphs

The next thing to look for is heart rhythm. The regular rhythm is set by the pacemaker, which produces the P wave. A constant distance between waves is good—and then each beat is examined. When there is a block in the pathway, it shows as a delay in the graph. Sometimes the pacemaker fires irregularly. Figure 3.10 shows sinus arrhythmia (fairly normal). The time between peaks is changing. In disease or emergency, there are potential pacemakers in all parts of the heart.

I should have pointed out the main parts. We have four chambers, an atrium-ventricle pair on the left and right. The SA node should be the pacemaker. The stimulus spreads from the atria to the ventricles—from the small chambers that "prime the pump" to the powerful chambers that drive blood through the body. The P wave comes with contraction of the atria. There is a pause of \( \frac{1}{2} \) second at the AV node. Then the big QRS wave starts contraction of the ventricles, and the T wave is when the ventricles relax. The cells switch back to negative charge and the heart cycle is complete.

![Diagram of heart with atria and ventricles](image)

**Fig. 3.9** Happy person with a heart and a normal electrocardiogram.

The ECG shows when the pacemaker goes wrong. Other pacemakers take over—the AV node will pace at 60/minute. An early firing in the ventricle can give a wide spike in the QRS complex, followed by a long pause. The impulses travel by a slow path. Also the pacemaker can suddenly speed up (paroxysmal tachycardia is 150-250/minute). But the most critical danger is fibrillation.

Figure 3.10b shows a dying heart. The ECG indicates irregular contractions—no normal PQRST sequence at all. What kind of heart would generate such a rhythm? The muscles are quivering or "fibrillating" independently. The pumping action is nearly gone, which means emergency care. The patient needs immediate CPR—someone to do the pumping that the heart can't do. Cardio-pulmonary resuscitation is a combination of chest pressure and air pressure (hand and mouth) to restart the rhythm. CPR can be done on the street. A hospital applies a defibrillator, which shocks the heart back to life. It depolarizes all the heart cells, so the timing can be reset. Then the charge spreads normally from SA node to atria to AV node to ventricles.

This discussion has not used all twelve graphs to locate the problem. That needs vectors. Look ahead at Section 11.1 for the heart vector, and especially at Section 11.2 for its twelve projections. Those readings distinguish between atrium and ventricle, left and right, forward and back. This information is of vital importance in the event of a heart attack. A "heart attack" is a myocardial infarction (MI).

An MI occurs when part of an artery to the heart is blocked (a coronary occlusion).
Applications of the Derivative

Fig. 3.10 Doubtful rhythm. Serious fibrillation. Signals of a heart attack.

An area is without blood supply—therefore without oxygen or glucose. Often the attack is in the thick left ventricle, which needs the most blood. The cells are first ischemic, then injured, and finally infarcted (dead). The classical ECG signals involve those three I's:

**Ischemia:** Reduced blood supply, upside-down T wave in the chest leads.

**Injury:** An elevated segment between S and T means a recent attack.

**Infarction:** The Q wave, normally a tiny dip or absent, is as wide as a small square (\( \frac{1}{2} \) second). It may occupy a third of the entire QRS complex.

The Q wave gives the diagnosis. You can find all three I's in Figure 3.10c.

It is absolutely amazing how much a good graph can do.

**THE MECHANICS OF GRAPHS**

From the meaning of graphs we descend to the mechanics. A formula is now given for \( f(x) \). The problem is to create the graph. It would be too old-fashioned to evaluate \( f(x) \) by hand and draw a curve through a dozen points. A computer has a much better idea of a parabola than an artist (who tends to make it asymptotic to a straight line). There are some things a computer knows, and other things an artist knows, and still others that you and I know—because we understand derivatives.

Our job is to apply calculus. We extract information from \( f' \) and \( f'' \) as well as \( f \). Small movements in the graph may go unnoticed, but the important properties come through. Here are the main tests:

1. The sign of \( f(x) \) (above or below axis: \( f = 0 \) at crossing point)
2. The sign of \( f'(x) \) (increasing or decreasing: \( f' = 0 \) at stationary point)
3. The sign of \( f''(x) \) (concave up or down: \( f'' = 0 \) at inflection point)
4. The behavior of \( f(x) \) as \( x \to \infty \) and \( x \to -\infty \)
5. The points at which \( f(x) \to \infty \) or \( f(x) \to -\infty \)
6. Even or odd? Periodic? Jumps in \( f \) or \( f' \)? Endpoints? \( f(0) \)?

**EXAMPLE 1**

\[ f(x) = \frac{x^2}{1 - x^2} \quad f'(x) = \frac{2x}{(1 - x^2)^2} \quad f''(x) = \frac{2 + 6x^2}{(1 - x^2)^3} \]

The sign of \( f(x) \) depends on \( 1 - x^2 \). Thus \( f(x) > 0 \) in the inner interval where \( x^2 < 1 \). The graph bends upwards \((f''(x) > 0)\) in that same interval. There are no inflection points, since \( f'' \) is never zero. The stationary point where \( f' \) vanishes is \( x = 0 \). We have a local minimum at \( x = 0 \).

The guidelines (or asymptotes) meet the graph at infinity. For large \( x \) the important terms are \( x^2 \) and \( -x^2 \). Their ratio is \( +x^2/-x^2 = -1 \)—which is the limit as \( x \to \infty \) and \( x \to -\infty \). The horizontal asymptote is the line \( y = -1 \).

The other infinities, where \( f \) blows up, occur when \( 1 - x^2 \) is zero. That happens at \( x = 1 \) and \( x = -1 \). The vertical asymptotes are the lines \( x = 1 \) and \( x = -1 \). The graph
in Figure 3.11a approaches those lines.

if \( f(x) \to b \) as \( x \to +\infty \) or \( -\infty \), the line \( y = b \) is a horizontal asymptote

if \( f(x) \to +\infty \) or \( -\infty \) as \( x \to a \), the line \( x = a \) is a vertical asymptote

if \( f(x) - (mx + b) \to 0 \) as \( x \to +\infty \) or \( -\infty \), the line \( y = mx + b \) is a sloping asymptote.

Finally comes the vital fact that this function is even: \( f(x) = f(-x) \) because squaring \( x \) obliterates the sign. The graph is symmetric across the \( y \) axis.

To summarize the effect of dividing by \( 1 - x^2 \): No effect near \( x = 0 \). Blowup at 1 and \( -1 \) from zero in the denominator. The function approaches \(-1\) as \( |x| \to \infty \).

**EXAMPLE 2**

\[
f(x) = \frac{x^2}{x - 1} \quad f'(x) = \frac{x^2 - 2x}{(x - 1)^2} \quad f''(x) = \frac{2}{(x - 1)^3}
\]

This example divides by \( x - 1 \). Therefore \( x = 1 \) is a vertical asymptote, where \( f(x) \) becomes infinite. Vertical asymptotes come mostly from zero denominators.

Look beyond \( x = 1 \). Both \( f(x) \) and \( f'(x) \) are positive for \( x > 1 \). The slope is zero at \( x = 2 \). That must be a local minimum.

What happens as \( x \to \infty \)? Dividing \( x^2 \) by \( x - 1 \), the leading term is \( x \). The function becomes large. It grows linearly—we expect a sloping asymptote. To find it, do the division properly:

\[
\frac{x^2}{x - 1} = x + 1 + \frac{1}{x - 1}.
\]

The last term goes to zero. The function approaches \( y = x + 1 \) as the asymptote.

This function is not odd or even. Its graph is in Figure 3.11b. With **zoom out** you see the asymptotes. **Zoom in** for \( f = 0 \) or \( f' = 0 \) or \( f'' = 0 \).

**EXAMPLE 3**

\( f(x) = \sin x + \frac{1}{3} \sin 3x \) has the slope \( f'(x) = \cos x + \cos 3x \).

Above all these functions are periodic. If \( x \) increases by \( 2\pi \), nothing changes. The graphs from \( 2\pi \) to \( 4\pi \) are repetitions of the graphs from \( 0 \) to \( 2\pi \). Thus \( f(x + 2\pi) = f(x) \) and the period is \( 2\pi \). Any interval of length \( 2\pi \) will show a complete picture, and Figure 3.11c picks the interval from \(-\pi\) to \( \pi \).

The second outstanding property is that \( f \) is odd. The sine functions satisfy \( f(-x) = -f(x) \). The graph is symmetric through the origin. By reflecting the right half through the origin, you get the left half. In contrast, the cosines in \( f'(x) \) are even.

To find the zeros of \( f(x) \) and \( f'(x) \) and \( f''(x) \), rewrite those functions as

\[
f(x) = 2 \sin x - \frac{1}{3} \sin 3x \quad f'(x) = -2 \cos x + 4 \cos^3 x \quad f''(x) = -10 \sin x + 12 \sin^3 x.
\]
We changed $\sin 3x$ to $3 \sin x - 4 \sin^3 x$. For the derivatives use $\sin^2 x = 1 - \cos^2 x$.

Now find the zeros—the crossing points, stationary points, and inflection points:

\[ f = 0 \quad 2 \sin x = \frac{3}{2} \sin^3 x \Rightarrow \sin x = 0 \text{ or } \sin^2 x = \frac{3}{4} \Rightarrow x = 0, \pm \pi \]
\[ f' = 0 \quad 2 \cos x = -4 \cos^3 x \Rightarrow \cos x = 0 \text{ or } \cos^2 x = \frac{1}{4} \Rightarrow x = \pm \pi/4, \pm \pi/2, \pm 3\pi/4 \]
\[ f'' = 0 \quad 5 \sin x = 6 \sin^3 x \Rightarrow \sin x = 0 \text{ or } \sin^2 x = \frac{5}{6} \Rightarrow x = 0, \pm 66^\circ, \pm 114^\circ, \pm \pi \]

That is more than enough information to sketch the graph. The stationary points $\pi/4, \pi/2, 3\pi/4$ are evenly spaced. At those points $f(x)$ is $\sqrt{8}/3$ (maximum), $2/3$ (local minimum), $\sqrt{8}/3$ (maximum). Figure 3.11c shows the graph.

I would like to mention a beautiful continuation of this same pattern:

\[ f(x) = \sin x + \frac{1}{3} \sin 3x + \frac{1}{5} \sin 5x + \cdots \quad f'(x) = \cos x + \cos 3x + \cos 5x + \cdots \]

If we stop after ten terms, $f(x)$ is extremely close to a step function. If we don't stop, the exact step function contains infinitely many sines. It jumps from $-\pi/4$ to $+\pi/4$ as $x$ goes past zero. More precisely it is a "square wave," because the graph jumps back down at $\pi$ and repeats. The slope $\cos x + \cos 3x + \cdots$ also has period $2\pi$. Infinitely many cosines add up to a delta function! (The slope at the jump is an infinite spike.) These sums of sines and cosines are Fourier series.

**GRAPHS BY COMPUTERS AND CALCULATORS**

We have come to a topic of prime importance. If you have graphing software for a computer, or if you have a graphing calculator, you can bring calculus to life. A graph presents $y(x)$ in a new way—different from the formula. Information that is buried in the formula is clear on the graph. But don't throw away $y(x)$ and $dy/dx$. The derivative is far from obsolete.

These pages discuss how calculus and graphs go together. We work on a crucial problem of applied mathematics—to find where $y(x)$ reaches its minimum. There is no need to tell you a hundred applications. Begin with the formula. How do you find the point $x^*$ where $y(x)$ is smallest?

First, draw the graph. That shows the main features. We should see (roughly) where $x^*$ lies. There may be several minima, or possibly none. But what we see depends on a decision that is ours to make—the range of $x$ and $y$ in the viewing window.

If nothing is known about $y(x)$, the range is hard to choose. We can accept a default range, and zoom in or out. We can use the autoscaling program in Section 1.7. Somehow $x^*$ can be observed on the screen. Then the problem is to compute it.

I would like to work with a specific example. We solved it by calculus—to find the best point $x^*$ to enter an expressway. The speeds in Section 3.2 were 30 and 60. The length of the fast road will be $b = 6$. The range of reasonable values for the entering point is $0 \leq x \leq 6$. The distance to the road in Figure 3.12 is $a = 3$. We drive a distance $\sqrt{3^2 + x^2}$ at speed 30 and the remaining distance $6 - x$ at speed 60:

\[ \text{driving time } \quad y(x) = \frac{1}{30} \sqrt{3^2 + x^2} + \frac{1}{60}(6 - x). \quad (2) \]

This is the function to be minimized. Its graph is extremely flat.

It may seem unusual for the graph to be so level. On the contrary, it is common. A flat graph is the whole point of $dy/dx = 0$.

The graph near the minimum looks like $y = Cx^2$. It is a parabola sitting on a horizontal tangent. At a distance of $\Delta x = .01$, we only go up by $C(\Delta x)^2 = .0001 C$. Unless $C$ is a large number, this $\Delta y$ can hardly be seen.
The solution is to change scale. **Zoom in on** \( x^* \). The tangent line stays flat, since \( dy/dx \) is still zero. But the bending from \( C \) is increased. Figure 3.12 shows the zoom box blown up into a new graph of \( y(x) \).

A calculator has one or more ways to find \( x^* \). With a TRACE mode, you direct a cursor along the graph. From the display of \( y \) values, read \( y_{\text{max}} \) and \( x^* \) to the nearest pixel. A zoom gives better accuracy, because it stretches the axes—each pixel represents a smaller \( \Delta x \) and \( \Delta y \). The TI-81 stretches by 4 as default. Even better, let the whole process be graphical—**draw the actual ZOOM BOX on the screen**. Pick two opposite corners, press ENTER, and the box becomes the new viewing window (Figure 3.12).

The first zoom narrows the search for \( x^* \). It lies between \( x = 1 \) and \( x = 3 \). We build a new ZOOM BOX and zoom in again. Now \( 1.5 \leq x^* \leq 2 \). Reasonable accuracy comes quickly. High accuracy does not come quickly. It takes time to create the box and execute the zoom.

**Question 1** What happens as we zoom in, if all boxes are square (equal scaling)?

**Answer** The picture gets flatter and flatter. We are zooming in to the tangent line. Changing \( x \) to \( X/4 \) and \( y \) to \( Y/4 \), the parabola \( y = x^2 \) flattens to \( Y = X^2/4 \). To see any bending, we must use a long thin zoom box.

I want to change to a totally different approach. Suppose we have a formula for \( dy/dx \). *That derivative was produced by an infinite zoom!* The limit of \( \Delta y/\Delta x \) came by brainpower alone:

\[
\frac{dy}{dx} = \frac{x}{30\sqrt{3^2 + x^2}} - \frac{1}{60}
\]

Call this \( f(x) \).

This function is zero at \( x^* \). The computing problem is completely changed: Solve \( f(x) = 0 \). It is easier to find a root of \( f(x) \) than a minimum of \( y(x) \). The graph of \( f(x) \) crosses the \( x \) axis. The graph of \( y(x) \) goes flat—this is harder to pinpoint.

Take the model function \( y = x^2 \) for \( |x| < .01 \). The slope \( f = 2x \) changes from \(-.02 \) to \(+.02 \). The value of \( x^2 \) moves only by \(.0001 \)—its minimum point is hard to see.

To repeat: Minimization is easier with \( dy/dx \). The screen shows an order of magnitude improvement, when we trace or zoom on \( f(x) = 0 \). In calculus, we have been taking the derivative for granted. It is natural to get blasé about \( dy/dx = 0 \). We forget how intelligent it is, to work with the slope instead of the function.

**Question 2** How do you get another order of magnitude improvement?

**Answer** Use the next derivative! With a formula for \( df/dx \), which is \( d^2y/dx^2 \), the convergence is even faster. In two steps the error goes from \(.01 \) to \(.0001 \) to \(.00000001 \). Another infinite zoom went into the formula for \( df/dx \), and **Newton's method** takes account of it. Sections 3.6 and 3.7 study \( f(x) = 0 \).
The expressway example allows perfect accuracy. We can solve \( \frac{dy}{dx} = 0 \) by algebra. The equation simplifies to \( 60x = 30\sqrt{3^2 + x^2} \). Dividing by 30 and squaring yields \( 4x^2 = 3^2 + x^2 \). Then \( 3x^2 = 3^2 \). The exact solution is \( x^* = \sqrt{3} = 1.73205... \).

A model like this is a benchmark, to test competing methods. It also displays what we never appreciated—the extreme flatness of the graph. The difference in driving time between entering at \( x^* = \sqrt{3} \) and \( x = 2 \) is one second.

THE CENTERING TRANSFORM AND ZOOM TRANSFORM

For a photograph we do two things—point the right way and stand at the right distance. Then take the picture. Those steps are the same for a graph. First we pick the new center point. The graph is shifted, to move that point from \((a, b)\) to \((0, 0)\). Then we decide how far the graph should reach. It fits in a rectangle, just like the photograph. Rescaling to \( x/c \) and \( y/d \) puts the desired section of the curve into the rectangle.

A good photographer does more (like an artist). The subjects are placed and the camera is focused. For good graphs those are necessary too. But an everyday calculator or computer or camera is built to operate without an artist—just aim and shoot. I want to explain how to aim at \( y = f(x) \).

We are doing exactly what a calculator does, with one big difference. It doesn't change coordinates. We do. When \( x = 1, y = -2 \) moves to the center of the viewing window, the calculator still shows that point as \((1, -2)\). When the centering transform acts on \( y + 2 = m(x - 1) \), those numbers disappear. This will be confusing unless \( x \) and \( y \) also change. The new coordinates are \( X = x - 1 \) and \( Y = y + 2 \). Then the new equation is \( Y = mX \).

The main point (for humans) is to make the algebra simpler. The computer has no preference for \( Y = mX \) over \( y - y_0 = m(x - x_0) \). It accepts \( 2x^2 - 4x \) as easily as \( x^2 \). But we do prefer \( Y = mX \) and \( y = x^2 \), partly because their graphs go through \((0, 0)\). Ever since zero was invented, mathematicians have liked that number best.

**EXAMPLE 4** The parabola \( y = 2x^2 - 4x \) has its minimum when \( dy/dx = 4x - 4 = 0 \). Thus \( x = 1 \) and \( y = -2 \). Move this bottom point to the center: \( y = 2x^2 - 4x \) is

\[
Y + 2 = 2(X - 1)^2 - 4(X - 1) \quad \text{or} \quad Y = 2X^2.
\]

The new parabola \( Y = 2X^2 \) has its bottom at \((0, 0)\). It is the same curve, shifted across and up. The only simpler parabola is \( y = x^2 \). This final step is the job of the zoom.

Next comes scaling. We may want more detail (zoom in to see the tangent line). We may want a big picture (zoom out to check asymptotes). We might stretch one axis more than the other, if the picture looks like a pancake or a skyscraper.

**3G** A zoom transform scales the \( X \) and \( Y \) axes by \( c \) and \( d \):

\[
x = cX \quad \text{and} \quad y = dY \quad \text{change} \quad Y = F(X) \quad \text{to} \quad y = dF(x/c).
\]

The new \( x \) and \( y \) are boldface letters, and the graph is rescaled. Often \( c = d \).
EXAMPLE 5 Start with $Y = 2X^2$. Apply a square zoom with $c = d$. In the new xy coordinates, the equation is $y/c = 2(x/c)^2$. The number 2 disappears if $c = d = 2$. With the right centering and the right zoom, every parabola that opens upward is $y = x^2$.

Question 3 What happens to the derivatives (slope and bending) after a zoom?
Answer The slope (first derivative) is multiplied by $d/c$. Apply the chain rule to $y = dF(x/c)$. A square zoom has $d/c = \text{lines keep their slope}$. The second derivative is multiplied by $d/c^2$, which changes the bending. A zoom out divides by small numbers $c = d$, so the big picture is more curved.

Combining the centering and zoom transforms, as we do in practice, gives $y$ in terms of $x$:

$$y = f(x) \quad \text{becomes} \quad Y = f(X + a) - b \quad \text{and then} \quad y = d \left( f \left( \frac{x}{c} + a \right) - b \right).$$

![Fig. 3.14 Change of coordinates by centering and zoom. Calculators still show $(x, y)$.](image)

Question 4 Find $x$ and $y$ ranges after two transforms. Start between $-1$ and 1.
Answer The window after centering is $-1 < x - a < 1$ and $-1 < y - b < 1$. The window after zoom is $-1 < c(x - a) < 1$ and $-1 < d(y - b) < 1$. The point $(1, 1)$ was originally in the corner. The point $(c^{-1} + a, d^{-1} + b)$ is now in the corner.

The numbers $a, b, c, d$ are chosen to produce a simpler function (like $y = x^2$). Or else—this is important in applied mathematics—they are chosen to make $x$ and $y$ “dimensionless.” An example is $y = \frac{1}{2} \cos \theta$. The frequency $\theta$ has dimension 1/time. The amplitude $\frac{1}{2}$ is a distance. With $d = 2\text{ cm}$ and $c = 8\text{ sec}$, the units are removed and $y = \cos t$.

May I mention one transform that does change the slope? It is a rotation. The whole plane is turned. A photographer might use it—but normally people are supposed to be upright. You use rotation when you turn a map or straighten a picture.

In the next section, an unrecognizable hyperbola is turned into $Y = 1/X$.

### 3.4 Exercises

**Read-through questions**

The position, slope, and bending of $y = f(x)$ are decided by $a, b, c, d$. If $|f(x)| \to \infty$ as $x \to a$, the line $x = a$ is a vertical asymptote. If $f(x) \to b$ for large $x$, then $y = b$ is a horizontal asymptote. The asymptotes of $y = x^2/(x^2 - 4)$ are $y = 0$ and $y = \pm 1$. This function is even because $f(-x) = f(x)$. The function $\sin kx$ has period $2\pi/k$.

Near a point where $dy/dx = 0$, the graph is extremely flat. For the model $y = Cx^2$, $x = 1$ gives $y = k$. A box around the graph looks long and slanted. We move in to that box for another digit of $x^n$. But solving $dy/dx = 0$ is more accurate, because its graph touches the x axis. The slope of $dy/dx$ is $0$. Each derivative is like an $n$-th power zoom.

To move $(a, b)$ to $(0, 0)$, shift the variables to $X = x - a$ and $Y = y - b$. This transform changes $y = f(x)$ to $Y = f(X)$. The original slope at $(a, b)$ equals the new slope at $(0, 0)$. To stretch the axes by $c$ and $d$, set $x = cX$ and $y = dY$. The transform changes $Y = f(X)$ to $y = dy/dx$. Slopes are multiplied by $c$. Second derivatives are multiplied by $d$. 


1. Find the pulse rate when heartbeats are 1/2 second or two dark lines or x seconds apart.

2. Another way to compute the heart rate uses marks for 6-second intervals. Doctors count the cycles in an interval.
   (a) How many dark lines in 6 seconds?
   (b) With 8 beats per interval, find the rate.
   (c) Rule: Heart rate = cycles per interval times ________.

Which functions in 3–18 are even or odd or periodic? Find all asymptotes: \( y = b \) or \( x = a \) or \( y = mx + b \). Draw roughly by hand or smoothly by computer.

3. \( f(x) = x - (9/x) \)
4. \( f(x) = x^n \) (any integer \( n \))
5. \( f(x) = \frac{1}{1-x^2} \)
6. \( f(x) = \frac{x^3}{4-x^2} \)
7. \( f(x) = \frac{x^2+3}{x^2+1} \)
8. \( f(x) = \frac{x^2+3}{x+1} \)
9. \( f(x) = (\sin x)(\sin 2x) \)
10. \( f(x) = \cos x + \cos 3x + \cos 5x \)
11. \( f(x) = \frac{x \sin x}{x^2-1} \)
12. \( f(x) = \frac{x}{\sin x} \)
13. \( f(x) = \frac{1}{x^3+x^2} \)
14. \( f(x) = \frac{1}{x-1} - 2x \)
15. \( f(x) = \frac{x^3+1}{x^3-1} \)
16. \( f(x) = \frac{\sin x + \cos x}{\sin x - \cos x} \)
17. \( f(x) = x - \sin x \)
18. \( f(x) = (1/x) - \sqrt{x} \)

In 19–24 construct \( f(x) \) with exactly these asymptotes.

19. \( x = 1 \) and \( y = 2 \)
20. \( x = 1, x = 2, y = 0 \)
21. \( y = x \) and \( x = 4 \)
22. \( y = 2x + 3 \) and \( x = 0 \)
23. \( y = x (x \to \infty), y = -x (x \to -\infty) \)
24. \( x = 1, x = 3, y = x \)
25. For \( P(x)/Q(x) \) to have \( y = 2 \) as asymptote, the polynomials \( P \) and \( Q \) must be ________.
26. For \( P(x)/Q(x) \) to have a sloping asymptote, the degrees of \( P \) and \( Q \) must be ________.
27. For \( P(x)/Q(x) \) to have the asymptote \( y = 0 \), the degrees of \( P \) and \( Q \) must _________. The graph of \( x^4/(1 + x^2) \) has what asymptotes?
28. Both \( 1/(x-1) \) and \( 1/(x-1)^2 \) have \( x = 1 \) and \( y = 0 \) as asymptotes. The most obvious difference in the graphs is ________.
29. If \( f'(x) \) has asymptotes \( x = 1 \) and \( y = 3 \) then \( f(x) \) has asymptotes _________.

30. True (with reason) or false (with example).
   (a) Every ratio of polynomials has asymptotes
   (b) If \( f(x) \) is even so is \( f'(x) \)
   (c) If \( f'(x) \) is even so is \( f(x) \)
   (d) Between vertical asymptotes, \( f'(x) \) touches zero.

31. Construct an \( f(x) \) that is "even around \( x = 3.\)"
32. Construct \( g(x) \) to be "odd around \( x = \pi.\)"

Create graphs of 33–38 on a computer or calculator.

33. \( y(x) = (1 + 1/x)^2, -3 \leq x \leq 3 \)
34. \( y(x) = x^{1/x}, 0.1 \leq x \leq 2 \)
35. \( y(x) = \sin(x/3) + \sin(x/5) \)
36. \( y(x) = (2 - x)/(2 + x), -3 \leq x \leq 3 \)
37. \( y(x) = 2x^3 + 3x^2 - 12x + 5 \) on \([-3, 3]\] and \([2.9, 3.1]\]
38. \( 100(\sin(x+.1) - 2 \sin x + \sin(x-.1)) \)

In 39–40 show the asymptotes on large-scale computer graphs.

39. (a) \( y = \frac{x^2 + 8x - 15}{x^2 - 2} \) \( (b) y = \frac{x^4 - 6x^3 + 1}{2x^4 + x^2} \)
40. (a) \( y = \frac{x^2 - 2}{x^3 + 8x - 15} \) \( (b) y = \frac{x^4 - x + 2}{x^3 - 2x + 1} \)

41. Rescale \( y = \sin x \) so \( X \) is in degrees, not radians, and \( Y \) changes from meters to centimeters.

Problems 42–46 minimize the driving time \( y(x) \) in the text. Some questions may not fit your software.

42. Trace along the graph of \( y(x) \) to estimate \( x^* \). Choose an \( xy \) range or use the default.
43. Zoom in by \( c = d = 4 \). How many zooms until you reach \( x^* = 1.73205 \) or \( 1.7320508 \)?
44. Ask your program for the minimum of \( y(x) \) and the solution of \( dy/dx = 0 \). Same answer?
45. What are the scaling factors \( c \) and \( d \) for the two zooms in Figure 3.12? They give the stretching of the \( x \) and \( y \) axes.
46. Show that \( dy/dx = -1/60 \) and \( d^2y/dx^2 = 1/90 \) at \( x = 0 \). Linear approximation gives \( dy/dx \approx -1/60 + x/90 \). So the slope is zero near \( x = \) _______. This is Newton's method, using the next derivative.

Change the function to \( y(x) = \sqrt{15 + x^2/30} + (10 - x)/60. \)
47. Find \( x^* \) using only the graph of \( y(x) \).
48. Find \( x^* \) using also the graph of \( dy/dx \).
49. What are the \( xy \) and \( XY \) and \( xy \) equations for the line in Figure 3.14?
50 Define \( f_n(x) = \sin x + \frac{1}{2} \sin 3x + \frac{1}{4} \sin 5x + \cdots \) \( (n \text{ terms}) \). Graph \( f_3 \) and \( f_{10} \) from \(-\pi\) to \(\pi\). Zoom in and describe the Gibbs phenomenon at \(x = 0\).

On the graphs of 51–56, zoom in to all maxima and minima (3 significant digits). Estimate inflection points.

51 \( y = 2x^5 - 16x^4 + 5x^3 - 37x^2 + 21x + 683 \)
52 \( y = x^5 - x^4 - \sqrt{3x^3 + 1} - 2 \)
53 \( y = x(x - 1)(x - 2)(x - 4) \)
54 \( y = 7 \sin 2x + 5 \cos 3x \)
55 \( y = (x^3 - 2x + 1)/(x^4 - 3x^2 - 15), -3 \leq x \leq 5 \)
56 \( y = x \sin (1/x), 0.1 \leq x \leq 1 \)
57 A 10-digit computer shows \( y = 0 \) and \( dy/dx = .01 \) at \( x^* = 1 \). This root should be correct to about (8 digits) (10 digits) (12 digits). \textit{Hint:} Suppose \( y = .01 \) \((x - 1 + \text{error})\). What errors don't show in 10 digits of \(y\)?
58 Which is harder to compute accurately: Maximum point or inflection point? First derivative or second derivative?

3.5 Parabolas, Ellipses, and Hyperbolas

Here is a list of the most important curves in mathematics, so you can tell what is coming. It is not easy to rank the top four:

1. \textit{straight lines}
2. \textit{sines and cosines} (oscillation)
3. \textit{exponentials} (growth and decay)
4. \textit{parabolas, ellipses, and hyperbolas} (using 1, \(x\), \(y\), \(x^2\), \(xy\), \(y^2\)).

The curves that I wrote last, the Greeks would have written first. It is so natural to go from linear equations to quadratic equations. Straight lines use 1, \(x\), \(y\). Second degree curves include \(x^2\), \(xy\), \(y^2\). If we go on to \(x^3\) and \(y^3\), the mathematics gets complicated. We now study equations of second degree, and the curves they produce.

It is quite important to see both the equations and the curves. This section connects two great parts of mathematics—analysis of the equation and geometry of the curve. Together they produce "analytic geometry." You already know about functions and graphs. Even more basic: Numbers correspond to points. We speak about "the point (5, 2)." Euclid might not have understood.

Where Euclid drew a 45° line through the origin, Descartes wrote down \( y = x \). Analytic geometry has become central to mathematics—we now look at one part of it.

Fig. 3.15 The cutting plane gets steeper: circle to ellipse to parabola to hyperbola.
The parabola and ellipse and hyperbola have absolutely remarkable properties. The Greeks discovered that all these curves come from slicing a cone by a plane. The curves are “conic sections.” A level cut gives a circle, and a moderate angle produces an ellipse. A steep cut gives the two pieces of a hyperbola (Figure 3.15d). At the borderline, when the slicing angle matches the cone angle, the plane carves out a parabola. It has one branch like an ellipse, but it opens to infinity like a hyperbola.

Throughout mathematics, parabolas are on the border between ellipses and hyperbolas.

To repeat: We can slice through cones or we can look for equations. For a cone of light, we see an ellipse on the wall. (The wall cuts into the light cone.) For an equation $Ax^2 + Bxy + Cy^2 + Dx + Ey + F = 0$, we will work to make it simpler. The graph will be centered and rescaled (and rotated if necessary), aiming for an equation like $y = x^2$. Eccentricity and polar coordinates are left for Chapter 9.

THE PARABOLA $y = ax^2 + bx + c$

You knew this function long before calculus. The graph crosses the $x$ axis when $y = 0$. The quadratic formula solves $y = 3x^2 - 4x + 1 = 0$, and so does factoring into $(x - 1)(3x - 1)$. The crossing points $x = 1$ and $x = \frac{1}{3}$ come from algebra.

The other important point is found by calculus. It is the minimum point, where $dy/dx = 6x - 4 = 0$. The $x$ coordinate is $\frac{2}{3}$, halfway between the crossing points. The height is $y_{\text{min}} = -\frac{1}{3}$. This is the vertex $V$ in Figure 3.16a—at the bottom of the parabola.

A parabola has no asymptotes. The slope $6x - 4$ doesn’t approach a constant.

To center the vertex Shift left by $\frac{1}{3}$ and up by $\frac{1}{3}$. So introduce the new variables $X = x - \frac{1}{3}$ and $Y = y + \frac{1}{3}$. Then $x = \frac{2}{3}$ and $y = -\frac{1}{3}$ correspond to $X = Y = 0$—which is the new vertex:

$$y = 3x^2 - 4x + 1 \quad \text{becomes} \quad Y = 3X^2.$$  \hspace{1cm} (1)

Check the algebra. $Y = 3X^2$ is the same as $y + \frac{1}{3} = 3(x - \frac{1}{3})^2$. That simplifies to the original equation $y = 3x^2 - 4x + 1$. The second graph shows the centered parabola $Y = 3X^2$, with the vertex moved to the origin.

To zoom in on the vertex Rescale $X$ and $Y$ by the zoom factor $a$:

$$Y = 3X^2 \quad \text{becomes} \quad \frac{Y}{a} = 3\left(\frac{X}{a}\right)^2.$$  

The final equation has $x$ and $y$ in boldface. With $a = 3$ we find $y = x^2$—the graph is magnified by 3. In two steps we have reached the model parabola opening upward.

![Figure 3.16](image-url)
A parabola has another important point—the **focus**. Its distance from the vertex is called $p$. The special parabola $y = x^2$ has $p = 1/4$, and other parabolas $Y = aX^2$ have $p = 1/4a$. You magnify by a factor $a$ to get $y = x^2$. The beautiful property of a parabola is that *every ray coming straight down is reflected to the focus*.

Problem 2.3.25 located the focus $F$—here we mention two applications. A solar collector and a TV dish are parabolic. They concentrate sun rays and TV signals onto a point—a heat cell or a receiver collects them at the focus. The 1982 *UMAP Journal* explains how radar and sonar use the same idea. Car headlights turn the idea around, and send the light outward.

Here is a classical fact about parabolas. From each point on the curve, the distance to the focus equals the distance to the “**directrix**.” The directrix is the line $y = -p$ below the vertex (so the vertex is halfway between focus and directrix). With $p = 1/4$, the distance down from any $(x, y)$ is $y + 1/4$. Match that with the distance to the focus at $(0, 1/4)$—this is the square root below. Out comes the special parabola $y = x^2$:

$$y + 1/4 = \sqrt{x^2 + (y - 1/4)^2} \quad \text{square both sides} \quad y = x^2. \quad (2)$$

The exercises give practice with all the steps we have taken—center the parabola to $Y = aX^2$, rescale it to $y = x^2$, locate the vertex and focus and directrix.

**Summary for other parabolas** $Y = ax^2 + bx + c$ has its vertex where $dy/dx$ is zero. Thus $2ax + b = 0$ and $x = -b/2a$. Shifting across to that point is “completing the square”:

$$ax^2 + bx + c = a \left( x + \frac{b}{2a} \right)^2 + C. \quad (3)$$

Here $C = c - (b^2/4a)$ is the height of the vertex. The centering transform $X = x + (b/2a)$, $Y = y - C$ produces $Y = aX^2$. It moves the vertex to $(0, 0)$, where it belongs.

For the ellipse and hyperbola, our plan of attack is the same:

1. Center the curve to remove any linear terms $Dx$ and $Ey$.
2. Locate each focus and discover the reflection property.
3. Rotate to remove $Bxy$ if the equation contains it.

**Ellipses** $\frac{x^2}{a^2} + \frac{y^2}{b^2} = 1$ (*circles have $a = b$*)

This equation makes the ellipse symmetric about $(0, 0)$—the center. Changing $x$ to $-x$ or $y$ to $-y$ leaves the same equation. No extra centering or rotation is needed.

The equation also shows that $x^2/a^2$ and $y^2/b^2$ cannot exceed one. (They add to one and can’t be negative.) Therefore $x^2 \leq a^2$, and $x$ stays between $-a$ and $a$. Similarly $y$ stays between $b$ and $-b$. The ellipse is inside a rectangle.

By solving for $y$ we get a function (or two functions!) of $x$:

$$\frac{y^2}{b^2} = 1 - \frac{x^2}{a^2} \quad \text{gives} \quad \frac{y}{b} = \pm \sqrt{1 - \frac{x^2}{a^2}} \quad \text{or} \quad y = \pm \frac{b}{a} \sqrt{a^2 - x^2}.$$  

The graphs are the top half (+) and bottom half (−) of the ellipse. To draw the ellipse, plot them together. They meet when $y = 0$, at $x = a$ on the far right of Figure 3.17 and at $x = -a$ on the far left. The maximum $y = b$ and minimum $y = -b$ are at the top and bottom of the ellipse, where we bump into the enclosing rectangle.

A circle is a special case of an ellipse, when $a = b$. The circle equation $x^2 + y^2 = r^2$ is the ellipse equation with $a = b = r$. This circle is centered at $(0, 0)$; other circles are...
Applications of the Derivative

centered at \( x = h, y = k \). The circle is determined by its radius \( r \) and its center \((h, k)\):

**Equation of circle:** \( (x - h)^2 + (y - k)^2 = r^2 \). \( \text{(4)} \)

In words, the distance from \((x, y)\) on the circle to \((h, k)\) at the center is \( r \). The equation has linear terms \(-2hx\) and \(-2ky\)—they disappear when the center is \((0, 0)\).

**EXAMPLE 1** Find the circle that has a diameter from \((1, 7)\) to \((5, 7)\).

**Solution** The center is halfway at \((3, 7)\). So \( r = 2 \) and \((x - 3)^2 + (y - 7)^2 = 2^2\).

**EXAMPLE 2** Find the center and radius of the circle \( x^2 - 6x + y^2 - 14y = -54 \).

**Solution** Complete \( x^2 - 6x \) to the square \((x - 3)^2\) by adding 9. Complete \( y^2 - 14y \) to \((y - 7)^2\) by adding 49. Adding 9 and 49 to both sides of the equation leaves \((x - 3)^2 + (y - 7)^2 = 4\)—the same circle as in Example 1.

**Quicker Solution** Match the given equation with \( (4) \). Then \( h = 3, k = 7, \) and \( r = 2 \):

\[ x^2 - 6x + y^2 - 14y = -54 \quad \text{must agree with} \quad x^2 - 2hx + h^2 + y^2 - 2ky + k^2 = r^2. \]

The change to \( X = x - h \) and \( Y = y - k \) moves the center of the circle from \((h, k)\) to \((0, 0)\). This is equally true for an ellipse:

The ellipse \( \frac{(x - h)^2}{a^2} + \frac{(y - k)^2}{b^2} = 1 \) becomes \( \frac{X^2}{a^2} + \frac{Y^2}{b^2} = 1 \).

When we rescale by \( x = X/a \) and \( y = Y/b \), we get the unit circle \( x^2 + y^2 = 1 \).

The unit circle has area \( \pi \). The ellipse has area \( \pi ab \) (proved later in the book). The distance around the circle is \( 2\pi \). The distance around an ellipse does not rescale—it has no simple formula.

![Figure 3.17](image)

**Fig. 3.17** Uncentered circle. Centered ellipse \( x^2/3^2 + y^2/2^2 = 1 \). The distance from center to far right is also \( a = 3 \). All rays from \( F_2 \) reflect to \( F_1 \).

Now we leave circles and concentrate on ellipses. They have two foci (pronounced fo-sigh). For a parabola, the second focus is at infinity. For a circle, both foci are at the center. The foci of an ellipse are on its longer axis (its major axis), one focus on each side of the center:

\[ F_1 \text{ is at } x = c = \sqrt{a^2 - b^2} \text{ and } F_2 \text{ is at } x = -c. \]

The right triangle in Figure 3.17 has sides \( a, b, c \). From the top of the ellipse, the distance to each focus is \( a \). From the endpoint at \( x = a \), the distances to the foci are \( a + c \) and \( a - c \). Adding \((a + c) + (a - c)\) gives \( 2a \). As you go around the ellipse, the distance to \( F_1 \) plus the distance to \( F_2 \) is constant (always \( 2a \)).
3.5 Parabolas, Ellipses, and Hyperbolas

At all points on the ellipse, the sum of distances from the foci is 2a. This is another equation for the ellipse:

from $F_1$ and $F_2$ to $(x,y)$: \[ \sqrt{(x-c)^2+y^2} + \sqrt{(x+c)^2+y^2} = 2a. \] (5)

To draw an ellipse, tie a string of length 2a to the foci. Keep the string taut and your moving pencil will create the ellipse. This description uses $a$ and $c$—the other form uses $a$ and $b$ (remember $b^2 + c^2 = a^2$). Problem 24 asks you to simplify equation (5) until you reach \[ \frac{x^2}{a^2} + \frac{y^2}{b^2} = 1. \]

The “whispering gallery” of the United States Senate is an ellipse. If you stand at one focus and speak quietly, you can be heard at the other focus (and nowhere else). Your voice is reflected off the walls to the other focus—following the path of the string. For a parabola the rays come in to the focus from infinity—where the second focus is.

A hospital uses this reflection property to split up kidney stones. The patient sits inside an ellipse with the kidney stone at one focus. At the other focus a lithotripter sends out hundreds of small shocks. You get a spinal anesthetic (I mean the patient) and the stones break into tiny pieces.

The most important focus is the Sun. The ellipse is the orbit of the Earth. See Section 12.4 for a terrible printing mistake by the Royal Mint, on England’s last pound note. They put the Sun at the center.

**Question 1** Why do the whispers (and shock waves) arrive together at the second focus?

**Answer** Whichever way they go, the distance is 2a. Exception: straight path is 2c.

**Question 2** Locate the ellipse with equation $4x^2 + 9y^2 = 36$.

**Answer** Divide by 36 to change the constant to 1. Now identify $a$ and $b$:
\[ \frac{x^2}{9} + \frac{y^2}{4} = 1 \] so $a = \sqrt{9}$ and $b = \sqrt{4}$. Foci at $\pm \sqrt{9-4} = \pm \sqrt{5}$.

**Question 3** Shift the center of that ellipse across and down to $x = 1$, $y = -5$.

**Answer** Change $x$ to $x - 1$. Change $y$ to $y + 5$. The equation becomes $(x - 1)^2/9 + (y + 5)^2/4 = 1$. In practice we start with this uncentered ellipse and go the other way to center it.

**HYPERBOLAS** \[ \frac{y^2}{a^2} - \frac{x^2}{b^2} = 1 \]

*Notice the minus sign for a hyperbola.* That makes all the difference. Unlike an ellipse, $x$ and $y$ can both be large. The curve goes out to infinity. It is still symmetric, since $x$ can change to $-x$ and $y$ to $-y$.

The center is at $(0, 0)$. Solving for $y$ again yields two functions (+ and −):
\[ \frac{y^2}{a^2} - \frac{x^2}{b^2} = 1 \] gives \[ \frac{y}{a} = \pm \sqrt{1 + \frac{x^2}{b^2}} \] or \[ y = \pm \frac{a}{b} \sqrt{b^2 + x^2}. \] (6)

The hyperbola has two branches that never meet. The upper branch, with a plus sign, has $y \geq a$. The vertex $V_1$ is at $x = 0$, $y = a$—the lowest point on the branch. Much further out, when $x$ is large, the hyperbola climbs up beside its *sloping asymptotes*:
\[ \text{if } \frac{x^2}{b^2} = 1000 \text{ then } \frac{y^2}{a^2} = 1001. \text{ So } \frac{y}{a} \text{ is close to } \frac{x}{b} \text{ or } -\frac{x}{b}. \]
The asymptotes are the lines $y/a = x/b$ and $y/a = -x/b$. Their slopes are $a/b$ and $-a/b$. You can't miss them in Figure 3.18.

For a hyperbola, the foci are inside the two branches. Their distance from the center is still called $c$. But now $c = \sqrt{a^2 + b^2}$, which is larger than $a$ and $b$. The vertex is a distance $c - a$ from one focus and $c + a$ from the other. The difference (not the sum) is $(c + a) - (c - a) = 2a$.

All points on the hyperbola have this property: The difference between distances to the foci is constantly $2a$. A ray coming in to one focus is reflected toward the other. The reflection is on the outside of the hyperbola, and the inside of the ellipse.

Here is an application to navigation. Radio signals leave two fixed transmitters at the same time. A ship receives the signals a millisecond apart. Where is the ship?

*Answer:* It is on a hyperbola with foci at the transmitters. Radio signals travel 186 miles in a millisecond, so $186 = 2a$. This determines the curve. In Long Range Navigation (LORAN) a third transmitter gives another hyperbola. Then the ship is located exactly.

**Question 4** How do hyperbolas differ from parabolas, far from the center?

*Answer:* Hyperbolas have asymptotes. Parabolas don't.

The hyperbola has a natural rescaling. The appearance of $x/b$ is a signal to change to $X$. Similarly $y/a$ becomes $Y$. Then $Y = 1$ at the vertex, and we have a standard hyperbola:

$$y^2/a^2 - x^2/b^2 = 1 \quad \text{becomes} \quad Y^2 - X^2 = 1.$$ 

A $90^\circ$ turn gives $X^2 - Y^2 = 1$—the hyperbola opens to the sides. A $45^\circ$ turn produces $2XY = 1$. We show below how to recognize $x^2 + xy + y^2 = 1$ as an ellipse and $x^2 + 3xy + y^2 = 1$ as a hyperbola. (They are not circles because of the $xy$ term.) When the $xy$ coefficient increases past 2, $x^2 + y^2$ no longer indicates an ellipse.

**Question 5** Locate the hyperbola with equation $9y^2 - 4x^2 = 36$.

*Answer:* Divide by 36. Then $y^2/4 - x^2/9 = 1$. Recognize $a = \sqrt{4}$ and $b = \sqrt{9}$.

**Question 6** Locate the uncentered hyperbola $9y^2 - 18y - 4x^2 - 4x = 28$.

*Answer:* Complete $9y^2 - 18y$ to $9(y - 1)^2$ by adding 9. Complete $4x^2 + 4x$ to $4(x + 1)^2$ by adding $4(1)^2 = 4$. The equation is rewritten as $9(y - 1)^2 - 4(x + 1)^2 = 28 + 9 - 1$. This is the hyperbola in Question 5—except its center is $(-\frac{1}{2}, 1)$.
3.5 Parabolas, Ellipses, and Hyperbolas

To summarize: Find the center by completing squares. Then read off \( a \) and \( b \).

**THE GENERAL EQUATION** \( Ax^2 + Bxy + Cy^2 + Dx + Ey + F = 0 \)

This equation is of second degree, containing any and all of \( 1, x, y, x^2, xy, y^2 \).

A plane is cutting through a cone. **Is the curve a parabola or ellipse or hyperbola?**

Start with the most important case \( Ax^2 + Bxy + Cy^2 = 1 \).

| 31 | The equation \( Ax^2 + Bxy + Cy^2 = 1 \) produces a hyperbola if \( B^2 > 4AC \) and an ellipse if \( B^2 < 4AC \). A parabola has \( B^2 = 4AC \). |

To recognize the curve, we remove \( Bxy \) by rotating the plane. This also changes \( A \) and \( C \)—but the combination \( B^2 - 4AC \) is not changed (proof omitted). An example is \( 2xy = 1 \), with \( B^2 = 4 \). It rotates to \( y^2 - x^2 = 1 \), with \( -4AC = 4 \). That positive number 4 signals a hyperbola—since \( A = -1 \) and \( C = 1 \) have opposite signs.

Another example is \( x^2 + y^2 = 1 \). It is a circle (a special ellipse). However we rotate, the equation stays the same. The combination \( B^2 - 4AC = 0 - 4 \cdot 1 \cdot 1 \) is negative, as predicted for ellipses.

To rotate by an angle \( \alpha \), change \( x \) and \( y \) to new variables \( x' \) and \( y' \):

\[
\begin{align*}
x' &= x \cos \alpha - y \sin \alpha \\
y' &= x \sin \alpha + y \cos \alpha
\end{align*}
\]

Substituting for \( x \) and \( y \) changes \( Ax^2 + Bxy + Cy^2 = 1 \) to \( A'x'^2 + B'x'y' + C'y'^2 = 1 \).

The formulas for \( A', B', C' \) are painful so I go to the key point:

**\( B' \) is zero if the rotation angle \( \alpha \) has \( \tan 2\alpha = B/(A - C) \).**

With \( B' = 0 \), the curve is easily recognized from \( A'x'^2 + C'y'^2 = 1 \). It is a hyperbola if \( A' \) and \( C' \) have opposite signs. Then \( B'^2 - 4A'C' \) is positive. The original \( B^2 - 4AC \) was also positive, because this special combination stays constant during rotation.

After the \( xy \) term is gone, we deal with \( x \) and \( y \)—by centering. To find the center, complete squares as in Questions 3 and 6. For total perfection, rescale to one of the model equations \( y = x^2 \) or \( x^2 + y^2 = 1 \) or \( y^2 - x^2 = 1 \).

The remaining question is about \( F = 0 \). What is the graph of \( Ax^2 + Bxy + Cy^2 = 0 \)?

The ellipse-hyperbola-parabola have disappeared. But if the Greeks were right, the cone is still cut by a plane. The degenerate case \( F = 0 \) occurs when the plane cuts right through the sharp point of the cone.

A level cut hits only that one point \((0, 0)\). The equation shrinks to \( x^2 + y^2 = 0 \), a circle with radius zero. A steep cut gives two lines. The hyperbola becomes \( y^2 - x^2 = 0 \), leaving only its asymptotes \( y = \pm x \). A cut at the exact angle of the cone gives only one line, as in \( x^2 = 0 \). A **single point**, **two lines**, and **one line** are very extreme cases of an ellipse, hyperbola, and parabola.

All these "conic sections" come from planes and cones. The beauty of the geometry, which Archimedes saw, is matched by the importance of the equations. Galileo discovered that projectiles go along parabolas (Chapter 12). Kepler discovered that the Earth travels on an ellipse (also Chapter 12). Finally Einstein discovered that light travels on hyperbolas. That is in four dimensions, and not in Chapter 12.
The graph of \(y = x^2 + 2x + 5\) is a \(\text{u}\). Its lowest point (the vertex) is \((x, y) = \left(\text{b}\right)\). Centering by \(X = x + 1\) and \(Y = \text{c}\) moves the vertex to \((0, 0)\). The equation becomes \(Y = \text{d}\). The focus of this centered parabola is \(\text{e}\). All rays coming straight down are \(f\) to the focus.

The graph of \(x^2 + 4y^2 = 16\) is an \(g\). Dividing by \(h\) leaves \(x^2/a^2 + y^2/b^2 = 1\) with \(a = \text{i}\) and \(b = \text{j}\). The graph lies in the rectangle whose sides are \(k\). The area is \(n\text{ab} = \text{l}\). The foci are at \(x = \pm c = \text{m}\). The sum of distances from the foci to a point on this ellipse is always \(n\). If we rescale to \(X = x/4\) and \(Y = y/2\) the equation becomes \(\text{o}\) and the graph becomes a \(p\).

The graph of \(y^2 - x^2 = 9\) is a \(q\). Dividing by 9 leaves \(y^2/a^2 - x^2/b^2 = 1\) with \(a = \text{t}\) and \(b = \text{u}\). On the upper branch \(y \geq 9\). The asymptotes are the lines \(w\). The foci are at \(y = \pm c = \text{v}\). The \(w\) of distances from the foci to a point on this hyperbola is \(x\).

All these curves are conic sections—the intersection of a \(y\) and a \(z\). A steep cutting angle yields a \(A\). At the borderline angle we get a \(B\). The general equation is \(Ax^2 + Cy^2 + Dxy + F = 0\). If \(D = E = 0\) the center of the graph is at \(D\). The equation \(4x^2 + 5xy + 6y^2 = 1\) gives an ellipse when \(E\). The graph of \(x^2 + 5x + 6y = 1\) is a \(F\).

1. The vertex of \(y = ax^2 + bx + c\) is at \(x = -b/2a\). What is special about this \(x\)? Show that it gives \(c = (b^2/4a)\).

2. The parabola \(y = 3x^2 - 12x\) has \(x_{\text{min}} = \text{w}\). At this minimum, \(3x^2\) is \(\text{w}\) as large as \(12x\). Introducing \(X = x - 2\) and \(Y = y + 12\) centers the equation to \(\text{w}\).

**Draw the curves 3–14 by hand or calculator or computer. Locate the vertices and foci.**

<table>
<thead>
<tr>
<th>equation</th>
<th>vertices</th>
<th>foci</th>
</tr>
</thead>
<tbody>
<tr>
<td>(P) (y = ax^2 + bx + c)</td>
<td>(\left(- \frac{b}{2a}, c - \frac{b^2}{4a}\right))</td>
<td>(\frac{1}{4a}) above vertex, also infinity</td>
</tr>
<tr>
<td>(E) (\frac{x^2}{a^2} + \frac{y^2}{b^2} = 1, a &gt; b)</td>
<td>((a, 0)) and ((-a, 0))</td>
<td>((c, 0)) and ((-c, 0)): (c = \sqrt{a^2 - b^2})</td>
</tr>
<tr>
<td>(H) (\frac{y^2}{a^2} - \frac{x^2}{b^2} = 1)</td>
<td>((0, a)) and ((0, -a))</td>
<td>((0, c)) and ((0, -c)): (c = \sqrt{a^2 + b^2})</td>
</tr>
</tbody>
</table>

**3.5 EXERCISES**

**Read-through questions**

The graph of \(y = x^2 + 2x + 5\) is \(\text{u}\). Its lowest point (the vertex) is \((x, y) = \left(\text{b}\right)\). Centering by \(X = x + 1\) and \(Y = \text{c}\) moves the vertex to \((0, 0)\). The equation becomes \(Y = \text{d}\). The focus of this centered parabola is \(\text{e}\). All rays coming straight down are \(f\) to the focus.

The graph of \(x^2 + 4y^2 = 16\) is \(\text{g}\). Dividing by \(h\) leaves \(x^2/a^2 + y^2/b^2 = 1\) with \(a = \text{i}\) and \(b = \text{j}\). The graph lies in the rectangle whose sides are \(k\). The area is \(n\text{ab} = \text{l}\). The foci are at \(x = \pm c = \text{m}\). The sum of distances from the foci to a point on this ellipse is always \(n\). If we rescale to \(X = x/4\) and \(Y = y/2\) the equation becomes \(\text{o}\) and the graph becomes a \(p\).

The graph of \(y^2 - x^2 = 9\) is \(\text{q}\). Dividing by 9 leaves \(y^2/a^2 - x^2/b^2 = 1\) with \(a = \text{t}\) and \(b = \text{u}\). On the upper branch \(y \geq 9\). The asymptotes are the lines \(w\). The foci are at \(y = \pm c = \text{v}\). The \(w\) of distances from the foci to a point on this hyperbola is \(x\).

All these curves are conic sections—the intersection of a \(y\) and a \(z\). A steep cutting angle yields a \(A\). At the borderline angle we get a \(B\). The general equation is \(Ax^2 + Cy^2 + Dxy + F = 0\). If \(D = E = 0\) the center of the graph is at \(D\). The equation \(4x^2 + 5xy + 6y^2 = 1\) gives an ellipse when \(E\). The graph of \(x^2 + 5x + 6y = 1\) is a \(F\).

1. The vertex of \(y = ax^2 + bx + c\) is at \(x = -b/2a\). What is special about this \(x\)? Show that it gives \(c = (b^2/4a)\).

2. The parabola \(y = 3x^2 - 12x\) has \(x_{\text{min}} = \text{w}\). At this minimum, \(3x^2\) is \(\text{w}\) as large as \(12x\). Introducing \(X = x - 2\) and \(Y = y + 12\) centers the equation to \(\text{w}\).

11. \(y^2 - 4x^2 = 1\)

12. \((y - 1)^2 - 4x^2 = 1\)

13. \(y^2 - x^2 = 0\)

14. \(xy = 0\)

**Problems 15–20 are about parabolas, 21–34 are about ellipses, 35–41 are about hyperbolas.**

15. Find the parabola \(y = ax^2 + bx + c\) that goes through \((0, 0)\) and \((1, 1)\) and \((2, 12)\).

16. \(y = x^2 - x\) has vertex at \(\text{w}\). To move the vertex to \((0, 0)\) set \(X = \text{w}\) and \(Y = \text{w}\). Then \(Y = X^2\).

17. (a) In equation (2) change \(\frac{1}{p}\) to \(p\). Square and simplify.

(b) Locate the focus and directrix of \(Y = 3X^2\). Which points are a distance 1 from the directrix and focus?

18. The parabola \(y = 9 - x^2\) opens \(\text{w}\) with vertex at \(\text{w}\). Centering by \(Y = y - 9\) yields \(Y = -x^2\).

19. Find equations for all parabolas which

(a) open to the right with vertex at \((0, 0)\)

(b) open upwards with focus at \((0, 0)\)

(c) open downwards and go through \((0, 0)\) and \((1, 0)\).

20. A projectile is at \(x = t, y = t - t^2\) at time \(t\). Find \(dx/dt\) and \(dy/dt\) at the start, the maximum height, and an \(xy\) equation for the path.

21. Find the equation of the ellipse with extreme points at \((\pm 2, 0)\) and \((0, \pm 1)\). Then shift the center to \((1, 1)\) and find the new equation.

22. On the ellipse \(x^2/a^2 + y^2/b^2 = 1\), solve for \(y\) when \(x = c = \sqrt{a^2 - b^2}\). This height above the focus will be valuable in proving Kepler's third law.

23. Find equations for the ellipses with these properties:

(a) through \((5, 0)\) with focus at \((\pm 4, 0)\)

(b) with sum of distances to \((1, 1)\) and \((5, 1)\) equal to \(12\)

(c) with both foci at \((0, 0)\) and sum of distances = \(2a = 10\).

24. Move a square root to the right side of equation (5) and square both sides. Then isolate the remaining square root and square again. Simplify to reach the equation of an ellipse.
25 Decide between circle-ellipse-parabola-hyperbola, based on the \( XY \) equation with \( X = x - 1 \) and \( Y = y + 3 \).
   
   (a) \( x^2 - 2x + y^2 + 6y = 6 \)
   
   (b) \( x^2 - 2x - y^2 - 6y = 6 \)
   
   (c) \( x^2 - 2x + 2y^2 + 12y = 6 \)
   
   (d) \( x^2 - 2x - y = 6 \).

26 A tilted cylinder has equation \((x - 2y - 2z)^2 + (y - 2x - 2z)^2 = 1\). Show that the water surface at \( z = 0 \) is an ellipse. What is its equation and what is \( B \)?

27 \((4, 9/5)\) is above the focus on the ellipse \( x^2/25 + y^2/9 = 1 \). Find \( dy/dx \) at that point and the equation of the tangent line.

28 (a) Check that the line \( xy_0 + yy_0 = r^2 \) is tangent to the circle \( x^2 + y^2 = r^2 \) at \((x_0, y_0)\).

   (b) For the ellipse \( x^2/a^2 + y^2/b^2 = 1 \) show that the tangent equation is \( xy_0/a^2 + yy_0/b^2 = 1 \). (Check the slope.)

29 The slope of the normal line in Figure A is \( s = -1/(\text{slope of tangent}) = \frac{1}{s} \). The slope of the line from \( F_2 \) is \( S = \frac{1}{s} \). By the reflection property,

\[
S = \cot 2\theta = \frac{1}{2} \left( \cot \theta + \tan \theta \right) = \frac{1}{2} \left( \frac{s - 1}{s} \right).
\]

Test your numbers \( s \) and \( S \) against this equation.

30 Figure B proves the reflecting property of an ellipse. \( R \) is the mirror image of \( F_1 \) in the tangent line; \( Q \) is any other point on the line. Deduce steps 2, 3, 4 from 1, 2, 3:

1. \( PF_1 + PF_2 < QF_1 + QF_2 \) (left side = \( 2a \), \( Q \) is outside)
2. \( PR + PF_2 < QR + QF_2 \)
3. \( P \) is on the straight line from \( F_2 \) to \( R \)
4. \( x = \beta \); the reflecting property is proved.

31 The ellipse \((x - 3)^2/4 + (y - 1)^2/4 = 1\) is really a _______ with center at _______ and radius _______. Choose \( X \) and \( Y \) to produce \( X^2 + Y^2 = 1 \).

32 Compute the area of a square that just fits inside the ellipse \( x^2/a^2 + y^2/b^2 = 1 \).

33 Rotate the axes of \( x^2 + xy + y^2 = 1 \) by using equation (7) with \( \sin \chi = \cos \chi = 1/\sqrt{2} \). The \( x'y' \) equation should show an ellipse.

34 What are \( a, b, c \) for the Earth's orbit around the sun?

35 Find an equation for the hyperbola with
   
   (a) vertices \((0, \pm 1)\), foci \((0, \pm 2)\)
   
   (b) vertices \((0, \pm 3)\), asymptotes \( y = \pm 2x \)
   
   (c) \((2, 3)\) on the curve, asymptotes \( y = \pm x \)

36 Find the slope of \( y^2 - x^2 = 1 \) at \((x_0, y_0)\). Show that \( yy_0 - xx_0 = 1 \) goes through this point with the right slope (it has to be the tangent line).

37 If the distances from \((x, y)\) to \((8, 0)\) and \((-8, 0)\) differ by 10, what hyperbola contains \((x, y)\)?

38 If a cannon was heard by Napoleon and one second later by the Duke of Wellington, the cannon was somewhere on a _______ with foci at _______.

39 \( y^2 - 4x \) is part of \((y - 2)^2 = \) _______ and \( 2x^2 + 12x \) is part of \((2x + 3)^2 = \) _______. Therefore \( y^2 - 4y - 2x^2 - 12x = 0 \) gives the hyperbola \((y - 2)^2 - 2(2x + 3)^2 = \) _______. Its center is _______ and it opens to the _______.

40 Following Problem 39 turn \( y^2 + 2y = x^2 + 10x \) into \( Y^2 = X^2 + C \) with \( X, Y, \) and \( C \) equal to _______.

41 Draw the hyperbola \( x^2 - 4y^2 = 1 \) and find its foci and asymptotes.

Problems 42–46 are about second-degree curves (conics).

42 For which \( A, C, F \) does \( Ax^2 + Cy^2 + F = 0 \) have no solution (empty graph)?

43 Show that \( x^2 + 2xy + y^2 + 2x + 2y + 1 = 0 \) is the equation (squared) of a single line.

44 Given any _______ points in the plane, a second-degree curve \( Ax^2 + \cdots + F = 0 \) goes through those points.

45 (a) When the plane \( z = ax + by + c \) meets the cone \( z^2 = x^2 + y^2 \), eliminate \( z \) by squaring the plane equation. Rewrite in the form \( Ax^2 + Bxy + Cy^2 + Dx + Ey + F = 0 \).

   (b) Compute \( B^2 - 4AC \) in terms of \( a \) and \( b \).

   (c) Show that the plane meets the cone in an ellipse if \( a^2 + b^2 < 1 \) and a hyperbola if \( a^2 + b^2 > 1 \) (steeper).

46 The roots of \( ax^2 + bx + c = 0 \) also involve the special combination \( b^2 - 4ac \). This quadratic equation has two real roots if _______ and no real roots if _______. The roots come together when \( b^2 = 4ac \), which is the borderline case like a parabola.
Iteration means repeating the same function. Suppose the function is \( F(x) = \cos x \). Choose any starting value, say \( x_0 = 1 \). Take its cosine: \( x_1 = \cos x_0 = .54 \). Then take the cosine of \( x_1 \). That produces \( x_2 = \cos .54 = .86 \). The iteration is \( x_{n+1} = \cos x_n \). I am in radian mode on a calculator, pressing “cos” each time. The early numbers are not important, what is important is the output after 12 or 30 or 100 steps:

**EXAMPLE 1**

\[ x_{12} = .75, \quad x_{13} = .73, \quad x_{14} = .74, \quad ... \quad x_{29} = .7391, \quad x_{30} = .7391. \]

The goal is to explain why the x’s approach \( x^* = .739085 \). Every starting value \( x_0 \) leads to this same number \( x^* \). What is special about .7391?

**Note on iterations**

Do \( x_1 = \cos x_0 \) and \( x_2 = \cos x_1 \) mean that \( x_2 = \cos^2 x_0 \)? Absolutely not! Iteration creates a new and different function \( \cos(\cos x) \). It uses the cos button, not the squaring button. The third step creates \( F(F(F(x))) \). As soon as you can, iterate with \( x_{n+1} = \frac{1}{2} \cos x_n \). What limit do the x’s approach? Is it \( \frac{1}{2}(.7931) \)?

Let me slow down to understand these questions. The central idea is expressed by the equation \( x_{n+1} = F(x_n) \). Substituting \( x_0 \) into \( F \) gives \( x_1 \). This output \( x_1 \) is the input that leads to \( x_2 \). In its turn, \( x_2 \) is the input and out comes \( x_3 = F(x_2) \). This is iteration, and it produces the sequence \( x_0, x_1, x_2, ... \).

The x’s may approach a limit \( x^* \), depending on the function \( F \). Sometimes \( x^* \) also depends on the starting value \( x_0 \). Sometimes there is no limit. Look at a second example, which does not need a calculator.

**EXAMPLE 2**

\[ x_{n+1} = F(x_n) = \frac{1}{2}x_n + 4. \]

Starting from \( x_0 = 0 \) the sequence is

\[ x_1 = \frac{1}{2} \cdot 0 + 4 = 4, \quad x_2 = \frac{1}{2} \cdot 4 + 4 = 6, \quad x_3 = \frac{1}{2} \cdot 6 + 4 = 7, \quad x_4 = \frac{1}{2} \cdot 7 + 4 = 7\frac{1}{2}, \quad ... \]

Those numbers 0, 4, 6, 7, 7\frac{1}{2}, ... seem to be approaching \( x^* = 8 \). A computer would convince us. So will mathematics, when we see what is special about 8:

When the x’s approach \( x^* \), the limit of \( x_{n+1} = \frac{1}{2}x_n + 4 \)

is \( x^* = \frac{1}{2}x^* + 4 \). This limiting equation yields \( x^* = 8 \).

8 is the “steady state” where input equals output: \( 8 = F(8) \). It is the fixed point.

If we start at \( x_0 = 8 \), the sequence is 8, 8, 8, .... When we start at \( x_0 = 12 \), the sequence goes back toward 8:

\[ x_1 = \frac{1}{2} \cdot 12 + 4 = 10, \quad x_2 = \frac{1}{2} \cdot 10 + 4 = 9, \quad x_3 = \frac{1}{2} \cdot 9 + 4 = 8.5, \quad ... \]

**Equation for limit:** If the iterations \( x_{n+1} = F(x_n) \) converge to \( x^* \), then \( x^* = F(x^*) \).

To repeat: 8 is special because it equals \( \frac{1}{2} \cdot 8 + 4 \). The number .7391... is special because it equals \( \cos .7391 \). The graphs of \( y = x \) and \( y = F(x) \) intersect at \( x^* \). To explain why the x’s converge (or why they don’t) is the job of calculus.

**EXAMPLE 3**

\[ x_{n+1} = x_n^2 \] has two fixed points: 0 = 0² and 1 = 1². Here \( F(x) = x^2 \).

Starting from \( x_0 = \frac{1}{2} \) the sequence \( \frac{1}{4}, \frac{1}{16}, \frac{1}{256}, ... \) goes quickly to \( x^* = 0 \). The only approaches to \( x^* = 1 \) are from \( x_0 = 1 \) (of course) and from \( x_0 = -1 \). Starting from \( x_0 = 2 \) we get 4, 16, 256, ... and the sequence diverges to \( +\infty \).

Each limit \( x^* \) has a “basin of attraction.” The basin contains all starting points \( x_0 \) that lead to \( x^* \). For Examples 1 and 2, every \( x_0 \) led to .7391 and 8. The basins were
3.6 Iterations \( x_{n+1} = F(x_n) \)

the whole line (that is still to be proved). Example 3 had three basins—the interval \(-1 < x_0 < 1\), the two points \(x_0 = \pm 1\), and all the rest. The outer basin \(|x_0| > 1\) led to \(\pm \infty\). I challenge you to find the limits and the basins of attraction (by calculator) for \(F(x) = x - \tan x\).

In Example 3, \(x^* = 0\) is attracting. Points near \(x^*\) move toward \(x^*\). The fixed point \(x^* = 1\) is repelling. Points near 1 move away. We now find the rule that decides whether \(x^*\) is attracting or repelling. The key is the slope \(dF/dx\) at \(x^*\).

3J Start from any \(x_0\) near a fixed point \(x^* = F(x^*)\):

- \(x^*\) is attracting if \(|dF/dx|\) is below 1 at \(x^*\);
- \(x^*\) is repelling if \(|dF/dx|\) is above 1 at \(x^*\).

First I will give a calculus proof. Then comes a picture of convergence, by "cobwebs." Both methods throw light on this crucial test for attraction: \(|dF/dx| < 1\).

First proof: Subtract \(x^* = F(x^*)\) from \(x_{n+1} = F(x_n)\). The difference \(x_{n+1} - x^*\) is the same as \(F(x_n) - F(x^*)\). This is \(\Delta F\). The basic idea of calculus is that \(\Delta F\) is close to \(F'(\Delta x)\):

\[ x_{n+1} - x^* = F(x_n) - F(x^*) \approx F'(x^*)(x_n - x^*). \]  

(1)

The "error" \(x_n - x^*\) is multiplied by the slope \(dF/dx\). The next error \(x_{n+1} - x^*\) is smaller or larger, based on \(|F'| < 1\) or \(|F'| > 1\) at \(x^*\). Every step multiplies approximately by \(F'(x^*)\). Its size controls the speed of convergence.

In Example 1, \(F(x) = \cos x\) and \(F'(x) = -\sin x\). There is attraction to \(.7391\) because \(|\sin x^*| < 1\). In Example 2, \(F = \frac{1}{2}x + 4\) and \(F' = \frac{1}{2}\). There is attraction to 8. In Example 3, \(F = x^2\) and \(F' = 2x\). There is superattraction to \(x^* = 0\) (where \(F' = 0\)). There is repulsion from \(x^* = 1\) (where \(F' = 2\)).

I admit one major difficulty. The approximation in equation (1) only holds near \(x^*\). If \(x_0\) is far away, does the sequence still approach \(x^*\)? When there are several attracting points, which \(x^*\) do we reach? This section starts with good iterations, which solve the equation \(x^* = F(x^*)\) or \(f(x) = 0\). At the end we discover Newton's method. The next section produces crazy but wonderful iterations, not converging and not blowing up. They lead to "fractals" and "Cantor sets" and "chaos."

The mathematics of iterations is not finished. It may never be finished, but we are converging on the answers. Please choose a function and join in.

THE GRAPH OF AN ITERATION: COBWEBS

The iteration \(x_{n+1} = F(x_n)\) involves two graphs at the same time. One is the graph of \(y = F(x)\). The other is the graph of \(y = x\) (the 45° line). The iteration jumps back and forth between these graphs. It is a very convenient way to see the whole process.

Example 1 was \(x_{n+1} = \cos x_n\). Figure 3.19 shows the graph of \(\cos x\) and the "cobweb." Starting at \((x_0, x_0)\) on the 45° line, the rule is based on \(x_1 = F(x_0)\):

- From \((x_0, x_0)\) go up or down to \((x_0, x_1)\) on the curve.
- From \((x_0, x_1)\) go across to \((x_1, x_1)\) on the 45° line.

These steps are repeated forever. From \(x_1\) go up to the curve at \(F(x_1)\). That height is \(x_2\). Now cross to the 45° line at \((x_2, x_2)\). The iterations are aiming for \((x^*, x^*) = (0.7391, 0.7391)\). This is the crossing point of the two graphs \(y = F(x)\) and \(y = x\).
Example 2 was $x_{n+1} = \frac{1}{2} x_n + 4$. Both graphs are straight lines. The cobweb is one-sided, from $(0, 0)$ to $(0, 4)$ to $(4, 4)$ to $(4, 6)$ to $(6, 6)$. Notice how $y$ changes (vertical line) and then $x$ changes (horizontal line). The slope of $F(x)$ is $\frac{1}{2}$, so the distance to 8 is multiplied by $\frac{1}{2}$ at every step.

Example 3 was $x_{n+1} = x_n^2$. The graph of $y = x^2$ crosses the 45° line at two fixed points: $0^2 = 0$ and $1^2 = 1$. Figure 3.20a starts the iteration close to 1, but it quickly goes away. This fixed point is repelling because $F'(1) = 2$. Distance from $x^* = 1$ is doubled (at the start). One path moves down to $x^* = 0$—which is superattractive because $F' = 0$. The path from $x_0 > 1$ diverges to infinity.

**EXAMPLE 4** $F(x)$ has two attracting points $x^*$ (a repelling $x^*$ is always between).

Figure 3.20b shows two crossings with slope zero. The iterations and cobwebs converge quickly. In between, the graph of $F(x)$ must cross the 45° line from below. That requires a slope greater than one. Cobwebs diverge from this unstable point, which separates the basins of attraction. The fixed point $x = \pi$ is in a basin by itself!

*Note 1* To draw cobwebs on a calculator, graph $y = F(x)$ on top of $y = x$. On a Casio, one way is to plot $(x_0, x_0)$ and give the command **LINE**: **PLOT** $X, Y$ followed by **EXE**. Now move the cursor vertically to $y = F(x)$ and press **EXE**. Then move horizontally to $y = x$ and press **EXE**. Continue. Each step draws a line.
For the TI-81 (and also the Casio) a short program produces a cobweb. Store $F(x)$ in the $Y = \text{function slot } Y_1$. Set the range (square window or autoscaling). Run the program and answer the prompt with $x_0$:

```
```

Note 2 The $x$'s approach $x^*$ from one side when $0 < \frac{dF}{dx} < 1$.

Note 3 A basin of attraction can include faraway $x_n$'s (basins can come in infinitely many pieces). This makes the problem interesting. If no fixed points are attracting, see Section 3.7 for “cycles” and “chaos.”

**THE ITERATION** $x_{n+1} = x_n - cf(x_n)$

At this point we offer the reader a choice. One possibility is to jump ahead to the next section on “Newton’s Method.” That method is an iteration to solve $f(x) = 0$. The function $F(x)$ combines $x_n$ and $f(x_n)$ and $f'(x_n)$ into an optimal formula for $x_{n+1}$. We will see how quickly Newton’s method works (when it works). It is the outstanding algorithm to solve equations, and it is totally built on tangent approximations.

The other possibility is to understand (through calculus) a whole family of iterations. This family depends on a number $c$, which is at our disposal. The best choice of $c$ produces Newton’s method. I emphasize that iteration is by no means a new and peculiar idea. It is a fundamental technique in scientific computing.

We start by recognizing that there are many ways to reach $f(x^*) = 0$. (I write $x^*$ for the solution.) A good algorithm may switch to Newton as it gets close. The iterations use $f(x_n)$ to decide on the next point $x_{n+1}$:

$$x_{n+1} = F(x_n) = x_n - cf(x_n).$$

(2)

Notice how $F(x)$ is constructed from $f(x)$ — they are different! We move $f$ to the right side and multiply by a “preconditioner” $c$. The choice of $c$ (or $c_n$, if it changes from step to step) is absolutely critical. The starting guess $x_0$ is also important — but its accuracy is not always under our control.

Suppose the $x_n$ converge to $x^*$. Then the limit of equation (2) is

$$x^* = x^* - cf(x^*).$$

(3)

That gives $f(x^*) = 0$. If the $x_n$'s have a limit, it solves the right equation. It is a fixed point of $F$ (we can assume $c_n \to c \neq 0$ and $f(x_n) \to f(x^*)$). There are two key questions, and both of them are answered by the slope $F'(x^*)$:

1. How quickly does $x_n$ approach $x^*$ (or do the $x_n$ diverge)?
2. What is a good choice of $c$ (or $c_n$)?

**EXAMPLE 5** $f(x) = ax - b$ is zero at $x^* = b/a$. The iteration $x_{n+1} = x_n - c(ax_n - b)$ intends to find $b/a$ without actually dividing. (Early computers could not divide; they used iteration.) Subtracting $x^*$ from both sides leaves an equation for the error:

$$x_{n+1} - x^* = x_n - x^* - c(ax_n - b).$$

Replace $b$ by $ax^*$. The right side is $(1 - ca)(x_n - x^*)$. This “error equation” is

$$(\text{error})_{n+1} = (1 - ca)(\text{error})_n.$$  

(4)
At every step the error is multiplied by \((1 - ca)\), which is \(F'\). The error goes to zero if \(|F'| < 1\). The absolute value \(|1 - ca|\) decides everything:

\[
x_n \text{ converges to } x^* \text{ if and only if } -1 < 1 - ca < 1.
\]  
(5)

The perfect choice (if we knew it) is \(c = 1/a\), which turns the multiplier \(1 - ca\) into zero. Then one iteration gives the exact answer: \(x_1 = x_0 - (1/a)(ax_0 - b) = b/a\). That is the horizontal line in Figure 3.21a, converging in one step. But look at the other lines.

This example did not need calculus. Linear equations never do. The key idea is that close to \(x^*\) the nonlinear equation \(f(x) = 0\) is nearly linear. We apply the tangent approximation. You are seeing how calculus is used, in a problem that doesn’t start by asking for a derivative.

**THE BEST CHOICE OF \(c\)**

The immediate goal is to study the errors \(x_n - x^*\). They go quickly to zero, if the multiplier is small. To understand \(x_{n+1} = x_n - cf(x_n)\), subtract the equation \(x^* = x^* - cf(x^*)\):

\[
x_{n+1} - x^* = x_n - x^* - c(f(x_n) - f(x^*)).
\]  
(6)

Now calculus enters. *When you see a difference of \(f's\ think of \(df/dx\).* Replace \(f(x_n) - f(x^*)\) by \(A(x_n - x^*)\), where \(A\) stands for the slope \(df/dx\) at \(x^*\):

\[
x_{n+1} - x^* \approx (1 - cA)(x_n - x^*).
\]  
(7)

This is the *error equation*. The new error at step \(n + 1\) is approximately the old error multiplied by \(m = 1 - cA\). This corresponds to \(m = 1 - ca\) in the linear example. We keep returning to the basic test \(|m| = |F'(x^*)| < 1\):

**3K** Starting near \(x^*\), the errors \(x_n - x^*\) go to zero if the multiplier has \(|m| < 1\). The perfect choice is \(c = 1/A = 1/f'(x^*). Then m = 1 - cA = 0.\)

There is only one difficulty: *We don't know \(x^*\). Therefore we don't know the perfect \(c\). It depends on the slope \(A = f'(x^*)\) at the unknown solution. However we can come close, by using the slope at \(x_0\):

*Choose \(c_n = 1/f'(x_n). Then x_{n+1} = x_n - f(x_n)/f'(x_n) = F(x_n).\)

This is *Newton’s method*. The multiplier \(m = 1 - cA\) is as near to zero as we can make it. By building \(df/dx\) into \(F(x)\), Newton speeded up the convergence of the iteration.

[Fig. 3.21] The error multiplier is \(m = 1 - cf'(x^*). Newton has \(c = 1/f'(x_n)\) and \(m \to 0.\)
3.6 Iterations \(x_{n+1} = f(x_n)\)

**Example 6** Solve \(f(x) = 2x - \cos x = 0\) with different iterations (different \(c\)'s).

The line \(y = 2x\) crosses the cosine curve somewhere near \(x = \frac{1}{2}\). The intersection point where \(2x^* = \cos x^*\) has no simple formula. We start from \(x_0 = \frac{1}{2}\) and iterate \(x_{n+1} = x_n - c(2x_n - \cos x_n)\) with three different choices of \(c\).

Take \(c = 1\) or \(c = 1/f'(x_0)\) or update \(c\) by Newton's rule \(c_n = 1/f''(x_n)\):

\[
\begin{array}{cccc}
  x_0 &=& 0.50 & c = 1 \\
  x_1 &=& 0.38 & c = 1/f'(x_0) \\
  x_2 &=& 0.55 & c = 1/f''(x_0) \\
  x_3 &=& 0.30 & c = 1/f''(x_0)
\end{array}
\]

The column with \(c = 1\) is diverging (repelled from \(x^*\)). The second column shows convergence (attracted to \(x^*\)). The third column (Newton's method) approaches \(x^*\) so quickly that \(.4501836\) and seven more digits are exact for \(x_3\).

How does this convergence match the prediction? Note that \(f'(x) = 2 + \sin x\) so \(A = 2.435\). Look to see whether the actual errors \(x_n - x^*\), going down each column, are multiplied by the predicted \(m\) below that column:

\[
\begin{array}{ccc}
  c = 1 & c = 1/(2 + \sin \frac{1}{2}) & c = 1/(2 + \sin x_n) \\
  x_0 - x^* &=& 0.05 & 4.98 \cdot 10^{-2} & 4.98 \cdot 10^{-2} \\
  x_1 - x^* &=& -0.07 & 4.43 \cdot 10^{-4} & 4.43 \cdot 10^{-4} \\
  x_2 - x^* &=& 0.10 & 7.88 \cdot 10^{-6} & 3.63 \cdot 10^{-8} \\
  x_3 - x^* &=& -0.15 & 1.41 \cdot 10^{-7} & 2.78 \cdot 10^{-16} \\
\end{array}
\]

The first column shows a multiplier below \(-1\). The errors grow at every step. Because \(m\) is negative the errors change sign—the cobweb goes outward.

The second column shows convergence with \(m = .018\). It takes one genuine Newton step, then \(c\) is fixed. After \(n\) steps the error is closely proportional to \(m^n = (.018)^n\)—that is "linear convergence" with a good multiplier.

The third column shows the "quadratic convergence" of Newton's method. Multiplying the error by \(m\) is more attractive than ever, because \(m \to 0\). In fact \(m\) itself is proportional to the error, so at each step the error is squared. Problem 3.8.31 will show that \((\text{error})_{n+1} \leq M(\text{error})^2\). This squaring carries us from \(10^{-2}\) to \(10^{-4}\) to \(10^{-8}\) to "machine \(\varepsilon\)" in three steps. The number of correct digits is doubled at every step as Newton converges.

\[\text{Note 1}\] The choice \(c = 1\) produces \(x_{n+1} = x_n - f(x_n)\). This is "successive substitution." The equation \(f(x) = 0\) is rewritten as \(x = x - f(x)\), and each \(x_n\) is substituted back to produce \(x_{n+1}\). Iteration with \(c = 1\) does not always fail!

\[\text{Note 2}\] Newton's method is successive substitution for \(f/f'\), not \(f\). Then \(m \approx 0\).

\[\text{Note 3}\] Edwards and Penney happened to choose the same example \(2x = \cos x\). But they cleverly wrote it as \(x_{n+1} = \frac{1}{2} \cos x_n\), which has \(|F'| = |\frac{1}{2} \sin x| < 1\). This iteration fits into our family with \(c = \frac{1}{2}\), and it succeeds. We asked earlier if its limit is \(\frac{1}{2}(.7391)\). No, it is \(x^* = .450\ldots\).
3 Applications of the Derivative

Note 4 The choice \( c = 1/f'(x_0) \) is “modified Newton.” After one step of Newton’s method, \( c \) is fixed. The steps are quicker, because they don’t require a new \( f'(x_0) \). But we need more steps. Millions of dollars are spent on Newton’s method, so speed is important. In all its forms, \( f(x) = 0 \) is the central problem of computing.

3.6 EXERCISES

Read-through questions

\[ x_{n+1} = x_n^2 \] describes an \( \alpha \)_. After one step \( x_1 = \beta \)_.

After two steps \( x_2 = F(x_1) = \gamma \)_. If it happens that input = output, or \( x^* = \delta \)_, then \( x^* \) is an \( \epsilon \) fixed point. \( F = x^3 \) has \( \zeta \) fixed points, at \( x^* = \theta \)_. Starting near a fixed point, the \( x_n \) will converge to it if \( \iota \) < 1. That is because \( x_{n+1} - x^* = F(x_n) - F(x^*) \approx \jmath \). The point is called \( \kappa \)_. The \( x_n \) are repelled if \( \lambda \)_. For \( F = x^3 \) the fixed points have \( F'(x^*) = \mu \). The cobweb goes from \( (x_0, x_n) \) to \( (\theta, \iota) \) to \( (\kappa, \lambda) \) and converges to \( (x^*, x^*) = \nu \). This is an intersection of \( \nu = x^3 \) and \( \nu = \omega \), and it is super-attracting because \( \alpha \)_.

\( f(x) = 0 \) can be solved iteratively by \( x_{n+1} = x_n - cf(x_n) \), in which case \( F'(x^*) = \psi \). Subtracting \( x^* = x^* - cf(x^*) \), the error equation is \( x_{n+1} - x^* = m \). The \( x_n \) are attracted too, \( \theta \).\( c \) produces Newton’s method. The choice \( c = 1 \) is “successive \( \omega \) ” and \( c = \upsilon \) is modified Newton. Convergence to \( x^* \) is \( \omega \) certain.

We have three ways to study iterations \( x_{n+1} = F(x_n) \):

1. Compute \( x_1, x_2, \ldots \) from different \( x_0 \) (2) find the fixed points \( x^* \) and test \( |dF/dx| < 1 \) (3) draw cobwebs.

In Problems 1–8 start from \( x_0 = 0.6 \) and \( x_0 = 2 \). Compute \( x_1, x_2, \ldots \) to test convergence:

1. \( x_{n+1} = x_n^2 - \frac{1}{2} \)
2. \( x_{n+1} = 2x_n(1 - x_n) \)
3. \( x_{n+1} = \sqrt{x_n} \)
4. \( x_{n+1} = \frac{1}{\sqrt{x_n}} \)
5. \( x_{n+1} = 3x_n(1 - x_n) \)
6. \( x_{n+1} = x_n^2 + x_n - 2 \)
7. \( x_{n+1} = \frac{1}{2} x_n - \frac{1}{2} \)
8. \( x_{n+1} = [x_n] \)

9. Check \( dF/dx \) at all fixed points in Problems 1–6. Are they attracting or repelling?

10. From \( x_0 = -1 \) compute the sequence \( x_{n+1} = -x_n^2 \). Draw the cobweb with its “cycle.” Two steps produce \( x_{n+2} = x_n^2 \), which has the fixed points ________.

11. Draw the cobwebs for \( x_{n+1} = \frac{1}{2} x_n - 1 \) and \( x_{n+1} = 1 - \frac{1}{2} x_n \) starting from \( x_0 = 2 \). Rule: Cobwebs are two-sided when \( dF/dx \) is ________.

12. Draw the cobweb for \( x_{n+1} = x_n^2 - 1 \) starting from the periodic point \( x_0 = 0 \). Another periodic point is ________.

13. Solve equations 13–16 within 1% by iteration.

14. \( x = \cos \frac{1}{3} x \)
15. \( x = \cos \frac{1}{2} x \)
16. \( x = 2x - 1 \) (?)

17. For which numbers \( a \) does \( x_{n+1} = a(x_n - x_n^2) \) converge to \( x^* = 0 \)?

18. For which numbers \( a \) does \( x_{n+1} = a(x_n - x_n^2) \) converge to \( x^* = (a - 1)/a \)?

19. Iterate \( x_{n+1} = 4(x_n - x_n^2) \) to see chaos. Why don’t the \( x_n \) approach \( x^* = \frac{5}{2} \)?

20. One fixed point of \( F(x) = x^2 - \frac{1}{2} \) is attracting, the other is repelling. By experiment or cobwebs, find the basin of \( x_0 \)’s that go to the attractor.

21. (Important) Find the fixed point for \( F(x) = ax + s \). When is it attracting?

22. What happens in the linear case \( x_{n+1} = ax_n + 4 \) when \( a = 1 \) and when \( a = -1 \)?

23. Starting with $1000, you spend half your money each year and a rich but foolish aunt gives you a new $1000. What is your steady state balance \( x^* \)? What is \( x^* \) if you start with a million dollars?

24. The US national debt was once $1 trillion. Inflation reduces its real value by 5% each year (so multiply by \( e^{0.05} \)), but overspending adds another $10 billion. What is the steady state debt \( x^* \)?

25. \( x_{n+1} = b/x_n \) has the fixed point \( x^* = \sqrt{b} \). Show that \( |dF/dx| = 1 \) at that point—we have a stable fixed point. Start the cobwebs for \( x_{n+1} = \sin x_n \) and \( x_{n+1} = \tan x_n \). In both cases \( dF/dx \) is 1 at \( x^* = 0 \). (a) Do the iterations converge? (b) Propose a theory based on \( F'' \) for cases when \( F'' = 1 \).

26. Solve \( f(x) = 0 \) in 29–32 by the iteration \( x_{n+1} = x_n - cf(x_n) \), to find a \( c \) that succeeds and a \( c \) that fails.

27. A $5 calculator takes square roots but not cube roots. Explain why \( x_{n+1} = \sqrt[3]{x_n} \) converges to \( \sqrt[3]{2} \).

28. A $5 calculator takes square roots but not cube roots. Explain why \( x_{n+1} = \sqrt[3]{x_n} \) converges to \( \sqrt[3]{2} \).

29. \( f(x) = x^2 - 4 \)
30. \( f(x) = x^2 - 4x + 3 \)
31. \( f(x) = (x - 2)^2 - 1 \)
32. \( f(x) = (1 - x)^{-1} - 3 \)
3.7 Newton's Method (and Chaos)

The equation to be solved is \( f(x) = 0 \). Its solution \( x^* \) is the point where the graph crosses the x axis. Figure 3.22 shows \( x^* \) and a starting guess \( x_0 \). Our goal is to come as close as possible to \( x^* \), based on the information \( f(x_0) \) and \( f'(x_0) \).

Section 3.6 reached Newton's formula for \( x_1 \) (the next guess). We now do that directly.

What do we see at \( x_0 \)? The graph has height \( f(x_0) \) and slope \( f'(x_0) \). We know where we are, and which direction the curve is going. We don't know if the curve bends (we don't have \( f'' \)). The best plan is to follow the tangent line, which uses all the information we have.

Newton replaces \( f(x) \) by its linear approximation (= tangent approximation):

\[
f(x) \approx f(x_0) + f'(x_0)(x - x_0).
\]

We want the left side to be zero. The best we can do is to make the right side zero! The tangent line crosses the axis at \( x_1 \), while the curve crosses at \( x^* \). The new guess \( x_1 \) comes from \( f(x_0) + f'(x_0)(x_1 - x_0) = 0 \). Dividing by \( f'(x_0) \) and solving for \( x_1 \), this is step 1 of Newton's method:

\[
x_1 = x_0 - \frac{f(x_0)}{f'(x_0)}.
\]

At this new point, compute \( f(x_1) \) and \( f'(x_1) \)—the height and slope at \( x_1 \). They give a new tangent line, which crosses at \( x_2 \). At every step we want \( f(x_{n+1}) = 0 \) and we settle for \( f(x_n) + f'(x_n)(x_{n+1} - x_n) = 0 \). After dividing by \( f'(x_n) \), the formula for \( x_{n+1} \) is Newton's method.

(b) Newton's iteration has \( F(x) = x - f(x)/f'(x) \). Show that \( F' = 0 \) when \( f(x) = 0 \). The multiplier for Newton is \( m = 0 \).

40 What are the solutions of \( f(x) = x^2 + 2 = 0 \) and why is Newton's method sure to fail? But carry out the iteration to see whether \( x_n \to \infty \).

41 Computer project \( F(x) = x - \tan x \) has fixed points where \( \tan x^* = 0 \). So \( x^* \) is any multiple of \( \pi \). From \( x_0 = 2.0 \) and \( 1.8 \) and \( 1.9 \), which multiple do you reach? Test points in \( 1.7 < x_0 < 1.9 \) to find basins of attraction to \( \pi \), \( 2\pi \), \( 3\pi \), \( 4\pi \).

Between any two basins there are basins for every multiple of \( \pi \). And more basins between these (a fractal). Mark them on the line from 0 to \( \pi \). Magnify the picture around \( x_0 = 1.9 \) (in color?).

42 Graph \( \cos x \) and \( \cos(\cos x) \) and \( \cos(\cos(\cos x)) \). Also \( \cos^k x \). What are these graphs approaching?

43 Graph \( \sin x \) and \( \sin(\sin x) \) and \( (\sin)^n x \). What are these graphs approaching? Why so slow?
Linear approximation involves three numbers. They are $\Delta x$ (across) and $\Delta f$ (up) and the slope $f'(x)$. If we know two of those numbers, we can estimate the third. It is remarkable to realize that calculus has now used all three calculations—they are the key to this subject:

1. Estimate the slope $f'(x)$ from $\Delta f/\Delta x$  
2. Estimate the change $\Delta f$ from $f'(x)\Delta x$  
3. Estimate the change $\Delta x$ from $\Delta f/f'(x)$

The desired $\Delta f$ is $-f(x_n)$. Formula (3) is exactly $\Delta x = -f(x_n)/f'(x_n)$.

**EXAMPLE 1** (Square roots) $f(x) = x^2 - b$ is zero at $x^* = \sqrt{b}$ and also at $-\sqrt{b}$. Newton's method is a quick way to find square roots—probably built into your calculator. The slope is $f'(x_n) = 2x_n$, and formula (3) for the new guess becomes

$$x_{n+1} = x_n - \frac{x_n^2 - b}{2x_n} = x_n - \frac{1}{2} x_n + \frac{b}{2x_n}.$$  

This simplifies to $x_{n+1} = \frac{1}{2}(x_n + b/x_n)$. **Guess the square root, divide into $b$, and average the two numbers.** The ancient Babylonians had this same idea, without knowing functions or slopes. They iterated $x_{n+1} = F(x_n)$:

$$F(x) = \frac{1}{2} \left(x + \frac{b}{x}\right) \quad \text{and} \quad F'(x) = \frac{1}{2} \left(1 - \frac{b}{x^2}\right).$$  

The Babylonians did exactly the right thing. The slope $F'$ is zero at the solution, when $x^2 = b$. That makes Newton's method converge at high speed. The convergence test is $|F'(x^*)| < 1$. Newton achieves $F'(x^*) = 0$—which is superconvergence.
To find $\sqrt{4}$, start the iteration $x_{n+1} = \frac{1}{2}(x_n + 4/x_n)$ at $x_0 = 1$. Then $x_1 = \frac{1}{2}(1 + 4)$:

$x_1 = 2.5 \quad x_2 = 2.05 \quad x_3 = 2.0006 \quad x_4 = 2.000000009.$

The wrong decimal is twice as far out at each step. The error is squared. Subtracting $x^* = 2$ from both sides of $x_{n+1} = F(x_n)$ gives an error equation which displays that square:

$$x_{n+1} - 2 = \frac{1}{2}(x_n + \frac{4}{x_n}) - 2 = \frac{1}{2x_n}(x_n - 2)^2.$$ (6)

This is $(\text{error})_{n+1} \approx \frac{1}{2}(\text{error})_n^2$. It explains the speed of Newton's method.

**Remark 1** You can't start this iteration at $x_0 = 0$. The first step computes $4/0$ and blows up. Figure 3.22a shows why—the tangent line at zero is horizontal. It will never cross the axis.

**Remark 2** Starting at $x_0 = -1$, Newton converges to $-\sqrt{2}$ instead of $+\sqrt{2}$. That is the other $x^*$. Often it is difficult to predict which $x^*$ Newton's method will choose. Around every solution is a "basin of attraction," but other parts of the basin may be far away. Numerical experiments are needed, with many starts $x_0$. Finding basins of attraction was one of the problems that led to fractals.

**EXAMPLE 2** Solve $\frac{1}{x} - a = 0$ to find $x^* = \frac{1}{a}$ without dividing by $a$.

Here $f(x) = (1/x) - a$. Newton uses $f'(x) = -1/x^2$. Surprisingly, we don't divide:

$$x_{n+1} = x_n - \frac{(1/x_n) - a}{-1/x_n^2} = x_n + x_n - ax_n^2.$$ (7)

Do these iterations converge? I will take $a = 2$ and aim for $x^* = \frac{1}{2}$. Subtracting $\frac{1}{2}$ from both sides of (7) changes the iteration into the error equation:

$$x_{n+1} = 2x_n - 2x_n^2 \quad \text{becomes} \quad x_{n+1} - \frac{1}{2} = -2(x_n - \frac{1}{2})^2.$$ (8)

At each step the error is squared. This is terrific if (and only if) you are close to $x^* = \frac{1}{2}$. Otherwise squaring a large error and multiplying by $-2$ is not good:

$x_0 = .70 \quad x_1 = .42 \quad x_2 = .487 \quad x_3 = .4997 \quad x_4 = .4999998$

$x_0 = 1.21 \quad x_1 = -.5 \quad x_2 = -1.5 \quad x_3 = -7.5 \quad x_4 = -127.5$

The algebra in Problem 18 confirms those experiments. There is fast convergence if $0 < x_0 < 1$. There is divergence if $x_0$ is negative or $x_0 > 1$. The tangent line goes to a negative $x_1$. After that Figure 3.22 shows a long trip backwards.

In the previous section we drew $F(x)$. The iteration $x_{n+1} = F(x_n)$ converged to the $45^\circ$ line, where $x^* = F(x^*)$. In this section we are drawing $f(x)$. Now $x^*$ is the point on the axis where $f(x^*) = 0$.

To repeat: It is $f(x^*) = 0$ that we aim for. But it is the slope $F'(x^*)$ that decides whether we get there. Example 2 has $F(x) = 2x - 2x^2$. The fixed points are $x^* = \frac{1}{2}$ (our solution) and $x^* = 0$ (not attractive). The slopes $F'(x^*)$ are zero (typical Newton) and 2 (typical repeller). The key to Newton's method is $F' = 0$ at the solution:

The slope of $F(x) = x - \frac{f(x)}{f'(x)}$ is $\frac{f(x)f''(x)}{(f'(x))^2}$. Then $F'(x) = 0$ when $f(x) = 0$. 

The examples $x^2 = b$ and $1/x = a$ show fast convergence or failure. In Chapter 13, and in reality, Newton's method solves much harder equations. Here I am going to choose a third example that came from pure curiosity about what might happen. The results are absolutely amazing. The equation is $x^2 = -1$.

**EXAMPLE 3** What happens to Newton's method if you ask it to solve $f(x) = x^2 + 1 = 0$?

The only solutions are the imaginary numbers $x^* = i$ and $x^* = -i$. There is no real square root of $-1$. Newton's method might as well give up. But it has no way to know that! The tangent line still crosses the axis at a new point $x_{n+1}$, even if the curve $y = x^2 + 1$ never crosses. Equation (5) still gives the iteration for $b = -1$:

$$x_{n+1} = \frac{1}{2} \left( x_n - \frac{1}{x_n} \right) = F(x_n).$$

(9)

The $x$'s cannot approach $i$ or $-i$ (nothing is imaginary). So what do they do?

The starting guess $x_0 = 1$ is interesting. It is followed by $x_1 = 0$. Then $x_2$ divides by zero and blows up. I expected other sequences to go to infinity. But the experiments showed something different (and mystifying). When $x_n$ is large, $x_{n+1}$ is less than half as large. After $x_n = 10$ comes $x_{n+1} = \frac{1}{2}(10 - \frac{1}{10}) = 4.95$. After much indecision and a long wait, a number near zero eventually appears. Then the next guess divides by that small number and goes far out again. This reminded me of “chaos.”

It is tempting to retreat to ordinary examples, where Newton's method is a big success. By trying exercises from the book or equations of your own, you will see that the fast convergence to $\sqrt{4}$ is very typical. The function can be much more complicated than $x^2 + 4$ (in practice it certainly is). The iteration for $2x = \cos x$ was in the previous section, and the error was squared at every step. If Newton's method starts close to $x^*$, its convergence is overwhelming. That has to be the main point of this section: Follow the tangent line.

Instead of those good functions, may I stay with this strange example $x^2 + 1 = 0$? It is not so predictable, and maybe not so important, but somehow it is more interesting. There is no real solution $x^*$, and Newton's method $x_{n+1} = \frac{1}{2}(x_n - 1/x_n)$ bounces around. We will now discover $x_n$.

**A FORMULA FOR $x_n$**

The key is an exercise from trigonometry books. Most of those problems just give practice with sines and cosines, but this one exactly fits $\frac{1}{2}(x_n - 1/x_n)$:

$$\frac{1}{2} \left( \cos \theta - \sin \theta \right) = \cos 2\theta \quad \text{or} \quad \frac{1}{2} \left( \cot \theta - \frac{1}{\cot \theta} \right) = \cot 2\theta$$

In the left equation, the common denominator is $2 \sin \theta \cos \theta$ (which is $\sin 2\theta$). The numerator is $\cos^2 \theta - \sin^2 \theta$ (which is $\cos 2\theta$). Replace cosine/sine by cotangent, and the identity says this:

If $x_0 = \cot \theta$ then $x_1 = \cot 2\theta$. Then $x_2 = \cot 4\theta$. Then $x_3 = \cot 8\theta$.

This is the formula. Our points are on the cotangent curve. Figure 3.23 starts from $x_0 = 2 = \cot \theta$, and every iteration doubles the angle.

**Example A** The sequence $x_0 = 1, x_1 = 0, x_2 = \infty$ matches the cotangents of $\pi/4, \pi/2$, and $\pi$. This sequence blows up because $x_2$ has a division by $x_1 = 0$. 


Example B The sequence \( \frac{1}{\sqrt{3}}, -\frac{1}{\sqrt{3}}, \frac{1}{\sqrt{3}} \) matches the cotangents of \( \pi/3, 2\pi/3, \) and \( 4\pi/3 \). This sequence cycles forever because \( x_0 = x_2 = x_4 = \ldots \).

Example C Start with a large \( x_0 \) (a small \( \theta \)). Then \( x_1 \) is about half as large (at \( 2\theta \)). Eventually one of the angles \( 4\theta, 8\theta, \ldots \) hits on a large cotangent, and the \( x \)'s go far out again. This is typical. Examples A and B were special, when \( \theta/\pi \) was \( \frac{1}{4} \) or \( \frac{1}{2} \).

What we have here is chaos. The \( x \)'s can't converge. They are strongly repelled by all points. They are also extremely sensitive to the value of \( \theta \). After ten steps \( \theta \) is multiplied by \( 2^{10} = 1024 \). The starting angles \( 60^\circ \) and \( 61^\circ \) look close, but now they are different by \( 1024^\circ \). If that were a multiple of \( 180^\circ \), the cotangents would still be close. In fact the \( x_{10} \)'s are 0.6 and 14.

This chaos in mathematics is also seen in nature. The most familiar example is the weather, which is much more delicate than you might think. The headline "Forecasting Pushed Too Far" appeared in Science (1989). The article said that the snowballing of small errors destroys the forecast after six days. We can't follow the weather equations for a month—the flight of a plane can change everything. This is a revolutionary idea, that a simple rule can lead to answers that are too sensitive to compute.

We are accustomed to complicated formulas (or no formulas). We are not accustomed to innocent-looking formulas like \( \cot 2^n \theta \), which are absolutely hopeless after 100 steps.

**CHAOS FROM A PARABOLA**

Now I get to tell you about new mathematics. First I will change the iteration \( x_{n+1} = \frac{1}{2}(x_n - 1/x_n) \) into one that is even simpler. By switching from \( x \) to \( z = 1/(1 + x^2) \), each new \( z \) turns out to involve only the old \( z \) and \( z^2 \):

\[
z_{n+1} = 4z_n - 4z_n^2.
\]

This is the most famous quadratic iteration in the world. There are books about it, and Problem 28 shows where it comes from. Our formula for \( x_n \) leads to \( z_n \):

\[
z_n = \frac{1}{1 + x_n^2} = \frac{1}{1 + (\cot 2^n \theta)^2} = (\sin 2^n \theta)^2.
\]
The sine is just as unpredictable as the cotangent, when $2^n \theta$ gets large. The new thing is to locate this quadratic as the last member (when $a = 4$) of the family

$$z_{n+1} = az_n - az_n^2, \quad 0 \leq a \leq 4.$$  \hfill (12)

Example 2 happened to be the middle member $a = 2$, converging to $\frac{1}{2}$. I would like to give a brief and very optional report on this iteration, for different $a$'s.

The general principle is to start with a number $z_0$ between 0 and 1, and compute $z_1, z_2, z_3, \ldots$. It is fascinating to watch the behavior change as $a$ increases. You can see it on your own computer. Here we describe some things to look for. All numbers stay between 0 and 1 and they may approach a limit. That happens when $a$ is small:

- for $0 \leq a \leq 1$ the $z_n$ approach $z^* = 0$
- for $1 \leq a \leq 3$ the $z_n$ approach $z^* = (a - 1)/a$

Those limit points are the solutions of $z = F(z)$. They are the fixed points where $z^* = az^* - a(z^*)^2$. But remember the test for approaching a limit: The slope at $z^*$ cannot be larger than one. Here $F = az - az^2$ has $F' = a - 2az$. It is easy to check $|F'| \leq 1$ at the limits predicted above. The hard problem—sometimes impossible—is to predict what happens above $a = 3$. Our case is $a = 4$.

The $z$'s cannot approach a limit when $|F'(z^*)| > 1$. Something has to happen, and there are at least three possibilities:

- The $z_n$'s can cycle or fill the whole interval $(0, 1)$ or approach a Cantor set.

I start with a random number $z_0$, take 100 steps, and write down steps 101 to 105:

<table>
<thead>
<tr>
<th>$a$</th>
<th>$z_{101}$</th>
<th>$z_{102}$</th>
<th>$z_{103}$</th>
<th>$z_{104}$</th>
<th>$z_{105}$</th>
</tr>
</thead>
<tbody>
<tr>
<td>3.4</td>
<td>.842</td>
<td>.452</td>
<td>.842</td>
<td>.452</td>
<td>.842</td>
</tr>
<tr>
<td>3.5</td>
<td>.875</td>
<td>.383</td>
<td>.827</td>
<td>.501</td>
<td>.875</td>
</tr>
<tr>
<td>3.8</td>
<td>.336</td>
<td>.848</td>
<td>.491</td>
<td>.950</td>
<td>.182</td>
</tr>
<tr>
<td>4.0</td>
<td>.169</td>
<td>.562</td>
<td>.985</td>
<td>.060</td>
<td>.225</td>
</tr>
</tbody>
</table>

The first column is converging to a "2-cycle." It alternates between $x = .842$ and $y = .452$. Those satisfy $y = F(x)$ and $x = F(y) = F(F(x))$. If we look at a double step when $a = 3.4$, $x$ and $y$ are fixed points of the double iteration $z_{n+2} = F(F(z_n))$. When $a$ increases past 3.45, this cycle becomes unstable.

At that point the period doubles from 2 to 4. With $a = 3.5$ you see a "4-cycle" in the table—it repeats after four steps. The sequence bounces from .875 to .383 to .827 to .501 and back to .875. This cycle must be attractive or we would not see it. But it also becomes unstable as $a$ increases. Next comes an 8-cycle, which is stable in a little window (you could compute it) around $a = 3.55$. The cycles are stable for shorter and shorter intervals of $a$'s. Those stability windows are reduced by the Feigenbaum shrinking factor 4.66922... Cycles of length 16 and 32 and 64 can be seen in physical experiments, but they are all unstable before $a = 3.57$. What happens then?

The new and unexpected behavior is between 3.57 and 4. Down each line of Figure 3.24, the computer has plotted the values of $z_{1001}$ to $z_{2000}$—omitting the first thousand points to let a stable period (or chaos) become established. No points appeared in the big white wedge. I don't know why. In the window for period 3, you
3.7 Newton's Method (and Chaos)

The period 2, 4, ..., is the number of z's in a cycle.

Fig. 3.24 Period doubling and chaos from iterating $F(z)$ (stolen by special permission from Introduction to Applied Mathematics by Gilbert Strang, Wellesley-Cambridge Press).

I can't tell what happens at $a = 3.8$. There may be a stable cycle of some long period. The z's may come close to every point between 0 and 1. A third possibility is to approach a very thin limit set, which looks like the famous Cantor set:

To construct the Cantor set, divide [0, 1] into three pieces and remove the open interval $(\frac{1}{3}, \frac{2}{3})$. Then remove $(\frac{1}{3}, \frac{2}{3})$ and $(\frac{2}{3}, \frac{1}{3})$ from what remains. At each step take out the middle thirds. The points that are left form the Cantor set.

All the endpoints $\frac{1}{3}, \frac{2}{3}, \frac{1}{3}, \frac{2}{3}, ...$ are in the set. So is $\frac{1}{3}$ (Problem 42). Nevertheless the lengths of the removed intervals add to 1 and the Cantor set has "measure zero."

What is especially striking is its self-similarity: Between 0 and $\frac{1}{3}$ you see the same Cantor set three times smaller. From 0 to $\frac{1}{3}$ the Cantor set is there again, scaled down by 3. Every section, when blown up, copies the larger picture.

Fractals That self-similarity is typical of a fractal. There is an infinite sequence of scales. A mathematical snowflake starts with a triangle and adds a bump in the middle of each side. At every step the bumps lengthen the sides by $4/3$. The final boundary is self-similar, like an infinitely long coastline.

The word "fractal" comes from fractional dimension. The snowflake boundary has dimension larger than 1 and smaller than 2. The Cantor set has dimension larger than 0 and smaller than 1. Covering an ordinary line segment with circles of radius $r$ would take $c/r$ circles. For fractals it takes $c/r^D$ circles—and $D$ is the dimension.
Fig. 3.25 Cantor set (middle thirds removed). Fractal snowflake (infinite boundary).

Our iteration $z_{n+1} = 4z_n - 4z_n^2$ has $a = 4$, at the end of Figure 3.24. The sequence $z_0, z_1, ...$ goes everywhere and nowhere. Its behavior is chaotic, and statistical tests find no pattern. For all practical purposes the numbers are random.

Think what this means in an experiment (or the stock market). If simple rules produce chaos, there is absolutely no way to predict the results. No measurement can ever be sufficiently accurate. The newspapers report that Pluto's orbit is chaotic—even though it obeys the law of gravity. The motion is totally unpredictable over long times. I don't know what that does for astronomy (or astrology).

The most readable book on this subject is Gleick's best-seller *Chaos: Making a New Science*. The most dazzling books are *The Beauty of Fractals* and *The Science of Fractal Images*, in which Peitgen and Richter and Saupe show photographs that have been in art museums around the world. The most original books are Mandelbrot's *Fractals* and *Fractal Geometry*. Our cover has a fractal from Figure 13.11.

We return to friendlier problems in which calculus is not helpless.

**NEWTON'S METHOD VS. SECANT METHOD: CALCULATOR PROGRAMS**

The hard part of Newton's method is to find $df/dx$. We need it for the slope of the tangent line. But calculus can approximate by $\Delta f/\Delta x$—using the values of $f(x)$ already computed at $x_n$ and $x_{n-1}$.

The secant method follows the secant line instead of the tangent line:

$$\text{Secant: } x_{n+1} = x_n - \frac{f(x_n)}{(\Delta f/\Delta x)_n} \text{ where } \left(\frac{\Delta f}{\Delta x}\right)_n = \frac{f(x_n) - f(x_{n-1})}{x_n - x_{n-1}}. \quad (13)$$

The secant line connects the two latest points on the graph of $f(x)$. Its equation is $y - f(x_n) = (\Delta f/\Delta x)(x - x_n)$. Set $y = 0$ to find equation (13) for the new $x = x_{n+1}$, where the line crosses the axis.

Prediction: Three secant steps are about as good as two Newton steps. Both should give four times as many correct decimals: $(\text{error}) \rightarrow (\text{error})^4$. Probably the secant method is also chaotic for $x^2 + 1 = 0$.

These Newton and secant programs are for the TI-81. Place the formula for $f(x)$ in slot $Y_1$ and the formula for $f'(x)$ in slot $Y_2$ on the $Y = $ function edit screen. Answer the prompt with the initial $x_0 = X0$. The programs pause to display each approximation $x_n$, the value $f(x_n)$, and the difference $x_n - x_{n-1}$. Press ENTER to continue or press ON and select item 2 : QUIT to break. If $f(x_n) = 0$, the programs display ROOT AT and the root $x_n$. 
3.7 Newton's Method (and Chaos)

Read-through questions

When \( f(x) = 0 \) is linearized to \( f(x_n) + f'(x_n)(x - x_n) = 0 \), the solution \( x = \cdots \) is Newton's \( x_{n+1} \). The \( \cdots \) to the curve crosses the axis at \( x_{n+1} \), while the \( \cdots \) crosses at \( x^* \). The errors at \( x_n \) and \( x_{n+1} \) are normally related by \( (\text{error})_{n+1} \approx M_{\cdots} \). This is \( \cdots \) convergence. The number of correct decimals \( \cdots \) at every step.

For \( f(x) = x^2 - b \), Newton's iteration is \( x_{n+1} = \cdots \). The \( \cdots \) to the curve if \( x_0 > 0 \) and to \( \cdots \) if \( x_0 < 0 \). For \( f(x) = x^3 + 1 \), the iteration becomes \( x_{n+1} = \cdots \). This cannot converge to \( \cdots \). Instead it leads to chaos. Changing to \( z = 1/(x^2 + 1) \) yields the parabolic iteration \( x_{n+1} = \cdots \).

For \( a < 3 \), \( z_{n+1} = az_n - az_n^2 \) converges to a single \( \cdots \). After \( a = 3 \) the limit is a 2-cycle, which means \( \cdots \). Later the limit is a Cantor set, which is a one-dimensional example of a \( \cdots \). The Cantor set is self- \( \cdots \).

1. To solve \( f(x) = x^3 - b = 0 \), what iteration comes from Newton's method?

2. For \( f(x) = (x - 1)/(x + 1) \), Newton's formula is \( x_{n+1} = F(x_n) = \cdots \). Solve \( x^* = F(x^*) \) and find \( F'(x^*) \). What limit do the \( x_n \)'s approach?

3. I believe that Newton only applied his method in public to one equation \( x^3 - 2x - 5 = 0 \). Raphson carried the idea forward but got partial credit at best. After two steps from \( x_0 = 2 \), how many decimals in \( x^* = 2.09455148 \) are correct?

4. Show that Newton's method for \( f(x) = x^{1/3} \) gives the strange formula \( x_{n+1} = -2x_n \). Draw a graph to show the iterations.

5. Find \( x_1 \) if (a) \( f(x_0) = 0 \); (b) \( f'(x_0) = 0 \).

6. Graph \( f(x) = x^3 - 3x - 1 \) and estimate its roots \( x^* \). Run Newton's method starting from 0, 1, \( \frac{1}{2} \), and 1.1. Experiment to decide which \( x_0 \) converge to which root.

7. Solve \( x^2 - 6x + 5 = 0 \) by Newton's method with \( x_0 = 2.5 \) and 3. Draw a graph to show which \( x_0 \) lead to which root.

8. If \( f(x) \) is increasing and concave up (\( f'' > 0 \) and \( f''' > 0 \)) show by a graph that Newton's method converges. From which side?

Solve 9–17 to four decimal places by Newton's method with a computer or calculator. Choose any \( x_0 \) except \( x^* \).

9. \( x^2 - 10 = 0 \)

10. \( x^4 - 100 = 0 \) (faster or slower than Problem 9?)

11. \( x^3 - x = 0 \) (which \( x_0 \) to which root?)

12. \( x^3 - x = 0 \) (which \( x_0 \) to which root?)

13. \( x + 5 \cos x = 0 \) (this has three roots?)

14. \( x + x + x \) (find two roots) (are there more?)

15. \( 1/(1 - x) = 2 \)

16. \( 1 + x + x^2 + x^3 + x^4 = 2 \)

17. \( x^3 + (x + 1)^3 = 10^3 \)

18. (a) Show that \( x_{n+1} = 2x_n - 2x_n^2 \) in Example 2 is the same as \( (1 - 2x_{n+1}) = (1 - 2x_n)^2 \).

(b) Prove divergence if \( |1 - 2x_0| > 1 \). Prove convergence if \( |1 - 2x_0| < 1 \) or \( x_0 \) to which root?

19. With \( a = 3 \) in Example 2, experiment with the Newton iteration \( x_{n+1} = 2x_n - 3x_n^2 \) to decide which \( x_0 \) lead to \( x^* \).

20. Rewrite \( x_{n+1} = 2x_n - ax_n^2 \) as \( (1 - ax_{n+1}) = (1 - ax_n)^2 \). For which \( x_0 \) do the sequence \( 1 - ax_n \) approach zero (so \( x_n \) to \( 1/a \))?

21. What is Newton's method to find the \( k \)th root of 7? Calculate \( \sqrt[3]{7} \) to 7 places.

22. Find all solutions of \( x^3 = 4x - 1 \) (5 decimals).
Problems 23–29 are about $x^2 + 1 = 0$ and chaos.
23 For $\theta = \pi/16$ when does $x_n = \cot 2^n \theta$ blow up? For $\theta = \pi/7$ when does $\cot 2^n \theta = \cot \theta$? (The angles $2^n \theta$ and $\theta$ differ by a multiple of $\pi$.)
24 For $\theta = \pi/9$ follow the sequence until $x_n = x_0$.
25 For $\theta = 1$, $x_n$ never returns to $x_0 = \cot 1$. The angles $2^k \theta$ and $1$ never differ by a multiple of $\pi$ because ________.
26 If $x_0$ equals $\sin\theta$, show that $z_1 = 4x_0 - 4x_0^2$ equals $\sin^2 \theta$.
27 If $y = x^2 + 1$, each new $y$ is
$$y_{n+1} = y_n^2 + 1 = \frac{1}{4} \left(x_n - \frac{1}{x_n}\right)^2 + 1.$$Show that this equals $y_n^2/4(y_n - 1)$.
28 Turn Problem 27 upside down, $1/y_{n+1} = 4(y_n - 1)/y_n^2$, to find the quadratic iteration (10) for $x_n = 1/y_n = 1/(1 + x_n^2)$.
29 If $F(x) = 4x - 4x^2$ what is $F(F(x))$? How many solutions to $z = F(F(z))$? How many are not solutions to $z = F(z)$?
30 Apply Newton's method to $x^3 - .64x - .36 = 0$ to find the basin of attraction for $x^* = 1$. Also find a pair of points for which $y = F(z)$ and $z = F(y)$. In this example Newton does not always find a root.
31 Newton's method solves $x/(1-x) = 0$ by $x_{n+1} = ____$. From which $x_0$ does it converge? The distance to $x^* = 0$ is exactly squared.

Problems 33–41 are about competitors of Newton.
32 At a double root, Newton only converges linearly. What is the iteration to solve $x^2 = 0$?
33 To speed up Newton's method, find the step $\Delta x$ from $f(x_0) + \Delta x f'(x_0) + \frac{1}{2}(\Delta x)^2 f''(x_0) = 0$. Test on $f(x) = x^3 - 1$ from $x_0 = 0$ and explain.
34 Halley's method uses $f_0 + \Delta x f'_0 + \frac{1}{2}(\Delta x)^2 f''_0 + \frac{1}{6}(\Delta x)^3 f'''_0 = 0$. For $f(x) = x^3 - 1$ and $x_0 = 1 + \epsilon$, show that $x_1 = 1 + O(\epsilon^3)$—which is cubic convergence.
35 Apply the secant method to $f(x) = x^3 - 4 = 0$, starting from $x_0 = 1$ and $x_1 = 2.5$. Find $\Delta f/\Delta x$ and the next point $x_2$ by hand. Newton uses $f'(x_1) = 5$ to reach $x_2 = 2.05$. Which is closer to $x^* = 2$?
36 Draw a graph of $f(x) = x^3 - 4$ to show the secant line in Problem 35 and the point $x_2$ where it crosses the axis.

3.8 The Mean Value Theorem and l'Hôpital's Rule

Now comes one of the cornerstones of calculus: the Mean Value Theorem. It connects the local picture (slope at a point) to the global picture (average slope across an interval). In other words it relates $df/dx$ to $\Delta f/\Delta x$. Calculus depends on this connec-
3.8 The Mean Value Theorem and l'Hôpital's Rule

The Mean Value Theorem and l'Hôpital's Rule

Fig. 3.26 (a) $v$ jumps over $v_{\text{average}}$. (b) $v$ equals $v_{\text{average}}$.

...(text continues)
3 Applications of the Derivative

Small problem: The maximum could be reached at the ends $a$ and $b$, if $f(x) < 0$ in between. At those endpoints $df/dx$ might not be zero. But in that case the minimum is reached at an interior point $c$, which is equally acceptable. The key to our proof is that a continuous function on $[a, b]$ reaches its maximum and minimum. This is the Extreme Value Theorem.†

It is ironic that Rolle himself did not believe the logic behind calculus. He may not have believed his own theorem! Probably he didn’t know what it meant—the language of “evanescent quantities” (Newton) and “infinitesimals” (Leibniz) was exciting but frustrating. Limits were close but never reached. Curves had infinitely many flat sides. Rolle didn’t accept that reasoning, and what was really serious, he didn’t accept the conclusions. The Académie des Sciences had to stop his battles (he fought against ordinary mathematicians, not Newton and Leibniz). So he went back to number theory, but his special case when $f(a) = f(b) = 0$ leads directly to the big one.

Fig. 3.27 Rolle’s theorem is when $f(a) = f(b) = 0$ in the Mean Value Theorem.

Proof of the Mean Value Theorem We are looking for a point where $df/dx$ equals $\Delta f/\Delta x$. The idea is to tilt the graph back to Rolle’s special case (when $\Delta f$ was zero). In Figure 3.27b, the distance $F(x)$ between the curve and the dotted secant line comes from subtraction:

$$F(x) = f(x) - \left[ f(a) + \frac{\Delta f}{\Delta x} (x - a) \right].$$

At $a$ and $b$, this distance is $F(a) = F(b) = 0$. Rolle’s theorem applies to $F(x)$. There is an interior point where $F'(c) = 0$. At that point take the derivative of equation (2):

$$0 = F'(c) = f'(c) - \left( \frac{\Delta f}{\Delta x} \right).$$

The desired point $c$ is found, proving the theorem.

EXAMPLE 1 The function $f(x) = \sqrt{x}$ goes from zero at $x = 0$ to ten at $x = 100$. Its average slope is $\Delta f/\Delta x = 10/100$. The derivative $f''(x) = 1/2 \sqrt{x}$ exists in the open interval $(0, 100)$, even though it blows up at the end $x = 0$. By the Mean Value Theorem there must be a point where $10/100 = f''(c) = 1/2 \sqrt{c}$. That point is $c = 25$.

The truth is that nobody cares about the exact value of $c$. Its existence is what matters. Notice how it affects the linear approximation $f(x) \approx f(a) + f'(a)(x - a)$, which was basic to this chapter. Close becomes exact (≈ becomes =) when $f'$ is computed at $c$ instead of $a$:

†If $f(x)$ doesn’t reach its maximum $M$, then $1/(M - f(x))$ would be continuous but also approach infinity. Essential fact: A continuous function on $[a, b]$ cannot approach infinity.
3.8 The Mean Value Theorem and l'Hôpital's Rule

**EXAMPLE 2** The function \( f(x) = \sin x \) starts from \( f(0) = 0 \). The linear prediction (tangent line) uses the slope \( \cos 0 = 1 \). The exact prediction uses the slope \( \cos c \) at an unknown point between 0 and \( x \):

\[
(\text{approximate}) \sin x \approx x \quad \quad (\text{exact}) \sin x = (\cos c)x.
\]  
(4)

The approximation is useful, because everything is computed at \( x = a = 0 \). The exact formula is interesting, because \( \cos c \leq 1 \) proves again that \( \sin x \leq x \). The slope is below 1, so the sine graph stays below the \( 45^\circ \) line.

**EXAMPLE 3** If \( f'(c) = 0 \) at all points in an interval then \( f(x) \) is constant.

Proof When \( f' \) is everywhere zero, the theorem gives \( \Delta f = 0 \). Every pair of points has \( f(b) = f(a) \). The graph is a horizontal line. That deceptively simple case is a key to the Fundamental Theorem of Calculus.

Most applications of \( \Delta f = f'(c)\Delta x \) do not end up with a number. They end up with another theorem (like this one). The goal is to connect derivatives (local) to differences (global). But the next application—l'Hôpital's Rule—manages to produce a number out of \( 0/0 \).

**L'HÔPITAL'S RULE**

When \( f(x) \) and \( g(x) \) both approach zero, what happens to their ratio \( f(x)/g(x) \)?

\[
\frac{f(x)}{g(x)} = \frac{x^2}{x} \quad \text{or} \quad \frac{\sin x}{x} \quad \text{or} \quad \frac{x - \sin x}{1 - \cos x} \quad \text{all become} \quad \frac{0}{0} \quad \text{at} \quad x = 0.
\]

Since \( 0/0 \) is meaningless, we cannot work separately with \( f(x) \) and \( g(x) \). This is a "race toward zero," in which two functions become small while their ratio might do anything. The problem is to find the limit of \( f(x)/g(x) \).

One such limit is already studied. It is the derivative! \( \Delta f/\Delta x \) automatically builds in a race toward zero, whose limit is \( df/dx \):

\[
\lim_{x \to a} \frac{f(x) - f(a)}{x - a} = f'(a). \quad (5)
\]

The idea of l'Hôpital is to use \( f'/g' \) to handle \( f/g \). The derivative is the special case \( g(x) = x - a \), with \( g' = 1 \). The Rule is followed by examples and proofs.

**3P l'Hôpital's Rule** Suppose \( f(x) \) and \( g(x) \) both approach zero as \( x \to a \). Then \( f(x)/g(x) \) approaches the same limit as \( f'(x)/g'(x) \), if that second limit exists:

\[
\lim_{x \to a} \frac{f(x)}{g(x)} = \lim_{x \to a} \frac{f'(x)}{g'(x)}. \quad \text{Normally this limit is} \quad \frac{f'(a)}{g'(a)}. \quad (6)
\]

This is not the quotient rule! The derivatives of \( f(x) \) and \( g(x) \) are taken separately. Geometrically, l'Hôpital is saying that when functions go to zero their slopes control their size. An easy case is \( f = 6(x - a) \) and \( g = 2(x - a) \). The ratio \( f/g \) is exactly 6/2,
the ratio of their slopes. Figure 3.28 shows these straight lines dropping to zero, controlled by $b$ and 2.

The next figure shows the same limit $6/2$, when the curves are tangent to the lines. That picture is the key to l'Hôpital's rule.

Generally the limit of $f/g$ can be a finite number $L$ or $+\infty$ or $-\infty$. (Also the limit point $x = a$ can represent a finite number or $+\infty$ or $-\infty$. We keep it finite.) The one absolute requirement is that $f(x)$ and $g(x)$ must separately approach zero—we insist on $0/0$. Otherwise there is no reason why equation (6) should be true. With $f(x) = x$ and $g(x) = x - 1$, don't use l'Hôpital:

$$\frac{f(x)}{g(x)} \rightarrow \frac{a}{a-1} \quad \text{but} \quad \frac{f'(x)}{g'(x)} = \frac{1}{1}.$$

Ordinary ratios approach $\lim f(x)$ divided by $\lim g(x)$. l'Hôpital enters only for $0/0$.

**EXAMPLE 4** (an old friend) $\lim_{x \to 0} \frac{1 - \cos x}{x}$ equals $\lim_{x \to 0} \frac{\sin x}{1}$. This equals zero.

**EXAMPLE 5**

If $f = \tan x$ leads to $f' = \frac{\sec^2 x}{\cos x}$. At $x = 0$ the limit is $\frac{1}{1}$.

**EXAMPLE 6**

If $g = \frac{x - \sin x}{1 - \cos x}$ leads to $f' = \frac{1 - \cos x}{\sin x}$. At $x = 0$ this is still $0/0$.

Solution Apply the Rule to $f'/g'$. It has the same limit as $f''/g''$:

$$\text{if } \frac{f}{g} \to \frac{0}{0} \quad \text{and} \quad \frac{f'}{g'} \to \frac{0}{0} \quad \text{then compute} \quad \frac{f''(x)}{g''(x)} = \frac{\sin x}{\cos x} \to \frac{0}{1} = 0.$$

The reason behind l'Hôpital's Rule is that the following fractions are the same:

$$\frac{f(x)}{g(x)} = \frac{f(x) - f(a)}{g(x) - g(a)} \quad \frac{f'(x) - f'(a)}{g'(x) - g'(a)} \quad \frac{x - a}{x - a} \quad (7)$$

That is just algebra; the limit hasn't happened yet. The factors $x - a$ cancel, and the numbers $f(a)$ and $g(a)$ are zero by assumption. Now take the limit on the right side of (7) as $x$ approaches $a$.

What normally happens is that one part approaches $f'$ at $x = a$. The other part approaches $g'(a)$. We hope $g'(a)$ is not zero. In this case we can divide one limit by
the other limit. That gives the "normal" answer
\[ \lim_{x \to a} \frac{f(x)}{g(x)} = \text{limit of } f'(a) g'(a). \] (8)

This is also l'Hôpital's answer. When \( f'(x) \to f'(a) \) and separately \( g'(x) \to g'(a) \), his overall limit is \( f'(a)/g'(a) \). He published this rule in the first textbook ever written on differential calculus. (That was in 1696—the limit was actually discovered by his teacher Bernoulli.) Three hundred years later we apply his name to other cases permitted in (6), when \( f'/g' \) might approach a limit even if the separate parts do not.

To prove this more general form of l'Hôpital's Rule, we need a more general Mean Value Theorem. I regard the discussion below as optional in a calculus course (but required in a calculus book). The important idea already came in equation (8).

**Remark** The basic "indeterminate" is \( \infty - \infty \). If \( f(x) \) and \( g(x) \) approach infinity, anything is possible for \( f(x) - g(x) \). We could have \( x^2 - x \) or \( x - x^2 \) or \( (x + 2) - x \). Their limits are \( \infty \) and \( -\infty \) and 2.

At the next level are \( 0/0 \) and \( \infty/\infty \) and \( 0 \cdot \infty \). To find the limit in these cases, try l'Hôpital's Rule. See Problem 24 when \( f(x)/g(x) \) approaches \( \infty/\infty \). When \( f(x) \to 0 \) and \( g(x) \to \infty \), apply the 0/0 rule to \( f(x)/(1/g(x)) \).

The next level has \( 0^0 \) and \( 1^\infty \) and \( \infty^0 \). Those come from limits of \( f(x)^{g(x)} \). If \( f(x) \) approaches 0, 1, or \( \infty \) while \( g(x) \) approaches 0, \( \infty \), or 0, we need more information. A really curious example is \( x^{1/\ln x} \), which shows all three possibilities \( 0^0 \) and \( 1^\infty \) and \( \infty^0 \). This function is actually a constant! It equals \( e \).

To go back down a level, take logarithms. Then \( g(x) \ln f(x) \) returns to \( 0/0 \) and \( 0 - \infty \) and l'Hôpital's Rule. But logarithms and \( e \) have to wait for Chapter 6.

THE GENERALIZED MEAN VALUE THEOREM

The MVT can be extended to two functions. The extension is due to Cauchy, who cleared up the whole idea of limits. You will recognize the special case \( g = x \) as the ordinary Mean Value Theorem.

**3Q Generalized MVT** If \( f(x) \) and \( g(x) \) are continuous on \([a, b]\) and differentiable on \((a, b)\), there is a point \( a < c < b \) where

\[ [f(b) - f(a)]g'(c) = [g(b) - g(a)]f'(c). \] (9)

The proof comes by constructing a new function that has \( F(a) = F(b) \):

\[ F(x) = [f(b) - f(a)]g(x) - [g(b) - g(a)]f(x). \]

The ordinary Mean Value Theorem leads to \( F'(c) = 0 \)—which is equation (9).

**Application 1** (Proof of l'Hôpital's Rule) The rule deals with \( f(a)/g(a) = 0/0 \). Inserting those zeros into equation (9) leaves \( f(b)g'(c) = g(b)f'(c) \). Therefore

\[ \frac{f(b)}{g(b)} = \frac{f'(c)}{g'(c)}. \] (10)

As \( b \) approaches \( a \), so does \( c \). The point \( c \) is squeezed between \( a \) and \( b \). The limit of equation (10) as \( b \to a \) and \( c \to a \) is l'Hôpital's Rule.
3 Applications of the Derivative

Application 2 (Error in linear approximation) Section 3.2 stated that the distance between a curve and its tangent line grows like $(x - a)^2$. Now we can prove this, and find out more. Linear approximation is

$$ f(x) = f(a) + f'(a)(x - a) + 	ext{error } e(x). $$

The pattern suggests an error involving $f''(x)$ and $(x - a)^2$. The key example $f = x^2$ shows the need for a factor $\frac{1}{2}$ (to cancel $f'' = 2$). The error in linear approximation is

$$ e(x) = \frac{1}{2} f''(c)(x - a)^2 \quad \text{with} \quad a < c < x. $$

Key idea Compare the error $e(x)$ to $(x - a)^2$. Both are zero at $x = a$:

$$ e = f(x) - f(a) - f'(a)(x - a) \quad e' = f'(x) - f'(a) \quad e'' = f''(x) $$

$$ g = (x - a)^2 \quad g' = 2(x - a) \quad g'' = 2 $$

The Generalized Mean Value Theorem finds a point $C$ between $a$ and $x$ where $e(x)/g(x) = e'(C)/g'(C)$. This is equation (10) with different letters. After checking $e'(a) = g'(a) = 0$, apply the same theorem to $e'(x)$ and $g'(x)$. It produces a point $c$ between $a$ and $x$—certainly between $a$ and $x$—where

$$ \frac{e'(C)}{g'(C)} = \frac{e'(c)}{g'(c)} $$

With $g = (x - a)^2$ and $g'' = 2$ and $e'' = f''$, the equation on the right is $e(x) = \frac{1}{2} f''(c)(x - a)^2$. The error formula is proved. A very good approximation is $\frac{1}{2} f''(a)(x - a)^2$.

**Example 7**

$$ f(x) = \sqrt{x} \quad \text{near } a = 100: \quad \sqrt{102} \approx 10 + \left( \frac{1}{20} \right) 2 + \frac{1}{2} \left( \frac{1}{4000} \right) 2^2. $$

That last term predicts $e = -.0005$. The actual error is $\sqrt{102} - 10.1 = - .000456$.

### 3.8 Exercises

Read-through questions

The Mean Value Theorem equates the average slope $\Delta f/\Delta x$ over an interval $[a, b]$ to the slope $df/dx$ at an unknown point $c$ in the interval $[a, b]$. The statement is $f(x) = f(a) + f'(c)(x - a)$. The key idea is $f'(a) = f'(b) = 0$, and the point $c$ satisfies $f'(c)$ and $f''(c)$.

Consequences of the Mean Value Theorem include:

- If $f'(x) = 0$ everywhere in an interval then $f(x) = f(a) + f'(a)(x - a)$ is exact for some $c$ between $a$ and $x$. The quadratic prediction $f(x) = f(a) + f'(a)(x - a) + \frac{1}{2} f''(c)(x - a)^2$ is exact for another $c$. The error in $f(x) = f(a) + f'(a)(x - a)$ is less than $\frac{1}{2} M(x - a)^2$ where $M$ is the maximum of $|f''|$. A chief consequence is l'Hôpital's Rule, which applies when $f(x)$ and $g(x) \to \infty$ as $x \to a$. In that case the limit of $f(x)/g(x)$ equals the limit of $f'(a)/g'(a)$, provided this limit exists. Normally this limit is $f'(a)/g'(a)$. If this is also 0/0, go on to the limit of $\frac{1}{x}$.

Find all points $0 < c < 2$ where $f(2) - f(0) = f'(c)(2 - 0)$.

1. $f(x) = x^3$
2. $f(x) = \sin \pi x$
3. $f(x) = \tan 2x$
4. $f(x) = 1 + x + x^2$
5. $f(x) = (x - 1)^{10}$
6. $f(x) = (x - 1)^9$

In 7–10 show that no point $c$ yields $f(1) - f(-1) = f'(c)(2)$. Explain why the Mean Value Theorem fails to apply.

7. $f(x) = |x - \frac{1}{2}|$
8. $f(x) = \text{unit step function}$
9. $f(x) = |x|^{1/2}$
10. $f(x) = 1/x^2$

11. Show that $\sec^2 x$ and $\tan^2 x$ have the same derivative, and draw a conclusion about $f(x) = \sec^2 x - \tan^2 x$.

12. Show that $\csc^2 x$ and $\cot^2 x$ have the same derivative and find $f(x) = \csc^2 x - \cot^2 x$.

Evaluate the limits in 13–22 by l'Hôpital's Rule.

13. $\lim_{x \to 3} \frac{x^2 - 9}{x - 3}$
14. $\lim_{x \to 3} \frac{x^2 - 9}{x + 3}$
15 \lim_{x \to 0} \frac{(1 + x)^2 - 1}{x} \\
16 \lim_{x \to 0} \frac{\sqrt{1 - \cos x}}{x} \\
17 \lim_{x \to \pi} \frac{x - \pi}{\sin x} \\
18 \lim_{x \to 1} \frac{x - 1}{\sin x} \\
19 \lim_{x \to 0} \frac{(1 + x)^2 - 1}{x} \\
20 \lim_{x \to 0} \frac{(1 + x)^2 - 1 - nx}{x^2} \\
21 \lim_{x \to 0} \frac{\sin x - \tan x}{x^3} \\
22 \lim_{x \to 0} \frac{\sqrt{1 + x} - \sqrt{1 - x}}{x} \\
23 \text{For } f = x^2 - 4 \text{ and } g = x + 2, \text{ the ratio } f'/g' \text{ approaches 4 as } x \to 2. \text{ What is the limit of } f(x)/g(x)? \text{ What goes wrong in l'Hôpital's Rule?} \\
24 \text{l'Hôpital's Rule still holds for } f(x)/g(x) \to \infty/\infty: L \text{ is} \\
\lim_{x \to 0} \frac{f(x)}{g(x)} = \lim_{x \to 0} \frac{f'(x)}{g'(x)} = \lim_{x \to 0} \frac{g'(x)/g(x)}{f'(x)/f(x)} = L^2 \lim_{x \to 0} \frac{g'(x)}{g(x)} \cdot \frac{f(x)}{f'(x)} \\
\text{Then } L \text{ equals } \lim \left[ \frac{f'(x)}{g'(x)} \right] \text{ if this limit exists. Where did we use the rule for 0/0? What other limit rule was used?} \\
25 \text{Compute } \lim_{x \to 0} \frac{1 + (1/x)}{1 - (1/x)}. \\
26 \text{Compute } \lim_{x \to \infty} \frac{x^2 + x}{2x^2}. \\
27 \text{Compute } \lim_{x \to \infty} \frac{x + \cos x}{x + \sin x} \text{ by common sense. Show that l'Hôpital gives no answer.} \\
28 \text{Compute } \lim_{x \to 0} \frac{\csc x}{\cot x} \text{ by common sense or trickery.} \\
29 \text{The Mean Value Theorem applied to } f(x) = x^3 \text{ guarantees that some number } c \text{ between 1 and 4 has a certain property. Say what the property is and find } c. \\
30 \text{If } |df/dx| \leq 1 \text{ at all points, prove this fact:} \\
|f(x) - f(y)| \leq |x - y| \text{ at all } x \text{ and } y. \\
31 \text{The error in Newton's method is squared at each step: } \\
|x_{n+1} - x^*| \leq M|x_n - x^*|^2. \text{ The proof starts from } 0 = f(x^*) = f(x_n) + f'(x_n)(x^* - x_n) + \frac{1}{2}f''(c)(x^* - x_n)^2. \text{ Divide by } f'(x_n), \text{ recognize } x_{n+1}, \text{ and estimate } M. \\
32 \text{(Rolle's theorem backward) Suppose } f'(c) = 0. \text{ Are there necessarily two points around } c \text{ where } f(a) = f(b)? \\
33 \text{Suppose } f(0) = 0. \text{ If } f(x)/x \text{ has a limit as } x \to 0, \text{ that limit is better known to us as } \lim_{x \to 0} f(x)/x. \text{ L'Hôpital's Rule looks instead at the limit of } \lim_{x \to 0} \frac{f'(x)}{g'(x)}. \\
\text{Conclusion from l'Hôpital: The limit of } f'(x), \text{ if it exists, agrees with } f'(0). \text{ Thus } f'(x) \text{ cannot have a "removable } \lim_{x \to 0} f(x)/x. \\
34 \text{It is possible that } f'(x)/g'(x) \text{ has no limit but } f(x)/g(x) \to L. \text{ This is why l'Hôpital included an "if."} \\
(a) \text{Find } L \text{ as } x \to 0 \text{ when } f(x) = x^2 \cos(1/x) \text{ and } g(x) = x. \text{ Remember that cosines are below 1.} \\
(b) \text{From the formula } f'(x) = \sin(1/x) + 2x \cos(1/x) \text{ show that } f'/g' \text{ has no limit as } x \to 0. \\
35 \text{Stein's calculus book asks for the limiting ratio of } f(x) = \text{triangular area } ABC \text{ to } g(x) = \text{curved area } ABC. \\
(a) \text{Guess the limit of } f/g \text{ as the angle } x \text{ goes to zero.} \\
(b) \text{Explain why } f(x) = \frac{1}{2}(\sin x - \sin x \cos x) \text{ and } g(x) = \frac{1}{2}(x - \sin x \cos x). \text{ (c) Compute the true limit of } f(x)/g(x). \\
36 \text{If you drive 3000 miles from New York to L.A. in 100 hours (sleeping and eating and going backwards are allowed) then at some moment your speed is } \infty. \\
37 \text{As } x \to \infty \text{ l'Hôpital's Rule still applies. The limit of } f(x)/g(x) \text{ equals the limit of } f'(x)/g'(x), \text{ if that limit exists. What is the limit as the graphs become parallel in Figure B?} \\
38 \text{Prove that } f(x) \text{ is increasing when its slope is positive: } \text{If } f'(c) > 0 \text{ at all points } c, \text{ then } f(b) > f(a) \text{ at all pairs of points } b > a.
# Contents

## CHAPTER 4  The Chain Rule
- 4.1 Derivatives by the Chain Rule 154
- 4.2 Implicit Differentiation and Related Rates 160
- 4.3 Inverse Functions and Their Derivatives 164
- 4.4 Inverses of Trigonometric Functions 171

## CHAPTER 5  Integrals
- 5.1 The Idea of the Integral 177
- 5.2 Antiderivatives 182
- 5.3 Summation vs. Integration 187
- 5.4 Indefinite Integrals and Substitutions 195
- 5.5 The Definite Integral 201
- 5.6 Properties of the Integral and the Average Value 206
- 5.7 The Fundamental Theorem and Its Consequences 213
- 5.8 Numerical Integration 220

## CHAPTER 6  Exponentials and Logarithms
- 6.1 An Overview 228
- 6.2 The Exponential $e^x$ 236
- 6.3 Growth and Decay in Science and Economics 242
- 6.4 Logarithms 252
- 6.5 Separable Equations Including the Logistic Equation 259
- 6.6 Powers Instead of Exponentials 267
- 6.7 Hyperbolic Functions 277

## CHAPTER 7  Techniques of Integration
- 7.1 Integration by Parts 283
- 7.2 Trigonometric Integrals 288
- 7.3 Trigonometric Substitutions 294
- 7.4 Partial Fractions 300
- 7.5 Improper Integrals 305

## CHAPTER 8  Applications of the Integral
- 8.1 Areas and Volumes by Slices 311
- 8.2 Length of a Plane Curve 320
- 8.3 Area of a Surface of Revolution 325
- 8.4 Probability and Calculus 328
- 8.5 Masses and Moments 336
- 8.6 Force, Work, and Energy 342
CHAPTER 4

Derivatives by the Chain Rule

4.1 The Chain Rule

You remember that the derivative of $f(x)g(x)$ is not $(df/dx)(dg/dx)$. The derivative of $\sin x$ times $x^2$ is not $\cos x$ times $2x$. The product rule gave two terms, not one term. But there is another way of combining the sine function $f$ and the squaring function $g$ into a single function. The derivative of that new function does involve the cosine times $2x$ (but with a certain twist). We will first explain the new function, and then find the "chain rule" for its derivative.

May I say here that the chain rule is important. It is easy to learn, and you will use it often. I see it as the third basic way to find derivatives of new functions from derivatives of old functions. (So far the old functions are $x^n$, $\sin x$, and $\cos x$. Still ahead are $e^x$ and $\log x$.) When $f$ and $g$ are added and multiplied, derivatives come from the sum rule and product rule. This section combines $f$ and $g$ in a third way.

The new function is $\sin(x^2)$—the sine of $x^2$. It is created out of the two original functions: if $x = 3$ then $x^2 = 9$ and $\sin(x^2) = \sin 9$. There is a "chain" of functions, combining $\sin x$ and $x^2$ into the composite function $\sin(x^2)$. You start with $x$, then find $g(x)$, then find $f(g(x))$:

The squaring function gives $y = x^2$. This is $g(x)$.

The sine function produces $z = \sin y = \sin(x^2)$. This is $f(g(x))$.

The "inside function" $g(x)$ gives $y$. This is the input to the "outside function" $f(y)$. That is called composition. It starts with $x$ and ends with $z$. The composite function is sometimes written $f \circ g$ (the circle shows the difference from an ordinary product $fg$). More often you will see $f(g(x))$:

$$z(x) = f(g(x)) = f(g(x)). \quad (1)$$

Other examples are $\cos 2x$ and $(2x)^3$, with $g = 2x$. On a calculator you input $x$, then push the "g" button, then push the "f" button:

From $x$ compute $y = g(x)$  \quad From $y$ compute $z = f(y)$.

There is not a button for every function! But the squaring function and sine function are on most calculators, and they are used in that order. Figure 4.1a shows how squaring will stretch and squeeze the sine function.
4.1 The Chain Rule

That graph of \( \sin x^2 \) is a crazy FM signal (the Frequency is Modulated). The wave goes up and down like \( \sin x \), but not at the same places. Changing to \( \sin g(x) \) moves the peaks left and right. Compare with a product \( g(x) \sin x \), which is an AM signal (the Amplitude is Modulated).

Remark \( f(g(x)) \) is usually different from \( g(f(x)) \). The order of \( f \) and \( g \) is usually important. For \( f(x) = \sin x \) and \( g(x) = x^2 \), the chain in the opposite order \( g(f(x)) \) gives something different:

First apply the sine function: \( y = \sin x \)
Then apply the squaring function: \( z = (\sin x)^2 \).

That result is often written \( \sin^2 x \), to save on parentheses. It is never written \( \sin x^2 \), which is totally different. Compare them in Figure 4.1.

\[ \begin{align*}
1 & \quad 1 \\
\text{y = sin} & \quad 1 \\
\text{sin}(x^2) & \quad y = (\sin x)^2 \\
\end{align*} \]

\[ \begin{align*}
1 & \quad 1 \\
\text{y = sin} & \quad 1 \\
\text{y = (sin x)^2} & \quad y = \sin x \\
\end{align*} \]

Fig. 4.1 \( f(g(x)) \) is different from \( g(f(x)) \). Apply \( g \) then \( f \), or \( f \) then \( g \).

**EXAMPLE 1** The composite function \( f \circ g \) can be deceptive. If \( g(x) = x^3 \) and \( f(y) = y^4 \), how does \( f(g(x)) \) differ from the ordinary product \( f(x)g(x) \)? The ordinary product is \( x^7 \). The chain starts with \( y = x^3 \), and then \( z = y^4 - x^{12} \). The composition of \( x^3 \) and \( y^4 \) gives \( f(g(x)) = x^{12} \).

**EXAMPLE 2** In Newton’s method, \( F(x) \) is composed with itself. This is iteration. Every output \( x_n \) is fed back as input, to find \( x_{n+1} = F(x_n) \). The example \( F(x) = \frac{1}{2}x + 4 \) has \( F(F(x)) = \frac{1}{2}(\frac{1}{2}x + 4) + 4 \). That produces \( z = \frac{1}{4}x + 6 \).

The derivative of \( F(x) \) is \( \frac{1}{2} \). The derivative of \( z = F(F(x)) \) is \( \frac{1}{2} \), which is \( \frac{1}{2} \) times \( \frac{1}{2} \).

We multiply derivatives. This is a special case of the chain rule.

An extremely special case is \( f(x) = x \) and \( g(x) = x \). The ordinary product is \( x^2 \). The chain \( f(g(x)) \) produces only \( x \) ! The output from the “identity function” is \( g(x) = x \).\footnote{A calculator has no button for the identity function. It wouldn't do anything.} When the second identity function operates on \( x \) it produces \( x \) again. The derivative is 1 times 1. I can give more composite functions in a table:

<table>
<thead>
<tr>
<th>( y = g(x) )</th>
<th>( z = f(y) )</th>
<th>( z = f(g(x)) )</th>
</tr>
</thead>
<tbody>
<tr>
<td>( x^2 - 1 )</td>
<td>( \sqrt{y} )</td>
<td>( \sqrt{x^2 - 1} )</td>
</tr>
<tr>
<td>( \cos x )</td>
<td>( y^3 )</td>
<td>( (\cos x)^3 )</td>
</tr>
<tr>
<td>( 2^x )</td>
<td>( 2^y )</td>
<td>( 2^x )</td>
</tr>
<tr>
<td>( x + 5 )</td>
<td>( y - 5 )</td>
<td>( x )</td>
</tr>
</tbody>
</table>

The last one adds 5 to get \( y \). Then it subtracts 5 to reach \( z \). So \( z = x \). Here output
equals input: \( f(g(x)) = x \). These "inverse functions" are in Section 4.3. The other examples create new functions \( z(x) \) and we want their derivatives.

**THE DERIVATIVE OF \( f(g(x)) \)**

What is the derivative of \( z = \sin x^2 \)? It is the limit of \( \Delta z/\Delta x \). Therefore we look at a nearby point \( x + \Delta x \). That change in \( x \) produces a change in \( y = x^2 \)—which moves to \( y + \Delta y = (x + \Delta x)^2 \). From this change in \( y \), there is a change in \( z = f(y) \). It is a "domino effect," in which each changed input yields a changed output: \( \Delta x \) produces \( \Delta y \) produces \( \Delta z \). We have to connect the final \( \Delta z \) to the original \( \Delta x \).

The key is to write \( \Delta z/\Delta x \) as \( \Delta z/\Delta y \) times \( \Delta y/\Delta x \). Then let \( \Delta x \) approach zero.

In the limit, \( \frac{dz}{dx} \) is given by the "chain rule":

\[
\frac{\Delta z}{\Delta x} = \frac{\Delta z}{\Delta y} \frac{\Delta y}{\Delta x} \quad \text{becomes the chain rule} \quad \frac{dz}{dx} = \frac{dz}{dy} \frac{dy}{dx}.
\]

(2)

As \( \Delta x \) goes to zero, the ratio \( \Delta y/\Delta x \) approaches \( dy/dx \). Therefore \( \Delta y \) must be going to zero, and \( \Delta z/\Delta y \) approaches \( dz/dy \). The limit of a product is the product of the separate limits (end of quick proof). We multiply derivatives:

**4A Chain Rule** Suppose \( g(x) \) has a derivative at \( x \) and \( f(y) \) has a derivative at \( y = g(x) \). Then the derivative of \( z = f(g(x)) \) is

\[
\frac{dz}{dx} = \frac{dz}{dy} \frac{dy}{dx} = f'(g(x)) \cdot g'(x).
\]

(3)

The slope at \( x \) is \( df/dy \) (at \( y \)) times \( dg/dx \) (at \( x \)).

**Caution** The chain rule does not say that the derivative of \( \sin x^2 \) is \( (\cos x)(2x) \). True, \( \cos y \) is the derivative of \( \sin y \). The point is that \( \cos y \) must be evaluated at \( y \) (not at \( x \)). We do not want \( df/dx \) at \( x \), we want \( df/dy \) at \( y = x^2 \):

The derivative of \( \sin x^2 \) is \( (\cos x^2)(2x) \).

(4)

**EXAMPLE 3** If \( z = (\sin x)^2 \) then \( dz/dx = (2 \sin x)(\cos x) \). Here \( y = \sin x \) is inside.

In this order, \( z = y^2 \) leads to \( dz/dy = 2y \). It does not lead to \( 2x \). The inside function \( \sin x \) produces \( dy/dx = \cos x \). The answer is \( 2y \cos x \). We have not yet found the function whose derivative is \( 2x \cos x \).

**EXAMPLE 4** The derivative of \( z = \sin 3x \) is \( \frac{dz}{dx} = \frac{dz}{dy} \frac{dy}{dx} = 3 \cos 3x \).
The outside function is \( z = \sin y \). The inside function is \( y = 3x \). Then \( dz/dy = \cos y \)—this is \( \cos 3x \), not \( \cos x \). Remember the other factor \( dy/dx = 3 \).

I can explain that factor 3, especially if \( x \) is switched to \( t \). The distance is \( z = \sin 3t \). That oscillates like \( \sin t \) except three times as fast. The speeded-up function \( \sin 3t \) completes a wave at time \( \frac{2\pi}{3} \) (instead of \( 2\pi \)). Naturally the velocity contains the extra factor 3 from the chain rule.

**EXAMPLE 5** Let \( z = f(y) = y^n \). Find the derivative of \( f(g(x)) = (g(x))^n \).

In this case \( dz/dy \) is \( n y^{n-1} \). The chain rule multiplies by \( dy/dx \):

\[
\frac{dz}{dx} = n y^{n-1} \frac{dy}{dx} \quad \text{or} \quad \frac{d}{dx} [(g(x))^n] = n[g(x)]^{n-1} \frac{dg}{dx}.
\]

This is the **power rule**! It was already discovered in Section 2.5. Square roots (when \( n = 1/2 \)) are frequent and important. Suppose \( y = x^2 - 1 \):

\[
\frac{d}{dx} \sqrt{x^2 - 1} = \frac{1}{2} (x^2 - 1)^{-1/2} (2x) = \frac{x}{\sqrt{x^2 - 1}}.
\]

**Question** A Buick uses 1/20 of a gallon of gas per mile. You drive at 60 miles per hour. How many gallons per hour?

**Answer** (Gallons/hour) = (gallons/mile)(miles/hour). The chain rule is \( (dy/dt) = (dy/dx)(dx/dt) \). The answer is \((1/20)(60) = 3 \text{ gallons/hour}\).

**Proof of the chain rule** The discussion above was correctly based on

\[
\frac{\Delta z}{\Delta x} = \frac{\Delta z}{\Delta y} \frac{\Delta y}{\Delta x} \quad \text{and} \quad \frac{dz}{dx} = \frac{dz}{dy} \frac{dy}{dx}.
\]

It was here, over the chain rule, that the “battle of notation” was won by Leibniz. His notation practically tells you what to do: Take the limit of each term. (I have to mention that when \( \Delta x \) is approaching zero, it is theoretically possible that \( \Delta y \) might hit zero. If that happens, \( \Delta z/\Delta y \) becomes \( 0/0 \). We have to assign it the correct meaning, which is \( dz/dy \).) As \( \Delta x \to 0 \),

\[
\frac{\Delta y}{\Delta x} \to g'(x) \quad \text{and} \quad \frac{\Delta z}{\Delta y} \to f'(y) = f'(g(x)).
\]

Then \( \Delta z/\Delta x \) approaches \( f'(y) \) times \( g'(x) \), which is the chain rule \( (dz/dy)(dy/dx) \). In the table below, the derivative of \( (\sin x)^3 \) is \( 3(\sin x)^2 \cos x \). That extra factor \( \cos x \) is easy to forget. It is even easier to forget the \(-1\) in the last example.

\[
\begin{align*}
z &= (x^3 + 1)^5 & dz/dx &= 5(x^3 + 1)^4 \times 3x^2 \\
z &= (\sin x)^3 & dz/dx &= 3 \sin^2 x \times \cos x \\
z &= (1 - x)^2 & dz/dx &= 2(1 - x) \times \cos x
\end{align*}
\]

**Important** All kinds of letters are used for the chain rule. We named the output \( z \). Very often it is called \( y \), and the inside function is called \( u \):

\[
\text{The derivative of } y = \sin u(x) \text{ is } \frac{dy}{dx} = \cos u \frac{du}{dx}.
\]

Examples with \( du/dx \) are extremely common. I have to ask you to accept whatever letters may come. What never changes is the key idea—derivative of outside function times derivative of inside function.
EXAMPLE 6 The chain rule is barely needed for $\sin(x-1)$. Strictly speaking the inside function is $u = x - 1$. Then $du/dx$ is just 1 (not $-1$). If $y = \sin(x-1)$ then $dy/dx = \cos(x-1)$. The graph is shifted and the slope shifts too.

Notice especially: The cosine is computed at $x-1$ and not at the unshifted $x$.

RECOGNIZING $f(y)$ AND $g(x)$

A big part of the chain rule is recognizing the chain. The table started with $(x^3 + 1)^5$.

You look at it for a second. Then you see it as $u^5$. The inside function is $u = x^3 + 1$.

With practice this decomposition (the opposite of composition) gets easy:

$$
\cos (2x+1) \text{ is } \cos u \quad \sqrt{1 + \sin t} \text{ is } \sqrt{u} \quad \sin x \text{ is ... (product rule!)}
$$

In calculations, the careful way is to write down all the functions:

$$
z = \cos u \quad u = 2x + 1 \quad dz/dx = (-\sin u)(2) = -2 \sin (2x + 1).
$$

The quick way is to keep in your mind “the derivative of what’s inside.” The slope of $\cos(2x+1)$ is $-\sin(2x+1)$, times 2 from the chain rule. The derivative of $2x + 1$ is remembered—without $z$ or $u$ or $f$ or $g$.

EXAMPLE 7 $\sin \sqrt{1-x}$ is a chain of $z = \sin y$, $y = \sqrt{u}$, $u = 1 - x$ (three functions).

With that triple chain you will have the hang of the chain rule:

The derivative of $\sin \sqrt{1-x}$ is $$(\cos \sqrt{1-x}) \left( \frac{1}{2\sqrt{1-x}} \right)(-1).$$

This is $(dz/dy)(dy/du)(du/dx)$. Evaluate them at the right places $y$, $u$, $x$.

Finally there is the question of second derivatives. The chain rule gives $dz/dx$ as a product, so $d^2z/dx^2$ needs the product rule:

$$
\frac{dz}{dx} = \frac{dz}{dy} \frac{dy}{dx} \quad \text{leads to} \quad \frac{d^2z}{dx^2} = \frac{dz}{dy} \frac{d^2y}{dx^2} + \frac{d}{dx} \left( \frac{dz}{dy} \right) \frac{dy}{dx}
$$

(8)

That last term needs the chain rule again. It becomes $d^2z/dy^2$ times $(dy/dx)^2$.

EXAMPLE 8 The derivative of $\sin x^2$ is $2x \cos x^2$. Then the product rule gives $d^2z/dx^2 = 2 \cos x^2 - 4x^2 \sin x^2$. In this case $y'' = 2$ and $(y')^2 = 4x^2$.

4.4 EXERCISES

Read-through questions

$z = f(g(x))$ comes from $z = f(y)$ and $y = ____$. At $x = 2$, the chain $(x^2-1)^3$ equals _____. Its inside function is $y = ____$, its outside function is $z = ____$. Then $dz/dx$ equals _____. The first factor is evaluated at $y = f(2)$. For $z = \sin(x^2-1)$ the derivative is _____. The triple chain $z = \cos(x+1)^2$ has a shift and a ____ and a cosine. Then $dz/dx = ____$. The proof of the chain rule begins with $\Delta z/\Delta x = (____)(____)$ and ends with _____. Changing letters, $y = ____$ has $dy/dx = ____$. The power rule for $y = [u(x)]^n$ is the chain rule $dy/dx = ____$. The slope of $5g(x)$ is ____ and the slope of $g(5x)$ is _____. When $f = \cos$ and $g = \sin$ and $x = 0$, the numbers $f(g(x))$ and $g(f(x))$ and $f(x)g(x)$ are ____.

In 1–10 identify $f(y)$ and $g(x)$. From their derivatives find $\frac{dz}{dx}$.

$1 \ z = (x^2 - 3)^3 \ 2 \ z = (x^3 - 3)^2$

$3 \ z = \cos(x^3) \ 4 \ z = \tan 2x$

$5 \ z = \sqrt{\sin x} \ 6 \ z = \sin \sqrt{x}$
4.1 The Chain Rule

7 $z = \tan(1/x) + 1/\tan x$ \hspace{1cm} 8 $z = \sin(\cos x)$
9 $z = \cos(x^2 + x + 1)$ \hspace{1cm} 10 $z = \sqrt{x^2}$

In 11–16 write down $dz/dx$. Don’t write down $f$ and $g$.
11 $z = \sin(17x)$ \hspace{1cm} 12 $z = \tan(x + 1)$
13 $z = \cos(\cos x)$ \hspace{1cm} 14 $z = (x^3)^{3/2}$
15 $z = x^2 \sin x$ \hspace{1cm} 16 $z = (9x + 4)^{3/2}$

Problems 17–22 involve three functions $z(y)$, $y(u)$, and $u(x)$. Find $dz/dx$ from $(dz/dy)(dy/du)(du/dx)$.
17 $z = \sin \sqrt{x + 1}$ \hspace{1cm} 18 $z = \sqrt{\sin(x + 1)}$
19 $z = \sqrt{1 + \sin x}$ \hspace{1cm} 20 $z = \sin(\sqrt{x + 1})$
21 $z = \sin(1/\sin x)$ \hspace{1cm} 22 $z = (\sin x^2)^2$

In 23–26 find $dz/dx$ by the chain rule and also by rewriting $z$.
23 $z = ((x^2)^3)^3$ \hspace{1cm} 24 $z = (3x)^3$
25 $z = (x + 1)^2 + \sin(x + 1)$ \hspace{1cm} 26 $z = \sqrt{1 - \cos^2 x}$

27 If $f(x) = x^2 + 1$ what is $f(f(x))$? If $U(x)$ is the unit step function (from 0 to 1 at $x = 0$) draw the graphs of $\sin U(x)$ and $U(\sin x)$. If $R(x)$ is the ramp function $(x + |x|)$, draw the graphs of $R(x)$ and $R(\sin x)$.

28 (Recommended) If $g(x) = x^3$ find $f(y)$ so that $f(g(x)) = x^3 + 1$. Then find $h(y)$ so that $h(g(x)) = x$. Then find $k(y)$ so that $k(g(x)) = 1$.
29 If $f(y) = y^2$ find $g(x)$ so that $f(g(x)) = x$. Then find $h(x)$ so that $h(f(x)) = x^2$. Then find $k(x)$ so that $f(k(x)) = 1$.
30 Find two different pairs $f(y)$, $g(x)$ so that $f(g(x)) = \sqrt{1 - x^2}$.

31 The derivative of $f(x)$ is $\ldots$. Is it $(df/dx)^2$? Test your formula on $f(x) = 1/x$.
32 If $f(3) = 3$ and $g(3) = 5$ and $f'(3) = 2$ and $g'(3) = 4$, find the derivative at $x = 3$ if possible for $a)$ $f(x)g(x)$ \hspace{1cm} $b)$ $f(g(x))$ \hspace{1cm} $c)$ $g(f(x))$ \hspace{1cm} $d)$ $f(f(x))$
33 For $F(x) = \frac{1}{2}x + 8$, show how iteration gives $F(F(x)) = \frac{1}{4}x + 12$. Find $F(F(F(x)))$—also called $F^3(x)$. The derivative of $F^4(x)$ is $\ldots$.
34 In Problem 33 the limit of $F^{10}(x)$ is a constant $C = \ldots$. From any start (try $x = 0$) the iterations $x_{n+1} = F(x_n)$ converge to $C$.
35 Suppose $g(x) = 3x + 1$ and $f(y) = \frac{1}{2}(y - 1)$. Then $f(g(x)) = \ldots$ and $g(f(y)) = \ldots$. These are inverse functions.
36 Suppose $g(x)$ is continuous at $x = 4$, say $g(4) = 7$. Suppose $f(y)$ is continuous at $y = 7$, say $f(7) = 9$. Then $f(g(x))$ is continuous at $x = 4$ and $f(g(4)) = 9$.

Proof is given. Because $\ldots$ is continuous, there is a $\delta$ such that $|f(g(x)) - 9| < \varepsilon$ whenever $|g(x) - 7| < \delta$. Then $\ldots$ is continuous, there is a $\theta$ such that $|g(x) - 7| < \delta$ whenever $|x - 4| < \theta$. Conclusion: If $|x - 4| < \theta$ then $\ldots$. This shows that $f(g(x))$ approaches $f(g(4))$.

37 Only six functions can be constructed by compositions (in any sequence) of $g(x) = 1 - x$ and $f(x) = 1/x$. Starting with $g$ and $f$, find the other four.
38 If $g(x) = 1 - x$ then $g(g(x)) = 1 - (1 - x) = x$. If $f(x) = 1/x$ then $g(g(x)) = 1/(1/x) = x$. Draw graphs of those $g's$ and explain from the graphs why $g(g(x)) = x$. Find two more $g's$ with this special property.

39 Construct functions so that $f(g(x))$ is always zero, but $f(y)$ is not always zero.

40 True or false
(a) If $f(x) = f(-x)$ then $f'(x) = f'(-x)$.
(b) The derivative of the identity function is zero.
(c) The derivative of $f(1/x)$ is $-1/(f(x))^2$.
(d) The derivative of $f(1 + x)$ is $f'(1 + x)$.
(e) The second derivative of $f(g(x))$ is $f''(g(x))g''(x)$.

41 On the same graph draw the parabola $y = x^2$ and the curve $z = \sin y$ (keep $y$ upwards, with $x$ and $z$ across). Starting at $x = 3$ find your way to $z = \sin 9$.
42 On the same graph draw $y = \sin x$ and $z = y^2$ ($y$ upwards for both). Starting at $x = \pi/2$ find $z = (\sin 2x)^2$ on the graph.
43 Find the second derivative of $a)$ $\sin(x^2 + 1)$ \hspace{1cm} $b)$ $\sqrt{x^2 - 1}$ \hspace{1cm} $c)$ $\cos \sqrt{x}$
44 Explain why $\frac{d}{dx}\left(\frac{dz}{dy}\frac{dy}{du}\frac{du}{dx}\right) = \left(\frac{dz}{dy}\right)^2\frac{dy}{dx}$ in equation (8).

Check this when $z = y^2$, $y = x^3$.

Final practice with the chain rule and other rules (and other letters!). Find the $x$ or $t$ derivative of $z$ or $y$.
45 $z = f(u(t))$ \hspace{1cm} 46 $z = u^3$, $u = x^3$
47 $y = \sin u(x)\cos u(x)$ \hspace{1cm} 48 $y = \sqrt{u(t)}$
49 $y = x^2 u(x)$ \hspace{1cm} 50 $y = f(x^2) + f(x)^2$
51 $z = \sqrt{1 - u}$, $u = \sqrt{1 - x}$ \hspace{1cm} 52 $z = 1/u^2(t)$
53 $z = f(u)$, $u = t^2$, $u = \sqrt{t}$ \hspace{1cm} 54 $y = u$, $u = x$, $x = 1/t$
55 If $f = x^4$ and $g = x^3$ then $f' = 4x^3$ and $g' = 3x^2$. The chain rule multiplies derivatives to get $12x^5$. But $f(g(x)) = x^{12}$ and its derivative is not $12x^5$. Where is the flaw?
56 The derivative of $y = \sin(\sin x)$ is $dy/dx = \cos(\cos x) - \sin(\sin x)\cos x$.
57 (a) A book has 400 words per page. There are 9 pages per section. So there are _____ words per section.
(b) You read 200 words per minute. So you read _____ pages per minute. How many minutes per section?
58 (a) You walk in a train at 3 miles per hour. The train moves at 50 miles per hour. Your ground speed is _____ miles per hour.

(b) You walk in a train at 3 miles per hour. The train is shown on TV (1 mile train = 20 inches on TV screen). Your speed across the screen is _____ inches per hour.

59 Coke costs 1/3 dollar per bottle. The buyer gets ______ bottles per dollar. If \( dy/dx = 1/3 \) then \( dx/dy = \) ______.

60 (Computer) Graph \( F(x) = \sin x \) and \( G(x) = \sin(\sin x) \)---not much difference. Do the same for \( F'(x) \) and \( G'(x) \). Then plot \( F''(x) \) and \( G''(x) \) to see where the difference shows up.

4.2 Implicit Differentiation and Related Rates

We start with the equations \( xy = 2 \) and \( y^5 + xy = 3 \). As \( x \) changes, these \( y \)'s will change—to keep \((x, y)\) on the curve. We want to know \( dy/dx \) at a typical point. For \( xy = 2 \) that is no trouble, but the slope of \( y^5 + xy = 3 \) requires a new idea.

In the first case, solve for \( y = 2/x \) and take its derivative: \( dy/dx = -2/x^2 \). The curve is a hyperbola. At \( x = 2 \) the slope is \( -2/4 = -1/2 \).

The problem with \( y^5 + xy = 3 \) is that it can’t be solved for \( y \). Galois proved that there is no solution formula for fifth-degree equations.† The function \( y(x) \) cannot be given explicitly. All we have is the implicit definition of \( y \), as a solution to \( y^5 + xy = 3 \). The point \( x = 2, y = 1 \) satisfies the equation and lies on the curve, but how to find \( dy/dx \)?

This section answers that question. It is a situation that often occurs. Equations like \( \sin y + \sin x = 1 \) or \( \sin y = x \) (maybe even \( \sin y = x \)) are difficult or impossible to solve directly for \( y \). Nevertheless we can find \( dy/dx \) at any point.

The way out is implicit differentiation. Work with the equation as it stands. Find the \( x \) derivative of every term in \( y^5 + xy = 3 \). That includes the constant term 3, whose derivative is zero.

**EXAMPLE 1** The power rule for \( y^5 \) and the product rule for \( xy \) yield

\[
5y^4 \frac{dy}{dx} + x \frac{dy}{dx} + y = 0. \tag{1}
\]

Now substitute the typical point \( x = 2 \) and \( y = 1 \), and solve for \( dy/dx \):

\[
5 \frac{dy}{dx} + 2 \frac{dy}{dx} + 1 = 0 \quad \text{produces} \quad \frac{dy}{dx} = -\frac{1}{7}. \tag{2}
\]

This is implicit differentiation (ID), and you see the idea: Include \( dy/dx \) from the chain rule, even if \( y \) is not known explicitly as a function of \( x \).

**EXAMPLE 2** \( \sin y + \sin x = 1 \) leads to \( \cos y \frac{dy}{dx} + \cos x = 0 \)

**EXAMPLE 3** \( y \sin y = x \) leads to \( y \cos y \frac{dy}{dx} + \sin y \frac{dy}{dx} = 1 \)

Knowing the slope makes it easier to draw the curve. We still need points \((x, y)\) that satisfy the equation. Sometimes we can solve for \( x \). Dividing \( y^5 + xy = 3 \) by \( y \)

†That was before he went to the famous duel, and met his end. Fourth-degree equations do have a solution formula, but it is practically never used.
4.2 Implicit Differentiation and Related Rates

gives \( x = 3/y - y^4 \). Now the derivative (the \( x \) derivative) is

\[
1 = \left( -\frac{3}{y^2} - 4y^3 \right) \frac{dy}{dx} = -7 \frac{dy}{dx} \text{ at } y = 1. \tag{3}
\]

Again \( \frac{dy}{dx} = -1/7 \). All these examples confirm the main point of the section:

4B (Implicit differentiation) An equation \( F(x, y) = 0 \) can be differentiated directly by the chain rule, without solving for \( y \) in terms of \( x \).

The example \( xy = 2 \), done implicitly, gives \( x \frac{dy}{dx} + y = 0 \). The slope \( \frac{dy}{dx} \) is \(-y/x\). That agrees with the explicit slope \(-2/x^2\).

ID is explained better by examples than theory (maybe everything is). The essential theory can be boiled down to one idea: "Go ahead and differentiate."

EXAMPLE 4 Find the tangent direction to the circle \( x^2 + y^2 = 25 \).

We can solve for \( y = \pm \sqrt{25 - x^2} \), or operate directly on \( x^2 + y^2 = 25 \):

\[
2x + 2y \frac{dy}{dx} = 0 \quad \text{or} \quad \frac{dy}{dx} = -\frac{x}{y}. \tag{4}
\]

Compare with the radius, which has slope \( y/x \). The radius goes across \( x \) and up \( y \). The tangent goes across \(-y\) and up \( x \). The slopes multiply to give \((-x/y)(y/x) = -1\).

To emphasize implicit differentiation, go on to the second derivative. The top of the circle is concave down, so \( d^2y/dx^2 \) is negative. Use the quotient rule on \(-x/y\):

\[
\frac{dy}{dx} = -\frac{x}{y} \quad \text{so} \quad \frac{d^2y}{dx^2} = \frac{-y \cdot \frac{dx}{dx} - x \cdot \frac{dy}{dx}}{y^2} = -\frac{x + (x^2/y)}{y^2} = -\frac{y^2 + x^2}{y^3}. \tag{5}
\]

RELATED RATES

There is a group of problems that has never found a perfect place in calculus. They seem to fit here—as applications of the chain rule. The problem is to compute \( df/dt \), but the odd thing is that we are given another derivative \( dg/dt \). To find \( df/dt \), we need a relation between \( f \) and \( g \).

The chain rule is \( df/dt = (df/dg)(dg/dt) \). Here the variable is \( t \) because that is typical in applications. From the rate of change of \( g \) we find the rate of change of \( f \). This is the problem of related rates, and examples will make the point.

EXAMPLE 5 The radius of a circle is growing by \( dr/dt = 7 \). How fast is the circumference growing? Remember that \( C = 2\pi r \) (this relates \( C \) to \( r \)).

Solution

\[
\frac{dC}{dt} = \frac{dC}{dr} \frac{dr}{dt} = (2\pi)(7) = 14\pi.
\]

That is pretty basic, but its implications are amazing. Suppose you want to put a rope around the earth that any 7-footer can walk under. If the distance is 24,000 miles, what is the additional length of the rope? Answer: Only 14\( \pi \) feet.

More realistically, if two lanes on a circular track are separated by 5 feet, how much head start should the outside runner get? Only 10\( \pi \) feet. If your speed around a turn is 55 and the car in the next lane goes 56, who wins? See Problem 14.

Examples 6–8 are from the 1988 Advanced Placement Exams (copyright 1989 by the College Entrance Examination Board). Their questions are carefully prepared.
EXAMPLE 6  The sides of the rectangle increase in such a way that $dz/dt = 1$ and $dx/dt = 3dy/dt$. At the instant when $x = 4$ and $y = 3$, what is the value of $dx/dt$?

Solution  The key relation is $x^2 + y^2 = z^2$. Take its derivative (implicitly):

$$\frac{2x}{dt} \frac{dx}{dt} + 2y \frac{dy}{dt} = 2z \frac{dz}{dt}$$

produces

$$8 \frac{dx}{dt} + 6 \frac{dy}{dt} = 10.$$ 

We used all information, including $z = 5$, except for $dx/dt = 3dy/dt$. The term $6dy/dt$ equals $2dx/dt$, so we have $10dx/dt = 10$. Answer: $dx/dt = 1$.

EXAMPLE 7  A person 2 meters tall walks directly away from a streetlight that is 8 meters above the ground. If the person's shadow is lengthening at the rate of $4/9$ meters per second, at what rate in meters per second is the person walking?

Solution  Draw a figure! You must relate the shadow length $s$ to the distance $x$ from the streetlight. The problem gives $ds/dt = 4/9$ and asks for $dx/dt$:

By similar triangles $\frac{x}{6} = \frac{s}{2}$, so

$$\frac{dx}{dt} = \frac{6}{2} \frac{ds}{dt} = (3) \left( \frac{4}{9} \right) = \frac{4}{3}.$$ 

Note  This problem was hard. I drew three figures before catching on to $x$ and $s$. It is interesting that we never knew $x$ or $s$ or the angle.

EXAMPLE 8  An observer at point $A$ is watching balloon $B$ as it rises from point $C$. (The figure is given.) The balloon is rising at a constant rate of 3 meters per second (this means $dy/dt = 3$) and the observer is 100 meters from point $C$.

(a) Find the rate of change in $z$ at the instant when $y = 50$. (They want $dz/dt$.)

$$z^2 = y^2 + 100^2 \Rightarrow 2z \frac{dz}{dt} = 2y \frac{dy}{dt}$$

$$z = \sqrt{50^2 + 100^2} = 50 \sqrt{5} \Rightarrow \frac{dz}{dt} = \frac{2 \cdot 50 \cdot 3}{2 \cdot 50 \sqrt{5}} = \frac{3 \sqrt{5}}{5}.$$ 

(b) Find the rate of change in the area of right triangle $BCA$ when $y = 50$.

$$A = \frac{1}{2} (100)(y) = 50y \quad \frac{dA}{dt} = 50 \frac{dy}{dt} = 50 \cdot 3 = 150.$$ 

(c) Find the rate of change in $\theta$ when $y = 50$. (They want $d\theta/dt$.)

$$y = 50 \Rightarrow \cos \theta = \frac{100}{50 \sqrt{5}} = \frac{2}{\sqrt{5}}$$

$$\tan \theta = \frac{y}{100} \Rightarrow \sec^2 \theta \frac{d\theta}{dt} = \frac{1}{100} \frac{dy}{dt} \Rightarrow \frac{d\theta}{dt} = \left( \frac{2}{\sqrt{5}} \right)^2 \frac{3}{100} = \frac{3}{125}.$$
In all problems I first wrote down a relation from the figure. Then I took its derivative. Then I substituted known information. (The substitution is after taking the derivative of tan θ = y/100. If we substitute y = 50 too soon, the derivative of 50/100 is useless.)

"Candidates are advised to show their work in order to minimize the risk of not receiving credit for it." 50% solved Example 6 and 21% solved Example 7. From 12,000 candidates, the average on Example 8 (free response) was 6.1 out of 9.

**Example 9**  
A is a lighthouse and BC is the shoreline (same figure as the balloon). The light at A turns once a second (dθ/dt = 2π radians/second). How quickly does the receiving point B move up the shoreline?

**Solution**  
The figure shows y = 100 tan θ. The speed dy/dt is 100 sec²θ dθ/dt. This is 200π sec²θ, so B speeds up as sec θ increases.

**Paradox**  
When θ approaches a right angle, sec θ approaches infinity. So does dy/dt. B moves faster than light! This contradicts Einstein's theory of relativity. The paradox is resolved (I hope) in Problem 18.

If you walk around a light at A, your shadow at B seems to go faster than light. Same problem. This speed is impossible—something has been forgotten.

**Smaller paradox** (not destroying the theory of relativity). The figure shows y = z sin θ. Apparently dy/dt = (dz/dt) sin θ. This is totally wrong. Not only is it wrong, the exact opposite is true: dz/dt = (dy/dt) sin θ. If you can explain that (Problem 15), then ID and related rates hold no terrors.

### 4.2 Exercises

**Read-through questions**

For x³ + y³ = 2 the derivative dy/dx comes from differentiation. We don’t have to solve for... Replace y by √x, this is dy/dx = __.

A second example is y² = x. The x derivative of this equation is ... Therefore dy/dx = __. Replacing y by √x, this is dy/dx = __.

In related rates, we are given dg/dt and we want df/dt. We need a relation between f and __. If f = g², then df/dt = __. If f² + g² = 1, then df/dt = __. If the sides of a cube grow by ds/dt = 2, then its volume grows by dV/dt = __. To find a number (8 is wrong), you also need to know __.

**By implicit differentiation find dy/dx in 1–10.**

1. y³ + x³ = 1  
2. 2x³y + y³x = 1  
3. (x − y)² = 4  
4. √x + √y = 3 at x = 4  
5. x = F(y)  
6. f(x) + F(y) = xy  
7. x²y = y²x  
8. x = sin y  
9. x = tan y  
10. y² = x at x = 1

**Problems 15–18 resolve the speed of light paradox in Example 9.**

15. (Small paradox first) The right triangle has z² = y² + 100². Take the t derivative to show that z' = y' sin θ.

16. (Even smaller paradox) As B moves up the line, why is dy/dt larger than dz/dt? Certainly z is larger than y. But as θ increases they become __.

17. (Faster than light) The derivative of y = 100 tan θ in Example 9 is y' = 100 sec²θ θ' = 200π sec²θ. Therefore y'
passes \(c\) (the speed of light) when \(\sec^2 \theta\) passes \(_______\). Such a speed is impossible—we forget that light takes time to reach \(B\).

18 (Explanation by ID) Light travels from \(A\) to \(B\) in time \(z/c\), distance over speed. Its arrival time is \(t = \theta/2\pi + z/c\) so \(\theta'/2\pi = 1 - z'/c\). Then \(z' = y' \sin \theta\) and \(y' = 100 \sec^2 \theta \theta'\) (all these are ID) lead to

\[
y' = 200nc(c \cos^2 \theta + 200x \sin \theta)
\]

As \(\theta\) approaches \(\pi/2\), this speed approaches \(_______\). Note: \(y'\) still exceeds \(c\) for some negative angle. That is for Einstein to explain. See the 1985 College Math Journal, page 186, and the 1960 Scientific American, "Things that go faster than light."

19 If a plane follows the curve \(y = f(x)\), and its ground speed is \(dx/dt = 500\) mph, how fast is the plane going up? How fast is the plane going?

20 Why can’t we differentiate \(x = 7\) and reach \(1 = 0\)??

Problems 21–29 are applications of related rates.

21 (Calculus classic) The bottom of a 10-foot ladder is going away from the wall at \(dx/dt = 2\) feet per second. How fast is the top going down the wall? Draw the right triangle to find \(dy/dt\) when the height \(y\) is (a) 6 feet (b) 5 feet (c) zero.

22 The top of the 10-foot ladder can go faster than light. At what height \(y\) does \(dy/dt = -c\)?

23 How fast does the level of a Coke go down if you drink a cubic inch a second? The cup is a cylinder of radius 2 inches—first write down the volume.

24 A jet flies at 8 miles up and 560 miles per hour. How fast is it approaching you when (a) it is 16 miles from you; (b) its shadow is 8 miles from you (the sun is overhead); (c) the plane is 8 miles from you (exactly above)?

25 Starting from a 3–4–5 right triangle, the short sides increase by 2 meters/second but the angle between them decreases by 1 radian/second. How fast does the area increase or decrease?

26 A pass receiver is at \(x = 4\), \(y = 8t\). The ball thrown at \(t = 3\) is at \(x = c(t - 3)\), \(y = 10c(t - 3)\).
   (a) Choose \(c\) so the ball meets the receiver.
   *(b) At that instant the distance \(D\) between them is changing at what rate?

27 A thief is 10 meters away (8 meters ahead of you, across a street 6 meters wide). The thief runs on that side at 7 meters/second, you run at 9 meters/second. How fast are you approaching if (a) you follow on your side; (b) you run toward the thief; (c) you run away on your side?

28 A spherical raindrop evaporates at a rate equal to twice its surface area. Find \(dr/dt\).

29 Starting from \(P = V = 5\) and maintaining \(PV = T\), find \(dV/dt\) if \(dP/dt = 2\) and \(dT/dt = 3\).

30 (a) The crankshaft \(AB\) turns twice a second so \(d\theta/dt = \_______\).
   (b) Differentiate the cosine law \(6^2 = 3^2 + x^2 - 2(3x \cos \theta)\) to find the piston speed \(dx/dt\) when \(\theta = \pi/2\) and \(\theta = \pi\).

31 A camera turns at \(C\) to follow a rocket at \(R\).
   (a) Relate \(dz/dt\) to \(dy/dt\) when \(y = 10\).
   (b) Relate \(d\theta/dt\) to \(dy/dt\) based on \(y = 10 \tan \theta\).
   (c) Relate \(d^2\theta/dt^2\) to \(d^2y/dt^2\) and \(dy/dt\).

---

3.3 Inverse Functions and Their Derivatives

There is a remarkable special case of the chain rule. It occurs when \(f(y)\) and \(g(x)\) are "inverse functions." That idea is expressed by a very short and powerful equation: \(f(g(x)) = x\). Here is what that means.

**Inverse functions:** Start with any input, say \(x = 5\). Compute \(y = g(x)\), say \(y = 3\). Then compute \(f(y)\), and the answer must be \(5\). What one function does, the inverse function...
undoes. If \( g(5) = 3 \) then \( f(3) = 5 \). The inverse function \( f \) takes the output \( y \) back to the input \( x \).

**EXAMPLE 1**  
\( g(x) = x - 2 \) and \( f(y) = y + 2 \) are inverse functions. Starting with \( x = 5 \), the function \( g \) subtracts 2. That produces \( y = 3 \). Then the function \( f \) adds 2. That brings back \( x = 5 \). To say it directly: The inverse of \( y = x - 2 \) is \( x = y + 2 \).

**EXAMPLE 2**  
\( y = g(x) = \frac{x}{9}(x - 32) \) and \( x = f(y) = \frac{9}{2}y + 32 \) are inverse functions (for temperature). Here \( x \) is degrees Fahrenheit and \( y \) is degrees Celsius. From \( x = 32 \) (freezing in Fahrenheit) you find \( y = 0 \) (freezing in Celsius). The inverse function takes \( y = 0 \) back to \( x = 32 \). Figure 4.4 shows how \( x = 50^\circ F \) matches \( y = 10^\circ C \).

Notice that \( \frac{x}{9}(x - 32) \) subtracts 32 first. The inverse \( \frac{9}{2}y + 32 \) adds 32 last. In the same way \( g \) multiplies last by \( \frac{9}{2} \) while \( f \) multiplies first by \( \frac{9}{2} \).

The inverse function is written \( f = g^{-1} \) and pronounced "\( g \) inverse." It is not \( 1/g(x) \).

If the demand \( y \) is a function of the price \( x \), then the price is a function of the demand. Those are inverse functions. Their derivatives obey a fundamental rule: \( dy/dx \times dx/dy \) equals 1. In Example 2, \( dy/dx \) is \( 5/9 \) and \( dx/dy \) is \( 9/5 \).

There is another important point. When \( f \) and \( g \) are applied in the opposite order, they still come back to the start. First \( f \) adds 2, then \( g \) subtracts 2. The chain \( g(f(y)) = (y + 2) - 2 \) brings back \( y \). If \( f \) is the inverse of \( g \) then \( g \) is the inverse of \( f \). The relation is completely symmetric, and so is the definition:

**Inverse function:** If \( y = g(x) \) then \( x = g^{-1}(y) \). If \( x = g^{-1}(y) \) then \( y = g(x) \).

The loop in the figure goes from \( x \) to \( y \) to \( x \). The composition \( g^{-1}(g(x)) \) is the "identity function." Instead of a new point \( z \) it returns to the original \( x \). This will make the chain rule particularly easy—leading to \( (dy/dx)(dx/dy) = 1 \).

**EXAMPLE 3**  
\( y = g(x) = \sqrt{x} \) and \( x = f(y) = y^2 \) are inverse functions.

Starting from \( x = 9 \) we find \( y = 3 \). The inverse gives \( 3^2 = 9 \). The square of \( \sqrt{x} \) is \( f(g(x)) = x \). In the opposite direction, the square root of \( y^2 \) is \( g(f(y)) = y \).

**Caution** That example does not allow \( x \) to be negative. The domain of \( g \)—the set of numbers with square roots—is restricted to \( x > 0 \). This matches the range of \( g^{-1} \).

The outputs \( y^2 \) are nonnegative. With domain of \( g = \text{range of } g^{-1} \), the equation \( x = (\sqrt{x})^2 \) is possible and true. The nonnegative \( x \) goes into \( g \) and comes out of \( g^{-1} \).

In this example \( y \) is also nonnegative. You might think we could square anything, but \( y \) must come back as the square root of \( y^2 \). So \( y \geq 0 \).

To summarize: The domain of a function matches the range of its inverse. The inputs to \( g^{-1} \) are the outputs from \( g \). The inputs to \( g \) are the outputs from \( g^{-1} \).
If \( g(x) = y \) then solving that equation for \( x \) gives \( x = g^{-1}(y) \):

- If \( y = 3x - 6 \) then \( x = \frac{1}{3}(y + 6) \) (this is \( g^{-1}(y) \))
- If \( y = x^2 + 1 \) then \( x = \pm \sqrt{y - 1} \) (this is \( g^{-1}(y) \))

In practice that is how \( g^{-1} \) is computed: Solve \( g(x) = y \). This is the reason inverses are important. Every time we solve an equation we are computing a value of \( g^{-1} \).

Not all equations have one solution. Not all functions have inverses. For each \( y \), the equation \( g(x) = y \) is only allowed to produce one \( x \). That solution is \( x = g^{-1}(y) \). If there is a second solution, then \( g^{-1} \) will not be a function—because a function cannot produce two \( x \)'s from the same \( y \).

**EXAMPLE 4** There is more than one solution to \( \sin x = \frac{1}{2} \). Many angles have the same sine. On the interval \( 0 \leq x \leq \pi \), the inverse of \( y = \sin x \) is not a function. Figure 4.5 shows how two \( x \)'s give the same \( y \).

Prevent \( x \) from passing \( \pi/2 \) and the sine has an inverse. Write \( x = \sin^{-1} y \).

The function \( g \) has no inverse if two points \( x_1 \) and \( x_2 \) give \( g(x_1) = g(x_2) \). Its inverse would have to bring the same \( y \) back to \( x_1 \) and \( x_2 \). No function can do that; \( g^{-1}(y) \) cannot equal both \( x_1 \) and \( x_2 \). There must be only one \( x \) for each \( y \).

To be invertible over an interval, \( g \) must be steadily increasing or steadily decreasing.

Fig. 4.5 Inverse exists (one \( x \) for each \( y \)). No inverse function (two \( x \)'s for one \( y \)).

**THE DERIVATIVE OF \( g^{-1} \)**

It is time for calculus. Forgive me for this very humble example.

**EXAMPLE 5** (ordinary multiplication) The inverse of \( y = g(x) = 3x \) is \( x = f(y) = \frac{1}{3}y \).

This shows with special clarity the rule for derivatives: The slopes \( dy/dx = 3 \) and \( dx/dy = \frac{1}{3} \) multiply to give 1. This rule holds for all inverse functions, even if their slopes are not constant. It is a crucial application of the chain rule to the derivative of \( f(g(x)) = x \).

<table>
<thead>
<tr>
<th>4C (Derivative of inverse function) From ( f(g(x)) = x ) the chain rule gives ( f'(g(x))g'(x) = 1 ). Writing ( y = g(x) ) and ( x = f(y) ), this rule looks better:</th>
</tr>
</thead>
<tbody>
<tr>
<td>( \frac{dx}{dy} ) or ( \frac{dx}{dy} = \frac{1}{dy/dx} ).</td>
</tr>
</tbody>
</table>

The slope of \( x = g^{-1}(y) \) times the slope of \( y = g(x) \) equals one.

This is the chain rule with a special feature. Since \( f(g(x)) = x \), the derivative of both sides is 1. If we know \( g' \) we now know \( f' \). That rule will be tested on a familiar example. In the next section it leads to totally new derivatives.
EXAMPLE 6  The inverse of $y = x^3$ is $x = y^{1/3}$. We can find $dx/dy$ two ways:

- **directly:** $\frac{dx}{dy} = \frac{1}{3}y^{-2/3}
- **indirectly:** $\frac{dx}{dy} = \frac{1}{dy/dx} = \frac{1}{3x^2} = \frac{1}{3y^{2/3}}$.

The equation $(dx/dy)(dy/dx) = 1$ is not ordinary algebra, but it is true. Those derivatives are limits of fractions. The fractions are $(\Delta x/\Delta y)(\Delta y/\Delta x) = 1$ and we let $\Delta x \to 0$.

---

**Fig. 4.6** Graphs of inverse functions: $x = \frac{1}{3}y$ is the mirror image of $y = 3x$.

---

Before going to new functions, I want to draw graphs. Figure 4.6 shows $y = \sqrt{x}$ and $y = 3x$. What is special is that the same graphs also show the inverse functions. The inverse of $y = \sqrt{x}$ is $x = y^2$. The pair $x = 4, y = 2$ is the same for both. That is the whole point of inverse functions—if $2 = g(4)$ then $4 = g^{-1}(2)$. Notice that the graphs go steadily up.

The only problem is, the graph of $x = g^{-1}(y)$ is on its side. To change the slope from 3 to $\frac{1}{3}$, you would have to turn the figure. After that turn there is another problem—the axes don’t point to the right and up. You also have to look in a mirror! (The typesetter refused to print the letters backward. He thinks it’s crazy but it’s not.) To keep the book in position, and the typesetter in position, we need a better idea.

The graph of $x = \frac{1}{3}y$ comes from turning the picture across the 45° line. The $y$ axis becomes horizontal and $x$ goes upward. The point (2, 6) on the line $y = 3x$ goes into the point (6, 2) on the line $x = \frac{1}{3}y$. The eyes see a reflection across the 45° line (Figure 4.6c). The mathematics sees the same pairs $x$ and $y$. The special properties of $g$ and $g^{-1}$ allow us to know two functions—and draw two graphs—at the same time.† The graph of $x = g^{-1}(y)$ is the mirror image of the graph of $y = g(x)$.

---

**EXponentials and Logarithms**

I would like to add two more examples of inverse functions, because they are so important. Both examples involve the exponential and the logarithm. One is made up of linear pieces that imitate $2^n$; it appeared in Chapter 1. The other is the true function $2^x$, which is not yet defined—and it is not going to be defined here. The functions $b^x$ and $\log_b y$ are so overwhelmingly important that they deserve and will get a whole chapter of the book (at least). But you have to see the graphs.

The slopes in the linear model are powers of 2. So are the heights $y$ at the start of each piece. The slopes $1, 2, 4, \ldots$ equal the heights $1, 2, 4, \ldots$ at those special points.

The inverse is a discrete model for the logarithm (to base 2). The logarithm of 1 is 0, because $2^0 = 1$. The logarithm of 2 is 1, because $2^1 = 2$. The logarithm of $2^j$ is the exponent $j$. Thus the model gives the correct $x = \log_2 y$ at the breakpoints $y = 1, 2, 4, 8, \ldots$. The slopes are $1, \frac{1}{2}, \frac{1}{4}, \frac{1}{8}, \ldots$ because $dx/dy = 1/(dy/dx)$.

†I have seen graphs with $y = g(x)$ and also $y = g^{-1}(x)$. For me that is wrong; it has to be $x = g^{-1}(y)$. If $y = \sin x$ then $x = \sin^{-1} y$. 
The model is good, but the real thing is better. The figure on the right shows the true exponential \( y = 2^x \). At \( x = 0, 1, 2, \ldots \) the heights \( y \) are the same as before. But now the height at \( x = \frac{1}{2} \) is the number \( 2^{1/2} \), which is \( \sqrt{2} \). The height at \( x = .10 \) is the tenth root \( 2^{1/10} = 1.07 \ldots \). The slope at \( x = 0 \) is no longer 1—it is closer to \( \Delta y/\Delta x = .07/.10 \). The exact slope is a number \( c \) (near .7) that we are not yet prepared to reveal.

The special property of \( y = 2^x \) is that the slope at all points is \( cy \). The slope is proportional to the function. The exponential solves \( dy/dx = cy \).

Now look at the inverse function—the logarithm. Its graph is the mirror image:

\[
\text{If } y = 2^x \text{ then } x = \log_2 y. \text{ If } 2^{1/10} \approx 1.07 \text{ then } \log_2 1.07 \approx 1/10
\]

What the exponential does, the logarithm undoes—and vice versa. The logarithm of \( 2^x \) is the exponent \( x \). Since the exponential starts with slope \( c \), the logarithm must start with slope \( 1/c \). Check that numerically. The logarithm of 1.07 is near 1/10. The slope is near .10/.07. The beautiful property is that \( dx/dy = 1/cy \).

\[
\begin{align*}
\text{Fig. 4.7} & \quad \text{Piecewise linear models and smooth curves: } y = 2^x \text{ and } x = \log_2 y. \text{ Base } b = 2.
\end{align*}
\]

I have to mention that calculus avoids logarithms to base 2. The reason lies in that mysterious number \( c \). It is the “natural logarithm” of 2, which is \( .693147 \ldots \)—and who wants that? Also \( 1/6.693147 \ldots \) enters the slope of \( \log_2 y \). Then \( (dx/dy)(dy/dx) = 1 \). The right choice is to use “natural logarithms” throughout. In place of 2, they are based on the special number \( e \):

\[
y = e^x \text{ is the inverse of } x = \ln y.
\]

The derivatives of those functions are sensational—they are saved for Chapter 6. Together with \( x^n \) and \( \sin x \) and \( \cos x \), they are the backbone of calculus.

\[\text{Note} \quad \text{It is almost possible to go directly to Chapter 6. The inverse functions } x = \sin^{-1} y \text{ and } x = \tan^{-1} y \text{ can be done quickly. The reason for including integrals first (Chapter 5) is that they solve differential equations with no guesswork:}
\]

\[
\frac{dy}{dx} = y \quad \text{or} \quad \frac{dx}{dy} = \frac{1}{y} \quad \text{leads to} \quad \int dx = \int \frac{dy}{y} \quad \text{or} \quad x = \ln y + C.
\]

Integrals have applications of all kinds, spread through the rest of the book. But do not lose sight of \( 2^x \) and \( e^x \). They solve \( dy/dx = cy \)—the key to applied calculus.

\section*{The Inverse of a Chain \( h(g(x)) \)}

The functions \( g(x) = x - 2 \) and \( h(y) = 3y \) were easy to invert. For \( g^{-1} \) we added 2, and for \( h^{-1} \) we divided by 3. Now the question is: If we create the composite function \( z = h(g(x)) \), or \( z = 3(x - 2) \), what is its inverse?
4.3 Inverse Functions and Their Derivatives

Virtually all known functions are created in this way, from chains of simpler functions. The problem is to invert a chain using the inverse of each piece. The answer is one of the fundamental rules of mathematics:

\[ z = h(g(x)) \]

The inverse of \( z = h(g(x)) \) is a chain of inverses in the opposite order:

\[ x = g^{-1}(h^{-1}(z)) \]

(3)

That last equation looks like a mess, but it holds the key. In the middle you see \( h^{-1} \) and \( h \). That part of the chain does nothing! The inverse functions cancel, to leave \( g^{-1}(h(g(x))) = x \).

**EXAMPLE 7** \( z = h(g(x)) = 3(x - 2) \) and \( x = g^{-1}(h^{-1}(z)) = \frac{1}{3}z + 2 \).

First \( h^{-1} \) divides by 3. Then \( g^{-1} \) adds 2. The inverse of \( h \circ g \) is \( g^{-1} \circ h^{-1} \). It can be found directly by solving \( z = 3(x - 2) \). A chain of inverses is like writing in prose—we do it without knowing it.

**EXAMPLE 8** Invert \( z = \sqrt{x - 2} \) by writing \( z^2 = x - 2 \) and then \( x = z^2 + 2 \).

The inverse adds 2 and takes the square—but not in that order. That would give \( (z + 2)^2 \), which is wrong. The correct order is \( z^2 + 2 \).

The domains and ranges are explained by Figure 4.8. We start with \( x \geq 2 \). Subtracting 2 gives \( y \geq 0 \). Taking the square root gives \( z \geq 0 \). Taking the square brings back \( y \geq 0 \). Adding 2 brings back \( x \geq 2 \)—which is in the original domain of \( g \).

**EXAMPLE 9** Inverse matrices \((AB)^{-1} = B^{-1}A^{-1}\) (this linear algebra is optional).

Suppose a vector \( x \) is multiplied by a square matrix \( B \): \( y = g(x) = Bx \). The inverse function multiplies by the inverse matrix: \( x = g^{-1}(y) = B^{-1}y \). It is like multiplication by \( B = 3 \) and \( B^{-1} = \frac{1}{3} \), except that \( x \) and \( y \) are vectors.

Now suppose a second function multiplies by another matrix \( A \): \( z = h(g(x)) = ABx \). The problem is to recover \( x \) from \( z \). The first step is to invert \( A \), because that came last: \( Bx = A^{-1}z \). Then the second step multiplies by \( B^{-1} \) and brings back \( x = B^{-1}A^{-1}z \). The product \( B^{-1}A^{-1} \) inverts the product \( AB \). The rule for matrix inverses is like the rule for function inverses—in fact it is a special case.

I had better not wander too far from calculus. The next section introduces the inverses of the sine and cosine and tangent, and finds their derivatives. Remember that the ultimate source is the chain rule.
4 Derivatives by the Chain Rule

4.3 Exercises

Read-through questions

The functions \( g(x) = x - 4 \) and \( f(y) = y + 4 \) are \( \text{a} \) functions, because \( f(g(x)) = \text{b} \). Also \( g(f(y)) = \text{c} \). The notation is \( g = g^{-1} \) and \( g = \text{d} \). The composition \( \text{e} \) is the identity function. By definition \( x = g^{-1}(y) \) if and only if \( y = \text{f} \). When \( y \) is in the range of \( g \), it is in the \( \text{g} \) of \( g^{-1} \). Similarly \( x \) is in the \( \text{h} \) of \( g \) when it is in the \( \text{i} \) of \( g^{-1} \). If \( g \) has an inverse then \( \text{j} \) \( g(x) \) by \( \text{k} \) \( g^{-1}(y) \) at any two points. The function \( g \) must be steadily \( \text{l} \) or steadily \( \text{m} \).

The chain rule applied to \( f(g(x)) = x \) gives \( (df/dy)(dy/dx) = \text{n} \). The slope of \( g^{-1} \) times the slope of \( g \) equals \( \text{o} \). More directly \( dx/dy = 1/\istration \). For \( y = 2x + 1 \) and \( x = (y - 1)/2 \), the slopes are \( dy/dx = \text{p} \) and \( dx/dy = \text{q} \). For \( y = x^2 \) and \( x = \text{r} \), the slopes are \( dy/dx = \text{s} \) and \( dx/dy = \text{t} \). Substituting \( x^2 \) for \( y \) gives \( dx/dy = \text{u} \). Then \( (dx/dy)(dy/dx) = \text{v} \).

The graph of \( y = g(x) \) is also the graph of \( x = \text{w} \), but with \( x \) across and \( y \) up. For an ordinary graph of \( g^{-1} \), take the reflection in the line \( \text{x} \). If \( (3, 8) \) is on the graph of \( g \), then its mirror image \( \text{y} \) is on the graph of \( g^{-1} \). Those particular points satisfy \( 8 = 2^3 \) and \( 3 = \text{z} \).

The inverse of the chain \( z = h(g(x)) \) is the chain \( x = \text{z} \). If \( g(x) = 3x \) and \( h(y) = y^3 \), then \( z = \text{y} \). Its inverse is \( x = \text{z} \), which is the composition of \( \text{t} \) and \( \text{f} \).

Solve equations 1–10 for \( x \), to find the inverse function \( x = g^{-1}(y) \). When more than one \( x \) gives the same \( y \), write "no inverse."

1. \( y = 3x - 6 \) \( \text{a} \)
2. \( y = Ax + b \)
3. \( y = x^2 - 1 \) \( \text{b} \)
4. \( y = x/(x - 1) \) \( \text{c} \)
5. \( y = 1 + x^{-1} \) \( \text{d} \)
6. \( y = |x| \)
7. \( y = x^3 - 1 \) \( \text{e} \)
8. \( y = 2x + |x| \)
9. \( y = \sin x \) \( \text{f} \)
10. \( y = x^{1/5} \) \( \text{g} \)

Solving \( y = \frac{1}{x-a} \) gives \( xy - ay = 1 \) or \( x = \frac{1+ay}{y} \). Now solve that equation for \( y \).

Solving \( y = \frac{x+1}{x-1} \) gives \( xy - y = x + 1 \) or \( x = \frac{y+1}{y-1} \). Draw the graph to see why \( f \) and \( f^{-1} \) are the same. Compute \( dy/dx \) and \( dx/dy \).

Suppose \( f \) is increasing and \( f(2) = 3 \) and \( f(3) = 5 \). What can you say about \( f^{-1}(4) \)?

Suppose \( f(2) = 3 \) and \( f(3) = 5 \) and \( f(5) = 5 \). What can you say about \( f^{-1} \)?

Suppose \( f(2) = 3 \) and \( f(3) = 5 \) and \( f(5) = 0 \). How do you know that there is no function \( f^{-1} \)?

**Vertical line test:** If no vertical line touches its graph twice then \( f(x) \) is a function (one \( y \) for each \( x \)). **Horizontal line test:** If no horizontal line touches its graph twice then \( f(x) \) is invertible because ________.

If \( f(x) \) and \( g(x) \) are increasing, which two of these might not be increasing?

\( f(x) + g(x) \) \( f(x)g(x) \) \( f(x)/g(x) \) \( f^{-1}(x) \) \( 1/f(x) \)

If \( y = 1/x \) then \( x = 1/y \). If \( y = 1 - x \) then \( x = 1 - y \). The graphs are their own mirror images in the 45° line. Construct two more functions with this property \( f = f^{-1} \) or \( f(f(x)) = x \).

For which numbers \( m \) are these functions invertible?

(a) \( y = mx + b \) \( \text{b} \)
(b) \( y = mx + x^3 \) \( \text{c} \)
(c) \( y = mx + \sin x \)

From its graph show that \( y = |x| + cx \) is invertible if \( c > 1 \) and also if \( c < -1 \). The inverse of a piecewise linear function is piecewise ________.

In 21–26 find \( dy/dx \) in terms of \( x \) and \( dx/dy \) in terms of \( y \).

21. \( y = x^5 \)
22. \( y = 1/(x - 1) \)
23. \( y = x^3 - 1 \)
24. \( y = 1/x^3 \)
25. \( y = x \)
26. \( y = \frac{ax + b}{cx + d} \)
27. If \( dy/dx = 1/y \) then \( dx/dy = \text{_____} \) and \( x = \text{_____} \).
28. If \( dx/dy = 1/y \) then \( dy/dx = \text{_____} \). (These functions are \( y = e^x \) and \( x = \ln y \), soon to be honored property).
29. The slopes of \( f(x) = \frac{1}{2}x^3 \) and \( g(x) = -1/x \) are \( x^2 \) and \( 1/x^2 \). Why isn't \( f = g^{-1} \)? What is \( g^{-1} \)? Show that \( g(g^{-1}) = 1 \).
30. At the points \( x_1, x_2, x_3 \) a piecewise constant function jumps to \( y_1, y_2, y_3 \). Draw its graph starting from \( y(0) = 0 \). The mirror image is piecewise constant with jumps at the points \( \text{_____} \) to the heights \( \text{_____} \). Why isn't this the inverse function?

In 31–38 draw the graph of \( y = g(x) \). Separately draw its mirror image \( x = g^{-1}(y) \).

31. \( y = 5x - 10 \)
32. \( y = \cos x, 0 \leq x \leq \pi \)
33. \( y = 1/(x + 1) \)
34. \( y = |x| - 2x \)
35. \( y = 10^x \)
36. \( y = \frac{1}{\sqrt{1-x^2}}, 0 \leq x \leq 1 \)
37. \( y = 2^{-x} \)
38. \( y = 1/\sqrt{1-x^2}, 0 \leq x < 1 \)

In 39–42 find \( dx/dy \) at the given point.

39. \( y = \sin x \) at \( x = \pi/6 \)
40. \( y = \tan x \) at \( x = \pi/4 \)
41. \( y = \sin x^2 \) at \( x = 3 \)
42. \( y = x - \sin x \) at \( x = 0 \)
43 If \( y \) is a decreasing function of \( x \), then \( x \) is a ______ function of \( y \). Prove by graphs and by the chain rule.

44 If \( f(x) > x \) for all \( x \), show that \( f^{-1}(y) < y \).

45 True or false, with example:
   (a) If \( f(x) \) is invertible so is \( h(x) = (f(x))^2 \).
   (b) If \( f(x) \) is invertible so is \( h(x) = f(f(x)) \).
   (c) \( f^{-1}(y) \) has a derivative at every \( y \).

In the chains 46–51 write down \( g(x) \) and \( f(y) \) and their inverses.
Then find \( x = g^{-1}(f^{-1}(z)) \).

46 \( z = 5(x - 4) \)
47 \( z = (x^2)^2 \)
48 \( z = (6 + x)^3 \)
49 \( z = 6 + x^3 \)
50 \( z = \sqrt{2x + 4} + 4 \)
51 \( z = \log(10^y) \)

52 Solving \( f(x) = 0 \) is a large part of applied mathematics.
Express the solution \( x^* \) in terms of \( f^{-1} \): \( x^* = \) ________.

53 (a) Show by example that \( d^2x/dy^2 \) is not \( 1/(d^2y/dx^2) \).
(b) If \( y \) is in meters and \( x \) is in seconds, then \( d^2y/dx^2 \) is in ______ and \( d^2x/dy^2 \) is in ______.

54 Newton’s method solves \( f(x^*) = 0 \) by applying a linear approximation to \( f^{-1} \):
\[
f^{-1}(0) \approx f^{-1}(y) + (df^{-1}/dy)(0 - y).
\]
For \( y = f(x) \) this is Newton’s equation \( x^* \approx x + \) ________.

55 If the demand is \( 1/(p + 1)^2 \) when the price is \( p \), then the demand is \( y \) when the price is ________. If the range of prices is \( p \geq 0 \), what is the range of demands?

56 If \( dF/dx = f(x) \) show that the derivative of \( G(y) = yf^{-1}(y) - F(f^{-1}(y)) \) is \( f^{-1}(y) \).

57 For each number \( y \) find the maximum value of \( yx - 2x^2 \).
This maximum is a function \( G(y) \). Verify that the derivatives of \( G(y) \) and \( 2x^2 \) are inverse functions.

58 (for professors only) If \( G(y) \) is the maximum value of \( yx - F(x) \), prove that \( F(x) \) is the maximum value of \( xy - G(y) \).
Assume that \( f(x) = dF/dx \) is increasing, like \( 8x^3 \) in Problem 57.

59 Suppose the richest \( x \) percent of people in the world have \( 10\sqrt{x} \) percent of the wealth. Then \( y \) percent of the wealth is held by ________ percent of the people.

4.4 Inverses of Trigonometric Functions

Mathematics is built on basic functions like the sine, and on basic ideas like the inverse. Therefore \textit{it is totally natural to invert the sine function.} The graph of \( x = \sin^{-1}y \) is a mirror image of \( y = \sin x \). This is a case where we pay close attention to the domains, since the sine goes up and down infinitely often. We only want one piece of that curve, in Figure 4.9.

For the bold line the domain is restricted. The angle \( x \) lies between \(-\pi/2 \) and \( +\pi/2 \). On that interval the sine is increasing, so each \( y \) comes from exactly one angle \( x \). If the whole sine curve is allowed, infinitely many angles would have \( \sin x = 0 \). The sine

\[ y = \sin x \]
\[ \frac{dy}{dx} = \cos x \]
\[ x = \pi \]
\[ \frac{dx}{dy} = \frac{1}{\sqrt{1 - y^2}} \]
\[ y = \sin^{-1} x \]
\[ \frac{dy}{dx} = \frac{1}{\sqrt{1 - x^2}} \]

**Fig. 4.9** Graphs of \( \sin x \) and \( \sin^{-1} y \). Their slopes are \( \cos x \) and \( 1/\sqrt{1 - y^2} \).
function could not have an inverse. By restricting to an interval where \( \sin x \) is increasing, we make the function invertible.

The inverse function brings \( y \) back to \( x \). It is \( x = \sin^{-1} y \) (the inverse sine):

\[
x = \sin^{-1} y \text{ when } y = \sin x \text{ and } |x| \leq \pi/2.
\]

(1)

The inverse starts with a number \( y \) between \(-1\) and \(1\). It produces an angle \( x = \sin^{-1}y \)—the angle whose sine is \( y \). The angle \( x \) is between \(-\pi/2\) and \(\pi/2\), with the required sine. Historically \( x \) was called the “arc sine” of \( y \), and \( \arcsin \) is used in computing. The mathematical notation is \( \sin^{-1} \). This has nothing to do with \( 1/\sin x \).

The figure shows the \( 30^\circ \) angle \( x = \pi/6 \). Its sine is \( y = 1/2 \). The inverse sine of \( y \) is \( x \). Again: The symbol \( \sin^{-1}(1) \) stands for the angle whose sine is \( 1 \) (this angle is \( x = \pi/2 \)). We are seeing \( g^{-1}(g(x)) = x \):

\[
\sin^{-1}(\sin x) = x \text{ for } -\pi/2 \leq x \leq \pi/2 \quad \text{and} \quad \sin(\sin^{-1}y) = y \text{ for } -1 \leq y \leq 1.
\]

EXAMPLE 1 (important) If \( \sin x = y \) find a formula for \( \cos x \).

Solution We are given the sine, we want the cosine. The key to this problem must be \( \cos^2 x = 1 - \sin^2 x \). When the sine is \( y \), the cosine is the square root of \( 1 - y^2 \):

\[
\cos x = \cos(\sin^{-1}y) = \sqrt{1 - y^2}.
\]

(2)

This formula is crucial for computing derivatives. We use it immediately.

THE DERIVATIVE OF THE INVERSE SINE

The calculus problem is to find the slope of the inverse function \( f(y) = \sin^{-1} y \). The chain rule gives (slope of inverse function) = \( 1/(\text{slope of original function}) \). Certainly the slope of \( \sin x \) is \( \cos x \). To switch from \( x \) to \( y \), use equation (2):

\[
y = \sin x \quad \text{gives} \quad \frac{dy}{dx} = \cos x \quad \text{so that} \quad \frac{dx}{dy} = \frac{1}{\cos x} = \frac{1}{\sqrt{1 - y^2}}.
\]

(3)

This derivative \( 1/\sqrt{1 - y^2} \) gives a new \( u-f \) pair that is extremely valuable in calculus:

velocity \( v(t) = 1/\sqrt{1 - t^2} \) \quad \text{distance} \quad f(t) = \sin^{-1} t.

Inverse functions will soon produce two more pairs, from the derivatives of \( \tan^{-1} y \) and \( \sec^{-1} y \). The table at the end lists all the essential facts.

EXAMPLE 2 The slope of \( \sin^{-1} y \) at \( y = 1 \) is infinite: \( 1/\sqrt{1 - y^2} = 1/0 \). Explain.

At \( y = 1 \) the graph of \( y = \sin x \) is horizontal. The slope is zero. So its mirror image is vertical. The slope \( 1/0 \) is an extreme case of the chain rule.

Question What is \( d/dx (\sin^{-1} x) \)? Answer \( 1/\sqrt{1 - x^2} \). I just changed letters.

THE INVERSE COSINE AND ITS DERIVATIVE

Whatever is done for the sine can be done for the cosine. But the domain and range have to be watched. The graph cannot be allowed to go up and down. Each \( y \) from \(-1\) to \(1\) should be the cosine of only one angle \( x \). That puts \( x \) between \(0\) and \( \pi \). Then the cosine is steadily decreasing and \( y = \cos x \) has an inverse:

\[
\cos^{-1}(\cos x) = x \text{ and } \cos(\cos^{-1}y) = y.
\]

(4)
4.4 Inverses of Trigonometric Functions

The cosine of the angle \( x = 0 \) is the number \( y = 1 \). The inverse cosine of \( y = 1 \) is the angle \( x = 0 \). Those both express the same fact, that \( \cos 0 = 1 \).

For the slope of \( \cos^{-1}y \), we could copy the calculation that succeeded for \( \sin^{-1}y \). The chain rule could be applied as in (3). But there is a faster way, because of a special relation between \( \cos^{-1}y \) and \( \sin^{-1}y \). \textit{Those angles always add to a right angle}:

\[
\cos^{-1}y + \sin^{-1}y = \pi/2. \tag{5}
\]

Figure 4.9c shows the angles and Figure 4.10c shows the graphs. The sum is \( \pi/2 \) (the dotted line), and its derivative is zero. So the derivatives of \( \cos^{-1}y \) and \( \sin^{-1}y \) must add to zero. \textit{Those derivatives have opposite sign.} There is a \textit{minus} for the inverse cosine, and its graph goes downward:

\begin{align*}
\text{The derivative of } x &= \cos^{-1}y \text{ is } \frac{dx}{dy} = -1/\sqrt{1-y^2}. \tag{6}
\end{align*}

\begin{figure}[h]
\centering
\includegraphics[width=\textwidth]{fig410.png}
\caption{The graphs of \( y = \cos x \) and \( x = \cos^{-1}y \). Notice the domain \( 0 \leq x \leq \pi \).
\end{figure}

\textbf{Question} How can two functions \( x = \sin^{-1}y \) and \( x = -\cos^{-1}y \) have the same derivative?

\textbf{Answer} \( \sin^{-1}y \) must be the same as \(-\cos^{-1}y + C \). Equation (5) gives \( C = \pi/2 \).

\begin{figure}[h]
\centering
\includegraphics[width=\textwidth]{fig411.png}
\caption{The graphs of \( y = \tan x \) and \( x = \tan^{-1}y \).
\end{figure}

\textbf{THE INVERSE TANGENT AND ITS DERIVATIVE}

The tangent is \( \sin x/\cos x \). The inverse tangent is \textit{not} \( \sin^{-1}y/\cos^{-1}y \). The inverse function produces \textit{the angle whose tangent is} \( y \). Figure 4.11 shows that angle, which is between \(-\pi/2 \) and \( \pi/2 \). The tangent can be any number, but the inverse tangent is in the \textit{open interval} \(-\pi/2 < x < \pi/2 \). (The interval is “open” because its endpoints are not included.) The tangents of \( \pi/2 \) and \( -\pi/2 \) are not defined.

The slope of \( y = \tan x \) is \( dy/dx = \sec^2x \). What is the slope of \( x = \tan^{-1}y \)?

By the chain rule
\[
\frac{dx}{dy} = \frac{1}{\sec^2x} = \frac{1}{1 + \tan^2x} = \frac{1}{1 + y^2}. \tag{7}
\]

\textbf{4E} The derivative of \( f(y) = \tan^{-1}y \) is \[
\frac{df}{dy} = \frac{1}{1 + y^2}. \tag{8}
\]
Derivatives by the Chain Rule

slope = \frac{1}{1 + y^2}

slope = \frac{1}{2}

slope = \frac{1}{2}

Fig. 4.11  

$y = \tan^{-1}y$ has slope $1/(1 + y^2)$. $y = \sec^{-1}y$ has slope $1/|y|\sqrt{y^2 - 1}$.

**EXAMPLE 3**  
The tangent of $x = \pi/4$ is $y = 1$. We check slopes. On the inverse tangent curve, $dx/dy = 1/(1 + y^2) = \frac{1}{2}$. On the tangent curve, $dy/dx = \sec^2x$. At $\pi/4$ the secant squared equals 2. The slopes $dx/dy = \frac{1}{2}$ and $dy/dx = 2$ multiply to give 1.

**Important**  
Soon will come the following question. *What function has the derivative $1/(1 + x^2)$?* One reason for reading this section is to learn the answer. The function is in equation (8)—if we change letters. It is $f(x) = \tan^{-1}x$ that has slope $1/(1 + x^2)$.

**Fig. 4.12**  

$\cos^2x + \sin^2x = 1$ and $1 + \tan^2x = \sec^2x$ and $1 + \cot^2x = \csc^2x$.

**INVERSE COTANGENT, INVERSE SECANT, INVERSE COSECANT**

There is no way we can avoid completing this miserable list! But it can be painless. The idea is to use $1/(dy/dx)$ for $y = \cot x$ and $y = \sec x$ and $y = \csc x$:

$$
\frac{dx}{dy} = -\frac{1}{\csc^2x} \quad \text{and} \quad \frac{dx}{dy} = \frac{1}{\sec x \tan x} \quad \text{and} \quad \frac{dx}{dy} = -\frac{1}{\csc x \cot x}.
$$

In the middle equation, replace $\sec x$ by $y$ and $\tan x$ by $\pm\sqrt{y^2 - 1}$. Choose the sign for positive slope (compare Figure 4.11). That gives the middle equation in (10):

The derivatives of $\cot^{-1}y$ and $\sec^{-1}y$ and $\csc^{-1}y$ are

$$
\frac{d}{dy}(\cot^{-1}y) = -\frac{1}{1 + y^2} \quad \frac{d}{dy}(\sec^{-1}y) = \frac{1}{|y|\sqrt{y^2 - 1}} \quad \frac{d}{dy}(\csc^{-1}y) = -\frac{1}{|y|\sqrt{y^2 - 1}}.
$$

**Note about the inverse secant**  
When $y$ is negative there is a choice for $x = \sec^{-1}y$. We selected the angle in the second quadrant (between $\pi/2$ and $\pi$). Its cosine is negative, so its secant is negative. This choice makes $\sec^{-1}y = \cos^{-1}(1/y)$, which matches $\sec x = 1/\cos x$. It also makes $\sec^{-1}y$ an increasing function, where $\cos^{-1}y$ is a decreasing function. So we needed the absolute value $|y|$ in the derivative.
4.4 Inverses of Trigonometric Functions

Some mathematical tables make a different choice. The angle $x$ could be in the third quadrant (between $-\pi$ and $-\pi/2$). Then the slope omits the absolute value.

**Summary** For the six inverse functions it is only necessary to learn three derivatives. The other three just have minus signs, as we saw for $\sin^{-1}y$ and $\cos^{-1}y$. Each inverse function and its "cofunction" add to $\pi/2$, so their derivatives add to zero. Here are the six functions for quick reference, with the three new derivatives.

<table>
<thead>
<tr>
<th>Function $f(y)$</th>
<th>Inputs $y$</th>
<th>Outputs $x$</th>
<th>Slope $dx/dy$</th>
</tr>
</thead>
<tbody>
<tr>
<td>$\sin^{-1}y$, $\cos^{-1}y$</td>
<td>$</td>
<td>y</td>
<td>\leq 1$</td>
</tr>
<tr>
<td>$\tan^{-1}y$, $\cot^{-1}y$</td>
<td>All $y$</td>
<td>$\left(-\frac{\pi}{2}, \frac{\pi}{2}\right)$</td>
<td>$(0, \pi)$</td>
</tr>
<tr>
<td>$\sec^{-1}y$, $\csc^{-1}y$</td>
<td>$</td>
<td>y</td>
<td>\geq 1$</td>
</tr>
</tbody>
</table>

If $y = \cos x$ or $y = \sin x$ then $|y| \leq 1$. For $y = \sec x$ and $y = \csc x$ the opposite is true; we must have $|y| \geq 1$. The graph of $\sec^{-1}y$ misses all the points $-1 < y < 1$. Also, that graph misses $x = \pi/2$—where the cosine is zero. The secant of $\pi/2$ would be $1/0$ (impossible). Similarly $\csc^{-1}y$ misses $x = 0$, because $y = \csc 0$ cannot be $1/\sin 0$.

The asterisks in the table are to remove those points $x = \pi/2$ and $x = 0$.

The column of derivatives is what we need and use in calculus.

### 4.4 Exercises

#### Read-through questions

The relation $x = \sin^{-1}y$ means that ___ is the sine of ___ . Thus $x$ is the angle whose sine is ___. The number $y$ lies between ___ and ___. The angle $x$ lies between ___ and ___. (If we want the inverse to exist, there cannot be two angles with the same sine.) The cosine of the angle $\sin^{-1}y$ is ___. The derivative of $x = \sin^{-1}y$ is $dx/dy = ___$.

The relation $x = \cos^{-1}y$ means that $y$ equals ___. Again the number $y$ lies between ___ and ___. This time the angle $x$ lies between ___ and ___. The sum $\sin^{-1}y + \cos^{-1}y = ___$.

The relation $x = \tan^{-1}y$ means that $y$ is the angle whose tangent is ___. The number $y$ lies between ___ and ___. The angle $x$ lies between ___ and ___. The derivative is $dx/dy = ___$.

The relation $x = \sec^{-1}y$ means that ___. The number $y$ never lies between ___ and ___. The angle $x$ lies between ___ and ___, but never at $x = ___$. The derivative of $x = \sec^{-1}y$ is $dx/dy = ___$.

The relation $x = \cot^{-1}y$ means that ___. The number $y$ lies between ___ and ___. The angle $x$ lies between ___ and ___. The sum $\tan^{-1}y + \cot^{-1}y = ___$.

The relation $x = \tan^{-1}y$ means that ___. The number $y$ lies between ___ and ___. The angle $x$ lies between ___ and ___. The derivative is $dx/dy = ___$.

#### In 1-4, find the angles $\sin^{-1}y$ and $\cos^{-1}y$ and $\tan^{-1}y$ in radians.

1. $y = 0$  
2. $y = -1$  
3. $y = 1$  
4. $y = \sqrt{3}$

5. We know that $\sin \pi = 0$. Why isn't $\pi = \sin^{-1}0$?

6. Suppose $\sin x = y$. Under what restriction is $x = \sin^{-1}y$?

7. Sketch the graph of $x = \sin^{-1}y$ and locate the points with slope $dx/dy = 2$.

8. Find $dx/dy$ if $x = \sin^{-1}y$. Draw the graph.

9. If $y = \cos x$ find a formula for $\sin x$. First draw a right triangle with angle $x$ and near side $y$—what are the other two sides?

10. If $y = \sin x$ find a formula for $\tan x$. First draw a right triangle with angle $x$ and far side $y$—what are the other sides?

11. Take the $x$ derivative of $\sin^{-1}(\sin x) = x$ by the chain rule. Check that $d(\sin^{-1}y)/dy = 1/\sqrt{1-y^2}$ gives a correct result.

12. Take the $y$ derivative of $\cos(\cos^{-1}y) = y$ by the chain rule. Check that $d(\cos^{-1}y)/dy = -1/\sqrt{1-y^2}$ gives a correct result.

13. At $y = 0$ and $y = 1$, find the slope $dx/dy$ of $x = \sin^{-1}y$ and $x = \cos^{-1}y$ and $x = \tan^{-1}y$.

14. At $x = 0$ and $x = 1$, find the slope $dx/dy$ of $x = \sin^{-1}y$ and $x = \cos^{-1}y$ and $x = \tan^{-1}y$. 
15 True or false, with reason:
(a) \((\sin^{-1}y)^2 + (\cos^{-1}y)^2 = 1\)
(b) \(\sin^{-1}y = \cos^{-1}y\) has no solution
(c) \(\sin^{-1}y\) is an increasing function
(d) \(\sin^{-1}y\) is an odd function
(e) \(\sin^{-1}y\) and \(-\cos^{-1}y\) have the same slope—so they are the same.
(f) \(\sin(\cos x) = \cos(\sin x)\)

16 Find \(\tan(\cos^{-1}(\sin x))\) by drawing a triangle with sides \(\sin x, \cos x, 1\).

Compute the derivatives in 17–28 (using the letters as given).

17 \(u = \sin^{-1}x\)
18 \(u = \tan^{-1}2x\)
19 \(z = \sin^{-1}(\sin 3x)\)
20 \(z = \sin^{-1}(\cos x)\)
21 \(z = (\sin^{-1}x)^2\)
22 \(z = (\sin^{-1}x)^{-1}\)
23 \(z = \sqrt{1 - y^2}\sin^{-1}y\)
24 \(z = (1 + x^2)\tan^{-1}x\)
25 \(x = \sec^{-1}(y + 1)\)
26 \(u = \sec^{-1}(\sec x^2)\)
27 \(u = \sin^{-1}y/\cos^{-1}\sqrt{1 - y^2}\)
28 \(u = \sin^{-1}y + \cos^{-1}y + \tan^{-1}y\)

29 Draw a right triangle to show why \(\tan^{-1}y + \cot^{-1}y = \pi/2\).
30 Draw a right triangle to show why \(\tan^{-1}y = \cot^{-1}(1/y)\).
31 If \(y = \tan x\) find \(\sec x\) in terms of \(y\).
32 Draw the graphs of \(y = \cot x\) and \(x = \cot^{-1}y\).
33 Find the slope \(dx/dy\) of \(x = \tan^{-1}y\) at
   (a) \(y = -3\)
   (b) \(x = 0\)
   (c) \(x = -\pi/4\)

34 Find a function \(u(t)\) whose slope satisfies \(u' + t^2u' = 1\).
35 What is the second derivative \(d^2x/dy^2\) of \(x = \sin^{-1}y\)?
36 What is \(d^2u/dy^2\) for \(u = \tan^{-1}y\)?

Find the derivatives in 37–44.

37 \(y = \sec \frac{1}{2}x\)
38 \(x = \sec^{-1}2y\)
39 \(u = \sec^{-1}(x^n)\)
40 \(u = \sec^{-1}(\tan x)\)
41 \(\tan y = (x - 1)/(x + 1)\)
42 \(z = (\sin x)(\sin^{-1}x)\)
43 \(y = \sec^{-1}\sqrt{x^2 + 1}\)
44 \(z = \sin(\cos^{-1}x) - \cos(\sin^{-1}x)\)

45 Differentiate \(\cos^{-1}(1/y)\) to find the slope of \(\sec^{-1}y\) in a new way.

46 The domain and range of \(x = \csc^{-1}y\) are ________.
47 Find a function \(u(y)\) such that \(du/dy = 4/(1 - y^2)\).
48 Solve the differential equation \(du/dx = 1/(1 + 4x^2)\).
49 If \(du/dx = 2/(\sqrt{1 - x^2})\) find \(u(1) - u(0)\).

50 (recommended) With \(u(x) = (x - 1)/(x + 1)\), find the derivative of \(\tan^{-1}u(x)\). This is also the derivative of _________. So the difference between the two functions is _________.

51 Find \(u(x)\) and \(\tan^{-1}u(x)\) and \(\tan^{-1}x\) at \(x = 0\) and \(x = \infty\). Conclusion based on Problem 50: \(\tan^{-1}u(x) - \tan^{-1}x\) equals the number _________.

52 Find \(u(x)\) and \(\tan^{-1}u(x)\) and \(\tan^{-1}x\) as \(x \to -\infty\). Now \(\tan^{-1}u(x) - \tan^{-1}x\) equals _________. Something has happened to \(\tan^{-1}u(x)\). At what \(x\) do \(u(x)\) and \(\tan^{-1}u(x)\) change instantly?
This chapter is about the idea of integration, and also about the technique of integration. We explain how it is done in principle, and then how it is done in practice. Integration is a problem of adding up infinitely many things, each of which is infinitely small. Doing the addition is not recommended. The whole point of calculus is to offer a better way.

The problem of integration is to find a limit of sums. The key is to work backward from a limit of differences (which is the derivative). We can integrate \( v(x) \) if it turns up as the derivative of another function \( f(x) \). The integral of \( v = \cos x \) is \( f = \sin x \). The integral of \( v = x \) is \( f = \frac{1}{2} x^2 \). Basically, \( f(x) \) is an "antiderivative". The list of \( f \)'s will grow much longer (Section 5.4 is crucial). A selection is inside the cover of this book. If we don't find a suitable \( f(x) \), numerical integration can still give an excellent answer.

I could go directly to the formulas for integrals, which allow you to compute areas under the most amazing curves. (Area is the clearest example of adding up infinitely many infinitely thin rectangles, so it always comes first. It is certainly not the only problem that integral calculus can solve.) But I am really unwilling just to write down formulas, and skip over all the ideas. Newton and Leibniz had an absolutely brilliant intuition, and there is no reason why we can't share it.

They started with something simple. We will do the same.

**SUMS AND DIFFERENCES**

Integrals and derivatives can be mostly explained by working (very briefly) with sums and differences. Instead of functions, we have \( n \) ordinary numbers. The key idea is nothing more than a basic fact of algebra. In the limit as \( n \to \infty \), it becomes the basic fact of calculus. The step of "going to the limit" is the essential difference between algebra and calculus! It has to be taken, in order to add up infinitely many infinitesimals—but we start out this side of it.

To see what happens before the limiting step, we need two sets of \( n \) numbers. The first set will be \( v_1, v_2, \ldots, v_n \), where \( v \) suggests velocity. The second set of numbers will be \( f_1, f_2, \ldots, f_n \), where \( f \) recalls the idea of distance. You might think \( d \) would be a better symbol for distance, but that is needed for the \( dx \) and \( dy \) of calculus.
A first example has $n = 4$:

\[ v_1, v_2, v_3, v_4 = 1, 2, 3, 4 \]

\[ f_1, f_2, f_3, f_4 = 1, 3, 6, 10. \]

The relation between the $v$'s and $f$'s is seen in that example. When you are given 1, 3, 6, 10, how do you produce 1, 2, 3, 4? By taking differences. The difference between 10 and 6 is 4. Subtracting 6 - 3 is 3. The difference $f_2 - f_1 = 3 - 1$ is $v_2 = 2$.

Each $v$ is the difference between two $f$'s:

\[ v_j \text{ is the difference } f_j - f_{j-1}. \]

This is the discrete form of the derivative. I admit to a small difficulty at $j = 1$, from the fact that there is no $f_0$. The first $v$ should be $f_1 - f_0$, and the natural idea is to agree that $f_0$ is zero. This need for a starting point will come back to haunt us (or help us) in calculus.

Now look again at those same numbers—but start with $v$. From $v = 1, 2, 3, 4$ how do you produce $f = 1, 3, 6, 10$? By taking sums. The first two $v$'s add to 3, which is $f_2$. The first three $v$'s add to $f_3 = 6$. The sum of all four $v$'s is $1 + 2 + 3 + 4 = 10$. Taking sums is the opposite of taking differences.

That idea from algebra is the key to calculus. The sum $f_j$ involves all the numbers $v_1 + v_2 + \cdots + v_j$. The difference $v_j$ involves only the two numbers $f_j - f_{j-1}$. The fact that one reverses the other is the "Fundamental Theorem." Calculus will change sums to integrals and differences to derivatives—but why not let the key idea come through now?

### 5A Fundamental Theorem of Calculus (before limits):

If each $v_j = f_j - f_{j-1}$, then $v_1 + v_2 + \cdots + v_n = f_n - f_0$.

The differences of the $f$'s add up to $f_n - f_0$. All $f$'s in between are canceled, leaving only the last $f_n$ and the starting $f_0$. The sum "telescopes":

\[ v_1 + v_2 + v_3 + \cdots + v_n = (f_1 - f_0) + (f_2 - f_1) + (f_3 - f_2) + \cdots + (f_n - f_{n-1}). \]

The number $f_1$ is canceled by $-f_1$. Similarly $-f_2$ cancels $f_2$ and $-f_3$ cancels $f_3$. Eventually $f_n$ and $-f_0$ are left. When $f_0$ is zero, the sum is the final $f_n$.

That completes the algebra. We add the $v$'s by finding the $f$'s.

**Question** How do you add the odd numbers 1 + 3 + 5 + \cdots + 99 (the $v$'s)?

**Answer** They are the differences between 0, 1, 4, 9, \ldots. These $f$'s are squares. By the Fundamental Theorem, the sum of 50 odd numbers is $(50)^2$.

The tricky part is to discover the right $f$'s! Their differences must produce the $v$'s. In calculus, the tricky part is to find the right $f(x)$. Its derivative must produce $v(x)$. It is remarkable how often $f$ can be found—more often for integrals than for sums. Our next step is to understand how the integral is a limit of sums.

### SUMS APPROACH INTEGRALS

Suppose you start a successful company. The rate of income is increasing. After $x$ years, the income per year is $\sqrt{x}$ million dollars. In the first four years you reach $\sqrt{1}, \sqrt{2}, \sqrt{3},$ and $\sqrt{4}$ million dollars. Those numbers are displayed in a bar graph (Figure 5.1a, for investors). I realize that most start-up companies make losses, but your company is an exception. If the example is too good to be true, please keep reading.
5.1 The Idea of the Integral

The graph shows four rectangles, of heights $\sqrt{1}$, $\sqrt{2}$, $\sqrt{3}$, $\sqrt{4}$. Since the base of each rectangle is one year, those numbers are also the areas of the rectangles. One investor, possibly weak in arithmetic, asks a simple question: What is the total income for all four years? There are two ways to answer, and I will give both.

The first answer is $\sqrt{1} + \sqrt{2} + \sqrt{3} + \sqrt{4}$. Addition gives 6.15 million dollars. Figure 5.1b shows this total—which is reached at year 4. This is exactly like velocities and distances, but now $v$ is the income per year and $f$ is the total income. Algebraically, $f_4$ is still $v_1 + \cdots + v_4$.

The second answer comes from geometry. The total income is the total area of the rectangles. We are emphasizing the correspondence between addition and area. That point may seem obvious, but it becomes important when a second investor (smarter than the first) asks a harder question.

Here is the problem. The incomes as stated are false. The company did not make a million dollars the first year. After three months, when $x$ was $1/4$, the rate of income was only $\sqrt{x} = 1/2$. The bar graph showed $\sqrt{1} = 1$ for the whole year, but that was an overstatement. The income in three months was not more than $1/2$ times $1/4$, the rate multiplied by the time.

All other quarters and years were also overstated. Figure 5.2a is closer to reality, with 4 years divided into 16 quarters. It gives a new estimate for total income.

Again there are two ways to find the total. We add $\sqrt{1/4} + \sqrt{2/4} + \cdots + \sqrt{16/4}$, remembering to multiply them all by $1/4$ (because each rate applies to $1/4$ year). This is also the area of the 16 rectangles. The area approach is better because the $1/4$ is automatic. Each rectangle has base $1/4$, so that factor enters each area. The total area is now 5.56 million dollars, closer to the truth.

You see what is coming. The next step divides time into weeks. After one week the rate $\sqrt{x}$ is only $\sqrt{1/52}$. That is the height of the first rectangle—its base is $Ax = 1/52$. There is a rectangle for every week. Then a hard-working investor divides time into days, and the base of each rectangle is $Ax = 1/365$. At that point there are $4 \times 365 = 1460$ rectangles, or 1461 because of leap year, with a total area below $5\frac{1}{2}$.
5 Integrals

Total income = area of rectangles

\[ I = \frac{1}{4} \left( \sqrt{\frac{1}{4}} + \sqrt{\frac{2}{4}} + \cdots + \sqrt{\frac{16}{4}} \right) \]

Fig. 5.2 Income = sum of areas (not heights)

million dollars. The calculation is elementary but depressing—adding up thousands of square roots, each multiplied by \( \Delta x \) from the base. There has to be a better way.

The better way, in fact the best way, is calculus. The whole idea is to allow for continuous change. The geometry problem is to find the area under the square root curve. That question cannot be answered by arithmetic, because it involves a limit. The rectangles have base \( \Delta x \) and heights \( \sqrt{\Delta x}, \sqrt{2\Delta x}, \ldots, \sqrt{4} \). There are \( 4/\Delta x \) rectangles—more and more terms from thinner and thinner rectangles. **The area is the limit of the sum as \( \Delta x \to 0 \).**

This limiting area is the "integral." We are looking for a number below 54.

**Algebra (area of \( n \) rectangles):** Compute \( v_1 + \cdots + v_n \) by finding \( f \)'s.

**Key idea:** If \( v_j = f_j - f_{j-1} \) then the sum is \( f_n - f_0 \).

**Calculus (area under curve):** Compute the limit of \( \Delta x [v(\Delta x) + v(2\Delta x) + \cdots] \).

**Key idea:** If \( v(x) = df/dx \) then area = integral to be explained next.

---

**5.1 Exercises**

**Read-through questions**

The problem of summation is to add \( v_1 + \cdots + v_n \). It is solved if we find \( f \)'s such that \( v_j = f_j - f_{j-1} \). Then \( v_1 + \cdots + v_n \) equals \( f_n - f_0 \). The cancellation in \((f_1 - f_0) + (f_2 - f_1) + \cdots + (f_n - f_{n-1})\) leaves only \( f_n - f_0 \). Taking sums is the **d** of taking differences.

The differences between 0, 1, 4, 9 are \( v_1, v_2, v_3 = e \).

For \( f_j - j^2 \) the difference between \( f_{10} \) and \( f_0 \) is \( v_{10} = f \).

From this pattern \( 1 + 3 + 5 + \cdots + 19 \) equals \( g \).

For functions, finding the integral is the reverse of **h**.

If the derivative of \( f(x) \) is \( v(x) \), then the **i** of \( v(x) \) is \( f(x) \).

If \( v(x) = 10x \) then \( f(x) = _j \). This is the **k** of a triangle with base \( x \) and height 10x.

Integrals begin with sums. The triangle under \( v = 10x \) out to \( x = 4 \) has area **l**. It is approximated by four rectangles of heights 10, 20, 30, 40 and area **m**. It is better approximated by eight rectangles of heights **n** and area **o**.

For \( n \) rectangles covering the triangle the area is the sum of **p**. As \( n \to \infty \) this sum should approach the number **q**. That is the integral of \( v = 10x \) from 0 to 4.
5.1 The Idea of the Integral

Problems 1–6 are about sums $f_j$ and differences $v_j$.

1. With $n = 1, 2, 4, 8$, the formula for $v_j$ is ________ (not 2).
   Find $f_1, f_2, f_3, f_4$ starting from $f_0 = 0$. What is $f_2$?

2. The same $c = 1, 2, 4, 8, ...$ are the differences between $f = 1, 2, 4, 8, 16, ...$. Now $f_0 = 1$ and $f_j = 2^j$. (a) Check that $2^5 - 2^3$ equals $v_5$. (b) What is $1 + 2 + 4 + 8 + 16$?

3. The differences between $f = 1, 1/2, 1/4, 1/8$ are $v = -1/2, -1/4, -1/8$. These negative $v$'s do not add up to these positive $f$'s. Verify that $v_1 + v_2 + v_3 = f_4 - f_0$ is still true.

4. Any constant $C$ can be added to the antiderivative $f(x)$ because the _______ of a constant is zero. Any $C$ can be added to $f_0, f_1, ...$ because the _______ between the $f$'s is not changed.

5. Show that $f_j = r^j (r - 1) + f_0$ has $f_j - f_{j-1} = r^{j-1}$. Therefore the geometric series $1 + r + r^2 + ... + r^{j-1}$ adds up to ________ (remember to subtract $f_0$).

6. The sums $f_j = r^j (r - 1)$ also have $f_j - f_{j-1} = r^{j-1}$. Now $f_0 = ________$. Therefore $1 + r + r^2 + ... + r^{j-1}$ adds up to $f_j$. The sum $1 + r + r^2 + ... + r^n$ equals ________.

7. Suppose $v(x) = 3$ for $x < 1$ and $v(x) = 7$ for $x > 1$. Find the area $f(x)$ from 0 to 1, under the graph of $v(x)$. (Two pieces.)

8. If $n = 1, -2, 3, -4, ...$, write down the $f$'s starting from $f_0 = 0$. Find formulas for $v_j$ and $f_j$ when $j$ is odd and $j$ is even.

Problems 9–16 are about the company earning $\sqrt{x}$ per year.

9. When time is divided into weeks there are $4 \times 52 = 208$ rectangles. Write down the first area, the 208th area, and the $j$th area.

10. How do you know that the sum over 208 weeks is smaller than the sum over 16 quarters?

11. A pessimist would use $\sqrt{x}$ at the beginning of each time period as the income rate for that period. Redraw Figure 5.1 (both parts) using heights $\sqrt{0}, \sqrt{1}, \sqrt{2}, \sqrt{3}$. How much lower is the estimate of total income?

12. The same pessimist would redraw Figure 5.1 with heights $0, \sqrt{1/4}, ...$. What is the height of the last rectangle? How much does this change reduce the total rectangular area 5.56?

13. At every step from years to weeks to days to hours, the pessimist's area goes ________ and the optimist's area goes ________. The difference between them is the area of the last ________.

14. The optimist and pessimist arrive at the same limit as years are divided into weeks, days, hours, seconds. Draw the $\sqrt{x}$ curve between the rectangles to show why the pessimist is always too low and the optimist is too high.

15. (Important) Let $f(x)$ be the area under the $\sqrt{x}$ curve, above the interval from 0 to $x$. The area to $x + \Delta x$ is $f(x + \Delta x)$. The extra area is $\Delta f = \sqrt{x}$. This is almost a rectangle with base _______ and height $\sqrt{x}$. So $\Delta f/\Delta x$ is close to _______. As $\Delta x \to 0$ we suspect that $df/dx = _______$. Draw the $\sqrt{x}$ curve from $x = 0$ to $x = 4$ and put triangles below to prove that the area under it is more than 5. Look left and right from the point where $\sqrt{1} = 1$.

Problems 17–22 are about a company whose expense rate $v(x) = 6 - x$ is decreasing.

17. The expenses drop to zero at $x = ________$. The total expense during those years equals ________. This is the area of ________.

18. The rectangles of heights 6, 5, 4, 3, 2, 1 give a total estimated expense of ________. Draw them enclosing the triangle to show why this total is too high.

19. How many rectangles (enclosing the triangle) would you need before their areas are within 1 of the correct triangular area?

20. The accountant uses 2-year intervals and computes $v = 5, 3, 1$ at the midpoints (the odd-numbered years). What is her estimate, how accurate is it, and why?

21. What is the area $f(x)$ under the line $v(x) = 6 - x$ above the interval from 2 to $x$? What is the derivative of this $f(x)$?

22. What is the area $f(x)$ under the line $v(x) = 6 - x$ above the interval from $x$ to 6? What is the derivative of this $f(x)$?

23. With $\Delta x = 1/3$, find the area of the three rectangles that enclose the graph of $v(x) = x^2$.

24. Draw graphs of $v = \sqrt{x}$ and $v = x^2$ from 0 to 1. Which areas add to 1? The same is true for $v = x^3$ and $v = _______$. From $x$ to $x + \Delta x$, the area under $v = x^2$ is $\Delta f$. This is almost a rectangle with base $\Delta x$ and height _______. So $\Delta f/\Delta x$ is close to _______. In the limit we find $df/dx = x^2$ and $f(x) = _______$. Compute the area of 208 rectangles under $v(x) = \sqrt{x}$ from $x = 0$ to $x = 4$. 


5.2 Antiderivatives

The symbol $\int$ was invented by Leibniz to represent the integral. It is a stretched-out $S$, from the Latin word for sum. This symbol is a powerful reminder of the whole construction: *Sum approaches integral, $S$ approaches $\int$, and rectangular area approaches curved area:*

\[
\text{curved area} = \int v(x) \, dx = \int \sqrt{x} \, dx. \tag{1}
\]

The rectangles of base $\Delta x$ lead to this limit—the integral of $\sqrt{x}$. The "$dx$" indicates that $\Delta x$ approaches zero. The heights $v_j$ of the rectangles are the heights $v(x)$ of the curve. The sum of $v_j$ times $\Delta x$ approaches "the integral of $v$ of $x \, dx$." You can imagine an infinitely thin rectangle above every point, instead of ordinary rectangles above special points.

We now find the area under the square root curve. The "limits of integration" are 0 and 4. The lower limit is $x = 0$, where the area begins. (*The start could be any point $x = a$.*) The upper limit is $x = 4$, since we stop after four years. (*The finish could be any point $x = b$.*) The area of the rectangles is a sum of base $\Delta x$ times heights $\sqrt{x}$. The curved area is the limit of this sum. *That limit is the integral of $\sqrt{x}$ from 0 to 4:*

\[
\lim_{\Delta x \to 0} \left[ (\sqrt{\Delta x})(\Delta x) + (\sqrt{2\Delta x})(\Delta x) + \cdots + (\sqrt{4\Delta x})(\Delta x) \right] = \int_{x=0}^{x=4} \sqrt{x} \, dx. \tag{2}
\]

The outstanding problem of integral calculus is still to be solved. *What is this limiting area?* We have a symbol for the answer, involving $\int$ and $\sqrt{x}$ and $dx$—but we don't have a number.

THE ANTIDERIVATIVE

I wish I knew who discovered the area under the graph of $\sqrt{x}$. It may have been Newton. The answer was available earlier, but the key idea was shared by Newton and Leibniz. They understood the parallels between sums and integrals, and between differences and derivatives. I can give the answer, by following that analogy. I can't give the proof (yet)—it is the Fundamental Theorem of Calculus.

In algebra the difference $f_j - f_{j-1}$ is $v_j$. When we add, the sum of the $v$'s is $f_n - f_0$. In calculus the derivative of $f(x)$ is $v(x)$. When we integrate, the *area under the $v(x)$ curve is $f(x)$ minus $f(0)$*. Our problem asks for the area out to $x = 4$:

\[
\text{5B (Discrete vs. continuous, rectangles vs. curved areas, addition vs. integration) } \quad \text{The integral of } v(x) \text{ is the difference in } f(x): \\
\text{If } df/dx = \sqrt{x} \text{ then area } = \int_{x=0}^{x=4} \sqrt{x} \, dx = f(4) - f(0). \tag{3}
\]

*What is $f(x)$?* Instead of the derivative of $\sqrt{x}$, we need its "antiderivative." We have to find a function $f(x)$ whose derivative is $\sqrt{x}$. It is the opposite of Chapters 2–4, and requires us to *work backwards*. The derivative of $x^n$ is $nx^{n-1}$—now we need the antiderivative. The quick formula is $f(x) = x^{n+1}/(n + 1)$—we aim to understand it.

**Solution** Since the derivative lowers the exponent, the antiderivative raises it. We go from $x^{1/2}$ to $x^{3/2}$. But then the derivative is $(3/2)x^{1/2}$. It contains an unwanted factor $3/2$. To cancel that factor, put $2/3$ into the antiderivative:

\[
f(x) = \frac{2}{3} x^{3/2} \text{ has the required derivative } v(x) = x^{1/2} = \sqrt{x}.
\]
Total income = $\frac{2}{3} \cdot 4^{3/2} = \frac{16}{3}$

Rate of income = $\sqrt{x}$

Fig. 5.3 The integral of $v(x) = \sqrt{x}$ is the exact area $16/3$ under the curve.

There you see the key to integrals: Work backward from derivatives (and adjust).

Now comes a number—the exact area. At $x = 4$ we find $x^{3/2} = 8$. Multiply by $2/3$ to get $16/3$. Then subtract $f(0)$:

$$f(x) = \int \sqrt{x} \, dx = \frac{2}{3} \cdot (4)^{3/2} - \frac{2}{3} \cdot (0)^{3/2} = \frac{2}{3} (8) = \frac{16}{3}$$

The total income over four years is $16/3 = 5\frac{1}{3}$ million dollars. This is $f(4) - f(0)$. The sum from thousands of rectangles was slowly approaching this exact area $5\frac{1}{3}$.

Other areas The income in the first year, at $x = 1$, is $\frac{2}{3}(1)^{3/2} = \frac{2}{3}$ million dollars. (The false income was 1 million dollars.) The total income after $x$ years is $\frac{2}{3}x^{3/2}$, which is the antiderivative $f(x)$. The square root curve covers $2/3$ of the overall rectangle it sits in. The rectangle goes out to $x$ and up to $\sqrt{x}$, with area $x^{3/2}$, and $2/3$ of that rectangle is below the curve. (1/3 is above.)

Other antiderivatives The derivative of $x^5$ is $5x^4$. Therefore the antiderivative of $x^4$ is $x^5/5$. Divide by 5 (or $n + 1$) to cancel the 5 (or $n + 1$) from the derivative. And don’t allow $n + 1 = 0$:

The derivative $v(x) = x^n$ has the antiderivative $f(x) = x^{n+1}/(n + 1)$.

EXAMPLE 1 The antiderivative of $x^2$ is $\frac{1}{3}x^3$. This is the area under the parabola $u(x) = x^2$. The area out to $x = 1$ is $\frac{1}{3}(1)^3 - \frac{1}{3}(0)^3$, or $1/3$.

Remark on $\sqrt{x}$ and $x^2$ The $2/3$ from $\sqrt{x}$ and the $1/3$ from $x^2$ add to 1. Those are the areas below and above the $\sqrt{x}$ curve, in the corner of Figure 5.3. If you turn the curve by 90°, it becomes the parabola. The functions $y = \sqrt{x}$ and $x = y^2$ are inverses! The areas for these inverse functions add to a square of area 1.

AREA UNDER A STRAIGHT LINE

You already know the area of a triangle. The region is below the diagonal line $v = x$ in Figure 5.4. The base is 4, the height is 4, and the area is $\frac{1}{2}(4)(4) = 8$. Integration is
Integrals

Exact area = 8

Area under \( v(x) = x \)

\[ f(x) = \int x\,dx = \frac{1}{2}x^2 \]

Fig. 5.4 Triangular area 8 as the limit of rectangular areas 10, 9, 8\(\frac{1}{2}\), ...,

not required! But if you allow calculus to repeat that answer, and build up the integral \( f(x) = \frac{1}{2}x^2 \) as the limiting area of many rectangles, you will have the beginning of something important.

The four rectangles have area \(1 + 2 + 3 + 4 = 10\). That is greater than 8, because the triangle is inside. 10 is a first approximation to the triangular area 8, and to improve it we need more rectangles.

The next rectangles will be thinner, of width \( \Delta x = 1/2 \) instead of the original \( \Delta x = 1 \). There will be eight rectangles instead of four. They extend above the line, so the answer is still too high. The new heights are \(1/2, 1, 3/2, 2, 5/2, 3, 7/2, 4\). The total area in Figure 5.4b is the sum of the base \( \Delta x = 1/2 \) times those heights:

\[
\text{area} = \frac{1}{2}(1 + \frac{1}{2} + 1 + \frac{3}{2} + 4) = 9 \quad \text{(which is closer to 8)}.
\]

**Question** What is the area of 16 rectangles? Their heights are \(1/4, \frac{1}{2}, ... , 4\).

**Answer** With base \( \Delta x = \frac{1}{4} \) the area is \(\frac{1}{2}(1 + \frac{1}{2} + \frac{3}{2} + 4) = 8\frac{1}{2}\).

The effort of doing the addition is increasing. A formula for the sums is needed, and will be established soon. (The next answer would be 8\(\frac{3}{4}\).) But more important than the formula is the idea. We are carrying out a limiting process, one step at a time. The area of the rectangles is approaching the area of the triangle, as \( \Delta x \) decreases. The same limiting process will apply to other areas, in which the region is much more complicated. Therefore we pause to comment on what is important.

**Area Under a Curve**

What requirements are imposed on those thinner and thinner rectangles? It is not essential that they all have the same width. And it is not required that they cover the triangle completely. The rectangles could lie below the curve. The limiting answer will still be 8, even if the widths \( \Delta x \) are unequal and the rectangles fit inside the triangle or across it. We only impose two rules:

1. The largest width \( \Delta x_{\text{max}} \) must approach zero.
2. The top of each rectangle must touch or cross the curve.

The area under the graph is defined to be the limit of these rectangular areas, if that limit exists. For the straight line, the limit does exist and equals 8. That limit is independent of the particular widths and heights—as we absolutely insist it should be.

Section 5.5 allows any continuous \( v(x) \). The question will be the same—*Does the limit exist?* The answer will be the same—*Yes*. That limit will be the integral of \( v(x) \), and it will be the area under the curve. It will be \( f(x) \).
**EXAMPLE 2** The triangular area from 0 to x is \( \frac{1}{2} \) (base)(height) = \( \frac{1}{2} x^2 \). Its derivative is \( v(x) = x \). But notice that \( \frac{1}{2} x^2 + 1 \) has the same derivative. So does \( f = \frac{1}{2} x^2 + C \), for any constant C. There is a "constant of integration" in \( f(x) \), which is wiped out in its derivative \( v(x) \).

**EXAMPLE 3** Suppose the velocity is decreasing: \( v(x) = 4 - x \). If we sample \( v \) at \( x = 1, 2, 3, 4 \), the rectangles lie under the graph. Because \( v \) is decreasing, the right end of each interval gives \( v_{\min} \). Then the rectangular area \( 3 + 2 + 1 + 0 = 6 \) is less than the exact area 8. The rectangles are inside the triangle, and eight rectangles with base \( \frac{1}{2} \) come closer:

\[
\text{rectangular area} = \frac{1}{2} \left( \frac{1}{2} + 3 + \ldots + \frac{1}{2} + 0 \right) = 7.
\]

Sixteen rectangles would have area \( 7\frac{1}{2} \). We repeat that the rectangles need not have the same widths \( \Delta x \), but it makes these calculations easier.

What is the area out to an arbitrary point (like \( x = 3 \) or \( x = 1 \))? We could insert rectangles, but the Fundamental Theorem offers a faster way. Any antiderivative of \( 4 - x \) will give the area. We look for a function whose derivative is \( 4 - x \). The derivative of \( 4x \) is 4, the derivative of \( \frac{1}{2} x^2 \) is \( x \), so work backward:

\[
to \, achieve \, df/dx = 4 - x \, choose \, f(x) = 4x - \frac{1}{2} x^2.
\]

Calculus skips past the rectangles and computes \( f(3) = 7\frac{1}{2} \). The area between \( x = 1 \) and \( x = 3 \) is the difference \( 7\frac{1}{2} - 3\frac{1}{2} = 4 \). In Figure 5.5, this is the area of the trapezoid.

*The f-curve flattens out when the v-curve touches zero. No new area is being added.*

![Graph](image)

**Fig. 5.5** The area is \( \Delta f = 7\frac{1}{2} - 3\frac{1}{2} = 4 \). Since \( v(x) \) decreases, \( f(x) \) bends down.

### INDEFINITE INTEGRALS AND DEFINITE INTEGRALS

We have to distinguish two different kinds of integrals. They both use the antiderivative \( f(x) \). The definite one involves the limits 0 and 4, the indefinite one doesn't:

*The indefinite integral is a function \( f(x) = 4x - \frac{1}{2} x^2 \).*

*The definite integral from \( x = 0 \) to \( x = 4 \) is the number \( f(4) - f(0) \).*

The definite integral is definitely 8. But the indefinite integral is not necessarily \( 4x - \frac{1}{2} x^2 \). *We can change \( f(x) \) by a constant without changing its derivative* (since the
derivative of a constant is zero). The following functions are also antiderivatives:

\[ f(x) = 4x - \frac{1}{2}x^2 + 1, \quad f(x) = 4x - \frac{1}{2}x^2 - 9, \quad f(x) = 4x - \frac{1}{2}x^2 + C. \]

The first two are particular examples. The last is the general case. The constant \( C \) can be anything (including zero), to give all functions with the required derivative. The theory of calculus will show that there are no others. The indefinite integral is the most general antiderivative (with no limits):

\[ \text{indefinite integral } f(x) = \int v(x) \, dx = 4x - \frac{1}{2}x^2 + C. \] (5)

By contrast, the definite integral is a number. It contains no arbitrary constant \( C \). More that that, it contains no variable \( x \). The definite integral is determined by the function \( v(x) \) and the limits of integration (also known as the endpoints). It is the area under the graph between those endpoints.

To see the relation of indefinite to definite, answer this question: What is the definite integral between \( x = 1 \) and \( x = 3 \)? The indefinite integral gives \( f(3) = 7\frac{1}{2} + C \) and \( f(1) = 3\frac{1}{2} + C \). To find the area between the limits, subtract \( f \) at one limit from \( f \) at the other limit:

\[ \int_{1}^{3} v(x) \, dx = f(3) - f(1) = (7\frac{1}{2} + C) - (3\frac{1}{2} + C) = 4. \] (6)

The constant cancels itself! The definite integral is the difference between the values of the indefinite integral. \( C \) disappears in the subtraction.

The difference \( f(3) - f(1) \) is like \( \Delta y \). The sum of \( v \) from 1 to \( n \) has become "the integral of \( v(x) \) from 1 to 3." Section 5.3 computes other areas from sums, and 5.4 computes many more from antiderivatives. Then we come back to the definite integral and the Fundamental Theorem:

\[ \int_{a}^{b} v(x) \, dx = \int_{a}^{b} \frac{df}{dx} \, dx = f(b) - f(a). \] (7)

### 5.2 Exercises

**Read-through questions**

Integration yields the \( \int \) under a curve \( y = v(x) \). It starts from rectangles with base \( b \) and heights \( v(x) \) and areas \( \int v(x) \, dx \). As \( \Delta x \to 0 \) the area \( v_1\Delta x + \cdots + v_n\Delta x \) becomes the \( \int \) of \( v(x) \). The symbol for the indefinite integral of \( v(x) \) is \( \int \). The problem of integration is solved if we find \( f(x) \) such that \( \int f(x) \, dx = v(x) \). Then \( f(x) \) is the \( \int \) of \( v(x) \), and \( \int_{a}^{b} v(x) \, dx \) equals \( b \), minus \( a \). The limits of integration are \( a \), \( b \). This is a \( \int \) integral, which is a \( \int \) and not a function \( f(x) \).

The example \( v(x) = x \) has \( f(x) = \frac{1}{2}x^2 \). It also has \( f(x) = -\frac{1}{2}x^2 \). The area under \( v(x) \) from 2 to 6 is \( \int_{2}^{6} v(x) \, dx \). The constant is canceled in computing the difference \( \int_{a}^{b} v(x) \, dx \). If \( v(x) = x^2 \) then \( f(x) = \frac{1}{3}x^3 \).

The sum \( v_1 + \cdots + v_n = f_n - f_0 \) leads to the Fundamental Theorem \( \int_{a}^{b} v(x) \, dx = \int_{a}^{b} \frac{df}{dx} \, dx = f(b) - f(a) \). Finding the \( \int \) under the \( v \)-graph is the opposite of finding the \( \int \) of the \( f \)-graph.

Find an antiderivative \( f(x) \) for \( v(x) \) in 1–14. Then compute the definite integral \( \int_{a}^{b} v(x) \, dx = f(b) - f(a) \).

1. \( 5x^4 + 4x^5 \)
2. \( x + 12x^2 \)
3. \( \int \sqrt{x} \, \text{for } x^{-1/2} \)
4. \( (\sqrt{x})^3 \text{ or } x^{3/2} \)
5. \( x^{-1/3} + (2x)^{1/3} \)
6. \( x^{1/3} \text{ or } x^{2/3} \)
7. \( 2 \sin x + 2x \)
8. \( \sec^2 x + 1 \)
9. \( x \cos x \) (by experiment)
10. \( x \sin x \) (by experiment)
11. \( \sin x \cos x \)
12. \( \sin^2 x \cos x \)
13. \( 0 \) (find all \( f \))
14. \( -1 \) (find all \( f \))
15. If \( df/dx = v(x) \) then the definite integral of \( v(x) \) from \( a \) to \( b \) is \( \int_{a}^{b} v(x) \, dx \). If \( f_1 - f_0 = v_1 \) then the definite sum of \( v_1 + \cdots + v_n \) is \( \int_{a}^{b} v(x) \, dx \).
16. The areas include a factor \( \Delta x \), the base of each rectangle. So the sum of \( v \)'s is multiplied by \( \Delta x \) to approach the integral. The difference of \( f \)'s is divided by \( \Delta x \) to approach the derivative.
5.3 Summation versus Integration

This section does integration the hard way. We find explicit formulas for $f_n = v_1 + \cdots + v_n$. From areas of rectangles, the limits produce the area $f(x)$ under a curve. According to the Fundamental Theorem, $df/dx$ should return us to $v(x)$ — and we verify in each case that it does.

May I recall that there is sometimes an easier way? If we can find an $f(x)$ whose derivative is $v(x)$, then the integral of $v$ is $f$. Sums and limits are not required, when $f$ is spotted directly. The next section, which explains how to look for $f(x)$, will displace this one. (If we can't find an antiderivative we fall back on summation.) Given a successful $f$, adding any constant produces another $f$ — since the derivative of the constant is zero. The right constant achieves $f(0) = 0$, with no extra effort.
This section constructs \( f(x) \) from sums. The next section searches for antiderivatives.

**THE SIGMA NOTATION**

In a section about sums, there has to be a decent way to express them. Consider \( 1^2 + 2^2 + 3^2 + 4^2 \). The individual terms are \( v_j = j^2 \). Their sum can be written in *summation notation*, using the capital Greek letter \( \Sigma \) (pronounced sigma):

\[
1^2 + 2^2 + 3^2 + 4^2 \text{ is written } \sum_{j=1}^{4} j^2.
\]

Spoken aloud, that becomes "**the sum of** \( j^2 \) **from** \( j = 1 \) **to** \( 4 \)." It equals 30. The limits on \( j \) (written below and above \( \Sigma \)) indicate where to start and stop:

\[
v_1 + \cdots + v_4 = \sum_{j=1}^{4} v_j \quad \text{and} \quad v_3 + \cdots + v_9 = \sum_{k=3}^{9} v_k.
\]

The \( k \) at the end of (1) makes an additional point. There is nothing special about the letter \( j \). That is a "dummy variable," no better and no worse than \( k \) (or \( i \)). Dummy variables are only on one side (the side with \( \Sigma \)), and they have no effect on the sum. *The upper limit \( n \) is on both sides.* Here are six sums:

\[
\begin{align*}
\sum_{k=1}^{n} k &= 1 + 2 + 3 + \cdots + n \\
\sum_{j=1}^{4} (j - 1)^4 &= -1 + 1 - 1 + 1 = 0 \\
\sum_{j=1}^{5} (2j - 1) &= 1 + 3 + 5 + 7 + 9 = 5^2 \\
\sum_{i=0}^{n} v_i &= v_0 \quad \text{[only one term]} \\
\sum_{j=1}^{4} j^2 &= \text{[meaningless?] } \\
\sum_{k=0}^{10} \frac{1}{2^k} &= 1 + \frac{1}{2} + \frac{1}{4} - \cdots = 2 \quad \text{[infinite series]} \end{align*}
\]

The numbers 1 and \( n \) or 1 and 4 (or 0 and \( \infty \)) are the *lower limit* and *upper limit*. The dummy variable \( i \) or \( j \) or \( k \) is the *index* of summation. I hope it seems reasonable that the infinite series \( 1 + \frac{1}{2} + \frac{1}{4} + \cdots \) adds to 2. We will come back to it in Chapter 10.

A sum like \( \sum_{j=1}^{6} b \) looks meaningless, but it is actually \( 6 + 6 + \cdots + 6 = 6n \). It follows the rules. In fact \( \sum_{j=1}^{4} j^2 \) is not meaningless either. Every term is \( j^2 \) and by the same rules, that sum is \( 4j^2 \). However the \( i \) was probably intended to be \( j \). Then the sum is \( 1 + 4 + 9 + 16 = 30 \).

**Question** What happens to these sums when the upper limits are changed to \( n \)?

**Answer** The sum depends on the stopping point \( n \). A formula is required (when possible). Integrals stop at \( x \), sums stop at \( n \), and we now look for special cases when \( f(x) \) or \( f_n \) can be found.

**A SPECIAL SUMMATION FORMULA**

How do you add the first 100 whole numbers? The problem is to compute

\[
\sum_{j=1}^{100} j = 1 + 2 + 3 + \cdots + 98 + 99 + 100 = ?
\]

\[\text{Zeno the Greek believed it was impossible to get anywhere, since he would only go halfway and then half again and half again. Infinite series would have changed his whole life.}\]
If you were Gauss, you would see the answer at once. (He solved this problem at a ridiculous age, which gave his friends the idea of getting him into another class.) His solution was to combine $1 + 100$, and $2 + 99$, and $3 + 98$, always adding to 101. There are fifty of those combinations. Thus the sum is $(50)(101) = 5050$.

The sum from 1 to $n$ uses the same idea. The first and last terms add to $n + 1$. The next terms $n - 1$ and 2 also add to $n + 1$. If $n$ is even (as 100 was) then there are $\frac{1}{2}n$ parts. Therefore the sum is $\frac{1}{2}n(n + 1)$:

\[
\sum_{j=1}^{n} j = 1 + 2 + \cdots + (n - 1) + n = \frac{1}{2}n(n + 1).
\] (2)

The important term is $\frac{1}{2}n^2$, but the exact sum is $\frac{1}{2}n^2 + \frac{1}{2}n$.

What happens if $n$ is an odd number (like $n = 99$)? Formula (2) remains true. The combinations $1 + 99$ and $2 + 98$ still add to $n + 1 = 100$. There are $\frac{1}{2}(99) = 49\frac{1}{2}$ such pairs, because the middle term (which is 50) has nothing to combine with. Thus $1 + 2 + \cdots + 99$ equals $49\frac{1}{2}$ times 100, or 4950.

Remark That sum had to be 4950, because it is 5050 minus 100. The sum up to 99 equals the sum up to 100 with the last term removed. Our key formula $f_n - f_{n-1} = v_n$ has turned up again!

**EXAMPLE** Find the sum $101 + 102 + \cdots + 200$ of the second hundred numbers.

First solution This is the sum from 1 to 200 minus the sum from 1 to 100:

\[
\sum_{j=1}^{200} j = \sum_{j=1}^{100} j + \sum_{j=101}^{200} j.
\] (3)

The middle sum is $\frac{1}{2}(200)(201)$ and the last is $\frac{1}{2}(100)(101)$. Their difference is 15050. **Note!** I left out $"j=\"$ in the limits. It is there, but not written.

Second solution The answer 15050 is exactly the sum of the first hundred numbers (which was 5050) plus an additional 10000. Believing that a number like 10000 can never turn up by accident, we look for a reason. It is found through changing the limits of summation:

\[
\sum_{j=1}^{200} j = \sum_{k=1}^{100} (k + 100).
\] (4)

This is important, to be able to shift limits around. Often the lower limit is moved to zero or one, for convenience. Both sums have 100 terms (that doesn’t change). The dummy variable $j$ is replaced by another dummy variable $k$. They are related by $j = k + 100$ or equivalently by $k = j - 100$.

The variable must change everywhere—in the lower limit and the upper limit as well as inside the sum. If $j$ starts at 101, then $k = j - 100$ starts at 1. If $j$ ends at 200, $k$ ends at 100. If $j$ appears in the sum, it is replaced by $k + 100$ (and if $j^2$ appeared it would become $(k + 100)^2$).

From equation (4) you see why the answer is 15050. The sum $1 + 2 + \cdots + 100$ is 5050 as before. 100 is added to each of those 100 terms. That gives 10000.

**EXAMPLES OF CHANGING THE VARIABLE (and the limits)**

\[
\sum_{i=0}^{2} 2^i equals \sum_{j=1}^{4} 2^{j-1} \, (here \, i = j - 1). \, Both \, sums \, are \, 1 + 2 + 4 + 8
\]

\[
\sum_{i=3}^{6} v_i equals \sum_{j=0}^{n-3} v_{j+3} \, (here \, i = j + 3 \, and \, j = i - 3). \, Both \, sums \, are \, v_3 + \cdots + v_n.
\]
Why change \( n \) to \( n - 3 \)? Because the upper limit is \( i = n \). So \( j + 3 = n \) and \( j = n - 3 \).

A final step is possible, and you will often see it. \textbf{The new variable} \( j \) \textbf{can be changed back to} \( i \). Dummy variables have no meaning of their own, but at first the result looks surprising:

\[
\sum_{i=0}^{n} \frac{2i}{i+1} \text{ equals } \sum_{j=1}^{n} 2^{j-1} \text{ equals } \sum_{i=1}^{6} 2^{i-1}.
\]

With practice you might do that in one step, skipping the temporary letter \( j \). Every \( i \) on the left becomes \( i - 1 \) on the right. Then \( i = 0, \ldots, 5 \) changes to \( i = 1, \ldots, 6 \). (At first two steps are safer.) This may seem a minor point, but soon we will be changing the limits on integrals instead of sums. Integration is parallel to summation, and it is better to see a "change of variable" here first.

\textbf{Note about} \( 1 + 2 + \cdots + n \). The good thing is that Gauss found the sum \( \frac{1}{2}n(n+1) \). The bad thing is that his method looked too much like a trick. I would like to show how this fits the fundamental rule connecting sums and differences:

\[
\text{if } v_1 + v_2 + \cdots + v_n = f_n \text{ then } v_n = f_n - f_{n-1}. \tag{5}
\]

Gauss says that \( f_n \) is \( \frac{1}{2}n(n+1) \). Reducing \( n \) by 1, his formula for \( f_{n-1} \) is \( \frac{1}{2}(n-1)n \). \textbf{The difference} \( f_n - f_{n-1} \text{ should be the last term in the sum:} \)

\[
f_n - f_{n-1} = \frac{1}{2}n(n+1) - \frac{1}{2}(n-1)n = \frac{1}{2}(n^2 + n - n^2 + n) = n. \tag{6}
\]

This is the one term \( v_n = n \) that is included in \( f_n \) but not in \( f_{n-1} \).

There is a deeper point here. For any sum \( f_n \), there are two things to check. The \( f \)'s must \textit{begin} correctly and they must \textit{change} correctly. The underlying idea is \textbf{mathematical induction}: \textit{Assume the statement is true below} \( n \). \textit{Prove it for} \( n \).

\textbf{Goal: } \textit{To prove that} \( 1 + 2 + \cdots + n = \frac{1}{2}n(n+1) \). \textit{This is the guess} \( f_n \).

\textbf{Proof by induction: } \textit{Check} \( f_1 \) (it equals 1). \textit{Check} \( f_n - f_{n-1} \) (it equals \( n \)).

For \( n = 1 \) the answer \( \frac{1}{2}n(n+1) = \frac{1}{2} \cdot 1 \cdot 2 \) is correct. For \( n = 2 \) this formula \( \frac{1}{2} \cdot 2 \cdot 3 \) agrees with \( 1 + 2 \). But that separate test is not necessary! \textit{If} \( f_1 \) \textit{is right, and if the change} \( f_n - f_{n-1} \textit{is right for every} n \), \textit{then} \( f_n \) \textit{must be right}. Equation (6) was the key test, to show that the change in \( f \)'s agrees with \( v \).

That is the logic behind mathematical induction, but I am not happy with most of the exercises that use it. There is absolutely no excitement. The answer is given by some higher power (like Gauss), and it is proved correct by some lower power (like us). It is much better when we lower powers find the answer for ourselves.† Therefore I will try to do that for the second problem, which is the \textbf{sum of squares}.

\section*{The Sum of \( j^2 \) and the Integral of \( x^2 \)}

An important calculation comes next. It is the area in Figure 5.6. One region is made up of rectangles, so its area is a sum of \( n \) pieces. The other region lies under the parabola \( v = x^2 \). It cannot be divided into rectangles, and calculus is needed.

The first problem is to find \( f_n = 1^2 + 2^2 + 3^2 + \cdots + n^2 \). This is a sum of squares, with \( f_1 = 1 \) and \( f_2 = 5 \) and \( f_3 = 14 \). The goal is to find the pattern in that sequence. By trying to guess \( f_n \), we are copying what will soon be done for integrals.

Calculus looks for an \( f(x) \) whose derivative is \( v(x) \). There \( f \) is an \textbf{antiderivative} (or

† \textit{The goal of real teaching is for the} student \textit{to find the answer. And also the problem.}
5.3 Summation versus Integration

an integral). Algebra looks for \( f_n \)'s whose differences produce \( v_n \). Here \( f_n \) could be called an antidifference (better to call it a sum).

The best start is a good guess. Copying directly from integrals, we might try \( f_n = \frac{1}{3} n^3 \). To test if it is right, check whether \( f_n - f_{n-1} \) produces \( v_n = n^2 \):

\[
\frac{1}{3} n^3 - \frac{1}{3} (n-1)^3 = \frac{1}{3} n^3 - \frac{1}{3} (n^3 - 3n^2 + 3n - 1) = n^2 - n + \frac{1}{3}.
\]

We see \( n^2 \), but also \( -n + \frac{1}{3} \). The guess \( \frac{1}{3} n^3 \) needs correction terms. To cancel \( \frac{1}{3} \) in the difference, I subtract \( \frac{1}{3} \) from the sum. To put back \( n \) in the difference, I add \( 1 + 2 + \cdots + n = \frac{1}{2} n(n+1) \) to the sum. The new guess (which should be right) is

\[
f_n = \frac{1}{3} n^3 + \frac{1}{3} n(n+1) - \frac{1}{3} n = \frac{1}{3} n^3 + \frac{1}{2} n^2 + \frac{1}{6} n. \tag{7}
\]

To check this answer, verify first that \( f_1 = 1 \). Also \( f_2 = 5 \) and \( f_3 = 14 \). To be certain, verify that \( f_n - f_{n-1} = n^2 \). For calculus the important term is \( \frac{1}{3} n^3 \):

The sum \( \sum_{j=1}^{n} j^2 \) of the first \( n \) squares is \( \frac{1}{3} n^3 \) plus corrections \( \frac{1}{2} n^2 \) and \( \frac{1}{6} n \).

In practice \( \frac{1}{3} n^3 \) is an excellent estimate. The sum of the first 100 squares is approximately \( \frac{1}{3} (100)^3 \), or a third of a million. If we need the exact answer, equation (7) is available: the sum is 338,350. Many applications (example: the number of steps to solve 100 linear equations) can settle for \( \frac{1}{3} n^3 \).

What is fascinating is the contrast with calculus. Calculus has no correction terms! They get washed away in the limit of thin rectangles. When the sum is replaced by the integral (the area), we get an absolutely clean answer:

The integral of \( v = x^2 \) from \( x = 0 \) to \( x = n \) is exactly \( \frac{1}{3} n^3 \).

The area under the parabola, out to the point \( x = 100 \), is precisely a third of a million. We have to explain why, with many rectangles.

The idea is to approach an infinite number of infinitely thin rectangles. A hundred rectangles gave an area of 338,350. Now take a thousand rectangles. Their heights are \( (\frac{1}{10})^2, (\frac{1}{10})^2, \ldots \) because the curve is \( v = x^2 \). The base of every rectangle is \( \Delta x = \frac{1}{10} \), and we add heights times base:

area of rectangles = \( \left( \frac{1}{10} \right)^2 \left( \frac{1}{10} \right) + \left( \frac{2}{10} \right)^2 \left( \frac{1}{10} \right) + \cdots + \left( \frac{1000}{10} \right)^2 \left( \frac{1}{10} \right) \).

Factor out \( (\frac{1}{10})^3 \). What you have left is \( 1^2 + 2^2 + \cdots + 1000^2 \), which fits the sum of squares formula. The exact area of the thousand rectangles is 333,833.5. I could try to guess ten thousand rectangles but I won't.

Main point: The area is approaching 333,333.333... But the calculations are getting worse. It is time for algebra—which means that we keep "\( \Delta x \)" and avoid numbers.
The interval of length 100 is divided into $n$ pieces of length $\Delta x$. (Thus $n = \frac{100}{\Delta x}$.)

The $j$th rectangle meets the curve $v = x^2$, so its height is $(j\Delta x)^2$. Its base is $\Delta x$, and we add areas:

$$
\text{area} = (\Delta x)(\Delta x) + (2\Delta x)(\Delta x) + \cdots + (n\Delta x)(\Delta x) = \sum_{j=1}^{n} (j\Delta x)^2(\Delta x). \quad (8)
$$

Factor out $(\Delta x)^3$, leaving a sum of $n$ squares. The area is $(\Delta x)^3$ times $\frac{f}{\Delta x}$, and $n = \frac{100}{\Delta x}$.

$$(\Delta x)^3 \left[ \frac{1}{3} \left( \frac{100}{\Delta x} \right)^3 + \frac{1}{2} \left( \frac{100}{\Delta x} \right)^2 + \frac{1}{6} \left( \frac{100}{\Delta x} \right) \right] = \frac{1}{3}100^3 + \frac{1}{2}100^2(\Delta x) + \frac{1}{6}100(\Delta x)^2. \quad (9)
$$

This equation shows what is happening. The leading term is a third of a million, as predicted. The other terms are approaching zero! They contain $\Delta x$, and as the rectangles get thinner they disappear. They only account for the small corners of rectangles that lie above the curve. The vanishing of those corners will eventually be proved for any continuous functions—the area from the correction terms goes to zero—but here in equation (9) you see it explicitly.

The area under the curve came from the central idea of integration: $\frac{100}{\Delta x}$ rectangles of width $\Delta x$ approach the limiting area $= \frac{1}{3}(100)^3$. The rectangular area is $\Sigma v_j \Delta x$. The exact area is $\int v(x) \, dx$. In the limit $\Sigma$ becomes $\int$ and $v_j$ becomes $v(x)$ and $\Delta x$ becomes $dx$.

That completes the calculation for a parabola. It used the formula for a sum of squares, which was special. But the underlying idea is much more general. The limit of the sums agrees with the antiderivative: The antiderivative of $v(x) = x^2$ is $f(x) = \frac{1}{3}x^3$.

According to the Fundamental Theorem, the area under $v(x)$ is $f(x)$:

$$
\int_{0}^{100} v(x) \, dx = f(100) - f(0) = \frac{1}{3}(100)^3.
$$

That Fundamental Theorem is not yet proved! I mean it is not proved by us. Whether Leibniz or Newton managed to prove it, I am not quite sure. But it can be done. Starting from sums of differences, the difficulty is that we have too many limits at once. The sums of $r_j \Delta x$ are approaching the integral. The differences $\Delta f/\Delta x$ approach the derivative. A real proof has to separate those steps, and Section 5.7 will do it.

Proved or not, you are seeing the main point. What was true for the numbers $f_j$ and $r_j$ is true in the limit for $v(x)$ and $f(x)$. Now $v(x)$ can vary continuously, but it is still the slope of $f(x)$. The reverse of slope is area.

![Proof without words by Roger Nelsen](Mathematics Magazine 1990)

Finally we review the area under $v = x$. The sum of $1 - 2 + \cdots + n$ is $\frac{1}{2}n^2 - \frac{1}{2}n$. This gives the area of $n = 4/\Delta x$ rectangles, going out to $x = 4$. The heights are $j\Delta x$, the bases are $\Delta x$, and we add areas:

$$
\sum_{j=1}^{4/\Delta x} (j\Delta x)(\Delta x) = (\Delta x)^2 \left[ \frac{1}{2} \left( \frac{4}{\Delta x} \right)^2 + \frac{1}{2} \left( \frac{4}{\Delta x} \right) \right] = 8 + 2\Delta x. \quad (10)
$$
5.3 Summation versus Integration

With $\Delta x = 1$ the area is $1 + 2 + 3 + 4 = 10$. With eight rectangles and $\Delta x = \frac{1}{2}$, the area was $8 + 2\Delta x = 9$. Sixteen rectangles of width $\frac{1}{8}$ brought the correction $2\Delta x$ down to $\frac{1}{2}$. The exact area is 8. The error is proportional to $\Delta x$.

**Important note** There you see a question in applied mathematics. If there is an error, what size is it? How does it behave as $\Delta x \to 0$? The $\Delta x$ term disappears in the limit, and $(\Delta x)^2$ disappears faster. But to get an error of $10^{-6}$ we need **eight million rectangles**:

$$2\Delta x = 2 \cdot 4/8,000,000 = 10^{-6}.$$  

That is horrifying! The numbers $10, 9, 8\frac{1}{2}, 8\frac{3}{4}, ...$ seem to approach the area 8 in a satisfactory way, but the convergence is **much too slow**. It takes twice as much work to get one more binary digit in the answer—which is absolutely unacceptable. Somehow the $\Delta x$ term must be removed. If the correction is $(\Delta x)^2$ instead of $\Delta x$, then a thousand rectangles will reach an accuracy of $10^{-6}$.

**The problem is that the rectangles are unbalanced.** Their right sides touch the graph of $v$, but their left sides are much too high. The best is to cross the graph in the **middle** of the interval—this is the **midpoint rule**. Then the rectangle sits halfway across the line $v = x$, and the error is zero. Section 5.8 comes back to this rule—and to Simpson’s rule that fits parabolas and removes the $(\Delta x)^2$ term and is built into many calculators.

Finally we try the quick way. The area under $v = x$ is $f = \frac{1}{2}x^2$, because $df/dx$ is $v$. The area out to $x = 4$ is $\frac{1}{2}(4)^2 = 8$. Done.

![Fig. 5.7](image)

**Fig. 5.7** Endpoint rules: error $\sim 1/(\text{work}) \sim 1/n$. Midpoint rule is better: error $\sim 1/(\text{work})^2$.

**Optional: pth powers** Our sums are following a pattern. First, $1 + \cdots + n$ is $\frac{1}{2}n^2$ plus $\frac{1}{2}n$. The sum of squares is $\frac{1}{3}n^3$ plus correction terms. **The sum of pth powers is**

$$1^p + 2^p + \cdots + n^p = \frac{1}{p+1} n^{p+1} + \text{correction terms.} \quad (11)$$

The correction involves lower powers of $n$, and you know what is coming. **Those corrections disappear in calculus.** The area under $v = x^p$ from 0 to $n$ is

$$\int_0^n x^p \, dx = \lim_{\Delta x \to 0} \sum_{j=1}^{n\Delta x} (j\Delta x)^p = \frac{1}{p+1} n^{p+1}. \quad (12)$$

Calculus doesn’t care if the upper limit $n$ is an integer, and it doesn’t care if the power $p$ is an integer. We only need $p+1 > 0$ to be sure $n^{p+1}$ is genuinely the leading term. **The antiderivative of** $v = x^p$ **is** $f = x^{p+1}/(p+1)$.

We are close to interesting experiments. The correction terms disappear and the sum approaches the integral. Here are actual numbers for $p = 1$, when the sum and integral are easy: $S_n = 1 + \cdots + n$ and $I_n = \int_0^n x \, dx = \frac{1}{2}n^2$. The difference is $D_n = \frac{1}{2}n$. The thing to watch is the relative error $E_n = D_n/I_n$:

<table>
<thead>
<tr>
<th>$n$</th>
<th>$S_n$</th>
<th>$I_n$</th>
<th>$D_n = S_n - I_n$</th>
<th>$E_n = D_n/I_n$</th>
</tr>
</thead>
<tbody>
<tr>
<td>100</td>
<td>5050</td>
<td>5000</td>
<td>50</td>
<td>0.010</td>
</tr>
<tr>
<td>200</td>
<td>20100</td>
<td>20000</td>
<td>100</td>
<td>0.005</td>
</tr>
</tbody>
</table>
The number 20100 is \( f(200)(201) \). Please write down the next line \( n = 400 \), and please find a formula for \( E_n \). You can guess \( E_n \) from the table, or you can derive it from knowing \( S_n \) and \( I_n \). The formula should show that \( E_n \) goes to zero. More important, it should show how quick (or slow) that convergence will be.

One more number—a third of a million—was mentioned earlier. It came from integrating \( x^2 \) from 0 to 100, which compares to the sum \( S_{100} \) of 100 squares:

\[
\begin{array}{cccc}
 n & p & S_n & I_n = \frac{1}{3}n^3 \\
100 & 2 & 338350 & 333333\frac{1}{3} \\
200 & 2 & 2686700 & 2666666\frac{2}{3} \\
\end{array}
\]

These numbers suggest a new idea, to keep \( n \) fixed and change \( p \). The computer can find sums without a formula! With its help we go to fourth powers and square roots:

\[
\begin{array}{cccc}
 n & p & S = 1^p + \cdots + n^p & I = n^{p+1}/(p + 1) \\
100 & 4 & 2050333330 & \frac{1}{6}(100)^5 \\
100^{1\frac{1}{2}} & 671.4629 & \frac{3}{2}(100)^{3/2} & 4.7963 \\
\end{array}
\]

In this and future tables we don't expect exact values. The last entries are rounded off, and the goal is to see the pattern. The errors \( E_{n,p} \) are sure to obey a systematic rule—they are proportional to \( 1/n \) and to an unknown number \( C(p) \) that depends on \( p \). I hope you can push the experiments far enough to discover \( C(p) \). This is not an exercise with an answer in the back of the book—it is mathematics.

### 5.3 Exercises

**Read-through questions**

The Greek letter \( \sum \) indicates summation. In \( \sum_i v_j \), the dummy variable is \( b \). The limits are \( c \), so the first term is \( d \) and the last term is \( e \). When \( v_j = j \) this sum equals \( f \). For \( n = 100 \) the leading term is \( g \).

The correction term is \( h \). The leading term equals the integral of \( v = x \) from 0 to 100, which is written \( i \). The sum is the total \( j \) of 100 rectangles. The correction term is the area between the \( k \) and the \( l \).

The sum \( \sum_{i=1}^{6} i^2 \) is the same as \( \sum_{i=1}^{4} m \) and equals \( n \). The sum \( \sum_{i=4}^{5} v_i \) is the same as \( o \).\( v_{i+4} \) and equals \( p \). For \( f_{i,n} = \sum_{i=1}^{n} v_j \) the difference \( f_n - f_{n-1} \) equals \( q \).

The formula for \( 1^2 + 2^2 + \cdots + n^2 \) is \( f_n = \frac{1}{3}r \). To prove it by mathematical induction, check \( f_1 = \frac{1}{3}s \) and check \( f_n - f_{n-1} = \frac{1}{3}t \). The area under the parabola \( v = x^2 \) from \( x = 0 \) to \( x = 9 \) is \( u \). This is close to the area of \( v \) rectangles of base \( \Delta x \). The correction terms approach zero very \( w \).

1. Compute the numbers \( \sum_{i=1}^{4} 1/n \) and \( \sum_{i=2}^{3} (2i - 3) \).
2. Compute \( \sum_{j=6}^{7} (j^2 - j) \) and \( \sum_{j=1}^{8} 1/2^j \).
3. Evaluate the sum \( \sum_{i=0}^{5} 2^i \) and \( \sum_{i=0}^{n} 2^i \).
4. Evaluate \( \sum_{j=1}^{n} (-1)^j \) and \( \sum_{j=1}^{n} (-1)^{2j} \).
5. Write these sums in sigma notation and compute them:

   \[
   2 + 4 + 6 + \cdots + 100 \quad 1 + 3 + 5 + \cdots + 199 \quad 1 - \frac{1}{2} + \frac{1}{3} - \frac{1}{4}
   \]

6. Express these sums in sigma notation:

   \[
   v_1 - v_2 + v_3 - v_4 + v_1w_1 + v_2w_2 + \cdots + v_nw_n \quad v_1 + v_3 + v_4
   \]

7. Convert these sums to sigma notation:

   \[
   a_0 + a_1x + \cdots + a_nx^n \quad \sin \frac{2\pi}{n} + \sin \frac{4\pi}{n} + \cdots + \sin 2\pi
   \]

8. The binomial formula uses coefficients \( \binom{n}{j} = \frac{n!}{j!(n-j)!} \):

   \[
   (a + b)^n = \binom{n}{0}a^n + \binom{n}{1}a^{n-1}b + \cdots + \binom{n}{n}b^n = \sum_{j=0}^{n} b^j
   \]

9. With electronic help compute \( \sum_{i=1}^{100} 1/i \) and \( \sum_{i=1}^{1000} 1/i \).
10. On a computer find \( \sum_{i=0}^{10} (-1)^{i+1} \) times \( \sum_{i=0}^{10} 1/i! \).
11 Simplify \( \sum_{i=1}^{n} (a_i + b_i)^2 + \sum_{i=1}^{n} (a_i - b_i)^2 \) to \( \sum_{i=1}^{n} \).  
12 Show that \( \left( \sum_{i=1}^{n} a_i \right)^2 \neq \sum_{i=1}^{n} a_i^2 \) and \( \sum_{i=1}^{n} a_i b_i \neq \sum_{i=1}^{n} a_i \sum_{i=1}^{n} b_i \).  
13 "Telescope" the sums \( \sum_{i=1}^{n} (2^i - 2^{i-1}) \) and \( \sum_{i=1}^{n} \left( \frac{1}{j+1} - \frac{1}{j} \right) \). All but two terms cancel.  
14 Simplify the sums \( \sum_{j=1}^{n} (f_j - f_{j-1}) \) and \( \sum_{j=3}^{n} (f_j + f_{j-2}) \).  
15 True or false: (a) \( \sum_{j=4}^{8} v_j = \sum_{j=2}^{6} v_{j-2} \) (b) \( \sum_{i=1}^{9} v_i = \sum_{i=3}^{11} v_{i-2} \).  
16 \( \sum_{i=1}^{n} v_i = \sum_{j=0}^{n-1} \) and \( \sum_{i=0}^{n} j^2 = \sum_{i=2}^{n} \).  
17 The antiderivative of \( \frac{d^2 f}{dx^2} \) is \( \frac{df}{dx} \). What is the sum \( \frac{d^2 f}{dx^2} + \frac{d f}{dx} + f \)?  
18 Induction: Verify that \( 1^2 + 2^2 + \cdots + n^2 = \frac{n(n+1)(2n+1)}{6} \) by checking that \( f_1 \) is correct and \( f_n - f_{n-1} = n^2 \).  
19 Prove by induction: \( 1 + 3 + \cdots + (2n-1) = n^2 \).  
20 Verify that \( 1^3 + 2^3 + \cdots + n^3 = \frac{n^2 (n+1)^2}{4} \) by checking \( f_1 \) and \( f_n - f_{n-1} \). The text has a proof without words.  
21 Suppose \( f_n \) has the form \( an + bn^2 + cn^3 \). If you know \( f_1 = 1, f_2 = 5, f_3 = 14 \), turn those into three equations for \( a, b, c \). The solutions \( a = \frac{1}{6}, b = \frac{1}{2}, c = \frac{1}{3} \) give what formula?  
22 Find \( q \) in the formula \( 1^p + \cdots + n^p = \binom{n+1}{p+1} + \text{correction} \).  
23 Add \( n = 400 \) to the table for \( S_n = 1 + \cdots + n \) and find the relative error \( E_n \). Guess and prove a formula for \( E_n \).  
24 Add \( n = 50 \) to the table for \( S_n = 1^2 + \cdots + n^2 \) and compute \( E_{50} \). Find an approximate formula for \( E_n \).  
25 Add \( p = \frac{1}{2} \) and \( p = 3 \) to the table for \( S_{100,p} = 1^p + \cdots + 100^p \). Guess an approximate formula for \( E_{100,p} \).  
26 Guess \( C(p) \) in the formula \( E_{n,p} \approx C(p)/n \).  
27 Show that \( |1-5| < |1| + |-5| \). Always \( |v_1 + v_2| < |v_1| + |v_2| \) unless \( \ldots \).  
28 Let \( S \) be the sum \( 1 + x + x^2 + \cdots \) of the (infinite) geometric series. Then \( xS = x + x^2 + x^3 + \cdots \) is the same as \( S \) minus \( \ldots \). Therefore \( S = \ldots \). None of this makes sense if \( x = 2 \) because \( \ldots \).  
29 The double sum \( \sum_{i=1}^{3} \left[ \sum_{j=1}^{3} (i+j) \right] = v_1 = \sum_{i=1}^{3} (1+j) \) plus \( v_2 = \sum_{j=1}^{3} (2+j) \). Compute \( v_1 \) and \( v_2 \) and the double sum.  
30 The double sum \( \sum_{i=1}^{3} \left( \sum_{j=1}^{3} w_{ij} \right) \) is \( (w_{1,1} + w_{1,2} + w_{1,3}) + \ldots \). The double sum \( \sum_{j=1}^{3} \left( \sum_{i=1}^{3} w_{ij} \right) \) is \( (w_{1,1} + w_{2,1} + w_{3,1}) + (w_{1,2} + w_{2,2} + \ldots) \). Compare.  
31 Find the flaw in the proof that \( 2^n = 1 \) for every \( n = 0, 1, 2, \ldots \). For \( n = 0 \) we have \( 2^0 = 1 \). If \( 2^n = 1 \) for every \( n < N \), then \( 2^n = 2^{N-1} \cdot 2^{N-1}/2^{N-2} = 1 \cdot 1/1 = 1 \).  
32 Write out all terms to see why the following are true: \( \sum_{i=1}^{3} \sum_{j=1}^{3} u_{ij} v_{ij} = \left( \sum_{i=1}^{3} u_{ij} \right) \left( \sum_{j=1}^{3} v_{ij} \right) \).  
33 The average of \( 6, 11, 4 \) is \( \bar{x} = \frac{1}{3}(6 + 11 + 4) \). Then \( (6 - \bar{x}) + (11 - \bar{x}) + (4 - \bar{x}) = \ldots \). The average of \( v_1, \ldots, v_{100} \) is \( \bar{v} = \ldots \). Prove that \( \sum (v_i - \bar{v}) = 0 \).  
34 The Schwarz inequality is \( \left( \sum_{i=1}^{n} a_i b_i \right)^2 \leq \left( \sum_{i=1}^{n} a_i^2 \right) \left( \sum_{i=1}^{n} b_i^2 \right) \). Compute both sides if \( a_1 = 2, a_2 = 3, b_1 = 1, b_2 = 4 \). Then compute both sides for any \( a_1, a_2, b_1, b_2 \). The proof in Section 11.1 uses vectors.  
35 Suppose \( n \) rectangles with base \( \Delta x \) touch the graph of \( v(x) \) at the points \( x = \Delta x, 2\Delta x, \ldots, n\Delta x \). Express the total rectangular area in sigma notation.  
36 If \( 1/\Delta x \) rectangles with base \( \Delta x \) touch the graph of \( v(x) \) at the left end of each interval (thus at \( x = 0, \Delta x, 2\Delta x, \ldots \)) express the total area in sigma notation.  
37 The sum \( \sum_{j=1}^{100} \frac{1}{j} dx = \frac{1}{\Delta x} \sum_{j=1}^{100} f((j-1)\Delta x) - f(j\Delta x) \) equals \(.\) In the limit this becomes \( \int_{0}^{100} \frac{1}{x} dx = \ldots \).
search for an antiderivative may not succeed. We may not find \( f \). In that case we go back to rectangles, or on to something better in Section 5.8.

A computer is ready to integrate \( v \), but not by discovering \( f \). It integrates between specified limits, to obtain a \textit{number} (the definite integral). Here we hope to find a \textit{function} (the indefinite integral). That requires a symbolic integration code like MACSYMA or \textit{Mathematica} or MAPLE, or a reasonably nice \( v(x) \), or both. An expression for \( f(x) \) can have tremendous advantages over a list of numbers.

Thus our goal is to find antiderivatives and use them. The techniques will be further developed in Chapter 7—this section is short but good. First we write down what we know. \textit{On each line,} \( f(x) \) is an antiderivative of \( v(x) \) because \( df/dx = v(x) \).

<table>
<thead>
<tr>
<th>Known pairs</th>
<th>Function ( v(x) )</th>
<th>Antiderivative ( f(x) )</th>
</tr>
</thead>
<tbody>
<tr>
<td>Powers of ( x )</td>
<td>( x^n )</td>
<td>( x^{n+1}/(n+1) + C )</td>
</tr>
</tbody>
</table>

\( n = -1 \) is not included, because \( n + 1 \) would be zero. \( v = x^{-1} \) will lead us to \( f = \ln x \).

\textbf{Trigonometric functions}

\( \cos x \) \hspace{2cm} \sin x + C
\( \sin x \) \hspace{2cm} -\cos x + C
\( \sec^2 x \) \hspace{2cm} \tan x + C
\( \csc^2 x \) \hspace{2cm} -\cot x + C
\( \sec x \tan x \) \hspace{2cm} \sec x + C
\( \csc x \cot x \) \hspace{2cm} -\csc x + C

\textbf{Inverse functions}

\( 1/\sqrt{1-x^2} \) \hspace{2cm} \sin^{-1} x + C
\( 1/(1+x^2) \) \hspace{2cm} \tan^{-1} x + C
\( 1/|x|\sqrt{x^2-1} \) \hspace{2cm} \sec^{-1} x + C

You recognize that each integration formula came directly from a differentiation formula. The integral of the cosine is the sine, because the derivative of the sine is the cosine. For emphasis we list three derivatives above three integrals:

\[
\frac{d}{dx} (\text{constant}) = 0 \quad \frac{d}{dx} (x) = 1 \quad \frac{d}{dx} \left( \frac{x^{n+1}}{n+1} \right) = x^n
\]

\[
\int 0 \ dx = C \quad \int 1 \ dx = x + C \quad \int \frac{x^{n+1}}{n+1} \ dx = C
\]

There are two ways to make this list longer. One is to find the derivative of a new \( f(x) \). Then \( f \) goes in one column and \( v = df/dx \) goes in the other column.\footnote{We will soon meet \( e^x \), which goes in both columns. It is \( f(x) \) and also \( v(x) \).} The other possibility is to use rules for derivatives to find rules for integrals. That is the way to extend the list, enormously and easily.

\textbf{RULES FOR INTEGRALS}

Among the rules for derivatives, three were of supreme importance. They were \textit{linearity}, the \textit{product rule}, and the \textit{chain rule}. Everything flowed from those three. In the
reverse direction (from \(v\) to \(f\)) this is still true. The three basic methods of differential calculus also dominate integral calculus:

- **Linearity of derivatives** \(\rightarrow\) **Linearity of integrals**
- **Product rule for derivatives** \(\rightarrow\) **Integration by parts**
- **Chain rule for derivatives** \(\rightarrow\) **Integrals by substitution**

The easiest is linearity, which comes first. Integration by parts will be left for Section 7.1. This section starts on substitutions, reversing the chain rule to make an integral simpler.

**LINEARITY OF INTEGRALS**

*What is the integral of \(v(x) + w(x)\)?* Add the two separate integrals. The graph of \(v + w\) has two regions below it, the area under \(v\) and the area from \(v\) to \(v + w\). Adding areas gives the sum rule. Suppose \(f\) and \(g\) are antiderivatives of \(v\) and \(w\):

- **Sum rule:** \(f + g\) is an antiderivative of \(v + w\)
- **Constant rule:** \(cf\) is an antiderivative of \(cv\)
- **Linearity:** \(af + bg\) is an antiderivative of \(av + bw\)

This is a case of overkill. The first two rules are special cases of the third, so logically the last rule is enough. However it is so important to deal quickly with constants—just "factor them out"—that the rule \(cv \leftrightarrow cf\) is stated separately. The proofs come from the linearity of derivatives: \([af + bg]' = af' + bg'\) which equals \(av + bw\). The rules can be restated with integral signs:

- **Sum rule:** \[
\int \left[ v(x) + w(x) \right] \, dx = \int v(x) \, dx + \int w(x) \, dx
\]
- **Constant rule:** \[
\int cv(x) \, dx = c \int v(x) \, dx
\]
- **Linearity:** \[
\int \left[ av(x) - bw(x) \right] \, dx = a \int v(x) \, dx + b \int w(x) \, dx
\]

*Note about the constant in \(f(x) - C\).* All antiderivatives allow the addition of a constant. For a combination like \(av(x) + bw(x)\), the antiderivative is \(af(x) + bg(x) + C\). The constants for each part combine into a single constant. To give all possible antiderivatives of a function, just remember to write "+ \(C\)" after one of them. The real problem is to find that one antiderivative.

**EXAMPLE 1** The antiderivative of \(v = x^2 + x^{-2}\) is \(f = x^3/3 + (x^{-1})/(-1) + C\).

**EXAMPLE 2** The antiderivative of \(6 \cos t + 7 \sin t\) is \(6 \sin t - 7 \cos t + C\).

**EXAMPLE 3** Rewrite \(\frac{1}{1 - \sin x}\) as \(\frac{1 - \sin x}{1 - \sin^2 x} = \frac{1 - \sin x}{\cos^2 x} = \sec^2 x - \sec x \tan x\).

The antiderivative is \(\tan x - \sec x + C\). That rewriting is done by a symbolic algebra code (or by you). Differentiation is often simple, so most people check that \(df \cdot dx = v(x)\).

**Question** How to integrate \(\tan^2 x\)?

**Method** Write it as \(\sec^2 x - 1\). **Answer** \(\tan x + x + C\).
We now present the most valuable technique in this section—substitution. To see the idea, you have to remember the chain rule:

\[ f(g(x)) \text{ has derivative } f'(g(x))(dg/dx) \]
\[ \sin x^2 \text{ has derivative } (\cos x^2)(2x) \]
\[ (x^3 + 1)^5 \text{ has derivative } 5(x^3 + 1)^4(3x^2) \]

If the function on the right is given, the function on the left is its antiderivative! There are two points to emphasize right away:

1. **Constants are no problem—they can always be fixed.** Divide by 2 or 15:

\[ \int x \cos(x^2) \, dx = \frac{1}{2} \sin(x^2) + C \]
\[ \int x^2(x^3 + 1)^4 \, dx = \frac{1}{15} (x^3 + 1)^5 + C \]

Notice the 2 from \( x^2 \), the 5 from the fifth power, and the 3 from \( x^3 \).

2. **Choosing the inside function \( g \) (or \( u \)) commits us to its derivative:**

the integral of \( 2x \cos x^2 \) is \( \sin x^2 + C \) \((g = x^2, \, dg/dx = 2x)\)

the integral of \( \cos x^2 \) is (failure) \((no \, dg/dx)\)

the integral of \( x^2 \cos x^2 \) is (failure) \((wrong \, dg/dx)\)

To substitute \( g \) for \( x^2 \), we need its derivative. The trick is to spot an inside function whose derivative is present. We can fix constants like 2 or 15, but otherwise \( dg/dx \) has to be there. **Very often the inside function \( g \) is written \( u \).** We use that letter to state the substitution rule, when \( f \) is the integral of \( v \):

\[ \int f(v(x)) \frac{du}{dx} \, dx = f(u(x)) + C. \tag{1} \]

**EXAMPLE 4** \[ \int \sin x \cos x \, dx = \frac{1}{2} \left( \sin x \right)^2 + C \quad u = \sin x \text{ (compare Example 6)} \]

**EXAMPLE 5** \[ \int \sin^2 x \cos x \, dx = \frac{3}{2} \left( \sin x \right)^3 + C \quad u = \sin x \]

**EXAMPLE 6** \[ \int \cos x \sin x \, dx = -\frac{1}{2} \left( \cos x \right)^2 + C \quad u = \cos x \text{ (compare Example 4)} \]

**EXAMPLE 7** \[ \int \tan^4 x \sec^2 x \, dx = \frac{1}{2} \left( \tan x \right)^5 + C \quad u = \tan x \]

The next example has \( u = x^2 - 1 \) and \( du/dx = 2x \). The key step is choosing \( u \):

**EXAMPLE 8** \[ \int x \, dx / \sqrt{x^2 - 1} = \sqrt{x^2 - 1} + C \quad \int x \sqrt{x^2 - 1} \, dx = \frac{1}{3} (x^2 - 1)^{3/2} + C \]

A shift of \( x \) (to \( x + 2 \)) or a multiple of \( x \) (rescaling to \( 2x \)) is particularly easy:

**EXAMPLES 9–10** \[ \int (x + 2)^3 \, dx = \frac{1}{4} (x + 2)^4 + C \quad \int \cos 2x \, dx = \frac{1}{2} \sin 2x + C \]

You will soon be able to do those in your sleep. Officially the derivative of \((x + 2)^4\) uses the chain rule. But the inside function \( u = x + 2 \) has \( du/dx = 1 \). The “1” is there automatically, and the graph shifts over—as in Figure 5.8b.

For Example 10 the inside function is \( u = 2x \). Its derivative is \( du/dx = 2 \). This
5.4 Indefinite Integrals and Substitutions

Fig. 5.8  Substituting \( u = x + 1 \) and \( u = 2x \) and \( u = x^2 \). The last graph has half of \( du/dx = 2x \).

required factor 2 is missing in \( \int \cos 2x \, dx \), but we put it there by multiplying and dividing by 2. Check the derivative of \( \frac{1}{2} \sin 2x \): the 2 from the chain rule cancels the \( \frac{1}{2} \). The rule for any nonzero constant is similar:

\[
\int v(x + c) \, dx = f(x + c) \quad \text{and} \quad \int v(cx) \, dx = \frac{1}{c} f(cx). \tag{2}
\]

Squeezing the graph by \( c \) divides the area by \( c \). Now \( 3x + 7 \) rescales and shifts:

**EXAMPLE 14** \( \int \cos(3x+7) \, dx = \frac{1}{3} \sin(3x+7) + C \)

**Remark on writing down the steps** When the substitution is complicated, it is a good idea to get \( du/dx \) where you need it. Here \( 3x^2 + 1 \) needs 6x:

\[
\int 7x(3x^2 + 1)^4 \, dx = \frac{7}{6} \int (3x^2 + 1)^4 6x \, dx = \frac{7}{6} \int u^4 \frac{du}{dx} \, dx
\]

Now integrate:

\[
\frac{7}{6} \frac{u^5}{5} + C = \frac{7}{6} \frac{(3x^2 + 1)^5}{5} + C. \tag{3}
\]

Check the derivative at the end. The exponent 5 cancels 5 in the denominator, 6x from the chain rule cancels 6, and 7x is what we started with.

**Remark on differentials** In place of \( (du/dx) \, dx \), many people just write \( du \):

\[
\int (3x^2 + 1)^4 6x \, dx = \int u^4 \, du = \frac{1}{5} u^5 + C. \tag{4}
\]

This really shows how substitution works. We switch from \( x \) to \( u \), and we also switch from \( dx \) to \( du \). The most common mistake is to confuse \( dx \) with \( du \). The factor \( du/dx \) from the chain rule is absolutely needed, to reach \( du \). The change of variables (dummy variables anyway!) leaves an easy integral, and then \( u \) turns back into \( 3x^2 + 1 \). Here are the four steps to substitute \( u \) for \( x \):

1. Choose \( u(x) \) and compute \( du/dx \)
2. Locate \( v(u) \) times \( du/dx \) times \( dx \), or \( v(u) \) times \( du \)
3. Integrate \( \int v(u) \, du \) to find \( f(u) + C \)
4. Substitute \( u(x) \) back into this antiderivative \( f \).

**EXAMPLE 12** \( \int (\cos \sqrt{x}) \, dx/2\sqrt{x} = \int \cos u \, du = \sin u + C = \sin \sqrt{x} + C \)

\( \text{(put in } u \text{) (integrate) (put back } x \text{)} \)

The choice of \( u \) must be right, to change everything from \( x \) to \( u \). With ingenuity, some remarkable integrals are possible. But most will remain impossible forever. The functions \( \cos x^2 \) and \( 1/\sqrt{4 - \sin^2 x} \) have no "elementary" antiderivative. Those integrals are well defined and they come up in applications—the latter gives the distance
Integrals around an ellipse. That can be computed to tremendous accuracy, but not to perfect accuracy.

The exercises concentrate on substitutions, which need and deserve practice. We give a nonexample—∫(x² + 1)² dx does not equal ∫(x² + 1)³—to emphasize the need for du/dx. Since 2x is missing, u = x² + 1 does not work. But we can fix up π:

\[ \int \sin \pi x \, dx = \int \sin u \frac{du}{\pi} = -\frac{1}{\pi} \cos u + C = -\frac{1}{\pi} \cos \pi x + C. \]

### 5.4 EXERCISES

#### Read-through questions

Finding integrals by substitution is the reverse of the __ rule. The derivative of (sin x)² is __. Therefore the antiderivative of __ is __. To compute ∫(1 + sin x)² cos x dx, substitute u = __. Then du/dx = __ so substitute du = __. In terms of u the integral is __ = __. Returning to x gives the final answer.

The best substitutions for ∫tan (x + 3) sec²(x + 3) dx and ∫(x² + 1)³ x dx are u = __ and u = __. Then du = __ and __. The antiderivative of v du/dx is __. ∫2x dx/(1 + x²)² leads to __, which we don't yet know. The integral ∫dx/(1 + x²) is known immediately as __.

Find the indefinite integrals in 1–20.

<table>
<thead>
<tr>
<th>Equation</th>
<th>Answer</th>
</tr>
</thead>
<tbody>
<tr>
<td>1 [ ∫ \sqrt{2} + x , dx ] (add + C)</td>
<td>[ ∫ √3 - x , dx ] (always + C)</td>
</tr>
<tr>
<td>3 [ ∫ (x + 1)^{1/2} , dx ]</td>
<td>[ ∫ (x + 1)^{-3} , dx ]</td>
</tr>
<tr>
<td>5 [ ∫ (x² + 1)^{1/2} , dx ]</td>
<td>[ ∫ \sqrt{1 - 3x} , dx ]</td>
</tr>
<tr>
<td>7 [ ∫ \cos³ x \sin x , dx ]</td>
<td>[ ∫ \cos x , dx ]</td>
</tr>
<tr>
<td>9 [ ∫ \cos³ 2x \sin 2x , dx ]</td>
<td>[ ∫ \cos³ x \sin 2x , dx ]</td>
</tr>
<tr>
<td>11 [ ∫ dt/\sqrt{1 - t²} ]</td>
<td>[ ∫ t\sqrt{1 - t²} , dt ]</td>
</tr>
<tr>
<td>13 [ ∫ t³ dt/\sqrt{1 + t²} ]</td>
<td>[ ∫ t³ \sqrt{1 + t²} , dt ]</td>
</tr>
<tr>
<td>15 [ ∫ (1 + \sqrt{x}) , dx/\sqrt{x} ]</td>
<td>[ ∫ (1 + x²)² \sqrt{x} , dx ]</td>
</tr>
<tr>
<td>17 [ ∫ \sec x \tan x , dx ]</td>
<td>[ ∫ \sec² x \tan² x , dx ]</td>
</tr>
<tr>
<td>19 [ ∫ \cos x \tan x , dx ]</td>
<td>[ ∫ \sin x , dx ]</td>
</tr>
</tbody>
</table>

In 21–32 find a function y(x) that solves the differential equation.

<table>
<thead>
<tr>
<th>Equation</th>
<th>Answer</th>
</tr>
</thead>
<tbody>
<tr>
<td>21 dy/dx = x² + \sqrt{x}</td>
<td>dy/dx = y² (try y = cx²)</td>
</tr>
<tr>
<td>23 dy/dx = \sqrt{1 - 2x}</td>
<td>dy/dx = 1/\sqrt{1 - 2x}</td>
</tr>
</tbody>
</table>

**True or false, when f is an antiderivative of v:**

(a) \[ ∫ f(x)(dv/dx) \, dx = f(v(x)) \] + C

(b) \[ ∫ v²(x) \, dx = \frac{1}{2} f^2(x) + C \]

(c) \[ ∫ v(x)(du/dx) \, dx = f(u(x)) + C \]

(d) \[ ∫ v(x)(dv/dx) \, dx = \frac{1}{4} f^3(x) + C \]

33 **True or false, when f is an antiderivative of v:**

(a) \[ ∫ f(v(x)) \, dx = f(u(x)) + C \]

(b) \[ ∫ v²(x) \, dx = \frac{1}{2} f^2(x) + C \]

(c) \[ ∫ f(u(x)) \, dx = f(u(x)) + C \]

(d) \[ ∫ v(x)(dv/dx) \, dx = \frac{1}{2} f^2(x) + C \]

34 **True or false, when f is an antiderivative of v:**

(a) \[ ∫ f(x)(dv/dx) \, dx = \frac{1}{2} f^2(x) + C \]

(b) \[ ∫ v(v(x))(dv/dx) \, dx = f(v(x)) + C \]

(c) \[ ∫ f(v(x)) \, dx = f(v(x)) + C \]

35 If df/dx = v(x) then \[ ∫ v(x - 1) \, dx = \] and \[ ∫ v(x/2) \, dx = \] .

36 If df/dx = v(x) then \[ ∫ v(2x - 1) \, dx = \] and \[ ∫ v(x²) \, dx = \] .

37 \[ \frac{x²}{1 + x²} = 1 - \frac{1}{1 + x²} \] so \[ ∫ \frac{x² \, dx}{1 + x²} = \] .

38 \[ ∫ (x² + 1)^{1/2} \, dx \] is not \[ \frac{3}{2}(x² + 1)^{3/2} \] but .

39 \[ ∫ 2x \, dx/(x² + 1) \] is \[ ∫ du \] which will soon be ln u.

40 Show that \[ ∫ 2x \, dx/(1 + x²)³ = \] .

41 The acceleration \[ d²f/dt² = 9.8 \] gives \[ f(t) = \] (two integration constants).

42 The solution to \[ d²y/dx² = 0 \] is \[ f(t) \] (four constants).

43 If \[ f(t) \] is an antiderivative of \[ v(t) \], find antiderivatives of

(a) \[ v(t + 3) \]  (b) \[ v(t) + 3 \]  (c) \[ 3v(t) \]  (d) \[ v(3t) \].
The integral of \( v(x) \) is an antiderivative \( f(x) \) plus a constant \( C \). This section takes two steps. First, we choose \( C \). Second, we construct \( f(x) \). The object is to define the integral—in the most frequent case when a suitable \( f(x) \) is not directly known.

The indefinite integral contains \( " + C. \)" The constant is not settled because \( f(x) + C \) has the same slope for every \( C \). When we care only about the derivative, \( C \) makes no difference. When the goal is a number—a definite integral—\( C \) can be assigned a definite value at the starting point.

For mileage traveled, we subtract the reading at the start. This section does the same for area. Distance is \( f(t) \) and area is \( f(x) \)—while the definite integral is \( f(b) - f(a) \). Don’t pay attention to \( t \) or \( x \), pay attention to the great formula of integral calculus:

\[
\int_a^b v(t) \, dt = \int_a^b v(x) \, dx = f(b) - f(a).
\]  

Viewpoint 1: When \( f \) is known, the equation gives the area from \( a \) to \( b \).

Viewpoint 2: When \( f \) is not known, the equation defines \( f \) from the area.

For a typical \( v(x) \), we can’t find \( f(x) \) by guessing or substitution. But still \( v(x) \) has an "area" under its graph—and this yields the desired integral \( f(x) \).

Most of this section is theoretical, leading to the definition of the integral. You may think we should have defined integrals before computing them—which is logically true. But the idea of area (and the use of rectangles) was already pretty clear in our first examples. Now we go much further. Every continuous function \( v(x) \) has an integral (also some discontinuous functions). Then the Fundamental Theorem completes the circle: The integral leads back to \( df/dx = v(x) \). The area up to \( x \) is the antiderivative that we couldn’t otherwise discover.

**THE CONSTANT OF INTEGRATION**

Our goal is to turn \( f(x) + C \) into a definite integral—the area between \( a \) and \( b \). The first requirement is to have area = zero at the start:

\[
f(a) + C = \text{starting area} = 0 \quad \text{so} \quad C = -f(a). \tag{2}
\]

For the area up to \( x \) (moving endpoint, indefinite integral), use \( t \) as the dummy variable:

the area from \( a \) to \( x \) is \( \int_a^x v(t) \, dt = f(x) - f(a) \) (indefinite integral)

the area from \( a \) to \( b \) is \( \int_a^b v(x) \, dx = f(b) - f(a) \) (definite integral)

**EXAMPLE 1** The area under the graph of \( 5(x + 1)^4 \) from \( a \) to \( b \) has \( f(x) = (x + 1)^5 \):

\[
\int_a^b 5(x + 1)^4 \, dx = (x + 1)^5 \bigg|_a^b = (b + 1)^5 - (a + 1)^5.
\]

The calculation has two separate steps—first find \( f(x) \), then substitute \( b \) and \( a \). After the first step, check that \( df/dx = v \). The upper limit in the second step gives plus \( f(b) \), the lower limit gives minus \( f(a) \). Notice the brackets (or the vertical bar):

\[
f(x) \bigg|_a^b = f(b) - f(a) \quad \text{and} \quad \Big[ \cos x \Big]_0^2 = \cos 2 - 1.
\]

Changing the example to \( f(x) = (x + 1)^5 - 1 \) gives an equally good antiderivative—
and now \( f(0) = 0 \). But \( f(b) - f(a) \) stays the same, because the \(-1\) disappears:

\[
\left[ (x + 1)^5 - 1 \right]_a^b = ((b + 1)^5 - 1) - ((a + 1)^5 - 1) = (b + 1)^5 - (a + 1)^5.
\]

**EXAMPLE 2**  When \( v = 2x \sin x^2 \) we recognize \( f = -\cos x^2 \). The area from 0 to 3 is

\[
\int_0^3 2x \sin x^2 \, dx = -\cos x^2 \bigg|_0^3 = -\cos 3 + \cos 0.
\]

The upper limit copies the minus sign. The lower limit gives \(-(-\cos 0)\), which is \(+\cos 0\). That example shows the right form for solving exercises on definite integrals.

Example 2 jumped directly to \( f(x) = -\cos x^2 \). But most problems involving the chain rule go more slowly—by substitution. Set \( u = x^2 \), with \( \frac{du}{dx} = 2x \):

\[
\int_0^3 2x \sin x^2 \, dx = \int_0^3 \sin u \, \frac{du}{dx} \, dx = \int_0^3 \sin u \, du. \tag{3}
\]

We need new limits when \( u \) replaces \( x^2 \). Those limits on \( u \) are \( a^2 \) and \( b^2 \). (In this case \( a^2 = 0^2 \) and \( b^2 = 3^2 = 9 \).) If \( x \) goes from \( a \) to \( b \), then \( u \) goes from \( u(a) \) to \( u(b) \).

\[
\int_a^b v(u(x)) \, \frac{du}{dx} \, dx = \int_{u(a)}^{u(b)} v(u) \, du = f(u(b)) - f(u(a)). \tag{4}
\]

**EXAMPLE 3**

\[
\int_{-1}^1 (x^2 + 5)^3 \, dx = \int_{-5}^{5} u^3 \, du = \frac{u^4}{4} \bigg|_{-5}^{5} = \frac{6^4 - 5^4}{8}.
\]

In this case \( u = x^2 + 5 \). Therefore \( du/dx = 2x \) (or \( du = 2x \, dx \) for differentials). We have to account for the missing 2. The integral is \( \frac{1}{2}u^4 \). The limits on \( u \) are \( u(0) = 0^2 + 5 \) and \( u(1) = 1^2 + 5 \). That is why the \( u \)-integral goes from 5 to 6. The alternative is to find \( f(x) = \frac{1}{2}(x^2 + 5)^4 \) in one jump (and check it).

**EXAMPLE 4**  \( \int_{-1}^{1} \sin x^2 \, dx = ??? \) (no elementary function gives this integral).

If we try \( \cos x^2 \), the chain rule produces an extra \( 2x \)—no adjustment will work. Does \( \sin x^2 \) still have an antiderivative? Yes! Every continuous \( v(x) \) has an \( f(x) \). Whether \( f(x) \) has an algebraic formula or not, we can write it as \( \int v(x) \, dx \). To define that integral, we now take the limit of rectangular areas.

**INTEGRALS AS LIMITS OF "RIEMANN SUMS"**

We have come to the definition of the integral. The chapter started with the integrals of \( x \) and \( x^2 \), from formulas for \( 1 + \cdots + n \) and \( 1^2 + \cdots + n^2 \). We will not go back to those formulas. But for other functions, too irregular to find exact sums, the rectangular areas also approach a limit.

That limit is the integral. This definition is a major step in the theory of calculus. It can be studied in detail, or understood in principle. The truth is that the definition is not so painful—you virtually know it already.

**Problem** Integrate the continuous function \( v(x) \) over the interval \([a, b]\).

**Step 1** Split \([a, b]\) into \( n \) subintervals \([a, x_1], [x_1, x_2], \ldots, [x_{n-1}, b]\).

The "meshpoints" \( x_1, x_2, \ldots \) divide up the interval from \( a \) to \( b \). The endpoints are \( x_0 = a \) and \( x_n = b \). The length of subinterval \( k \) is \( \Delta x_k = x_k - x_{k-1} \). In that smaller interval, the minimum of \( v(x) \) is \( m_k \). The maximum is \( M_k \).
5.5 The Definite Integral

Now construct rectangles. The "lower rectangle" over interval \( k \) has height \( m_k \). The "upper rectangle" reaches to \( M_k \). Since \( v \) is continuous, there are points \( x_{\min} \) and \( x_{\max} \) where \( v = m_k \) and \( v = M_k \) (extreme value theorem). The graph of \( v(x) \) is in between.

Important: The area under \( v(x) \) contains the area "s" of the lower rectangles:

\[
\int_{a}^{b} v(x) \, dx \geq m_1 \Delta x_1 + m_2 \Delta x_2 + \cdots + m_n \Delta x_n = s. \tag{5}
\]

The area under \( v(x) \) is contained in the area "S" of the upper rectangles:

\[
\int_{a}^{b} v(x) \, dx \leq M_1 \Delta x_1 + M_2 \Delta x_2 + \cdots + M_n \Delta x_n = S. \tag{6}
\]

The lower sum \( s \) and the upper sum \( S \) were computed earlier in special cases—when \( v \) was \( x \) or \( x^2 \) and the spacings \( \Delta x \) were equal. Figure 5.9a shows why \( s \leq \text{area} \leq S \).

![Fig. 5.9 Area of lower rectangles = s. Upper sum S includes top pieces. Riemann sum S* is in between.](image)

Notice an important fact. When a new dividing point \( x' \) is added, the lower sum increases. The minimum in one piece can be greater (see second figure) than the original \( m_k \). Similarly the upper sum decreases. The maximum in one piece can be below the overall maximum. As new points are added, \( s \) goes up and \( S \) comes down. So the sums come closer together:

\[
s \leq s' \leq S' \leq S. \tag{7}
\]

I have left space in between for the curved area—the integral of \( v(x) \).

Now add more and more meshpoints in such a way that \( \Delta x_{\max} \to 0 \). The lower sums increase and the upper sums decrease. They never pass each other. If \( v(x) \) is continuous, those sums close in on a single number \( A \). That number is the definite integral—the area under the graph.

**DEFINITION** The area \( A \) is the common limit of the lower and upper sums:

\[
s \to A \text{ and } S \to A \text{ as } \Delta x_{\max} \to 0. \tag{8}
\]

This limit \( A \) exists for all continuous \( v(x) \), and also for some discontinuous functions. When it exists, \( A \) is the "Riemann integral" of \( v(x) \) from \( a \) to \( b \).

**REMARKS ON THE INTEGRAL**

As for derivatives, so for integrals: The definition involves a limit. Calculus is built on limits, and we always add "if the limit exists." That is the delicate point. I hope the next five remarks (increasingly technical) will help to distinguish functions that are Riemann integrable from functions that are not.

**Remark 1** The sums \( s \) and \( S \) may fail to approach the same limit. A standard example has \( V(x) = 1 \) at all fractions \( x = p/q \), and \( V(x) = 0 \) at all other points. Every
interval contains rational points (fractions) and irrational points (nonrepeating decimals). Therefore \( m_k = 0 \) and \( M_k = 1 \). The lower sum is always \( s = 0 \). The upper sum is always \( S = b - a \) (the sum of 1's times \( \Delta x \)'s). The gap in equation (7) stays open. This function \( V(x) \) is not Riemann integrable. The area under its graph is not defined (at least by Riemann—see Remark 5).

**Remark 2** The step function \( U(x) \) is discontinuous but still integrable. On every interval the minimum \( m_k \) equals the maximum \( M_k \)—except on the interval containing the jump. That jump interval has \( m_k = 0 \) and \( M_k = 1 \). But when we multiply by \( \Delta x_k \), and require \( \Delta x_{\text{max}} \to 0 \), the difference between \( s \) and \( S \) goes to zero. The area under a step function is clear—the rectangles fit exactly.

**Remark 3** With patience another key step could be proved: If \( s \to A \) and \( S \to A \) for one sequence of meshpoints, then this limit \( A \) is approached by every choice of meshpoints with \( \Delta x_{\text{max}} \to 0 \). The integral is the lower bound of all upper sums \( S \), and it is the upper bound of all possible \( s \)—provided those bounds are equal. The gap must close, to define the integral.

The same limit \( A \) is approached by “in-between rectangles.” The height \( v(x^*_k) \) can be computed at any point \( x^*_k \) in subinterval \( k \). See Figures 5.9c and 5.10. Then the total rectangular area is a “Riemann sum” between \( s \) and \( S \):

\[
S^* = v(x^*_1)\Delta x_1 + v(x^*_2)\Delta x_2 + \cdots + v(x^*_n)\Delta x_n.
\]  
(9)

We cannot tell whether the true area is above or below \( S^* \). Very often \( A \) is closer to \( S^* \) than to \( s \) or \( S \). The midpoint rule takes \( x^*_k \) in the middle of its interval (Figure 5.10), and Section 5.8 will establish its extra accuracy. The extreme sums \( s \) and \( S \) are used in the definition while \( S^* \) is used in computation.

![Fig. 5.10 Various positions for \( x^*_k \) in the base. The rectangles have height \( v(x^*_k) \).](image)

**Remark 4** Every continuous function is Riemann integrable. The proof is optional (in my class), but it belongs here for reference. It starts with continuity at \( x^* \): “For any \( \varepsilon \) there is a \( \delta \). . . .” When the rectangles sit between \( x^* - \delta \) and \( x^* + \delta \), the bounds \( M_k \) and \( m_k \) differ by less than \( 2\varepsilon \). Multiplying by the base \( \Delta x_k \), the areas differ by less than \( 2\varepsilon(\Delta x_k) \). Combining all rectangles, the upper and lower sums differ by less than \( 2\varepsilon(\Delta x_1 + \Delta x_2 + \cdots + \Delta x_n) = 2\varepsilon(b-a) \).

As \( \varepsilon \to 0 \) we conclude that \( S \) comes arbitrarily close to \( s \). They squeeze in on a single number \( A \). The Riemann sums approach the Riemann integral, if \( v \) is continuous.

Two problems are hidden by that reasoning. One is at the end, where \( S \) and \( s \) come together. We have to know that the line of real numbers has no “holes,” so there is a number \( A \) to which these sequences converge. That is true.

Any increasing sequence, if it is bounded above, approaches a limit.

The decreasing sequence \( S \), bounded below, converges to the same limit. So \( A \) exists.

The other problem is about continuity. We assumed without saying so that the
width $2\delta$ is the same around every point $x^*$. We did not allow for the possibility that $\delta$ might approach zero where $v(x)$ is rapidly changing—in which case an infinite number of rectangles could be needed. Our reasoning requires that

$v(x)$ is uniformly continuous: $\delta$ depends on $\epsilon$ but not on the position of $x^*$.

For each $\epsilon$ there is a $\delta$ that works at all points in the interval. A continuous function on a closed interval is uniformly continuous. This fact (proof omitted) makes the reasoning correct, and $v(x)$ is integrable.

On an infinite interval, even $v = x^2$ is not uniformly continuous. It changes across a subinterval by $(x^* + \delta)^2 - (x^* - \delta)^2 = 4x^*\delta$. As $x^*$ gets larger, $\delta$ must get smaller—to keep $4x^*\delta$ below $\epsilon$. No single $\delta$ succeeds at all $x^*$. But on a finite interval $[0, b]$, the choice $\delta = \epsilon/4b$ works everywhere—so $v = x^2$ is uniformly continuous.

**Remark 5** If those four remarks were fairly optional, this one is totally at your discretion. Modern mathematics needs to integrate the zero-one function $V(x)$ in the first remark. Somehow $V$ has more 0's than 1's. The fractions (where $V(x) = 1$) can be put in a list, but the irrational numbers (where $V(x) = 0$) are “uncountable.” The integral ought to be zero, but Riemann's upper sums all involve $M_k = 1$.

Lebesgue discovered a major improvement. He allowed infinitely many subintervals (smaller and smaller). Then all fractions can be covered with intervals of total width $\epsilon$. (Amazing, when the fractions are packed so densely.) The idea is to cover $1/q, 2/q, \ldots, q/q$ by narrow intervals of total width $\epsilon/2^q$. Combining all $q = 1, 2, 3, \ldots$, the total width to cover all fractions is no more than $\epsilon(1 + 1/2 + 1/4 + \ldots) = \epsilon$. Since $V(x) = 0$ everywhere else, the upper sum $S$ is only $\epsilon$. And since $\epsilon$ was arbitrary, the “Lebesgue integral” is zero as desired.

That completes a fair amount of theory, possibly more than you want or need—but it is satisfying to get things straight. The definition of the integral is still being studied by experts (and so is the derivative, again to allow more functions). By contrast, the properties of the integral are used by everybody. Therefore the next section turns from definition to properties, collecting the rules that are needed in applications. They are very straightforward.

---

### Exercises

**Read-through questions**

In $\int_a^b v(t) \, dt = f(x) + C$, the constant $C$ equals ___a___. Then at $x = a$ the integral is ___b___. At $x = b$ the integral becomes ___c___. The notation $f(x)^x_a$ means ___d___. Thus $\cos x^x_a$ equals ___e___. Also $[\cos x + 3]^x_a$ equals ___f___, which shows why the antiderivative includes an arbitrary ___g___. Substituting $u = 2x - 1$ changes $\int_1^h \sqrt{2x - 1} \, dx$ into ___h___ (with limits on $u$).

The integral $\int_a^b v(x) \, dx$ can be defined for any ___i___ function $v(x)$, even if we can't find a simple ___j___. First the meshpoints $x_1, x_2, \ldots$ divide $[a, b]$ into subintervals of length $\Delta x = k$. The upper rectangle with base $\Delta x_k$ has height $M_k = l$. The upper sum $S$ is equal to ___m___. The lower sum $s$ is _n_. The ___o___ is between $s$ and $S$. As more meshpoints are added, $S$ ___p___ and $s$ ___q___. If $S$ and $s$ approach the same ___r___, that defines the integral. The intermediate sums $S^*$, named after ___t___, use rectangles of height $v(x^*)$. Here $x^*_n$ is any point between ___u___, and $S^* = v(x)$ approaches the area.

If $v(x) = df/dx$, what constants $C$ make $1-10$ true?

1. $\int_2^1 v(x) \, dx = f(b) + C$
2. $\int_b^a v(x) \, dx = f(a) + C$
3. $\int_b^a v(t) \, dt = -f(x) + C$
4. $\int_{\sin x}^x v(x) \cos x \, dx = f(\sin b) + C$
5. $\int_1^t v(t) \, dt = f(t) + C$ (careful)
6. $df/dx = v(x) + C$
7. $\int_0^1 (x^2 - 1)^3 2x \, dx = \int_{x^*}^c u^3 \, du$. 

---
8 \[ \int_0^1 v(t) \, dt = f(x^2) + C \]
9 \[ \int_0^1 v(-x) \, dx = C \text{ change } -x \text{ to } t; \text{ also } dx \text{ and limits} \]
10 \[ \int_0^1 v(x) \, dx = C \int_0^1 v(2t) \, dt. \]

Choose \( u(x) \) in 11–18 and change limits. Compute the integral in 11–16.

11 \[ \int_0^1 (x^2 + 1)^4 \, dx \]
12 \[ \int_0^{\pi/2} \sin^9 x \cos x \, dx \]
13 \[ \int_0^1 \tan x \sec^2 x \, dx \]
14 \[ \int_0^1 x^{x+1} \, dx \text{ (take } u = x^2) \]
15 \[ \int_0^1 \sec^2 x \tan x \, dx \]
16 \[ \int_0^1 x \, dx / \sqrt{1-x^2} \]
17 \[ \int_1^2 dx/x \text{ (take } u = 1/x) \]
18 \[ \int_0^1 x^3(1-x)^3 \, dx \text{ (take } u = 1-x) \]

With \( \Delta x = \frac{1}{4} \) in 19–22, find the maximum \( M_k \) and minimum \( m_k \) and upper and lower sums \( S \) and \( s \).

19 \[ \int_0^1 (x^2 + 1)^4 \, dx \]
20 \[ \int_0^1 \sin 2\pi x \, dx \]
21 \[ \int_0^1 x^3 \, dx \]
22 \[ \int_{-1}^1 x \, dx. \]

23 Repeat 19 and 20 with \( \Delta x = \frac{1}{4} \) and compare with the correct answer.

24 The difference \( S - s \) in 21 is the area \( 2^\Delta x \) of the far right rectangle. Find \( \Delta x \) so that \( S < 4.001 \).

25 If \( v(x) \) is increasing for \( a \leq x \leq b \), the difference \( S - s \) is the area of the rectangle minus the area of the rectangle. Those areas approach zero. So every increasing function on \([a, b]\) is Riemann integrable.

26 Find the Riemann sum \( S^* \) for \( V(x) \) in Remark 1, when \( \Delta x = 1/n \) and each \( x_n^* \) is the midpoint. This \( S^* \) is well-behaved but still \( V(x) \) is not Riemann integrable.

27 \( W(x) \) equals \( 1 \) at \( x = \frac{1}{2}, \frac{1}{4}, \ldots, \) and elsewhere \( W(x) = 0 \). For \( \Delta x = 0.01 \) find the upper sum \( S \). Is \( W(x) \) integrable?

28 Suppose \( M(x) \) is a multistep function with jumps of \( \frac{1}{2}, \frac{1}{4}, \ldots, \) at the points \( x = \frac{1}{2}, \frac{1}{4}, \ldots \). Draw a rough graph with \( M(0) = 0 \) and \( M(1) = 1 \). With \( \Delta x = \frac{1}{2} \) find \( S \) and \( s \).

29 For \( M(x) \) in Problem 28 find the difference \( S - s \) (which approaches zero as \( \Delta x \to 0 \)). What is the area under the graph?

30 If \( df/dx = -v(x) \) and \( f(1) = 0 \), explain \( f(x) = \int_1^x v(t) \, dt \).

31 (a) If \( df/dx = +v(x) \) and \( f(1) = 3 \), find \( f(x) \).
(b) If \( df/dx = +v(x) \) and \( f(3) = 0 \), find \( f(x) \).

32 In your own words define the integral of \( v(x) \) from \( a \) to \( b \).

33 True or false, with reason or example.
(a) Every continuous \( v(x) \) has an antiderivative \( f(x) \).
(b) If \( v(x) \) is not continuous, \( S \) and \( s \) approach different limits.
(c) If \( S \) and \( s \) approach \( A \) as \( \Delta x \to 0 \), then all Riemann sums \( S^* \) in equation (9) also approach \( A \).
(d) If \( v_1(x) + v_2(x) = v_3(x) \), their upper sums satisfy \( S_1 + S_2 = S_3 \).
(e) If \( v_1(x) + v_2(x) = v_3(x) \), their Riemann sums at the midpoints \( x_n^* \) satisfy \( S_1^* + S_2^* = S_3^* \).
(f) The midpoint sum is the average of \( S \) and \( s \).
(g) One \( x_n^* \) in Figure 5.10 gives the exact area.

### 5.6 Properties of the Integral and Average Value

The previous section reached the definition of \( \int_a^b v(x) \, dx \). But the subject cannot stop there. The integral was defined in order to be used. Its properties are important, and its applications are even more important. The definition was chosen so that the integral has properties that make the applications possible.

One direct application is to the average value of \( v(x) \). The average of \( n \) numbers is clear, and the integral extends that idea—it produces the average of a whole continuum of numbers \( v(x) \). This develops from the last rule in the following list (Property 7). We now collect together seven basic properties of definite integrals.

The addition rule for \( \int [u(x) + w(x)] \, dx \) will not be repeated—even though this property of linearity is the most fundamental. We start instead with a different kind of addition. There is only one function \( v(x) \), but now there are two intervals.

The integral from \( a \) to \( b \) is added to its neighbor from \( b \) to \( c \). Their sum is the integral from \( a \) to \( c \). That is the first (not surprising) property in the list.
5.6 Properties of the Integral and Average Value

Property 1 Areas over neighboring intervals add to the area over the combined interval:

\[ \int_a^b v(x) \, dx + \int_b^c v(x) \, dx = \int_a^c v(x) \, dx. \]  

This sum of areas is graphically obvious (Figure 5.11a). It also comes from the formal definition of the integral. Rectangular areas obey (1)—with a meshpoint at \( x = b \) to make sure. When \( \Delta x_{\text{max}} \) approaches zero, their limits also obey (1). All the normal rules for rectangular areas are obeyed in the limit by integrals.

Property 1 is worth pursuing. It indicates how to define the integral when \( a = b \). The integral "from \( b \) to \( b \)" is the area over a point, which we expect to be zero. It is.

Property 2 \( \int_a^b v(x) \, dx = 0. \)

That comes from Property 1 when \( c = b \). Equation (1) has two identical integrals, so the one from \( b \) to \( b \) must be zero. Next we see what happens if \( c = a \)—which makes the second integral go from \( b \) to \( a \).

What happens when an integral goes backward? The "lower limit" is now the larger number \( b \). The "upper limit" \( a \) is smaller. Going backward reverses the sign:

Property 3 \( \int_a^b v(x) \, dx = -\int_b^a v(x) \, dx = f(a) - f(b). \)

Proof When \( c = a \) the right side of (1) is zero. Then the integrals on the left side must cancel, which is Property 3. In going from \( b \) to \( a \) the steps \( \Delta x \) are negative. That justifies a minus sign on the rectangular areas, and a minus sign on the integral (Figure 5.11b). Conclusion: Property 1 holds for any ordering of \( a, b, c \).

EXAMPLES \[ \int_x^0 t^2 \, dt = -\frac{x^3}{3} \quad \int_1^0 dt = -1 \quad \int_2^0 \frac{dt}{x} = 0 \]

Property 4 For odd functions \( \int_{-a}^a v(x) \, dx = 0. \) "Odd" means that \( v(-x) = -v(x) \).

For even functions \( \int_{-a}^a v(x) \, dx = 2 \int_0^a v(x) \, dx. \) "Even" means that \( v(-x) = +v(x) \).

The functions \( x, x^3, x^5, \ldots \) are odd. If \( x \) changes sign, these powers change sign. The functions \( \sin x \) and \( \tan x \) are also odd, together with their inverses. This is an important family of functions, and the integral of an odd function from \( -a \) to \( a \) equals zero. Areas cancel:

\[ \int_{-a}^a 6x^5 \, dx = x^6 \bigg|_{-a}^a = a^6 - (-a)^6 = 0. \]

If \( v(x) \) is odd then \( f(x) \) is even! All powers \( 1, x^2, x^4, \ldots \) are even functions. Curious fact: Odd function times even function is odd, but odd number times even number is even.

For even functions, areas add: \( \int_{-a}^a \cos x \, dx = \sin a - \sin(-a) = 2 \sin a. \)

Fig. 5.11 Properties 1–4: Add areas, change sign to go backward, odd cancels, even adds.
The next properties involve inequalities. If $v(x)$ is positive, the area under its graph is positive (not surprising). Now we have a proof: The lower sums $s$ are positive and they increase toward the area integral. So the integral is positive:

**Property 5** If $v(x) > 0$ for $a < x < b$ then $\int_a^b v(x) \, dx > 0$.

A positive velocity means a positive distance. A positive $v$ lies above a positive area. A more general statement is true. Suppose $v(x)$ stays between a lower function $l(x)$ and an upper function $u(x)$. Then the rectangles for $v$ stay between the rectangles for $l$ and $u$. In the limit, the area under $v$ (Figure 5.12) is between the areas under $l$ and $u$:

**Property 6** If $l(x) \leq v(x) \leq u(x)$ for $a \leq x \leq b$ then

$$\int_a^b l(x) \, dx \leq \int_a^b v(x) \, dx \leq \int_a^b u(x) \, dx.$$  \hspace{1cm} (2)

**EXAMPLE 1** \[ \cos t \leq 1 \quad \Rightarrow \quad \int_0^x \cos t \, dt \leq \int_0^x 1 \, dt \quad \Rightarrow \quad \sin x \leq x \]

**EXAMPLE 2** \[ 1 \leq \sec^2 t \quad \Rightarrow \quad \int_0^x 1 \, dt \leq \int_0^x \sec^2 t \, dt \quad \Rightarrow \quad x \leq \tan x \]

**EXAMPLE 3** Integrating $\frac{1}{1 + x^2} \leq 1$ leads to $\tan^{-1} x \leq x$.

All those examples are for $x > 0$. You may remember that Section 2.4 used geometry to prove $\sin h < h < \tan h$. Examples 1–2 seem to give new and shorter proofs. But I think the reasoning is doubtful. The inequalities were needed to compute the derivatives (therefore the integrals) in the first place.

![Fig. 5.12 Properties 5–7: $v$ above zero, $v$ between $l$ and $u$, average value (+ balances −).](image)

**Property 7** (**Mean Value Theorem for Integrals**) If $v(x)$ is continuous, there is a point $c$ between $a$ and $b$ where $v(c)$ equals the average value of $v(x)$:

$$v(c) = \frac{1}{b-a} \int_a^b v(x) \, dx = "\text{average value of } v(x)."$$ \hspace{1cm} (3)

This is the same as the ordinary Mean Value Theorem (for the derivative of $f(x)$):

$$f'(c) = \frac{f(b) - f(a)}{b-a} = "\text{average slope of } f."$$ \hspace{1cm} (4)

With $f' = v$, (3) and (4) are the same equation. But honesty makes me admit to a flaw in the logic. We need the Fundamental Theorem of Calculus to guarantee that $f(x) = \int_a^x v(t) \, dt$ really gives $f' = v$.

A direct proof of (3) places one rectangle across the interval from $a$ to $b$. Now raise the top of that rectangle, starting at $v_{\min}$ (the bottom of the curve) and moving up to $v_{\max}$ (the top of the curve). At some height the area will be just right—equal to the area under the curve. Then the rectangular area, which is $(b-a)$ times $v(c)$, equals the curved area $\int_a^b v(x) \, dx$. This is equation (3).
5.6 Properties of the Integral and Average Value

\[ u(x) = x \]
\[ u(x) = x^2 \]
\[ u(x) = \sin^2 x \]

Fig. 5.13 Mean Value Theorem for integrals: area \((b - a) = \text{average height} = u(c)\) at some \(c\).

That direct proof uses the **intermediate value theorem**: A continuous function \(v(x)\) takes on every height between \(v_{\min}\) and \(v_{\max}\). At some point (at two points in Figure 5.12) the function \(v(x)\) equals its own average value.

**Figure 5.13 shows equal areas above and below the average height** \(v(c) = u_{\text{ave}}\).

**EXAMPLE 4** The average value of an odd function is zero (between \(-1\) and 1):

\[
\frac{1}{2} \int_{-1}^{1} x \, dx = \frac{x^2}{4} \bigg|_{-1}^{1} = \frac{1}{4} - \frac{1}{4} = 0
\]

For once we know \(c\). It is the center point \(x = 0\), where \(v(c) = v_{\text{ave}} = 0\).

**EXAMPLE 5** The average value of \(x^2\) is \(\frac{1}{3}\) (between 1 and \(-1\)):

\[
\frac{1}{2} \int_{-1}^{1} x^2 \, dx = \frac{x^3}{6} \bigg|_{-1}^{1} = \frac{1}{6} - \left(-\frac{1}{6}\right) = \frac{1}{3}
\]

Where does this function \(x^2\) equal its average value \(\frac{1}{3}\)? That happens when \(c^2 = \frac{1}{3}\), so \(c\) can be either of the points \(1/\sqrt{3}\) and \(-1/\sqrt{3}\) in Figure 5.13b. Those are the **Gauss points**, which are terrific for numerical integration as Section 5.8 will show.

**EXAMPLE 6** The average value of \(\sin^2 x\) over a period (zero to \(\pi\)) is \(\frac{1}{4}\):

\[
\frac{1}{\pi} \int_{0}^{\pi} \sin^2 x \, dx = \frac{x - \sin x \cos x}{2\pi} \bigg|_{0}^{\pi} = \frac{1}{2}
\]

The point \(c\) is \(\pi/4\) or \(3\pi/4\), where \(\sin^2 c = \frac{1}{2}\). The graph of \(\sin^2 x\) oscillates around its average value \(\frac{1}{4}\). See the figure or the formula:

\[ \sin^2 x = \frac{1}{2} - \frac{1}{2} \cos 2x. \] (5)

The steady term is \(\frac{1}{2}\), the oscillation is \(-\frac{1}{2} \cos 2x\). The integral is \(f(x) = \frac{1}{2}x - \frac{1}{4} \sin 2x\), which is the same as \(\frac{1}{2}x - \frac{1}{2} \sin x \cos x\). This integral of \(\sin^2 x\) will be seen again. Please verify that \(df/dx = \sin^2 x\).

**THE AVERAGE VALUE AND EXPECTED VALUE**

The "average value" from \(a\) to \(b\) is the integral divided by the length \(b - a\). This was computed for \(x\) and \(x^2\) and \(\sin^2 x\), but not explained. It is a major application of the integral, and it is guided by the ordinary average of \(n\) numbers:

\[ u_{\text{ave}} = \frac{1}{b - a} \int_{a}^{b} v(x) \, dx \quad \text{comes from} \quad u_{\text{ave}} = \frac{1}{n} (v_1 + v_2 + \cdots + v_n) . \]

Integration is parallel to summation! Sums approach integrals. Discrete averages
approach continuous averages. The average of $\frac{1}{4}, \frac{1}{2}, \frac{3}{4} \text{ is } \frac{1}{2}$. The average of $\frac{1}{4}, \frac{1}{2}, \frac{1}{4}, \frac{3}{4}$, $\frac{1}{2}$ is $\frac{1}{2}$. The average of $n$ numbers from $1/n$ to $n/n$ is

$$v_{ave} = \frac{1}{n} \left( \frac{1}{n} + \frac{2}{n} + \cdots + \frac{n}{n} \right) = \frac{n + 1}{2n}.$$  \hfill (7)

The middle term gives the average, when $n$ is odd. Or we can do the addition. As $n \to \infty$ the sum approaches an integral (do you see the rectangles?). The ordinary average of numbers becomes the continuous average of $v(x) = x$:

$$\frac{n + 1}{2n} \to \frac{1}{2} \quad \text{and} \quad \int_0^1 x \, dx = \frac{1}{2} \left( \text{note } \frac{1}{b-a} = 1 \right)$$

In ordinary language: "The average value of the numbers between 0 and 1 is $\frac{1}{2}$." Since a whole continuum of numbers lies between 0 and 1, that statement is meaningless until we have integration.

The average value of the squares of those numbers is $(x^2)_{ave} = \int_0^1 x^2 \, dx/(b-a) = \frac{1}{3}$. If you pick a number randomly between 0 and 1, its expected value is $\frac{1}{3}$ and its expected square is $\frac{1}{9}$.

To me that sentence is a puzzle. First, we don't expect the number to be exactly $\frac{1}{3}$—so we need to define "expected value." Second, if the expected value is $\frac{1}{3}$, why is the expected square equal to $\frac{1}{9}$ instead of $\frac{1}{3}$? The ideas come from probability theory, and calculus is leading us to continuous probability. We introduce it briefly here, and come back to it in Chapter 8.

**PREDICTABLE AVERAGES FROM RANDOM EVENTS**

Suppose you throw a pair of dice. The outcome is not predictable. Otherwise why throw them? **But the average over more and more throws is totally predictable.** We don't know what will happen, but we know its probability.

For dice, we are adding two numbers between 1 and 6. The outcome is between 2 and 12. The probability of 2 is the chance of two ones: $(1/6)(1/6) = 1/36$. Beside each outcome we can write its probability:

$$2 \left( \frac{1}{36} \right) 3 \left( \frac{2}{36} \right) 4 \left( \frac{3}{36} \right) 5 \left( \frac{4}{36} \right) 6 \left( \frac{5}{36} \right) 7 \left( \frac{6}{36} \right) 8 \left( \frac{5}{36} \right) 9 \left( \frac{4}{36} \right) 10 \left( \frac{3}{36} \right) 11 \left( \frac{2}{36} \right) 12 \left( \frac{1}{36} \right)$$

To repeat, one roll is unpredictable. Only the probabilities are known, and they add to 1. (Those fractions add to $36/36$; all possibilities are covered.) The total from a million rolls is even more unpredictable—it can be anywhere between 2,000,000 and 12,000,000. Nevertheless the average of those million outcomes is almost completely predictable. This expected value is found by adding the products in that line above:

**Expected value:** multiply (outcome) times (probability of outcome) and add:

$$\frac{2}{36} + \frac{6}{36} + \frac{12}{36} + \frac{20}{36} + \frac{30}{36} + \frac{42}{36} + \frac{40}{36} + \frac{36}{36} + \frac{30}{36} + \frac{22}{36} + \frac{12}{36} = \frac{7}{1}.$$  \hfill (7)

If you throw the dice 1000 times, and the average is not between 6.9 and 7.1, you get an A. Use the random number generator on a computer and round off to integers.

**Now comes continuous probability.** Suppose all numbers between 2 and 12 are equally probable. This means all numbers—not just integers. What is the probability of hitting the particular number $x = \pi$? **It is zero!** By any reasonable measure, $\pi$ has
Properties of the Integral and Average Value

no chance to occur. In the continuous case, every \( x \) has probability zero. But an interval of \( x \)'s has a nonzero probability:

- the probability of an outcome between 2 and 3 is \( 1/10 \)
- the probability of an outcome between \( x \) and \( x + \Delta x \) is \( \Delta x/10 \)

To find the average, add up each outcome times the probability of that outcome. First divide 2 to 12 into intervals of length \( \Delta x = 1 \) and probability \( p = 1/10 \). If we round off \( x \), the average is 6.5:

\[
2 \left( \frac{1}{10} \right) + 3 \left( \frac{1}{10} \right) + \ldots + 11 \left( \frac{1}{10} \right) = 6.5.
\]

Here all outcomes are integers (as with dice). It is more accurate to use 20 intervals of length 1/2 and probability 1/20. The average is 6.4, and you see what is coming. These are rectangular areas (Riemann sums). As \( \Delta x \to 0 \) they approach an integral.

The probability of an outcome between \( x \) and \( x + dx \) is \( p(x) \, dx \), and this problem has \( p(x) = 1/10 \). The average outcome in the continuous case is not a sum but an integral:

\[
E(x) = \int_{2}^{12} x \, p(x) \, dx = \int_{2}^{12} \frac{x \, dx}{10} = \frac{x^2}{20} \bigg|_{2}^{12} = 7.
\]

That is a big jump. From the point of view of integration, it is a limit of sums. From the point of view of probability, the chance of each outcome is zero but the probability density at \( x \) is \( p(x) = 1/10 \). The integral of \( p(x) \) is 1, because some outcome must happen. The integral of \( xp(x) \) is \( x_{\text{ave}} = 7 \), the expected value. Each choice of \( x \) is random, but the average is predictable.

This completes a first step in probability theory. The second step comes after more calculus. Decaying probabilities use \( e^{-x} \) and \( e^{-x^2} \) — then the chance of a large \( x \) is very small. Here we end with the expected values of \( x^n \) and \( 1/\sqrt{x} \) and \( 1/x \), for a random choice between 0 and 1 (so \( p(x) = 1 \)):

\[
E(x^n) = \int_{0}^{1} x^n \, dx = \frac{1}{n + 1} \quad E \left( \frac{1}{\sqrt{x}} \right) = \int_{0}^{1} \frac{dx}{\sqrt{x}} = 2 \quad E \left( \frac{1}{x} \right) = \int_{0}^{1} \frac{dx}{x} = \infty (!)
\]

A confusion about "expected" class size

A college can advertise an average class size of 29, while most students are in large classes most of the time. I will show quickly how that happens.

Suppose there are 95 classes of 20 students and 5 classes of 200 students. The total enrollment in 100 classes is 1900 + 1000 = 2900. A random professor has expected class size 29. But a random student sees it differently. The probability is 1900/2900 of being in a small class and 1000/2900 of being in a large class. Adding class size times probability gives the expected class size for the student:

\[
(20) \left( \frac{1900}{2900} \right) + (200) \left( \frac{1000}{2900} \right) = 82 \text{ students in the class.}
\]

Similarly, the average waiting time at a restaurant seems like 40 minutes (to the customer). To the hostess, who averages over the whole day, it is 10 minutes. If you came at a random time it would be 10. But if you are a random customer it is 40.

Traffic problems could be eliminated by raising the average number of people per car to 2.5, or even 2. But that is virtually impossible. Part of the problem is the
difference between (a) the percentage of cars with one person and (b) the percentage of people alone in a car. Percentage (b) is smaller. In practice, most people would be in crowded cars. See Problems 37–38.

5.6 EXERCISES

Read-through questions

The integrals \( \int_0^a v(x) \, dx \) and \( \int_0^b v(x) \, dx \) add to \( \boxed{a} \). The integral \( \int_0^1 v(x) \, dx \) equals \( \boxed{b} \). The reason is \( \boxed{c} \). If \( v(x) \leq x \) then \( \int_0^1 v(x) \, dx \leq \boxed{d} \). The average value of \( v(x) \) on the interval \( 1 \leq x \leq 9 \) is defined by \( \boxed{e} \). It is equal to \( v(c) \) at a point \( x = c \) which is \( \boxed{f} \). The rectangle across this interval with height \( v(c) \) has the same area as \( \boxed{g} \).

If \( V(x) < x \) then \( v(x) \, dx \leq \boxed{h} \).

The average value of \( v(x) \) on the interval \( 1 \leq x \leq 8 \) is defined by \( \boxed{i} \).

It is equal to \( v(c) \) at a point \( x = c \) which is \( \boxed{j} \).

The rectangle across this interval with height \( v(c) \) has the same area as \( \boxed{k} \).

The average value of \( u(x) = x + 1 \) on the interval \( 1 \leq x \leq 9 \) is \( \boxed{l} \).

If \( x \) is chosen from \( 1, 3, 5, 7 \) with equal probabilities $, its expected value (average) is \( \boxed{m} \).

The expected value of \( x^2 \) is \( \boxed{n} \).

If \( x \) is chosen from \( 1, 2, \ldots, 8 \) with probabilities \( \boxed{o} \), its expected value is \( \boxed{p} \).

If \( x \) is chosen from \( 1 \leq x \leq 9 \), the chance of hitting an integer is \( \boxed{q} \).

The chance of falling between \( x \) and \( x + dx \) is \( p(x) \, dx = \boxed{r} \).

The expected value \( E(x) \) is the integral \( \boxed{s} \).

It equals \( \boxed{t} \).

In 1–6 find the average value of \( v(x) \) between \( a \) and \( b \), and find all points \( c \) where \( v_{ave} = v(c) \).

1. \( v = x^4 \), \( a = -1 \), \( b = 1 \)  
   \( v = x^3 \), \( a = -1 \), \( b = 1 \)  
   \( v = \cos^2 x \), \( a = 0 \), \( b = \pi \)  
   \( v = \sqrt{x} \), \( a = 0 \), \( b = 4 \)  
   \( v = 1/x^2 \), \( a = 1 \), \( b = 2 \)  
   \( v = (\sin x)^2 \), \( a = -\pi \), \( b = \pi \).

2. \( v = x^2 \), \( a = 0 \), \( b = 2 \)  
   \( v = x^2/(1 + x^2) \)

3. \( v_{ave} \) on \( 0 \leq x \leq 1 \) and \( \boxed{t} \) on \( 1 \leq x \leq 3 \).

4. The ratio \( (f(b) - f(a))/(b - a) \) is the average value of \( f(x) \) on \( a \leq x \leq b \).

5. On the symmetric interval \( -1 \leq x \leq 1 \), \( v(x) - v_{ave} \) is an odd function.

6. If \( h(x) \leq v(x) \leq u(x) \) then \( dh/dx \leq dv/dx \leq du/dx \).

7. The average of \( v(x) \) from 0 to 2 plus the average from 2 to 4 equals the average from 0 to 4.

8. (a) Antiderivatives of even functions are odd functions.  
(b) Squares of odd functions are odd functions.

9. What number \( \bar{v} \) gives \( \int_0^b (v(x) - \bar{v}) \, dx = 0 \)?

10. If \( f(2) = 6 \) and \( f(6) = 2 \) then the average of \( df/dx \) from \( x = 2 \) to \( x = 6 \) is \( \boxed{u} \).

11. (a) The averages of \( \cos x \) and \( |\cos x| \) from 0 to \( \pi \) are \( \boxed{v} \).

12. The average of the numbers \( v_1, \ldots, v_n \) is \( \boxed{w} \) than the average of \( |v_1|, \ldots, |v_n| \).

13. Which property of integrals proves \( \int_0^1 v(x) \, dx \leq \int_0^1 |v(x)| \, dx \)?

14. Which property proves \( -\int_0^1 v(x) \, dx \leq \int_0^1 |v(x)| \, dx \)?

Together these are Property 8: \( \int_0^1 v(x) \, dx \leq \int_0^1 |v(x)| \, dx \).

15. What function has \( v_{ave} \) from 0 to \( x \) equal to \( \frac{1}{2} v(x) \) at all \( x \)? What functions have \( v_{ave} = v(x) \) at all \( x \)?

16. (a) If \( v(x) \) is increasing, explain from Property 6 why \( \int_0^x v(t) \, dt \leq x v(x) \) for \( x > 0 \).

(b) Take derivatives of both sides for a second proof.

17. The average of \( v(x) = 1/(1 + x^2) \) on the interval \( [0, b] \) approaches \( \boxed{v} \) as \( b \to \infty \). The average of \( V(x) = x^2/(1 + x^2) \) approaches \( \boxed{w} \).

18. If the positive numbers \( v_n \) approach zero as \( n \to \infty \) prove that their average \( (v_1 + \cdots + v_n)/n \) also approaches zero.

19. Find the average distance from \( x = a \) to points in the interval \( 0 < x < 2 \). Is the formula different if \( a < 2 \)?

20. (Computer experiment) Choose random numbers \( x \) between 0 and 1 until the average value of \( x^2 \) is between .333 and .334. How many values of \( x^2 \) are above and below? If possible repeat ten times.

21. A point \( P \) is chosen randomly along a semicircle (see figure: equal probability for equal arcs). What is the average distance \( y \) from the \( x \) axis? The radius is 1.

22. A point \( Q \) is chosen randomly between -1 and 1.

(a) What is the average distance \( Y \) up to the semicircle?

(b) Why is this different from Problem 27?

**Buffon needle**
29 (A classic way to compute π) A 2" needle is tossed onto a floor with boards 2" wide. Find the probability of falling across a crack. (This happens when \( \cos \theta > y \) = distance from midpoint of needle to nearest crack. In the rectangle \( 0 \leq \theta \leq \pi/2, 0 \leq y \leq 1 \), shade the part where \( \cos \theta > y \) and find the fraction of area that is shaded.)

30 If Buffon's needle has length 2x instead of 2, find the probability \( P(x) \) of falling across the same cracks.

31 If you roll three dice at once, what are the probabilities of each outcome between 3 and 18? What is the expected value?

32 If you choose a random point in the square \( 0 \leq x \leq 1, 0 \leq y \leq 1 \), what is the chance that its coordinates have \( y^2 \leq x? \)

33 The voltage \( V(t) = 220 \cos 2\pi t/60 \) has frequency 60 hertz and amplitude 220 volts. Find \( V_{\text{ave}} \) from 0 to \( t \).

34 (a) Show that \( u_{\text{even}}(x) = \frac{1}{2}(u(x) + u(-x)) \) is always even.
   (b) Show that \( u_{\text{odd}}(x) = \frac{1}{2}(u(x) - u(-x)) \) is always odd.

35 By Problem 34 or otherwise, write \((x + 1)^3 \) and \( 1/(x + 1) \) as an even function plus an odd function.

36 Prove from the definition of \( df/dx \) that it is an odd function if \( f(x) \) is even.

37 Suppose four classes have 6, 8, 10, and 40 students, averaging _____. The chance of being in the first class is _______. The expected class size (for the student) is \( E(x) = 6 \left( \frac{6}{64} \right) + 8 \left( \frac{8}{64} \right) + 10 \left( \frac{10}{64} \right) + 40 \left( \frac{40}{64} \right) = _____. \)

38 With groups of sizes \( x_1, \ldots, x_n \) adding to \( G \), the average size is _______. The chance of an individual belonging to group 1 is _______. The expected size of his or her group is \( E(x) = x_1(x_1/G) + \ldots + x_n(x_n/G) \). Prove \( \sum x_i^2/G \geq G/n \).

39 True or false, 15 seconds each:
   (a) If \( f(x) \leq g(x) \) then \( df/dx \leq dg/dx \).
   (b) If \( df/dx \leq dg/dx \) then \( f(x) \leq g(x) \).
   (c) \( x^2 \) is odd if \( v(x) \) is even.
   (d) If \( v_{\text{ave}} \leq w_{\text{ave}} \) on all intervals then \( v(x) \leq w(x) \) at all points.

40 If \( v(x) = \begin{cases} 2x & \text{for } x < 3 \\ -x^2 & \text{for } x > 3 \end{cases} \) then \( f(x) = \begin{cases} x^2 & \text{for } x < 3 \\ -x^2 & \text{for } x > 3 \end{cases} \)

Thus \( \int_0^3 v(x) \, dx = f(4) - f(0) = -16 \). Correct the mistake.

41 If \( v(x) = |x - 2| \) find \( f(x) \). Compute \( \int_0^3 v(x) \, dx \).

42 Why are there equal areas above and below \( v_{\text{ave}}? \)

---

### 5.7 The Fundamental Theorem and Its Applications

When the endpoints are fixed at \( a \) and \( b \), we have a *definite* integral. When the upper limit is a variable point \( x \), we have an *indefinite* integral. More generally: When the endpoints depend in any way on \( x \), the integral is a function of \( x \). Therefore we can find its derivative. This requires the Fundamental Theorem of Calculus.

The essence of the Theorem is: *Derivative of integral of \( v \) equals \( v \).* We also compute the derivative when the integral goes from \( a(x) \) to \( b(x) \)—both limits variable.

Part 2 of the Fundamental Theorem reverses the order: *Integral of derivative of \( f \) equals \( f + C \).* That will follow quickly from Part 1, with help from the Mean Value Theorem. It is Part 2 that we use most, since integrals are harder than derivatives.

After the proofs we go to new applications, beyond the standard problem of area under a curve. Integrals can add up rings and triangles and shells—not just rectangles. The answer can be a volume or a probability—not just an area.

**THE FUNDAMENTAL THEOREM, PART 1**

Start with a continuous function \( v \). Integrate it from a fixed point \( a \) to a variable point \( x \). For each \( x \), this integral \( f(x) \) is a number. We do not require or expect a formula for \( f(x) \)—it is the area out to the point \( x \). It is a function of \( x! \) The Fundamental Theorem says that this area function has a derivative (another limiting process). *The derivative \( df/dx \) equals the original \( v(x) \).*
5C (Fundamental Theorem, Part 1) Suppose \( v(x) \) is a continuous function:

\[
\text{If } f(x) = \int_a^x v(t) \, dt \text{ then } \frac{df}{dx} = v(x).
\]

The dummy variable is written as \( t \), so we can concentrate on the limits. The value of the integral depends on the limits \( a \) and \( x \), not on \( t \).

To find \( df/dx \), start with \( \Delta f = f(x + \Delta x) - f(x) = \text{difference of areas} \):

\[
\Delta f = \int_x^{x+\Delta x} v(t) \, dt - \int_x^x v(t) \, dt = \int_x^{x+\Delta x} v(t) \, dt.
\]

Officially, this is Property 1. The area out to \( x + \Delta x \) minus the area out to \( x \) equals the small part from \( x \) to \( x + \Delta x \). Now divide by \( \Delta x \):

\[
\frac{\Delta f}{\Delta x} = \frac{1}{\Delta x} \int_x^{x+\Delta x} v(t) \, dt = \text{average value} = v(c).
\]

This is Property 7, the Mean Value Theorem for integrals. The average value on this short interval equals \( v(c) \). This point \( c \) is somewhere between \( x \) and \( x + \Delta x \) (exact position not known), and we let \( \Delta x \) approach zero. That squeezes \( c \) toward \( x \), so \( v(c) \) approaches \( v(x) \)—remember that \( v \) is continuous. The limit of equation (2) is the Fundamental Theorem:

\[
\frac{\Delta f}{\Delta x} \rightarrow \frac{df}{dx} \quad \text{and} \quad v(c) \rightarrow v(x) \quad \text{so} \quad \frac{df}{dx} = v(x).
\]

If \( \Delta x \) is negative the reasoning still holds. Why assume that \( v(x) \) is continuous? Because if \( v \) is a step function, then \( f(x) \) has a corner where \( df/dx \) is not \( v(x) \).

We could skip the Mean Value Theorem and simply bound \( v \) above and below:

for \( t \) between \( x \) and \( x + \Delta x \):

\[
v_{\min} \leq v(t) \leq v_{\max}
\]

integrate over that interval:

\[
v_{\min} \Delta x \leq \Delta f \leq v_{\max} \Delta x
\]

divide by \( \Delta x \):

\[
v_{\min} \leq \frac{\Delta f}{\Delta x} \leq v_{\max}
\]

As \( \Delta x \rightarrow 0 \), \( v_{\min} \) and \( v_{\max} \) approach \( v(x) \). In the limit \( df/dx \) again equals \( v(x) \).

**Graphical meaning** The \( f \)-graph gives the area under the \( v \)-graph. The thin strip in Figure 5.14 has area \( \Delta f \). That area is approximately \( v(x) \) times \( \Delta x \). Dividing by its base, \( \Delta f/\Delta x \) is close to the height \( v(x) \). When \( \Delta x \rightarrow 0 \) and the strip becomes infinitely thin, the expression "close to" converges to "equals." Then \( df/dx \) is the height at \( v(x) \).

**Derivatives with Variable Endpoints**

When the upper limit is \( x \), the derivative is \( v(x) \). Suppose the lower limit is \( x \). The integral goes from \( x \) to \( b \), instead of \( a \) to \( x \). When \( x \) moves, the lower limit moves.
The change in area is on the left side of Figure 5.15. As \( x \) goes forward, area is removed. So there is a minus sign in the derivative of area:

\[
\text{The derivative of } g(x) = \int_x^b v(t) \, dt \quad \text{is} \quad \frac{dg}{dx} = -v(x). \tag{5}
\]

The quickest proof is to reverse \( b \) and \( x \), which reverses the sign (Property 3):

\[
g(x) = -\int_b^x v(t) \, dt \quad \text{so by Part 1} \quad \frac{dg}{dx} = -v(x).
\]

![Fig. 5.15 Area from \( x \) to \( b \) has \( \frac{dg}{dx} = -v(x) \). Area \( v(b) \Delta b \) is added, area \( v(a) \Delta a \) is lost](image)

The general case is messier but not much harder (it is quite useful). Suppose both limits are changing. The upper limit \( b(x) \) is not necessarily \( x \), but it depends on \( x \). The lower limit \( a(x) \) can also depend on \( x \) (Figure 5.15b). The area \( A \) between those limits changes as \( x \) changes, and we want \( \frac{dA}{dx} \):

\[
\text{If } A = \int_{a(x)}^{b(x)} v(t) \, dt \quad \text{then} \quad \frac{dA}{dx} = v(b(x)) \frac{db}{dx} - v(a(x)) \frac{da}{dx}. \tag{6}
\]

The figure shows two thin strips, one added to the area and one subtracted.

First check the two cases we know. When \( a = 0 \) and \( b = x \), we have \( \frac{da}{dx} = 0 \) and \( \frac{db}{dx} = 1 \). The derivative according to (6) is \( v(x) \) times 1 — the Fundamental Theorem. The other case has \( a = x \) and \( b = \) constant. Then the lower limit in (6) produces \( -v(x) \).

When the integral goes from \( a = 2x \) to \( b = x^3 \), its derivative is new:

**EXAMPLE 1**

\[
A = \int_{2x}^{3x} \cos t \, dt = \sin x^3 - \sin 2x
\]

\[dA/dx = (\cos x^3)(3x^2) - (\cos 2x)(2).\]

That fits with (6), because \( db/dx \) is \( 3x^2 \) and \( da/dx \) is 2 (with minus sign). It also looks like the chain rule — which it is! To prove (6) we use the letters \( v \) and \( f \):

\[
A = \int_{a(x)}^{b(x)} v(t) \, dt = f(b(x)) - f(a(x)) \quad \text{(by Part 2 below)}
\]

\[
\frac{dA}{dx} = f'(b(x)) \frac{db}{dx} - f'(a(x)) \frac{da}{dx} \quad \text{(by the chain rule)}
\]

Since \( f' = v \), equation (6) is proved. In the next example the area turns out to be constant, although it seems to depend on \( x \). Note that \( v(t) = 1/t \) so \( v(3x) = 1/3x \).

**EXAMPLE 2**

\[
A = \int_{2x}^{3x} \frac{1}{t} \, dt \quad \text{has} \quad \frac{dA}{dx} = \left(\frac{1}{3x}\right) (3) - \left(\frac{1}{2x}\right) (2) = 0.
\]
Question \[ A = \int_{-x}^{x} v(t) \, dt \] has \[ \frac{dA}{dx} = v(x) + v(-x). \] Why does \[ v(-x) \] have a plus sign?

THE FUNDAMENTAL THEOREM, PART 2

We have used a hundred times the Theorem that is now to be proved. It is the key to integration. "The integral of \[ \frac{df}{dx} \] is \[ f(x) + C.\]" The application starts with \( v(x). \) We search for an \( f(x) \) with this derivative. If \( df/dx = v(x), \) the Theorem says that

\[ \int v(x) \, dx = \int \frac{df}{dx} \, dx = f(x) + C. \]

We can't rely on knowing formulas for \( v \) and \( f \) only the definitions of \( \int \) and \( d/dx. \)

The proof rests on one extremely special case: \( df/dx \) is the \textbf{zero function}. We easily find \( f(x) = \text{constant}. \) The problem is to prove that there are no other possibilities: \( f \) must be constant. When the slope is zero, the graph must be flat. Everybody knows this is true, but intuition is not the same as proof.

Assume that \( df/dx = 0 \) \textbf{in an interval}. If \( f(x) \) is not constant, there are points where \( f(a) \neq f(b). \) By the Mean Value Theorem, there is a point \( c \) where

\[ f'(c) = \frac{f(b) - f(a)}{b - a} \quad \text{(this is not zero because } f(a) \neq f(b)). \]

But \( f'(c) \neq 0 \) directly contradicts \( df/dx = 0. \) Therefore \( f(x) \) must be constant.

Note the crucial role of the Mean Value Theorem. A \textit{local} hypothesis \( (df/dx = 0 \text{ at each point}) \) yields a \textit{global} conclusion \( (f = \text{constant in the whole interval}). \) The derivative narrows the field of view, the integral widens it. The Mean Value Theorem connects instantaneous to average, local to global, points to intervals. This special case (the zero function) applies when \( A(x) \) and \( f(x) \) have the same derivative:

\[ \text{If } df/dx = 0 \text{ on an interval, then } A(x) = f(x) + C. \quad (7) \]

Reason: The derivative of \( A(x) - f(x) \) is zero. So \( A(x) - f(x) \) must be constant.

Now comes the big theorem. It assumes that \( v(x) \) is continuous, and integrates using \( f(x): \)

\begin{center}
\textbf{5D (Fundamental Theorem, Part 2)} If \( v(x) = \frac{df}{dx} \) then \( \int_{a}^{b} v(x) \, dx = f(b) - f(a). \)
\end{center}

Proof\footnote{The antiderivative is \( f(x). \) But Part I gave another antiderivative for the same \( v(x). \) It was the integral — constructed from rectangles and now called \( A(x): \)

\[ A(x) = \int_{a}^{x} v(t) \, dt \quad \text{also has } \frac{dA}{dx} = v(x). \]

Since \( A' = v \) and \( f' = v, \) the special case in equation (7) states that \( A(x) = f(x) + C. \) That is the essential point: \textbf{The integral from rectangles equals } \( f(x) + C. \)

At the lower limit the area integral is \( A = 0. \) So \( f(a) + C = 0. \) At the upper limit \( f(b) + C = A(b). \) Subtract to find \( A(b), \) the definite integral:

\[ A(b) = \int_{a}^{b} v(x) \, dx = f(b) - f(a). \]

Calculus is beautiful its Fundamental Theorem is also its most useful theorem.
Another proof of Part 2 starts with $f' = v$ and looks at subintervals:

\[
\begin{align*}
&f(x_1) - f(a) = v(x_1^*)(x_1 - a) \quad \text{(by the Mean Value Theorem)} \\
&f(x_2) - f(x_1) = v(x_2^*)(x_2 - x_1) \quad \text{(by the Mean Value Theorem)} \\
&\quad \quad \cdots = \cdots \\
&f(b) - f(x_{n-1}) = v(x_n^*)(b - x_{n-1}) \quad \text{(by the Mean Value Theorem)}.
\end{align*}
\]

The left sides add to $f(b) - f(a)$. The sum on the right, as $\Delta x \to 0$, is $\int_a^b v(x) \, dx$.

**APPLICATIONS OF INTEGRATION**

Up to now the integral has been the area under a curve. There are many other applications, quite different from areas. *Whenever addition becomes "continuous," we have integrals instead of sums.* Chapter 8 has space to develop more applications, but four examples can be given immediately—which will make the point.

We stay with geometric problems, rather than launching into physics or engineering or biology or economics. All those will come. The goal here is to take a first step away from rectangles.

**EXAMPLE 3** (for circles) The area $A$ and circumference $C$ are related by $dA/dr = C$.

The question is why. The area is $\pi r^2$. Its derivative $2\pi r$ is the circumference. By the Fundamental Theorem, the integral of $C$ is $A$. What is missing is the geometrical reason. Certainly $\pi r^2$ is the integral of $2\pi r$, but what is the real explanation for $A = \int C(r) \, dr$?

My point is that the pieces are not rectangles. We could squeeze rectangles under a circular curve, but their heights would have nothing to do with $C$. Our intuition has to take a completely different direction, and add up the thin rings in Figure 5.16.

![Figure 5.16](image)

**Fig. 5.16** Area of circle = integral over rings. Volume of sphere = integral over shells.

Suppose the ring thickness is $\Delta r$. Then the ring area is close to $C$ times $\Delta r$. This is precisely the kind of approximation we need, because its error is of higher order $(\Delta r)^2$. The integral adds ring areas just as it added rectangular areas:

\[
A = \int_0^r C \, dr = \int_0^r 2\pi r \, dr = \pi r^2.
\]

That is our first step toward freedom, away from rectangles to rings.
The ring area $\Delta A$ can be checked exactly—it is the difference of circles:

$$\Delta A = \pi (r + \Delta r)^2 - \pi r^2 = 2\pi r \Delta r + \pi (\Delta r)^2.$$  

This is $C \Delta r$ plus a correction. Dividing both sides by $\Delta r \to 0$ leaves $dA/dr = C$.

Finally there is a geometrical reason. The ring unwinds into a thin strip. Its width is $\Delta r$ and its length is close to $C$. The inside and outside circles have different perimeters, so this is not a true rectangle—but the area is near $C \Delta r$.

**EXAMPLE 4** For a sphere, surface area and volume satisfy $A = dV/dr$. What worked for circles will work for spheres. The thin rings become thin shells. A shell goes from radius $r$ to radius $r + \Delta r$, so its thickness is $\Delta r$. We want the volume of the shell, but we don’t need it exactly. The surface area is $4\pi r^2$, so the volume is about $4\pi r^2 \Delta r$. That is close enough!

Again we are correct except for $(\Delta r)^2$. Infinitesimally speaking $dV = A \, dr$:

$$V = \int_0^r A \, dr = \int_0^r 4\pi r^2 \, dr = \frac{4}{3} \pi r^3.$$  

This is the volume of a sphere. The derivative of $V$ is $A$, and the shells explain why. Main point: *Integration is not restricted to rectangles.*

**EXAMPLE 5** The distance around a square is $4s$. Why does the area have $dA/ds = 2s$? The side is $s$ and the area is $s^2$. Its derivative $2s$ goes only half way around the square. I tried to understand that by drawing a figure. Normally this works, but in the figure $dA/ds$ looks like $4s$. Something is wrong. The bell is ringing so I leave this as an exercise.

**EXAMPLE 6** Find the area under $v(x) = \cos^{-1} x$ from $x = 0$ to $x = 1$.

That is a conventional problem, but we have no antiderivative for $\cos^{-1} x$. We could look harder, and find one. However there is another solution—unconventional but correct. *The region can be filled with horizontal rectangles* (not vertical rectangles). Figure 5.17 shows a typical strip of length $x = \cos v$ (the curve has $v = \cos^{-1} x$). As the thickness $\Delta v$ approaches zero, the total area becomes $\int x \, dv$. We are integrating upward, so the limits are on $v$ not on $x$:

$$\text{area} = \int_0^{\pi/2} \cos v \, dv = \sin v \Big|_0^{\pi/2} = 1.$$  

The exercises ask you to set up other integrals—not always with rectangles. Archimedes used triangles instead of rings to find the area of a circle.
5.7 The Fundamental Theorem and Its Applications

5.7 Exercises

Read-through questions

The area \( f(x) = \int_t^x v(t) \, dt \) is a function of \( x \). By Part I of the Fundamental Theorem, its derivative is \( \frac{df}{dx} \). In the proof, a small change \( \Delta x \) produces the area of a thin \( \Delta x \). This area \( \Delta f \) is approximately \( \Delta x \) times \( \frac{df}{dx} \). So the derivative of \( \int_a^x t^2 \, dt \) is \( \frac{x^3}{3} \).

The integral \( \int_0^x t^2 \, dt \) has derivative \( \frac{x^3}{3} \). The minus sign is because \( \frac{df}{dx} \). When both limits \( a(x) \) and \( b(x) \) depend on \( x \), the formula for \( \frac{df}{dx} \) becomes \( \frac{df}{dx} \) minus \( \frac{df}{dx} \). In the example \( \int_2^x t \, dt \), the derivative is \( \frac{x^2}{2} \).

By Part II of the Fundamental Theorem, the integral of \( \frac{df}{dx} \) is \( F(x) \). In the special case when \( \frac{df}{dx} = 0 \), this says that \( \frac{df}{dx} \). From this special case we conclude: If \( \frac{df}{dx} = \frac{db}{dx} \) then \( F(x) = \frac{b}{n} \). If an antiderivative of \( 1/x \) is \( \ln \) (whatever that is), then automatically \( \int_t^x \frac{dx}{x} = \frac{\ln x}{\ln x} \).

The square \( 0 \leq x \leq s \), \( 0 \leq y \leq s \) has area \( A = \frac{p}{q} \). If \( s \) is increased by \( \Delta s \), the extra area has the shape of \( q \). That area \( \Delta A \) is approximately \( \frac{1}{x} \). So \( \frac{dA}{ds} = \frac{1}{x} \).

Find the derivatives of the following functions \( F(x) \).

1. \( \int_0^x \cos t \, dt \)
2. \( \int_0^x \cos 3t \, dt \)
3. \( \int_0^x t^4 \, dt \)
4. \( \int_0^x x^3 \, dt \)
5. \( \int_0^x u^3 \, du \)
6. \( \int_0^x v(u) \, du \)
7. \( \int_0^x \left( \int_0^v t \, dt \right) \) (a "running average" of \( v \))
8. \( \frac{1}{x} \int_0^x v(t) \, dt \) (the average of \( v \); use product rule)
9. \( \frac{1}{x} \int_0^x \sin t \, dt \)
10. \( \frac{1}{x} \int_0^x \frac{\sin t}{t} \, dt \)
11. \( \int_0^x \left[ \int_0^v u \, du \right] \, dt \)
12. \( \int_0^x \left( \int_0^x \frac{df}{dt} \right) \, dt \)
13. \( \int_0^x v(t) \, dt + \int_0^x v(t) \, dt \)
14. \( \int_0^x v(t) \, dt \)
15. \( \int_0^x \sin t \, dt \)
16. \( \int_0^x \sin t \, dt \)
17. \( \int_0^x (u(t) \, dt) \)
18. \( \int_0^x (u(t) \, dt) \)
19. \( \int_0^x \sin^{-1} t \, dt \)
20. \( \int_0^x \frac{d^2f}{dt^2} \, dt \)
21. True or false
   (a) If \( df/dx = dg/dx \) then \( f(x) = g(x) \).
   (b) If \( d^2f/dx^2 = d^2g/dx^2 \) then \( f(x) = g(x) + C \).
   (c) If \( 3 > x \) then the derivative of \( \int_3^x \) \( v(t) \, dt \) is \( -v(x) \).
   (d) The derivative of \( \int_1^x v(x) \, dx \) is zero.
22. For \( F(x) = \int_0^x \sin t \, dt \), locate \( F(\pi + \Delta x) - F(\pi) \) on a sine graph. Where is \( F(\Delta x) - F(0) \)?
23. Find the function \( v(x) \) whose average value between 0 and \( x \) is \( \cos x \). Start from \( \int_0^x v(t) \, dt = x \cos x \).
24. Suppose \( df/dx = 2x \). We know that \( d(x^2)/dx = 2x \). How do we prove that \( f(x) = x^2 + C \)?
25. If \( \int_0^x v(t) \, dt \) then \( \int_0^x v(t) \, dt \) (equal areas left and right of zero), then \( v(x) \) is an \( \frac{df}{dx} \) function. Take derivatives to prove it.
26. Example 2 said that \( \int_0^x t \, dt \) does not really depend on \( x \) (or \( t \)). Substitute \( xu \) for \( t \) and watch the limits on \( u \).
27. True or false, with reason:
   (a) All continuous functions have derivatives.
   (b) All continuous functions have antiderivatives.
   (c) All antiderivatives have derivatives.
   (d) \( A(x) = \int_0^x \frac{dx}{dx^2} \) has \( A(x) = \frac{dx}{dx} \).

Find \( \int_0^x \frac{v(t)}{dt} \) from the facts in 28-29.

28. \( \frac{d(x^2)}{dx} = 2x \)
29. \( \int_0^x \frac{v(t) \, dt}{dt} = \frac{x}{x + 2} \).
30. What is wrong with Figure 5.17? It seems to show that \( A = 4s ds \), which would mean \( A = \int 4s ds = 2s^2 \).
31. The cube \( 0 \leq x, y, z \leq s \) has volume \( V = \frac{4}{3}s^3 \). The three square faces with \( x = s \) or \( y = s \) or \( z = s \) have total area \( A = 6s^2 \). If \( s \) is increased by \( \Delta s \), the extra volume has the shape of \( \Delta s \). That volume \( \Delta V \) is approximately \( \frac{1}{2}s \). So \( \Delta V/\Delta s = \frac{1}{2}s \).
32. The four-dimensional cube \( 0 \leq x, y, z, t \leq s \) has hypervolume \( H = \frac{1}{4}r^4 \). The face with \( x = s \) is really a \( \frac{1}{4}r^4 \). Its volume is \( V = \frac{1}{4}r^4 \). The total volume of the four faces with \( x = s, y = s, z = s, t = s \) is \( \frac{1}{4}r^4 \). When \( s \) is increased by \( \Delta s \), the extra hypervolume is \( \Delta H \approx \frac{1}{4}s \). So \( \Delta H/\Delta s = \frac{1}{4}s \).
33. The hypervolume of a four-dimensional sphere is \( H = \frac{1}{32}r^4 \). Therefore the area (volume?) of its three-dimensional surface \( x^2 + y^2 + z^2 + t^2 = r^2 \) is \( \frac{1}{2}r^2 \).
34. The area above the parabola \( y = x^2 \) from \( x = 0 \) to \( x = 1 \) is \( \frac{1}{2} \). Draw a figure with horizontal strips and integrate.
35. The wedge in Figure (a) has area \( \frac{1}{2}r^2 \). One reason: It is a fraction \( \frac{1}{2}r^2 \). Another reason: It is close to a triangle with small base \( r \delta \) and height \( \delta \). Integrating \( \frac{1}{2}r^2 \delta \) from \( \delta = 0 \) to \( \delta = \pi \) gives the area \( \frac{1}{4}r^2 \) of a quarter-circle.
36. \( A = \int_0^\pi \sqrt{r^2 - x^2} \, dx \) is also the area of a quarter-circle. Show why, with a graph and thin rectangles. Calculate this integral by substituting \( x = r \sin \theta \) and \( dx = r \cos \theta \, d\theta \).
37. The distance $r$ in Figure (b) is related to $\theta$ by $r = \ldots$. Therefore the area of the thin triangle is $\frac{1}{2}r^2d\theta = \ldots$. Integration to $\theta = \ldots$ gives the total area $\frac{1}{2}$.

38. The $x$ and $y$ coordinates in Figure (c) add to $r \cos \theta + r \sin \theta = \ldots$. Without integrating explain why $\int_0^{\pi/2} \frac{d\theta}{(\cos \theta + \sin \theta)^2} = 1$.

39. The horizontal strip at height $y$ in Figure (d) has width $dy$ and length $x = \ldots$. So the area up to $y = 2$ is $\ldots$. What length are the vertical strips that give the same area?

40. Use thin rings to find the area between the circles $r = 2$ and $r = 3$. Draw a picture to show why thin rectangles would be extra difficult.

41. The length of the strip in Figure (e) is approximately $\ldots$. The width is $\ldots$. Therefore the triangle has area $\int_0^1 \ldots \, da$ (do you get $\frac{1}{2}$?)

42. The area of the ellipse in Figure (f) is $2\pi r^2$. Its derivative is $4\pi r$. But this is not the correct perimeter. Where does the usual reasoning go wrong?

43. The derivative of the integral of $v(x)$ is $v(x)$. What is the corresponding statement for sums and differences of the numbers $v_j$? Prove that statement.

44. The integral of the derivative of $f(x)$ is $f(x) + C$. What is the corresponding statement for sums of differences of $f_j$? Prove that statement.

45. Does $\frac{d^2f}{dx^2} = a(x)$ lead to $\int_a^b (\int_0^t a(t) \, dt) \, dx = f(1) - f(0)$?

46. The mountain $y = -x^2 + t$ has an area $A(t)$ above the $x$ axis. As $t$ increases so does the area. Draw an $xy$ graph of the mountain at $t = 1$. What line gives $dA/dt$? Show with words or derivatives that $d^2A/dt^2 > 0$.

### 5.8 Numerical Integration

This section concentrates on definite integrals. The inputs are $y(x)$ and two endpoints $a$ and $b$. The output is the integral $I$. Our goal is to find that number $\int_a^b y(x) \, dx = I$, accurately and in a short time. Normally this goal is achievable—as soon as we have a good method for computing integrals.

Our two approaches so far have weaknesses. The search for an antiderivative succeeds in important cases, and Chapter 7 extends that range—but generally $f(x)$ is not available. The other approach (by rectangles) is in the right direction but too crude. The height is set by $y(x)$ at the right and left end of each small interval. The **right and left rectangle rules** add the areas ($\Delta x$ times $y$):

$$R_n = (\Delta x)(y_1 + y_2 + \ldots + y_n) \quad \text{and} \quad L_n = (\Delta x)(y_0 + y_1 + \ldots + y_{n-1}).$$

The value of $y(x)$ at the end of interval $j$ is $y_j$. The extreme left value $y_0 = y(a)$ enters $L_n$. With $n$ equal intervals of length $\Delta x = (b-a)/n$, the extreme right value is $y_n = y(b)$. It enters $R_n$. Otherwise the sums are the same—simple to compute, easy to visualize, but very inaccurate.

This section goes from slow methods (rectangles) to better methods (trapezoidal and midpoint) to good methods (Simpson and Gauss). Each improvement cuts down the error. You could discover the formulas without the book, by integrating $x$ and
5.8 Numerical Integration

$x^2$ and $x^4$. The rule $R_n$ would come out on one side of the answer, and $L_n$ would be on the other side. You would figure out what to do next, to come closer to the exact integral. The book can emphasize one key point:

The quality of a formula depends on how many integrals

$\int 1 \, dx, \int x \, dx, \int x^2 \, dx, \ldots$, it computes exactly. If $\int x^p \, dx$

is the first to be wrong, the order of accuracy is $p$.

By testing the integrals of $1, x, x^2, \ldots$, we decide how accurate the formulas are.

Figure 5.18 shows the rectangle rules $R_n$ and $L_n$. They are already wrong when $y = x$. These methods are first-order: $p = 1$. The errors involve the first power of $\Delta x$—where we would much prefer a higher power. A larger $p$ in $(\Delta x)^p$ means a smaller error.

When the graph of $y(x)$ is a straight line, the integral $I$ is known. The error triangles $E$ and $e$ have base $\Delta x$. Their heights are the differences $y_{j+1} - y_j$. The areas are $\frac{1}{2}$(base)(height), and the only difference is a minus sign. ($L$ is too low, so the error $L - I$ is negative.) The total error in $R_n$ is the sum of the $E$'s:

$$R_n - I = \frac{1}{2} \Delta x(y_1 - y_0) + \cdots + \frac{1}{2} \Delta x(y_n - y_{n-1}) = \frac{1}{2} \Delta x(y_n - y_0). \quad (1)$$

All $y$'s between $y_0$ and $y_n$ cancel. Similarly for the sum of the $e$'s:

$$L_n - I = -\frac{1}{2} \Delta x(y_n - y_0) = -\frac{1}{2} \Delta x[y(b) - y(a)]. \quad (2)$$

The greater the slope of $y(x)$, the greater the error—since rectangles have zero slope.

Formulas (1) and (2) are nice—but those errors are large. To integrate $y = x$ from $a = 0$ to $b = 1$, the error is $\frac{1}{2} \Delta x(1 - 0)$. It takes 500,000 rectangles to reduce this error to $1/1,000,000$. This accuracy is reasonable, but that many rectangles is unacceptable.

The beauty of the error formulas is that they are “asymptotically correct” for all functions. When the graph is curved, the errors don’t fit exactly into triangles. But the ratio of predicted error to actual error approaches 1. As $\Delta x \to 0$, the graph is almost straight in each interval—this is linear approximation.

The error prediction $\frac{1}{2} \Delta x[y(b) - y(a)]$ is so simple that we test it on $y(x) = \sqrt{x}$:

$$I = \int_0^1 \sqrt{x} \, dx = \frac{2}{3}$$

$$n = 1 \quad 10 \quad 100 \quad 1000$$

error $R_n - I = \quad .33 \quad .044 \quad .0048 \quad .00049$

error $L_n - I = \quad -.67 \quad -.056 \quad -.0052 \quad -.00051$

The error decreases along each row. So does $\Delta x = .1, .01, .001, .0001$. Multiplying $n$ by 10 divides $\Delta x$ by 10. The error is also divided by 10 (almost). The error is nearly proportional to $\Delta x$—typical of first-order methods.

The predicted error is $\frac{1}{2} \Delta x$, since here $y(1) = 1$ and $y(0) = 0$. The computed errors in the table come closer and closer to $\frac{1}{2} \Delta x = .5, .05, .005, .0005$. The prediction is the “leading term” in the actual error.
The table also shows a curious fact. Subtracting the last row from the row above gives exact numbers 1, .1, .01, and .001. This is \((R_n - I) - (L_n - I)\), which is \(R_n - L_n\). It comes from an extra rectangle at the right, included in \(R_n\) but not \(L_n\). Its height is 1 and its area is 1, .1, .01, .001.

**The errors in \(R_n\) and \(L_n\) almost cancel.** The average \(T_n = \frac{1}{2}(R_n + L_n)\) has less error—it is the "trapezoidal rule." First we give the rectangle prediction two final tests:

\[
\begin{array}{c|c|c|c|c}
 n & \text{errors} & n & \text{errors} & n & \text{errors} \\
 1 & 1.7 \cdot 10^{-1} & 10 & 1.7 \cdot 10^{-3} & 100 & 1.7 \cdot 10^{-5} \\
 1000 & 1.7 \cdot 10^{-7} &
\end{array}
\]

Those errors are falling faster than \(\Delta x\). For \(y = x^2 - x\) the prediction explains why: \(y(0)\) equals \(y(1)\). The leading term, with \(y(b)\) minus \(y(a)\), is zero. The exact errors are \(\frac{1}{2}(\Delta x)^2\), dropping from \(10^{-1}\) to \(10^{-3}\) to \(10^{-5}\) to \(10^{-7}\). In these examples \(L_n\) is identical to \(R_n\) (and also to \(T_n\)), because the end rectangles are the same. We will see these \(\frac{1}{2}(\Delta x)^2\) errors in the trapezoidal rule.

The last row in the table is more unusual. It shows practically no error. Why do the rectangle rules suddenly achieve such an outstanding success?

The reason is that \(y(x) = 1/(10 + \cos 2\pi x)\) is periodic. The leading term in the error is zero, because \(y(0) = y(1)\). Also the next term will be zero, because \(y'(0) = y'(1)\). Every power of \(\Delta x\) is multiplied by zero, when we integrate over a complete period. So the errors go to zero exponentially fast.

**Personal note** I tried to integrate \(1/(10 + \cos 2\pi x)\) by hand and failed. Then I was embarrassed to discover the answer in my book on applied mathematics. The method was a special trick using complex numbers, which applies over an exact period. Finally I found the antiderivative (quite complicated) in a handbook of integrals, and verified the area \(1/\sqrt{99}\).

### THE TRAPEZOIDAL AND MIDPOINT RULES

We move to integration formulas that are exact when \(y = x\). They have **second-order accuracy**. The \(\Delta x\) error term disappears. The formulas give the correct area under straight lines. The predicted error is a multiple of \((\Delta x)^2\). That multiple is found by testing \(y = x^2\)—for which the answers are not exact.

The first formula combines \(R_n\) and \(L_n\). To cancel as much error as possible, take the average \(\frac{1}{2}(R_n + L_n)\). This yields the **trapezoidal rule**, which approximates \(\int y(x) dx\) by \(T_n\):

\[
T_n = \frac{1}{2} \sum_{j=0}^{n-1} y_j + \frac{1}{2} y_n = \Delta x \left( y_0 + y_1 + y_2 + \cdots + y_{n-1} + \frac{1}{2} y_n \right).
\]

Another way to find \(T_n\) is from the area of the "trapezoid" below \(y = x\) in Figure 5.19a.

\[
T_n = \Delta x \left[ \frac{1}{2} (y_0 + y_1) + \frac{1}{2} (y_1 + y_2) + \cdots \right]
\]

\[
E = \frac{1}{12} (\Delta x)^2 (y'_{j+1} - y'_j)
\]

\[
e = -\frac{1}{2} E
\]

Fig. 5.19 **Second-order accuracy:** The error prediction is based on \(v = x^2\).
The base is \( \Delta x \) and the sides have heights \( y_{j-1} \) and \( y_j \). Adding those areas gives \( \frac{1}{2}(L_n + R_n) \) in formula (3)—the coefficients of \( y_j \) combine into \( \frac{1}{2} + \frac{1}{2} = 1 \). Only the first and last intervals are missing a neighbor, so the rule has \( \frac{1}{2} y_0 \) and \( \frac{1}{2} y_n \). Because trapezoids (unlike rectangles) fit under a sloping line, \( T_n \) is exact when \( y = x \).

What is the difference from rectangles? The trapezoidal rule gives weight \( \frac{1}{2} \Delta x \) to \( y_0 \) and \( y_n \). The rectangle rule \( R_n \) gives full weight \( \Delta x \) to \( y_n \) (and no weight to \( y_0 \)). \( R_n - T_n \) is exactly the leading error \( \frac{1}{2} y_n - \frac{1}{2} y_0 \). The change to \( T_n \) knocks out that error.

Another important formula is exact for \( y(x) = x \). A rectangle has the same area as a trapezoid, if the height of the rectangle is halfway between \( y_j \) and \( y_{j+1} \). On a straight line graph that is achieved at the midpoint of the interval. By evaluating \( y(x) \) at the halfway points \( \frac{1}{2} \Delta x, \frac{3}{2} \Delta x, \frac{5}{2} \Delta x, \ldots \), we get much better rectangles. This leads to the midpoint rule \( M_n \):

\[
M_n = \Delta x (y_{1/2} + y_{3/2} + \cdots + y_{n-1/2}).
\]  

For \( \int_a^b x \, dx \), trapezoids give \( \frac{1}{2}(0) + 1 + 2 + 3 + \frac{1}{2}(4) = 8 \). The midpoint rule gives \( \frac{1}{2} + \frac{3}{2} + \frac{5}{2} + \frac{7}{2} = 8 \), again correct. The rules become different when \( y = x^2 \), because \( y_{1/2} \) is no longer the average of \( y_0 \) and \( y_1 \). Try both second-order rules on \( x^2 \):

<table>
<thead>
<tr>
<th>( n )</th>
<th>( T_n - I )</th>
<th>( M_n - I )</th>
</tr>
</thead>
<tbody>
<tr>
<td>1</td>
<td>1/6</td>
<td>-1/12</td>
</tr>
<tr>
<td>10</td>
<td>1/600</td>
<td>-1/1200</td>
</tr>
<tr>
<td>100</td>
<td>1/60000</td>
<td>-1/12000</td>
</tr>
</tbody>
</table>

The errors fall by 100 when \( n \) is multiplied by 10. The midpoint rule is twice as good (-1/12 vs. 1/6). Since all smooth functions are close to parabolas (quadratic approximation in each interval), the leading errors come from Figure 5.19. The trapezoidal error is exactly \( \frac{1}{6}(\Delta x)^2 \) when \( y(x) = x^2 \) (the 12 in the formula divides the 2 in \( y' \)):

\[
T_n - I \approx \frac{1}{12} (\Delta x)^2 \left[ (y'_1 - y'_0) + \cdots + (y'_n - y'_{n-1}) \right] = \frac{1}{12} (\Delta x)^2 \left[ y'_n - y'_0 \right]
\]  

\[
M_n - I \approx -\frac{1}{24} (\Delta x)^2 \left[ y'_n - y'_0 \right] = -\frac{1}{24} (\Delta x)^2 \left[ y'(b) - y'(a) \right]
\]

For exact error formulas, change \( y'(b) - y'(a) \) to \( (b - a)y''(c) \). The location of \( c \) is unknown (as in the Mean Value Theorem). In practice these formulas are not much used—they involve the \( p \)th derivative at an unknown location \( c \). The main point about the error is the factor \( (\Delta x)^p \).

One crucial fact is easy to overlook in our tests. Each value of \( y(x) \) can be extremely hard to compute. Every time a formula asks for \( y_j \), a computer calls a subroutine. The goal of numerical integration is to get below the error tolerance, while calling for a minimum number of values of \( y \). Second-order rules need about a thousand values for a typical tolerance of \( 10^{-6} \). The next methods are better.

**FOURTH-ORDER RULE: SIMPSON**

The trapezoidal error is nearly twice the midpoint error (1/6 vs. -1/12). So a good combination will have twice as much of \( M_n \) as \( T_n \). That is Simpson's rule:

\[
S_n = \frac{1}{3} T_n + \frac{2}{3} M_n = \frac{1}{6} \Delta x \left[ y_0 + 4 y_{1/2} + 2 y_1 + 4 y_{3/2} + 2 y_2 + \cdots + 4 y_{n-1/2} + y_n \right].
\]  

Multiply the midpoint values by 2/3 = 4/6. The endpoint values are multiplied by
2/6, except at the far ends \( a \) and \( b \) (with heights \( y_0 \) and \( y_n \)). This 1–4–2–4–2–4–1 pattern has become famous.

Simpson's rule goes deeper than a combination of \( T \) and \( M \). It comes from a parabolic approximation to \( y(x) \) in each interval. When a parabola goes through \( y_0, y_1/2, y_1 \), the area under it is \( \frac{1}{3} \Delta x (y_0 + 4y_{1/2} + y_1) \). This is \( S \) over the first interval. **All our rules are constructed this way: Integrate correctly as many powers \( 1, x, x^2, \ldots \) as possible.** Parabolas are better than straight lines, which are better than flat pieces. \( S \) beats \( M \), which beats \( R \). Check Simpson's rule on powers of \( x \), with \( \Delta x = 1/n \):

\[
\begin{array}{ccc}
n = 1 & n = 10 & n = 100 \\
\text{error if } y = x^2 & 0 & 0 \\
\text{error if } y = x^3 & 0 & 0 \\
\text{error if } y = x^4 & 8.33 \times 10^{-3} & 8.33 \times 10^{-7} & 8.33 \times 10^{-11}
\end{array}
\]

Exact answers for \( x^2 \) are no surprise. \( S \) was selected to get parabolas right. But the zero errors for \( x^3 \) were not expected. The accuracy has jumped to fourth order, with errors proportional to \( (\Delta x)^4 \). That explains the popularity of Simpson's rule.

To understand why \( x^3 \) is integrated exactly, look at the interval \([-1, 1]\). The odd function \( x^3 \) has zero integral, and Simpson agrees by symmetry:

\[
\int_{-1}^{1} x^3 \, dx = \frac{1}{4} x^4 \bigg|_{-1}^{1} = 0 \quad \text{and} \quad \frac{2}{6} \left[ (-1)^3 + 4(0)^3 + 1^3 \right] = 0.
\]

**THE GAUSS RULE (OPTIONAL)**

We need a competitor for Simpson, and Gauss can compete with anybody. He calculated integrals in astronomy, and discovered that two points are enough for a fourth-order method. From \(-1\) to \(1\) (a single interval) his rule is

\[
\int_{-1}^{1} y(x) \, dx \approx y(-1/\sqrt{3}) + y(1/\sqrt{3}).
\]

Those "Gauss points" \( x = -1/\sqrt{3} \) and \( x = 1/\sqrt{3} \) can be found directly. By placing them symmetrically, all odd powers \( x, x^3, \ldots \) are correctly integrated. The key is in \( y = x^2 \), whose integral is \( 2/3 \). The Gauss points \(-x_G\) and \(+x_G\) get this integral right:

\[
\frac{2}{3} = (-x_G)^2 + (x_G)^2, \quad \text{so } x_G^2 = \frac{1}{3} \quad \text{and} \quad x_G = \pm \frac{1}{\sqrt{3}}.
\]

Figure 5.20c shifts to the interval from \( 0 \) to \( \Delta x \). The Gauss points are \((1 \pm 1/\sqrt{3}) \Delta x/2\). They are not as convenient as Simpson's (which hand calculators prefer). Gauss is good for thousands of integrations over one interval. Simpson is
5.8 Numerical Integration

good when intervals go back to back—then Simpson also uses two y's per interval. For \( y = x^4 \), you see both errors drop by \( 10^{-4} \) in comparing \( n = 1 \) to \( n = 10 \):

\[
I = \int_0^1 x^4 \, dx \quad \text{Simpson error} \quad 8.33 \cdot 10^{-3} \quad 8.33 \cdot 10^{-7}
\]

\[
\text{Gauss error} \quad -5.56 \cdot 10^{-3} \quad -5.56 \cdot 10^{-7}
\]

DEFINITE INTEGRALS ON A CALCULATOR

It is fascinating to know how numerical integration is actually done. The points are not equally spaced! For an integral from 0 to 1, Hewlett-Packard machines might internally replace \( x \) by \( 3u^2 - 2u^3 \) (the limits on \( u \) are also 0 and 1). The machine remembers to change \( dx \). For example,

\[
\int_0^1 \frac{dx}{\sqrt{x}} \quad \text{becomes} \quad \int_0^1 \frac{6(1-u)}{\sqrt{3u^2 - 2u^3}} \, du - \int_0^1 \frac{6u}{\sqrt{3 - 2u}} \, du.
\]

Algebraically that looks worse—but the infinite value of \( 1/\sqrt{x} \) at \( x = 0 \) disappears at \( u = 0 \). The differential \( 6(u - u^2) \, du \) was chosen to vanish at \( u = 0 \) and \( u = 1 \). We don't need \( y(x) \) at the endpoints—where infinity is most common. In the \( u \) variable the integration points are equally spaced—therefore in \( x \) they are not.

When a difficult point is inside \([a, b] \), break the interval in two pieces. And chop off integrals that go out to infinity. The integral of \( e^{x^2} \) should be stopped by \( x = 10 \), since the tail is so thin. (It is bad to go too far.) Rapid oscillations are among the toughest—the answer depends on cancellation of highs and lows, and the calculator requires many integration points.

The change from \( x \) to \( u \) affects periodic functions. I thought equal spacing was good, since \( 1/(10 + \cos 2\pi x) \) was integrated above to enormous accuracy. But there is a danger called aliasing. If \( \sin 8\pi x \) is sampled with \( \Delta x = 1/8 \), it is always zero. A high frequency 8 is confused with a low frequency 0 (its "alias" which agrees at the sample points). With unequal spacing the problem disappears. Notice how any integration method can be deceived:

Ask for the integral of \( y = 0 \) and specify the accuracy. The calculator samples \( y \) at \( x_1, \ldots, x_k \). (With a PAUSE key, the \( x \)'s may be displayed.) Then integrate \( Y(x) = (x - x_1)^2 \cdots (x - x_k)^2 \). That also returns the answer zero (now wrong), because the calculator follows the same steps.

On the HP-28S you enter the function, the endpoints, and the accuracy. The variable \( x \) can be named or not (see the margin). The outputs 4.67077 and 4.7E-5 are the requested integral \( \int_0^1 e^x \, dx \) and the estimated error bound. Your input accuracy .00001 guarantees

\[
3: '\text{EXP}(X)' \quad 2: \{X 1 2\} \quad 1: .00001
\]

The machine estimates accuracy based on its experience in sampling \( y(x) \). If you guarantee \( e^x \) within .0000000000, it thinks you want high accuracy and takes longer.

In consulting for HP, William Kahan chose formulas using 1, 3, 7, 15, ... sample points. Each new formula uses the samples in the previous formula. The calculator stops when answers are close. The last paragraphs are based on Kahan's work.
TI-81 Program to Test the Integration Methods \( L, R, T, M, S \)

```plaintext
Prgm1:NUM INT
:Disp "A=
:Input A
:Disp "B=
:Input B
:Lbl N
:Disp "N=
:Input N
:((B-A)/N)+D/2+H
:Input X:YpL:
@+R
8+M
:Lbl I
:X+H+X
M+Yl+R
IS>(J,N)
:Goto 1
:(L+R-Yl)D+L
:RD+R
:MD+M
:(L+R)/2+T
:(2MtT)/3+S
:Disp "L, R, M, T, S"
:Disp L
:Disp R
:Disp M
:Disp T
:Disp S
:Pause
:Goto N
```

Place the integrand \( y(x) \) in the \( Y1 \) position on the \( Y = \) function edit screen. Execute this program, indicating the interval \([A, B]\) and the number of subintervals \( N \). Rules \( L \) and \( R \) and \( M \) use \( N \) evaluations of \( y(x) \). The trapezoidal rule uses \( N+1 \) and Simpson's rule uses \( 2N+1 \). The program pauses to display the results. Press ENTER to continue by choosing a different \( N \). The program never terminates (only pauses). You break out by pressing ON. Don't forget that \( IS, G o t o, ... \) are on menus.

### 5.8 EXERCISES

#### Read-through questions

To integrate \( y(x) \), divide \([a, b]\) into \( n \) pieces of length \( \Delta x = \frac{b-a}{n} \). \( R_n \) and \( L_n \) place a \( b \) over each piece, using the height at the right or c endpoint: \( R_n = \Delta x(y_1 + \cdots + y_n) \) and \( L_n = \Delta x \). These are \( n \) order methods, because they are incorrect for \( y = x \). The total error on \([0,1]\) is approximately \( g \). For \( y = \cos \pi x \) the error is very small because \( [0, 1] \) is a complete interval.

A much better method is \( T_n = \frac{1}{2}R_n + \frac{1}{2}L_n = \Delta x[y_0 + \frac{1}{2}y_1 + \cdots + \frac{1}{2}y_n]. \) This \( n \)-order rule is \( n \)-order because the error for \( y = x \) is \( \frac{1}{2} \). The error for \( y = x^2 \) from \( a \) to \( b \) is \( \frac{1}{2} \). The \( \frac{1}{2} \) rule is twice as accurate, using \( M_n = \Delta x[\frac{1}{2}]. \)

Simpson's method is \( S_n = \frac{2}{3}M_n + \frac{1}{3} \). It is \( \frac{1}{3} \)-order, because the powers \( \frac{1}{3} \) are integrated correctly. The coefficients of \( y_0, y_1/2, y_1 \) are \( \frac{1}{3} \) times \( \Delta x \). Over three intervals the weights are \( \Delta x/6 \) times \( 1-4-\ldots-1 \). Gauss uses \( \frac{1}{3} \) points in each interval, separated by \( \Delta x/\sqrt{3} \). For a method of order \( p \) the error is nearly proportional to \( \frac{1}{2} \).

1. What is the difference \( L_n - T_n \)? Compare with the leading error term in (2).

2. If you cut \( \Delta x \) in half, by what factor is the trapezoidal error reduced (approximately)? By what factor is the error in Simpson's rule reduced?

3. Compute \( R_n \) and \( L_n \) for \( \int_0^1 x^3 \, dx \) and \( n = 1, 2, 10 \). Either verify (with computer) or use (without computer) the formula \( 1^3 + 2^3 + \cdots + n^3 = \frac{1}{4}n^2(n + 1)^2 \).

4. One way to compute \( T_n \) is by averaging \( \frac{1}{2}(L_n + R_n) \). Another way is to add \( \frac{1}{6}y_0 + \frac{1}{3}y_1 + \cdots + \frac{1}{6}y_n \). Which is more efficient? Compare the number of operations.

5. Test three different rules on \( I = \int_0^1 x^4 \, dx \) for \( n = 3, 6, 12, \ldots \).

6. Compute \( \pi \) to six places as \( 4 \int_0^1 \frac{x}{1 + x^2} \, dx \), using any rule.

7. Change Simpson's rule to \( \Delta x(y_0 + \frac{1}{3}y_1 + \frac{1}{3}y_2) \) in each interval and find the order of accuracy \( p \).

8. Demonstrate superdecay of the error when \( 1/(3 + \sin x) \) is integrated from \( 0 \) to \( 2\pi \).

9. Check that \( (\Delta x)^3(y_{j+1} - y_j)/12 \) is the correct error for \( y = 1 \) and \( y = x \) and \( y = x^2 \) from the first trapezoid \( (j = 0) \). Then it is correct for every parabola over every interval.

10. Repeat Problem 9 for the midpoint error \(-(\Delta x)^2(y_{j+1} - y_j)/24 \). Draw a figure to show why the rectangle \( M \) has the same area as any trapezoid through the midpoint (including the trapezoid tangent to \( y(x) \)).

11. In principle \( \int_0^\infty \frac{\sin^2 x}{x} \, dx = \pi \) With a symbolic algebra code or an HP-28S, how many decimal places do you get? Cut off the integral to \( \int_0^a \) and test large and small \( a \).

12. These four integrals all equal \( \pi \):

\[
\int_0^1 \frac{dx}{\sqrt{x(1-x)}} = \int_0^\infty \frac{\sin x}{x} \, dx = 8 \int_0^\pi \sin^4 x \, dx = 3 \int_0^\pi \sin^1 x \, dx = \int_0^\infty \frac{x^{-1/2} \, dx}{1 + x}
\]

(a) Apply the midpoint rule to two of them until \( \pi \approx 3.1416 \).

(b) Optional: Pick the other two and find \( \pi \approx 3 \).
13 To compute \( \int_0^1 \frac{dx}{x} = .69315 \) with error less than .001, how many intervals should \( T_n \) need? Its leading error is \((\Delta x)^2 [y'(b) - y'(a)]/12\). Test the actual error with \( y = 1/x \).

14 Compare \( T_n \) with \( M_n \) for \( \int_0^1 \sqrt{x} \, dx \) and \( n = 1, 10, 100 \). The error prediction breaks down because \( y'(0) = \infty \).

15 Take \( f(x) = \int_0^x y(x) \, dx \) in error formula 3R to prove that \( \int_0^x y(x) \, dx - y(x) \Delta x \) is exactly \((\Delta x)^2 y'(c)\) for some point \( c \).

16 For the periodic function \( y(x) = 1/(2 + \cos 6\pi x) \) from \(-1\) to \(1\), compare \( T \) and \( S \) and \( G \) for \( n = 2 \).

17 For \( I = \int_0^1 \sqrt{1 - x^2} \, dx \), the leading error in the trapezoidal rule is \( \frac{1}{12} \). Try \( n = 2, 4, 8 \) to defy the prediction.

18 Change to \( x = \sin \theta, \sqrt{1 - x^2} = \cos \theta, dx = \cos \theta \, d\theta \), and repeat \( T_n \) on \( \int_0^{\pi/2} \cos^2 \theta \, d\theta \). What is the predicted error after the change to \( \theta \)?

19 Write down the three equations \( Ay(0) + By(\frac{1}{2}) + Cy(1) = I \) for the three integrals \( I = \int_0^1 1 \, dx, \int_0^1 x \, dx, \int_0^1 x^2 \, dx \). Solve for \( A, B, C \) and name the rule.

20 Can you invent a rule using \( Ay(0) + By_{1/4} + Cy_{1/2} + Dy_{3/4} + Ey_1 \) to reach higher accuracy than Simpson's?

21 Show that \( T_n \) is the only combination of \( L_n \) and \( R_n \) that has second-order accuracy.

22 Calculate \( \int e^{-x^2} \, dx \) with ten intervals from 0 to 5 and 0 to 20 and 0 to 400. The integral from 0 to \( \infty \) is \( \frac{\sqrt{\pi}}{2} \). What is the best point to chop off the infinite integral?

23 The graph of \( y(x) = 1/(x^2 + 10^{-10}) \) has a sharp spike and a long tail. Estimate \( \int_0^1 y \, dx \) from \( T_{100} \) and \( T_{200} \) (don't expect much). Then substitute \( x = 10^{-5} \tan \theta, dx = 10^{-5} \sec^2 \theta \, d\theta \) and integrate \( 10^5 \) from 0 to \( \pi/4 \).

24 Compute \( \int_0^1 |x - \pi| \, dx \) from \( T_n \) and compare with the divide and conquer method of separating \( \int_0^1 |x - \pi| \, dx \) from \( \int_0^1 |x - \pi| \, dx \).

25 Find \( a, b, c \) so that \( y = ax^2 + bx + c \) equals 1, 3, 7 at \( x = 0, 1, 2 \) (three equations). Check that \( \frac{1}{7} + \frac{3}{5} + \frac{7}{3} \) equals \( \int_0^3 y \, dx \).

26 Find \( c \) in \( S - I = c(\Delta x)^2 [y''(1) - y''(0)] \) by taking \( y = x^4 \) and \( \Delta x = 1 \).

27 Find \( c \) in \( G - I = c(\Delta x)^3 [y'''(1) - y'''(-1)] \) by taking \( y = x^4, \Delta x = 2, \) and \( G = (-1/\sqrt{3})^4 + (1/\sqrt{3})^4 \).

28 What condition on \( y(x) \) makes \( L_n = R_n = T_n \) for the integral \( \int_a^b y(x) \, dx \)?

29 Suppose \( y(x) \) is concave up. Show from a picture that the trapezoidal answer is too high and the midpoint answer is too low. How does \( y'' > 0 \) make equation (5) positive and (6) negative?
<table>
<thead>
<tr>
<th>CHAPTER 4</th>
<th>The Chain Rule</th>
</tr>
</thead>
<tbody>
<tr>
<td>4.1 Derivatives by the Chain Rule</td>
<td>154</td>
</tr>
<tr>
<td>4.2 Implicit Differentiation and Related Rates</td>
<td>160</td>
</tr>
<tr>
<td>4.3 Inverse Functions and Their Derivatives</td>
<td>164</td>
</tr>
<tr>
<td>4.4 Inverses of Trigonometric Functions</td>
<td>171</td>
</tr>
</tbody>
</table>

<table>
<thead>
<tr>
<th>CHAPTER 5</th>
<th>Integrals</th>
</tr>
</thead>
<tbody>
<tr>
<td>5.1 The Idea of the Integral</td>
<td>177</td>
</tr>
<tr>
<td>5.2 Antiderivatives</td>
<td>182</td>
</tr>
<tr>
<td>5.3 Summation vs. Integration</td>
<td>187</td>
</tr>
<tr>
<td>5.4 Indefinite Integrals and Substitutions</td>
<td>195</td>
</tr>
<tr>
<td>5.5 The Definite Integral</td>
<td>201</td>
</tr>
<tr>
<td>5.6 Properties of the Integral and the Average Value</td>
<td>206</td>
</tr>
<tr>
<td>5.7 The Fundamental Theorem and Its Consequences</td>
<td>213</td>
</tr>
<tr>
<td>5.8 Numerical Integration</td>
<td>220</td>
</tr>
</tbody>
</table>

<table>
<thead>
<tr>
<th>CHAPTER 6</th>
<th>Exponentials and Logarithms</th>
</tr>
</thead>
<tbody>
<tr>
<td>6.1 An Overview</td>
<td>228</td>
</tr>
<tr>
<td>6.2 The Exponential $e^x$</td>
<td>236</td>
</tr>
<tr>
<td>6.3 Growth and Decay in Science and Economics</td>
<td>242</td>
</tr>
<tr>
<td>6.4 Logarithms</td>
<td>252</td>
</tr>
<tr>
<td>6.5 Separable Equations Including the Logistic Equation</td>
<td>259</td>
</tr>
<tr>
<td>6.6 Powers Instead of Exponentials</td>
<td>267</td>
</tr>
<tr>
<td>6.7 Hyperbolic Functions</td>
<td>277</td>
</tr>
</tbody>
</table>

<table>
<thead>
<tr>
<th>CHAPTER 7</th>
<th>Techniques of Integration</th>
</tr>
</thead>
<tbody>
<tr>
<td>7.1 Integration by Parts</td>
<td>283</td>
</tr>
<tr>
<td>7.2 Trigonometric Integrals</td>
<td>288</td>
</tr>
<tr>
<td>7.3 Trigonometric Substitutions</td>
<td>294</td>
</tr>
<tr>
<td>7.4 Partial Fractions</td>
<td>300</td>
</tr>
<tr>
<td>7.5 Improper Integrals</td>
<td>305</td>
</tr>
</tbody>
</table>

<table>
<thead>
<tr>
<th>CHAPTER 8</th>
<th>Applications of the Integral</th>
</tr>
</thead>
<tbody>
<tr>
<td>8.1 Areas and Volumes by Slices</td>
<td>311</td>
</tr>
<tr>
<td>8.2 Length of a Plane Curve</td>
<td>320</td>
</tr>
<tr>
<td>8.3 Area of a Surface of Revolution</td>
<td>325</td>
</tr>
<tr>
<td>8.4 Probability and Calculus</td>
<td>328</td>
</tr>
<tr>
<td>8.5 Masses and Moments</td>
<td>336</td>
</tr>
<tr>
<td>8.6 Force, Work, and Energy</td>
<td>342</td>
</tr>
</tbody>
</table>
CHAPTER 6

Exponentials and Logarithms

This chapter is devoted to exponentials like $2^x$ and $10^x$ and above all $e^x$. The goal is to understand them, differentiate them, integrate them, solve equations with them, and invert them (to reach the logarithm). The overwhelming importance of $e^x$ makes this a crucial chapter in pure and applied mathematics.

In the traditional order of calculus books, $e^x$ waits until other applications of the integral are complete. I would like to explain why it is placed earlier here. I believe that the equation $dy/dx = y$ has to be emphasized above techniques of integration. The laws of nature are expressed by differential equations, and at the center is $e^x$. Its applications are to life sciences and physical sciences and economics and engineering (and more—wherever change is influenced by the present state). The model produces a differential equation and I want to show what calculus can do.

The key is always $b^{m+n} = (b^m)(b^n)$. Section 6.1 applies that rule in three ways:

1. to understand the logarithm as the exponent;
2. to draw graphs on ordinary and semilog and log-log paper;
3. to find derivatives. The slope of $b^x$ will use $b^{x+\Delta x} = (b^x)(b^\Delta x)$.

6.1 An Overview

There is a good chance you have met logarithms. They turn multiplication into addition, which is a lot simpler. They are the basis for slide rules (not so important) and for graphs on log paper (very important). Logarithms are mirror images of exponentials—and those I know you have met.

Start with exponentials. The numbers 10 and $10^2$ and $10^3$ are basic to the decimal system. For completeness I also include $10^0$, which is “ten to the zeroth power” or 1. The logarithms of those numbers are the exponents. The logarithms of 1 and 10 and 100 and 1000 are 0 and 1 and 2 and 3. These are logarithms “to base 10,” because the powers are powers of 10.

Question When the base changes from 10 to $b$, what is the logarithm of 1?
Answer Since $b^0 = 1$, $\log_b 1$ is always zero. To base $b$, the logarithm of $b^n$ is $n$. 
Negative powers are also needed. The number $10^x$ is positive, but its exponent $x$ can be negative. The first examples are $1/10$ and $1/100$, which are the same as $10^{-1}$ and $10^{-2}$. The logarithms are the exponents $-1$ and $-2$:

\[
\begin{align*}
1000 &= 10^3 & \text{and} & \log 1000 &= 3 \\
1/1000 &= 10^{-3} & \text{and} & \log 1/1000 &= -3.
\end{align*}
\]

Multiplying $1000$ times $1/1000$ gives $1 = 10^0$. Adding logarithms gives $3 + (-3) = 0$. Always $10^m$ times $10^n$ equals $10^{m+n}$. In particular $10^3$ times $10^2$ produces five tens:

\[
\]

The law for $b^m$ times $b^n$ extends to all exponents, as in $10^{4.6}$ times $10^5$. Furthermore the law applies to all bases (we restrict the base to $b > 0$ and $b \neq 1$). In every case multiplication of numbers is addition of exponents.

6A \[ b^m \text{ times } b^n \text{ equals } b^{m+n}, \text{ so logarithms (exponents) add } b^m \text{ divided by } b^n \text{ equals } b^{m-n}, \text{ so logarithms (exponents) subtract } \]

\[
\log_b(yz) = \log_b y + \log_b z \quad \text{and} \quad \log_b(y/z) = \log_b y - \log_b z. \tag{1}
\]

**Historical note**  In the days of slide rules, $1.2$ and $1.3$ were multiplied by sliding one edge across to $1.2$ and reading the answer under $1.3$. A slide rule made in Germany would give the third digit in $1.56$. Its photograph shows the numbers on a log scale. The distance from 1 to 2 equals the distance from 2 to 4 and from 4 to 8. By sliding the edges, you add distances and multiply numbers.

Division goes the other way. Notice how $1000/10 = 100$ matches $3 - 1 = 2$. To divide $1.56$ by $1.3$, look back along line D for the answer $1.2$.

The second figure, though smaller, is the important one. **When $x$ increases by 1, $2^x$ is multiplied by 2. Adding to $x$ multiplies $y$.** This rule easily gives $y = 1, 2, 4, 8$, but look ahead to calculus—which doesn’t stay with whole numbers.

Calculus will add $\Delta x$. Then $y$ is multiplied by $2^{\Delta x}$. This number is near 1. If $\Delta x = \frac{1}{10}$ then $2^{\Delta x} \approx 1.07$ — the tenth root of 2. **To find the slope, we have to consider** $(2^{\Delta x} - 1)/\Delta x$. The limit is near $(1.07 - 1)/\frac{1}{10} = .7$, but the exact number will take time.

![](image) **Fig. 6.1** An ancient relic (the slide rule). When exponents $x$ add, powers $2^x$ multiply.
6 Exponentials and Logarithms

**Base Change**  Bases other than 10 and exponents other than 1, 2, 3, ... are needed for applications. The population of the world x years from now is predicted to grow by a factor close to 1.02^x. Certainly x does not need to be a whole number of years. And certainly the base 1.02 should not be 10 (or we are in real trouble). This prediction will be refined as we study the differential equations for growth. It can be rewritten to base 10 if that is preferred (but look at the exponent):

\[ 1.02^x \text{ is the same as } 10^{\log_{1.02} x}. \]

When the base changes from 1.02 to 10, the exponent is multiplied—as we now see. For practice, start with base \( b \) and change to base \( a \). The logarithm to base \( a \) will be written "\( \log_a \)." Everything comes from the rule that logarithm = exponent:

**base change for numbers:**  \( b = a^{\log_b a} \).

Now raise both sides to the power \( x \). You see the change in the exponent:

**base change for exponentials:**  \( b^x = a^{(\log_b a)x} \).

Finally set \( y = b^x \). Its logarithm to base \( b \) is \( x \). Its logarithm to base \( a \) is the exponent on the right hand side: \( \log_a y = (\log_a b)x \). Now replace \( x \) by \( \log_b y \):

**base change for logarithms:**  \( \log_a y = (\log_b a)(\log_b y) \).

We absolutely need this ability to change the base. An example with \( a = 2 \) is

\[ b = 8 = 2^3 \quad 8^2 = (2^3)^2 = 2^6 \quad \log_2 64 = 3 \cdot 2 = (\log_2 8)(\log_2 64). \]

**The rule behind base changes is**  \( (a^m)^r = a^{mr} \). When the \( m \)th power is raised to the \( x \)th power, the exponents multiply. The square of the cube is the sixth power:

\( (a)(a)(a) \text{ times } (a)(a)(a) \text{ equals } (a)(a)(a)(a)(a)(a): \quad (a^3)^2 = a^6 \).

Another base will soon be more important than 10—here are the rules for base changes:

6B  **To change numbers, powers, and logarithms from base \( b \) to base \( a \), use**

\[
\begin{align*}
 b &= a^{\log_b a} & b^x &= a^{(\log_b a)x} & \log_a y &= (\log_b a)(\log_b y)
\end{align*}
\]

The first is the definition. The second is the \( x \)th power of the first. The third is the logarithm of the second (remember \( y = b^x \)). An important case is \( y = a \):

\[ \log_a a = (\log_b a)(\log_b a) = 1 \text{ so } \log_a b = 1/\log_a a. \]

**EXAMPLE**  \( 8 = 2^3 \) means \( 8^{1/3} = 2 \). Then \( (\log_2 8)(\log_2 2) = 3(1/3) = 1 \).

This completes the algebra of logarithms. The addition rules 6A came from \( (b^m)(b^n) = b^{m+n} \). The multiplication rule 6B came from \( (a^m)^r = a^{mr} \). We still need to define \( b^x \) and \( a^x \) for all real numbers \( x \). When \( x \) is a fraction, the definition is easy. The square root of \( a^b \) is \( a^{b/2} \). When \( x \) is not a fraction, as in \( 2^x \), the graph suggests one way to fill in the hole. We could define \( 2^x \) as the limit of \( 2^{1/2}, 2^{1/10}, 2^{1/100}, ..., \) As the fractions \( r \) approach \( \pi \), the powers \( 2^r \) approach \( 2^x \). This makes \( y = 2^x \) into a continuous function, with the desired properties \( (2^m)(2^n) = 2^{m+n} \) and \( (2^m)^r = 2^{mr} \)—whether \( m \) and \( n \) and \( x \) are integers or not. But the \( e \)'s and \( \delta \)'s of continuity are not attractive, and we eventually choose (in Section 6.4) a smoother approach based on integrals.
GRAPHS OF $b^x$ AND $\log_b y$

It is time to draw graphs. In principle one graph should do the job for both functions, because $y = b^x$ means the same as $x = \log_b y$. These are inverse functions. What one function does, its inverse undoes. The logarithm of $g(x) = b^x$ is $x$:

$$g^{-1}(g(x)) = \log_b(b^x) = x.$$  

In the opposite direction, the exponential of the logarithm of $y$ is $y$:

$$g(g^{-1}(y)) = b^{(\log_b y)} = y.$$  

This holds for every base $b$, and it is valuable to see $b = 2$ and $b = 4$ on the same graph. Figure 6.2a shows $y = 2^x$ and $y = 4^x$. Their mirror images in the 45° line give the logarithms to base 2 and base 4, which are in the right graph.

When $x$ is negative, $y = b^x$ is still positive. If the first graph is extended to the left, it stays above the x axis. Sketch it in with your pencil. Also extend the second graph down, to be the mirror image. Don’t cross the vertical axis.

There are interesting relations within the left figure. All exponentials start at 1, because $b^0$ is always 1. At the height $y = 16$, one graph is above $x = 2$ (because $4^2 = 16$). The other graph is above $x = 4$ (because $2^4 = 16$). Why does $4^x$ in one graph equal $2^{2x}$ in the other? This is the base change for powers, since $4 = 2^2$.

The figure on the right shows the mirror image—the logarithm. All logarithms start from zero at $y = 1$. The graphs go down to $-\infty$ at $y = 0$. (Roughly speaking $2^{-\infty}$ is zero.) Again $x$ in one graph corresponds to $2x$ in the other (base change for logarithms). Both logarithms climb slowly, since the exponentials climb so fast.

The number $\log_2 10$ is between 3 and 4, because 10 is between $2^3$ and $2^4$. The slope of $2^x$ is proportional to $2^x$—which never happened for $x^a$. But there are two practical difficulties with those graphs:

1. $2^x$ and $4^x$ increase too fast. The curves turn virtually straight up.
2. The most important fact about $Ah^x$ is the value of $h$—and the base doesn’t stand out in the graph.

There is also another point. In many problems we don’t know the function $y = f(x)$. We are looking for it! All we have are measured values of $y$ (with errors mixed in). When the values are plotted on a graph, we want to discover $f(x)$.

Fortunately there is a solution. Scale the $y$ axis differently. On ordinary graphs, each unit upward adds a fixed amount to $y$. On a log scale each unit multiplies $y$ by
a fixed amount. The step from \( y = 1 \) to \( y = 2 \) is the same length as the step from 3 to 6 or 10 to 20.

On a log scale, \( y = 10 \) is not halfway between 10 and 12. And \( y = 0 \) is not there at all. Each step down divides by a fixed amount—we never reach zero. This is completely satisfactory for \( Ab^x \), which also never reaches zero.

Figure 6.3 is on semilog paper (also known as log-linear), with an ordinary \( x \) axis. **The graph of \( y = Ab^x \) is a straight line.** To see why, take logarithms of that equation:

\[
\log y = \log A + x \log b.
\]

**(6)**

The relation between \( x \) and \( \log y \) is linear. It is really \( \log y \) that is plotted, so the graph is straight. The markings on the \( y \) axis allow you to enter \( y \) without looking up its logarithm—you get an ordinary graph of \( \log y \) against \( x \).

Figure 6.3 shows two examples. One graph is an exact plot of \( y = 2 \cdot 10^x \). It goes upward with slope 1, because a unit across has the same length as multiplication by 10 going up. \( 10^x \) has slope 1 and \( 10^{\log b} \) (which is \( b^x \)) will have slope \( \log b \). The crucial number \( \log b \) can be measured directly as the slope.

![Fig. 6.3 2 \cdot 10^x \text{ and } 4 \cdot 10^{-x/2} \text{ on semilog paper.}](image)

**Fig. 6.4** Graphs of \( Ax^b \) on log-log paper.

The second graph in Figure 6.3 is more typical of actual practice, in which we start with measurements and look for \( f(x) \). Here are the data points:

<table>
<thead>
<tr>
<th>( x )</th>
<th>0.0</th>
<th>0.2</th>
<th>0.4</th>
<th>0.6</th>
<th>0.8</th>
<th>1.0</th>
</tr>
</thead>
<tbody>
<tr>
<td>( y )</td>
<td>4.0</td>
<td>3.2</td>
<td>2.4</td>
<td>2.0</td>
<td>1.6</td>
<td>1.3</td>
</tr>
</tbody>
</table>

We don't know in advance whether these values fit the model \( y = Ab^x \). The graph is strong evidence that they do. The points lie close to a line with negative slope—indicating \( \log b < 0 \) and \( b < 1 \). The slope down is half of the earlier slope up, so the
model is consistent with
\[
y = A \cdot 10^{-x/2} \quad \text{or} \quad \log y = \log A - \frac{1}{2}x.
\] (7)

When \(x\) reaches 2, \(y\) drops by a factor of 10. At \(x = 0\) we see \(A \approx 4\).

Another model—a power \(y = Ax^k\) instead of an exponential—also stands out with logarithmic scaling. This time we use log–log paper, with both axes scaled. The logarithm of \(y = Ax^k\) gives a linear relation between \(\log y\) and \(\log x\):
\[
\log y = \log A + k \log x.
\] (8)

The exponent \(k\) becomes the slope on log–log paper. The base \(b\) makes no difference. We just measure the slope, and a straight line is a lot more attractive than a power curve.

The graphs in Figure 6.4 have slopes 3 and \(\frac{1}{2}\) and \(-1\). They represent \(Ax^3\) and \(A\sqrt{x}\) and \(A/x\). To find the \(A\)'s, look at one point on the line. At \(x = 4\) the height is 8, so adjust the \(A\)'s to make this happen: The functions are \(x^3/8\) and \(4\sqrt{x}\) and \(32/x\). On semilog paper those graphs would not be straight!

You can buy log paper or create it with computer graphics.

**The Derivatives of \(y = b^x\) and \(x = \log_b y\)**

This is a calculus book. We have to ask about slopes. The algebra of exponents is done, the rules are set, and on log paper the graphs are straight. Now come limits.

The central question is the derivative. What is \(dy/dx\) when \(y = b^x\)? What is \(dx/dy\) when \(x\) is the logarithm \(\log_b y\)? Those questions are closely related, because \(b^x\) and \(\log_b y\) are inverse functions. If one slope can be found, the other is known from \(dx/dy = 1/(dy/dx)\). The problem is to find one of them, and the exponential comes first.

You will now see that those questions have quick (and beautiful) answers, except for a mysterious constant. There is a multiplying factor \(c\) which needs more time. I think it is worth separating out the part that can be done immediately, leaving \(c\) in \(dy/dx\) and \(1/c\) in \(dx/dy\). Then Section 6.2 discovers \(c\) by studying the special number called \(e\) (but \(c \neq e\)).

### 6C The derivative of \(b^x\) is a multiple \(cb^x\). The number \(c\) depends on the base \(b\).

The product and power and chain rules do not yield this derivative. We are pushed all the way back to the original definition, the limit of \(\Delta y/\Delta x\):
\[
\frac{dy}{dx} = \lim_{h \to 0} \frac{y(x + h) - y(x)}{h} = \lim_{h \to 0} \frac{b^{x+h} - b^x}{h}.
\] (9)

**Key idea:** Split \(b^{x+h}\) into \(b^x\) times \(b^h\). Then the crucial quantity \(b^x\) factors out. More than that, \(b^x\) comes outside the limit because it does not depend on \(h\). The remaining limit, inside the brackets, is the number \(c\) that we don’t yet know:
\[
\frac{dy}{dx} = \lim_{h \to 0} \frac{b^x b^h - b^x}{h} = b^x \left[ \lim_{h \to 0} \frac{b^h - 1}{h} \right] = cb^x.
\] (10)

This equation is central to the whole chapter: \(dy/dx equals cb^x which equals cy\). The rate of change of \(y\) is proportional to \(y\). The slope increases in the same way that \(b^x\) increases (except for the factor \(c\)). A typical example is money in a bank, where
6 Exponentials and Logarithms

interest is proportional to the principal. The rich get richer, and the poor get slightly richer. We will come back to compound interest, and identify \( b \) and \( c \).

The inverse function is \( x = \log_by \). Now the unknown factor is \( 1/c \):

### 6D The slope of \( \log_by \) is \( 1/cy \) with the same \( c \) (depending on \( b \)).

**Proof** If \( dy/dx = cb^x \) then \( dx/dy = 1/cb^x = 1/cy \).

That proof was like a Russian toast, powerful but too quick! We go more carefully:

\[
\begin{align*}
  \frac{d}{dx} f(b^x) &= x \quad \text{(logarithm of exponential)} \\
  f'(b^x)(cb^x) &= 1 \quad \text{(x derivative by chain rule)} \\
  f'(b^x) &= 1/cb^x \quad \text{(divide by \( cb^x \))} \\
  f'(y) &= 1/cy \quad \text{(identify \( b^x \) as \( y \))}
\end{align*}
\]

The logarithm gives another way to find \( c \). From its slope we can discover \( 1/c \). This is the way that finally works (next section).

**Fig. 6.5** The slope of \( 2^x \) is about \(.7 \cdot 2^x\). The slope of \( \log_2 y \) is about \( 1/7y \).

**Final remark** It is extremely satisfying to meet an \( f(y) \) whose derivative is \( 1/cy \). At last the "\(-1\) power" has an antiderivative. Remember that \( \int x^n dx = x^{n+1}/(n+1) \) is a failure when \( n = -1 \). The derivative of \( x^0 \) (a constant) does not produce \( x^{-1} \).

**We had no integral for \( x^{-1} \), and the logarithm fills that gap.** If \( y \) is replaced by \( x \) or \( t \) (all dummy variables) then

\[
\begin{align*}
  \frac{d}{dx} \log_b x &= \frac{1}{cx} \\
  \frac{d}{dt} \log_b t &= \frac{1}{ct}
\end{align*}
\]

The base \( b \) can be chosen so that \( c = 1 \). Then the derivative is \( 1/x \). This final touch comes from the magic choice \( b = e \)—the highlight of Section 6.2.

### 6.1 EXERCISES

**Read-through questions**

In \( 10^4 = 10,000 \), the exponent 4 is the ___ of 10,000. The base is \( b = \) ___ . The logarithm of \( 10^4 \) times \( 10^8 \) is ___ .

The logarithm of \( 10^9/10^6 \) is ___ . The logarithm of 10,000 is ___ . If \( y = b^x \) then \( x = \) ___ . Here \( x \) is any number, and \( y \) is always ___ .

A base change gives \( b = a^{h} \) and \( b^x = a^{x} \). Then \( 8^3 \) is \( 2^{15} \). In other words \( \log_2 y \) is ___ times \( \log_8 y \). When \( y = 2 \) it follows that \( \log_2 8 \) times \( \log_8 2 \) equals ___ .

On ordinary paper the graph of \( y = \) ___ is a straight line. Its slope is ___ . On semilog paper the graph of \( y = \) ___ \( n \) is a straight line. Its slope is ___ . On log-log paper the graph of \( y = \) ___ \( n \) is a straight line. Its slope is ___ .

The slope of \( y = b^x \) is \( dy/dx = \) ___ , where \( c \) depends on \( b \). The number \( c \) is the limit as \( h \to 0 \) of ___ . Since \( x = \log_b y \) is the inverse, \( dx/dy)(dy/dx) = \) ___ . Knowing \( dy/dx = cb^x \) yields \( dx/dy = \) ___ . Substituting \( b^x \) for \( y \), the slope of \( \log_b y \) is ___ . With a change of letters, the slope of \( \log_b x \) is ___ .
6.1 An Overview

Problems 1–10 use the rules for logarithms.

1. Find these logarithms (or exponents):
   (a) \( \log_2 32 \)  
   (b) \( \log_3 (1/3) \)  
   (c) \( \log_2 (1/32) \)  
   (d) \( \log_3 2 \)  
   (e) \( \log_10 (\sqrt{10}) \)  
   (f) \( \log_2 (\log_2 16) \)

2. Without a calculator find the values of
   (a) \( 3^{\log_3 5} \)  
   (b) \( 3^{2\log_3 5} \)  
   (c) \( \log_{10} 5 + \log_{10} 2 \)  
   (d) \( \log_{10} (\log_{10} 9) \)  
   (e) \( 10^{3\times 10^{-4}} \times 10^{3} \)  
   (f) \( \log_2 56 - \log_2 7 \)

3. Sketch \( y = 2^{-x} \) and \( y = \frac{1}{2}(4^x) \) from -1 to 1 on the same graph. Put their mirror images \( x = -\log_2 y \) and \( x = \log_4 2y \) on a second graph.

4. Following Figure 6.2 sketch the graphs of \( y = (\frac{1}{2})^x \) and \( x = \log_{1/2} y \). What are \( \log_{1/2} 2 \) and \( \log_{1/2} 4 \)?

5. Compute without a computer:
   (a) \( \log_2 3 + \log_2 3 \)  
   (b) \( \log_2 (\frac{1}{2})^{10} \)  
   (c) \( \log_{10} 100^{10} \)  
   (d) \( \log_{10} (\log_{10} 9) \)  
   (e) \( 2^{3/2} / (2^3) \)  
   (f) \( \log_2 (1/\sqrt{2}) \)

6. Solve the following equations for \( x \):
   (a) \( \log_{10} (10^{2x}) = 7 \)  
   (b) \( \log_4 4x - \log_4 4 = \log_4 3 \)  
   (c) \( \log_{10} 10 = 2 \)  
   (d) \( \log_2 (1/x) = 2 \)  
   (e) \( \log x + \log x = \log 8 \)  
   (f) \( \log_2 (x^2) = 5 \)

7. The logarithm of \( y = x^x \) is \( \log_{10} y = \) _________.

*8. Prove that \( (\log_{10} a)(\log_{10} c) = (\log_{10} c)(\log_{10} a) \).

9. \( 2^{10} \) is close to \( 10^3 \) (1024 versus 1000). If they were equal then \( \log_{10} 2 \) would be _______. Also \( \log_{10} 2 \) would be _______ instead of 0.301.

10. The number \( 2^{10^{100}} \) has approximately how many (decimal) digits?

Questions 11–19 are about the graphs of \( y = b^x \) and \( x = \log_{b} y \).

11. By hand draw the axes for semilog paper and the graphs of \( y = 1.1^x \) and \( y = 10(1.1)^x \).

12. Display a set of axes on which the graph of \( y = \log_{10} x \) is a straight line. What other equations give straight lines on those axes?

13. When noise is measured in decibels, amplifying by a factor \( A \) increases the decibel level by \( 10 \log A \). If a whisper is 20dB and a shout is 70dB then \( 10 \log A = 50 \) and \( A = \) ________.

14. Draw semilog graphs of \( y = 10^{1-x} \) and \( y = \frac{1}{2}(\sqrt{10})^x \).

15. The Richter scale measures earthquakes by \( \log_{10} (I/I_0) = R \). What is \( R \) for the standard earthquake of intensity \( I_0 \)? If the 1989 San Francisco earthquake measured \( R = 7 \), how did its intensity \( I \) compare to \( I_0 \)? The 1906 San Francisco quake had \( R = 8.3 \). The record quake was four times as intense with \( R = \) ________.

16. The frequency of \( A \) above middle \( C \) is \( 440/\text{second} \). The frequency of the next higher \( A \) is _______. Since \( 2^{1/12} \approx 1.5 \), the note with frequency \( 660/\text{sec} \) is ________.

17. Draw your own semilog paper and plot the data 
   \[ y = 7, 11, 16, 28, 44 \] at \( x = 0, 1/2, 1, 3/2, 2 \).

Estimate \( A \) and \( b \) in \( y = Ab^x \).

18. Sketch log–log graphs of \( y = x^2 \) and \( y = \sqrt{x} \).

19. On log–log paper, printed or homemade, plot \( y = 4, 11, 21, 32, 45 \) at \( x = 1, 2, 3, 4, 5 \). Estimate \( A \) and \( k \) in \( y = Ax^k \).

Questions 20–29 are about the derivative \( \frac{dy}{dx} = cb^x \).

20. \( g(x) = b^x \) has slope \( g'(x) = cb^x \). Apply the chain rule to \( g(f(y)) = y \) to prove that \( \frac{df}{dy} = 1/ey \).

21. If the slope of \( \log x \) is \( 1/\text{cx} \), find the slopes of \( \log(2x) \) and \( \log(x^2) \) and \( \log(2^x) \).

22. What is the equation (including \( c \)) for the tangent line to \( y = 10^x \) at \( x = 0 \)? Find also the equation at \( x = 1 \).

23. What is the equation for the tangent line to \( x = \log_{10} y \) at \( y = 1? \) Find also the equation at \( y = 10 \).

24. With \( b = 10 \), the slope of \( 10^x \) is \( c10^x \). Use a calculator for small \( h \) to estimate \( c = \lim (10^x - 1)/x \).

25. The unknown constant in the slope of \( y = (.1)^y \) is \( L = \lim (1^{y} - 1)/y \). (a) Estimate \( L \) by choosing a small \( h \).  
(b) Change \( h \) to \( -h \) to show that \( L = -c \) from Problem 24.

26. Find a base \( b \) for which \( (b^h - 1)/h \) is close to \( 1 \). Use \( h = 1/4 \) by hand or \( h = 1/10 \) and \( 1/100 \) by calculator.

27. Find the second derivative of \( y = b^x \) and also of \( x = \log_{b} y \).

28. Show that \( C = \lim (100^y - 1)/y \) is twice as large as \( c = \lim (10^x - 1)/x \). (Replace the last \( h \) 's by \( 2h \).)

29. In 28, the limit for \( b = 100 \) is twice as large as for \( b = 10 \). So \( c \) probably involves the ______ of \( b \).
The last section discussed $b^x$ and $\log_b y$. The base $b$ was arbitrary—it could be 2 or 6 or 9.3 or any positive number except 1. But in practice, only a few bases are used. I have never met a logarithm to base 6 or 9.3. Realistically there are two leading candidates for $b$, and 10 is one of them. This section is about the other one, which is an extremely remarkable number. This number is not seen in arithmetic or algebra or geometry, where it looks totally clumsy and out of place. In calculus it comes into its own.

The number is $e$. That symbol was chosen by Euler (initially in a fit of selfishness, but he was a wonderful mathematician). It is the base of the natural logarithm. It also controls the exponential $e^x$, which is much more important than $\ln x$. Euler also chose $\pi$ to stand for perimeter—anyway, our first goal is to find $e$.

Remember that the derivatives of $b^x$ and $\log_b y$ include a constant $c$ that depends on $b$. Equations (10) and (11) in the previous section were

$$\frac{d}{dx} b^x = cb^x \quad \text{and} \quad \frac{d}{dy} \log_b y = \frac{1}{cy}.$$ (1)

At $x = 0$, the graph of $b^x$ starts from $b^0 = 1$. The slope is $c$. At $y = 1$, the graph of $\log_b y$ starts from $\log_b 1 = 0$. The logarithm has slope $1/c$. With the right choice of the base $b$ those slopes will equal 1 (because $c$ will equal 1).

For $y = 2^x$ the slope $c$ is near .7. We already tried $\Delta x = .1$ and found $\Delta y \approx .07$. The base has to be larger than 2, for a starting slope of $c = 1$.

We begin with a direct computation of the slope of $\log_b y$ at $y = 1$:

$$\frac{1}{c} = \text{slope at } 1 = \lim_{h \to 0} \frac{1}{h} [\log_b (1 + h) - \log_b 1] = \lim_{h \to 0} \log_b [(1 + h)^{1/h}].$$ (2)

Always $\log_b 1 = 0$. The fraction in the middle is $\log_b (1 + h)$ times the number $1/h$. This number can go up into the exponent, and it did.

The quantity $(1 + h)^{1/h}$ is unusual, to put it mildly. As $h \to 0$, the number $1 + h$ is approaching 1. At the same time, $1/h$ is approaching infinity. In the limit we have $1^\infty$. But that expression is meaningless (like $0/0$). Everything depends on the balance between "nearly 1" and "nearly $\infty$." This balance produces the extraordinary number $e$:

**DEFINITION** The number $e$ is equal to $\lim_{h \to 0} (1 + h)^{1/h}$. Equivalently $e = \lim_{n \to \infty} \left(1 + \frac{1}{n}\right)^n$.

Before computing $e$, look again at the slope $1/c$. At the end of equation (2) is the logarithm of $e$:

$$1/c = \log_b e.$$ (3)

When the base is $b = e$, the slope is $\log_e e = 1$. That base $e$ has $c = 1$ as desired:

The derivative of $e^x$ is $1 \cdot e^x$ and the derivative of $\log_e y$ is $\frac{1}{1 \cdot y}$. (4)

This is why the base $e$ is all-important in calculus. It makes $c = 1$.

To compute the actual number $e$ from $(1 + h)^{1/h}$, choose $h = 1, 1/10, 1/100, \ldots$. Then the exponents $1/h$ are $n = 1, 10, 100, \ldots$ (All limits and derivatives will become official in Section 6.4.) The table shows $(1 + h)^{1/h}$ approaching $e$ as $h \to 0$ and $n \to \infty$.
The last column is converging to \( e \) (not quickly). There is an infinite series that converges much faster. We know 125,000 digits of \( e \) (and a billion digits of \( \pi \)). There are no definite patterns, although you might think so from the first sixteen digits:

\[
e = 2.71828 1828 45 90 43 43 26 
\]

The powers of \( e \) produce \( y = e^x \). At \( x = 2.3 \) and \( 5 \), we are close to \( y = 10 \) and \( 150 \).

The logarithm is the inverse function. The logarithms of 150 and 10, to the base \( e \), are close to \( x = 5 \) and \( x = 2.3 \). There is a special name for this logarithm—the natural logarithm. There is also a special notation “\( \ln \)” to show that the base is \( e \):

\[
\ln y \text{ means the same as } \log_e y. \text{ The natural logarithm is the exponent in } e^x = y. 
\]

The notation \( \ln y \) (or \( \ln x \)—it is the function that matters, not the variable) is standard in calculus courses. After calculus, the base is generally assumed to be \( e \). In most of science and engineering, the natural logarithm is the automatic choice. The symbol “\( \exp (x) \)” means \( e^x \), and the truth is that the symbol “\( \log x \)” generally means \( \ln x \). Base \( e \) is understood even without the letters \( \ln \). But in any case of doubt—on a calculator key for example—the symbol “\( \ln x \)” emphasizes that the base is \( e \).

**THE DERIVATIVES OF \( e^x \) AND \( \ln x \)**

Come back to derivatives and slopes. The derivative of \( b^x \) is \( cb^x \), and the derivative of \( \log_b y \) is \( \frac{1}{cy} \). If \( b = e \) then \( c = 1 \). For all bases, equation (3) is \( \frac{1}{c} = \log_b e \). This gives \( c \)—the slope of \( b^x \) at \( x = 0 \):

\[
c = \ln b \text{ is the mysterious constant that was not available earlier. The slope of } 2^x \text{ is } \ln 2 \text{ times } 2^x. \text{ The slope of } e^x \text{ is } \ln e \text{ times } e^x \text{ (but } \ln e = 1 \text{). We have the derivatives on which this chapter depends:}
\]

\[
6E \text{ The number } c \text{ is } \frac{1}{\log_b e} = \log_e b. \text{ Thus } c \text{ equals } \ln b. \quad (5)
\]

\[
e = \ln b \text{ is the mysterious constant that was not available earlier. The slope of } 2^x \text{ is } \ln 2 \text{ times } 2^x. \text{ The slope of } e^x \text{ is } \ln e \text{ times } e^x \text{ (but } \ln e = 1 \text{). We have the derivatives on which this chapter depends:}
\]

\[
6F \text{ The derivatives of } e^x \text{ and } \ln y \text{ are } e^x \text{ and } 1/y. \text{ For other bases}
\]

\[
\frac{d}{dx} b^x = (\ln b) b^x \quad \text{and} \quad \frac{d}{dy} \log_b y = \frac{1}{(\ln b)y}. \quad (6)
\]

To make clear that those derivatives come from the functions (and not at all from the dummy variables), we rewrite them using \( t \) and \( x \):

\[
\frac{d}{dt} e^t = e^t \quad \text{and} \quad \frac{d}{dx} \ln x = \frac{1}{x}. \quad (7)
\]
Remark on slopes at $x = 0$: It would be satisfying to see directly that the slope of $2^x$ is below 1, and the slope of $4^x$ is above 1. Quick proof: $e$ is between 2 and 4. But the idea is to see the slopes graphically. This is a small puzzle, which is fun to solve but can be skipped.

$2^x$ rises from 1 at $x = 0$ to 2 at $x = 1$. On that interval its average slope is 1. Its slope at the beginning is smaller than average, so it must be less than 1—as desired. On the other hand $4^x$ rises from $\frac{1}{2}$ at $x = -\frac{1}{2}$ to 1 at $x = 0$. Again the average slope is $\frac{1}{2}$, which comes at the end of this new interval, the slope of $4^x$ at that point exceeds 1. Somewhere between $2^x$ and $4^x$ is $e^x$, which starts out with slope 1.

This is the graphical approach to $e$. There is also the infinite series, and a fifth definition through integrals which is written here for the record:

1. $e$ is the number such that $e^x$ has slope 1 at $x = 0$
2. $e$ is the base for which $\ln y = \log_e y$ has slope 1 at $y = 1$
3. $e$ is the limit of $(1 + \frac{1}{n})^n$ as $n \to \infty$
4. $e = \frac{1}{0!} + \frac{1}{1!} + \frac{1}{2!} + \frac{1}{3!} + \cdots = 1 + \frac{1}{2} + \frac{1}{6} + \cdots$
5. the area $\int_1^e x^{-1} \, dx$ equals 1.

The connections between 1, 2, and 3 have been made. The slopes are 1 when $e$ is the limit of $(1 + 1/n)^n$. Multiplying this out will lead to 4, the infinite series in Section 6.6. The official definition of $\ln x$ comes from $\int dx/x$, and then 5 says that $\ln e = 1$. This approach to $e$ (Section 6.4) seems less intuitive than the others.

Figure 6.6b shows the graph of $e^{-x}$. It is the mirror image of $e^x$ across the vertical axis. Their product is $e^x e^{-x} = 1$. Where $e^x$ grows exponentially, $e^{-x}$ decays exponentially—or it grows as $x$ approaches $-\infty$. Their growth and decay are faster than any power of $x$. Exponential growth is more rapid than polynomial growth, so that $e^x/x^n$ goes to infinity (Problem 59). It is the fact that $e^x$ has slope $e^x$ which keeps the function climbing so fast.

The other curve is $y = e^{-x^{2/3}}$. This is the famous “bell-shaped curve” of probability theory. After dividing by $\sqrt{2\pi}$, it gives the normal distribution, which applies to so many averages and so many experiments. The Gallup Poll will be an example in Section 8.4. The curve is symmetric around its mean value $x = 0$, since changing $x$ to $-x$ has no effect on $x^2$.

About two thirds of the area under this curve is between $x = -1$ and $x = 1$. If you pick points at random below the graph, 2/3 of all samples are expected in that interval. The points $x = -2$ and $x = 2$ are “two standard deviations” from the center,
enclosing 95% of the area. There is only a 5% chance of landing beyond. The decay is even faster than an ordinary exponential, because $\frac{1}{2}x^2$ has replaced $x$.

**THE DERIVATIVES OF $e^{ux}$ AND $e^{u(x)}$**

The slope of $e^x$ is $e^x$. This opens up a whole world of functions that calculus can deal with. The chain rule gives the slope of $e^{3x}$ and $e^{\sin x}$ and every $e^{u(x)}$.

6G The derivative of $e^{u(x)}$ is $e^{u(x)}$ times $du/dx$. (8)

Special case $u = cx$: The derivative of $e^{cx}$ is $ce^{cx}$. (9)

**EXAMPLE 1** The derivative of $e^{3x}$ is $3e^{3x}$ (here $c = 3$). The derivative of $e^{\sin x}$ is $e^{\sin x} \cos x$ (here $u = \sin x$). The derivative of $f(u(x))$ is $df/du$ times $du/dx$. Here $f = e^u$ so $df/du = e^u$. The chain rule demands that second factor $du/dx$.

**EXAMPLE 2** $e^{\ln(2)x}$ is the same as $2^x$. Its derivative is $\ln 2$ times $2^x$. The chain rule redisovers our constant $c = \ln 2$. In the slope of $b^x$ it redisovers the factor $c = \ln b$.

Generally $e^{cx}$ is preferred to the original $b^x$. The derivative just brings down the constant $c$. It is better to agree on $e$ as the base, and put all complications (like $c = \ln b$) up in the exponent. The second derivative of $e^{cx}$ is $c^2e^{cx}$.

**EXAMPLE 3** The derivative of $e^{-x^2/2}$ is $-xe^{-x^2/2}$ (here $u = -x^2/2$ so $du/dx = -x$).

**EXAMPLE 4** The second derivative of $f = e^{-x^2/2}$, by the chain rule and product rule, is

$$f'' = (-1)\cdot e^{-x^2/2} + (-x)^2e^{-x^2/2} = (x^2 - 1)e^{-x^2/2}.$$ (10)

Notice how the exponential survives. With every derivative it is multiplied by more factors, but it is still there to dominate growth or decay. The points of inflection, where the bell-shaped curve has $f'' = 0$ in equation (10), are $x = 1$ and $x = -1$.

**EXAMPLE 5** $(u = n \ln x)$. Since $e^{n\ln x}$ is $x^n$ in disguise, its slope must be $nx^{n-1}$:

$$\text{slope} = e^{n\ln x} \frac{d}{dx} (n \ln x) = x^n \frac{n}{x} = nx^{n-1}.$$ (11)

This slope is correct for all $n$, integer or not. Chapter 2 produced $3x^2$ and $4x^3$ from the binomial theorem. Now $nx^{n-1}$ comes from $\ln$ and exp and the chain rule.

**EXAMPLE 6** An extreme case is $x^x = (e^{\ln x})^x$. Here $u = x \ln x$ and we need $du/dx$:

$$\frac{d}{dx} (x^x) = e^{x \ln x} \left( \ln x + x \cdot \frac{1}{x} \right) = x^x (\ln x + 1).$$

**INTEGRALS OF $e^{ux}$ AND $e^u$ $du/dx$**

The integral of $e^x$ is $e^x$. The integral of $e^{ux}$ is not $e^{ux}$. The derivative multiplies by $c$ so the integral divides by $c$. The integral of $e^{cx}$ is $e^{cx}/c$ (plus a constant).

**EXAMPLES**

$$\int e^{2x}dx = \frac{1}{2}e^{2x} + C \quad \int b^x dx = \frac{b^x}{\ln b} + C$$
Exponentials and Logarithms

The first one has \( c = 2 \). The second has \( c = \ln b \)—remember again that \( b^x = e^{(\ln b)x} \). The integral divides by \( \ln b \). In the third one, \( e^{3(x + 1)} \) is \( e^{3x} \) times the number \( e^3 \) and that number is carried along. Or more likely we see \( e^{3(x + 1)} \) as \( e^3 \). The missing \( du/dx \) is fixed by dividing by 3. The last example fails because \( du/dx \) is not there. **We cannot integrate without \( du/dx \):**

### 6H The indefinite integral \( \int e^x \frac{du}{dx} \, dx = e^{ax} + C \).

Here are three examples with \( du/dx \) and one without it:

\[
\int e^{\sin x} \cos x \, dx = e^{\sin x} + C \quad \int xe^{x/2} \, dx = e^{x/2} + C \\
\int e^{\sqrt{x}} \, dx = 2e^{\sqrt{x}} + C \quad \int \frac{e^x \, dx}{(1 + e^x)^2} = \frac{-1}{1 + e^x} + C
\]

The first is a pure \( e^x \). So is the second. The third has \( u = \sqrt{x} \) and \( du/dx = 1/2\sqrt{x} \), so only the factor 2 had to be fixed. The fourth example does not belong with the others. It is the integral of \( du/u^2 \), not the integral of \( e^x \). I don't know any way to tell you which substitution is best—except that the complicated part is \( 1 + e^x \) and it is natural to substitute \( u \). If it works, good.

Without an extra \( e^x \) for \( du/dx \), the integral \( \int dx/(1 + e^x) \) looks bad. But \( u = 1 + e^x \) is still worth trying. It has \( du = e^x \, dx = (u - 1) \, dx \):

\[
\int \frac{dx}{(1 + e^x)^2} = \int \frac{du}{(u - 1)u^2} = \int \frac{1}{u-1} \left( \frac{1}{u} - \frac{1}{u^2} \right) \, du.
\]

That last step is "partial fractions." The integral splits into simpler pieces (explained in Section 7.4) and we integrate each piece. Here are three other integrals:

\[
\int e^{1/x} \, dx \quad \int e^{x(4 + e^x)} \, dx \quad \int e^{-x(4 + e^x)} \, dx
\]

The first can change to \(- \int e^x du/u^2 \), which is not much better. (It is just as impossible.)

The second is actually \( \int u \, du \), but I prefer a split: \( \int 4e^x \) and \( \int e^{2x} \) are safer to do separately. The third is \( \int (4e^{-x} + 1) \, dx \), which also separates. The exercises offer practice in reaching \( e^x \, dx \) — ready to be integrated.

**Warning about definite integrals**  When the lower limit is \( x = 0 \), there is a natural tendency to expect \( f(0) = 0 \)—in which case the lower limit contributes nothing. For a power \( f = x^3 \) that is true. For an exponential \( f = e^{3x} \) it is definitely not true, because \( f(0) = 1 \):

\[
\int_0^1 e^{3x} \, dx = \frac{1}{3} e^3 \bigg|_0^1 = \frac{1}{3} (e^3 - 1) \quad \int_0^1 xe^{2x} \, dx = \frac{1}{2} e^2 \bigg|_0^1 = \frac{1}{2} (e - 1).
\]
6.2 The Exponential e^x

**Read-through questions**

The number e is approximately \( \approx 2.71828 \). It is the limit of \((1 + h)^{1/h}\) to the power \( \lim_{h \to 0} \). This gives \( e = \lim (1 + h)^{1/h} \). An equivalent form is \( e = \lim (1 + h)^{1/h} \).

When the base is \( b = e \), the constant \( c \) in Section 6.1 is \( c = \ln b \). Therefore the derivative of \( y = e^x \) is \( dy/dx = e^x \). The derivative of \( x = \log_y y \) is \( dx/dy = -\ln y \). The slopes at \( y = 0 \) and \( y = 1 \) are both \( \ln y \). The notation for \( \log_n y \) is \( \log_b y \) for some constant \( b \).

The constant \( c \) in the slope of \( b^x \) is \( c = \ln b \). The function \( b^x \) can be rewritten as \( b^{x+1} / b^x \). Its derivative is \( b^x \ln b \). The derivative of \( e^{bx} \) is \( xe^x \). The derivative of \( e^{ax} \) is \( e^{ax} \).

The integral of \( e^x \) is \( e^x \). The integral of \( e^{-x} \) is \( -e^{-x} \). The integral of \( e^{ax} \) by itself is \( \frac{1}{a} e^{ax} \) to find.

**Find the derivatives of the functions in 1–18.**

1 \( 7e^{2x} \)
2 \(-7e^{-7x} \)
3 \((e^x)^3 \)
4 \((e^{-x})^9 \)
5 \(3x \)
6 \(e^x \ln x \)
7 \((2/3)x^5 \)
8 \(4xe^x \)
9 \(1/(1 + e^x) \)
10 \(e^{1/(1 + x)} \)
11 \(e^{2x} + xe^x \)
12 \(xe^{1/2} \)
13 \(xe^x - e^x \)
14 \(x^2 e^x - 2xe^x + 2e^x \)
15 \(\frac{e^x - e^{-x}}{e^x + e^{-x}} \)
16 \(e^{\cos x} + \ln(e^x) \)
17 \(e^{\sin x} + \sin e^x \)
18 \(x^{-1/3} \) (which is \( 1/x^3 \))

19 The difference between \( e \) and \((1 + 1/n)^n\) is approximately \( Ce/n \). Subtract the calculated values for \( n = 10, 100, 1000 \) from 2.7183 to discover the number \( C \).

20 By algebra or a calculator find the limits of \((1 + 1/n)^{2n}\) and \((1 + 1/n)^{3n}\).

21 The limit of \((11/10)^{10}, (101/100)^{100}, \ldots \) is \( e \). So the limit of \((10/11)^{10}, (100/110)^{100}, \ldots \) is \( e \). So the limit of \((10/11)^{10}, (100/101)^{10}, \ldots \) is \( e \). The last sequence is \((1 - 1/n)^{n} \).

22 Compare the number of correct decimals of \( e \) for \((1.001)^{10000} \) and \((1.00001)^{100000} \) and if possible \((1.000001)^{100000} \). Which power \( n \) would give all the decimals in 2.71828?

23 The function \( y = e^x \) solves \( dy/dx = y \). Approximate this equation by \( \Delta Y/X \), which is \( Y(x + h) - Y(x) = hY(x) \). With \( h = \tau \) find \( Y(h) \) after one step starting from \( Y(0) = 1 \). What is \( Y(1) \) after ten steps?

24 The function that solves \( dy/dx = -y \) starting from \( y = 1 \) at \( x = 0 \) is \( \approx \). Approximate by \( Y(x + h) = Y(x) = -hY(x) \). If \( h = \frac{1}{10} \) what is \( Y(h) \) after one step and what is \( Y(1) \) after four steps?

25 Invent three functions \( f, g, h \) such that for \( x > 10 \)

\( (1 + 1/x)^x < f(x) < e^x < g(x) < e^{1/x} < h(x) < x^x \).

26 Graph \( e^x \) and \( \sqrt{e^x} \) at \( x = -2, -1, 0, 1, 2 \). Another form of \( \sqrt{e^x} \) is \( \approx \).

**Find antiderivatives for the functions in 27–36.**

27 \( e^{3x} + e^{7x} \)
28 \((e^{2x})(e^{7x}) \)
29 \(1^x + 2^x + 3^x \)
30 \(2^{-x} \)
31 \((2e)^x + 2e^x \)
32 \((1/e)^x + (1/x^7) \)
33 \(xe^x + xe^{-x^2} \)
34 \((\sin x)e^{\cos x} + (\cos x)e^{\sin x} \)
35 \(\sqrt{e^x + (e^x)^2} \)
36 \(xe^x \) (trial and error)

37 Compare \( e^{-x} \) with \( e^{-x^2} \). Which one decreases faster near \( x = 0 \)? Where do the graphs meet again? When is the ratio of \( e^{-x^2} \) to \( e^{-x} \) less than \( 1/100 \)?

38 Compare \( e^x \) with \( x^x \). Where do the graphs meet? What are their slopes at that point? Divide \( x^x \) by \( e^x \) and show that the ratio approaches infinity.

39 Find the minima of \( y = x^x \) for \( x > 0 \). Show from \( d^2y/dx^2 \) that the curve is concave upward.

40 By comparing slopes, prove that if \( x > 0 \) then

(a) \( e^x > 1 + x \)
(b) \( e^{-x} > 1 - x \).

41 Find the minimum value of \( y = x^x \) for \( x > 0 \). Show from \( d^2y/dx^2 \) that the curve is concave upward.

42 Find the slope of \( y = x^{1/x} \) and the point where \( dy/dx = 0 \). Check \( d^2y/dx^2 \) to show that the maximum of \( x^{1/x} \) is \( \approx \).

43 If \( dy/dx = y \) find the derivative of \( e^{-y}y \) by the product rule. Deduce that \( y(x) = Ce^x \) for some constant \( C \).

44 Prove that \( x^x = e^x \) has only one positive solution.

**Evaluate the integrals in 45–54.** With infinite limits, 49–50 are "improper."

45 \( \int_0^1 e^{2x} \) dx
46 \( \int_0^\infty \sin x e^{\cos x} \) dx
47 \( \int_{-1}^1 2^x \) dx
48 \( \int_{-1}^1 2^{-x} \) dx
49 \( \int_0^\infty e^{-x} \) dx
50 \( \int_0^\infty xe^{-x^2} \) dx
51 \( \int_0^1 e^{1+x} \) dx
52 \( \int_0^1 e^{1+x} \) x dx
53 $\int_0^x 2^{\text{int}} \cos x \, dx$
54 $\int_0^1 (1 - e^x)^{10} e^x \, dx$

55 Integrate the integrals that can be integrated:

$\int \frac{e^x}{du/\, \, dx} \, dx$ 
$\int \left( \frac{du}{e^x} \right)^2 \, dx$ 
$\int (e^x)^2 \frac{du}{dx} \, dx$.

56 Find a function that solves $y'(x) = 5y(x)$ with $y(0) = 2$.

57 Find a function that solves $y'(x) = 1/y(x)$ with $y(0) = 2$.

58 With electronic help graph the function $(1 + 1/x)^x$. What are its asymptotes? Why?

59 This exercise shows that $F(x) = x^4/e^x \to 0$ as $x \to \infty$.
(a) Find $dF/\, \, dx$. Notice that $F(x)$ decreases for $x > n > 0$.
The maximum of $x^n/e^n$, at $x = n$, is $n^n/e^n$.
(b) $F(2x) = (2x)^n/e^{2x} = 2^n x^n/e^n \cdot e^x \leq 2^n n^n/e^n \cdot e^x$.

Deduce that $F(2x) \to 0$ as $x \to \infty$. Thus $F(x) \to 0$.

60 With $n = 6$, graph $F(x) = x^6/e^x$ on a calculator or computer. Estimate its maximum. Estimate $x$ when you reach $F(x) = 1$. Estimate $x$ when you reach $F(x) = \frac{1}{2}$.

61 Stirling’s formula says that $n! \approx \sqrt{2\pi n} n^n/e^n$. Use it to estimate $6^6/e^6$ to the nearest whole number. Is it correct? How many decimal digits in $10!$?

62 $x^6/e^x \to 0$ is also proved by l’Hôpital’s rule (at $x = \infty$):

$$\lim_{x \to \infty} x^6/e^x = \lim_{x \to \infty} 6x^5/e^x = \text{fill this in} = 0.$$ 

6.3 Growth and Decay in Science and Economics

The derivative of $y = e^x$ has taken time and effort. The result was $y' = ce^x$, which means that $y' = cy$. That computation brought others with it, virtually for free—the derivatives of $bx$ and $x^a$ and $e^{a(x)}$. But I want to stay with $y' = cy$—which is the most important differential equation in applied mathematics.

Compare $y' = x$ with $y' = y$. The first only asks for an antiderivative of $x$. We quickly find $y = \frac{1}{2} x^2 + C$. The second has $dy/\, \, dx$ equal to $y$ itself—which we rewrite as $dy/y = dx$. The integral is $\ln y = x + C$. Then $y$ itself is $e^x e^x$. Notice that the first solution is $\frac{1}{2} x^2$ plus a constant, and the second solution is $e^x$ times a constant.

There is a way to graph slope $x$ versus slope $y$. Figure 6.7 shows “tangent arrows,” which give the slope at each $x$ and $y$. For parabolas, the arrows grow steeper as $x$

![Fig. 6.7](image-url)  

The slopes are $y' = x$ and $y' = y$. The solution curves fit those slopes.
grows—because \( y' = \text{slope} = x \). For exponentials, the arrows grow steeper as \( y \) grows—the equation is \( y' = \text{slope} = y \). Now the arrows are connected by \( y = Ae^x \).

A differential equation gives a field of arrows (slopes). Its solution is a curve that stays tangent to the arrows—then the curve has the right slope.

A field of arrows can show many solutions at once (this comes in a differential equations course). Usually a single \( y_0 \) is not sacred. To understand the equation we start from many \( y_0 \)—on the left the parabolas stay parallel, on the right the heights stay proportional. For \( y' = -y \) all solution curves go to zero.

From \( y' = y \) it is a short step to \( y' = cy \). To make \( c \) appear in the derivative, put \( c \) into the exponent. The derivative of \( y = e^{cx} \) is \( ce^{cx} \), which is \( c \) times \( y \). We have reached the key equation, which comes with an initial condition—a starting value \( y_0 \):

\[
\frac{dy}{dt} = cy \quad \text{with} \quad y = y_0 \quad \text{at} \quad t = 0. \tag{1}
\]

A small change: \( x \) has switched to \( t \). In most applications time is the natural variable, rather than space. The factor \( c \) becomes the “growth rate” or “decay rate”—and \( e^{ct} \) converts to \( e^t \).

The last step is to match the initial condition. The problem requires \( y = y_0 \) at \( t = 0 \). Our \( e^t \) starts from \( e^0 = 1 \). The constant of integration is needed now—the solutions are \( y = Ae^{ct} \). By choosing \( A = y_0 \), we match the initial condition and solve equation (1). The formula to remember is \( y_0 e^{ct} \).

\[
\text{The exponential law} \quad y = y_0 e^{ct} \text{ solves } y' = cy \text{ starting from } y_0.
\]

The rate of growth or decay is \( c \). May I call your attention to a basic fact? The formula \( y_0 e^{ct} \) contains three quantities \( y_0, c, t \). If two of them are given, plus one additional piece of information, the third is determined. Many applications have one of these three forms: find \( t \), find \( c \), find \( y_0 \).

1. Find the doubling time \( T \) if \( c = 1/10 \). At that time \( y_0 e^{cT} \) equals \( 2y_0 \):

\[
e^{cT} = 2 \quad \text{yields} \quad cT = \ln 2 \quad \text{so that} \quad T = \frac{\ln 2}{c} \approx \frac{.7}{.1}. \tag{2}
\]

The question asks for an exponent \( T \). The answer involves logarithms. If a cell grows at a continuous rate of \( c = 10\% \) per day, it takes about \( .7/.1 = 7 \) days to double in size. (Note that \(.7 \) is close to \( \ln 2 \).) If a savings account earns 10% continuous interest, it doubles in 7 years.

In this problem we knew \( c \). In the next problem we know \( T \).

2. Find the decay constant \( c \) for carbon-14 if \( y = \frac{1}{2}y_0 \) in \( T = 5568 \) years.

\[
e^{cT} = \frac{1}{2} \quad \text{yields} \quad cT = \ln \frac{1}{2} \quad \text{so that} \quad c \approx (\ln \frac{1}{2})/5568. \tag{3}
\]

After the half-life \( T = 5568 \), the factor \( e^{cT} \) equals \( \frac{1}{2} \). Now \( c \) is negative (\( \ln \frac{1}{2} = -\ln 2 \)).

Question 1 was about growth. Question 2 was about decay. Both answers found \( e^{cT} \) as the ratio \( y(T)/y(0) \). Then \( cT \) is its logarithm. Note how \( c \) sticks to \( T \). \( T \) has the units of time, \( c \) has the units of “1/time.”

Main point: The doubling time is \( (\ln 2)/c \), because \( cT = \ln 2 \). The time to multiply by \( e \) is \( 1/c \). The time to multiply by 10 is \( (\ln 10)/c \). The time to divide by \( e \) is \( -1/c \), when a negative \( c \) brings decay.

3. Find the initial value \( y_0 \) if \( c = 2 \) and \( y(1) = 5 \):

\[
y(t) = y_0 e^{ct} \quad \text{yields} \quad y_0 = y(t)e^{-ct} = 5e^{-2}.
\]
All we do is run the process backward. Start from 5 and go back to $y_0$. With time reversed, $e^t$ becomes $e^{-t}$. The product of $e^2$ and $e^{-2}$ is 1—growth forward and decay backward.

Equally important is $T + t$. **Go forward to time $T$ and go on to $T + t$:**

$$y(T + t) = y_0 e^{(T+t)}$$

which is $$(y_0 e^T) e^t.$$ (4)

Every step $t$, at the start or later, multiplies by the same $e^t$. This uses the fundamental property of exponentials, that $e^{t+u} = e^t e^u$.

**EXAMPLE 1** *Population growth from birth rate $b$ and death rate $d$ (both constant):*

$$\frac{dy}{dt} = by - dy = cy$$ (the net rate is $c = b - d$).

The population in this model is $y_0 e^t = y_0 e^{b-(-d)t}$. It grows when $b > d$ (which makes $c > 0$). One estimate of the growth rate is $c = 0.02$/year:

The earth's population doubles in about $T = \frac{\ln 2}{c} \approx \frac{.7}{.02} = 35$ years.

First comment: We predict the future based on $c$. We count the past population to find $c$. Changes in $c$ are a serious problem for this model.

Second comment: $y_0 e^t$ is not a whole number. You may prefer to think of bacteria instead of people. *(This section begins a major application of mathematics to economics and the life sciences.)* Malthus based his theory of human population on this equation $y' = cy$—and with large numbers a fraction of a person doesn't matter so much. To use calculus we go from discrete to continuous. The theory must fail when $t$ is very large, since populations cannot grow exponentially forever. Section 6.5 introduces the *logistic equation* $y' = cy - by^2$, with a competition term $-by^2$ to slow the growth.

Third comment: The dimensions of $b$, $c$, $d$ are "1/time." The dictionary gives birth rate = number of births per person in a unit of time. It is a *relative* rate—people divided by people and time. The product $ct$ is dimensionless and $e^{ct}$ makes sense (also dimensionless). Some texts replace $c$ by $\lambda$ (lambda). Then $1/\lambda$ is the growth time or decay time or drug elimination time or diffusion time.

**EXAMPLE 2** *Radioactive dating* A gram of charcoal from the cave paintings in France gives 0.97 disintegrations per minute. A gram of living wood gives 6.68 disintegrations per minute. Find the age of those Lascaux paintings.

The charcoal stopped adding radiocarbon when it was burned (at $t = 0$). The amount has decayed to $y_0 e^{-t}$. In living wood this amount is still $y_0$, because cosmic
6.3 Growth and Decay in Science and Economics

rays maintain the balance. Their ratio is $e^d = 0.97/6.68$. Knowing the decay rate $c$ from Question 2 above, we know the present time $t$:

$$ct = \ln \left( \frac{0.97}{6.68} \right) \quad \text{yields} \quad t = \frac{5568}{-0.7 \ln(0.97)} = 14,400 \text{ years.}$$

Here is a related problem—the age of uranium. Right now there is 140 times as much U-238 as U-235. Nearly equal amounts were created, with half-lives of $(4.5)10^9$ and $(0.7)10^9$ years. Question: How long since uranium was created? Answer: Find $t$ by substituting $c = (\ln \frac{1}{2})/(4.5)10^9$ and $C = (\ln \frac{1}{2})/(0.7)10^9$:

$$e^{ct/C} = 140 \Rightarrow ct - C = \ln 140 \Rightarrow t = \frac{\ln 140}{c - C} = 6(10^9) \text{ years.}$$

**EXAMPLE 3 Calculus in Economics: price inflation and the value of money**

We begin with two inflation rates—a continuous rate and an annual rate. For the price change $\Delta y$ over a year, use the annual rate:

$$\Delta y = (\text{annual rate}) \times (y) \times (\Delta t). \quad (5)$$

*Calculus applies the continuous rate to each instant $dt$. The price change is $dy$.*

$$dy = (\text{continuous rate}) \times (y) \times (dt). \quad (6)$$

Dividing by $dt$, this is a differential equation for the price:

$$\frac{dy}{dt} = (\text{continuous rate}) \times (y) = .05y.$$

The solution is $y_0e^{.05t}$. Set $t = 1$. Then $e^{.05} \approx 1.0513$ and the annual rate is 5.13%.

When you ask a bank what interest they pay, they give both rates: 8% and 8.33%. The higher one they call the “effective rate.” It comes from compounding (and depends how often they do it). If the compounding is continuous, every $dt$ brings an increase of $dy$—and $e^{.08}$ is near 1.0833.

Section 6.6 returns to compound interest. The interval drops from a month to a day to a second. That leads to $(1 + 1/n)^n$, and in the limit to $e$. Here we compute the effect of 5% continuous interest:

**Future value** A dollar now has the same value as $e^{.05T}$ dollars in $T$ years.

**Present value** A dollar in $T$ years has the same value as $e^{-0.05T}$ dollars now.

**Doubling time** Prices double ($e^{.05T} = 2$) in $T = \ln 2/0.05 \approx 14$ years.

With no compounding, the doubling time is 20 years. Simple interest adds on 20 times 5% = 100%. With continuous compounding the time is reduced by the factor $\ln 2 \approx .7$, regardless of the interest rate.

**EXAMPLE 4** In 1626 the Indians sold Manhattan for $24. Our calculations indicate that they knew what they were doing. Assuming 8% compound interest, the original $24 is multiplied by $e^{.08t}$. After $t = 365$ years the multiplier is $e^{29.2}$ and the $24 has grown to 115 trillion dollars. With that much money they could buy back the land and pay off the national debt.

This seems farfetched. Possibly there is a big flaw in the model. It is absolutely true that Ben Franklin left money to Boston and Philadelphia, to be invested for 200 years. In 1990 it yielded millions (not trillions, that takes longer). Our next step is a new model.
Question  How can you estimate $e^{29.2}$ with a $24$ calculator (log but not ln)?

Answer  Multiply $29.2$ by log base $10$ of $e = .434$ to get $12.7$. This is the exponent to base $10$. After that base change, we have $10^{12.7}$ or more than a trillion.

**GROWTH OR DECAY WITH A SOURCE TERM**

The equation $y' = y$ will be given a new term. Up to now, all growth or decay has started from $y_0$. No deposit or withdrawal was made later. The investment grew by itself—a pure exponential. The new term $s$ allows you to add or subtract from the account. It is a "source"—or a "sink" if $s$ is negative. The source $s = 5$ adds $5dt$, proportional to $dt$ but not to $y$:

**Constant source:** $dy/dt = y + 5$ starting from $y = y_0$.

Notice $y$ on both sides! My first guess $y = e^{t+5}$ failed completely. Its derivative is $e^{t+5}$ again, which is not $y + 5$. The class suggested $y = e^t + 5t$. But its derivative $e^t + 5$ is still not $y + 5$. We tried other ways to produce $5$ in $dy/dt$. This idea is doomed to failure. Finally we thought of $y = Ae^t - 5$. That has $y' = Ae^t = y + 5$ as required.

Important: $A$ is not $y_0$. Set $t = 0$ to find $y_0 = A - 5$. The source contributes $5e^t - 5$:

The solution is $(y_0 + 5)e^t - 5$. That is the same as $y_0e^t + 5(e^t - 1)$.

$s = 5$ multiplies the growth term $e^t - 1$ that starts at zero. $y_0e^t$ grows as before.

**EXAMPLE 5  $dy/dt = -y + 5$ has $y = (y_0 - 5)e^{-t} + 5$. This is $y_0e^{-t} + 5(1 - e^{-t})$.**

That final term from the source is still positive. The other term $y_0e^{-t}$ decays to zero. The limit as $t \to \infty$ is $y_\infty = 5$. A negative $c$ leads to a steady state $y_\infty$.

Based on these examples with $c = 1$ and $c = -1$, we can find $y$ for any $c$ and $s$.

**EQUATION WITH SOURCE  $dy/dt = cy + s$ starts from $y = y_0$ at $t = 0$. (7)**

The source could be a deposit of $s = $1000/year, after an initial investment of $y_0 = $8000. Or we can withdraw funds at $s = -$200/year. The units are “dollars per year” to match $dy/dt$. The equation feeds in $1000$ or removes $200$ continuously—not all at once.

Note again that $y = e^{(c+s)t}$ is not a solution. Its derivative is $(c+s)y$. The combination $y = e^{ct} + s$ is also not a solution (but closer). The analysis of $y' = cy + s$ will be our main achievement for first-order differential equations (in this section). The equation is not restricted to finance—far from it—but that produces excellent examples.

I propose to find $y$ in four ways. You may feel that one way is enough.† The first way is the fastest—only three lines—but please give the others a chance. There is no point in preparing for real problems if we don’t solve them.

Solution by Method 1  (fast way) Substitute the combination $y = Ae^t + B$. The solution has this form—exponential plus constant. From two facts we find $A$ and $B$:

the equation $y' = cy + s$ gives $cAe^t = c(Ae^t + B) + s$

the initial value at $t = 0$ gives $A + B = y_0$.

†My class says one way is more than enough. They just want the answer. Sometimes I cave in and write down the formula: $y$ is $y_0e^c$ plus $s(e^t - 1)/c$ from the source term.
The first line has $cAe^{ct}$ on both sides. Subtraction leaves $cB + s = 0$, or $B = -s/c$. Then the second line becomes $A = y_0 - B = y_0 + (s/c)$.

**KEY FORMULA**

$$y = \left(y_0 + \frac{s}{c}\right)e^{ct} - \frac{s}{c} \quad \text{or} \quad y = y_0e^{ct} + \frac{s}{c}(e^{ct} - 1).$$  \hspace{1cm} (8)

With $s = 0$ this is the old solution $y_0e^{ct}$ (no source). The example with $c = 1$ and $s = 5$ produced $(y_0 + 5)e^t - 5$. Separating the source term gives $y_0e^t + 5(e^t - 1)$.

**Solution by Method 2** (slow way) The input $y_0$ produces the output $y_0e^{ct}$. After $t$ years any deposit is multiplied by $e^t$. That also applies to deposits made after the account is opened. If the deposit enters at time $T$, the growing time is only $t - T$. Therefore the multiplying factor is only $e^{c(t-T)}$. This growth factor applies to the small deposit (amount $s\,dT$) made between time $T$ and $T + dT$.

Now add up all outputs at time $t$. The output from $y_0$ is $y_0e^{ct}$. The small deposit $s\,dT$ near time $T$ grows to $e^{c(t-T)}s\,dT$. The total is an integral:

$$y(t) = y_0e^{ct} + \int_{T = 0}^{t} e^{c(t-T)}s\,dT.$$  \hspace{1cm} (9)

This principle of Duhamel would still apply when the source $s$ varies with time. Here $s$ is constant, and the integral divides by $c$:

$$s \int_{T = 0}^{t} e^{c(t-T)}dT = \frac{se^{c(t-T)}}{e^{ct}} \bigg|_{0}^{t} = -\frac{s}{c} + \frac{s}{c}e^{ct}.$$  \hspace{1cm} (10)

That agrees with the source term from Method 1, at the end of equation (8). There we looked for “exponential plus constant,” here we added up outputs.

Method 1 was easier. It succeeded because we knew the form $Ae^{ct} + B$—with “undetermined coefficients.” Method 2 is more complete. The form for $y$ is part of the output, not the input. The source $s$ is a continuous supply of new deposits, all growing separately. Section 6.5 starts from scratch, by directly integrating $y' = cy + s$.

**Remark** Method 2 is often described in terms of an integrating factor. First write the equation as $y' - cy = s$. Then multiply by a magic factor that makes integration possible:

$$(y' - cy)e^{-ct} = se^{-ct} \quad \text{multiply by the factor } e^{-ct}$$

$$ye^{-ct} \bigg|_{0}^{t} = -\frac{s}{c}e^{-ct} \bigg|_{0}^{t} \quad \text{integrate both sides}$$

$$ye^{-ct} - y_0 = -\frac{s}{c}(e^{-ct} - 1) \quad \text{substitute 0 and } t$$

$$y = e^{ct}y_0 + \frac{s}{c}(e^{ct} - 1) \quad \text{isolate } y \text{ to reach formula (8)}$$

The integrating factor produced a perfect derivative in line 1. I prefer Duhamel's idea, that all inputs $y_0$ and $s$ grow the same way. Either method gives formula (8) for $y$.

**THE MATHEMATICS OF FINANCE (AT A CONTINUOUS RATE)**

The question from finance is this: What inputs give what outputs? The inputs can come at the start by $y_0$, or continuously by $s$. The output can be paid at the end or continuously. There are six basic questions, two of which are already answered.

The future value is $y_0e^{ct}$ from a deposit of $y_0$. To produce $y$ in the future, deposit the present value $ye^{-ct}$. Questions 3–6 involve the source term $s$. We fix the continuous
rate at 5% per year \( c = .05 \), and start the account from \( y_0 = 0 \). The answers come fast from equation (8).

**Question 3** With deposits of \( s = $1000/\text{year} \), how large is \( y \) after 20 years?

\[
y = \frac{s}{c} (e^{ct} - 1) = \frac{1000}{.05} (e^{(.05)(20)} - 1) = 20,000(e - 1) \approx $34,400.
\]

One big deposit yields \( 20,000e \approx $54,000 \). The same 20,000 via \( s \) yields \( $34,400 \).

Notice a small by-product (for mathematicians). When the interest rate is \( c = 0 \), our formula \( s(e^t - 1)/c \) turns into 0/0. We are absolutely sure that depositing \$1000/year with no interest produces \$20,000 after 20 years. But this is not obvious from 0/0. By l'Hôpital's rule we take \( c \)-derivatives in the fraction:

\[
\lim_{c \to 0} \frac{s(e^t - 1)}{c} = \lim_{c \to 0} \frac{s te^t}{1} = st. \text{ This is } (1000)(20) = 20,000. \tag{11}
\]

**Question 4** What continuous deposit of \( s \) per year yields \$20,000 after 20 years?

\[
20,000 = \frac{s}{.05} (e^{(.05)(20)} - 1) \text{ requires } s = \frac{1000}{e - 1} \approx 582.
\]

Deposits of \$582 over 20 years total \$11,640. A single deposit of \( y_0 = 20,000/e = $7,360 \) produces the same \$20,000 at the end. Better to be rich at \( t = 0 \).

Questions 1 and 2 had \( s = 0 \) (no source). Questions 3 and 4 had \( y_0 = 0 \) (no initial deposit). Now we come to \( y = 0 \). In 5, everything is paid out by an annuity. In 6, everything is paid up on a loan.

**Question 5** What deposit \( y_0 \) provides \$1000/year for 20 years? End with \( y = 0 \).

\[
y = y_0 e^{ct} + \frac{s}{c} (e^{ct} - 1) = 0 \text{ requires } y_0 = -\frac{s}{c} (1 - e^{-ct}).
\]

Substituting \( s = -1000, c = .05, t = 20 \) gives \( y_0 \approx 12,640 \). If you win \$20,000 in a lottery, and it is paid over 20 years, the lottery only has to put in \$12,640. Even less if the interest rate is above 5%.

**Question 6** What payments \( s \) will clear a loan of \( y_0 = $20,000 \) in 20 years?

Unfortunately, \( s \) exceeds \$1000 per year. The bank gives up more than the \$20,000 to buy your car (and pay tuition). It also gives up the interest on that money. You pay that back too, but you don't have to stay even at every moment. Instead you repay at a constant rate for 20 years. Your payments mostly cover interest at the start and principal at the end. After \( t = 20 \) years you are even and your debt is \( y = 0 \).

This is like Question 5 (also \( y = 0 \)), but now we know \( y_0 \) and we want \( s \):

\[
y = y_0 e^{ct} + \frac{s}{c} (e^{ct} - 1) = 0 \text{ requires } s = -cy_0 e^{ct}/(e^{ct} - 1).
\]

The loan is \( y_0 = $20,000 \), the rate is \( c = .05/\text{year} \), the time is \( t = 20 \) years. Substituting in the formula for \( s \), your payments are \$1582 per year.

**Puzzle** How is \( s = $1582 \) for loan payments related to \( s = $582 \) for deposits?

\[
0 \rightarrow $582 \text{ per year} \rightarrow $20,000 \text{ and } $20,000 \rightarrow -$1582 \text{ per year} \rightarrow 0.
\]
That difference of exactly 1000 cannot be an accident. 1582 and 582 came from
\[
\frac{e^{1000}}{e-1} \text{ and } \frac{1000}{e-1} \text{ with difference } 1000 \frac{e^1 - 1}{e-1} = 1000.
\]

**Why?** Here is the real reason. Instead of repaying 1582 we can pay only 1000 (to keep even with the interest on 20,000). The other 582 goes into a separate account. After 20 years the continuous 582 has built up to 20,000 (including interest as in Question 4). From that account we pay back the loan.

Section 6.6 deals with daily compounding—which differs from continuous compounding by only a few cents. Yearly compounding differs by a few dollars.

**TRANSIENTS VS. STEADY STATE**

Suppose there is decay instead of growth. The constant \( c \) is negative and \( y_0 e^{ct} \) dies out. That is the "transient" term, which disappears as \( t \to \infty \). What is left is the "steady state." We denote that limit by \( y_\infty \).

Without a source, \( y_\infty \) is zero (total decay). When \( s \) is present, \( y_\infty = -s/c \):

\[
6J \text{ The solution } y = \left( y_0 + \frac{s}{c} \right) e^{ct} - \frac{s}{c} \text{ approaches } y_\infty = -\frac{s}{c} \text{ when } e^{ct} \to 0.
\]

At this steady state, the source \( s \) exactly balances the decay \( cy \). In other words \( cy + s = 0 \). From the left side of the differential equation, this means \( dy/dt = 0 \). There is no change. That is why \( y_\infty \) is steady.

Notice that \( y_\infty \) depends on the source and on \( c \)—but not on \( y_0 \).

**EXAMPLE 6** Suppose Bermuda has a birth rate \( b = .02 \) and death rate \( d = .03 \). The net decay rate is \( c = -.01 \). There is also immigration from outside, of \( s = 1200/\text{year} \). The initial population might be \( y_0 = 5 \text{ thousand or } y_0 = 5 \text{ million, but that number has no effect on } y_\infty \). The steady state is independent of \( y_0 \).

In this case \( y_\infty = -s/c = 1200/.01 = 120,000 \). The population grows to 120,000 if \( y_0 \) is smaller. It decays to 120,000 if \( y_0 \) is larger.

**EXAMPLE 7** Newton's Law of Cooling: \[
dy/dt = c(y - y_\infty).
\]

This is back to physics. The temperature of a body is \( y \). The temperature around it is \( y_\infty \). Then \( y \) starts at \( y_0 \) and approaches \( y_\infty \), following Newton's rule: The rate is proportional to \( y - y_\infty \). The bigger the difference, the faster heat flows.

The equation has \( -cy_\infty \) where before we had \( s \). That fits with \( y_\infty = -s/c \). For the solution, replace \( s \) by \( -cy_\infty \) in formula (8). Or use this new method:
Solution by Method 3   The new idea is to look at the difference $y - y_\infty$. Its derivative is $dy/dt$, since $y_\infty$ is constant. But $dy/dt$ is $c(y - y_\infty)$—this is our equation. The difference starts from $y_0 - y_\infty$, and grows or decays as a pure exponential:

$$\frac{d}{dt}(y - y_\infty) = c(y - y_\infty) \quad \text{has the solution} \quad (y - y_\infty) = (y_0 - y_\infty)e^{ct}. \quad (13)$$

This solves the law of cooling. We repeat Method 3 using the letters $s$ and $c$:

$$\frac{d}{dt}(y + \frac{s}{c}) = c(y + \frac{s}{c}) \quad \text{has the solution} \quad (y + \frac{s}{c}) = (y_0 + \frac{s}{c})e^{ct}. \quad (14)$$

Moving $s/c$ to the right side recovers formula (8). There is a constant term and an exponential term. In a differential equations course, those are the "particular solution" and the "homogeneous solution." In a calculus course, it's time to stop.

EXAMPLE 8   In a $70^\circ$ room, Newton's corpse is found with a temperature of $90^\circ$. A day later the body registers $80^\circ$. When did he stop integrating (at $98.6^\circ$)?

Solution   Here $y_\infty = 70$ and $y_0 = 90$. Newton's equation (13) is $y = 20e^{ct} + 70$. Then $y = 80$ at $t = 1$ gives $20e^{c} = 10$. The rate of cooling is $c = \ln \frac{1}{2}$. Death occurred when $20e^{c} + 70 = 98.6$ or $e^{c} = 1.43$. The time was $t = \ln 1.43/\ln \frac{1}{2} = \frac{1}{2}$ day earlier.

6.3 Exercises

Read-through exercises

If $y' = cy$ then $y(t) = _{\text{exponential}}$. If $dy/dt = 7y$ and $y_0 = 4$ then $y(t) = _{\text{exponential}}$. This solution reaches 8 at $t = _{\text{exponential}}$. If the doubling time is $T$ then $c = _{\text{exponential}}$. If $y' = 3y$ and $y(1) = 9$ then $y_0$ was _constant_. When $c$ is negative, the solution approaches _zero_ as $t \to \infty$.

The constant solution to $dy/dt = y + 6$ is $y = _{\text{exponential}}$. The general solution is $y = A e^{c} - 6$. If $y_0 = 4$ then $A = _{\text{exponential}}$. The solution of $dy/dt = cy + s$ starting from $y_0$ is $y = Ae^{ct} + B = _{\text{exponential}}$. The output from the source $s$ is _exponential_. An input at time $T$ grows by the factor _exponential_ at time $t$.

At $c = 10\%$, the interest in time $dt$ is $dy = _{\text{exponential}}$. This equation yields $y(t) = _{\text{exponential}}$. With a source term instead of $y_0$, a continuous deposit of $s = 4000$/year yields $y = _{\text{exponential}}$. The deposit required to produce 10,000 in 10 years is $s = _{\text{exponential}}$ (exactly or approximately). An income of 4000/year forever (it) comes from $y_0 = _{\text{exponential}}$. The deposit to give 4000/year for 20 years is $y_0 = _{\text{exponential}}$. The payment rate $s$ to clear a loan of 10,000 in 10 years is _exponential_.

The solution to $y' = -3y + s$ approaches $y_\infty = _{\text{exponential}}$.

Solve 1–4 starting from $y_0 = 1$ and from $y_0 = -1$. Draw both solutions on the same graph.

1 $\frac{dy}{dt} = 2t$  2 $\frac{dy}{dt} = -t$  3 $\frac{dy}{dt} = 2y$  4 $\frac{dy}{dt} = -y$

Solve 5–8 starting from $y_0 = 10$. At what time does $y$ increase to 100 or drop to 1?

5 $\frac{dy}{dt} = 4y$  6 $\frac{dy}{dt} = 4t$  7 $\frac{dy}{dt} = e^{4t}$  8 $\frac{dy}{dt} = e^{-4t}$

9 Draw a field of "tangent arrows" for $y' = -y$, with the solution curves $y = e^{-x}$ and $y = -e^{-x}$.

10 Draw a direction field of arrows for $y' = y - 1$, with solution curves $y = e^{x} + 1$ and $y = 1$.

Problems 11–27 involve $y_0e^{ct}$. They ask for $c$ or $t$ or $y_0$.

11 If a culture of bacteria doubles in two hours, how many hours to multiply by 10? First find $c$.

12 If bacteria increase by factor of ten in ten hours, how many hours to increase by 100? What is $c$?

13 How old is a skull that contains $\frac{1}{2}$ as much radiocarbon as a modern skull?

14 If a relic contains 90% as much radiocarbon as new material, could it come from the time of Christ?

15 The population of Cairo grew from 5 million to 10 million in 20 years. From $y' = cy$ find $c$. When was $y = 8$ million?

16 The populations of New York and Los Angeles are growing at 1% and 1.4% a year. Starting from 8 million (NY) and 6 million (LA), when will they be equal?
17 Suppose the value of $1 in Japanese yen decreases at 2% per year. Starting from $1 = Y240, when will 1 dollar equal 1 yen?

18 The effect of advertising decays exponentially. If 40% remember a new product after three days, find c. How long will 20% remember it?

19 If \( y = 1000 \) at \( t = 3 \) and \( y = 3000 \) at \( t = 4 \) (exponential growth), what is \( y_0 \) at \( t = 0 \)?

20 If \( y = 100 \) at \( t = 4 \) and \( y = 10 \) at \( t = 8 \) (exponential decay) when will \( y = 1 \)? What was \( y_0 \)?

21 Atmospheric pressure decreases with height according to \( dp/dh = -cp \). The pressures at \( h = 0 \) (sea level) and \( h = 20 \) km are 1013 and 50 millibars. Find \( c \). Explain why \( p = \sqrt{1013 \times 50} \) halfway up at \( h = 10 \).

22 For exponential decay show that \( y(t) \) is the square root of \( y(0) \times y(2t) \). How could you find \( y(3t) \) from \( y(t) \) and \( y(2t) \)?

23 Most drugs in the bloodstream decay by \( y' = ky \) (first-order kinetics). (a) The half-life of morphine is 3 hours. Find its decay constant \( c \) (with units). (b) The half-life of nicotine is 2 hours. After a six-hour flight what fraction remains?

24 How often should a drug be taken if its dose is 3 mg, it is cleared at \( c = .01 \) hour, and 1 mg is required in the bloodstream at all times? (The doctor decides this level based on body size.)

25 The antiseizure drug dilantin has constant clearance rate \( y' = -a \) until \( y = y_1 \). Then \( y' = -ay(y_1) \). Solve for \( y(t) \) in two pieces from \( y_0 \). When does \( y \) reach \( y_1 \)?

26 The actual elimination of nicotine is multiexponential: \( y = Ae^{-t} + Ce^{Ct} \). The first-order equation \((d/dt - c)y = 0 \) changes to the second-order equation \((d/dt - c)(d/dt - C)y = 0 \). Write out this equation starting with \( y' \), and show that it is satisfied by the given \( y \).

27 True or false. If false, say what's true.
(a) The time for \( y = e^t \) to double is \((\ln 2)/\ln c)\)
(b) If \( y' = cy \) and \( x' = cx \) then \((y + x)' = 2c(y + z)\)
(c) If \( y' = cy \) and \( x' = cz \) then \((y/z)' = 0\)
(d) If \( y' = cy \) and \( x' = Cz \) then \((yz)' = (c + C)yz\).

28 A rocket has velocity \( v \). Burnt fuel of mass \( \Delta m \) leaves at velocity \( v - 7 \). Total momentum is constant:
\[
m_0 = (m - \Delta m)(v + \Delta v) + \Delta m(v - 7).
\]
What differential equation connects \( m \) to \( v \)? Solve for \( \Delta m \) not \( v(t) \), starting from \( v_0 = 20 \) and \( m_0 = 4 \).

Problems 29–36 are about solutions of \( y' = cy + s \).

29 Solve \( y' = 3y + 1 \) with \( y_0 = 0 \) by assuming \( y = Ae^{yt} + B \) and determining \( A \) and \( B \).

30 Solve \( y' = 8 - y \) starting from \( y_0 \) and \( y = Ae^{-t} + B \).

Solve 31–34 with \( y_0 = 0 \) and graph the solution.

31 \( \frac{dy}{dt} = y + 1 \)

32 \( \frac{dy}{dt} = y - 1 \)

33 \( \frac{dy}{dt} = -y + 1 \)

34 \( \frac{dy}{dt} = -y - 1 \)

35 (a) What value \( y = \) constant solves \( dy/dt = -2y + 12 \)?

(b) Find the solution with an arbitrary constant \( A \).

(c) What solutions start from \( y_0 = 0 \) and \( y_0 = 10 \)?

(d) What is the steady state \( y_0 \)?

36 Choose \( \pm \) signs in \( dy/dt = \pm 3y + 6 \) to achieve the following results starting from \( y_0 = 1 \). Draw graphs.
(a) \( y \) increases to \( \infty \)
(b) \( y \) increases to \( 2 \)
(c) \( y \) decreases to \( -2 \)
(d) \( y \) decreases to \( -\infty \)

37 What value \( y = \) constant solves \( dy/dt = 4 - y \)? Show that \( y(t) = Ae^{-t} + 4 \) is also a solution. Find \( y(1) \) and \( y_0 \) if \( y_0 = 3 \).

38 Solve \( y' = y + e^t \) from \( y_0 = 0 \) by Method 2, where the deposit \( e^t \) at time \( T \) is multiplied by \( e^{-t} \). The total output at time \( t \) is \( y(t) = \int_0^t e^t e^{-t} e^{-t} \). Substitute back to check \( y' = y + e^t \).

39 Rewrite \( y' = y + e^t \) as \( y' - y = e^t \). Multiplying by \( e^{-t} \), the left side is the derivative of \( \int e^{-t} \). Integrate both sides from \( y_0 = 0 \) to find \( y(t) \).

40 Solve \( y' = y + 1 \) from \( y_0 = 0 \) by rewriting as \( y' + y = 1 \), multiplying by \( e^t \), and integrating both sides.

41 Solve \( y' = y + t \) from \( y_0 = 0 \) by assuming \( y = Ae^t + Bt + C \).

Problems 42–57 are about the mathematics of finance.

42 Dollar bills decrease in value at \( c = -0.04/\text{year} \) because of inflation. If you hold \$1000, what is the decrease in \$ at years? At what rate \( s \) should you print money to keep even?

43 If a bank offers annual interest of \( 7\% \) or continuous interest of \( 7\% \), which is better?

44 What continuous interest rate is equivalent to an annual rate of \( 9\% \)? Extra credit: Telephone a bank for both rates and check their calculation.

45 At \( 100\% \) interest \((c = 1)\) how much is a continuous deposit of \$ per year worth after one year? What initial deposit \( y_0 \) would have produced the same output?

46 To have \$50,000 for college tuition in 20 years, what gift \( y_0 \) should a grandparent make now? Assume \( c = 10\% \). What continuous deposit should a parent make during 20 years? If the parent saves \$1000 per year, when does he or she reach \$50,000 and retire?
We have given first place to $e^x$ and a lower place to $\ln x$. In applications that is absolutely correct. But logarithms have one important theoretical advantage (plus many applications of their own). The advantage is that the derivative of $\ln x$ is $1/x$, whereas the derivative of $e^x$ is $e^x$. We can't define $e^x$ as its own integral, without circular reasoning. But we can and do define $\ln x$ (the natural logarithm) as the integral of the $"-1\ power"$ which is $1/x$:

$$\ln x = \int_1^x \frac{1}{t} \, dt \quad \text{or} \quad \ln y = \int_1^y \frac{1}{u} \, du. \quad (1)$$

Note the dummy variables, first $x$ then $u$. Note also the live variables, first $x$ then $y$. Especially note the lower limit of integration, which is 1 and not 0. The logarithm is the area measured from 1. Therefore $\ln 1 = 0$ at that starting point—as required.
Earlier chapters integrated all powers except this "-1 power." The logarithm is that missing integral. The curve in Figure 6.11 has height \( y = 1/x \)—it is a hyperbola. At \( x = 0 \) the height goes to infinity and the area becomes infinite: \( \log 0 = -\infty \). The minus sign is because the integral goes backward from 1 to 0. The integral does not extend past zero to negative \( x \). We are defining \( \ln x \) only for \( x > 0 \).†

![Figure 6.11 Logarithm as area. Neighbors \( \ln a + \ln b = \ln ab \). Equal areas: \(-\ln \frac{1}{2} = \ln 2 = \frac{1}{2} \ln 4\).](image)

With this new approach, \( \ln x \) has a direct definition. It is an integral (or an area). Its two key properties must follow from this definition. That step is a beautiful application of the theory behind integrals.

**Property 1:** \( \ln ab = \ln a + \ln b \). The areas from 1 to \( a \) and from \( a \) to \( ab \) combine into a single area (1 to \( ab \) in the middle figure):

\[
\text{Neighboring areas: } \int_1^a \frac{1}{x} \, dx + \int_a^{ab} \frac{1}{x} \, dx = \int_1^{ab} \frac{1}{x} \, dx.
\]

The right side is \( \ln ab \), from definition (1). The first term on the left is \( \ln a \). The problem is to show that the second integral (\( a \) to \( ab \)) is \( \ln b \):

\[
\int_a^{ab} \frac{1}{x} \, dx = \int_1^b \frac{1}{u} \, du = \ln b.
\]

We need \( u = 1 \) when \( x = a \) (the lower limit) and \( u = b \) when \( x = ab \) (the upper limit). The choice \( u = x/a \) satisfies these requirements. Substituting \( x = au \) and \( dx = a \, du \) yields \( dx/x = du/u \). Equation (3) gives \( \ln b \), and equation (2) is \( \ln a + \ln b = \ln ab \).

**Property 2:** \( \ln b^n = n \ln b \). These are the left and right sides of

\[
\int_1^{b^n} \frac{1}{x} \, dx = n \int_1^b \frac{1}{u} \, du.
\]

This comes from the substitution \( x = u^n \). The lower limit \( x = 1 \) corresponds to \( u = 1 \), and \( x = b^n \) corresponds to \( u = b \). The differential \( dx \) is \( mu^{n-1} \, du \). Dividing by \( x = u^n \) leaves \( dx/x = n \, du/du \). Then equation (4) becomes \( \ln b^n = n \ln b \).

Everything comes logically from the definition as an area. Also definite integrals:

**EXAMPLE 1** Compute \( \int_x^{3x} \frac{1}{t} \, dt \). Solution: \( \ln 3x - \ln x = \ln \frac{3x}{x} = \ln 3 \).

**EXAMPLE 2** Compute \( \int_1^{x^{1/3}} \frac{1}{x} \, dx \). Solution: \( \ln 1 - \ln .1 = \ln 10 \). (Why?)

†The logarithm of \(-1\) is \( \pi i \) (an imaginary number). That is because \( e^{\pi i} = -1 \). The logarithm of \( i \) is also imaginary—it is \( \frac{1}{2} \pi i \). In general, logarithms are complex numbers.
EXAMPLE 3  Compute \( \int_{1}^{e^2} \frac{1}{u} \, du \). Solution: \( \ln e^2 = 2 \). The area from 1 to \( e^2 \) is 2.

Remark  While working on the theory this is a chance to straighten out old debts. The book has discussed and computed (and even differentiated) the functions \( e^x \) and \( b^x \) and \( x^n \), without defining them properly. When the exponent is an irrational number like \( \pi \), how do we multiply \( e \) by itself \( \pi \) times? One approach (not taken) is to come closer and closer to \( \pi \) by rational exponents like \( 22/7 \). Another approach (taken now) is to determine the number \( e^\pi = 23.1 \ldots \) by its logarithm.† Start with \( e \) itself:

\[
e \text{ is (by definition) the number whose logarithm is 1}
\]

\[
e^\pi \text{ is (by definition) the number whose logarithm is } \pi.
\]

When the area in Figure 6.12 reaches 1, the basepoint is \( e \). When the area reaches \( \pi \), the basepoint is \( e^\pi \). We are constructing the inverse function (which is \( e^x \)). But how do we know that the area reaches \( \pi \) or 1000 or \(-1000\) at exactly one point? (The area is 1000 far out at \( e^{1000} \). The area is \(-1000\) very near zero at \( e^{-1000} \).) To define \( e \) we have to know that somewhere the area equals 1!

For a proof in two steps, go back to Figure 6.11c. The area from 1 to 2 is more than \( \frac{1}{2} \) (because \( 1/x \) is more than \( \frac{1}{2} \) on that interval of length one). The combined area from 1 to 4 is more than 1. We come to area = 1 before reaching 4. (Actually at \( e = 2.718 \ldots \) ) Since \( 1/x \) is positive, the area is increasing and never comes back to 1.

To double the area we have to square the distance. The logarithm creeps upwards:

\[
\ln x \to \infty \quad \text{but} \quad \frac{\ln x}{x} \to 0.
\]

The logarithm grows slowly because \( e^x \) grows so fast (and vice versa—they are inverses). Remember that \( e^x \) goes past every power \( x^n \). Therefore \( \ln x \) is passed by every root \( x^{1/n} \). Problems 60 and 61 give two proofs that \( (\ln x)/x^{1/n} \) approaches zero.

We might compare \( \ln x \) with \( \sqrt{x} \). At \( x = 10 \) they are close (2.3 versus 3.2). But out at \( x = e^{10} \) the comparison is 10 against \( e^5 \), and \( \ln x \) loses to \( \sqrt{x} \).

†Chapter 9 goes on to imaginary exponents, and proves the remarkable formula \( e^{\pi i} = -1 \).
The limiting cases $\ln 0 = -\infty$ and $\ln \infty = +\infty$ are important. More important are logarithms near the starting point $\ln 1 = 0$. Our question is: What is $\ln (1 + x)$ for $x$ near zero? The exact answer is an area. The approximate answer is much simpler. If $x$ (positive or negative) is small, then

$$\ln (1 + x) \approx x \quad \text{and} \quad e^x \approx 1 + x.$$ (6)

The calculator gives $\ln 1.01 = .0099503$. This is close to $x = .01$. Between 1 and $1 + x$ the area under the graph of $1/x$ is nearly a rectangle. Its base is $x$ and its height is 1. So the curved area $\ln (1 + x)$ is close to the rectangular area $x$. Figure 6.14 shows how a small triangle is chopped off at the top.

The difference between .0099503 (actual) and .01 (linear approximation) is $-.0000497$. That is predicted almost exactly by the second derivative: $\frac{1}{2}$ times $(\Delta x)^2$ times $(\ln x)''$ is $\frac{1}{2}(.01)^2(-1) = -.00005$. This is the area of the small triangle!

$$\ln(1 + x) \approx \text{rectangular area minus triangular area} = x - \frac{1}{2}x^2.$$  

The remaining mistake of .0000003 is close to $\frac{1}{2}x^3$ (Problem 65).

May I switch to $e^x$? Its slope starts at $e^0 = 1$, so its linear approximation is $1 + x$. Then $\ln (e^x) = \ln (1 + x) \approx x$. Two wrongs do make a right: $\ln (e^x) = x$ exactly.

The calculator gives $e^{01}$ as 1.0100502 (actual) instead of 1.01 (approximation). The second-order correction is again a small triangle: $\frac{1}{2}x^2 = .00005$. The complete series for $\ln (1 + x)$ and $e^x$ are in Sections 10.1 and 6.6:

$$\ln (1 + x) = x - x^2/2 + x^3/3 - \ldots \quad e^x = 1 + x + x^2/2 + x^3/6 + \ldots.$$ 

**DERIVATIVES BASED ON LOGARITHMS**

Logarithms turn up as antiderivatives very often. To build up a collection of integrals, we now differentiate $\ln u(x)$ by the chain rule.

6K The derivative of $\ln x$ is $\frac{1}{x}$. The derivative of $\ln u(x)$ is $\frac{1}{u \, dx}$. 

The slope of $\ln x$ was hard work in Section 6.2. With its new definition (the integral of $1/x$) the work is gone. By the Fundamental Theorem, the slope must be $1/x$.

For $\ln u(x)$ the derivative comes from the chain rule. The inside function is $u$, the outside function is $\ln$. (Keep $u > 0$ to define $\ln u$.) The chain rule gives

$$\frac{d}{dx} \ln cx = \frac{1}{cx} \cdot c = \frac{1}{x} \quad \frac{d}{dx} \ln x^3 = 3x^2/x^3 = \frac{3}{x}$$

$$\frac{d}{dx} \ln (x^2 + 1) = 2x/(x^2 + 1) \quad \frac{d}{dx} \ln \cos x = -\sin x/\cos x = -\tan x$$

$$\frac{d}{dx} \ln e^x = e^x/e^x = 1 \quad \frac{d}{dx} \ln (\ln x) = \frac{1}{\ln x \cdot x}.$$ 

Those are worth another look, especially the first. Any reasonable person would expect the slope of $\ln 3x$ to be $3/x$. Not so. The 3 cancels, and $\ln 3x$ has the same slope as $\ln x$. (The real reason is that $\ln 3x = \ln 3 + \ln x$.) The antiderivative of $3/x$ is not $\ln 3x$ but $3 \ln x$, which is $\ln x^3$. 

Before moving to integrals, here is a new method for derivatives: logarithmic differentiation or LD. It applies to products and powers. The product and power rules are always available, but sometimes there is an easier way.

Main idea: The logarithm of a product \( p(x) \) is a sum of logarithms. Switching to \( \ln p \), the sum rule just adds up the derivatives. But there is a catch at the end, as you see in the example.

**EXAMPLE 4** Find \( dp/dx \) if \( p(x) = x^{\sqrt{x+1}} \). Here \( \ln p(x) = x \ln x + \frac{1}{2} \ln(x+1) \).

Take the derivative of \( \ln p \):

\[
\frac{dp}{dx} = \frac{1}{p} \cdot \frac{d}{dx} \left( x \ln x + \frac{1}{2} \ln(x+1) \right).
\]

Now multiply by \( p(x) \):

\[
\frac{dp}{dx} = p \left( 1 + \ln x + \frac{1}{2(x+1)} \right).
\]

The catch is that last step. Multiplying by \( p \) complicates the answer. This can't be helped—logarithmic differentiation contains no magic. The derivative of \( p = fg \) is the same as from the product rule: \( \ln p = \ln f + \ln g \) gives

\[
\frac{p'}{p} = \frac{f'}{f} + \frac{g'}{g} \quad \text{and} \quad p' = p \left( \frac{f'}{f} + \frac{g'}{g} \right) = f'g + fg'.
\]

For \( p = xe^x \sin x \), with three factors, the sum has three terms:

\[
\ln p = \ln x + x + \ln \sin x \quad \text{and} \quad p' = x + 1 + \frac{\cos x}{\sin x}.
\]

We multiply \( p \) times \( p'/p \) (the derivative of \( \ln p \)). Do the same for powers:

**EXAMPLE 5** \( p = x^{1/x} \Rightarrow \ln p = \frac{1}{x} \ln x \Rightarrow \frac{dp}{dx} = p \left[ -\frac{1}{x^2} + \frac{\ln x}{x^2} \right] \).

**EXAMPLE 6** \( p = x^{\ln x} \Rightarrow \ln p = (\ln x)^2 \Rightarrow \frac{dp}{dx} = p \left[ \frac{2 \ln x}{x} \right] \).

**EXAMPLE 7** \( p = x^{\ln x} \Rightarrow \ln p = \frac{1}{\ln x} \ln x = 1 \Rightarrow \frac{dp}{dx} = 0 \) (!)

**INTEGRALS BASED ON LOGARITHMS**

Now comes an important step. Many integrals produce logarithms. The foremost example is \( \frac{1}{x} \), whose integral is \( \ln x \). In a certain way that is the only example, but its range is enormously extended by the chain rule. The derivative of \( \ln u(x) \) is \( u'/u \), so the integral goes from \( u'/u \) back to \( \ln u \):

\[
\int \frac{du}{u} \frac{dx}{dt} = \ln u(x) \quad \text{or equivalently} \quad \int \frac{du}{u} = \ln u.
\]

Try to choose \( u(x) \) so that the integral contains \( du/dx \) divided by \( u \).

**EXAMPLES**

\[
\int \frac{dx}{x+7} = \ln|x + 7| \quad \int \frac{dx}{cx + 7} = \frac{1}{c} \ln|cx + 7|.
\]
Final remark  When $u$ is negative, $\ln u$ cannot be the integral of $1/u$. The logarithm is not defined when $u < 0$. But the integral can go forward by switching to $-u$:

$$\int \frac{du}{-u} \, dx = \int -\frac{du}{u} \, dx = \ln(-u). \quad (8)$$

Thus $\ln(-u)$ succeeds when $\ln u$ fails.$^\dagger$ The forbidden case is $u = 0$. The integrals $\ln u$ and $\ln(-u)$, on the plus and minus sides of zero, can be combined as $\ln|u|$. Every integral that gives a logarithm allows $u < 0$ by changing to the absolute value $|u|$:

$$\int_{-e}^{e} \frac{dx}{x} = \left[ \ln|x| \right]_{-e}^{e} = \ln 1 - \ln e \quad \int_{2}^{5} \frac{dx}{x-5} = \left[ \ln|x-5| \right]_{2}^{5} = \ln 1 - \ln 3.$$  

The areas are $-1$ and $-\ln 3$. The graphs of $1/x$ and $1/(x-5)$ are below the $x$ axis. We do not have logarithms of negative numbers, and we will not integrate $1/(x-5)$ from 2 to 6. That crosses the forbidden point $x = 5$, with infinite area on both sides.

The ratio $du/u$ leads to important integrals. When $u = \cos x$ or $u = \sin x$, we are integrating the tangent and cotangent. When there is a possibility that $u < 0$, write the integral as $\ln|u|$.

$$\int \tan x \, dx = \int \frac{\sin x}{\cos x} \, dx = -\ln |\cos x| \quad \int \frac{x \, dx}{x^2 + 7} = \frac{1}{2} \ln(x^2 + 7)$$

$$\int \cot x \, dx = \int \frac{\cos x}{\sin x} \, dx = \ln |\sin x| \quad \int \frac{dx}{x \ln x} = \ln |\ln x|$$

Now we report on the secant and cosecant. The integrals of $1/\cos x$ and $1/\sin x$ also surrender to an attack by logarithms — based on a crazy trick:

$$\int \sec x \, dx = \int \sec x \left( \frac{\sec x + \tan x}{\sec x + \tan x} \right) \, dx = \ln |\sec x + \tan x|. \quad (9)$$

$$\int \csc x \, dx = \int \csc x \left( \frac{\csc x - \cot x}{\csc x - \cot x} \right) \, dx = \ln |\csc x - \cot x|. \quad (10)$$

Here $u = \sec x + \tan x$ is in the denominator; $du/dx = \sec x \tan x + \sec^2 x$ is above it. The integral is $\ln |u|$. Similarly (10) contains $du/dx$ over $u = \csc x - \cot x$.

In closing we integrate $\ln x$ itself. The derivative of $x \ln x$ is $\ln x + 1$. To remove the extra 1, subtract $x$ from the integral: $\int \ln x \, dx = x \ln x - x$.

In contrast, the area under $1/(\ln x)$ has no elementary formula. Nevertheless it is the key to the greatest approximation in mathematics—the prime number theorem. The area $\int_{a}^{b} \frac{dx}{\ln x}$ is approximately the number of primes between $a$ and $b$. Near $e^{1000}$, about 1/1000 of the integers are prime.

### 6.4 Exercises

**Read-through questions**

The natural logarithm of $x$ is $\int_{1}^{x} \frac{dx}{x}$. This definition leads to $\ln xy = d$ and $\ln x^e = e$. Then $e$ is the number whose logarithm (area under $1/x$ curve) is $a$. Similarly $e^x$ is now defined as the number whose natural logarithm is $c$. As $x \to \infty$, $\ln x$ approaches $f$. But the ratio $(\ln x)/\sqrt{x}$ approaches $g$. The domain and range of $\ln x$ are $h$.

The derivative of $\ln x$ is $i$. The derivative of $\ln(1 + x)$ is $j$.

$^\dagger$The integral of $1/x$ (odd function) is $\ln |x|$ (even function). Stay clear of $x = 0$. 

is \_\_\_\_. The tangent approximation to $\ln(1 + x)$ at $x = 0$ is \_\_\_. The quadratic approximation is \_\_\_. The quadratic approximation to $e^x$ is \_\_\_.

The derivative of $\ln u(x)$ by the chain rule is \_\_\_. Thus $[\ln \cos x]' = \_\_\_\_. An antiderivative of $\tan x$ is \_\_\_. The product $p = xe^{5x}$ has $p' = \_\_\_\_. The derivative of this equation is \_\_\_. Multiplying by $p$ gives $p' = \_\_\_\_\_\_\_\_. Which is LD or logarithmic differentiation.

The integral of $u(x)/u(x)$ is \_\_\_. The integral of $2x/(x^2 + 4)$ is \_\_\_. The integral of $1/cx$ is \_\_\_. The integral of $1/(ct + s)$ is \_\_\_. The integral of $1/\cos x$, after a trick, is \_\_\_. We should write $\ln 1x$ for the antiderivative of $1/x$, since this allows \_\_\_. Similarly $\int du/u$ should be written \_\_\_\_\_\_\_\_.

Find the derivative $dy/dx$ in 1–10.

1. $y = \ln(2x)$
2. $y = \ln(2x + 1)$
3. $y = (\ln x)^{-1}$
4. $y = (\ln x)/x$
5. $y = x \ln x - x$
6. $y = \log_{10} x$
7. $y = \ln(\sin x)$
8. $y = \ln(\ln x)$
9. $y = 7 \ln 4x$
10. $y = \ln((4x)^3)$

Find the indefinite (or definite) integral in 11–24.

11. $\int \frac{dt}{3t}$
12. $\int \frac{dx}{1 + x}$
13. $\int_0^1 \frac{dx}{3 + x}$
14. $\int_0^1 \frac{dt}{3 + 2t}$
15. $\int_0^x \frac{xdx}{x^2 + 1}$
16. $\int_0^x \frac{x^2 dx}{x^2 + 1}$
17. $\int_0^x \frac{dx}{2x(\ln x)}$
18. $\int_0^x \frac{e^{x dx}}{1 \times (\ln x)^2}$
19. $\int_0^n \frac{\cos x dx}{\sin x}$
20. $\int_0^n \frac{\tan x dx}{\cot 3x dx}$
21. $\int_0^n \frac{\tan 3x dx}{\cot 3x dx}$
22. $\int_0^n \frac{\cot 3x dx}{\tan 3x dx}$
23. $\int_0^n \frac{(\ln x)^2 dx}{x}$
24. $\int_0^n \frac{dx}{x(\ln x)(\ln \ln x)}$
25. Graph $y = \ln(1 + x)$
26. Graph $y = \ln(\sin x)$

Compute $dy/dx$ by differentiating $\ln y$. This is LD:

27. $y = \sqrt{x^2 + 1}$
28. $y = \sqrt{x^2 + 1} \sqrt{x^2 - 1}$
29. $y = e^{\ln x}$
30. $y = x^{-1/x}$
31. $y = e^{e^x}$
32. $y = x^e$
33. $y = x^{e^x}$
34. $y = (\sqrt{x}) \sqrt[3]{x} \sqrt[5]{x}$
35. $y = x^{-1/ln x}$
36. $y = e^{-\ln x}$

Evaluate 37–42 by any method.

37. $\int_0^{10} \frac{dt}{s} - \int_{5x}^{10} \frac{dt}{s}$
38. $\int_1^x \frac{dx}{1 + x} + \int_{-2}^0 \frac{dx}{x}$
39. $\int_1^x \frac{dt}{t}$
40. $\int_1^x \frac{dt}{t}$
41. $\int_x \ln(\sec x + \tan x)$
42. $\int_{x}^y \frac{dx}{\cos x}$

Verify the derivatives 43–46, which give useful antiderivatives:

43. $\frac{d}{dx} \ln(x + \sqrt{x^2 + 1}) = \frac{1}{\sqrt{1 + x^2}}$
44. $\frac{d}{dx} \ln(x - a) = \frac{2a}{(x^2 - a^2)}$
45. $\frac{d}{dx} \ln \left(\frac{1 + \sqrt{1 - x^2}}{x}\right) = \frac{1}{x \sqrt{1 - x^2}}$
46. $\frac{d}{dx} \ln(x + \sqrt{x^2 - a^2}) = \frac{1}{\sqrt{x^2 - a^2}}$

Estimate 47–50 to linear accuracy, then quadratic accuracy, by $e^x \approx 1 + x + \frac{1}{2} x^2$. Then use a calculator.

47. $\ln(1.1)$
48. $e^x$
49. $\ln(0.99)$
50. $e^2$

51. Compute $\lim_{x \to 0} \frac{\ln(1 + x)}{x}$
52. Compute $\lim_{x \to 0} \frac{e^x - 1}{x}$
53. Compute $\lim_{x \to 0} \frac{\log(1 + x)}{x}$
54. Compute $\lim_{x \to 0} \frac{b^x - 1}{x}$

55. Find the area of the “hyperbolic quarter-circle” enclosed by $x = 2$ and $y = 2$ above $y = 1/x$.

56. Estimate the area under $y = 1/x$ from 4 to 8 by four upper rectangles and four lower rectangles. Then average the answers (trapezoidal rule). What is the exact area?

57. Why is $1/2 + 1/3 + \cdots + 1/n$ near $\ln n$? Is it above or below?

58. Prove that $\ln x \leq 2(\sqrt{x} - 1)$ for $x > 1$. Compare the integrals of $1/t$ and $1/\sqrt{t}$, from 1 to $x$.

59. Dividing by $x$ in Problem 58 gives $(\ln x)/x \leq 2(\sqrt{x} - 1)/x$. Deduce that $(\ln x)/x \to 0$ as $x \to \infty$. Where is the maximum of $(\ln x)/x$?

60. Prove that $(\ln x)/x^{1/n}$ also approaches zero. (Start with $(\ln x^{1/n})/x^{1/n} \to 0$.) Where is its maximum?
61 For any power $n$, Problem 6.2.59 proved $e^x > x^n$ for large $x$. Then by logarithms, $x > n \ln x$. Since $(\ln x)x$ goes below $\ln x$ and stays below, it converges to ______.

62 Prove that $y \ln y$ approaches zero as $y \to 0$, by changing $y$ to $1/x$. Find the limit of $y^2$ (take its logarithm as $y \to 0$). What is $e^{10}$ on your calculator?

63 Find the limit of $\ln x / \log_{10} x$ as $x \to \infty$.

64 We know the integral $\int_1^x e^{t-1} \, dt = [t^{x-1}]_1 = (x-1) / h$. Its limit as $h \to 0$ is ______.

65 Find linear approximations near $x = 0$ for $e^{-x}$ and $2^x$.

66 The $x^3$ correction to $\ln(1 + x)$ yields $x - \frac{1}{2} x^3 + \frac{1}{3} x^3$. Check that $\ln 1.01$ is 0.0099503 and find $\ln 1.02$.

67 An ant crawls at 1 foot/second along a rubber band whose original length is 2 feet. The band is being stretched at 1 foot/second by pulling the other end. At what time $T$, if ever, does the ant reach the other end?

One approach: The band's length at time $t$ is $t + 2$. Let $y(t)$ be the fraction of that length which the ant has covered, and explain

(a) $y' = 1/(t + 2)$  
(b) $y = \ln(t + 2) - \ln 2$  
(c) $T = 2e - 2$.

68 If the rubber band is stretched at 8 feet/second, when if ever does the same ant reach the other end?

69 A weaker ant slows down to $2/(t + 2)$ feet/second, so $y' = 2/(t + 2)^2$. Show that the other end is never reached.

70 The slope of $p = x^2$ comes two ways from $\ln p = x \ln x$.


2. Exponential differentiation (ED): Write $x^2$ as $e^{a\ln x}$, take its derivative, and put back $x^2$.

71 If $p = 2^x$ then $\ln p = ______$. LD gives $p' = (p)(\ln p)' = ______$. ED gives $p = e^{2^x}$ and then $p' = ______$.

72 Compute $\ln 2$ by the trapezoidal rule and/or Simpson's rule, to get five correct decimals.

73 Compute $\ln 10$ by either rule with $\Delta x = 1$, and compare with the value on your calculator.

74 Estimate $\ln 90,000$, the fraction of numbers near 90,000 that are prime. (879 of the next 10,000 numbers are actually prime.)

75 Find a pair of positive integers for which $x^2 = y^3$. Show how to change this equation to $(\ln x)/(x) = (\ln y)/(y)$. So look for two points at the same height in Figure 6.13. Prove that you have discovered all the integer solutions.

*76 Show that $(\ln x)/x = (\ln y)/y$ is satisfied by

$$x = \left(\frac{t + 1}{t}\right)^{t+1} \quad \text{and} \quad y = \left(\frac{t + 1}{t}\right)^{t+1}$$

with $t \neq 0$. Graph those points to show the curve $x^2 = y^3$. It crosses the line $y = x$ at $x = ______$, where $t \to \infty$.

### 6.5 Separable Equations Including the Logistic Equation

This section begins with the integrals that solve two basic differential equations:

$$\frac{dy}{dt} = cy \quad \text{and} \quad \frac{dy}{dt} = cy + s.$$  \hspace{1cm} (1)

We already know the solutions. What we don't know is how to discover those solutions, when a suggestion "try $e^{ct}$" has not been made. Many important equations, including these, separate into a $y$-integral and a $t$-integral. The answer comes directly from the two separate integrations. When a differential equation is reduced that far— to integrals that we know or can look up—it is solved.

One particular equation will be emphasized. The logistic equation describes the speedup and slowdown of growth. Its solution is an S-curve, which starts slowly, rises quickly, and levels off. (The 1990's are near the middle of the S, if the prediction is correct for the world population.) S-curves are solutions to nonlinear equations, and we will be solving our first nonlinear model. It is highly important in biology and all life sciences.
SEPARABLE EQUATIONS

The equations $\frac{dy}{dt} = cy$ and $\frac{dy}{dt} = cy + s$ (with constant source $s$) can be solved by a direct method. The idea is to separate $y$ from $t$:

$$\frac{dy}{y} = c \, dt \quad \text{and} \quad \frac{dy}{y + (s/c)} = c \, dt.$$  \hspace{1cm} (2)

All $y$'s are on the left side. All $t$'s are on the right side (and $c$ can be on either side). This separation would not be possible for $\frac{dy}{dt} = y + t$.

Equation (2) contains differentials. They suggest integrals. The $t$-integrals give $ct$ and the $y$-integrals give logarithms:

$$\ln y = ct + \text{constant} \quad \text{and} \quad \ln \left( y + \frac{s}{c} \right) = ct + \text{constant}.$$  \hspace{1cm} (3)

The constant is determined by the initial condition. At $t = 0$ we require $y = y_0$, and the right constant will make that happen:

$$\ln y = ct + \ln y_0 \quad \text{and} \quad \ln \left( y + \frac{s}{c} \right) = ct + \ln \left( y_0 + \frac{s}{c} \right).$$  \hspace{1cm} (4)

Then the final step isolates $y$. The goal is a formula for $y$ itself, not its logarithm, so take the exponential of both sides ($e^{\ln y}$ is $y$):

$$y = y_0 e^t \quad \text{and} \quad y + \frac{s}{c} = \left( y_0 + \frac{s}{c} \right) e^t.$$  \hspace{1cm} (5)

It is wise to substitute $y$ back into the differential equation, as a check.

This is our fourth method for $y' = cy + s$. Method 1 assumed from the start that $y = Ae^{rt} + B$. Method 2 multiplied all inputs by their growth factors $e^{rt}$ and added up outputs. Method 3 solved for $y - y_m$. Method 4 is separation of variables (and all methods give the same answer). This separation method is so useful that we repeat its main idea, and then explain it by using it.

To solve $\frac{dy}{dt} = u(y)v(t)$, separate $\frac{dy}{u(y)}$ from $v(t)dt$ and integrate both sides:

$$\int \frac{dy}{u(y)} = \int v(t)dt + C.$$  \hspace{1cm} (6)

Then substitute the initial condition to determine $C$, and solve for $y(t)$.

EXAMPLE 1 $\frac{dy}{dt} = y^2$ separates into $\frac{dy}{y^2} = dt$. Integrate to reach $-1/y = t + C$. Substitute $t = 0$ and $y = y_0$ to find $C = -1/y_0$. Now solve for $y$:

$$-\frac{1}{y} = t - \frac{1}{y_0} \quad \text{and} \quad y = \frac{y_0}{1 - ty_0}.$$  

This solution blows up (Figure 6.15a) when $t$ reaches $1/y_0$. If the bank pays interest on your deposit squared ($y' = y^2$), you soon have all the money in the world.

EXAMPLE 2 $\frac{dy}{dt} = ty$ separates into $dy/y = t \, dt$. Then by integration $\ln y = \frac{1}{2}t^2 + C$. Substitute $t = 0$ and $y = y_0$ to find $C = \ln y_0$. The exponential of $\frac{1}{2}t^2 + \ln y_0$ gives $y = y_0 e^{t^2/2}$. When the interest rate is $c = t$, the exponent is $t^2/2$.

EXAMPLE 3 $\frac{dy}{dt} = y + t$ is not separable. Method 1 survives by assuming $y =$
6.5 Separable Equations Including the Logistic Equation

Separable Equations Including the Logistic Equation

\[ Ae^t + B + Dt \] — with an extra coefficient \( D \) in Problem 23. Method 2 also succeeds—but not the separation method.

**EXAMPLE 4** Separate \( dy/dt = ny/t \) into \( dy/y = n \, dt/t \). By integration \( \ln y = n \ln t + C \). Substituting \( t = 0 \) produces \( \ln 0 \) and disaster. This equation cannot start from time zero (it divides by \( t \)). However \( y \) can start from \( y_1 \) at \( t = 1 \), which gives \( C = \ln y_1 \). The solution is a power function \( y = y_1 t^n \).

This was the first differential equation in the book (Section 2.2). The ratio of \( dy/y \) to \( dt/t \) is the “elasticity” in economics. These relative changes have units like dollars/dollars—they are dimensionless, and \( y = t^n \) has constant elasticity \( n \).

On log–log paper the graph of \( \ln y = n \ln t + C \) is a straight line with slope \( n \).

**THE LOGISTIC EQUATION**

The simplest model of population growth is \( dy/dt = cy \). The growth rate \( c \) is the birth rate minus the death rate. If \( c \) is constant the growth goes on forever—beyond the point where the model is reasonable. A population can’t grow all the way to infinity! Eventually there is competition for food and space, and \( y = e^t \) must slow down.

The true rate \( c \) depends on the population size \( y \). It is a function \( c(y) \) not a constant. The choice of the model is at least half the problem:

- **Problem in biology or ecology:** Discover \( c(y) \).
- **Problem in mathematics:** Solve \( dy/dt = c(y)y \).

Every model looks linear over a small range of \( y \)'s—but not forever. When the rate drops off, two models are of the greatest importance. The *Michaelis–Menten* equation has \( c(y) = c/(y + K) \). The logistic equation has \( c(y) = c - by \). It comes first.

The nonlinear effect is from “interaction.” For two populations of size \( y \) and \( z \), the number of interactions is proportional to \( y \) times \( z \). The Law of Mass Action produces a quadratic term \( byz \). It is the basic model for interactions and competition. Here we have one population competing within itself, so \( z \) is the same as \( y \). This competition slows down the growth, because \(-by^2\) goes into the equation.

The basic model of growth versus competition is known as the **logistic equation**:

\[ dy/dt = cy - by^2. \] (7)

Normally \( b \) is very small compared to \( c \). The growth begins as usual (close to \( e^t \)). The competition term \( by^2 \) is much smaller than \( cy \), until \( y \) itself gets large. Then \( by^2 \)
(with its minus sign) slows the growth down. The solution follows an S-curve that we can compute exactly.

What are the numbers $b$ and $c$ for human population? Ecologists estimate the natural growth rate as $c = 0.029$/year. That is not the actual rate, because of $b$. About 1930, the world population was 3 billion. The $cy$ term predicts a yearly increase of $(0.029)(3 \text{ billion}) = 87$ million. The actual growth was more like $dy/dt = 60$ million/year. That difference of 27 million/year was $by^2$:

$$27 \text{ million/year} = b(3 \text{ billion})^2 \text{ leads to } b = 3 \cdot 10^{-12}/\text{year}.$$

Certainly $b$ is a small number (three trillionths) but its effect is not small. It reduces 87 to 60. What is fascinating is to calculate the *steady state*, when the new term $by^2$ equals the old term $cy$. When these terms cancel each other, $dy/dt = cy - by^2$ is zero. The loss from competition balances the gain from new growth: $cy = by^2$ and $y = c/b$. The growth stops at this equilibrium point—the top of the S-curve:

$$y_e = \frac{c}{b} = \frac{0.029}{3} \approx 10^{12} \approx 10 \text{ billion people}.$$

According to Verhulst's logistic equation, the *world population is converging to 10 billion*. That is from the model. From present indications we are growing much faster. We will very probably go beyond 10 billion. The United Nations report in Section 3.3 predicts 11 billion to 14 billion.

*Notice a special point halfway to $y_e = c/b$. (In the model this point is at 5 billion.) It is the inflection point where the S-curve begins to bend down. The second derivative $d^2y/dt^2$ is zero. The slope $dy/dt$ is a maximum. It is easier to find this point from the differential equation (which gives $dy/dt$) than from $y$. Take one more derivative:*

$$y'' = (cy - by^2)' = cy' - 2byy' = (c - 2by)y'. \quad (8)$$

The factor $c - 2by$ is zero at the inflection point $y = c/2b$, halfway up the S-curve.

**THE S-CURVE**

*The logistic equation is solved by separating variables $y$ and $t$:*

$$dy/dt = cy - by^2 \text{ becomes } \int dy/(cy - by^2) = \int dt. \quad (9)$$

The first question is whether we recognize this $y$-integral. No. The second question is whether it is listed in the cover of the book. No. The nearest is $\int dx/(a^2 - x^2)$, which can be reached with considerable manipulation (Problem 21). The third question is whether a general method is available. Yes. "Partial fractions" is perfectly suited to $1/(cy - by^2)$, and Section 7.4 gives the following integral of equation (9):

$$\ln \frac{y}{c - by} = ct + C \quad \text{and then} \quad \ln \frac{y_0}{c - by_0} = C. \quad (10)$$

That constant $C$ makes the solution correct at $t = 0$. The logistic equation is integrated, but the solution can be improved. Take exponentials of both sides to remove the logarithms:

$$\frac{y}{c - by} = e^t \frac{y_0}{c - by_0}. \quad (11)$$

This contains the same growth factor $e^t$ as in linear equations. But the logistic
equation is not linear—it is not $y$ that increases so fast. According to (11), it is $y/(c - by)$ that grows to infinity. This happens when $c - by$ approaches zero.

The growth stops at $y = c/b$. That is the final population of the world (10 billion?). We still need a formula for $y$. The perfect S-curve is the graph of $y = 1/(1 + e^{-t})$. It equals 1 when $t = \infty$, it equals $\frac{1}{2}$ when $t = 0$, it equals 0 when $t = -\infty$. It satisfies $y' = y - y^2$, with $c = b = 1$. The general formula cannot be so beautiful, because it allows any $c$, $b$, and $y_0$. To find the S-curve, multiply equation (11) by $c - by$ and solve for $y$:

$$
y = \frac{c}{b + e^{-\alpha(c - by_0)/y_0}} \text{ or } y = \frac{c}{b + de^{-\alpha t}}.
$$

When $t$ approaches infinity, $e^{-\alpha t}$ approaches zero. The complicated part of the formula disappears. Then $y$ approaches its steady state $c/b$, the asymptote in Figure 6.16. The S-shape comes from the inflection point halfway up.

Surprising observation: $z = 1/y$ satisfies a linear equation. By calculus $z' = -y'/y^2$. So

$$
z' = \frac{-cy + by^2}{y^2} = -\frac{c}{y} + b = -cz + b.
$$

This equation $z' = -cz + b$ is solved by an exponential $e^{-\alpha t}$ plus a constant:

$$
z = Ae^{-\alpha t} + \frac{b}{c} = \left(\frac{1}{y_0} - \frac{b}{c}\right)e^{-\alpha t} + \frac{b}{c}.
$$

Turned upside down, $y = 1/z$ is the S-curve (12). As $z$ approaches $b/c$, the S-curve approaches $c/b$. Notice that $z$ starts at $1/y_0$.

**EXAMPLE 1 (United States population)** The table shows the actual population and the model. Pearl and Reed used census figures for 1790, 1850, and 1910 to compute $c$ and $b$. In between, the fit is good but not fantastic. One reason is war—another is depression. Probably more important is immigration.† In fact the Pearl-Reed steady state $c/b$ is below 200 million, which the US has already passed. Certainly their model can be and has been improved. The 1990 census predicted a stop before 300 million.

For constant immigration $s$ we could still solve $y' = cy - by^2 + s$ by partial fractions—but in practice the computer has taken over. The table comes from Braun's book *Differential Equations* (Springer 1975).

† Immigration does not enter for the world population model (at least not yet).
Remark. For good science the $y^2$ term should be explained and justified. It gave a nonlinear model that could be completely solved, but simplicity is not necessarily truth. The basic justification is this: In a population of size $y$, the number of encounters is proportional to $y^2$. If those encounters are fights, the term is $-by^2$. If those encounters increase the population, as some like to think, the sign is changed. There is a cooperation term $+by^2$, and the population increases very fast.

**EXAMPLE 5**

$y' = cy + by^2$: $y$ goes to infinity in a finite time.

**EXAMPLE 6**

$y' = -dy + by^2$: $y$ dies to zero if $y_0 < d/b$.

In Example 6 death wins. A small population dies out before the cooperation $by^2$ can save it. A population below $d/b$ is an endangered species.

The logistic equation can't predict oscillations—those go beyond $dy/dt = f(y)$.

The $y$ line. Here is a way to understand every nonlinear equation $y' = f(y)$. Draw a "$y$ line." Add arrows to show the sign of $f(y)$. When $y' = f(y)$ is positive, $y$ is increasing (it follows the arrow to the right). When $f$ is negative, $y$ goes to the left. When $f$ is zero, the equation is $y' = 0$ and $y$ is stationary:

The arrows take you left or right, to the steady state or to infinity. Arrows go toward stable steady states. The arrows go away, when the stationary point is unstable. The $y$ line shows which way $y$ moves and where it stops.

The terminal velocity of a falling body is $v_w = \sqrt{g}$ in Problem 6.7.54. For $f(y) = \sin y$ there are several steady states:

**EXAMPLE 7** Kinetics of a chemical reaction $mA + nB \rightarrow pC$.

The reaction combines $m$ molecules of $A$ with $n$ molecules of $B$ to produce $p$ molecules of $C$. The numbers $m, n, p$ are 1, 1, 2 for hydrogen chloride: $\text{H}_2 + \text{Cl}_2 = 2 \text{HCl}$. The Law of Mass Action says that the reaction rate is proportional to the product of the concentrations $[A]$ and $[B]$. Then $[A]$ decays as $[C]$ grows:

$$d[A]/dt = -r[A][B] \quad \text{and} \quad d[C]/dt = +k[A][B]. \quad (15)$$

Chemistry measures $r$ and $k$. Mathematics solves for $[A]$ and $[C]$. Write $y$ for the concentration $[C]$, the number of molecules in a unit volume. Forming those $y$ molecules drops the concentration $[A]$ from $a_0$ to $a_0 - (m/p)y$. Similarly $[B]$ drops from $b_0$ to $b_0 - (n/p)y$. The mass action law (15) contains $y^2$:

$$\frac{dy}{dt} = k\left(a_0 - \frac{m}{p}y\right)\left(b_0 - \frac{n}{p}y\right). \quad (16)$$
This fits our nonlinear model (Problem 33–34). We now find this same mass action in biology. You recognize it whenever there is a product of two concentrations.

**THE MM EQUATION** \( \frac{dy}{dt} = -cy(y + K) \)

Biochemical reactions are the keys to life. They take place continually in every living organism. Their mathematical description is not easy! Engineering and physics go far with linear models, while biology is quickly nonlinear. It is true that \( y' = cy \) is extremely effective in first-order kinetics (Section 6.3), but nature builds in a nonlinear regulator.

It is enzymes that speed up a reaction. Without them, your life would be in slow motion. Blood would take years to clot. Steaks would take decades to digest. Calculus would take centuries to learn. The whole system is awesomely beautiful—DNA tells amino acids how to combine into useful proteins, and we get enzymes and elephants and Isaac Newton.

Briefly, the enzyme enters the reaction and comes out again. It is the catalyst. Its combination with the substrate is an unstable intermediate, which breaks up into a new product and the enzyme (which is ready to start over).

Here are examples of catalysts, some good and some bad.

1. The platinum in a catalytic converter reacts with pollutants from the car engine. (But platinum also reacts with lead—ten gallons of leaded gasoline and you can forget the platinum.)
2. Spray propellants (CFC's) catalyze the change from ozone \((O_3)\) into ordinary oxygen \((O_2)\). This wipes out the ozone layer—our shield in the atmosphere.
3. Milk becomes yoghurt and grape juice becomes wine.
4. Blood clotting needs a whole cascade of enzymes, amplifying the reaction at every step. In hemophilia—the “Czar's disease”—the enzyme called Factor VIII is missing. A small accident is disaster; the bleeding won't stop.
5. Adolph's Meat Tenderizer is a protein from papayas. It predigests the steak.
6. The same enzyme (chymopapain) is injected to soften herniated disks.
7. Yeast makes bread rise. Enzymes put the sour in sourdough.

Of course, it takes enzymes to make enzymes. The maternal egg contains the material for a cell, and also half of the DNA. The fertilized egg contains the full instructions.

We now look at the Michaelis–Menten (MM) equation, to describe these reactions. It is based on the Law of Mass Action. An enzyme in concentration \( z \) converts a substrate in concentration \( y \) by \( \frac{dy}{dt} = -byz \). The rate constant is \( b \), and you see the product of “enzyme times substrate.” A similar law governs the other reactions (some go backwards). The equations are nonlinear, with no exact solution. It is typical of applied mathematics (and nature) that a pattern can still be found.

What happens is that the enzyme concentration \( z(t) \) quickly drops to \( z_0K/(y + K) \). The Michaelis constant \( K \) depends on the rates (like \( b \)) in the mass action laws. Later the enzyme reappears \((z_\infty = z_0)\). But by then the first reaction is over. Its law of mass action is effectively

\[
\frac{dy}{dt} = -byz = -\frac{cy}{y + K}
\]

with \( c = bz_0K \). This is the **Michaelis–Menten equation**—basic to biochemistry.

The rate \( dy/dt \) is all-important in biology. Look at the function \( cy/(y + K) \):

when \( y \) is large, \( dy/dt \approx -c \)  
when \( y \) is small, \( dy/dt \approx -cy/K \).
The start and the finish operate at different rates, depending whether \( y \) dominates \( K \) or \( K \) dominates \( y \). The fastest rate is \( c \).

A biochemist solves the MM equation by separating variables:

\[
\int \frac{y + K}{y} \, dy = - \int c \, dt \quad \text{gives} \quad y + K \ln y = -ct + C. \tag{18}
\]

Set \( t = 0 \) as usual. Then \( C = y_0 + K \ln y_0 \). The exponentials of the two sides are

\[
e^{y}y^{K} = e^{-ct}e^{y_0y^{K}}. \tag{19}
\]

We don't have a simple formula for \( y \). We are lucky to get this close. A computer can quickly graph \( y(t) \)—and we see the dynamics of enzymes.

Problems 27–32 follow up the Michaelis–Menten theory. In science, concentrations and rate constants come with units. In mathematics, variables can be made dimensionless and constants become 1. We solve \( dy/dT = Y/(Y + 1) \) and then switch back to \( y, t, c, K \). This idea applies to other equations too.

**Essential point:** Most applications of calculus come through differential equations. That is the language of mathematics—with populations and chemicals and epidemics obeying the same equation. Running parallel to \( dy/dt = cy \) are the difference equations that come next.

### 6.5 EXERCISES

**Read-through questions**

The equations \( dy/dt = cy \) and \( dy/dt = cy + s \) and \( dy/dt = u(y)s(t) \) are called _____ because we can separate \( y \) from \( t \).

Integration of \( \int dy/y = \int c \, dt \) gives ____. Integration of \( \int dy/(y + s/c) = \int c \, dt \) gives ____. The equation \( dy/dx = -x/y \) leads to _____. Then \( y^2 + x^2 = \) ____ and the solution stays on a circle.

The logistic equation is \( dy/dt = ____ \). The new term \(-by^2\) represents _____ when \( cy \) represents growth. Separation gives \( \int dy/(cy - by^2) = \int dt \), and the \( y \)-integral is \( 1/c \) times \( \ln ____ \). Substituting \( y_0 \) at \( t = 0 \) and taking exponentials produces \( y/c - byy = e^{t(____)} \). As \( t \to \infty \), \( y \) approaches _____. That is the steady state where \( cy - by^2 = ____ \). The graph of \( y \) looks like an _____ because it has an inflection point at \( y = ____ \).

In biology and chemistry, concentrations \( y \) and \( z \) react at a rate proportional to \( y \) times _____. This is the Law of _____. In a model equation \( dy/dt = c(y) \), the rate \( c \) depends on _____. The MM equation is \( dy/dt = ____ \). Separating variables yields \( \int ____\) \( dy = ____\) \( = -ct + C \).

Separate, integrate, and solve equations 1–8.

1. \( dy/dt = y + 5 \), \( y_0 = 2 \)
2. \( dy/dt = 1/y \), \( y_0 = 1 \)
3. \( dy/dx = x/y^2 \), \( y_0 = 1 \)
4. \( dy/dx = y^2 + 1 \), \( y_0 = 0 \)
5. \( dy/dx = (y + 1)/(x + 1) \), \( y_0 = 0 \)
6. \( dy/dx = \tan y \cos x \), \( y_0 = 1 \)
7. \( dy/dt = y \sin t \), \( y_0 = 1 \)
8. \( dy/dt = e^{-ct} \), \( y_0 = e \)
9. Suppose the rate of growth is proportional to \( \sqrt{y} \) instead of \( y \). Solve \( dy/dt = cy \sqrt{y} \) starting from \( y_0 \).
10. The equation \( dy/dx = ny/x \) for constant elasticity is the same as \( d(ln y)/(d(ln x)) = ____ \). The solution is \( ln y = ____ \).
11. When \( c = 0 \) in the logistic equation, the only term is \( y' = -by^2 \). What is the steady state \( y_s \)? How long until \( y \) drops from \( y_0 \) to \( \frac{1}{2} y_0 \)?
12. Reversing signs in Problem 11, suppose \( y' = +by^2 \). At what time does the population explode to \( y = \infty \), starting from \( y_0 = 2 \) (Adam + Eve)?

Problems 13–26 deal with logistic equations \( y' = cy - by^2 \).

13. Show that \( y = 1/(1 + e^{-t}) \) solves the equation \( y' = y - y^2 \). Draw the graph of \( y \) from starting values \( \frac{1}{2} \) and \( \frac{3}{2} \).
14. (a) What logistic equation is solved by \( y = 2(1 - e^{-t}) \)?
   (b) Find \( c \) and \( b \) in the equation solved by \( y = 1/(1 + e^{-3t}) \).
15. Solve \( z' = -z + 1 \) with \( z_0 = 2 \). Turned upside down as in (13), what is \( y = 1/z \)?
6.6 Powers Instead of Exponentials

You may remember our first look at \( e \). It is the special base for which \( e^x \) has slope 1 at \( x = 0 \). That led to the great equation of exponential growth: The derivative of \( e^x \) equals \( e^x \). But our look at the actual number \( e = 2.71828... \) was very short.
It appeared as the limit of \((1 + 1/n)^n\). This seems an unnatural way to write down such an important number.

I want to show how \((1 + 1/n)^n\) and \((1 + x/n)^n\) arise naturally. They give **discrete growth in finite steps**—with applications to compound interest. Loans and life insurance and money market funds use the discrete form of \(y' = cy + s\). (We include extra information about bank rates, hoping this may be useful some day.) The applications in science and engineering are equally important. Scientific computing, like accounting, has difference equations in parallel with differential equations.

Knowing that this section will be full of formulas, I would like to jump ahead and tell you the best one. It is an infinite series for \(e^x\). What makes the series beautiful is that its derivative is itself.

Start with \(y = 1 + x\). This has \(y = 1\) and \(y' = 1\) at \(x = 0\). But \(y''\) is zero, not one. Such a simple function doesn't stand a chance! No polynomial can be its own derivative, because the highest power \(x^n\) drops down to \(nx^{n-1}\). The only way is to have no highest power. We are forced to consider infinitely many terms—a power series—to achieve "derivative equals function."

To produce the derivative \(1 + x\), we need \(1 + x + \frac{1}{2}x^2\). Then \(\frac{1}{2}x^2\) is the derivative of \(\frac{1}{2}x^3\), which is the derivative of \(\frac{1}{2}x^4\). The best way is to write the whole series at once:

\[\text{Infinite series} \quad e^x = 1 + x + \frac{1}{2}x^2 + \frac{1}{6}x^3 + \frac{1}{24}x^4 + \cdots. \tag{1}\]

This must be the greatest power series ever discovered. Its derivative is itself:

\[\frac{de^x}{dx} = 0 + 1 + x + \frac{1}{2}x^2 + \frac{1}{6}x^3 + \cdots = e^x. \tag{2}\]

The derivative of each term is the term before it. The integral of each term is the one after it (so \(\int e^x\,dx = e^x + C\)). The approximation \(e^x \approx 1 + x\) appears in the first two terms. Other properties like \((e^x)'(e^x) = e^{2x}\) are not so obvious. (Multiplying series is hard but interesting.) It is not even clear why the sum is \(2.718\ldots\) when \(x = 1\).

Notice that \(x^n\) is divided by the product \(1 \cdot 2 \cdot 3 \cdots n\). This is "n factorial." Thus \(x^4\) is divided by \(1 \cdot 2 \cdot 3 \cdot 4 = 4! = 24\), and \(x^5\) is divided by \(5! = 120\). The derivative of \(x^5/120\) is \(x^4/24\), because 5 from the derivative cancels 5 from the factorial. In general \(x^n/n!\) has derivative \(x^{n-1}/(n-1)!\). Surprisingly 0! is 1.

Chapter 10 emphasizes that \(x^n/n!\) becomes extremely small as \(n\) increases. The infinite series adds up to a finite number—which is \(e^x\). We turn now to discrete growth, which produces the same series in the limit.

This headline was on page one of the New York Times for May 27, 1990.

**213 Years After Loan, Uncle Sam is Dunned**

San Antonio, May 26—More than 200 years ago, a wealthy Pennsylvania merchant named Jacob DeHaven lent $450,000 to the Continental Congress to rescue the troops at Valley Forge. That loan was apparently never repaid.

So Mr. DeHaven's descendants are taking the United States Government to court to collect what they believe they are owed. The total: $141 billion if the interest is compounded daily at 6 percent, the going rate at the time. If compounded yearly, the bill is only $98 billion.

The thousands of family members scattered around the country say they are not being greedy. "It's not the money—it's the principle of the thing," said Carolyn Cokerham, a DeHaven on her father's side who lives in San Antonio.
"You have to wonder whether there would even be a United States if this man had not made the sacrifice that he did. He gave everything he had."

The descendants say that they are willing to be flexible about the amount of settlement. But they also note that interest is accumulating at $190 a second.

"None of these people have any intention of bankrupting the Government," said Jo Beth Kloecker, a lawyer from Stafford, Texas. Fresh out of law school, Ms. Kloecker accepted the case for less than the customary 30 percent contingency.

It is unclear how many descendants there are. Ms. Kloecker estimates that based on 10 generations with four children in each generation, there could be as many as half a million.

The initial suit was dismissed on the ground that the statute of limitations is six years for a suit against the Federal Government. The family's appeal asserts that this violates Article 6 of the Constitution, which declares as valid all debts owed by the Government before the Constitution was adopted.

Mr. DeHaven died penniless in 1812. He had no children.

**COMPOUND INTEREST**

The idea of compound interest can be applied right away. Suppose you invest $1000 at a rate of 100% (hard to do). If this is the annual rate, the interest after a year is another $1000. You receive $2000 in all. But if the interest is compounded you receive more:

after six months: Interest of $500 is reinvested to give $1500

end of year: New interest of $750 (50% of 1500) gives $2250 total.

The bank multiplied twice by 1.5 (1000 to 1500 to 2250). Compounding quarterly multiplies four times by 1.25 (1 for principal, .25 for interest):

after one quarter the total is $1000 + (.25)(1000) = 1250$

after two quarters the total is $1250 + (.25)(1250) = 1562.50$

after nine months the total is $1562.50 + (.25)(1562.50) = 1953.12$

after a full year the total is $1953.12 + (.25)(1953.12) = 2441.41$

Each step multiplies by $1 + (1/n)$, to add one $n$th of a year's interest—still at 100%:

quarterly conversion: $(1 + 1/4)^4 \times 1000 = 2441.41$

monthly conversion: $(1 + 1/12)^{12} \times 1000 = 2613.04$

daily conversion: $(1 + 1/365)^{365} \times 1000 = 2714.57$.

Many banks use 360 days in a year, although computers have made that obsolete. Very few banks use minutes (525,600 per year). Nobody compounds every second ($n = 31,536,000$). But some banks offer *continuous compounding*. This is the limiting case ($n \to \infty$) that produces $e$:

$$\left(1 + \frac{1}{n}\right)^n \times 1000 \text{ approaches } e \times 1000 = 2718.28.$$

1. **Quick method for $(1 + 1/n)^n$**: *Take its logarithm. Use* $\ln(1 + x) \approx x$ *with* $x = \frac{1}{n}$;
6 Exponentials and Logarithms

\[ \ln \left( 1 + \frac{1}{n} \right)^n = n \ln \left( 1 + \frac{1}{n} \right) \approx n \left( \frac{1}{n} \right) = 1. \] (3)

As \(1/n\) gets smaller, this approximation gets better. The limit is 1. Conclusion: \((1 + 1/n)^n\) approaches the number whose logarithm is 1. Sections 6.2 and 6.4 define the same number (which is \(e\)).

2. Slow method for \((1 + 1/n)^n\): Multiply out all the terms. Then let \(n \to \infty\).

This is a brutal use of the binomial theorem. It involves nothing smart like logarithms, but the result is a fantastic new formula for \(e\).

Practice for \(n = 3\):
\[(1 + \frac{1}{3})^3 = 1 + 3 \left( \frac{1}{3} \right) + \frac{3 \cdot 2 \cdot 1}{1 \cdot 2 \cdot 3} \left( \frac{1}{3} \right)^2 + \frac{3 \cdot 2 \cdot 1}{1 \cdot 2 \cdot 3} \left( \frac{1}{3} \right)^3. \]

Binomial theorem for any positive integer \(n\):
\[(1 + \frac{1}{n})^n = 1 + n \left( \frac{1}{n} \right) + \frac{n(n-1)}{1 \cdot 2} \left( \frac{1}{n} \right)^2 + \frac{n(n-1)(n-2)}{1 \cdot 2 \cdot 3} \left( \frac{1}{n} \right)^3 + \ldots + \left( \frac{1}{n} \right)^n. \] (4)

Each term in equation (4) approaches a limit as \(n \to \infty\). Typical terms are
\[ \frac{n(n-1)}{1 \cdot 2} \left( \frac{1}{n} \right)^2 \to \frac{1}{1 \cdot 2} \quad \text{and} \quad \frac{n(n-1)(n-2)}{1 \cdot 2 \cdot 3} \left( \frac{1}{n} \right)^3 \to \frac{1}{1 \cdot 2 \cdot 3}. \]

Next comes \(1/1 \cdot 2 \cdot 3 \cdot 4\). The sum of all those limits in (4) is our new formula for \(e\):
\[ \lim \left( 1 + \frac{1}{n} \right)^n = 1 + 1 + \frac{1}{1 \cdot 2} + \frac{1}{1 \cdot 2 \cdot 3} + \frac{1}{1 \cdot 2 \cdot 3 \cdot 4} + \ldots = e. \] (5)

In summation notation this is \(\sum_{k=0}^{\infty} 1/k! = e\). The factorials give fast convergence:
\[ 1 + 1 + .5 + .16667 + .04167 + .00833 + .00139 + .00020 + .00002 = 2.71828. \]

Those nine terms give an accuracy that was not reached by \(n = 365\) compoundings. A limit is still involved (to add up the whole series). You never see \(e\) without a limit! It can be defined by derivatives or integrals or powers \((1 + 1/n)^n\) or by an infinite series. Something goes to zero or infinity, and care is required.

All terms in equation (4) are below (or equal to) the corresponding terms in (5). The power \((1 + 1/n)^n\) approaches \(e\) from below. There is a steady increase with \(n\). Faster compounding yields more interest. Continuous compounding at 100% yields \(e\), as each term in (4) moves up to its limit in (5).

Remark Change \((1 + 1/n)^n\) to \((1 + x/n)^n\). Now the binomial theorem produces \(e^x\):
\[ \left( 1 + \frac{x}{n} \right)^n = 1 + n \left( \frac{x}{n} \right) + \frac{n(n-1)}{1 \cdot 2} \left( \frac{x}{n} \right)^2 + \ldots \text{ approaches } 1 + x + \frac{x^2}{1 \cdot 2} + \ldots. \] (6)

Please recognize \(e^x\) on the right side! It is the infinite power series in equation (1). The next term is \(x^3/6\) (\(x\) can be positive or negative). This is a final formula for \(e^x\):

\[ \text{6. The limit of } (1 + x/n)^n \text{ is } e^x. \] At \(x = 1\) we find \(e\).

The logarithm of that power is \(n \ln(1 + x/n) \approx n(x/n) = x\). The power approaches \(e^x\).

To summarize: The quick method proves \((1 + 1/n)^n \to e\) by logarithms. The slow method (multiplying out every term) led to the infinite series. Together they show the agreement of all our definitions of \(e\).
DIFFERENCE EQUATIONS VS. DIFFERENTIAL EQUATIONS

We have the chance to see an important part of applied mathematics. This is not a course on differential equations, and it cannot become a course on difference equations. But it is a course with a purpose—we aim to use what we know. Our main application of $e$ was to solve $y' = cy$ and $y' = cy + s$. Now we solve the corresponding difference equations.

Above all, the goal is to see the connections. The purpose of mathematics is to understand and explain patterns. The path from “discrete to continuous” is beautifully illustrated by these equations. Not every class will pursue them to the end, but I cannot fail to show the pattern in a difference equation:

$$y(t + 1) = ay(t).$$

(7)

Each step multiplies by the same number $a$. The starting value $y_0$ is followed by $ay_0$, $a^2y_0$, and $a^3y_0$. The solution at discrete times $t = 0, 1, 2, \ldots$ is $y(t) = a^t y_0$.

This formula $a^t y_0$ replaces the continuous solution $e^{ct} y_0$ of the differential equation.

![Fig. 6.17 Growth for $|a| > 1$, decay for $|a| < 1$. Growth factor $a$ compares to $e^c$.](image)

A source or sink (birth or death, deposit or withdrawal) is like $y' = cy + s$:

$$y(t + 1) = ay(t) + s.$$  

(8)

Each step multiplies by $a$ and adds $s$. The first outputs are

$$y(1) = ay_0 + s, \quad y(2) = a^2y_0 + as + s, \quad y(3) = a^3y_0 + a^2s + as + s.$$

We saw this pattern for differential equations—every input $s$ becomes a new starting point. It is multiplied by powers of $a$. Since $s$ enters later than $y_0$, the powers stop at $t - 1$. Algebra turns the sum into a clean formula by adding the geometric series:

$$y(t) = a^t y_0 + s \left[ a^{t-1} + a^{t-2} + \cdots + a + 1 \right] = a^t y_0 + s(a^t - 1)/(a-1).$$

(9)

EXAMPLE 1  Interest at 8% from annual IRA deposits of $s = $2000 (here $y_0 = 0$).

The first deposit is at year $t = 1$. In a year it is multiplied by $a = 1.08$, because 8% is added. At the same time a new $s = 2000$ goes in. At $t = 3$ the first deposit has been multiplied by $(1.08)^2$, the second by 1.08, and there is another $s = 2000$. After year $t$,

$$y(t) = 2000(1.08^2 - 1)/(1.08 - 1).$$

(10)

With $t = 1$ this is 2000. With $t = 2$ it is 2000 $(1.08 + 1)$—two deposits. Notice how $a - 1$ (the interest rate .08) appears in the denominator.

EXAMPLE 2  Approach to steady state when $|a| < 1$. Compare with $c < 0$.

With $a > 1$, everything has been increasing. That corresponds to $c > 0$ in the differential equation (which is growth). But things die, and money is spent, so $a$ can be smaller than one. In that case $a^t y_0$ approaches zero—the starting balance disappears. What happens if there is also a source? Every year half of the balance $y(t)$ is
spent and a new $2000 is deposited. Now $a = \frac{1}{2}$:

\[ y(t + 1) = \frac{1}{2} y(t) + 2000 \Rightarrow y(t) = \left(\frac{1}{2}\right)^t y_0 + 2000 \left[ \left(\frac{1}{2}\right)^t - 1 \right] \left(\frac{1}{2} - 1\right) \]

The limit as $t \to \infty$ is an equilibrium point. As $\left(\frac{1}{2}\right)^t$ goes to zero, $y(t)$ stabilizes to

\[ y_\infty = 2000(0 - 1)/(\frac{1}{2} - 1) = 4000 \text{ steady state.} \quad (11) \]

Why is 4000 steady? Because half is lost and the new 2000 makes it up again. The iteration is $y_{n+1} = \frac{1}{2} y_n + 2000$. Its fixed point is where $y_\infty = \frac{1}{2} y_\infty + 2000$.

In general the steady equation is $y_\infty = ay_\infty + s$. Solving for $y_\infty$ gives $s/(1 - a)$.

\[ y_\infty = \frac{s}{1 - a} \text{(differential equation)} \quad \text{vs.} \quad y_\infty = \frac{s}{1 - a} \text{(difference equation).} \quad (12) \]

**EXAMPLE 3** Demand equals supply when the price is right.

Difference equations are basic to economics. Decisions are made every year (by a farmer) or every day (by a bank) or every minute (by the stock market). There are three assumptions:

1. Supply next time depends on price this time: $S(t + 1) = cP(t)$.
2. Demand next time depends on price next time: $D(t + 1) = -dP(t + 1) + b$.
3. Demand next time equals supply next time: $D(t + 1) = S(t + 1)$.

Comment on 3: the price sets itself to make demand = supply. The demand slope $-d$ is negative. The supply slope $c$ is positive. Those lines intersect at the competitive price, where supply equals demand. To find the difference equation, substitute 1 and 2 into 3:

Difference equation: $-dP(t + 1) + b = cP(t)$

Steady state price: $-dP_\infty + b = cP_\infty$. Thus $P_\infty = b/(c + d)$.

If the price starts above $P_\infty$, the difference equation brings it down. If below, the price goes up. When the price is $P_\infty$, it stays there. This is not news-economic theory depends on approach to a steady state. But convergence only occurs if $c < d$.

If supply is less sensitive than demand, the economy is stable.

Blow-up example: $c = 2$, $b = d = 1$. The difference equation is $-P(t + 1) + 1 = 2P(t)$. From $P(0) = 1$ the price oscillates as it grows: $P = 1, 3, -5, 11, \ldots$.

Stable example: $c = 1/2$, $b = d = 1$. The price moves from $P(0) = 1$ to $P(\infty) = 2/3$:

$-P(t + 1) + 1 = \frac{1}{2} P(t) \Rightarrow P = 1, \frac{3}{2}, \frac{5}{3}, \ldots$, approaching $\frac{2}{3}$.

Increasing $d$ gives greater stability. That is the effect of price supports. For $d = 0$ (fixed demand regardless of price) the economy is out of control.

**THE MATHEMATICS OF FINANCE**

It would be a pleasure to make this supply-demand model more realistic—with curves, not straight lines. Stability depends on the slope—calculus enters. But we also have to be realistic about class time. I believe the most practical application is to solve the fundamental problems of finance. Section 6.3 answered six questions about continuous interest. We now answer the same six questions when the annual rate is $x = .05 = 5\%$ and interest is compounded $n$ times a year.
6.6 Powers Instead of Exponentials

First we compute effective rates, higher than .05 because of compounding:

\[
\text{compounded quarterly } \left(1 + \frac{.05}{4}\right) = 1.0509 \quad \text{[effective rate .0509 = 5.09%]}
\]

\[
\text{compounded continuously } e^{.05} = 1.0513 \quad \text{[effective rate 5.13%]}
\]

Now come the six questions. Next to the new answer (discrete) we write the old answer (continuous). One is algebra, the other is calculus. The time period is 20 years, so simple interest on \( y_0 \) would produce \((.05)(20)(y_0)\). That equals \( y_0 \) — money doubles in 20 years at 5% simple interest.

Questions 1 and 2 ask for the future value \( y \) and present value \( y_0 \) with compound interest \( n \) times a year:

1. \( y \) growing from \( y_0 \):
   \[
y = \left(1 + \frac{.05}{n}\right)^{20n} y_0 \quad y = e^{(.05)(20)} y_0
\]

2. Deposit \( y_0 \) to reach \( y \):
   \[
y_0 = \left(1 + \frac{.05}{n}\right)^{-20n} y \quad y_0 = e^{-.05)(20)} y
\]

Each step multiplies by \( a = (1 + .05/n) \). There are \( 20n \) steps in 20 years. Time goes backward in Question 2. We divide by the growth factor instead of multiplying. The future value is greater than the present value (unless the interest rate is negative!). As \( n \to \infty \) the discrete \( y \) on the left approaches the continuous \( y \) on the right.

Questions 3 and 4 connect \( y \) to \( s \) (with \( y_0 = 0 \) at the start). As soon as each \( s \) is deposited, it starts growing. Then \( y = s + as + a^2s + \cdots \).

3. \( y \) growing from deposits \( s \):
   \[
y = s \left[\left(1 + \frac{.05}{n}\right)^{20n} - 1\right] \quad y = s \left[\frac{e^{(.05)(20)} - 1}{05/} \right]
\]

4. Deposits \( s \) to reach \( y \):
   \[
s = y \left[\frac{.05}{(1 + .05/n)^{20n} - 1}\right] \quad s = y \left[\frac{.05}{e^{.05)(20)} - 1}\right]
\]

Questions 5 and 6 connect \( y_0 \) to \( s \). This time \( y \) is zero — there is nothing left at the end. Everything is paid. The deposit \( y_0 \) is just enough to allow payments of \( s \). This is an annuity, where the bank earns interest on your \( y_0 \) while it pays you \( s \) (\( n \) times a year for 20 years). So your deposit in Question 5 is less than \( 20ns \).

Question 6 is the opposite—a loan. At the start you borrow \( y_0 \) (instead of giving the bank \( y_0 \)). You can earn interest on it as you pay it back. Therefore your payments have to total more than \( y_0 \). This is the calculation for car loans and mortgages.

5. Annuity: Deposit \( y_0 \) to receive \( 20n \) payments of \( s \):
   \[
y_0 = s \left[\frac{1 - (1 + .05/n)^{-20n}}{.05/n}\right] \quad y_0 = s \left[\frac{1 - e^{-.05)(20)}{05/}\right]
\]

6. Loan: Repay \( y_0 \) with \( 20n \) payments of \( s \):
   \[
s = y_0 \left[\frac{.05/n}{1 - (1 + .05/n)^{-20n}}\right] \quad s = y_0 \left[\frac{.05}{1 - e^{-.05)(20)}}\right]
\]

Questions 2, 4, 6 are the inverses of 1, 3, 5. Notice the pattern: There are three numbers \( y \), \( y_0 \), and \( s \). One of them is zero each time. If all three are present, go back to equation (9).

The algebra for these lines is in the exercises. It is not calculus because \( \Delta t \) is not \( dt \). All factors in brackets \( [\quad] \) are listed in tables, and the banks keep copies. It might
also be helpful to know their symbols. If a bank has interest rate $i$ per period over $N$ periods, then in our notation $a = 1 + i = 1 + .05/n$ and $t = N = 20n$:

- future value of $y_0 = \$1$ (line 1): $y(N) = (1 + i)^N$
- present value of $y = \$1$ (line 2): $y_0 = (1 + i)^{-N}$
- future value of $s = \$1$ (line 3): $y(N) = s_{\text{r}} = \left[ (1 + i)^N - 1 \right]/i$
- present value of $s = \$1$ (line 5): $y_0 = a_{\text{r}} = \left[ 1 - (1 + i)^{-N} \right]/i$

To tell the truth, I never knew the last two formulas until writing this book. The mortgage on my home has $N = (12)(25)$ monthly payments with interest rate $i = .07/12$. In 1972 the present value was $\$42,000 = amount borrowed. I am now going to see if the bank is honest.†

**Remark** In many loans, the bank computes interest on the amount paid back instead of the amount received. This is called discounting. A loan of $\$1000$ at $5\%$ for one year costs $\$50$ interest. Normally you receive $\$1000$ and pay back $\$1050$. **With discounting you receive** $\$950$ (called the proceeds) **and you pay back** $\$1000$. The true interest rate is higher than $5\%$—because the $\$50$ interest is paid on the smaller amount $\$950$. In this case the "discount rate" is $50/950 = 5.26\%$.

### SCIENTIFIC COMPUTING: DIFFERENTIAL EQUATIONS BY DIFFERENCE EQUATIONS

In biology and business, most events are discrete. In engineering and physics, time and space are continuous. Maybe at some quantum level it's all the same, but the equations of physics (starting with Newton's law $F = ma$) are differential equations. The great contribution of calculus is to model the rates of change we see in nature. But to solve that model with a computer, it needs to be made digital and discrete.

These paragraphs work with $dy/dt = cy$. It is the test equation that all analysts use, as soon as a new computing method is proposed. Its solution is $y = e^{ct}$, starting from $y_0 = 1$. Here we test Euler's method (nearly ancient, and not well thought of). He replaced $dy/dt$ by $\Delta y/\Delta t$:

$$\text{Euler's Method} \quad \frac{y(t + \Delta t) - y(t)}{\Delta t} = cy(t). \tag{13}$$

The left side is $dy/dt$, in the limit $\Delta t \to 0$. We stop earlier, when $\Delta t > 0$.

The problem is to solve (13). Multiplying by $\Delta t$, the equation is

$$y(t + \Delta t) = (1 + c\Delta t)y(t) \quad (\text{with } y(0) = 1).$$

Each step multiplies by $a = 1 + c\Delta t$, so $n$ steps multiply by $a^n$:

$$y = a^n = (1 + c\Delta t)^n \text{ at time } n\Delta t. \tag{14}$$

This is growth or decay, depending on $a$. The correct $e^t$ is growth or decay, depending on $c$. **The question is whether $a^n$ and $e^{nt}$ stay close.** Can one of them grow while the other decays? We expect the difference equation to copy $y' = cy$, but we might be wrong.

A good example is $y' = -y$. Then $c = -1$ and $y = e^{-t}$—the true solution decays.

†It's not. $s$ is too big. I knew it.
The calculator gives the following answers \( a^n \) for \( n = 2, 10, 20 \):

<table>
<thead>
<tr>
<th>( \Delta t )</th>
<th>( a = 1 + c\Delta t )</th>
<th>( a^2 )</th>
<th>( a^{10} )</th>
<th>( a^{20} )</th>
</tr>
</thead>
<tbody>
<tr>
<td>3</td>
<td>-2</td>
<td>4</td>
<td>1024</td>
<td>1048576</td>
</tr>
<tr>
<td>1</td>
<td>0</td>
<td>0</td>
<td>0</td>
<td>0</td>
</tr>
<tr>
<td>1/10</td>
<td>.90</td>
<td>.81</td>
<td>.35</td>
<td>.12</td>
</tr>
<tr>
<td>1/20</td>
<td>.95</td>
<td>.90</td>
<td>.60</td>
<td>.36</td>
</tr>
</tbody>
</table>

The big step \( \Delta t = 3 \) shows total instability (top row). The numbers blow up when they should decay. The row with \( \Delta t = 1 \) is equally useless (all zeros). In practice the magnitude of \( c\Delta t \) must come down to .10 or .05. For accurate calculations it would have to be even smaller, unless we change to a better difference equation. That is the right thing to do.

Notice the two reasonable numbers. They are .35 and .36, approaching \( e^{-1} = .37 \). They come from \( n = 10 \) (with \( \Delta t = 1/10 \)) and \( n = 20 \) (with \( \Delta t = 1/20 \)). Those have the same clock time \( n\Delta t = 1 \):

\[
\left(1 - \frac{1}{10}\right)^{10} = .35 \quad \left(1 - \frac{1}{20}\right)^{20} = .36 \quad \left(1 - \frac{1}{n}\right)^n \to e^{-1} = .37.
\]

The main diagonal of the table is executing \( (1 + x/n)^n \to e^x \) in the case \( x = -1 \).

Final question: How quickly are .35 and .36 converging to \( e^{-1} = .37 \)? With \( \Delta t = .10 \) the error is .02. With \( \Delta t = .05 \) the error is .01. Cutting the time step in half cuts the error in half. We are not keeping enough digits to be sure, but the error seems close to \( \frac{1}{2}\Delta t \). To test that, apply the "quick method" and estimate \( a^n = (1 - \Delta t)^n \) from its logarithm:

\[
\ln((1 - \Delta t)^n) = n \ln(1 - \Delta t) \approx n\left[-\Delta t - \frac{1}{2}(\Delta t)^2\right] = -1 - \frac{1}{2}\Delta t.
\]

The clock time is \( n\Delta t = 1 \). Now take exponentials of the far left and right:

\[
a^n = (1 - \Delta t)^n \approx e^{-1} e^{-\Delta t/2} \approx e^{-1}(1 - \frac{1}{2}\Delta t).
\]

The difference between \( a^n \) and \( e^{-1} \) is the last term \( \frac{1}{2}\Delta t e^{-1} \). Everything comes down to one question: Is that error the same as \( \frac{1}{2}\Delta t \)? The answer is yes, because \( e^{-1/2} \) is 1/5. If we keep only one digit, the prediction is perfect!

That took an hour to work out, and I hope it takes longer than \( \Delta t \) to read. I wanted you to see in use the properties of \( \ln x \) and \( e^x \). The exact property \( \ln a^n = n \ln a \) came first. In the middle of (15) was the key approximation \( \ln(1 + x) \approx x - \frac{1}{2}x^2 \), with \( x = -\Delta t \). That \( x^2 \) term uses the second derivative (Section 6.4). At the very end came \( e^x \approx 1 + x \).

A linear approximation shows convergence: \( (1 + x/n)^n \to e^x \). A quadratic shows the error: proportional to \( \Delta t = 1/n \). It is like using rectangles for areas, with error proportional to \( \Delta x \). This minimal accuracy was enough to define the integral, and here it is enough to define \( e \). It is completely unacceptable for scientific computing.

The trapezoidal rule, for integrals or for \( y' = cy \), has errors of order \( (\Delta x)^2 \) and \( (\Delta t)^2 \). All good software goes further than that. Euler's first-order method could not predict the weather before it happens.

\[
\text{Euler's Method for } \frac{dy}{dt} = F(y, t); \quad \frac{y(t + \Delta t) - y(t)}{\Delta t} = F(y(t), t).
\]
Read-through questions

The infinite series for \( e^x \) is \( \sum_{n=0}^{\infty} \frac{x^n}{n!} \). Its derivative is \( e^x \). The denominator \( n! \) is called \( \frac{1}{n} \). It equals \( 9 \). At \( x = 9 \) the series for \( e^x \) is \( 9 \).

6.6 Exercises

The difference equation \( y(t+1) = ay(t) \) yields \( y(t) = y(0) \) when \( a = 1 \).

1 Write down a power series \( y = 1 - x + \cdots \) whose derivative is \( -y \).
2 Write down a power series \( y = 1 + 2x + \cdots \) whose derivative is \( 2y \).
3 Find two series that are equal to their second derivatives.
4 By comparing \( e = 1 + 1 + \frac{1}{2} + \frac{1}{3} + \cdots \) with a larger series (whose sum is easier) show that \( e < 3 \).
5 At 5% interest compute the output from \$1000 in a year with 6-month and 3-month and weekly compounding.

6 With the quick method \( \ln(1 + x) \approx x \), estimate \( \ln(1 - 1/n)^n \) and \( \ln(1 + 2/n)^n \). Then take exponentials to find the two limits.
7 With the slow method multiply out the three terms of \( (1 - \frac{1}{3})^3 \) and the five terms of \( (1 - \frac{1}{2})^5 \). What are the first three terms of \( (1 - 1/n)^n \), and what are their limits as \( n \to \infty \)?
8 The slow method leads to \( 1 + 1/2! - 1/3! + \cdots \) for the limit of \( (1 - 1/n)^n \). What is the sum of this infinite series — the exact sum and the sum after five terms?
9 Knowing that \( (1 + 1/n)^n \to e \), explain \( (1 + 1/n)^{2n} \to e^2 \) and \( (1 + 2/n)^n \to e^2 \).
10 What are the limits of \( (1 + 1/n)^n \) and \( (1 + 1/n)^n \)? OK to use a calculator to guess these limits.
11 (a) The power \( (1 + 1/n)^n \) (decreases) (increases) with \( n \), as we compound more often. (b) The derivative of \( f(x) = x \ln(1 + 1/x) \), which is \( \frac{1}{x} \), should be \((<0)(>0)\). This is confirmed by Problem 12.
12 Show that \( \ln(1 + 1/x) > 1/(x + 1) \) by drawing the graph of \( 1/x \). The area from \( t = 1 \) to \( 1 + 1/x \) is \( \frac{1}{x} \). The rectangle inside it has area \( \frac{1}{x} \).
13 Take three steps of \( y(t+1) = 2y(t) \) from \( y_0 = 1 \).
14 Take three steps of \( y(t+1) = 2y(t) + 1 \) from \( y_0 = 0 \).

Solve the difference equations 15–22.

15 \( y(t+1) = 3y(t), \ y_0 = 4 \)
16 \( y(t+1) = \frac{1}{2}y(t), \ y_0 = 1 \)
17 \( y(t+1) = y(t) + 1, \ y_0 = 0 \)
18 \( y(t+1) = y(t) - 1, \ y_0 = 0 \)
19 \( y(t+1) = 3y(t) + 1, \ y_0 = 0 \)
20 \( y(t+1) = 3y(t) + 1, \ y_0 = 1 \)
21 \( y(t+1) = ay(t) + 1, \ y_0 = 0 \)
22 \( y(t+1) = ay(t) + 1, \ y_0 = 0 \)

In 23–26, which initial value produces \( y_1 = y_0 \) (steady state)?
23 \( y(t+1) = 2y(t) - 6 \)
24 \( y(t+1) = \frac{1}{2}y(t) - 6 \)
25 \( y(t+1) = -y(t) + 6 \)
26 \( y(t+1) = -\frac{1}{2}y(t) + 6 \)

27 In Problems 23 and 24, start from \( y_0 = 2 \) and take three steps to reach \( y_1 \). Is this approaching a steady state?
28 For which numbers \( a \) does \( (1 - a)/(1 - a) \) approach a limit as \( t \to \infty \) and what is the limit?
29 The price \( P \) is determined by supply = demand or \( -dP(t + 1) + b = cP(t) \). Which price \( P \) is not changed from one year to the next?
30 Find \( P(t) \) from the supply-demand equation with \( c = 1 \), \( d = 2 \), \( b = 8 \), \( P(0) = 0 \). What is the steady state as \( t \to \infty \)?

Assume 10% interest (so \( a = 1 + i = 1.1 \)) in Problems 31–38.
31 At 10% interest compounded quarterly, what is the effective rate?
32 At 10% interest compounded daily, what is the effective rate?
33 Find the future value in 20 years of \$1000 deposited now.
34 Find the present value of \$1000 promised in twenty years.
6.7 Hyperbolic Functions

This section combines $e^x$ with $e^{-x}$. Up to now those functions have gone separate ways—one increasing, the other decreasing. But two particular combinations have earned names of their own ($\cosh x$ and $\sinh x$):

$$\text{hyperbolic cosine} \quad \cosh x = \frac{e^x + e^{-x}}{2}$$

$$\text{hyperbolic sine} \quad \sinh x = \frac{e^x - e^{-x}}{2}$$

The first name rhymes with "gosh". The second is usually pronounced "cinch".

The graphs in Figure 6.18 show that $\cosh x > \sinh x$. For large $x$ both hyperbolic functions come extremely close to $\frac{1}{2}e^x$. When $x$ is large and negative, it is $e^{-x}$ that dominates. $\cosh x$ still goes up to $+\infty$ while $\sinh x$ goes down to $-\infty$ (because $\sinh x$ has a minus sign sign front of $e^{-x}$).

![Cosh x and sinh x. The hyperbolic functions combine \(\frac{1}{2}e^x\) and \(\frac{1}{2}e^{-x}\).](image1)

![Gateway Arch courtesy of the St. Louis Visitors Commission.](image2)

The following facts come directly from $\frac{1}{2}(e^x + e^{-x})$ and $\frac{1}{2}(e^x - e^{-x})$:

- $\cosh(-x) = \cosh x$ and $\cosh 0 = 1$ (cosh is even like the cosine)
- $\sinh(-x) = -\sinh x$ and $\sinh 0 = 0$ (sinh is odd like the sine)

---

35 For a mortgage of $100,000 over 20 years, what is the monthly payment?

36 For a car loan of $10,000 over 6 years, what is the monthly payment?

37 With annual compounding of deposits $s = 1000$, what is the balance in 20 years?

38 If you repay $s = 1000$ annually on a loan of $8000$, when are you paid up? (Remember interest.)

39 Every year two thirds of the available houses are sold, and 1000 new houses are built. What is the steady state of the housing market—how many are available?

40 If a loan shark charges 5% interest a month on the $1000 you need for blackmail, and you pay $60 a month, how much do you still owe after one month (and after a year)?

41 Euler charges $c = 100\%$ interest on his $1 fee for discovering $e$. What do you owe (including the $1$) after a year with (a) no compounding; (b) compounding every week; (c) continuous compounding?

42 Approximate $(1 + 1/n)^n$ as in (15) and (16) to show that you owe Euler about $e - e/2n$. Compare Problem 6.2.5.

43 My Visa statement says monthly rate = 1.42% and yearly rate = 17%. What is the true yearly rate, since Visa compounds the interest? Give a formula or a number.

44 You borrow $y_0 = 80,000$ at 9% to buy a house.
   (a) What are your monthly payments $s$ over 30 years?
   (b) How much do you pay altogether?
The graph of cosh \(x\) corresponds to a *hanging cable* (hanging under its weight). Turned upside down, it has the shape of the Gateway Arch in St. Louis. That must be the largest upside-down cosh function ever built. A cable is easier to construct than an arch, because gravity does the work. With the right axes in Problem 55, the height of the cable is a stretched-out cosh function called a *catenary*:

\[
y = a \cosh \left(\frac{x}{a}\right) \quad \text{(cable tension/cable density = a)}.
\]

Busch Stadium in St. Louis has 96 catenary curves, to match the Arch.

The properties of the hyperbolic functions come directly from the definitions. There are too many properties to memorize—and no reason to do it! One rule is the most important. *Every fact about sines and cosines is reflected in a corresponding fact about sinh \(x\) and cosh \(x\).* Often the only difference is a minus sign. Here are four properties:

1. \((\cosh x)^2 - (\sinh x)^2 = 1\) [instead of \((\cos x)^2 + (\sin x)^2 = 1\)]
   
   Check:
   \[
   \left[\frac{e^x + e^{-x}}{2}\right]^2 - \left[\frac{e^x - e^{-x}}{2}\right]^2 = \frac{e^{2x} + 2 + e^{-2x} - e^{2x} + 2 - e^{-2x}}{4} = 1
   \]

2. \(\frac{d}{dx}(\cosh x) = \sinh x\) [instead of \(\frac{d}{dx}(\cos x) = -\sin x\)]

3. \(\frac{d}{dx}(\sinh x) = \cosh x\) [like \(\frac{d}{dx}(\sin x) = \cos x\)]

4. \(\int \sinh x \, dx = \cosh x + C\) and \(\int \cosh x \, dx = \sinh x + C\)

**Fig. 6.20** The unit circle \(\cos^2 t + \sin^2 t = 1\) and the unit hyperbola \(\cosh^2 t - \sinh^2 t = 1\).

Property 1 is the connection to hyperbolas. It is responsible for the "h" in cosh and sinh. Remember that \((\cos x)^2 + (\sin x)^2 = 1\) puts the point \((\cos x, \sin x)\) onto a *unit circle.* As \(x\) varies, the point goes around the circle. The ordinary sine and cosine are "circular functions." Now look at \((\cosh x, \sinh x)\). Property 1 is \((\cosh x)^2 - (\sinh x)^2 = 1\), so this point travels on the *unit hyperbola* in Figure 6.20.

You will guess the definitions of the other four hyperbolic functions:

\[
\begin{align*}
\tanh x &= \frac{\sinh x}{\cosh x} = \frac{e^x - e^{-x}}{e^x + e^{-x}} \\
\coth x &= \frac{\cosh x}{\sinh x} = \frac{e^x + e^{-x}}{e^x - e^{-x}} \\
\sech x &= \frac{1}{\cosh x} = \frac{2}{e^x + e^{-x}} \\
\csch x &= \frac{1}{\sinh x} = \frac{2}{e^x - e^{-x}}
\end{align*}
\]

I think "tanh" is pronounceable, and "sech" is easy. The others are harder. Their
6.7 Hyperbolic Functions

Properties come directly from \( \cosh^2 x - \sinh^2 x = 1 \). Divide by \( \cosh^2 x \) and \( \sinh^2 x \):

\[
1 - \tanh^2 x = \text{sech}^2 x \quad \text{and} \quad \coth^2 x - 1 = \text{csch}^2 x
\]

\( (\tanh x)' = \sec^2 x \quad \text{and} \quad (\text{sech} x)' = -\text{sech} x \tanh x \)

\[
\int \tanh x \, dx = \int \frac{\sinh x}{\cosh x} \, dx = \ln(\cosh x) + C.
\]

**Inverse Hyperbolic Functions**

You remember the angles \( \sin^{-1} x \) and \( \tan^{-1} x \) and \( \sec^{-1} x \). In Section 4.4 we differentiated those inverse functions by the chain rule. The main application was to **integrals**. If we happen to meet \( \int \frac{dx}{1 + x^2} \), it is \( \tan^{-1} x + C \). The situation for \( \sinh^{-1} x \) and \( \tanh^{-1} x \) and \( \text{sech}^{-1} x \) is the same except for sign changes — which are expected for hyperbolic functions. We write down the **three new derivatives**:

\[
y = \sinh^{-1} x \quad \text{(meaning } x = \sinh y) \quad \text{has } \frac{dy}{dx} = \frac{1}{\sqrt{x^2 + 1}} \quad (1)
\]

\[
y = \tanh^{-1} x \quad \text{(meaning } x = \tanh y) \quad \text{has } \frac{dy}{dx} = \frac{1}{1 - x^2} \quad (2)
\]

\[
y = \text{sech}^{-1} x \quad \text{(meaning } x = \text{sech} y) \quad \text{has } \frac{dy}{dx} = \frac{-1}{x\sqrt{1 - x^2}} \quad (3)
\]

Problems 44–46 compute \( dy/dx \) from \( 1/(dx/dy) \). The alternative is to use logarithms. Since \( \ln x \) is the inverse of \( e^x \), we can express \( \sinh^{-1} x \) and \( \tanh^{-1} x \) and \( \text{sech}^{-1} x \) as logarithms. Here is \( y = \tanh^{-1} x \):

\[
y = \frac{1}{2} \ln \left[ \frac{1 + x}{1 - x} \right] \quad \text{has slope} \quad \frac{dy}{dx} = \frac{1}{2} \frac{1 + x}{1 + x} - \frac{1}{2} \frac{1 - x}{1 - x} = \frac{1}{1 - x^2}.
\]

The last step is an ordinary derivative of \( \frac{1}{2} \ln(1 + x) - \frac{1}{2} \ln(1 - x) \). Nothing is new except the answer. But where did the logarithms come from? In the middle of the following identity, multiply above and below by \( \cosh y \):

\[
1 + x = \frac{1 + \tanh y}{1 - \tanh y} = \frac{\cosh y + \sinh y}{\cosh y - \sinh y} = \frac{e^y}{e^{-y}} = e^{2y}.
\]

Then \( 2y \) is the logarithm of the left side. This is the first equation in \( (4) \), and it is the third formula in the following list:

\[
\sinh^{-1} x = \ln \left[ x + \sqrt{x^2 + 1} \right] \quad \cosh^{-1} x = \ln \left[ x + \sqrt{x^2 - 1} \right]
\]

\[
\tanh^{-1} x = \frac{1}{2} \ln \left[ \frac{1 + x}{1 - x} \right] \quad \text{sech}^{-1} x = \ln \left[ \frac{1 + \sqrt{1 - x^2}}{x} \right]
\]

**Remark 1** Those are listed only for reference. If possible do not memorize them. The derivatives in equations (1), (2), (3) offer a choice of antiderivatives — either inverse functions or logarithms (most tables prefer logarithms). The inside cover of the book has

\[
\int \frac{dx}{1 - x^2} = \frac{1}{2} \ln \left[ \frac{1 + x}{1 - x} \right] + C \quad \text{(in place of } \tanh^{-1} x + C \text{)}.
\]

**Remark 2** Logarithms were not seen for \( \sin^{-1} x \) and \( \tan^{-1} x \) and \( \sec^{-1} x \). You might
wonder why. How does it happen that \( \tanh^{-1}x \) is expressed by logarithms, when the parallel formula for \( \tan^{-1}x \) was missing? Answer: There must be a parallel formula. To display it I have to reveal a secret that has been hidden throughout this section.

The secret is one of the great equations of mathematics. What formulas for \( \cos x \) and \( \sin x \) correspond to \( e^x + e^{-x} \) and \( e^x - e^{-x} \)? With so many analogies (circular vs. hyperbolic) you would expect to find something. The formulas do exist, but they involve imaginary numbers. Fortunately they are very simple and there is no reason to withhold the truth any longer:

\[
\cos x = \frac{1}{2}(e^{ix} + e^{-ix}) \quad \text{and} \quad \sin x = \frac{1}{2i}(e^{ix} - e^{-ix}).
\] (5)

It is the imaginary exponents that kept those identities hidden. Multiplying \( \sin x \) by \( i \) and adding to \( \cos x \) gives Euler’s unbelievably beautiful equation

\[
\cos x + i \sin x = e^{ix}.
\] (6)

That is parallel to the non-beautiful hyperbolic equation \( \cosh x + \sinh x = e^x \).

I have to say that (6) is infinitely more important than anything hyperbolic will ever be. The sine and cosine are far more useful than the sinh and cosh. So we end our record of the main properties, with exercises to bring out their applications.

**Read-through questions**

1. Find \( \cosh x + \sinh x \), \( \cosh x - \sinh x \), and \( \cosh x \sinh x \).
2. From the definitions of \( \cosh x \) and \( \sinh x \), find their derivatives.
3. Show that both functions satisfy \( y'' = y \).
4. By the quotient rule, verify \( (\tanh x)' = \text{sech}^2 x \).
5. Derive \( \cosh^2 x + \sinh^2 x = \cosh 2x \), from the definitions.
6. From the derivative of Problem 5 find \( \sinh 2x \).
7. The parallel to \( (\cos x + i \sin x)^n = \cos nx + i \sin nx \) is a hyperbolic formula \( (\cosh x + \sinh x)^n = \cosh nx + \ldots \).
8. Prove \( \sinh(x + y) = \sinh x \cosh y + \cosh x \sinh y \) by changing to exponentials. Then the \( x \)-derivative gives \( \cosh(x + y) = \ldots \).

**Exercises**

**Find the derivatives of the functions 9–18:**

9. \( \cosh(3x + 1) \)
10. \( \sinh x^2 \)
11. \( 1/\cosh x \)
12. \( \sinh(\ln x) \)
13. \( \cosh^2 x + \sinh^2 x \)
14. \( \cosh^2 x - \sinh^2 x \)
15. \( \tanh \sqrt{x^2 + 1} \)
16. \( (1 + \tanh x)/(1 - \tanh x) \)
17. \( \sinh^2 x \)
18. \( \ln(\text{sech} x + \tanh x) \)
19. Find the minimum value of \( \cosh(\ln x) \) for \( x > 0 \).
20. From \( \tanh x = \frac{1}{2} \) find \( \text{sech} x, \cosh x, \sinh x, \coth x, \text{csch} x \).
21. Do the same if \( \tanh x = -\frac{12}{13} \).
22. Find the other five values if \( \sinh x = 2 \).
23. Find the other five values if \( \cosh x = 1 \).
24. Compute \( \sinh(\ln 5) \) and \( \tanh(2 \ln 4) \).

**Find antiderivatives for the functions in 25–32:**

25. \( \cosh(2x + 1) \)
26. \( x \cosh(x^2) \)
27. \( \cosh^2 x \sinh x \)
28. \( \tanh^2 x \text{sech}^2 x \)
29. \( \frac{\sinh x}{1 + \cosh x} \)
30. \( \coth x = \frac{e^x + e^{-x}}{e^x - e^{-x}} \)
31. \( \sinh x + \cosh x \)
32. \( (\sinh x + \cosh x)^r \)
The triangle in Figure 6.20 has area \( \frac{1}{2} \cosh t \sinh t \).

(a) Integrate to find the shaded area below the hyperbola

(b) For the area \( A \) in red verify that \( dA/dt = \frac{1}{2} \)

(c) Conclude that \( A = \frac{1}{2}t + C \) and show \( C = 0 \).

Sketch graphs of the functions in 34–40.

34 \( y = \tanh x \) (with inflection point)

35 \( y = \coth x \) (in the limit as \( x \to \infty \))

36 \( y = \sech x \)

37 \( y = \sinh^{-1} x \)

38 \( y = \cosh^{-1} x \) for \( x \geq 1 \)

39 \( y = \sech^{-1} x \) for \( 0 < x < 1 \)

40 \( y = \tanh^{-1} x = \frac{1}{2} \ln \left( \frac{1 + x}{1 - x} \right) \) for \( |x| < 1 \)

41 (a) Multiplying \( x = \sinh y = \frac{1}{2}(e^y - e^{-y}) \) by \( 2e^y \) gives \( (e^y)^2 - 2x(e^y) - 1 = 0 \). Solve as a quadratic equation for \( e^y \).

(b) Take logarithms to find \( y = \sinh^{-1} x \) and compare with the text.

42 (a) Multiplying \( x = \cosh y = \frac{1}{2}(e^y + e^{-y}) \) by \( 2e^y \) gives \( (e^y)^2 - 2x(e^y) + 1 = 0 \). Solve for \( e^y \).

(b) Take logarithms to find \( y = \cosh^{-1} x \) and compare with the text.

43 Turn (4) upside down to prove \( y' = -1/(1 - x^2) \), if \( y = \coth^{-1} x \).

44 Compute \( dy/dx = 1/\sqrt{x^2 + 1} \) by differentiating \( x = \sinh y \) and using \( \cosh^2 y - \sinh^2 y = 1 \).

45 Compute \( dy/dx = 1/(1 - x^2) \) if \( y = \tanh^{-1} x \) by differentiating \( x = \tanh y \) and using \( \sech^2 y + \tanh^2 y = 1 \).

46 Compute \( dy/dx = -1/x\sqrt{1 - x^2} \) for \( y = \sech^{-1} x \), by differentiating \( x = \sech y \).

From formulas (1), (2), (3) or otherwise, find antiderivatives in 47–52:

47 \( \int dx/(4 - x^2) \)

48 \( \int dx/(a^2 - x^2) \)

49 \( \int dx/\sqrt{x^2 + 1} \)

50 \( \int x \, dx/\sqrt{x^2 + 1} \)

51 \( \int dx/x\sqrt{1 - x^2} \)

52 \( \int dx/\sqrt{1 - x^2} \)

53 Compute \( \int_0^{1/2} \frac{dx}{1-x^2} \) and \( \int_0^1 \frac{dx}{1-x^2} \).

54 A falling body with friction equal to velocity squared obeys \( dv/dt = g - v^2 \).

(a) Show that \( v(t) = \sqrt{g} \tanh \sqrt{gt} \) satisfies the equation.

(b) Derive this \( v \) yourself, by integrating \( dv/(g - v^2) = dt \).

(c) Integrate \( v(t) \) to find the distance \( f(t) \).

55 A cable hanging under its own weight has slope \( S = \frac{dy}{dx} \) that satisfies \( dS/dx = c/\sqrt{1 + S^2} \). The constant \( c \) is the ratio of cable density to tension.

(a) Show that \( S = \sinh cx \) satisfies the equation.

(b) Integrate \( dy/dx = \sinh cx \) to find the cable height \( y(x) \), if \( y(0) = 1/c \).

(c) Sketch the cable hanging between \( x = -L \) and \( x = L \) and find how far it sags down at \( x = 0 \).

56 The simplest nonlinear wave equation (Burgers' equation) yields a waveform \( W(x) \) that satisfies \( W'' = WW' - W' \). One integration gives \( W' = \frac{1}{2}W^2 - W \).

(a) Separate variables and integrate: \( dx = dW/(2W - W) = -dW/(2 - W) - dW/W \).

(b) Check \( W' = \frac{1}{2}W^2 - W \).

57 A solitary water wave has a shape satisfying the KdV equation \( y'' = y - 6yy' \).

(a) Integrate once to find \( y'' \). Multiply the answer by \( y' \).

(b) Integrate again to find \( y' \) (all constants of integration are zero).

(c) Show that \( y = \frac{1}{2} \sech^2(x/2) \) gives the shape of the "soliton."
## Contents

### CHAPTER 4  The Chain Rule
- 4.1 Derivatives by the Chain Rule 154
- 4.2 Implicit Differentiation and Related Rates 160
- 4.3 Inverse Functions and Their Derivatives 164
- 4.4 Inverses of Trigonometric Functions 171

### CHAPTER 5  Integrals
- 5.1 The Idea of the Integral 177
- 5.2 Antiderivatives 182
- 5.3 Summation vs. Integration 187
- 5.4 Indefinite Integrals and Substitutions 195
- 5.5 The Definite Integral 201
- 5.6 Properties of the Integral and the Average Value 206
- 5.7 The Fundamental Theorem and Its Consequences 213
- 5.8 Numerical Integration 220

### CHAPTER 6  Exponentials and Logarithms
- 6.1 An Overview 228
- 6.2 The Exponential $e^x$ 236
- 6.3 Growth and Decay in Science and Economics 242
- 6.4 Logarithms 252
- 6.5 Separable Equations Including the Logistic Equation 259
- 6.6 Powers Instead of Exponentials 267
- 6.7 Hyperbolic Functions 277

### CHAPTER 7  Techniques of Integration
- 7.1 Integration by Parts 283
- 7.2 Trigonometric Integrals 288
- 7.3 Trigonometric Substitutions 294
- 7.4 Partial Fractions 300
- 7.5 Improper Integrals 305

### CHAPTER 8  Applications of the Integral
- 8.1 Areas and Volumes by Slices 311
- 8.2 Length of a Plane Curve 320
- 8.3 Area of a Surface of Revolution 325
- 8.4 Probability and Calculus 328
- 8.5 Masses and Moments 336
- 8.6 Force, Work, and Energy 342
Chapter 5 introduced the integral as a limit of sums. The calculation of areas was started—by hand or computer. Chapter 6 opened a different door. Its new functions $e^x$ and $\ln x$ led to differential equations. You might say that all along we have been solving the special differential equation $df/dx = u(x)$. The solution is $f = \int u(x)dx$. But the step to $dy/dx = cy$ was a breakthrough.

The truth is that we are able to do remarkable things. *Mathematics has a language, and you are learning to speak it.* A short time ago the symbols $dy/dx$ and $\int u(x)dx$ were a mystery. (My own class was not too sure about $u(x)$ itself—the symbol for a function.) It is easy to forget how far we have come, in looking ahead to what is next.

I do want to look ahead. For integrals there are two steps to take—more functions and more applications. *By using mathematics we make it live.* The applications are most complete when we know the integral. This short chapter will widen (very much) the range of functions we can integrate. A computer with symbolic algebra widens it more.

Up to now, integration depended on recognizing derivatives. If $u(x) = \sec^2x$ then $f(x) = \tan x$. To integrate $\tan x$ we use a substitution:

$$\int \frac{\sin x}{\cos x} dx = -\int \frac{du}{u} = -\ln u = -\ln \cos x.$$  

What we need now are *techniques for other integrals*, to change them around until we can attack them. Two examples are $\int x \cos x \, dx$ and $\int \sqrt{1-x^2} \, dx$, which are not immediately recognizable. With integration by parts, and a new substitution, they become simple.

Those examples indicate where this chapter starts and stops. With reasonable effort (and the help of tables, which is fair) you can integrate important functions. With intense effort you could integrate even more functions. In older books that extra exertion was made—it tended to dominate the course. They had integrals like $\int (x + 1)dx/\sqrt{2x^2 - 6x + 4}$, which we could work on if we had to. *Our time is too valuable for that!* Like long division, the ideas are for us and their intricate elaboration is for the computer.

Integration by parts comes first. Then we do new substitutions. Partial fractions is a useful idea (already applied to the logistic equation $y' = cy - by^2$). In the last section $x$ goes to infinity or $y(x)$ goes to infinity—but the area stays finite. These improper integrals are quite common. Chapter 8 brings the applications.
There are two major ways to manipulate integrals (with the hope of making them easier). Substitutions are based on the chain rule, and more are ahead. Here we present the other method, based on the \textit{product rule}. The reverse of the product rule, to find integrals not derivatives, is \textit{integration by parts}.

We have mentioned $\int \cos^2 x \, dx$ and $\int \ln x \, dx$. Now is the right time to compute them (plus more examples). You will see how $\int \ln x \, dx$ is exchanged for $\int 1 \, dx$ — a definite improvement. Also $\int x e^x \, dx$ is exchanged for $\int e^x \, dx$. The difference between the harder integral and the easier integral is a known term — that is the point.

One note before starting: Integration by parts is \textit{not just a trick} with no meaning. On the contrary, it expresses basic physical laws of equilibrium and force balance. It is a foundation for the theory of differential equations (and even delta functions). The final paragraphs, which are completely optional, illustrate those points too.

We begin with the product rule for the derivative of $u(x)$ times $v(x)$:

\begin{equation}
\frac{d}{dx} u(x)v(x) = \frac{du}{dx} v(x) + u(x) \frac{dv}{dx}.
\end{equation}

\textit{Integrate both sides.} On the right, integration brings back $u(x)v(x)$. On the left are two integrals, and one of them moves to the other side (with a minus sign):

\begin{equation}
\int u(x) \frac{dv}{dx} \, dx = u(x)v(x) - \int v(x) \frac{du}{dx} \, dx.
\end{equation}

That is the key to this section — not too impressive at first, but very powerful. It is \textit{integration by parts} ($u$ and $v$ are the parts). In practice we write it without $x$'s:

\[ \int u \, dv = uv - \int v \, du. \]  

The problem of integrating $u \, dv/dx$ is changed into the problem of integrating $v \, du/dx$. There is a minus sign to remember, and there is the "integrated term" $u(x)v(x)$. In the \textit{definite integral}, that product $u(x)v(x)$ is evaluated at the endpoints $a$ and $b$:

\begin{equation}
\int_a^b u \, \frac{dv}{dx} \, dx = u(b)v(b) - u(a)v(a) - \int_a^b v \, \frac{du}{dx} \, dx.
\end{equation}

\textit{The key is in choosing $u$ and $v$.} The goal of that choice is to make $\int v \, du$ easier than $\int u \, dv$. This is best seen by examples.

\textbf{EXAMPLE 1} For $\int \ln x \, dx$ choose $u = \ln x$ and $dv = dx$ (so $v = x$):

\[ \int \ln x \, dx = uv - \int v \, du = x \ln x - \int x \frac{1}{x} \, dx. \]

I used the basic formula (3). Instead of working with $\ln x$ (searching for an antiderivative), we now work with the right hand side. There $x$ times $1/x$ is 1. The integral of 1 is $x$. Including the minus sign and the integrated term $uv = x \ln x$ and the constant $C$, the answer is

\[ \int \ln x \, dx = x \ln x - x + C. \]
To repeat: We exchanged the integral of \( \ln x \) for the integral of \( 1 \).

**EXAMPLE 2**  
For \( \int x \cos x \, dx \) choose \( u = x \) and \( dv = \cos x \, dx \) (so \( v(x) = \sin x \)):  
\[
\int x \cos x \, dx = uv - \int v \, du = x \sin x - \int \sin x \, dx.
\]  
(6)  
Again the right side has a simple integral, which completes the solution:  
\[
\int x \cos x \, dx = x \sin x + \cos x + C.
\]  
(7)  
*Note*  
The new integral is not always simpler. We could have chosen \( u = \cos x \) and \( dv = x \, dx \). Then \( v(x) = \frac{1}{2} x^2 \). Integration using those parts give the true but useless result  
\[
\int x \cos x \, dx = uv - \int v \, du = \frac{1}{2} x^2 \cos x + \frac{1}{2} x^2 \sin x \, dx.
\]  
The last integral is harder instead of easier (\( x^2 \) is worse than \( x \)). In the forward direction this is no help. But in the opposite direction it simplifies \( \int \frac{1}{2} x^2 \sin x \, dx \). The idea in choosing \( u \) and \( v \) is this: *Try to give \( u \) a nice derivative and \( du \) a nice integral.*

**EXAMPLE 3**  
For \( \int (\cos x)^2 \, dx \) choose \( u = \cos x \) and \( dv = \cos x \, dx \) (so \( v = \sin x \)):  
\[
\int (\cos x)^2 \, dx = uv - \int v \, du = x \sin x + \int (\sin x)^2 \, dx.
\]  
(8)  
The integral of \( (\sin x)^2 \) is no better and no worse than the integral of \( (\cos x)^2 \). But we never see \( (\sin x)^2 \) without thinking of \( 1 - (\cos x)^2 \). So substitute for \( (\sin x)^2 \):  
\[
\int (\cos x)^2 \, dx = \cos x \sin x + \int 1 \, dx - \int (\cos x)^2 \, dx.
\]  
The last integral on the right joins its twin on the left, and \( \int 1 \, dx = x \):  
\[
2 \int (\cos x)^2 \, dx = \cos x \sin x + x.
\]  
Dividing by 2 gives the answer, which is definitely not \( \frac{1}{2} (\cos x)^3 \). Add any \( C \):  
\[
\int (\cos x)^2 \, dx = \frac{1}{2} (\cos x \sin x + x) + C.
\]  
(9)  
*Question*  
Integrate \( (\cos x)^2 \) from 0 to \( 2\pi \). Why should the area be \( \pi \)?

*Answer*  
The definite integral is \( \frac{1}{2} [\cos x \sin x + x]_0^{2\pi} \). This does give \( \pi \). That area can also be found by common sense, starting from \( (\cos x)^2 + (\sin x)^2 = 1 \). The area under 1 is \( 2\pi \). The areas under \( (\cos x)^2 \) and \( (\sin x)^2 \) are the same. So each one is \( \pi \).

**EXAMPLE 4**  
Evaluate \( \int \tan^{-1} x \, dx \) by choosing \( u = \tan^{-1} x \) and \( v = x \):  
\[
\int \tan^{-1} x \, dx = uv - \int v \, du = x \tan^{-1} x - \int x \frac{dx}{1 + x^2}.
\]  
(9)  
The last integral has \( w = 1 + x^2 \) below and almost has \( dw = 2x \, dx \) above:  
\[
\int \frac{x \, dx}{1 + x^2} = \frac{1}{2} \int \frac{dw}{w} = \frac{1}{2} \ln w = \frac{1}{2} \ln(1 + x^2).
\]  
Substituting back into (9) gives \( \int \tan^{-1} x \, dx = x \tan^{-1} x - \frac{1}{2} \ln(1 + x^2) \). All the familiar inverse functions can be integrated by parts (take \( v = x \), and add " + \( C \)" at the end).  
Our final example shows how *two integrations by parts* may be needed, when the first one only simplifies the problem half way.

**EXAMPLE 5**  
For \( \int x^2 e^x \, dx \) choose \( u = x^2 \) and \( dv = e^x \, dx \) (so \( v = e^x \)):  
\[
\int x^2 e^x \, dx = uv - \int v \, du = x^2 e^x - \int e^x (2x \, dx).
\]  
(10)
The last integral involves $xe^x$. This is better than $x^2e^x$, but it still needs work:
\[ \int xe^x\,dx = uv - \int v\,du = xe^x - \int e^x\,dx \quad \text{(now } u = x). \tag{11} \]
Finally $e^x$ is alone. After two integrations by parts, we reach $\int e^x\,dx$. In equation (11), the integral of $xe^x$ is $xe^x - e^x$. Substituting back into (10),
\[ \int x^2e^x\,dx = x^2e^x - 2[xe^x - e^x] + C. \tag{12} \]
These five examples are in the list of prime candidates for integration by parts:

$x^ne^x$, $x^n\sin x$, $x^n\cos x$, $x^n\ln x$, $e^x\sin x$, $e^x\cos x$, $\sin^{-1}x$, $\tan^{-1}x$, …

This concludes the presentation of the method—brief and straightforward. Figure 7.1a shows how the areas $\int u\,dv$ and $\int v\,du$ fill out the difference between the big area $u(b)v(b)$ and the smaller area $u(a)v(a)$.

In the movie *Stand and Deliver*, the Los Angeles teacher Jaime Escalante computed $\int x^2\sin x\,dx$ with two integrations by parts. His success was through exercises—plus insight in choosing $u$ and $v$. (Notice the difference from $\int x\sin x^2\,dx$. That falls the other way—to a substitution.) The class did extremely well on the Advanced Placement Exam. If you saw the movie, you remember that the examiner didn’t believe it was possible. I spoke to him long after, and he confirms that practice was the key.

**THE DELTA FUNCTION**

From the most familiar functions we move to the least familiar. The delta function is the derivative of a step function. The step function $U(x)$ jumps from 0 to 1 at $x = 0$. We write $\delta(x) = dU/dx$, recognizing as we do it that there is no genuine derivative at the jump. The delta function is the limit of higher and higher spikes—from the "burst of speed" in Section 1.2. They approach an infinite spike concentrated at a single point (where $U$ jumps). This "non-function" may be unconventional—it is certainly optional—but it is important enough to come back to.

The slope $dU/dx$ is zero except at $x = 0$, where the step function jumps. Thus $\delta(x) = 0$ except at that one point, where the delta function has a "spike." We cannot give a value for $\delta$ at $x = 0$, but we know its integral across the jump. On every interval from $-A$ to $A$, the integral of $dU/dx$ brings back $U$:
\[ \int_{-A}^{A} \delta(x)\,dx = \int_{-A}^{A} \frac{dU}{dx} \,dx = U(x) \bigg|_{-A}^{A} = 1. \tag{13} \]
"The area under the infinitely tall and infinitely thin spike $\delta(x)$ equals 1."

So far so good. The integral of $\delta(x)$ is $U(x)$. We now integrate by parts for a crucial purpose—to find the area under $v(x)\delta(x)$. This is an ordinary function times the delta function. In some sense $v(x)$ times $\delta(x)$ equals $v(0)$ times $\delta(x)$—because away from $x = 0$ the product is always zero. Thus $e^x\delta(x)$ equals $\delta(x)$, and $\sin x \, \delta(x) = 0$. 

The area under $v(x)\delta(x)$ is $v(0)$—which integration by parts will prove:

7B The integral of $v(x)$ times $\delta(x)$ is $\int_{-A}^{A} v(x)\delta(x)dx = v(0)$.

The area is $v(0)$ because the spike is multiplied by $v(0)$—the value of the smooth function $v(x)$ at the spike. But multiplying infinity is dangerous, to say the least. (Two times infinity is infinity). We cannot deal directly with the delta function. It is only known by its integrals! As long as the applications produce integrals (as they do), we can avoid the fact that $\delta$ is not a true function.

The integral of $v(x)\delta(x) = v(x)dU/dx$ is computed "by parts:"

$$\int_{-A}^{A} v(x)\delta(x) dx = v(x)U(x) \bigg|_{-A}^{A} - \int_{-A}^{A} U(x) \frac{dv}{dx} dx.$$  \hspace{1cm} (14)

Remember that $U = 0$ or $U = 1$. The right side of (14) is our area $v(0)$:

$$v(A) \cdot 1 - \int_{0}^{A} 1 \frac{dv}{dx} dx = v(A) - (v(A) - v(0)) = v(0).$$  \hspace{1cm} (15)

When $v(x) = 1$, this answer matches $\int \delta dx = 1$. We give three examples:

$$\int_{-2}^{2} \cos x \delta(x) dx = 1 \quad \int_{6}^{5} (U(x) + \delta(x)) dx = 7 \quad \int_{-1}^{1} (\delta(x))^2 dx = \infty.$$  

A nightmare question occurs to me. What is the derivative of the delta function?

**INTEGRATION BY PARTS IN ENGINEERING**

Physics and engineering and economics frequently involve products. Work is force times distance. Power is voltage times current. Income is price times quantity. When there are several forces or currents or sales, we add the products. When there are infinitely many, we integrate (probably by parts).

I start with differential equations for the displacement $u$ at point $x$ in a bar:

$$-\frac{dv}{dx} = f(x) \text{ with } v(x) = k \frac{du}{dx}.  \hspace{1cm} (16)$$

This describes a hanging bar pulled down by a force $f(x)$. Each point $x$ moves through a distance $u(x)$. The top of the bar is fixed, so $u(0) = 0$. The stretching in the bar is $du/dx$. The internal force created by stretching is $v = k \frac{du}{dx}$. (This is Hooke’s law.) Equation (16) is a balance of forces on the small piece of the bar in Figure 7.2.

---

Fig. 7.2 Difference in internal force balances external force

$-\Delta v = f\Delta x$ or $-dv/dx = f(x)$

$v = W$ at $x = 1$ balances hanging weight
7.1 Integration by Parts

EXAMPLE 6 Suppose \( f(x) = F \), a constant force per unit length. We can solve (16):
\[
u(x) = -Fx + C \quad \text{and} \quad ku(x) = -\frac{1}{2}Fx^2 + Cx + D.
\]
(17)
The constants \( C \) and \( D \) are settled at the endpoints (as usual for integrals). At \( x = 0 \) we are given \( u = 0 \) so \( D = 0 \). At \( x = 1 \) we are given \( v = W \) so \( C = W + F \). Then \( u(x) \) and \( u(x) \) give force and displacement in the bar.

To see integration by parts, multiply \(-\frac{dv}{dx} = f(x)\) by \( u(x)\) and integrate:
\[
\int_0^1 f(x)u(x) \, dx = -\int_0^1 u(x) \, dv = -\left[u(1) - u(0)\right] + \int_0^1 u(x) \frac{dv}{dx} \, dx.
\]
(18)
The left side is force times displacement, or external work. The last term is internal force times stretching—or internal work. The integrated term has \( u(0) = 0 \)—the fixed support does no work. It also has \(-u(1)W\), the work by the hanging weight. The balance of forces has been replaced by a balance of work.

This is a touch of engineering mathematics, and here is the main point. Integration by parts makes physical sense! When \(-\frac{dv}{dx} = f\) is multiplied by other functions—called test functions or virtual displacements—then equation (18) becomes the principle of virtual work. It is absolutely basic to mechanics.

7.1 Exercises

Read-through questions

Integration by parts is the reverse of the \( \int uv \) rule. It changes \( \int u \, dv \) into \( \int \) minus \( \int \). In case \( u = x \) and \( dv = e^{2x} \, dx \), it changes \( \int xe^{2x} \, dx \) to \( \int d \) minus \( \int \). The definite integral \( \int_0^1 xe^{2x} \, dx \) becomes \( \int \) minus \( \int \).

In choosing \( u \) and \( dv \), the \( h \) of \( u \) and the \( i \) of \( dv \) should be as simple as possible. Normally ln \( x \) goes into \( \int \) and \( e^x \) goes into \( \int \). Prime candidates are \( u = x \) or \( x^2 \) and \( v = \sin x \) or \( \int \) or \( \int \). When \( u = x^2 \) we need \( \int \) integrations by parts. For \( \int \sin^n x \, dx \), the choice \( dv = dx \) leads to \( \int \) minus \( \int \).

If \( U \) is the unit step function, \( \frac{dU}{dx} = \delta \) is the unit \( \int \) function. The integral from \(-A\) to \( A \) is \( U(A) - U(-A) = \int \). The integral of \( \sin(x)\delta(x) \) equals \( \int \). The integral \( \int_0^1 \cos x \, \delta(x) \, dx \) equals \( \int \). In engineering, the balance of forces \(-\frac{dv}{dx} = f\) is multiplied by a displacement \( u(x)\) and integrated to give a balance of \( \int \).

Integrate 1–16, usually by parts (sometimes twice).

1 \( \int x \sin x \, dx \)
2 \( \int x e^{ax} \, dx \)
3 \( \int x e^{-ax} \, dx \)
4 \( \int x \cos 3x \, dx \)
5 \( \int x^2 \cos x \, dx \) (use Problem 1)
6 \( \int x \ln x \, dx \)
7 \( \int \ln(2x + 1) \, dx \)
8 \( \int x^2 e^{ax} \, dx \) (use Problem 2)

9 \( \int e^{x} \sin x \, dx \)
10 \( \int e^{x} \cos x \, dx \)

[9 and 10 need two integrations. I think \( e^x \) can be \( u \) or \( v \).]
11 \( \int e^{x} \sin bx \, dx \)
12 \( \int e^{x} \cos bx \, dx \)
13 \( \int \sin(ln x) \, dx \)
14 \( \int \cos(ln x) \, dx \)
15 \( \int (\ln x)^2 \, dx \)
16 \( \int x^2 \ln x \, dx \)
17 \( \int \sin^{-1} x \, dx \)
18 \( \int \cos^{-1}(2x) \, dx \)
19 \( \int x \tan^{-1} x \, dx \)
20 \( \int x^2 \sin x \, dx \) (from the movie)
21 \( \int x^2 \cos x \, dx \)
22 \( \int x^2 \sin x \, dx \)
23 \( \int x^2 e^{ax} \, dx \)
24 \( \int x \sec^{-1} x \, dx \)
25 \( \int x \tan x \, dx \)
26 \( \int x \cosh x \, dx \)

Compute the definite integrals 27–34.

27 \( \int_0^1 \ln x \, dx \)
28 \( \int_0^1 e^{x} \, dx \) (let \( u = \sqrt{x} \))
29 \( \int_0^1 e^{-2x} \, dx \)
30 \( \int_0^1 \ln(x^2) \, dx \)
31 \( \int_0^1 \cos x \, dx \)
32 \( \int_{-\infty}^{\infty} x \sin x \, dx \)
33 \( \int_0^1 \ln(x^2 + 1) \, dx \)
34 \( \int_0^{$1$/2} x^2 \sin x \, dx \)

In 35–40 derive "reduction formulas" from higher to lower powers.

35 \( \int x^n e^x \, dx = x^n e^x - n \int x^{n-1} e^x \, dx \)
36 \int x^n e^x dx = ________

37 \int x^n \cos x dx = x^n \sin x - n \int x^{n-1} \sin x dx

38 \int x^n \sin x dx = ________

39 \int (\ln x)^n dx = x(\ln x)^n - n \int (\ln x)^{n-1} dx

40 \int x(\ln x)^n dx = ________

41 How would you compute \( \int x \sin x e^x dx \) using Problem 9? Not necessary to do it.

42 How would you compute \( \int x e^x \tan^{-1} x dx \)? Don't do it.

43 (a) Integrate \( \int x^3 \sin x^2 dx \) by substitution and parts.

(b) The integral \( \int x^n \sin x^2 dx \) is possible if \( n \) is ________.

44-54 are about optional topics at the end of the section.

44 For the delta function \( \delta(x) \) find these integrals:

(a) \( \int_{-1}^{1} e^{x} \delta(x) dx \)

(b) \( \int_{-1}^{1} v(x) \delta(x) dx \)

(c) \( \int_{-1}^{1} \cos x \delta(x) dx \)

45 Solve \( dy/dx = 3\delta(x) \) and \( dy/dx = 3\delta(x) + y(x) \).

46 Strange fact: \( \delta(2x) \) is different from \( \delta(x) \). Integrate them both from \(-1\) to 1.

47 The integral of \( \delta(x) \) is the unit step \( U(x) \). Graph the next integrals \( R(x) = \int U(x) dx \) and \( Q(x) = \int R(x) dx \). The ramp \( R \) and quadratic spline \( Q \) are zero at \( x = 0 \).

48 In \( \delta(x - \frac{1}{2}) \), the spike shifts to \( x = \frac{1}{2} \). It is the derivative of the shifted step \( U(x - \frac{1}{2}) \). The integral of \( \delta(x)(x - \frac{1}{2}) \) equals the value of \( v \) at \( x = \frac{1}{2} \). Compute

(a) \( \int_{-\frac{1}{2}}^{\frac{1}{2}} \delta(x - \frac{1}{2}) dx \)

(b) \( \int_{-\frac{1}{2}}^{\frac{1}{2}} e^x \delta(x - \frac{1}{2}) dx \)

(c) \( \int_{-\frac{1}{2}}^{\frac{1}{2}} \cos x \delta(x - \frac{1}{2}) dx \)

49 The derivative of \( \delta(x) \) is extremely singular. It is a “dipole” known by its integrals. Integrate by parts in (b) and (c):

(a) \( \int_{-\frac{1}{2}}^{\frac{1}{2}} \frac{d}{dx} \delta(x) dx \)

(b) \( \int_{-\frac{1}{2}}^{\frac{1}{2}} x \frac{d}{dx} \delta(x) dx \)

(c) \( \int_{-\frac{1}{2}}^{\frac{1}{2}} v(x) \frac{d}{dx} \delta(x) dx = -v(0) \)

50 Why is \( \int_{-\frac{1}{2}}^{\frac{1}{2}} U(x) \delta(x) dx \) equal to \( \frac{1}{2} \)? (By parts.)

51 Choose limits of integration in \( v(x) = \int f(x) dx \) so that \( dv/dx = -f(x) \) and \( v = 0 \) at \( x = 1 \).

52 Draw the graph of \( v(x) \) if \( v(1) = 0 \) and \( -dv/dx = f(x) \):

(a) \( f = x \)

(b) \( f = U(x - \frac{1}{2}) \)

(c) \( f = \delta(x - \frac{1}{2}) \)

53 What integral \( u(x) \) solves \( k du/dx = u(x) \) with end condition \( u(0) = 0 \)? Find \( u(x) \) for the three \( v \)'s (not \( f \)'s) in Problem 52, and graph the three \( v \)'s.

54 Draw the graph of \( \Delta U/\Delta x = [U(x + \Delta x) - U(x)]/\Delta x \). What is the area under this graph?

Problems 55-62 need more than one integration.

55 Two integrations by parts lead to \( V = \int u v' dx \):

\[ \int u v' dx = uv - \int v u' dx \]

Test this rule on \( \int x^3 \sin x dx \).

56 After \( n \) integrations by parts, \( \int u dv dx \) becomes

\[ uv - u^{(1)} v^{(1)} + u^{(2)} v^{(2)} - \cdots + (-1)^n u^{(n)} v^{(n)} \]

\( u^{(n)} \) is the \( n \)th derivative of \( u \), and \( v^{(n)} \) is the \( n \)th integral of \( v \). Integrate the last term by parts to stretch this formula to \( n + 1 \) integrations.

57 Use Problem 56 to find \( \int x^2 e^x dx \).

58 From \( f(x) = f(0) = \int_{0}^{x} f(t) dt \); integrate by parts (notice \( dt \) not \( dx \)) to reach \( f(x) = f(0) + f'(0)x + \int_{0}^{x} f''(t)(x-t) dt \). Continuing as in Problem 56 produces Taylor’s formula:

\[ f(x) = f(0) + f'(0)x + \frac{1}{2!} f''(0)x^2 + \cdots + \frac{f^{(n)}(0)}{n!} \int_{0}^{x} f^{(n+1)}(t)(x-t)^n dt \]

59 What is the difference between \( \int_{0}^{x} u^{(n)} v^{(n)} dx \) and \( \int_{0}^{x} u^{(n)} v^{(n)} dx \)?

60 Compute the areas \( A = \int_{0}^{x} \ln x dx \) and \( B = \int_{0}^{x} e^x dx \). Mark them on the rectangle with corners \((0, 0), (x, 0), (x, e), (0, 1)\).

61 Find the mistake. I don’t believe \( e^x \cosh x = e^x \sinh x \):

\[ \int e^x \sinh x dx = e^x \cosh x - \int e^x \cosh x dx \]

\[ = e^x \cosh x - e^x \sinh x + \int e^x \sinh x dx \]

62 Choose \( C \) and \( D \) to make the derivative of \( C e^x \cos bx + D e^x \sin bx \) equal to \( e^x \cos bx \). Is this easier than integrating \( e^x \cos bx \) twice by parts?

### 7.2 Trigonometric Integrals

The next section will put old integrals into new forms. For example \( \int x^2 \sqrt{1 - x^2} dx \) will become \( \int \sin^2 \theta \cos^2 \theta d\theta \). That looks simpler because the square root is gone. But still \( \sin^2 \theta \cos^2 \theta \) has to be integrated. This brief section integrates any product of sines and cosines and secants and tangents.

There are two methods to choose from. One uses integration by parts, the other is based on trigonometric identities. Both methods try to make the integral easy (but that may take time). We follow convention by changing the letter \( \theta \) back to \( x \).
7.2 Trigonometric Integrals

Notice that \( \sin^4x \cos x \, dx \) is easy to integrate. It is \( u^4 \, du \). This is the goal in Example 1—to separate out \( \cos x \, dx \). It becomes \( du \), and \( \sin x \) is \( u \).

**EXAMPLE 1** \( \int \sin^2x \cos^3x \, dx \) (the exponent 3 is odd)

Solution Keep \( \cos x \, dx \) as \( du \). Convert the other \( \cos^2x \) to \( 1 - \sin^2x \):

\[
\int \sin^2x \cos^3x \, dx = \int \sin^2x(1 - \sin^2x) \cos x \, dx = \frac{\sin^3x}{3} - \frac{\sin^5x}{5} + C.
\]

**EXAMPLE 2** \( \int \sin^5x \, dx \) (the exponent 5 is odd)

Solution Keep \( \sin x \, dx \) and convert everything else to cosines. The conversion is always based on \( \sin^2x + \cos^2x = 1 \):

\[
\int (1 - \cos^2x)^2 \sin x \, dx = \int (1 - 2 \cos^2x + \cos^4x) \sin x \, dx.
\]

Now \( \cos x \) is \( u \) and \( -\sin x \, dx \) is \( du \). We have \( \int (-1 + 2u^2 - u^4) \, du \).

**General method for \( \int \sin^m x \cos^n x \, dx \), when \( m \) or \( n \) is odd**

If \( n \) is odd, separate out a single \( \cos x \, dx \). That leaves an even number of cosines.

Convert them to sines. Then \( \cos x \, dx \) is \( du \) and the sines are \( u \)'s.

If \( m \) is odd, separate out a single \( \sin x \, dx \) as \( du \). Convert the rest to cosines.

If \( m \) and \( n \) are both odd, use either method.

If \( m \) and \( n \) are both even, a new method is needed. Here are two examples.

**EXAMPLE 3** \( \int \cos^2x \, dx \) (\( m = 0, n = 2 \), both even)

There are two good ways to integrate \( \cos^2x \), but substitution is not one of them. If \( u \) equals \( \cos x \), then \( du \) is not here. The successful methods are integration by parts and double-angle formulas. Both answers are in equation (2) below—I don't see either one as the obvious winner.

Integrating \( \cos^2x \) by parts was Example 3 of Section 7.1. The other approach, by double angles, is based on these formulas from trigonometry:

\[
\cos^2x = \frac{1}{2}(1 + \cos 2x) \quad \sin^2x = \frac{1}{2}(1 - \cos 2x)
\]

(1)

The integral of \( \cos 2x \) is \( \frac{1}{2} \sin 2x \). So these formulas can be integrated directly. They give the only integrals you should memorize—either the integration by parts form, or the result from these double angles:

\[
\int \cos^2x \, dx \text{ equals } \frac{1}{2}(x + \sin x \cos x) \quad \text{or} \quad \frac{1}{2}x + \frac{1}{4} \sin 2x \quad \text{(plus C)}, \quad (2)
\]

\[
\int \sin^2x \, dx \text{ equals } \frac{1}{2}(x - \sin x \cos x) \quad \text{or} \quad \frac{1}{2}x - \frac{1}{4} \sin 2x \quad \text{(plus C)}. \quad (3)
\]

**EXAMPLE 4** \( \int \cos^4x \, dx \) (\( m = 0, n = 4 \), both are even)

Changing \( \cos^2x \) to \( 1 - \sin^2x \) gets us nowhere. All exponents stay even. Substituting \( u = \sin x \) won't simplify \( \sin^4x \, dx \), without \( du \). Integrate by parts or switch to \( 2x \).

**First solution Integrate by parts.** Take \( u = \cos^3x \) and \( dv = \cos x \, dx \):

\[
\int (\cos^3x)(\cos x \, dx) = uv - \int v \, du = \cos^3x \sin x - \int (\sin x)(-3 \cos^2x \sin x \, dx).
\]

The last integral has even powers \( \sin^2x \) and \( \cos^2x \). This looks like no progress. Replacing \( \sin^2x \) by \( 1 - \cos^2x \) produces \( \cos^4x \) on the right-hand side also:

\[
\int \cos^4x \, dx = \cos^3x \sin x + 3 \int \cos^2x(1 - \cos^2x) \, dx.
\]
Always even powers in the integrals. But now move $3 \int \cos^4 x \, dx$ to the left side:

**Reduction**

$$4 \int \cos^4 x \, dx = \cos^3 x \sin x + 3 \int \cos^2 x \, dx. \quad (4)$$

Partial success—the problem is reduced from $\cos^4 x$ to $\cos^2 x$. Still an even power, but a lower power. The integral of $\cos^2 x$ is already known. Use it in equation (4):

$$\int \cos^4 x \, dx = \frac{1}{4} \cos^3 x \sin x + \frac{3}{2} \left( x + \sin x \cos x \right) + C. \quad (5)$$

**Second solution** Substitute the double-angle formula $\cos^2 x = \frac{1}{2} + \frac{1}{2} \cos 2x$:

$$\int \cos^4 x \, dx = \int \left( \frac{1}{2} + \frac{1}{2} \cos 2x \right)^2 \, dx = \frac{1}{4} \int \left( 1 + 2 \cos 2x + \cos^2 2x \right) \, dx.$$

Certainly $\int dx = x$. Also $2 \int \cos 2x \, dx = \sin 2x$. That leaves the cosine squared:

$$\int \cos^2 2x \, dx = \int \frac{1}{2} (1 + \cos 4x) \, dx = \frac{1}{8} x + \frac{1}{8} \sin 4x + C.$$

The integral of $\cos^4 x$ using double angles is

$$\frac{1}{4} \left[ x + \sin 2x + \frac{1}{2} x + \frac{1}{4} \sin 4x \right] + C. \quad (6)$$

That solution looks different from equation (5), but it can't be. There all angles were $x$, here we have $2x$ and $4x$. We went from $\cos^4 x$ to $\cos^2 2x$ to $\cos 4x$, which was integrated immediately. The powers were cut in half as the angle was doubled.

**Double-angle method for $\int \sin^m x \cos^n x \, dx$, when $m$ and $n$ are even.**

Replace $\sin^2 x$ by $\frac{1}{2} (1 - \cos 2x)$ and $\cos^2 x$ by $\frac{1}{2} (1 + \cos 2x)$. The exponents drop to $m/2$ and $n/2$. If those are even, repeat the idea ($2x$ goes to $4x$). If $m/2$ or $n/2$ is odd, switch to the “general method” using substitution. With an odd power, we have $du$.

**EXAMPLE 5** (Double angle) $\int \sin^2 x \cos^2 x \, dx = \int \frac{1}{2} (1 - \cos 2x)(1 + \cos 2x) \, dx$.

This leaves $1 - \cos^2 2x$ in the last integral. That is familiar but not necessarily easy. We can look it up (safest) or remember it (quickest) or use double angles again:

$$\frac{1}{4} \int (1 - \cos^2 2x) \, dx = \frac{1}{4} \int \left( 1 - \frac{1}{2} - \frac{1}{2} \cos 4x \right) \, dx = \frac{x}{8} - \frac{\sin 4x}{32} + C.$$

**Conclusion** Every $\sin^m x \cos^n x$ can be integrated. This includes negative $m$ and $n$—see tangents and secants below. Symbolic codes like MACSYMA or *Mathematica* give the answer directly. Do they use double angles or integration by parts?

You may prefer the answer from integration by parts (I usually do). It avoids $2x$ and $4x$. But it makes no sense to go through every step every time. Either a computer does the algebra, or we use a “reduction formula” from $n$ to $n - 2$:

**Reduction**

$$n \int \cos^n x \, dx = \cos^{n-1} x \sin x + (n - 1) \int \cos^{n-2} x \, dx. \quad (7)$$

For $n = 2$ this is $\int \cos^2 x \, dx$—the integral to learn. For $n = 4$ the reduction produces $\cos^2 x$. The integral of $\cos^4 x$ goes to $\cos^2 x$. There are similar reduction formulas for $\sin^n x$ and also for $\sin^n x \cos^n x$. I don't see a good reason to memorize them.

**INTEGRALS WITH ANGLES $\pi x$ AND $\pi x$**

Instead of $\sin^n x$ times $\cos^m x$, suppose you have $\sin 8x$ times $\cos 6x$. How do you integrate? Separately a sine and cosine are easy. The new question is the integral of the product:
EXAMPLE 6 Find \( \int_0^{2\pi} \sin 8x \cos 6x \, dx \). More generally find \( \int_0^{2\pi} \sin px \cos qx \, dx \).

This is not for the sake of making up new problems. I believe these are the most important definite integrals in this chapter (\( p \) and \( q \) are 0, 1, 2, \( \ldots \)). They may be the most important in all of mathematics, especially because the question has such a beautiful answer. The integrals are zero. On that fact rests the success of Fourier series, and the whole industry of signal processing.

One approach (the slow way) is to replace \( \sin 8x \) and \( \cos 6x \) by powers of cosines. That involves \( \cos^{14}x \). The integration is not fun. A better approach, which applies to all angles \( px \) and \( qx \), is to use the identity

\[
\sin px \cos qx = \frac{1}{2} [\sin((p + q)x) + \sin((p - q)x)] .
\]

Thus \( \sin 8x \cos 6x = \frac{1}{2} \sin 14x + \frac{1}{2} \sin 2x \). Separated like that, sines are easy to integrate:

\[
\int_0^{2\pi} \sin 8x \cos 6x \, dx = \left[ -\frac{1}{2} \cos 14x + \frac{1}{2} \cos 2x \right]_0^{2\pi} = 0 .
\]

Since \( \cos 14x \) is periodic, it has the same value at 0 and \( 2\pi \). Subtraction gives zero. The same is true for \( \cos 2x \). The integral of sine times cosine is always zero over a complete period (like 0 to \( 2\pi \)).

What about \( \sin px \sin qx \) and \( \cos px \cos qx \)? Their integrals are also zero, provided \( p \) is different from \( q \). When \( p = q \) we have a perfect square. There is no negative area to cancel the positive area. The integral of \( \cos^2 px \) or \( \sin^2 px \) is \( \pi \).

EXAMPLE 7 \( \int_0^{2\pi} \sin 8x \sin 7x \, dx = 0 \) and \( \int_0^{2\pi} \sin^2 8x \, dx = \pi \).

With two sines or two cosines (instead of sine times cosine), we go back to the addition formulas of Section 1.5. Problem 24 derives these formulas:

\[
\sin px \sin qx = -\frac{1}{2} [\cos((p + q)x) + \cos((p - q)x)] \quad \text{and} \quad \cos px \cos qx = -\frac{1}{2} [\cos((p + q)x) + \cos((p - q)x)].
\]

With \( p = 8 \) and \( q = 7 \), we get \( \cos 15x \) and \( \cos x \). Their definite integrals are zero. With \( p = 8 \) and \( q = 8 \), we get \( \cos 16x \) and \( \cos 0x \) (which is 1). Formulas (9) and (10) also give a factor \( \frac{1}{2} \). The integral of \( \frac{1}{2} \) is \( \pi \):

\[
\int_0^{2\pi} \sin 8x \sin 7x \, dx = -\frac{1}{2} \left[ \frac{1}{8} \cos 15x \right]_0^{2\pi} + \frac{1}{2} \left[ \cos x \right]_0^{2\pi} = 0 + 0
\]

\[
\int_0^{2\pi} \sin 8x \sin 8x \, dx = -\frac{1}{2} \left[ \frac{1}{8} \cos 16x \right]_0^{2\pi} + \frac{1}{2} \left[ \cos 0x \right]_0^{2\pi} = 0 + \pi
\]

The answer zero is memorable. The answer \( \pi \) appears constantly in Fourier series. No ordinary numbers are seen in these integrals. The case \( p = q = 1 \) brings back \( \int \cos^2 x \, dx = \frac{1}{2} + \frac{1}{2} \sin 2x \).

SECANTS AND TANGENTS

When we allow negative powers \( m \) and \( n \), the main fact remains true. All integrals \( \int \sin^m x \cos^n x \, dx \) can be expressed by known functions. The novelty for negative powers is that logarithms appear. That happens right at the start, for \( \sin x/\cos x \) and for \( 1/\cos x \) (tangent and secant):

\[
\int \tan x \, dx = -\int \frac{du}{u} = -\ln|\cos x| \quad \text{(here } u = \cos x) \]

\[
\int \sec x \, dx = \int \frac{du}{u} = \ln|\sec x + \tan x| \quad \text{(here } u = \sec x + \tan x) .
\]
For higher powers there is one key identity: 

\[ 1 + \tan^2 x = \sec^2 x. \]

That is the old identity

\[ \cos^2 x + \sin^2 x = 1 \]

in disguise (just divide by \( \cos^2 x \)). We switch tangents to secants just as we switched sines to cosines. Since \((\tan x)' = \sec^2 x\) and \((\sec x)' = \sec x \tan x\), nothing else comes in.

**EXAMPLE 8** 

\[ \int \tan^2 x \, dx = \int (\sec^2 x - 1) \, dx = \tan x - x + C. \]

**EXAMPLE 9** 

\[ \int \tan^3 x \, dx = \int \tan x (\sec^2 x - 1) \, dx. \]

The first integral on the right is \( \int u \, du = \frac{1}{2} u^2 \), with \( u = \tan x \). The last integral is \(- \int \tan x \, dx\). The complete answer is \( \frac{1}{2} (\tan x)^2 + \ln |\cos x| + C \). By taking absolute values, a negative cosine is also allowed. Avoid \( \cos x = 0 \).

**EXAMPLE 10** *Reduction*

\[ \int (\tan x)^m \, dx = \frac{(\tan x)^{m-1}}{m-1} - \int (\tan x)^{m-2} \, dx \]

Same idea—separate off \((\tan x)^2\) as \( \sec^2 x - 1 \). Then integrate \((\tan x)^{m-2} \sec^2 x \, dx\), which is \( u^{m-2} \, du \). This leaves the integral on the right, with the exponent lowered by \( 2 \). Every power \((\tan x)^n\) is eventually reduced to Example 8 or 9.

**EXAMPLE 11** 

\[ \int \sec^3 x \, dx = vu - \int v \, du = \sec x \tan x - \int \tan^2 x \sec x \, dx \]

This was integration by parts, with \( u = \sec x \) and \( v = \tan x \). In the integral on the right, replace \( \tan^2 x \) by \( \sec^2 x - 1 \) (this identity is basic):

\[ \int \sec^3 x \, dx = \sec x \tan x - \int \sec^3 x \, dx + \int \sec x \, dx. \]

Bring \( \int \sec^3 x \, dx \) to the left side. That reduces the problem from \( \sec^3 x \) to \( \sec x \).

I believe those examples make the point—*trigonometric integrals are computable.* Every product \( \tan^m x \sec^n x \) can be reduced to one of these examples. If \( n \) is even we substitute \( u = \tan x \). If \( m \) is odd we set \( u = \sec x \). If \( m \) is even and \( n \) is odd, use a reduction formula (and always use \( \tan^2 x = \sec^2 x - 1 \)).

I mention very briefly a completely different substitution \( u = \tan \frac{x}{2} \). This seems to all students and instructors (quite correctly) to come out of the blue:

\[
\begin{align*}
\sin x &= \frac{2u}{1 + u^2} \\
\cos x &= \frac{1 - u^2}{1 + u^2} \\
dx &= \frac{2\,du}{1 + u^2}.
\end{align*}
\]

The \( x \)-integral can involve sums as well as products—not only \( \sin^m x \cos^n x \) but also \( 1/(5 + \sin x - \tan x) \). (No square roots.) The \( u \)-integral is a *ratio of ordinary polynomials*. It is done by partial fractions.

*Application of \( \int \sec x \, dx \) to distance on a map* (Mercator projection)

The strange integral \( \ln(\sec x + \tan x) \) has an everyday application. It measures the distance from the equator to latitude \( x \), on a Mercator map of the world.

All mapmakers face the impossibility of putting part of a sphere onto a flat page. You can't preserve distances, when an orange peel is flattened. But angles can be preserved, and Mercator found a way to do it. His map came before Newton and Leibniz. Amazingly, and accidentally, somebody matched distances on the map with a table of logarithms—and discovered \( \int \sec x \, dx \) before calculus. You would not be surprised to meet \( \sin x \), but who would recognize \( \ln(\sec x + \tan x) \)?

The map starts with strips at all latitudes \( x \). The heights are \( dx \), the lengths are proportional to \( \cos x \). *We stretch the strips by \( 1/\cos x \)—then Figure 7.3c lines up evenly on the page.* When \( dx \) is also divided by \( \cos x \), angles are preserved—a small
Fig. 7.3 Strips at latitude $x$ are scaled by $\sec x$, making Greenland too large.

To integrate $\sin^4 x \cos^3 x$, replace $\cos^2 x$ by $\frac{a}{\cos^2 x}$. Then $(\sin^4 x - \sin^6 x) \cos x \, dx$ is $\frac{b}{\cos^4 x} \, du$. In terms of $u = \sin x$ the integral is $\frac{c}{\cos^2 x}$. This work for $\sin^m x \cos^n x$ if either $m$ or $n$ is $\frac{d}{\cos^2 x}$. If both $m$ and $n$ are $\frac{e}{\cos^2 x}$, one method is integration by $\frac{f}{\cos^2 x}$. For $\int \sin^4 x \, dx$, split off $du = \sin x \, dx$. Then $\int \cos^n x \, dx$ is $\frac{g}{\cos^2 x}$. Replacing $\cos^2 x$ by $\frac{h}{\cos^2 x}$ creates a new $\sin^4 x \, dx$ that combines with the original one. The result is a reduction to $\int \sin^2 x \, dx$, which is known to equal $\frac{i}{\cos^2 x}$.

The second method uses the double-angle formula $\sin^2 x = \frac{l}{\cos^2 x}$. Then $\sin^4 x$ involves $\cos^2 x = \frac{m}{\cos^2 x}$. Another doubling comes from $\cos^2 2x = \frac{n}{\cos^2 x}$. The integral contains the sine of $\frac{o}{\cos^2 x}$.

To integrate $6 \sin^6 x \cos^4 x$, rewrite it as $\frac{p}{\cos^4 x} \sin 10x + \frac{q}{\cos^4 x}$. The indefinite integral is $\frac{r}{\cos^4 x}$. The definite integral from $0$ to $2\pi$ is $\frac{s}{\cos^4 x}$. The product $\cos px \cos qx$ is written as $\frac{t}{\cos (p + q)x + \frac{u}{\cos (p + q)x}}$. Its integral is also zero, except if $\frac{v}{\cos (p + q)x}$ when the answer is $\frac{w}{\cos (p + q)x}$.

With $u = \tan x$, the integral of $\tan^6 x \sec^2 x$ is $\frac{x}{\cos^2 x}$. Similarly $\int \sec^8 x \, dx$ is $\frac{x}{\cos^2 x}$. For the combination $\tan^m x \sec^2 x$ we apply the identity $\tan^2 x = \frac{v}{\cos^2 x}$. After reduction we may need $\int \tan x \, dx = \frac{w}{\cos^2 x}$ and $\int \sec x \, dx = \frac{x}{\cos^2 x}$.

Compute 1–8 by the "general method," when $m$ or $n$ is odd.

1 $\int \sin^3 x \, dx$  
2 $\int \cos^3 x \, dx$
3 $\int \sin x \cos x \, dx$  
4 $\int \cos x \, dx$
5 $\int \sin^2 x \cos^2 x \, dx$  
6 $\int \sin^3 x \cos^3 x \, dx$
7 $\int \sqrt{\sin x} \cos x \, dx$  
8 $\int \sqrt{\sin x} \cos^3 x \, dx$
9 Repeat Problem 6 starting with $\sin x \cos x = \frac{1}{2} \sin 2x$.
10 Find $\int \sin^2 x \cos x \, dx$ and $\int \sin x \cos ax \, dx$.

In 11–16 use the double-angle formulas ($m, n$ even).
11 $\int \sin^2 x \, dx$  
12 $\int \sin^4 x \, dx$
13 $\int \cos^2 x \, dx$  
14 $\int \sin^2 x \cos^2 x \, dx$
15 $\int \sin^2 x \, dx + \int \cos^2 x \, dx$  
16 $\int \sin^2 x \cos^2 x \, dx$

17 Use the reduction formula (7) to integrate $\cos^6 x$.
18 For $n > 1$ use formula (7) to prove
$$\int_0^{\pi/2} \cos^n x \, dx = \frac{n - 1}{n} \int_0^{\pi/2} \cos^{n-2} x \, dx.$$  
19 For $n = 2, 4, 6, \ldots$ deduce from Problem 18 that
$$\int_0^{\pi/2} \cos^n x \, dx = \frac{\pi (1)(3)(n-1)}{2^{\frac{n}{2}} (2)(4) \cdots (n)}.$$  
20 For $n = 3, 5, 7, \ldots$ deduce from Problem 18 that
$$\int_0^{\pi/2} \cos^n x \, dx = \frac{(2)(4) \cdots (n-1)}{(3)(5) \cdots (n)}.$$  
21 (a) Separate $dv = \sin x \, dx$ from $u = \sin^{n-1} x$ and integrate $\int \sin^n x \, dx$ by parts.
(b) Substitute $1 - \sin^2 x$ for $\cos^2 x$ to find a reduction formula like equation (7).
22 For which $n$ does symmetry give $\int_0^\pi \cos^2 x \, dx = 0$?
23 Are the integrals (a)–(f) positive, negative, or zero?
(a) $\int_0^\pi \cos 3x \sin 3x \, dx$  
(b) $\int_0^\pi \cos x \sin 2x \, dx$
(c) $\int_0^\pi \cos x \sin x \, dx$  
(d) $\int_0^\pi (\cos^2 x - \sin^2 x) \, dx$
(e) $\int_0^\pi \sin x \cos x \sin x \, dx$  
(f) $\int_0^\pi \cos^2 x \, dx$
294

7 Techniques of Integration

24 Write down equation (9) for \( p = q = 1 \), and (10) for \( p = 2, q = 1 \). Derive (9) from the addition formulas for \( \cos(s + t) \) and \( \cos(s - t) \) in Section 1.5.

In 25–32 compute the indefinite integrals first, then the definite integrals.

25 \( \int_0^{2x} \cos x \sin 2x \, dx \)  
26 \( \int_0^x \sin 3x \sin 5x \, dx \)  
27 \( \int_0^x \cos 99x \sin 101x \, dx \)  
28 \( \int_0^x \cos 3x \, dx \)  
29 \( \int_0^x \cos 99x \cos 101x \, dx \)  
30 \( \int_0^x \sin x \sin 2x \sin 3x \, dx \)  
31 \( \int_0^x \cos x/2 \sin x/2 \, dx \)  
32 \( \int_0^x \cos x \cos 3x \, dx \) (by parts)

33 Suppose a Fourier sine series \( A \sin x + B \sin 2x + C \sin 3x + \cdots \) adds up to \( x \) on the interval from 0 to \( \pi \). Find \( A \) by multiplying all those functions (including \( x \)) by \( \sin x \) and integrating from 0 to \( \pi \). (\( B \) and \( C \) will disappear.)

34 Suppose a Fourier sine series \( A \sin x + B \sin 2x + C \sin 3x + \cdots \) adds up to 1 on the interval from 0 to \( \pi \). Find \( C \) by multiplying all functions (including 1) by \( \sin 3x \) and integrating from 0 to \( \pi \). (\( A \) and \( B \) will disappear.)

35 In 33, the series also equals \( x \) from \( -\pi \) to 0, because all functions are odd. Sketch the “sawtooth function,” which equals \( x \) from \( -\pi \) to \( \pi \) and then has period \( 2\pi \). What is the sum of the sine series at \( x = \pi \)?

36 In 34, the series equals \( -1 \) from \( -\pi \) to 0, because sines are odd functions. Sketch the “square wave,” which is alternately \( -1 \) and \( +1 \), and find \( A \) and \( B \).

37 The area under \( y = \sin x \) from 0 to \( \pi \) is positive. Which frequencies \( p \) have \( \int_0^\pi \sin px \, dx = 0 \)?

38 Which frequencies \( q \) have \( \int_0^\pi \cos qx \, dx = 0 \)?

39 For which \( p, q \) is \( \int_0^\pi \sin px \cos qx \, dx = 0 \)?

40 Show that \( \int_0^\pi \sin px \cos qx \, dx \) is always zero.

Compute the indefinite integrals 41–52.

41 \( \int \sec x \tan x \, dx \)  
42 \( \int \tan 5x \, dx \)  
43 \( \int \tan^2 x \sec^2 x \, dx \)  
44 \( \int \tan^2 x \sec x \, dx \)  
45 \( \int \tan x \sec^3 x \, dx \)  
46 \( \int \sec^4 x \, dx \)  
47 \( \int \tan^4 x \, dx \)  
48 \( \int \tan^3 x \, dx \)  
49 \( \int \cot x \, dx \)  
50 \( \int \csc x \, dx \)  
51 \( \int \frac{\sin x}{\cos^3 x} \, dx \)  
52 \( \int \frac{\sin^5 x}{\cos^3 x} \, dx \)

53 Choose \( A \) so that \( \cos x - \sin x = A \cos(x - \pi/4) \). Then integrate \( 1/\cos x - \sin x \).

54 Choose \( A \) so that \( \cos x - \sqrt{3} \sin x = A \cos(x + \pi/3) \). Then integrate \( 1/\cos x - \sqrt{3} \sin x \).

55 Evaluate \( \int_0^\pi \cos x - \sin x \, dx \).

56 Show that \( a \cos x + b \sin x = \sqrt{a^2 + b^2} \cos (x - \alpha) \) and find the correct phase angle \( \alpha \).

57 If a square Mercator map shows 1000 miles at latitude 30°, how many miles does it show at latitude 60°?

58 When lengths are scaled by \( \sec x \), area is scaled by \( \tan x \). Why is the area from the equator to latitude \( x \) proportional to \( \tan x \)?

59 Use substitution (11) to find \( \int dx/(1 + \cos x) \).

60 Explain from areas why \( \int_0^\pi \sin^2 x \, dx = \int_0^\pi \cos^2 x \, dx \). These integrals add to \( \int_0^\pi 1 \, dx \), so they both equal \( \frac{\pi}{2} \).

61 What product \( \sin px \sin qx \) is graphed below? Check that \( (p \cos px \sin qx - q \sin px \cos qx)/(q^2 - p^2) \) has this derivative.

62 Finish \( \int \sec^3 x \, dx \) in Example 11. This is needed for the length of a parabola and a spiral (Problem 7.3.8 and Sections 8.2 and 9.3).

The most powerful tool we have, for integrating with pencil and paper and brain, is the method of substitution. To make it work, we have to think of good substitutions—which make the integral simpler. This section concentrates on the single most valuable collection of substitutions. They are the only ones you should memorize, and two examples are given immediately.
To integrate $\sqrt{1-x^2}$, substitute $x = \sin \theta$. Do not set $u = 1-x^2 \left( \frac{du}{dx} \right)$ is missing.

$$\int \sqrt{1-x^2} \, dx \to \int (\cos \theta)(\cos \theta \, d\theta) \quad \int \frac{dx}{\sqrt{1-x^2}} \to \int \frac{\cos \theta \, d\theta}{\cos \theta}$$

The expression $\sqrt{1-x^2}$ is awkward as a function of $x$. It becomes graceful as a function of $\theta$. We are practically invited to use the equation $1 - (\sin \theta)^2 = (\cos \theta)^2$.

Then the square root is simply $\cos \theta$—provided this cosine is positive.

Notice the change in $dx$. When $x$ is $\sin \theta$, $dx$ is $\cos \theta \, d\theta$. Figure 7.4a shows the original area with new letters. Figure 7.4b shows an equal area, after rewriting $\int (\cos \theta)(\cos \theta \, d\theta)$ as $\int (\cos^2 \theta) \, d\theta$. Changing from $x$ to $\theta$ gives a new height and a new base. There is no change in area—that is the point of substitution.

To put it bluntly: If we go from $\sqrt{1-x^2}$ to $\cos \theta$, and forget the difference between $dx$ and $d\theta$, and just compute $\int \cos \theta \, d\theta$, the answer is totally wrong.

![Fig. 7.4 Same area for $\sqrt{1-x^2} \, dx$ and $\cos^2 \theta \, d\theta$. Third area is wrong: $dx \neq d\theta$](image)

We still need the integral of $\cos^2 \theta$. This was Example 3 of integration by parts, and also equation 7.2.6. It is worth memorizing. The example shows this $\theta$ integral, and returns to $x$:

**EXAMPLE 1**

$$\int \cos^2 \theta \, d\theta = \frac{1}{2} \sin \theta \cos \theta + \frac{1}{2} \theta$$ is after substitution

$$\int \sqrt{1-x^2} \, dx = \frac{1}{2} x \sqrt{1-x^2} + \frac{1}{2} \sin^{-1} x$$ is the original problem.

We changed $\sin \theta$ back to $x$ and $\cos \theta$ to $\sqrt{1-x^2}$. Notice that $\theta$ is $\sin^{-1} x$. The answer is trickier than you might expect for the area under a circular arc. Figure 7.5 shows how the two pieces of the integral are the areas of a pie-shaped wedge and a triangle.

**EXAMPLE 2**

$$\int \frac{dx}{\sqrt{1-x^2}} = \int \frac{\cos \theta \, d\theta}{\cos \theta} = \theta + C = \sin^{-1} x + C.$$  

Remember: We already know $\sin^{-1} x$. Its derivative $1/\sqrt{1-x^2}$ was computed in Section 4.4. That solves the example. But instead of matching this special problem

![Fig. 7.5 $\int \sqrt{1-x^2} \, dx$ is a sum of simpler areas. Infinite graph but finite area.](image)
with a memory from Chapter 4, the substitution \( x = \sin \theta \) makes the solution automatic. From \( \int d\theta = \theta \) we go back to \( \sin^{-1}x \).

The rest of this section is about other substitutions. They are more complicated than \( x = \sin \theta \) (but closely related). A table will display the three main choices—\( \sin \theta \), \( \tan \theta \), \( \sec \theta \)—and their uses.

**TRIGONOMETRIC SUBSTITUTIONS**

After working with \( \sqrt{1 - x^2} \), the next step is \( \sqrt{4 - x^2} \). The change \( x = \sin \theta \) simplified the first, but it does nothing for the second: \( 4 - \sin^2 \theta \) is not familiar. Nevertheless a factor of 2 makes everything work. Instead of \( x = \sin \theta \), the idea is to substitute \( x = 2 \sin \theta \):

\[
\sqrt{4 - x^2} = \sqrt{4 - 4 \sin^2 \theta} = 2 \cos \theta \quad \text{and} \quad dx = 2 \cos \theta \, d\theta.
\]

Notice both 2's. The integral is \( 4 \int \cos^2 \theta \, d\theta = 2 \sin \theta + 2 \theta \). But watch closely. This is not 4 times the previous \( \int \cos^2 \theta \, d\theta \)! Since \( x \) is \( 2 \sin \theta \), \( \theta \) is now \( \sin^{-1}\left(\frac{x}{2}\right) \).

**EXAMPLE 3**

\[
\int \sqrt{4 - x^2} \, dx = 4 \int \cos^2 \theta \, d\theta = x\sqrt{1 - (x/2)^2} + 2 \sin^{-1}(x/2).
\]

Based on \( \sqrt{1 - x^2} \) and \( \sqrt{4 - x^2} \), here is the general rule for \( \sqrt{a^2 - x^2} \). Substitute \( x = a \sin \theta \). Then the a's separate out:

\[
\sqrt{a^2 - x^2} = \sqrt{a^2 - a^2 \sin^2 \theta} = a \cos \theta \quad \text{and} \quad dx = a \cos \theta \, d\theta.
\]

That is the automatic substitution to try, whenever the square root appears.

**EXAMPLE 4**

\[
\int_{x=0}^{4} d\frac{x}{\sqrt{16 - x^2}} = \int_{\theta=0}^{\pi/2} \frac{4 \cos \theta \, d\theta}{\sqrt{4^2 - 4^2 \sin^2 \theta}} = \int_{\theta=0}^{\pi/2} \frac{4 \cos \theta \, d\theta}{\frac{8}{2}} = \frac{\pi}{2}.
\]

Here \( a^2 = 16 \). Then \( a = 4 \) and \( x = 4 \sin \theta \). The integral has \( 4 \cos \theta \) above and below, so it is \( \int d\theta \). The antiderivative is just \( \theta \). For the definite integral notice that \( x = 4 \) means \( \sin \theta = 1 \), and this means \( \theta = \pi/2 \).

A table of integrals would hide that substitution. The table only gives \( \sin^{-1}(x/4) \). There is no mention of \( \int d\theta = \theta \). But what if \( 16 - x^2 \) changes to \( x^2 - 16 \)?

**EXAMPLE 5**

\[
\int_{x=4}^{8} d\frac{x}{\sqrt{x^2 - 16}} = ?
\]

Notice the two changes—the sign in the square root and the limits on \( x \). Example 4 stayed inside the interval \( |x| \leq 4 \), where \( 16 - x^2 \) has a square root. Example 5 stays outside, where \( x^2 - 16 \) has a square root. The new problem cannot use \( x = 4 \sin \theta \), because we don’t want the square root of \( -\cos^2 \theta \).

The new substitution is \( x = 4 \sec \theta \). This turns the square root into \( 4 \tan \theta \):

\[
x = 4 \sec \theta \quad \text{gives} \quad dx = 4 \sec \theta \tan \theta \, d\theta \quad \text{and} \quad x^2 - 16 = 16 \sec^2 \theta - 16 = 16 \tan^2 \theta.
\]

This substitution solves the example, when the limits are changed to \( \theta \):

\[
\int_{\theta=0}^{\pi/3} 4 \sec \theta \tan \theta \, d\theta = \int_{\theta=0}^{\pi/3} \sec \theta \, d\theta = \ln(\sec \theta + \tan \theta) \bigg|_{0}^{\pi/3} = \ln(2 + \sqrt{3}).
\]

I want to emphasize the three steps. First came the substitution \( x = 4 \sec \theta \). An unrecognizable integral became \( \int \sec \theta \, d\theta \). Second came the new limits (\( \theta = 0 \) when \( x = 4 \), \( \theta = \pi/3 \) when \( x = 8 \)). Then I integrated \( \sec \theta \).
Example 6 has the same $x^2 - 16$. So the substitution is again $x = 4 \sec \theta$:

**EXAMPLE 6** \[ \int_{x=8}^{\infty} \frac{16 \, dx}{(x^2 - 16)^{3/2}} = \int_{\theta=\pi/3}^{\pi/2} \frac{64 \sec \theta \tan \theta \, d\theta}{(4 \tan \theta)^3} = \int_{\pi/3}^{\pi/2} \frac{16 \cos \theta \, d\theta}{\sin^2 \theta}. \]

Step one substitutes $x = 4 \sec \theta$. Step two changes the limits to $\theta$. The upper limit $x = \infty$ becomes $\theta = \pi/2$, where the secant is infinite. The limit $x = 8$ again means $\theta = \pi/3$. To get a grip on the integral, I also changed to sines and cosines.

The integral of $\cos \theta/\sin^2 \theta$ needs another substitution! (Or else recognize $\cot \theta \csc \theta$.) With $u = \sin \theta$ we have $\int \frac{du}{u^2} = -\frac{1}{u} = -\frac{1}{\sin \theta}$:

**Solution** \[ \int_{\pi/3}^{\pi/2} \frac{16 \cos \theta \, d\theta}{\sin^2 \theta} = \left. \left( -\frac{1}{\sin \theta} \right) \right|_{\pi/3}^{\pi/2} = -1 + \frac{2}{\sqrt{3}}. \]

**Warning** With lower limit $\theta = 0$ (or $x = 4$) this integral would be a disaster. It divides by $\sin 0$, which is zero. This area is infinite.

(Warning) Example 5 also blew up at $x = 4$, but the area was **not** infinite. To make the point directly, compare $x^{-1/2}$ to $x^{-3/2}$. Both blow up at $x = 0$, but the first one has finite area:

\[ \int_0^1 \frac{1}{\sqrt{x}} \, dx = 2 \int_0^1 x^{-1/2} \, dx = 2 \left( \frac{-2}{-2/2} \right) = \infty. \]

Section 7.5 separates finite areas (slow growth of $1/\sqrt{x}$) from infinite areas (fast growth of $x^{-3/2}$).

**Last substitution** Together with $16 - x^2$ and $x^2 - 16$ comes the possibility $16 + x^2$. (You might ask about $-16 - x^2$, but for obvious reasons we don't take its square root.) This third form $16 + x^2$ requires a third substitution $x = 4 \tan \theta$. Then $16 + x^2 = 16 + 16 \tan^2 \theta = 16 \sec^2 \theta$. Here is an example:

**EXAMPLE 7** \[ \int_{x=0}^{\infty} \frac{dx}{16 + x^2} = \int_{\theta=0}^{\pi/2} \frac{4 \sec^2 \theta \, d\theta}{16 \sec^2 \theta} = \frac{1}{4} \int_{\theta=0}^{\pi/2} \frac{\pi}{2}. \]

**Table of substitutions** for $a^2 - x^2$, $a^2 + x^2$, $x^2 - a^2$

- $x = a \sin \theta$ replaces $a^2 - x^2$ by $a^2 \cos^2 \theta$ and $dx$ by $a \cos \theta \, d\theta$
- $x = a \tan \theta$ replaces $a^2 + x^2$ by $a^2 \sec^2 \theta$ and $dx$ by $a \sec^2 \theta \, d\theta$
- $x = a \sec \theta$ replaces $x^2 - a^2$ by $a^2 \tan^2 \theta$ and $dx$ by $a \sec \theta \tan \theta \, d\theta$

**Note** There is a subtle difference between changing $x$ to $\sin \theta$ and changing $\sin \theta$ to $u$:

- In Example 1, $dx$ was replaced by $\cos \theta \, d\theta$ (new method)
- In Example 6, $\cos \theta \, d\theta$ was already there and became $du$ (old method).

The combination $\cos \theta \, d\theta$ was put into the first and pulled out of the second.

My point is that Chapter 5 needed $du/dx$ inside the integral. Then $(du/dx)dx$ became $du$. **Now it is not necessary to see so far ahead.** We can try any substitution. If it works, we win. In this section, $x = \sin \theta$ or $\sec \theta$ or $\tan \theta$ is bound to succeed.

**NEW** \[ \int \frac{dx}{1 + x^2} = \int d\theta \text{ by trying } x = \tan \theta \]

**OLD** \[ \int \frac{x \, dx}{1 + x^2} = \int \frac{du}{2u} \text{ by seeing } du \]
We mention the hyperbolic substitutions \( \tanh \theta, \sinh \theta, \) and \( \cosh \theta \). The table below shows their use. They give new forms for the same integrals. If you are familiar with hyperbolic functions the new form might look simpler—as it does in Example 8.

\[
\begin{align*}
\sinh \theta \, d\theta &= \cosh \theta + C \\
\cosh \theta \, d\theta &= \cosh \theta + C \\
\tanh \theta \, d\theta &= \ln \left| \cosh \theta + 1 \right| + C.
\end{align*}
\]

**EXAMPLE 8**

\[
\int \frac{dx}{\sqrt{x^2 - 1}} = \int \frac{\sinh \theta \, d\theta}{\sinh \theta} = \cosh^{-1} x + C.
\]

This way looks harder, but most tables prefer that final logarithm. It is clearer than \( \cosh^{-1} x \), even if it takes more space. All answers agree if Problem 35 is correct.

**COMPLETING THE SQUARE**

We have not said what to do for \( \sqrt{x^2 - 2x + 2} \) or \( \sqrt{-x^2 + 2x} \). Those square roots contain a linear term—a multiple of \( x \). The device for removing linear terms is worth knowing. It is called completing the square, and two examples will begin to explain it:

\[
\begin{align*}
x^2 - 2x + 2 &= (x - 1)^2 + 1 = u^2 + 1 \\
x^2 + 2x &= -(x - 1)^2 + 1 = 1 - u^2.
\end{align*}
\]

The idea has three steps. First, get the \( x^2 \) and \( x \) terms into one square. Here that square was \( (x - 1)^2 = x^2 - 2x + 1 \). Second, fix up the constant term. Here we recover the original functions by adding 1. Third, set \( u = x - 1 \) to leave no linear term. Then the integral goes forward based on the substitutions of this section:

\[
\int \frac{dx}{\sqrt{x^2 - 2x + 2}} = \int \frac{du}{\sqrt{u^2 + 1}} \\
\int \frac{dx}{\sqrt{x^2 - x^2}} = \int \frac{du}{\sqrt{1 - u^2}}
\]

The same idea applies to any quadratic that contains a linear term \( 2bx \):

\[
\begin{align*}
\text{rewrite} \quad x^2 + 2bx + c &= (x + b)^2 + C, \text{ with } C = c - b^2 \\
\text{rewrite} \quad x^2 + 2bx + c &= -(x - b)^2 + C, \text{ with } C = c + b^2
\end{align*}
\]

To match the quadratic with the square, we fix up the constant:

\[
\begin{align*}
x^2 + 10x + 16 &= (x + 5)^2 + C \text{ leads to } C = 16 - 25 = -9 \\
-x^2 + 10x + 16 &= -(x - 5)^2 + C \text{ leads to } C = 16 + 25 = 41.
\end{align*}
\]

**EXAMPLE 9**

\[
\int \frac{dx}{x^2 + 10x + 16} = \int \frac{dx}{(x - 5)^2 - 9} = \int \frac{du}{u^2 - 9}.
\]

Here \( u = x + 5 \) and \( du = dx \). Now comes a choice—struggle on with \( u = 3 \sec \theta \) or look for \( \int \frac{du}{(u^2 - a^2)} \) inside the front cover. Then set \( a = 3 \):

\[
\int \frac{du}{u^2 - 9} = \frac{1}{6} \ln \left| \frac{u - 3}{u + 3} \right| = \frac{1}{6} \ln \left| \frac{x + 2}{x + 8} \right|.
\]
Note If the quadratic starts with $5x^2$ or $-5x^2$, factor out the 5 first:

$$5x^2 - 10x + 25 = 5(x^2 - 2x + 5) = 5[(x - 1)^2 + 4].$$

Now $u = x - 1$ produces $5[u^2 + 4]$. This is ready for table lookup or $u = 2 \tan \theta$.

EXAMPLE 10

$$\int \frac{dx}{5x^2 - 10x + 25} - \int \frac{du}{5[u^2 + 4]} = \int \frac{2\sec^2 \theta \, d\theta}{5[4 \sec^2 \theta]} = \frac{1}{10} \int d\theta.$$

This answer is $\theta/10 + C$. Now go backwards: $\theta/10 = (\tan^{-1} \frac{1}{2}u)/10 = (\tan^{-1} \frac{1}{2}(x - 1))/10$. Nobody could see that from the start. A double substitution takes practice, from $x$ to $u$ to $\theta$. Then go backwards from $\theta$ to $u$ to $x$.

**Final remark** For $u^2 + a^2$ we substitute $u = a \tan \theta$. For $u^2 - a^2$ we substitute $u = a \sec \theta$. This big dividing line depends on whether the constant $C$ (after completing the square) is positive or negative. We either have $C = a^2$ or $C = -a^2$. The same dividing line in the original $x^2 + 2bx + c$ is between $c > b^2$ and $c < b^2$. In between, $c = b^2$ yields the perfect square $(x + b)^2$—and no trigonometric substitution at all.

### 7.3 EXERCISES

**Read-through questions**

The function $\sqrt{1 - x^2}$ suggests the substitution $x = \sin \theta$. The square root becomes $\frac{1}{\sqrt{1 - x^2}}$ and $dx$ changes to $\frac{d\theta}{\sqrt{1 - \sin^2 \theta}}$. The integral $\int (1 - x^2)^{3/2} \, dx$ becomes $\int \frac{d\theta}{\sqrt{1 - \sin^2 \theta}}$. The interval $\frac{\pi}{2} < x < 1$ changes to $0 < \theta < \frac{\pi}{2}$.

For $\sqrt{a^2 - x^2}$ the substitution is $x = a \sin \theta$ with $dx = a \cos \theta \, d\theta$. For $x^2 - a^2$ we use $x = \sqrt{a^2 - 1}$ with $dx = \frac{1}{\sqrt{a^2 - 1}} \, d\theta$. Then $\int dx/(1 + x^2)$ becomes $\int d\theta$, because $1 + \tan^2 \theta = \sec^2 \theta$. The answer is $\theta - \tan^{-1} x$. We already knew that $\frac{1}{\sqrt{1 - x^2}}$ is the derivative of $\sin^{-1} x$.

The quadratic $x^2 + 2bx + c$ contains a $2b$ term $2bx$. To remove it we subtract the square. This gives $(x + b)^2 + C$ with $C = a^2$. The example $x^2 + 4x + 9$ becomes $x^2 + 1$. Then $u = x + 2$. In case $x^2$ enters with a minus sign, $-x^2 + 4x + 9$ becomes $-x^2 + 1$. When the quadratic contains $4x^2$, start by factoring out $4$.

Integrate 1–20 by substitution. Change $\theta$ back to $x$.

1. $\int \frac{dx}{\sqrt{4 - x^2}}$
2. $\int \frac{dx}{\sqrt{x^2 - a^2}}$
3. $\int \sqrt{4 - x^2} \, dx$
4. $\int \frac{dx}{1 + 9x^2}$
5. $\int \frac{x^2 \, dx}{\sqrt{1 - x^2}}$
6. $\int \frac{dx}{x^2 \sqrt{1 - x^2}}$
7. $\int \frac{dx}{(1 + x^2)^2}$
8. $\int \sqrt{x^2 + a^2} \, dx$ (see 7.2.62)
9. $\int \frac{x^2 - 25}{x} \, dx$
10. $\int \frac{x^3 \, dx}{\sqrt{9 - x^2}}$
11. $\int \frac{dx}{\sqrt{x^6 - x^2}}$
12. $\int \frac{dx}{\sqrt{x^6 - x^2}}$
13. $\int \frac{dx}{(1 + x^2)^{3/2}}$
14. $\int \frac{dx}{(1 - x^2)^{3/2}}$
15. $\int \frac{dx}{(x^2 - 9)^{3/2}}$
16. $\int \frac{dx}{(1 + x^2)^{3/2}}$
17. $\int \frac{x^2 \, dx}{\sqrt{x^2 - 1}}$
18. $\int \frac{x^2 \, dx}{\sqrt{x^2 - 1}}$
19. $\int \frac{dx}{x^2 \sqrt{x^2 + 1}}$
20. $\int \frac{x^2 \, dx}{\sqrt{x^2 + 1}}$
21. (Important) This section started with $x = \sin \theta$ and

$$\int dx \sqrt{1 - x^2} = \int d\theta = \theta = \sin^{-1} x.$$ (a) Use $x = \cos \theta$ to get a different answer.
(b) How can the same integral give two answers?

22. Compute $\int dx / x \sqrt{x^2 - 1}$ with $x = \sec \theta$. Recompute with $x = \csc \theta$. How can both answers be correct?

23. Integrate $x/(x^2 + 1)$ with $x = \tan \theta$, and also directly as a logarithm. Show that the results agree.

24. Show that $\int dx / x \sqrt{x^2 - 1} = \frac{1}{2} \sec^{-1} (x^2)$.

Calculate the definite integrals 25–32.

25. $\int_{-a}^{a} \sqrt{a^2 - x^2} \, dx = \text{area of}$
26. $\int_{-1}^{1} (1 - x^2)^{3/2} \, dx$
27. $\int_{-1}^{1} \frac{dx}{s \sqrt{1 - x^2}}$
In 37–42 substitute \( x = \sin \theta, \cosh \theta, \) or \( \tanh \theta. \) After integration change back to \( x. \)
Our plan is to split $P/Q$ into a sum like this—and integrate each piece.

Which rational function produced that particular sum? It was

$$\frac{1}{x-2} + \frac{3}{x+2} - \frac{4}{x} = \frac{(x+2)(x) + 3(x-2)(x) - 4(x-2)(x+2)}{(x-2)(x+2)(x)} = \frac{-4x + 16}{(x-2)(x+2)(x)}$$

This is $P/Q$. It is a ratio of polynomials, degree 1 over degree 3. The pieces of $P$ are collected into $-4x + 16$. The common denominator $(x-2)(x+2)(x) = x^3 - 4x$ is $Q$. But I kept these factors separate, for the following reason. When we start with $P/Q$, and break it into a sum of pieces, the first things we need are the factors of $Q$. 

In the standard problem $P/Q$ is given. To integrate it, we break it up. The goal of partial fractions is to find the pieces—to prepare for integration. That is the technique to learn in this section, and we start right away with examples.

**EXAMPLE 1** Suppose $P/Q$ has the same $Q$ but a different numerator $P$:

$$\frac{P}{Q} = \frac{3x^2 + 8x - 4}{(x-2)(x+2)(x)} = \frac{A}{x-2} + \frac{B}{x+2} + \frac{C}{x}. \tag{1}$$

Notice the form of those pieces! They are the "partial fractions" that add to $P/Q$. Each one is a constant divided by a factor of $Q$. We know the factors $x-2$ and $x+2$ and $x$. We don't know the constants $A$, $B$, $C$. In the previous case they were 1, 3, $-4$. In this and other examples, there are two ways to find them.

**Method 1** (slow) Put the right side of (1) over the common denominator $Q$:

$$\frac{3x^2 + 8x - 4}{Q} = \frac{A(x+2)(x) + B(x-2)(x) + C(x-2)(x+2)}{(x-2)(x+2)(x)} \tag{2}$$

Why is $A$ multiplied by $(x+2)(x)$? Because canceling those factors will leave $A/(x-2)$ as in equation (1). Similarly we have $B/(x+2)$ and $C/x$. Choose the numbers $A$, $B$, $C$ so that the numerators match. As soon as they agree, the splitting is correct.

**Method 2** (quicker) Multiply equation (1) by $x-2$. That leaves a space:

$$\frac{3x^2 + 8x - 4}{(x+2)(x)} = A + \frac{B(x-2)}{x+2} + \frac{C(x-2)}{x}. \tag{3}$$

Now set $x = 2$ and immediately you have $A$. The last two terms of (3) are zero, because $x-2$ is zero when $x = 2$. On the left side, $x = 2$ gives

$$\frac{3(2)^2 + 8(2) - 4}{(2+2)(2)} = \frac{24}{8} = 3 \quad \text{(which is } A).$$

Notice how multiplying by $x-2$ produced a hole on the left side. Method 2 is the "cover-up method." Cover up $x-2$ and then substitute $x = 2$. The result is $3 = A + 0 + 0$, just what we wanted.

In Method 1, the numerators of equation (2) must agree. The factors that multiply $B$ and $C$ are again zero at $x = 2$. That leads to the same $A$—but the cover-up method avoids the unnecessary step of writing down equation (2).
**Calculation of B** Multiply equation (1) by \(x + 2\), which covers up the \((x + 2)\):

\[
\frac{3x^2 + 8x - 4}{x - 2} = \frac{A(x + 2)}{x - 2} + B + \frac{C(x + 2)}{x}.
\]

(4)

Now set \(x = -2\), so \(A\) and \(C\) are multiplied by zero:

\[
\frac{3(-2)^2 + 8(-2) - 4}{(-2 - 2)} = -\frac{8}{8} = -1 = B.
\]

This is almost full speed, but (4) was not needed. *Just cover up in \(Q\) and give \(x\) the right value (which makes the covered factor zero).*

**Calculation of C (quickest)** In equation (1), cover up the factor \((x)\) and set \(x = 0\):

\[
\frac{3(0)^2 + 8(0) - 4}{(0 - 2)(0 + 2)} = -\frac{4}{-4} = 1 = C.
\]

(5)

To repeat: The same result \(A = 3\), \(B = -1\), \(C = 1\) comes from Method 1.

**EXAMPLE 2**

\[
\frac{x + 2}{(x - 1)(x + 3)} = \frac{A}{x - 1} + \frac{B}{x + 3}.
\]

First cover up \((x - 1)\) on the left and set \(x = 1\). Next cover up \((x + 3)\) and set \(x = -3\):

\[
\frac{1 + 2}{(1 - 1)(1 + 3)} = \frac{3}{4} = A \quad \frac{-3 + 2}{(-3 - 1)(-3 + 3)} = -\frac{1}{4} = B.
\]

The integral is \(\frac{3}{4}\ln|x - 1| + \frac{1}{4}\ln|x + 3| + C\).

**EXAMPLE 3** This was needed for the logistic equation in Section 6.5:

\[
\frac{1}{y(c - by)} = \frac{A}{y} + \frac{B}{c - by}.
\]

(6)

First multiply by \(y\). That covers up \(y\) in the first two terms and changes \(B\) to \(by\). Then set \(y = 0\). The equation becomes \(1/c = A\).

To find \(B\), multiply by \(c - by\). That covers up \(c - by\) in the outside terms. In the middle, \(A\) times \(c - by\) will be zero at \(y = c/b\). That leaves \(B\) on the right equal to \(1/y = b/c\) on the left. Then \(A = 1/c\) and \(B = b/c\) give the integral announced in Equation 6.5.9:

\[
\int \frac{dy}{cy - by^2} = \int \frac{dy}{cy} + \int \frac{b\ dy}{c(c - by)} = \frac{\ln y}{c} - \frac{\ln(c - by)}{c}.
\]

(7)

*It is time to admit that the general method of partial fractions can be very awkward.* First of all, it requires the factors of the denominator \(Q\). When \(Q\) is a quadratic \(ax^2 + bx + c\), we can find its roots and its factors. In theory a cubic or a quartic can also be factored, but in practice only a few are possible—for example \(x^4 - 1\) is \((x^2 - 1)(x^2 + 1)\). Even for this good example, *two of the roots are imaginary.* We can split \(x^2 - 1\) into \((x + 1)(x - 1)\). We cannot split \(x^2 + 1\) without introducing \(i\).

The method of partial fractions can work directly with \(x^2 + 1\), as we now see.

**EXAMPLE 4**

\[
\int \frac{3x^2 + 2x + 7}{x^2 + 1} \, dx \quad \text{(a quadratic over a quadratic)}.
\]

This has another difficulty. The degree of \(P\) equals the degree of \(Q\) (= 2). *Partial*
Partial Fractions cannot start until \( P \) has lower degree. Therefore I divide the leading term \( x^2 \) into the leading term \( 3x^2 \). That gives 3, which is separated off by itself:

\[
\frac{3x^2 + 2x + 7}{x^2 + 1} = 3 + \frac{2x + 4}{x^2 + 1}.
\]  (8)

Note how 3 really used \( 3x^2 + 3 \) from the original numerator. That left \( 2x + 4 \). Partial fractions will accept a linear factor \( 2x + 4 \) (or \( Ax + B \), not just \( A \)) above a quadratic.

This example contains \( 2x/(x^2 + 1) \), which integrates to \( \ln(x^2 + 1) \). The final \( 4/(x^2 + 1) \) integrates to \( 4 \tan^{-1}x \). When the denominator is \( x^2 + x + 1 \) we complete the square before integrating. The point of Sections 7.2 and 7.3 was to make that integration possible. This section gets the fraction ready—in parts.

The essential point is that we never have to go higher than quadratics. Every denominator \( Q \) can be split into linear factors and quadratic factors. There is no magic way to find those factors, and most examples begin by giving them. They go into their own fractions, and they have their own numerators—which are the \( A \) and \( B \) and \( 2x + 4 \) we have been computing.

The one remaining question is what to do if a factor is repeated. This happens in Example 5.

**Example 5**

\[
\frac{2x + 3}{(x - 1)^2} = \frac{A}{x - 1} + \frac{B}{(x - 1)^2}.
\]

The key is the new term \( B/(x - 1)^2 \). That is the right form to expect. With \((x - 1)(x - 2)\) this term would have been \( B/(x - 2) \). But when \( (x - 1) \) is repeated, something new is needed. To find \( B \), multiply through by \( (x - 1)^2 \) and set \( x = 1 \):

\[
2x + 3 = A(x - 1) + B \quad \text{becomes} \quad 5 = B \quad \text{when} \quad x = 1.
\]

This cover-up method gives \( B \). Then \( A = 2 \) is easy, and the integral is \( 2 \ln|x - 1| - 5/(x - 1) \). The fraction \( 5/(x - 1)^3 \) has an integral without logarithms.

**Example 6**

\[
\frac{2x^3 + 9x^2 + 4}{x^2(x^2 + 4)(x - 1)} = \frac{A}{x} + \frac{B}{x^2} + \frac{Cx + D}{x^2 + 4} + \frac{E}{x - 1}.
\]

This final example has almost everything! It is more of a game than a calculus problem. In fact calculus doesn't enter until we integrate (and nothing is new there). Before computing \( A, B, C, D, E \), we write down the overall rules for partial fractions:

1. The degree of \( P \) must be less than the degree of \( Q \). Otherwise divide their leading terms as in equation (8) to lower the degree of \( P \). Here \( 3 < 5 \).
2. Expect the fractions illustrated by Example 6. The linear factors \( x \) and \( x + 1 \) (and the repeated \( x^2 \)) are underneath constants. The quadratic \( x^2 + 4 \) is under a linear term. A repeated \( (x^2 + 4)^2 \) would be under a new \( Fx + G \).
3. Find the numbers \( A, B, C, \ldots \) by any means, including cover-up.
4. Integrate each term separately and add.

We could prove that this method always works. It makes better sense to show that it works once, in Example 6.

To find \( E \), cover up \( (x - 1) \) on the left and substitute \( x = 1 \). Then \( E = 3 \). To find \( B \), cover up \( x^2 \) on the left and set \( x = 0 \). Then \( B = 4(0 + 4)(0 - 1) = -1 \). The cover-up method has done its job, and there are several ways to find \( A, C, D \).
7 Techniques of Integration

Compare the numerators, after multiplying through by the common denominator $Q$:

$$2x^3 + 9x^2 + 4 = Ax(x^2 + 4)(x - 1) - (x^2 + 4)(x - 1) + (Cx + D)(x^2)(x - 1) + 3x^2(x^2 + 4).$$

The known terms on the right, from $B = -1$ and $E = 3$, can move to the left:

$$-3x^2 + 3x^2 - 4x^2 + 4x = Ax(x^2 + 4)(x - 1) + (Cx + D)x^2(x - 1).$$

We can divide through by $x$ and $x - 1$, which checks that $B$ and $E$ were correct:

$$-3x^2 - 4 = A(x^2 + 4) + (Cx + D)x.$$

Finally $x = 0$ yields $A = -1$. This leaves $-2x^2 = (Cx + D)x$. Then $C = -2$ and $D = 0$.

You should never have to do such a problem! I never intend to do another one. It completely depends on expecting the right form and matching the numerators. They could also be matched by comparing coefficients of $x^4$, $x^3$, $x^2$, $x$, $1$—to give five equations for $A$, $B$, $C$, $D$, $E$. That is an invitation to human error. Cover-up is the way to start, and usually the way to finish. With repeated factors and quadratic factors, match numerators at the end.

7.4 Exercises

Read-through questions

The idea of a fraction is to express $P(x)/Q(x)$ as a sum of simpler terms, each one easy to integrate. To begin, the degree of $P$ should be less than the degree of $Q$. Then $Q$ is split into factors like $x - 5$ (possibly repeated) and quadratic factors like $x^2 + x + 1$ (possibly repeated). The quadratic factors have two complex roots, and do not allow real linear factors.

A factor like $x - 5$ contributes a fraction $A/\underbrace{x - 5\cdots}$. Its integral is $\underbrace{\text{a}}$. To compute $A$, cover up $\underbrace{\text{b}}$ in the denominator of $P/Q$. Then set $x = \underbrace{\text{c}}$, and the rest of $P/Q$ becomes $A$. An equivalent method puts all fractions over a common denominator (which is $1$). Then match the $\underbrace{\text{k}}$. At the same point $x = \underbrace{\text{i}}$ this matching gives $A$.

A repeated linear factor $(x - 5)^2$ contributes not only $A/(x - 5)$ but also $B/(x - 5)$. A quadratic factor like $x^2 + x + 1$ contributes a fraction $\underbrace{\text{n}}/(x^2 + x + 1)$ involving $C$ and $D$. A repeated quadratic factor or a triple linear factor would bring in $(Ex + F)/(x^2 + x + 1)^2$ or $G/(x - 5)^3$. The conclusion is that any $P/Q$ can be split into partial fractions, which can always be integrated.

1 Find the numbers $A$ and $B$ to split $1/(x^2 - x)$:

$$\frac{1}{x(x - 1)} = \frac{A}{x} + \frac{B}{x - 1}.$$

Cover up $x$ and set $x = 0$ to find $A$. Cover up $x - 1$ and set $x = 1$ to find $B$. Then integrate.

2 Find the numbers $A$ and $B$ to split $1/(2x - 1)$:

$$\frac{1}{x^2 - 1} = \frac{A}{x - 1} + \frac{B}{x + 1}.$$

Multiply by $x - 1$ and set $x = 1$. Multiply by $x + 1$ and set $x = -1$. Integrate. Then find $A$ and $B$ again by method 1—

with numerator $A(x + 1) + B(x - 1)$ equal to 1.

Express the rational functions 3–16 as partial fractions:

$$3 \quad \frac{1}{(x - 3)(x - 2)} \quad 4 \quad \frac{x}{(x - 3)(x - 2)}$$

$$5 \quad \frac{x^2 + 1}{(x)(x + 1)(x + 2)} \quad 6 \quad \frac{1}{x^2 - x}$$

$$7 \quad \frac{3x + 1}{x^2} \quad 8 \quad \frac{3x + 1}{(x - 1)^2}$$

$$9 \quad \frac{3x^2}{x^2 + 1} \quad \text{(divide first)} \quad 10 \quad \frac{1}{(x - 1)(x^2 + 1)}$$

$$11 \quad \frac{1}{x^2(x - 1)} \quad 12 \quad \frac{x}{x^2 - 4}$$

$$13 \quad \frac{1}{x(x - 1)(x - 2)(x - 3)} \quad 14 \quad \frac{x^2 + 1}{x + 1} \quad \text{(divide first)}$$

$$15 \quad \frac{x^4 - 1}{(x + 1)(x - 1)(x^2 + 1)}$$

$$16 \quad \frac{1}{x^2(x - 1)} \quad \text{(remember the $A/x$ term)}$$

17 Apply Method 1 (matching numerators) to Example 3:

$$\frac{1}{cy - by^2} = \frac{A}{y} + \frac{B}{c - by} = \frac{A(c - by) + B y}{y(c - by)}.$$

Match the numerators on the far left and far right. Why does $Ac = 1$? Why does $-bA + B = 0$? What are $A$ and $B$?
18 What goes wrong if we look for $A$ and $B$ so that

$$\frac{x^2}{(x-3)(x+3)} = \frac{A}{x-3} + \frac{B}{x+3}.$$  

Over a common denominator, try to match the numerators. What to do first?

19 Split $\frac{3x^2}{x^3-1} = \frac{3x^2}{(x-1)(x^2+x+1)}$ into $\frac{A}{x-1} + \frac{Bx+C}{x^2+x+1}$.

(a) Cover up $x-1$ and set $x = 1$ to find $A$.
(b) Subtract $\frac{A}{x-1}$ from the left side. Find $Bx+C$.
(c) Integrate all terms. Why do we already know $\ln(x^3-1) = \ln(x-1) + \ln(x^2+x+1)$?

20 Solve $dy/dt = 1 - y^2$ by separating $\int dy/(1 - y^2) = \int dt$. Then

$$\frac{1}{1 - y^2} = \frac{1}{(1-y)(1+y)} = \frac{1/2}{1-y} + \frac{1/2}{1+y}$$

Integration gives $\frac{1}{2}\ln\frac{1+y}{1-y} = t + C$. With $y_0 = 0$ the constant is $C = \ldots$. Taking exponentials gives $\ldots$. The solution is $y = \ldots$. This is the S-curve.

21 By substitution change 21–28 to integrals of rational functions. Problem 23 integrates $1/\sin \theta$ with no special trick.

22 $\int \frac{1}{e^{2x} - e^x} \, dx$

23 $\int \frac{\sin \theta \, d\theta}{1 - \cos^2 \theta}$

24 $\int \frac{dt}{(e^t - e^{-t})^2}$

25 $\int \frac{1 + e^x}{1 - e^x} \, dx$

26 $\int \frac{\sqrt{x} - 8}{x} \, dx$

27 $\int \frac{dx}{1 + \sqrt{x} + 1}$

28 $\int \frac{dx}{\sqrt{x} - 1 + \sqrt{x}}$

29 Multiply this partial fraction by $x - a$. Then let $x \to a$:

$$\frac{1}{Q(x)} = \frac{A}{x-a} + \ldots.$$  

Show that $A = 1/Q'(a)$. When $x = a$ is a double root this fails because $Q'(a) = \ldots$.

30 Find $A$ in $\frac{1}{x^3 - 1} = \frac{A}{x-1} + \ldots$. Use Problem 29.

31 (for instructors only) Which rational functions $P/Q$ are the derivatives of other rational functions (no logarithms)?

---

7.5 Improper Integrals

"Improper" means that some part of $\int_a^b y(x) \, dx$ becomes infinite. It might be $b$ or $a$ or the function $y$. The region under the graph reaches infinitely far—to the right or left or up or down. (Those come from $b = \infty$ and $a = -\infty$ and $y \to \infty$ and $y \to -\infty$.) Nevertheless the integral may "converge." Just because the region is infinite, it is not automatic that the area is infinite. That is the point of this section—to decide when improper integrals have proper answers.

The first examples show finite area when $b = \infty$, then $a = -\infty$, then $y = 1/\sqrt{x}$ at $x = 0$. The areas in Figure 7.6 are 1, 1, 2:

$$\int_1^\infty \frac{dx}{x^2} = -\left[ \frac{1}{x} \right]_1^\infty = 1$$

$$\int_{-\infty}^0 e^x \, dx = e^0 - e^{-\infty} = 1$$

$$\int_0^1 \frac{dx}{\sqrt{x}} = 2\sqrt{x}\bigg|_0^1 = 2.$$  

Fig. 7.6 The shaded areas are finite but the regions go to infinity.
In practice we substitute the dangerous limits and watch what happens. When the integral is \(-1/x\), substituting \(b = \infty\) gives \(\int_{-1}^\infty \frac{1}{x} \ dx = \ln x \bigg|_{-\infty}^{\infty} = 0\). When the integral is \(e^x\), substituting \(a = -\infty\) gives \(\int_{-\infty}^{\infty} e^x \ dx = \lim_{a \to -\infty} \int_a^b y(x) \ dx = \frac{1}{p-1}\).

But it is not completely precise. The strict rules involve a limit. Calculus sneaks up on \(1/\infty\) and \(e^{-\infty}\) just as it sneaks up on \(0/0\). Instead of swallowing an infinite region all at once, the formal definitions push out to the limit:

\[
\int_a^b y(x) \ dx = \lim_{b \to \infty} \int_a^b y(x) \ dx
\]

The conclusion is the same. The first examples converged to 1, 1, 2. Now come two more examples going out to \(b = \infty\):

The area under \(1/x\) is infinite:

\[
\int_1^\infty \frac{dx}{x} = \ln x \bigg|_1^\infty = \infty
\]

The area under \(1/x^p\) is finite if \(p > 1\):

\[
\int_1^\infty \frac{dx}{x^p} = \frac{x^{1-p}}{1-p} \bigg|_1^\infty = \frac{1}{p-1}.
\]

The area under \(1/x\) is like \(1 + \frac{1}{2} + \frac{1}{3} + \cdots\), which is also infinite. In fact the sum approximates the integral—the curved area is close to the rectangular area. They go together (slowly to infinity).

A larger \(p\) brings the graph more quickly to zero. Figure 7.7a shows a finite area \(1/(p - 1) = 100\). The region is still infinite, but we can cover it with strips cut out of a square! The borderline for finite area is \(p = 1\). I call it the borderline, but \(p = 1\) is strictly on the side of divergence.

The borderline is also \(p = 1\) when the function climbs the \(y\) axis. At \(x = 0\), the graph of \(y = 1/x^p\) goes to infinity. For \(p = 1\), the area under \(1/x\) is again infinite. But at \(x = 0\) it is a small \(p\) (meaning \(p < 1\)) that produces finite area:

\[
\int_0^1 \frac{dx}{x} = \ln x \bigg|_0^1 = \infty
\]

Loosely speaking "\(-\ln 0 = \infty\)." Strictly speaking we integrate from the point \(x = a\) near zero, to get \(\int_a^1 \frac{dx}{x} = -\ln a\). As \(a\) approaches zero, the area shows itself as infinite. For \(y = 1/x^p\), which blows up faster, the area \(-1/x\) is again infinite.

For \(y = 1/\sqrt{x}\), the area from 0 to 1 is 2. In that case \(p = \frac{1}{2}\). For \(p = 99/100\) the area is \(1/(1 - p) = 100\). Approaching \(p = 1\) the borderline in Figure 7.7 seems clear. But that cutoff is not as sharp as it looks.

Fig. 7.7 Graphs of \(1/x^p\) on both sides of \(p = 1\). I drew the same curves!
Narrower borderline  Under the graph of $1/x$, the area is infinite. When we divide by $\ln x$ or $(\ln x)^2$, the borderline is somewhere in between. One has infinite area (going out to $x = \infty$), the other area is finite:

$$
\int_0^\infty \frac{dx}{x(\ln x)} = \ln(\ln x) \bigg|_0^\infty = \infty \quad \text{and} \quad \int_0^\infty \frac{dx}{x(\ln x)^2} = -\frac{1}{\ln x} \bigg|_0^\infty = 1. \quad (4)
$$

The first is $\int du/u$ with $u = \ln x$. The logarithm of $\ln x$ does eventually make it to infinity. At $x = 10^{10}$, the logarithm is near 23 and $\ln(\ln x)$ is near 3. That is slow! Even slower is $\ln(\ln(\ln x))$ in Problem 11. No function is exactly on the borderline.

The second integral in equation (4) is convergent (to 1). It is $\int du/u^2$ with $u = \ln x$. At first I wrote it with $x$ going from zero to infinity. That gave an answer I couldn't believe:

$$
\int_0^\infty \frac{dx}{x(\ln x)^3} = -\frac{1}{\ln x} \bigg|_0^\infty = 0 \,(?)
$$

There must be a mistake, because we are integrating a positive function. The area can't be zero. It is true that $1/\ln b$ goes to zero as $b \to \infty$. It is also true that $1/\ln a$ goes to zero as $a \to 0$. But there is another infinity in this integral. The trouble is at $x = 1$, where $\ln x$ is zero and the area is infinite.

EXAMPLE 1  The factor $e^{-x}$ overrides any power $x^p$ (but only as $x \to \infty$).

$$
\int_0^\infty x^5 e^{-x} \, dx = 50! \quad \text{but} \quad \int_0^\infty x^{-1} e^{-x} \, dx = \infty.
$$

The first integral is $(50)(49)(48)\cdots(1)$. It comes from fifty integrations by parts (not recommended). Changing 50 to $\frac{1}{2}$, the integral defines "half factorial." The product $\frac{1}{2}(-\frac{1}{2})(-\frac{3}{2})\cdots$ has no way to stop, but somehow $\frac{1}{2}!$ is $\frac{1}{2}\sqrt{\pi}$. See Problem 28. The integral $\int_0^\infty x^0 e^{-x} \, dx = 1$ is the reason behind "zero factorial" = 1. That seems the most surprising of all.

The area under $e^{-x}/x$ is $(-1)! = \infty$. The factor $e^{-x}$ is absolutely no help at $x = 0$. That is an example (the first of many) in which we do not know an antiderivative—but still we get a decision. To integrate $e^{-x}/x$ we need a computer. But to decide that an improper integral is infinite (in this case) or finite (in other cases), we rely on the following comparison test:

\begin{center}
\begin{tabular}{|c|}
\hline
\textbf{7C} (Comparison test)  Suppose that $0 \leq u(x) \leq v(x)$. Then the area under $u(x)$ is smaller than the area under $v(x)$:
\end{tabular}
\begin{tabular}{|c|}
\hline
$\int u(x) \, dx < \infty$ if $\int u(x) \, dx < \infty$ \quad if $\int u(x) \, dx = \infty$ then $\int v(x) \, dx = \infty$.
\hline
\end{tabular}
\end{center}

Comparison can decide if the area is finite. We don't get the exact area, but we learn about one function from the other. The trick is to construct a simple function (like $1/x^p$) which is on one side of the given function—and stays close to it:

EXAMPLE 2  \quad \int_1^\infty \frac{dx}{x^2 + 4x} \ \text{converges by comparison with} \ \int_1^\infty \frac{dx}{x^2} = 1.

EXAMPLE 3  \quad \int_1^\infty \frac{dx}{\sqrt{x} + 1} \ \text{diverges by comparison with} \ \int_1^\infty \frac{dx}{2\sqrt{x}} = \infty.
EXAMPLE 4 \[ \int_0^1 \frac{dx}{x^2 + 4x} \] diverges by comparison with \[ \int_0^1 \frac{dx}{5x} = \infty. \]

EXAMPLE 5 \[ \int_0^1 \frac{dx}{\sqrt{x} + 1} \] converges by comparison with \[ \int_0^1 \frac{dx}{1} = 1. \]

In Examples 2 and 5, the integral on the right is larger than the integral on the left. Removing \(4x\) and \(\sqrt{x}\) increased the area. Therefore the integrals on the left are somewhere between 0 and 1.

In Examples 3 and 4, we increased the denominators. The integrals on the right are smaller, but still they diverge. So the integrals on the left diverge. The idea of comparing functions is seen in the next examples and Figure 7.8.

EXAMPLE 6 \[ \int_0^\infty e^{-x^2} dx \] is below \[ \int_0^1 dx + \int_1^\infty e^{-x} dx = 1 + 1. \]

EXAMPLE 7 \[ \int_1^e \frac{dx}{\ln x} \] is above \[ \int_1^\infty \frac{dx}{x \ln x} = \infty. \]

EXAMPLE 8 \[ \int_0^1 \frac{dx}{\sqrt{x-x^2}} \] is below \[ \int_0^1 \frac{dx}{\sqrt{x}} + \int_0^1 \frac{dx}{\sqrt{1-x}} = 2 + 2. \]

Fig. 7.8 Comparing \(u(x)\) to \(v(x)\): \(\int_1^1 \frac{dx}{\ln x} = \infty\) and \(\int_0^1 \frac{dx}{\sqrt{x-x^2}} < 4.\) But \(\infty - \infty \neq 0.\)

There are two situations not yet mentioned, and both are quite common. The first is an integral all the way from \(a = -\infty\) to \(b = +\infty.\) That is split into two parts, and each part must converge. By definition, the limits at \(-\infty\) and \(+\infty\) are kept separate: \[ \int_{-\infty}^{\infty} y(x) \, dx = \int_{-\infty}^{0} y(x) \, dx + \int_{0}^{\infty} y(x) \, dx = \lim_{a \to -\infty} \int_{a}^{0} y(x) \, dx + \lim_{b \to +\infty} \int_{0}^{b} y(x) \, dx. \]

The bell-shaped curve \(y = e^{-x^2}\) covers a finite area (exactly \(\pi\)). The region extends to infinity in both directions, and the separate areas are \(\frac{1}{2}\pi.\) But notice: \[ \int_{-\infty}^{\infty} x \, dx \text{ is not defined even though } \int_{-b}^{b} x \, dx = 0 \text{ for every } b. \]

The area under \(y = x\) is \(+\infty\) on one side of zero. The area is \(-\infty\) on the other side. We cannot accept \(\infty - \infty = 0.\) The two areas must be separately finite, and in this case they are not.
EXAMPLE 9 \( \frac{1}{x} \) has balancing regions left and right of \( x = 0 \). Compute \( \int_{-1}^{1} \frac{dx}{x} \).

This integral does not exist. There is no answer, even for the region in Figure 7.8c. (They are mirror images because \( \frac{1}{x} \) is an odd function.) You may feel that the combined integral from \(-1\) to \(1\) should be zero. Cauchy agreed with that—his “principal value integral” is zero. But the rules say no: \( -\infty \) to \( \infty \) is not zero.

7.5 EXERCISES

Read-through questions

An improper integral \( \int_{a}^{b} y(x) \, dx \) has lower limit \( a = \_ \_ \_ \) or upper limit \( b = \_ \_ \_ \_ \) or \( y \) becomes \_ \_ \_ \_ \_ \_ in the interval \( a < x < b \). The example \( \int_{1}^{\infty} \frac{dx}{x^3} \) is improper because \_ \_ \_ \_ \_ \_ \_ \_ \_.

We should study the limit of \( \int_{a}^{b} \frac{dx}{x^3} \) as \_ \_ \_ \_ \_. In practice we work directly with \( -\frac{1}{2} x^{-2} \). For \( p > 1 \) the improper integral \( \_ \_ \_ \_ \_ \_ \) is finite. For \( p < 1 \) the improper integral \( \_ \_ \_ \_ \_ \) is finite. For \( y = e^{-x} \) the integral from \( 0 \) to \( \infty \) is \_ \_ \_ \_ \_.

Suppose \( 0 < u(x) \leq v(x) \) for all \( x \). The convergence of \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_ \_
34 The escape velocity for a black hole exceeds the speed of light: $v_0 > 3 \cdot 10^8 \text{ m/sec}$. The Earth has $GM = 4 \cdot 10^{10} \text{m}^3/\text{sec}^2$. If it were compressed to radius $R = \ldots$, the Earth would be a black hole.

35 Show how the area under $y = 1/2^x$ can be covered (draw a graph) by rectangles of area $1 + \frac{1}{2} + \frac{1}{4} + \cdots = 2$. What is the exact area from $x = 0$ to $x = \infty$?

36 Explain this paradox:
\[ \int_{-b}^{b} \frac{x \, dx}{1 + x^2} = 0 \text{ for every } b \text{ but } \int_{-\infty}^{\infty} \frac{x \, dx}{1 + x^2} \text{ diverges.} \]

37 Compute the area between $y = \sec x$ and $y = \tan x$ for $0 \leq x \leq \pi/2$. What is improper?

38 Compute any of these integrals found by geniuses:
\[
\begin{align*}
\int x^{-1/2} \, dx &= \pi \\
\int_0^\infty \frac{e^{-x} - e^{-2x}}{x} \, dx &= \ln 2 \\
\int_0^\infty x e^{-x} \cos x \, dx &= 0 \\
\int_0^\infty \cos x^2 \, dx &= \sqrt{\pi/8}.
\end{align*}
\]

39 For which $p$ is $\int_0^\infty \frac{dx}{x^p + x^{-p}} = \infty$?

40 Explain from Figure 7.6c why the red area is 2, when Figure 7.6a has red area 1.
# Contents

**CHAPTER 4**

| 4.1 | Derivatives by the Chain Rule | 154 |
| 4.2 | Implicit Differentiation and Related Rates | 160 |
| 4.3 | Inverse Functions and Their Derivatives | 164 |
| 4.4 | Inverses of Trigonometric Functions | 171 |

**CHAPTER 5**

| 5.1 | The Idea of the Integral | 177 |
| 5.2 | Antiderivatives | 182 |
| 5.3 | Summation vs. Integration | 187 |
| 5.4 | Indefinite Integrals and Substitutions | 195 |
| 5.5 | The Definite Integral | 201 |
| 5.6 | Properties of the Integral and the Average Value | 206 |
| 5.7 | The Fundamental Theorem and Its Consequences | 213 |
| 5.8 | Numerical Integration | 220 |

**CHAPTER 6**

| 6.1 | An Overview | 228 |
| 6.2 | The Exponential $e^x$ | 236 |
| 6.3 | Growth and Decay in Science and Economics | 242 |
| 6.4 | Logarithms | 252 |
| 6.5 | Separable Equations Including the Logistic Equation | 259 |
| 6.6 | Powers Instead of Exponentials | 267 |
| 6.7 | Hyperbolic Functions | 277 |

**CHAPTER 7**

| 7.1 | Integration by Parts | 283 |
| 7.2 | Trigonometric Integrals | 288 |
| 7.3 | Trigonometric Substitutions | 294 |
| 7.4 | Partial Fractions | 300 |
| 7.5 | Improper Integrals | 305 |

**CHAPTER 8**

| 8.1 | Areas and Volumes by Slices | 311 |
| 8.2 | Length of a Plane Curve | 320 |
| 8.3 | Area of a Surface of Revolution | 325 |
| 8.4 | Probability and Calculus | 328 |
| 8.5 | Masses and Moments | 336 |
| 8.6 | Force, Work, and Energy | 342 |
CHAPTER 8

Applications of the Integral

We are experts in one application of the integral—to find the area under a curve. The curve is the graph of \( y = v(x) \), extending from \( x = a \) at the left to \( x = b \) at the right. The area between the curve and the \( x \) axis is the definite integral.

I think of that integral in the following way. The region is made up of thin strips. Their width is \( dx \) and their height is \( v(x) \). The area of a strip is \( v(x) \times dx \). The area of all the strips is \( \int_a^b v(x) \, dx \). Strictly speaking, the area of one strip is meaningless—genuine rectangles have width \( \Delta x \). My point is that the picture of thin strips gives the correct approach.

We know what function to integrate (from the picture). We also know how (from this course or a calculator). The new applications to volume and length and surface area cut up the region in new ways. Again the small pieces tell the story. In this chapter, what to integrate is more important than how.

8.1 Areas and Volumes by Slices

This section starts with areas between curves. Then it moves to volumes, where the strips become slices. We are weighing a loaf of bread by adding the weights of the slices. The discussion is dominated by examples and figures—the theory is minimal. The real problem is to set up the right integral. At the end we look at a different way of cutting up volumes, into thin shells. All formulas are collected into a final table.

Figure 8.1 shows the area between two curves. The upper curve is the graph of \( y = v(x) \). The lower curve is the graph of \( y = w(x) \). The strip height is \( v(x) - w(x) \), from one curve down to the other. The width is \( dx \) (speaking informally again). The total area is the integral of “top minus bottom”:

\[
\text{area between two curves} = \int_a^b \left[ v(x) - w(x) \right] \, dx. \tag{1}
\]

EXAMPLE 1 The upper curve is \( y = 6x \) (straight line). The lower curve is \( y = 3x^2 \) (parabola). The area lies between the points where those curves intersect.

To find the intersection points, solve \( v(x) = w(x) \) or \( 6x = 3x^2 \).
One crossing is at $x = 0$, the other is at $x = 2$. The area is an integral from 0 to 2:

$$\text{area} = \int_0^2 (v - w) \, dx = \int_0^2 (6x - 3x^2) \, dx = 3x^2 - x^3 \bigg|_0^2 = 4.$$  

**EXAMPLE 2** Find the area between the circle $v = \sqrt{1 - x^2}$ and the $45^\circ$ line $w = x$.

First question: Which area and what limits? Start with the pie-shaped wedge in Figure 8.1b. The area begins at the y axis and ends where the circle meets the line. At the intersection point we have $v(x) = w(x)$:

$$\sqrt{1 - x^2} = x$$

Squaring gives $1 - x^2 = x^2$ and then $2x^2 = 1$, thus $x^2 = \frac{1}{2}$. The endpoint is at $x = \frac{1}{\sqrt{2}}$. Now integrate the strip height $v - w$:

$$\int_0^{1/\sqrt{2}} (\sqrt{1 - x^2} - x) \, dx = \frac{1}{2} \sin^{-1} x + \frac{1}{2} x \sqrt{1 - x^2} - \frac{1}{2} x^2 \bigg|_0^{1/\sqrt{2}}$$

$$= \frac{1}{2} \sin^{-1} \left( \frac{1}{\sqrt{2}} \right) + \frac{1}{4} - \frac{1}{4} = \frac{1}{2} \left( \frac{\pi}{4} \right).$$

The area is $\pi/8$ (one eighth of the circle). To integrate $\sqrt{1 - x^2} \, dx$ we apply the techniques of Chapter 7: Set $x = \sin \theta$, convert to $\int \cos^2 \theta \, d\theta = \frac{1}{2} (\theta + \sin \theta \cos \theta)$, convert back using $\theta = \sin^{-1} x$. It is harder than expected, for a familiar shape.

**Remark** Suppose the problem is to find the whole area between the circle and the line. The figure shows $v = w$ at two points, which are $x = 1/\sqrt{2}$ (already used) and also $x = -1/\sqrt{2}$. Instead of starting at $x = 0$, which gave $\frac{1}{8}$ of a circle, we now include the area to the left.

Main point: Integrating from $x = -1/\sqrt{2}$ to $x = 1/\sqrt{2}$ will give the wrong answer. It misses the part of the circle that bulges out over itself, at the far left. In that part, the strips have height $2v$ instead of $v - w$. The figure is essential, to get the correct area of this half-circle.

**HORIZONTAL STRIPS INSTEAD OF VERTICAL STRIPS**

There is more than one way to slice a region. Vertical slices give x integrals. Horizontal slices give y integrals. We have a free choice, and sometimes the y integral is better.
Figure 8.2 shows a unit parallelogram, with base 1 and height 1. To find its area from vertical slices, three separate integrals are necessary. You should see why! With horizontal slices of length 1 and thickness $dy$, the area is just $\int_0^1 dy = 1$.

**EXAMPLE 3** Find the area under $y = \ln x$ (or beyond $x = e^y$) out to $x = e$.

The $x$ integral from vertical slices is in Figure 8.2c. The $y$ integral is in 8.2d. The area is a choice between two equal integrals (I personally would choose $y$):

$$\int_0^1 \ln x \, dx = \left[ x \ln x - x \right]_1^e = 1 \quad \text{or} \quad \int_0^1 (e - e^y) \, dy = \left[ ey - e^y \right]_0^1 = 1.$$

**VOLUMES BY SLICES**

For the first time in this book, we now look at **volumes**. The regions are *three-dimensional solids*. There are three coordinates $x, y, z$—and many ways to cut up a solid.

Figure 8.3 shows one basic way—using *slices*. The slices have thickness $dx$, like strips in the plane. Instead of the height $y$ of a strip, we now have the **area** $A$ of a **cross-section**. This area is different for different slices: $A$ depends on $x$. The volume of the slice is its area times its thickness: $dV = A(x) \, dx$. The **volume of the whole solid is the integral**:

$$\text{volume} = \int_A(x) \, dx.$$  \hspace{1cm} (2)

**Note** An actual slice does not have the same area on both sides! Its thickness is $\Delta x$ (not $dx$). Its volume is approximately $A(x) \Delta x$ (but not exactly). In the limit, the thickness approaches zero and the sum of volumes approaches the integral.

For a cylinder all slices are the same. Figure 8.3b shows a cylinder—*not circular*. The area is a fixed number $A$, so integration is trivial. The **volume is $A$ times $h$**. The
letter \( h \), which stands for \textit{height}, reminds us that the cylinder often stands on its end. Then the slices are horizontal and the \( y \) integral or \( z \) integral goes from 0 to \( h \).

When the cross-section is a circle, the cylinder has volume \( \pi r^2 h \).

**EXAMPLE 4**  The \textit{triangular wedge} in Figure 8.3b has constant cross-sections with area \( A = \frac{1}{2}(3)(4) = 6 \). The volume is \( 6h \).

**EXAMPLE 5**  For the \textit{triangular pyramid} in Figure 8.3c, the area \( A(x) \) drops from 6 to 0. It is a general rule for pyramids or cones that their volume has an extra factor \( \frac{1}{3} \) (compared to cylinders). The volume is now \( 2h \) instead of \( 6h \). For a cone with base area \( \pi r^2 \), the volume is \( \frac{1}{3} \pi r^2 h \). \textit{Tapering the area to zero leaves only} \( \frac{1}{3} \) \textit{of the volume}.

Why the \( \frac{1}{3} \)? Triangles sliced from the pyramid have shorter sides. Starting from 3 and 4, the side lengths \( 3(1 - x/h) \) and \( 4(1 - x/h) \) drop to zero at \( x = h \). The area is \( A = 6(1 - x/h)^2 \). Notice: The side lengths go down linearly, the area drops quadratically. The factor \( \frac{1}{3} \) really comes from integrating \( x^2 \) to get \( \frac{1}{3} x^3 \):

\[
\int_0^h A(x) \, dx = \int_0^h 6 \left( 1 - \frac{x}{h} \right)^2 \, dx = -2h \left( 1 - \frac{x}{h} \right)^3 \bigg|_0^h = 2h.
\]

**EXAMPLE 6**  A half-sphere of radius \( R \) has known volume \( \frac{1}{2}(\frac{4}{3} \pi R^3) \). Its cross-sections are \textit{semicircles}. The key relation is \( x^2 + r^2 = R^2 \), for the right triangle in Figure 8.4a. The area of the semicircle is \( A = \frac{1}{2} \pi r^2 = \frac{1}{2} \pi (R^2 - x^2) \). So we integrate \( A(x) \):

\[
\text{volume} = \int_{-R}^R A(x) \, dx = \frac{1}{2} \pi \int_{-R}^R (R^2 - x^2) \, dx = \frac{2}{3} \pi R^3.
\]

**EXAMPLE 7**  Find the volume of the same half-sphere using horizontal slices (Figure 8.4b). The sphere still has radius \( R \). The new right triangle gives \( y^2 + r^2 = R^2 \). Since we have full circles the area is \( \pi r^2 = \pi (R^2 - y^2) \). Notice that this is \( A(y) \) not \( A(x) \). But the \( y \) integral starts at zero:

\[
\text{volume} = \int_0^R A(y) \, dy = \pi \int_0^R \left( R^2 - y^2 \right) \, dy = \frac{2}{3} \pi R^3 \quad \text{(as before)}.
\]

---

**Fig. 8.4** A half-sphere sliced vertically or horizontally. Washer area \( \pi f^2 - \pi g^2 \).

---

**SOLIDS OF REVOLUTION**

Cones and spheres and circular cylinders are \textit{solids of revolution.} Rotating a horizontal line around the \( x \) axis gives a cylinder. Rotating a sloping line gives a cone. Rotating a semicircle gives a sphere. If a circle is moved away from the axis, rotation produces a torus (a doughnut). The rotation of any curve \( y = f(x) \) produces a \textit{solid of revolution}. 
The volume of that solid is made easier because every cross-section is a circle. All slices are pancakes (or pizzas). Rotating the curve \( y = f(x) \) around the \( x \) axis gives disks of radius \( y \), so the area is \( A = \pi y^2 = \pi [f(x)]^2 \). We add the slices:

\[
\text{volume of solid of revolution} = \int_a^b \pi y^2 \, dx = \int_a^b \pi [f(x)]^2 \, dx.
\]

**EXAMPLE 8** Rotating \( y = \sqrt{x} \) with \( A = \pi (\sqrt{x})^2 \) produces a "headlight" (Figure 8.5a):

\[
\text{volume of headlight} = \int_0^2 A \, dx = \int_0^2 \pi x \, dx = \left[ \frac{1}{2} \pi x^2 \right]_0^2 = 2\pi.
\]

If the same curve is rotated around the \( y \) axis, it makes a champagne glass. The slices are horizontal. The area of a slice is \( \pi x^2 \) not \( \pi y^2 \). When \( y = \sqrt{x} \) this area is \( \pi y^4 \).

Integrating from \( y = 0 \) to \( \sqrt{2} \) gives the champagne volume \( \pi (\sqrt{2})^4/5 \).

\[
\text{revolution around the } y \text{ axis: volume} = \int \pi x^2 \, dy.
\]

**EXAMPLE 9** The headlight has a hole down the center (Figure 8.5b). Volume = ?

The hole has radius 1. All of the \( \sqrt{x} \) solid is removed, up to the point where \( \sqrt{x} \) reaches 1. After that, from \( x = 1 \) to \( x = 2 \), each cross-section is a disk with a hole. The disk has radius \( f = \sqrt{x} \) and the hole has radius \( g = 1 \). The slice is a flat ring or a "washer." Its area is the full disk minus the area of the hole:

\[
\text{area of washer} = \pi f^2 - \pi g^2 = \pi (\sqrt{x})^2 - \pi (1)^2 = \pi x - \pi.
\]

This is the area \( A(x) \) in the method of washers. Its integral is the volume:

\[
\int_1^2 A \, dx = \int_1^2 (\pi x - \pi) \, dx = \left[ \frac{1}{2} \pi x^2 - \pi x \right]_1^2 = \frac{1}{2} \pi.
\]

Please notice: The washer area is not \( \pi(f - g)^2 \). It is \( A = \pi f^2 - \pi g^2 \).

**EXAMPLE 10** (Doughnut sliced into washers) Rotate a circle of radius \( a \) around the \( x \) axis. The center of the circle stays out at a distance \( b > a \). Show that the volume of the doughnut (or torus) is \( 2\pi^2 a^2 b \).
The outside half of the circle rotates to give the outside of the doughnut. The inside half gives the hole. The biggest slice (through the center plane) has outer radius \( b + a \) and inner radius \( b - a \).

Shifting over by \( x \), the outer radius is \( f = b + \sqrt{a^2 - x^2} \) and the inner radius is \( g = b - \sqrt{a^2 - x^2} \). Figure 8.5c shows a slice (a washer) with area \( \pi f^2 - \pi g^2 \).

\[
\text{area } A = \pi(b + \sqrt{a^2 - x^2})^2 - \pi(b - \sqrt{a^2 - x^2})^2 = 4\pi b \sqrt{a^2 - x^2}.
\]

Now integrate over the washers to find the volume of the doughnut:

\[
\int_{-a}^{a} A(x)\,dx = 4\pi b \int_{-a}^{a} \sqrt{a^2 - x^2} \,dx = (4\pi b)(\frac{1}{2}\pi a^2) = 2\pi^2 a^2 b.
\]

That integral \( \frac{1}{2}\pi a^2 \) is the area of a semicircle. When we set \( x = a \sin \theta \) the area is \( \int a^2 \cos^2 \theta \,d\theta \). Not for the last time do we meet \( \cos^2 \theta \).

The hardest part is visualizing the washers, because a doughnut usually breaks the other way. A better description is a bagel, sliced the long way to be buttered.

**VOLUMES BY CYLINDRICAL SHELLS**

Finally we look at a different way of cutting up a solid of revolution. So far it was cut into slices. The slices were perpendicular to the axis of revolution. Now the cuts are parallel to the axis, and each piece is a thin cylindrical shell. The new formula gives the same volume, but the integral to be computed might be easier.

Figure 8.6a shows a solid cone. A shell is inside it. The inner radius is \( x \) and the outer radius is \( x + dx \). The shell is an outer cylinder minus an inner cylinder:

\[
\text{shell volume } \pi(x + dx)^2 h - \pi x^2 h = \pi x^2 h + 2\pi x (dx)h + \pi(dx)^2 h = \pi x^2 h.
\]

The term that matters is \( 2\pi x (dx)h \). The shell volume is essentially \( 2\pi x \) (the distance around) \( \times dx \) (the thickness) \( \times h \) (the height). The volume of the solid comes from putting together the thin shells:

\[
\text{solid volume } = \text{integral of shell volumes } = \int 2\pi x h \,dx.
\]

This is the central formula of the shell method. The rest is examples.

**Remark on this volume formula** It is completely typical of integration that \( (dx)^2 \) and \( (\Delta x)^2 \) disappear. The reason is this. The number of shells grows like \( 1/\Delta x \). Terms of order \( (\Delta x)^2 \) add up to a volume of order \( \Delta x \) (approaching zero). The linear term involving \( \Delta x \) or \( dx \) is the one to get right. Its limit gives the integral \( \int 2\pi x h \,dx \). The key is to build the solid out of shells—and to find the area or volume of each piece.

**EXAMPLE 11** Find the volume of a cone (base area \( \pi r^2 \), height \( b \)) cut into shells.

A tall shell at the center has \( h \) near \( b \). A short shell at the outside has \( h \) near zero. In between the shell height \( h \) decreases linearly, reaching zero at \( x = r \). The height in Figure 8.6a is \( h = b - bx/r \). Integrating over all shells gives the volume of the cone (with the expected \( \frac{1}{3} \)):

\[
\int_{0}^{r} 2\pi x \left(b - \frac{x}{r}\right) dx = \left[ \pi x^2 b - \frac{2\pi x^3 b}{3r} \right]_{0}^{r} = \frac{1}{3} \pi r^2 b.
\]
8.1 Areas and Volumes by Slices

EXAMPLE 12  Bore a hole of radius $a$ through a sphere of radius $b > a$.

The hole removes all points out to $x = a$, where the shells begin. The height of the shell is $h = 2\sqrt{b^2 - x^2}$. (The key is the right triangle in Figure 8.6b. The height upward is $\sqrt{b^2 - x^2}$—this is half the height of the shell.) Therefore the sphere-with-hole has

$$\text{volume} = \int_a^b 2\pi x h \, dx = \int_a^b 4\pi x \sqrt{b^2 - x^2} \, dx.$$  

With $u = b^2 - x^2$ we almost see $du$. Multiplying $du = -2x \, dx$ is an extra factor $-2\pi$:

$$\text{volume} = -2\pi \int u \, du = -2\pi \left( \frac{3}{2} u^{3/2} \right).$$

We can find limits on $u$, or we can put back $u = b^2 - x^2$:

$$\text{volume} = -\frac{4\pi}{3} (b^2 - x^2)^{3/2} \bigg|_a^b = \frac{4\pi}{3} (b^2 - a^2)^{3/2}.$$  

If $a = b$ (the hole is as big as the sphere) this volume is zero. If $a = 0$ (no hole) we have $4\pi b^3/3$ for the complete sphere.

**Question** What if the sphere-with-hole is cut into slices instead of shells?

**Answer** Horizontal slices are washers (Problem 66). Vertical slices are not good.

EXAMPLE 13  Rotate the parabola $y = x^2$ around the $y$ axis to form a bowl.

We go out to $x = \sqrt{2}$ (and up to $y = 2$). The shells in Figure 8.6c have height $h = 2 - x^2$. The bowl (or paraboloid) is the same as the headlight in Example 8, but we have shells not slices:

$$\int_0^{\sqrt{2}} 2\pi x(2 - x^2) \, dx = 2\pi x^2 - \frac{2\pi x^4}{4} \bigg|_0^{\sqrt{2}} = 2\pi.$$  

**TABLE**

| **area between curves:** $A = \int (w(x) - v(x)) \, dx$ |
| **solid volume cut into slices:** $V = \int A(x) \, dx$ or $\int A(y) \, dy$ |
| **solid of revolution:** cross-section $A = \pi y^2$ or $\pi x^2$ |
| **solid with hole:** washer area $A = \pi f^2 - \pi g^2$ |
| **solid of revolution cut into shells:** $V = \int 2\pi x h \, dx$. |
**Which to use, slices or shells?** Start with a vertical line going up to \( y = \cos x \). Rotating the line around the \( x \) axis produces a slice (a circular disk). The radius is \( \cos x \). Rotating the line around the \( y \) axis produces a shell (the outside of a cylinder). The height is \( \cos x \). See Figure 8.7 for the slice and the shell. For volumes we just integrate \( \pi \cos^2 x\, dx \) (the slice volume) or \( 2\pi \cos x\, dx \) (the shell volume).

This is the normal choice—slices through the \( x \) axis and shells around the \( y \) axis. Then \( y = f(x) \) gives the disk radius and the shell height. The slice is a washer instead of a disk if there is also an inner radius \( g(x) \). No problem—just integrate small volumes.

What if you use slices for rotation around the \( y \) axis? The disks are in Figure 8.7b, and their radius is \( x \). This is \( x = \cos^{-1} y \) in the example. You have to solve \( y = f(x) \) to find \( x \) in terms of \( y \). Similarly for shells around the \( x \) axis: The length of the shell is \( x = f^{-1}(y) \). Integrating may be difficult or impossible.

When \( y = \cos x \) is rotated around the \( x \) axis, here are the choices for volume:

\[
\text{(good by slices)} \int \pi \cos^2 x\, dx \quad \text{(bad by shells)} \int 2\pi y \cos^{-1} y\, dy.
\]

![Figure 8.7 Slices through x axis and shells around y axis (good). The opposite way needs \( f^{-1}(y) \).](image)

### 8.1 EXERCISES

#### Read-through questions

The area between \( y = x^2 \) and \( y = x^4 \) equals the integral of \[ a \]. If the region ends where the curves intersect, we find the limits on \( x \) by solving \[ b \]. Then the area equals \[ c \]. When the area between \( y = \sqrt{x} \) and the \( y \) axis is sliced horizontally, the integral to compute is \[ d \].

In three dimensions the volume of a slice is its thickness \( dx \) times its \[ e \]. If the cross-sections are squares of side \( 1 - x \), the volume comes from \[ f \]. From \( x = 0 \) to \( x = 1 \), this gives the volume \[ g \] of a square \[ h \]. If the cross-sections are circles of radius \( 1 - x \), the volume comes from \[ i \]. This gives the volume \[ j \] of a circular \[ k \].

For a solid of revolution, the cross-sections are \[ l \]. Rotating the graph of \( y = f(x) \) around the \( x \) axis gives a solid volume \[ m \]. Rotating around the \( y \) axis leads to \[ n \]. Rotating the area between \( y = f(x) \) and \( y = g(x) \) around the \( x \) axis, the slices look like \[ o \]. Their areas are \[ p \] so the volume is \[ q \].

Another method is to cut the solid into thin cylindrical \[ r \]. Revolving the area under \( y = f(x) \) around the \( y \) axis, a shell has height \[ s \] and thickness \( dx \) and volume \[ t \]. The total volume is \[ u \].

Find where the curves in 1–12 intersect, draw rough graphs, and compute the area between them.

1. \( y = x^2 - 3 \) and \( y = 1 \)
2. \( y = x^2 - 2 \) and \( y = 0 \)
3. \( y^2 = x \) and \( x = 9 \)
4. \( y^2 = x \) and \( x = y + 2 \)
5. \( y = x^4 - 2x^2 \) and \( y = 2x^2 \)
6. \( x = y^5 \) and \( y = x^4 \)
7. \( y = x^2 \) and \( y = -x^2 + 18x \)
8. \( y = 1/x \) and \( y = 1/x^2 \) and \( x = 3 \)
9. \( y = \cos x \) and \( y = \cos^2 x \)
10. \( y = \sin \pi x \) and \( y = 2x \) and \( x = 0 \)
11. \( y = e^x \) and \( y = e^{2x-1} \) and \( x = 0 \)
12. \( y = e^x \) and \( y = e^x \) and \( y = e^{-x} \)
13. Find the area inside the three lines \( y = 4 - x \), \( y = 3x \), and \( y = x \).
14. Find the area bounded by \( y = 12 - x \), \( y = \sqrt{x} \), and \( y = 1 \).
15. Does the parabola \( y = 1 - x^2 \) out to \( x = 1 \) sit inside or outside the unit circle \( x^2 + y^2 = 1 \)? Find the area of the "skin" between them.
Find the area of the largest triangle with base on the x axis that fits (a) inside the unit circle (b) inside that parabola.

Rotate the ellipse \( \frac{x^2}{a^2} + \frac{y^2}{b^2} = 1 \) around the x axis to find the volume of a football. What is the volume around the y axis? If \( a = 2 \) and \( b = 1 \), locate a point \((x, y, z)\) that is in one football but not the other.

What is the volume of the loaf of bread which comes from rotating \( y = \sin x \) (0 < x < \( \pi \)) around the x axis? What is the volume around the y axis?

What is the volume common to two circular cylinders of radius \( a \). One eighth of the region is shown (axes are perpendicular and horizontal slices are squares).

In 29–34 find the volume and draw a typical slice.

A cap of height \( h \) is cut off the top of a sphere of radius \( R \). Slice the sphere horizontally starting at \( y = R - h \).

A pyramid \( P \) has height 6 and square base of side 2. Its volume is \( \frac{1}{3}(6)(2)^2 = 8 \).

(a) Find the volume up to height 3 by horizontal slices. What is the length of a side at height \( y \)?

(b) Recompute by removing a smaller pyramid from \( P \).

The base is a disk of radius \( a \). Slices perpendicular to the base are squares.

The base is the region under the parabola \( y = 1 - x^2 \). Slices perpendicular to the x axis are squares.

The base is the region under the parabola \( y = 1 - x^2 \). Slices perpendicular to the y axis are squares.

The base is the triangle with corners \((0, 0), (1, 0), (0, 1)\). Slices perpendicular to the x axis are semicircles.

Cavalieri's principle for areas: If two regions have strips of equal length, then the regions have the same area. Draw a parallelogram and a curved region, both with the same strips as the unit square. Why are the areas equal?

Cavalieri's principle for volumes: If two solids have slices of equal area, the solids have the same volume. Find the volume of the tilted cylinder in the figure.

Draw another region with the same slice areas as the tilted cylinder. When all areas \( A(x) \) are the same, the volumes \( \int A(x) \) are the same.

Find the volume of the flying saucer that comes from rotating \( y = \sin x \) (0 < x < \( \pi \)) around the y axis.

Find the volume of the galaxy that comes from rotating \( y = \sin x \) (0 < x < \( \pi \)) around the x axis and then rotating the whole thing around the y axis.

Draw the region bounded by the curves in 21–28. Find the volume when the region is rotated (a) around the x axis (b) around the y axis.

In 29–34 find the volume and draw a typical slice.

A cap of height \( h \) is cut off the top of a sphere of radius \( R \). Slice the sphere horizontally starting at \( y = R - h \).

A pyramid \( P \) has height 6 and square base of side 2. Its volume is \( \frac{1}{3}(6)(2)^2 = 8 \).

(a) Find the volume up to height 3 by horizontal slices. What is the length of a side at height \( y \)?

(b) Recompute by removing a smaller pyramid from \( P \).

The base is a disk of radius \( a \). Slices perpendicular to the base are squares.

The base is the region under the parabola \( y = 1 - x^2 \). Slices perpendicular to the x axis are squares.

The base is the region under the parabola \( y = 1 - x^2 \). Slices perpendicular to the y axis are squares.

The base is the triangle with corners \((0, 0), (1, 0), (0, 1)\). Slices perpendicular to the x axis are semicircles.

Cavalieri's principle for areas: If two regions have strips of equal length, then the regions have the same area. Draw a parallelogram and a curved region, both with the same strips as the unit square. Why are the areas equal?

Cavalieri's principle for volumes: If two solids have slices of equal area, the solids have the same volume. Find the volume of the tilted cylinder in the figure.

Draw another region with the same slice areas as the tilted cylinder. When all areas \( A(x) \) are the same, the volumes \( \int A(x) \) are the same.

Find the volume of the galaxy that comes from rotating \( y = \sin x \) (0 < x < \( \pi \)) around the x axis and then rotating the whole thing around the y axis.

In 29–34 find the volume and draw a typical slice.

A cap of height \( h \) is cut off the top of a sphere of radius \( R \). Slice the sphere horizontally starting at \( y = R - h \).

A pyramid \( P \) has height 6 and square base of side 2. Its volume is \( \frac{1}{3}(6)(2)^2 = 8 \).

(a) Find the volume up to height 3 by horizontal slices. What is the length of a side at height \( y \)?

(b) Recompute by removing a smaller pyramid from \( P \).

The base is a disk of radius \( a \). Slices perpendicular to the base are squares.

The base is the region under the parabola \( y = 1 - x^2 \). Slices perpendicular to the x axis are squares.

The base is the region under the parabola \( y = 1 - x^2 \). Slices perpendicular to the y axis are squares.

The base is the triangle with corners \((0, 0), (1, 0), (0, 1)\). Slices perpendicular to the x axis are semicircles.

Cavalieri's principle for areas: If two regions have strips of equal length, then the regions have the same area. Draw a parallelogram and a curved region, both with the same strips as the unit square. Why are the areas equal?

Cavalieri's principle for volumes: If two solids have slices of equal area, the solids have the same volume. Find the volume of the tilted cylinder in the figure.

Draw another region with the same slice areas as the tilted cylinder. When all areas \( A(x) \) are the same, the volumes \( \int A(x) \) are the same.

Find the volume of the galaxy that comes from rotating \( y = \sin x \) (0 < x < \( \pi \)) around the x axis and then rotating the whole thing around the y axis.

In 29–34 find the volume and draw a typical slice.

A cap of height \( h \) is cut off the top of a sphere of radius \( R \). Slice the sphere horizontally starting at \( y = R - h \).

A pyramid \( P \) has height 6 and square base of side 2. Its volume is \( \frac{1}{3}(6)(2)^2 = 8 \).

(a) Find the volume up to height 3 by horizontal slices. What is the length of a side at height \( y \)?

(b) Recompute by removing a smaller pyramid from \( P \).

The base is a disk of radius \( a \). Slices perpendicular to the base are squares.

The base is the region under the parabola \( y = 1 - x^2 \). Slices perpendicular to the x axis are squares.

The base is the region under the parabola \( y = 1 - x^2 \). Slices perpendicular to the y axis are squares.

The base is the triangle with corners \((0, 0), (1, 0), (0, 1)\). Slices perpendicular to the x axis are semicircles.

Cavalieri's principle for areas: If two regions have strips of equal length, then the regions have the same area. Draw a parallelogram and a curved region, both with the same strips as the unit square. Why are the areas equal?

Cavalieri's principle for volumes: If two solids have slices of equal area, the solids have the same volume. Find the volume of the tilted cylinder in the figure.

Draw another region with the same slice areas as the tilted cylinder. When all areas \( A(x) \) are the same, the volumes \( \int A(x) \) are the same.

Find the volume of the galaxy that comes from rotating \( y = \sin x \) (0 < x < \( \pi \)) around the x axis and then rotating the whole thing around the y axis.

In 29–34 find the volume and draw a typical slice.

A cap of height \( h \) is cut off the top of a sphere of radius \( R \). Slice the sphere horizontally starting at \( y = R - h \).

A pyramid \( P \) has height 6 and square base of side 2. Its volume is \( \frac{1}{3}(6)(2)^2 = 8 \).

(a) Find the volume up to height 3 by horizontal slices. What is the length of a side at height \( y \)?

(b) Recompute by removing a smaller pyramid from \( P \).

The base is a disk of radius \( a \). Slices perpendicular to the base are squares.

The base is the region under the parabola \( y = 1 - x^2 \). Slices perpendicular to the x axis are squares.

The base is the region under the parabola \( y = 1 - x^2 \). Slices perpendicular to the y axis are squares.

The base is the triangle with corners \((0, 0), (1, 0), (0, 1)\). Slices perpendicular to the x axis are semicircles.

Cavalieri's principle for areas: If two regions have strips of equal length, then the regions have the same area. Draw a parallelogram and a curved region, both with the same strips as the unit square. Why are the areas equal?
The slices in 43–46 are washers. Find the slice area and volume.

43 The rectangle with sides $x = 1, x = 3, y = 2, y = 5$ is rotated around the $x$ axis.

44 The same rectangle is rotated around the $y$ axis.

45 The same rectangle is rotated around the line $y = 1$.

46 Draw the triangle with corners $(1, 0), (1, 1), (0, 1)$. After rotation around the $x$ axis, describe the solid and find its volume.

47 Bore a hole of radius $a$ down the axis of a cone and through the base of radius $b$. If it is a 45° cone (height also $b$), what volume is left? Check a $a$ and $b$.

48 Find the volume common to two spheres of radius $r$ if the $y$ axis. (a) Find the volume by cutting into shells. (b) Find the volume by slicing into washers.

49 (Shells vs. disks) Rotate $y = 3 - x$ around the $x$ axis from $x = 0$ to $x = 2$. Write down the volume integral by disks and then by shells.

50 (Shells vs. disks) Rotate $y = x^3$ around the $y$ axis from $y = 0$ to $y = 8$. Write down the volume integral by shells and disks and compute both ways.

51 Yogurt comes in a solid of revolution. Rotate the line $y = mx$ around the $y$ axis to find the volume between $y = a$ and $y = b$.

52 Suppose $y = f(x)$ decreases from $f(0) = b$ to $f(1) = 0$. The curve is rotated around the $y$ axis. Compare shells to disks:

$$\int_0^b 2\pi xy(x)\, dx = \int_0^1 \pi (f^{-1}(y))^2 \, dy.$$  

Substitute $y = f(x)$ in the second. Also substitute $dy = f'(x)\, dx$. Integrate by parts to reach the first.

53 If a roll of paper with inner radius 2 cm and outer radius 10 cm has about 10 thicknesses per centimeter, approximately how long is the paper when unrolled?

54 Find the approximate volume of your brain. OK to include everything above your eyes (skull too).

Use shells to find the volumes in 55–63. The rotated regions lie between the curve and $x$ axis.

55 $y = 1 - x^2, 0 \leq x \leq 1$ (around the $y$ axis)

56 $y = 1/x, 1 \leq x \leq 100$ (around the $y$ axis)

57 $y = \sqrt{1 - x^2}, 0 \leq x \leq 1$ (around either axis)

58 $y = 1/(1 + x^2), 0 \leq x \leq 3$ (around the $y$ axis)

59 $y = \sin(x^2), 0 \leq x \leq \sqrt{\pi}$ (around the $y$ axis)

60 $y = 1/\sqrt{1 - x^2}, 0 \leq x \leq 1$ (around the $y$ axis)

61 $y = x^2, 0 \leq x \leq 2$ (around the $x$ axis)

62 $y = e^x, 0 \leq x \leq 1$ (around the $x$ axis)

63 $y = \ln x, 1 \leq x \leq e$ (around the $x$ axis)

64 The region between $y = x^2$ and $y = x$ is revolved around the $y$ axis. (a) Find the volume by cutting into shells. (b) Find the volume by slicing into washers.

65 The region between $y = f(x)$ and $y = 1 + f(x)$ is revolved around the $y$ axis. The shells have height _______ . The volume out to $x = a$ is _______ because the shells are the same.

66 A horizontal slice of the sphere-with-hole in Figure 8.6b is a washer. Its area is $\pi(x^2 - a^2) = \pi(b^2 - y^2 - a^2)$.

(a) Find the upper limit on $y$ (the top of the hole).

(b) Integrate the area to verify the volume in Example 12.

67 If the hole in the sphere has length 2, show that the volume is $4\pi/3$ regardless of the radii $a$ and $b$.

68 An upright cylinder of radius $r$ is sliced by two parallel planes at angle $\alpha$. One is a height $h$ above the other.

(a) Draw a picture to show that the volume between the planes is $\pi r^2 h$.

(b) Draw the picture by $\alpha$, so the base and top are flat. What is the shape of the base? What is its area $A$? What is the height $H$ of the tilted cylinder?

69 True or false, with a reason.

(a) A cube can only be sliced into squares.

(b) A cube cannot be cut into cylindrical shells.

(c) The washer with radii $r$ and $R$ has area $\pi(R^2 - r^2)$.

(d) The plane $w = \frac{1}{2}$ slices a 3-dimensional sphere out of a 4-dimensional sphere $x^2 + y^2 + z^2 + w^2 = 1$.

---

### 8.2 Length of a Plane Curve

The graph of $y = x^{3/2}$ is a curve in the $x$-$y$ plane. **How long is that curve?** A definite integral needs endpoints, and we specify $x = 0$ and $x = 4$. The first problem is to know what "length function" to integrate.

The distance along a curve is the **arc length**. To set up an integral, we break the
problem into small pieces. Roughly speaking, *small pieces of a smooth curve are nearly straight*. We know the exact length $\Delta s$ of a straight piece, and Figure 8.8 shows how it comes close to a curved piece.

\[ (\Delta s)^2 = (\Delta x)^2 + (\Delta y)^2 \]

\[ ds^2 = (dx)^2 + \left( \frac{dy}{dx} \right)^2 (dx)^2 \]

\[ ds = \sqrt{1 + \left( \frac{dy}{dx} \right)^2} \, dx \]

![Fig. 8.8 Length $\Delta s$ of short straight segment. Length $ds$ of very short curved segment.](image)

Here is the unofficial reasoning that gives the length of the curve. A straight piece has $(\Delta s)^2 = (\Delta x)^2 + (\Delta y)^2$. Within that right triangle, the height $\Delta y$ is the slope $(\Delta y/\Delta x)$ times $\Delta x$. This secant slope is close to the slope of the curve. Thus $\Delta y$ is approximately $(dy/dx) \Delta x$.

\[ \Delta s \approx \sqrt{(\Delta x)^2 + (dy/dx)^2(\Delta x)^2} = \sqrt{1 + (dy/dx)^2} \, \Delta x. \]  
(1)

Now add these pieces and make them smaller. The infinitesimal triangle has $(ds)^2 = (dx)^2 + (dy)^2$. Think of $ds$ as $\sqrt{1 + (dy/dx)^2} \, dx$ and integrate:

\[ \text{length of curve} = \int ds = \int \sqrt{1 + (dy/dx)^2} \, dx. \]  
(2)

**EXAMPLE 1** Keep $y = x^{3/2}$ and $dy/dx = \frac{3}{2}x^{1/2}$. Watch out for $\frac{3}{2}$ and $\frac{9}{4}$:

\[ \text{length} = \int_0^4 \sqrt{1 + \frac{9}{4}x} \, dx = \left( \frac{3}{2} \right) \left( \frac{9}{4} \right) \int_0^4 \left( 1 + \frac{9}{4}x \right)^{1/2} \, dx = \frac{9}{10} \left( 10^{3/2} - 1^{3/2} \right). \]  
(3)

This answer is just above $9$. A straight line from $(0,0)$ to $(4,8)$ has exact length $\sqrt{80}$. Note $4^2 + 8^2 = 80$. Since $\sqrt{80}$ is just below $9$, the curve is surprisingly straight.

You may not approve of those numbers (or the reasoning behind them). We can fix the reasoning, but nothing can be done about the numbers. This example $y = x^{3/2}$ had to be chosen carefully to make the integration possible at all. The length integral is difficult because of the square root. In most cases we integrate numerically.

**EXAMPLE 2** The straight line $y = 2x$ from $x = 0$ to $x = 4$ has $dy/dx = 2$:

\[ \text{length} = \int_0^4 \sqrt{1 + 4} \, dx = 4\sqrt{5} = \sqrt{80} \text{ as before} \quad \text{(just checking).} \]

We return briefly to the reasoning. The curve is the graph of $y = f(x)$. Each piece contains at least one point where secant slope equals tangent slope: $\Delta y/\Delta x = f'(c)$. The Mean Value Theorem applies when the slope is continuous—this is required for a smooth curve. The straight length $\Delta s$ is exactly $\sqrt{(\Delta x)^2 + (f'(c)\Delta x)^2}$. Adding
the \( n \) pieces gives the length of the broken line (close to the curve):

\[
\Delta y = f'(c) \Delta x \\
\sum_{i=1}^{n} \Delta x_i = \sum_{i=1}^{n} \sqrt{1 + [f'(c)]^2} \Delta x_i.
\]

As \( n \to \infty \) and \( \Delta x_{\text{max}} \to 0 \) this approaches the integral that gives arc length.

**8A** The length of the curve \( y = f(x) \) from \( x = a \) to \( x = b \) is

\[
s = \int_{a}^{b} \sqrt{1 + [f'(x)]^2} \, dx = \int_{a}^{b} \sqrt{1 + (dy/dx)^2} \, dx.
\]

**EXAMPLE 3** Find the length of the first quarter of the circle \( y = \sqrt{1 - x^2} \).

Here \( dy/dx = -x/\sqrt{1 - x^2} \). From Figure 8.9a, the integral goes from \( x = 0 \) to \( x = 1 \):

\[
\text{length} = \int_{0}^{1} \sqrt{1 + (dy/dx)^2} \, dx = \int_{0}^{1} \sqrt{1 + \frac{x^2}{1-x^2}} \, dx = \int_{0}^{1} \frac{dx}{\sqrt{1-x^2}}.
\]

The antiderivative is \( \sin^{-1} x \). It equals \( \frac{\pi}{2} \) at \( x = 1 \). This length \( \frac{\pi}{2} \) is a quarter of the full circumference \( 2\pi \).

**EXAMPLE 4** Compute the distance around a quarter of the ellipse \( y^2 + 2x^2 = 2 \).

The equation is \( y = \sqrt{2 - 2x^2} \) and the slope is \( dy/dx = -2x/\sqrt{2 - 2x^2} \). So \( \int ds \) is

\[
\int_{0}^{1} \sqrt{1 + \frac{4x^2}{2 - 2x^2}} \, dx = \int_{0}^{1} \sqrt{2 + 2x^2} \, dx = \int_{0}^{1} \sqrt{\frac{1 + x^2}{1 - x^2}} \, dx.
\]

That integral can't be done in closed form. The length of an ellipse can only be computed numerically. The denominator is zero at \( x = 1 \), so a blind application of the trapezoidal rule or Simpson's rule would give length = \( \infty \). The midpoint rule gives length = 1.91 with thousands of intervals.

**Fig. 8.9** Circle and ellipse, directly by \( y = f(x) \) or parametrically by \( x(t) \) and \( y(t) \).

**LENGTH OF A CURVE FROM PARAMETRIC EQUATIONS: \( x(t) \) AND \( y(t) \)**

We have met the unit circle in two forms. One is \( x^2 + y^2 = 1 \). The other is \( x = \cos t, \\
y = \sin t \). Since \( \cos^2 t + \sin^2 t = 1 \), this point goes around the correct circle. One advantage of the "parameter" \( t \) is to give extra information—it tells *where* the point is and
also when. In Chapter 1, the parameter was the time and also the angle—because we moved around the circle with speed 1.

Using \( t \) is a natural way to give the position of a particle or a spacecraft. We can recover the velocity if we know \( x \) and \( y \) at every time \( t \). An equation \( y = f(x) \) tells the shape of the path, not the speed along it.

Chapter 12 deals with parametric equations for curves. Here we concentrate on the path length—which allows you to see the idea of a parameter \( t \) without too much detail. We give \( x \) as a function of \( t \) and \( y \) as a function of \( t \). The curve is still approximated by straight pieces, and each piece has \((\Delta s)^2 = (\Delta x)^2 + (\Delta y)^2 \). But instead of using \( \Delta y \approx (dy/dx) \Delta x \), we approximate \( \Delta x \) and \( \Delta y \) separately:

\[
\Delta x \approx (dx/dt) \Delta t, \quad \Delta y \approx (dy/dt) \Delta t, \quad \Delta s \approx \sqrt{(dx/dt)^2 + (dy/dt)^2} \Delta t.
\]

**8.2 Length of a Plane Curve**

The length of a parametric curve is an integral with respect to \( t \):

\[
\int ds = \int (ds/dt) dt = \int \sqrt{(dx/dt)^2 + (dy/dt)^2} dt.
\]  

**EXAMPLE 5** Find the length of the quarter-circle using \( x = \cos t \) and \( y = \sin t \):

\[
\int_0^{\pi/2} \sqrt{(dx/dt)^2 + (dy/dt)^2} dt = \int_0^{\pi/2} \sqrt{\sin^2 t + \cos^2 t} dt = \int_0^{\pi/2} dt = \frac{\pi}{2}.
\]

The integral is simpler than \( 1/\sqrt{1-x^2} \), and there is one new advantage. We can integrate around a whole circle with no trouble. Parametric equations allow a path to close up or even cross itself. The time \( t \) keeps going and the point \((x(t), y(t))\) keeps moving. In contrast, curves \( y = f(x) \) are limited to one \( y \) for each \( x \).

**EXAMPLE 6** Find the length of the quarter-ellipse: \( x = \cos t \) and \( y = \sqrt{2} \sin t \):

On this path \( y^2 + 2x^2 = 2 \sin^2 t + 2 \cos^2 t = 2 \) (same ellipse). The non-parametric equation \( y = \sqrt{2 - 2x^2} \) comes from eliminating \( t \). We keep \( t \):

\[
\text{length} = \int_0^{\pi/2} \sqrt{(dx/dt)^2 + (dy/dt)^2} dt = \int_0^{\pi/2} \sqrt{\sin^2 t + 2 \cos^2 t} dt.
\]  

This integral (7) must equal (5). If one cannot be done, neither can the other. They are related by \( x = \cos t \), but (7) does not blow up at the endpoints. The trapezoidal rule gives 1.9101 with less than 100 intervals. Section 5.8 mentioned that calculators automatically do a substitution that makes (5) more like (7).

**EXAMPLE 7** The path \( x = t^2, y = t^3 \) goes from (0, 0) to (4, 8). Stop at \( t = 2 \).

To find this path without the parameter \( t \), first solve for \( t = x^{1/2} \). Then substitute into the equation for \( y \): \( y = t^3 = x^{3/2} \). The non-parametric form (with \( t \) eliminated) is the same curve \( y = x^{3/2} \) as in Example 1.

The length from the \( t \)-integral equals the length from the \( x \)-integral. This is Problem 22.

**EXAMPLE 8** Special choice of parameter: \( t \) is \( x \). The curve becomes \( x = t, y = t^{3/2} \).

If \( x = t \) then \( dx/dt = 1 \). The square root in (6) is the same as the square root in (4). Thus the non-parametric form \( y = f(x) \) is a special case of the parametric form—just take \( t = x \).

Compare \( x = t, y = t^{3/2} \) with \( x = t^2, y = t^3 \). Same curve, same length, different speed.
EXAMPLE 9  Define "speed" by
\[
\frac{ds}{dt} = \sqrt{\left(\frac{dx}{dt}\right)^2 + \left(\frac{dy}{dt}\right)^2}
\]

When a ball is thrown straight upward, \(dx/dt\) is zero. But the speed is not \(dy/dt\).
It is \(|dy/dt|\). The speed is positive downward as well as upward.

8.2 EXERCISES

Read-through questions
The length of a straight segment (\(\Delta x\) across, \(\Delta y\) up) is
\(\Delta s = \sqrt{\Delta x^2 + \Delta y^2}\). Between two points of the graph of \(y(x)\), \(\Delta y\) is
approximately \(dy/dx\) times \(\Delta x\). The length of that piece
is approximately \(\sqrt{1 + (dy/dx)^2}\). An infinitesimal piece of
the curve has length \(ds = \sqrt{1 + (dy/dx)^2}\). Then the arc length integral
is \(\int ds\).

For \(y = 4 - x\) from \(x = 0\) to \(x = 3\) the arc length is
\(\int_0^3 \sqrt{1 + (-1)^2} \, dx = 3\). For \(y = x^3\) the arc length integral is \(\int_0^1 \sqrt{1 + (3x^2)^2} \, dx\).

The curve \(x = \cos t, y = \sin t\) is the same as \(i\). The length of a curve given by
\((x(t), y(t))\) is \(\int \sqrt{\left(\frac{dx}{dt}\right)^2 + \left(\frac{dy}{dt}\right)^2} \, dt\). For example
\(x = \cos t, \ y = \sin t\) from \(t = \pi/3\) to \(t = \pi/2\) has length
\(\frac{\sqrt{2}}{2}\). The speed is \(ds/dt = \frac{\sqrt{2}}{2}\). For the special case
\(x = t, \ y = f(t)\) the length formula goes back to \(\int \sqrt{1 + f'(x)^2} \, dx\).

Find the lengths of the curves in Problems 1-8.
1) \(y = \frac{1}{3}x^{1/2}\) from \((0, 0)\) to \((1, 1)\)
2) \(y = \frac{1}{3}x^{1/3}\) from \((0, 0)\) to \((1, 1)\) (compare with Problem 1 or
put \(u = \frac{1}{3}x^{1/3}\) in the length integral)
3) \(y = \frac{1}{3}(x^2 + 2)^{1/2}\) from \(x = 0\) to \(x = 1\)
4) \(y = \frac{1}{3}(x^2 - 2)^{1/2}\) from \(x = 2\) to \(x = 4\)
5) \(y = \frac{x^3}{3} + \frac{1}{4x}\) from \(x = 1\) to \(x = 3\)
6) \(y = \frac{x^3}{4} + \frac{1}{8x^2}\) from \(x = 1\) to \(x = 2\)
7) \(y = \frac{1}{2}x^{1/2} - \frac{1}{2}x^{1/2}\) from \(x = 1\) to \(x = 4\)
8) \(y = x^3\) from \((0, 0)\) to \((1, 1)\)

The curve given by \(x = \cos^3 t, \ y = \sin^3 t\) is an astroid (a
hypocycloid). Its non-parametric form is \(x^{2/3} + y^{2/3} = 1\).
Sketch the curve from \(t = 0\) to \(t = \pi/2\) and find its length.

10) Find the length from \(t = 0\) to \(t = \pi\) of the curve given by
\(x = \cos t + \sin t, \ y = \cos t - \sin t\). Show that the curve is a
circle (of what radius)?

11) Find the length from \(t = 0\) to \(t = \pi/2\) of the curve given by
\(x = \cos t, \ y = 1 - \sin t\).

12) What integral gives the length of Archimedes' spiral
\(x = t \cos t, \ y = t \sin t\)?

13) Find the distance traveled in the first second (to \(t = 1\)) if
\(x = \frac{1}{2}t^2, \ y = \frac{1}{4}(2t + 1)^{3/2}\).

14) \(x = (1 - \frac{1}{2} \cos 2t) \cos t + y = (1 + \frac{1}{2} \cos 2t) \sin t\) lead to
\(4(1 - x^3 - y^3) = 27(x^2 - y^2)^2\). Find the arc length from \(t = 0\) to
\(\pi/4\).

Find the arc lengths in 15-18 by numerical integration.
15) One arch of \(y = \sin x\), from \(x = 0\) to \(x = \pi\).
16) \(y = e^x\) from \(x = 0\) to \(x = 1\).
17) \(y = \ln x\) from \(x = 1\) to \(x = e\).
18) \(x = \cos t, \ y = 3 \sin t\), \(0 \leq x \leq 2\pi\).
19) Draw a rough picture of \(y = x^{10}\). Without computing the
length of \(y = x^{10}\) from \((0, 0)\) to \((1, 1)\), find the limit as \(n \to \infty\).

20) Which is longer between \((1, 1)\) and \((2, \frac{1}{2})\), the hyperbola
\(y = 1/x\) or the graph of \(x + 2y = 3\)?

21) Find the speed \(ds/dt\) on the circle \(x = 2 \cos 3t, \ y = 2 \sin 3t\).

22) Examples 1 and 7 were \(y = x^{3/2}\) and \(x = t^2, \ y = t^3\):

\[\text{length} = \int_0^1 \sqrt{1 + \frac{3}{2}x} \, dx, \text{length} = \int_0^1 \sqrt{4t^2 + 9t^4} \, dt.\]

Show by substituting \(x = \frac{1}{2}t\) that these integrals agree.

23) Instead of \(y = f(x)\) a curve can be given as \(x = g(y)\). Then
\(ds = \sqrt{(dx/dx)^2 + (dy/dy)^2} = \sqrt{(dx/dy)^2 + 1} \, dy\).

Draw \(x = 5y\) from \(y = 0\) to \(y = 1\) and find its length.

24) The length of \(x = y^{3/2}\) from \((0, 0)\) to \((1, 1)\) is
\(\int ds = \sqrt{2} + 1 \, dy\). Compare with Problem 1: Same length?
Same curve?

25) Find the length of \(x = \frac{1}{2}(e^x + e^{-x})\) from \(y = -1\) to \(y = 1\)
and draw the curve.

26) The length of \(x = g(y)\) is a special case of equation (6) with
\(y = t\) and \(x = g(t)\). The length integral becomes ________.

27) Plot the point \(x = 3 \cos t, \ y = 4 \sin t\) at the five times
\(t = 0, \pi/2, \pi, 3\pi/2, 2\pi\). The equation of the curve is
\((x/3)^2 + (y/4)^2 = 1\), not a circle but an ________.

28) (a) Find the length of \(x = \cos^2 t, \ y = \sin^2 t\), \(0 \leq y \leq \pi\).
(b) Why does this path stay on the line \(x + y = 1\)?
(c) Why isn't the path length equal to \(\sqrt{2}\)?
29. (important) The line $y = x$ is close to a staircase of pieces that go straight across or straight up. With 100 pieces of length $\Delta x = 1/100$ or $\Delta y = 1/100$, find the length of carpet on the staircase. (The length of the 45° line is $\sqrt{2}$. The staircase can be close when its length is not close.)

30. The area of an ellipse is $\pi ab$. The area of a strip around it (width $\Delta$) is $\pi(a + \Delta)(b + \Delta) - \pi ab \approx \pi(a + b)\Delta$. The distance around the ellipse seems to be $\pi(a + b)$. But this distance is impossible to find—what is wrong?

31. The point $x = \cos t, y = \sin t, z = t$ moves on a space curve.
   
   (a) In three-dimensional space $(dx)^2$ equals $(dx)^2 + \ldots$. In equation (6), $ds$ is now _______ $dt$.

(b) This particular curve has $ds = \ldots$. Find its length from $t = 0$ to $t = 2\pi$.

(c) Describe the curve and its shadow in the $xy$ plane.

32. Explain in 50 words the difference between a non-parametric equation $y = f(x)$ and two parametric equations $x = x(t), y = y(t)$.

33. Write down the integral for the length $L$ of $y = x^2$ from $(0, 0)$ to $(1, 1)$. Show that $y = \frac{1}{4}x^2$ from $(0, 0)$ to $(2, 2)$ is exactly twice as long. If possible give a reason using the graphs.

34. (for professors) Compare the lengths of the parabola $y = x^2$ and the line $y = bx$ from $(0, 0)$ to $(b, b^2)$. Does the difference approach a limit as $b \to \infty$?

---

### 8.3 Area of a Surface of Revolution

This section starts by constructing surfaces. A **curve** $y = f(x)$ is revolved around an **axis**. That produces a "**surface of revolution**," which is symmetric around the axis. If we revolve a sloping line, the result is a cone. When the line is parallel to the axis we get a cylinder (a pipe). By revolving a curve we might get a lamp or a lamp shade (or even the light bulb).

Section 8.1 computed the volume inside that surface. **This section computes the surface area**. Previously we cut the solid into slices or shells. Now we need a good way to cut up the surface.

The key idea is **to revolve short straight line segments**. Their slope is $\Delta y/\Delta x$. They can be the same pieces of length $\Delta s$ that were used to find length—now we compute area. When revolved, a straight piece produces a **"thin band"** (Figure 8.10). The curved surface, from revolving $y = f(x)$, is close to the bands. The first step is to compute **the surface area of a band**.

A small comment: Curved surfaces can also be cut into tiny patches. Each patch is nearly flat, like a little square. The sum of those patches leads to a double integral (with $dx\,dy$). Here the integral stays one-dimensional ($dx$ or $dy$ or $dt$). Surfaces of revolution are special—we approximate them by bands that go all the way around. A band is just a belt with a slope, and its slope has an effect on its area.

---

Fig. 8.10  Revolving a straight piece and a curve around the $y$ axis and $x$ axis.
Revolve a small straight piece (length \( \Delta s \) not \( \Delta x \)). The center of the piece goes around a circle of radius \( r \). The band is a slice of a cone. When we flatten it out (Problems 11–13) we discover its area. The area is the side length \( \Delta s \) times the middle circumference \( 2\pi r \):

**The surface area of a band is**

\[
2\pi r \Delta s = 2\pi r \sqrt{1 + (\Delta y/\Delta x)^2} \Delta x.
\]

For revolution around the \( y \) axis, the radius is \( r = x \). For revolution around the \( x \) axis, the radius is the height: \( r = y = f(x) \). Figure 8.10 shows both bands—the problem tells us which to use. The sum of band areas \( 2\pi r \Delta s \) is close to the area \( S \) of the curved surface. In the limit we integrate \( 2\pi r \, ds \):

\[
\text{(8C)} \quad \text{The surface area generated by revolving the curve } y = f(x) \text{ between } x = a \text{ and } x = b \text{ is}
\]

\[
S = \int_a^b 2\pi y \sqrt{1 + (dy/dx)^2} \, dx \quad \text{around the } x \text{ axis (} r = y \) \quad (1)
\]

\[
S = \int_a^b 2\pi x \sqrt{1 + (dy/dx)^2} \, dx \quad \text{around the } y \text{ axis (} r = x \). \quad (2)
\]

**EXAMPLE 1** Revolve a complete semicircle \( y = \sqrt{R^2 - x^2} \) around the \( x \) axis.

The surface of revolution is a sphere. Its area (known!) is \( 4\pi R^2 \). The limits on \( x \) are \( -R \) and \( R \). The slope of \( y = \sqrt{R^2 - x^2} \) is \( dy/dx = -x/\sqrt{R^2 - x^2} \):

\[
\text{area } S = \int_{-R}^R 2\pi \sqrt{R^2 - x^2} \sqrt{1 + \frac{x^2}{R^2 - x^2}} \, dx = \int_{-R}^R 2\pi R \, dx = 4\pi R^2.
\]

**EXAMPLE 2** Revolve a piece of the straight line \( y = 2x \) around the \( x \) axis.

The surface is a cone with \( (dy/dx)^2 = 4 \). The band from \( x = 0 \) to \( x = 1 \) has area \( 2\pi \sqrt{5} \):

\[
S = \int 2\pi y \, ds = \int_0^1 2\pi(2x) \sqrt{1 + 4} \, dx = 2\pi \sqrt{5}.
\]

This answer must agree with the formula \( 2\pi r \Delta s \) (which it came from). The line from \((0, 0)\) to \((1, 2)\) has length \( \Delta s = \sqrt{5} \). Its midpoint is \((\frac{1}{2}, 1)\). Around the \( x \) axis, the middle radius is \( r = 1 \) and the area is \( 2\pi \sqrt{5} \).

**EXAMPLE 3** Revolve the same straight line segment around the \( y \) axis. Now the radius is \( x \) instead of \( y = 2x \). The area in Example 2 is cut in half:

\[
S = \int 2\pi x \, ds = \int_0^1 2\pi x \sqrt{1 + 4} \, dx = \pi \sqrt{5}.
\]

For surfaces as for arc length, only a few examples have convenient answers. Watermelons and basketballs and light bulbs are in the exercises. Rather than stretching out this section, we give a final area formula and show how to use it.

The formula applies when there is a parameter \( t \). Instead of \((x, f(x))\) the points on the curve are \((x(t), y(t))\). As \( t \) varies, we move along the curve. The length formula \( (ds)^2 = (dx)^2 + (dy)^2 \) is expressed in terms of \( t \).

For the surface of revolution around the \( x \) axis, the area becomes a \( t \)-integral:

\[
\text{(8D)} \quad \text{The surface area is } S = \int 2\pi y(t) \sqrt{(dx/dt)^2 + (dy/dt)^2} \, dt. \quad (3)
\]
8.3 Area of a Surface of Revolution

EXAMPLE 4 The point \( x = \cos t, \ y = 5 + \sin t \) travels on a circle with center at \((0, 5)\). Revolving that circle around the x axis produces a doughnut. Find its surface area.

Solution \( \left( \frac{dx}{dt} \right)^2 + \left( \frac{dy}{dt} \right)^2 = \sin^2 t + \cos^2 t = 1 \). The circle is complete at \( t = 2\pi \):

\[
\int 2\pi y \, ds = \int_{0}^{2\pi} 2\pi (5 + \sin t) \, dt = [2\pi (5t - \cos t)]_{0}^{2\pi} = 20\pi^2.
\]

8.3 EXERCISES

Read-through questions

A surface of revolution comes from revolving a _______ around a _______. This section computes the _______. When the curve is a short straight piece (length \( \Delta s \)), the surface is a _______. Its area is \( \Delta S = \) _______. In that formula (Problem 13) \( r \) is the radius of _______. The line from \((0, 0)\) to \((1, 1)\) has length _______. For the curve given by \( x = 2t, \ y = t^2 \), change \( ds \) to _______.

- Find the surface area when curves 1–6 revolve around the x axis.
  1. \( y = \sqrt{x}, \ 2 \leq x \leq 6 \)
  2. \( y = x^3, \ 0 \leq x \leq 1 \)
  3. \( y = 7x, \ -1 \leq x \leq 1 \) (watch sign)
  4. \( y = \sqrt{4 - x^2}, \ 0 \leq x \leq 2 \)
  5. \( y = \sqrt{4 - x^2}, \ -1 \leq x \leq 1 \)
  6. \( y = \cosh x, \ 0 \leq x \leq 1 \)

- In 7–10 find the area of the surface of revolution around the y axis.
  7. \( y = x^2, \ 0 \leq x \leq 2 \)
  8. \( y = \frac{1}{2} x^2 + \frac{1}{2}, \ 0 \leq x \leq 1 \)
  9. \( y = x + 1, \ 0 \leq x \leq 3 \)
  10. \( y = x^3, \ 0 \leq x \leq 1 \)

- A cone with base radius \( R \) and slant height \( s \) is laid out flat. Explain why the angle (in radians) is \( \theta = 2\pi R/s \). Then the surface area is a fraction of a circle:

\[
\text{area} = \left( \frac{\theta}{2\pi} \right) \pi s^2 = \left( \frac{R}{s} \right) \pi s^2 = \pi Rs.
\]

- 12 A band with slant height \( \Delta s = s - s' \) and radii \( R \) and \( R' \) is laid out flat. Explain in one line why its surface area is \( \pi Rs - \pi R's' \).

- 13 By similar triangles \( R/s = R'/s' \) or \( Rs' = R's \). The middle radius \( r \) is \( \frac{1}{2}(R + R') \). Substitute for \( r \) and \( \Delta s \) in the proposed area formula \( 2\pi \Delta s \), to show that this gives the correct area \( \pi Rs - \pi R's' \).

- 14 Slices of a basketball all have the same area of cover, if they have the same thickness.

  (a) Rotate \( y = \sqrt{1 - x^2} \) around the x axis. Show that \( ds = 2\pi dx \).

  (b) The area between \( x = a \) and \( x = a + h \) is _______.

  (c) \( \frac{1}{4} \) of the Earth's area is above latitude _______.

- 15 Change the circle in Example 4 to \( x = a \cos t, \ y = b + a \sin t \). Its radius is _______ and its center is _______. Find the surface area of a torus by revolving this circle around the x axis.

- 16 What part of the circle \( x = R \cos t, \ y = R \sin t \) should rotate around the y axis to produce the top half of a sphere? Choose limits on \( t \) and verify the area.

- 17 The base of a lamp is constructed by revolving the quarter-circle \( y = \sqrt{2x - x^2} (x = 1 \rightarrow x = 2) \) around the y axis. Draw the quarter-circle, find the area integral, and compute the area.

- 18 The light bulb is a sphere of radius \( 1/2 \) with its bottom sliced off to fit onto a cylinder of radius \( 1/4 \) and length \( 1/3 \). Draw the light bulb and find its surface area (ends of the cylinder not included).

- 19 The lampshade is constructed by rotating \( y = 1/x \) around the y axis, and keeping the part from \( y = 1 \) to \( y = 2 \). Set up the definite integral that gives its surface area.

- 20 Compute the area of that lampshade.
21. Explain why the surface area is infinite when \( y = \frac{1}{x} \) is rotated around the \( x \)-axis \((1 < x < \infty)\). But the volume of “Gabriel’s horn” is _________. It can’t hold enough paint to paint its surface.

22. A disk of radius 1” can be covered by four strips of tape (width \( \frac{1}{2} \)). If the strips are not parallel, prove that they can’t cover the disk. Hint: Change to a unit sphere sliced by planes \( \frac{1}{2} \) apart. Problem 14 gives surface area \( \pi \) for each slice.

23. A watermelon (maybe a football) is the result of rotating half of the ellipse \( x = \sqrt{2} \cos t, \ y = \sin t \) (which means \( x^2 + 2y^2 = 2 \)). Find the surface area, parametrically or not.

24. Estimate the surface area of an egg.

8.4 Probability and Calculus

Discrete probability usually involves careful counting. Not many samples are taken and not many experiments are made. There is a list of possible outcomes, and a known probability for each outcome. But probabilities go far beyond red cards and black cards. The real questions are much more practical:

1. How often will too many passengers arrive for a flight?
2. How many random errors do you make on a quiz?
3. What is the chance of exactly one winner in a big lottery?

Those are important questions and we will set up models to answer them.

There is another point. Discrete models do not involve calculus. The number of errors or bumped passengers or lottery winners is a small whole number. Calculus enters for continuous probability. Instead of results that exactly equal 1 or 2 or 3, calculus deals with results that fall in a range of numbers. Continuous probability comes up in at least two ways:

(A) An experiment is repeated many times and we take averages.
(B) The outcome lies anywhere in an interval of numbers.

In the continuous case, the probability \( p_x \) of hitting a particular value \( x = n \) becomes zero. Instead we have a probability density \( p(x) \)—which is a key idea. The chance that a random \( X \) falls between \( a \) and \( b \) is found by integrating the density \( p(x) \):

\[
\text{Prob} : a \leq X \leq b = \int_a^b p(x) \, dx.
\]

Roughly speaking, \( p(x) \, dx \) is the chance of falling between \( x \) and \( x + dx \). Certainly \( p(x) \geq 0 \). If \( a \) and \( b \) are the extreme limits \(-\infty\) and \( \infty \), including all possible outcomes, the probability is necessarily one:

\[
\text{Prob} : -\infty < X < +\infty = \int_{-\infty}^{\infty} p(x) \, dx = 1.
\]

This is a case where infinite limits of integration are natural and unavoidable. In studying probability they create no difficulty—areas out to infinity are often easier.

Here are typical questions involving continuous probability and calculus:

4. How conclusive is a 53%–47% poll of 2500 voters?
5. Are 16 random football players safe on an elevator with capacity 3600 pounds?
6. How long before your car is in an accident?

It is not so traditional for a calculus course to study these questions. They need extra thought, beyond computing integrals (so this section is harder than average). But probability is more important than some traditional topics, and also more interesting.
Drug testing and gene identification and market research are major applications. Comparing Questions 1–3 with 4–6 brings out the relation of discrete to continuous—the differences between them, and the parallels.

It would be impossible to give here a full treatment of probability theory. I believe you will see the point (and the use of calculus) from our examples. Frank Morgan’s lectures have been a valuable guide.

DISCRETE RANDOM VARIABLES

A discrete random variable \( X \) has a list of possible values. For two dice the outcomes are \( X = 2, 3, \ldots, 12 \). For coin tosses (see below), the list is infinite: \( X = 1, 2, 3, \ldots \).

A continuous variable lies in an interval \( a \leq X \leq b \).

EXAMPLE 1 Toss a fair coin until heads come up. The outcome \( X \) is the number of tosses. The value of \( X \) is 1 or 2 or 3 or \( \ldots \), and the probability is \( \frac{1}{2} \) that \( X = 1 \) (heads on the first toss). The probability of tails then heads is \( p_2 = \frac{1}{4} \). The probability that \( X = n \) is \( p_n = (\frac{1}{2})^n \)—this is the chance of \( n - 1 \) tails followed by heads. The sum of all probabilities is necessarily 1:

\[
p_1 + p_2 + p_3 + \ldots = \frac{1}{2} + \frac{1}{4} + \frac{1}{8} + \ldots = 1.
\]

EXAMPLE 2 Suppose a student (not you) makes an average of 2 unforced errors per hour exam. The number of actual errors on the next exam is \( X = 0 \) or 1 or 2 or \( \ldots \). A reasonable model for the probability of \( n \) errors—when they are random and independent—is the Poisson model (pronounced Pwason):

\[
p_n = \text{probability of } n \text{ errors } = \frac{n^n}{n!} e^{-2}.
\]

The probabilities of no errors, one error, and two errors are \( p_0, p_1, \) and \( p_2 \):

\[
p_0 = \frac{2^0}{0!} e^{-2} = \frac{1}{1} e^{-2} \approx .135 \quad p_1 = \frac{2^1}{1!} e^{-2} \approx .27 \quad p_2 = \frac{2^2}{2!} e^{-2} \approx .27.
\]

The probability of more than two errors is \( 1 - .135 - .27 - .27 = .325 \).

This Poisson model can be derived theoretically or tested experimentally. The total probability is again 1, from the infinite series (Section 6.6) for \( e^2 \):

\[
p_0 + p_1 + p_2 + \ldots = \left( \frac{2^0}{0!} + \frac{2^1}{1!} + \frac{2^2}{2!} + \ldots \right) e^{-2} = e^2 e^{-2} = 1.
\]

EXAMPLE 3 Suppose on average 3 out of 100 passengers with reservations don’t show up for a flight. If the plane holds 98 passengers, what is the probability that someone will be bumped?

If the passengers come independently to the airport, use the Poisson model with 2 changed to 3. \( X \) is the number of no-shows, and \( X = n \) happens with probability \( p_n \):

\[
p_n = \frac{3^n}{n!} e^{-3} \quad p_0 = \frac{3^0}{0!} e^{-3} = e^{-3} \quad p_1 = \frac{3^1}{1!} e^{-3} = 3e^{-3}.
\]

There are 98 seats and 100 reservations. Someone is bumped if \( X = 0 \) or \( X = 1 \):

\[
\text{chance of bumping } = p_0 + p_1 = e^{-3} + 3e^{-3} \approx 4/20.
\]

We will soon define the average or expected value or mean of \( X \)—this model has \( \mu = 3 \).
CONTINUOUS RANDOM VARIABLES

If $X$ is the lifetime of a VCR, all numbers $X \geq 0$ are possible. If $X$ is a score on the SAT, then $200 \leq X \leq 800$. If $X$ is the fraction of computer owners in a poll of 600 people, $X$ is between 0 and 1. You may object that the SAT score is a whole number and the fraction of computer owners must be 0 or 1/600 or 2/600 or ..., but it is completely impractical to work with 601 discrete possibilities. Instead we take $X$ to be a continuous random variable, falling anywhere in the range $X \geq 0$ or $[200, 800]$ or $0 < X < 1$. Of course the various values of $X$ are not equally probable.

EXAMPLE 4 The average lifetime of a VCR is 4 years. A reasonable model for breakdown time is an exponential random variable. Its probability density is

$$p(x) = \frac{1}{4} e^{-x/4} \quad \text{for} \quad 0 \leq x < \infty.$$  

The probability that the VCR will eventually break is 1:

$$\int_{0}^{\infty} \frac{1}{4} e^{-x/4} \, dx = \left[ -e^{-x/4} \right]_{0}^{\infty} = 0 - (-1) = 1.$$  

(4)

The probability of breakdown within 12 years ($X$ from 0 to 12) is .95:

$$\int_{0}^{12} \frac{1}{4} e^{-x/4} \, dx = \left[ -e^{-x/4} \right]_{0}^{12} = -e^{-3} + 1 \approx .95.$$  

(5)

An exponential distribution has $p(x) = ae^{-ax}$. Its integral from 0 to $x$ is $F(x) = 1 - e^{-ax}$. Figure 8.11 is the graph for $a = 1$. It shows the area up to $x = 1$.

To repeat: The probability that $a < X < b$ is the integral of $p(x)$ from $a$ to $b$.

Fig. 8.11 Probabilities add to $\Sigma p_n = 1$. Continuous density integrates to $\int p(x) \, dx = 1$.

EXAMPLE 5 We now define the most important density function. Suppose the average SAT score is 500, and the standard deviation (defined below—it measures the spread around the average) is 200. Then the normal distribution of grades has

$$p(x) = \frac{1}{200\sqrt{2\pi}} e^{-\left(x - 500\right)^2/2(200)^2} \quad \text{for} \quad -\infty < x < \infty.$$  

This is the normal (or Gaussian) distribution with mean 500 and standard deviation 200. The graph of $p(x)$ is the famous bell-shaped curve in Figure 8.12.

A new objection is possible. The actual scores are between 200 and 800, while the density $p(x)$ extends all the way from $-\infty$ to $\infty$. I think the Educational Testing Service counts all scores over 800 as 800. The fraction of such scores is pretty small—in fact the normal distribution gives

$$\text{Prob}(X \geq 800) = \int_{800}^{\infty} \frac{1}{200\sqrt{2\pi}} e^{-\left(x - 500\right)^2/2(200)^2} \, dx \approx .0013.$$  

(6)
Regrettably, $e^{-x^2}$ has no elementary antiderivative. We need numerical integration. But there is nothing the matter with that! The integral is called the "error function," and special tables give its value to great accuracy. The integral of $e^{-x^2/2}$ from $-\infty$ to $\infty$ is exactly $\sqrt{2\pi}$. Then division by $\sqrt{2\pi}$ keeps \( \int p(x) \, dx = 1 \).

Notice that the normal distribution involves two parameters. They are the mean value (in this case $\mu = 500$) and the standard deviation (in this case $\sigma = 200$). Those numbers $\mu$ and $\sigma$ are often given the "normalized" values $\mu = 0$ and $\sigma = 1$:

$$ p(x) = \frac{1}{\sigma \sqrt{2\pi}} e^{-((x-\mu)/\sigma)^2} \quad \text{becomes} \quad p(x) = \frac{1}{\sqrt{2\pi}} e^{-x^2/2}. $$

The bell-shaped graph of $p$ is symmetric around the middle point $x = \mu$. The width of the graph is governed by the second parameter $\sigma$—which stretches the $x$ axis and shrinks the $y$ axis (leaving total area equal to 1). The axes are labeled to show the standard case $\mu = 0, \sigma = 1$ and also the graph for any other $\mu$ and $\sigma$.

We now give a name to the integral of $p(x)$. The limits will be $-\infty$ and $x$, so the integral $F(x)$ measures the probability that a random sample is below $x$:

$$ \text{Prob} \{ X \leq x \} = \int_{-\infty}^{x} p(x) \, dx = \text{cumulative density function} \ F(x). \quad (7) $$

$F(x)$ accumulates the probabilities given by $p(x)$, so $dF/dx = p(x)$. The total probability is $F(\infty) = 1$. This integral from $-\infty$ to $\infty$ covers all outcomes.

Figure 8.12b shows the integral of the bell-shaped normal distribution. The middle point $x = \mu$ has $F = \frac{1}{2}$. By symmetry there is a 50-50 chance of an outcome below the mean. The cumulative density $F(x)$ is near .16 at $\mu - \sigma$ and near .84 at $\mu + \sigma$. The chance of falling in between is .84 - .16 = .68. Thus 68% of the outcomes are less than one deviation $\sigma$ away from the center $\mu$.

Moving out to $\mu - 2\sigma$ and $\mu + 2\sigma$, 95% of the area is in between. With 95% confidence $X$ is less than two deviations from the mean. Only one sample in 20 is further out (less than one in 40 on each side).

Note that $\sigma = 200$ is not the precise value for the SAT!

**MEAN, VARIANCE, AND STANDARD DEVIATION**

In Example 1, $X$ was the number of coin tosses until the appearance of heads. The probabilities were $p_1 = \frac{1}{2}, p_2 = \frac{1}{4}, p_3 = \frac{1}{8}, \ldots$. What is the average number of tosses? We now find the "mean" $\mu$ of any distribution $p(x)$—not only the normal distribution, where symmetry guarantees that the built-in number $\mu$ is the mean.

*To find $\mu$, multiply outcomes by probabilities and add:*

$$ \mu = \text{mean} = \sum n p_n = 1(p_1) + 2(p_2) + 3(p_3) + \cdots. \quad (8) $$
The average number of tosses is \(1(\frac{1}{2}) + 2(\frac{1}{4}) + 3(\frac{1}{8}) + \cdots\). This series adds up (in Section 10.1) to \(\mu = 2\). Please do the experiment 10 times. I am almost certain that the average will be near 2.

When the average is \(\lambda = 2\) quiz errors or \(\lambda = 3\) no-shows, the Poisson probabilities are \(p_n = \lambda^n e^{-\lambda}/n!\). Check that the formula \(\mu = \sum np_n\) does give \(\lambda\) as the mean:

\[
\left[ 1 \cdot \frac{\lambda}{1!} + 2 \cdot \frac{\lambda^2}{2!} + 3 \cdot \frac{\lambda^3}{3!} + \cdots \right] e^{-\lambda} = \lambda \left[ 1 + \frac{\lambda}{1!} + \frac{\lambda^2}{2!} + \cdots \right] e^{-\lambda} = \lambda e^\lambda e^{-\lambda} = \lambda.
\]

For continuous probability, the sum \(\mu = \sum np_n\) changes to \(\mu = \int xp(x) \, dx\). We multiply outcome \(x\) by probability \(p(x)\) and integrate. In the VCR model, integration by parts gives a mean breakdown time of \(\mu = 4\) years:

\[
\int x \, p(x) \, dx = \int_0^\infty x(\frac{1}{4}e^{-x/4}) \, dx = \left[ -xe^{-x/4} - 4e^{-x/4} \right]_0^\infty = 4.
\]

Together with the mean we introduce the variance. It is always written \(\sigma^2\), and in the normal distribution that measured the "width" of the curve. When \(\sigma^2\) was 200\(^2\), SAT scores spread out pretty far. If the testing service changed to \(\sigma^2 = 1\), the scores would be a disaster. 95% of them would be within \(\pm 2\) of the mean. When a teacher announces an average grade of 72, the variance should also be announced—if it is big then those with 60 can relax. At least they have company.

### 8E

The mean \(\mu\) is the expected value of \(X\). The variance \(\sigma^2\) is the expected value of \((X - \text{mean})^2 = (X - \mu)^2\). Multiply outcome times probability and add:

\[
\begin{align*}
\mu &= \sum np_n \\
\sigma^2 &= \sum (n - \mu)^2 p_n \\
\mu &= \int_0^\infty x \, p(x) \, dx \\
\sigma^2 &= \int_0^\infty (x - \mu)^2 \, p(x) \, dx
\end{align*}
\]

The standard deviation (written \(\sigma\)) is the square root of \(\sigma^2\).

**EXAMPLE 6** (Yes-no poll, one person asked) The probabilities are \(p\) and \(1 - p\).

A fraction \(p = \frac{1}{3}\) of the population thinks yes, the remaining fraction \(1 - p = \frac{2}{3}\) thinks no. Suppose we only ask one person. If \(X = 1\) for yes and \(X = 0\) for no, the expected value of \(X\) is \(\mu = p = \frac{1}{3}\). The variance is \(\sigma^2 = p(1 - p) = \frac{2}{9}\):

\[
\mu = 0 \left(\frac{2}{3}\right) + 1 \left(\frac{1}{3}\right) = \frac{1}{3} \quad \text{and} \quad \sigma^2 = \left(0 - \frac{1}{3}\right)^2 \left(\frac{2}{3}\right) + \left(1 - \frac{1}{3}\right)^2 \left(\frac{1}{3}\right) = \frac{2}{9}.
\]

The standard deviation is \(\sigma = \sqrt{2/9}\). When the fraction \(p\) is near one or near zero, the spread is smaller—and one person is more likely to give the right answer for everybody. The maximum of \(\sigma^2 = p(1 - p)\) is at \(p = \frac{1}{2}\), where \(\sigma = \frac{1}{2}\).

The table shows \(\mu\) and \(\sigma^2\) for important probability distributions.

<table>
<thead>
<tr>
<th>Model</th>
<th>Mean</th>
<th>Variance</th>
<th>Application</th>
</tr>
</thead>
<tbody>
<tr>
<td>(p_1 = p, p_0 = 1 - p)</td>
<td>(p)</td>
<td>(p(1 - p))</td>
<td>yes-no</td>
</tr>
<tr>
<td>Poisson (p_n = \lambda^n e^{-\lambda}/n!)</td>
<td>(\lambda)</td>
<td>(\lambda)</td>
<td>random occurrence</td>
</tr>
<tr>
<td>Exponential (p(x) = ae^{-ax})</td>
<td>(1/a)</td>
<td>(1/a^2)</td>
<td>waiting time</td>
</tr>
<tr>
<td>Normal (p(x) = \frac{1}{\sqrt{2\pi} \sigma} e^{-(x-\mu)^2/2\sigma^2})</td>
<td>(\mu)</td>
<td>(\sigma^2)</td>
<td>distribution around mean</td>
</tr>
</tbody>
</table>
THE LAW OF AVERAGES AND THE CENTRAL LIMIT THEOREM

We come to the center of probability theory (without intending to give proofs). The key idea is to repeat an experiment many times—poll many voters, or toss many dice, or play considerable poker. Each independent experiment produces an outcome $X$, and the average from $N$ experiments is $\bar{X}$. It is called "$X$ bar":

$$\bar{X} = \frac{X_1 + X_2 + \cdots + X_N}{N} = \text{average outcome.}$$

All we know about $p(x)$ is its mean $\mu$ and variance $\sigma^2$. It is amazing how much information that gives about the average $\bar{X}$:

| Law of Averages: $\bar{X}$ is almost sure to approach $\mu$ as $N \to \infty$. |
| Central Limit Theorem: The probability density $p_\mu(x)$ for $\bar{X}$ approaches a normal distribution with the same mean $\mu$ and with variance $\sigma^2/N$. |

No matter what the probabilities for $X$, the probabilities for $\bar{X}$ move toward the normal bell-shaped curve. The standard deviation is close to $\sigma/\sqrt{N}$ when the experiment is repeated $N$ times. In the Law of Averages, "almost sure" means that the chance of $\bar{X}$ not approaching $\mu$ is zero. It can happen, but it won't.

**Remark 1** The Boston Globe doesn't understand the Law of Averages. I quote from September 1988: "What would happen if a giant Red Sox slump arrived? What would happen if the fabled Law of Averages came into play, reversing all those can't miss decisions during the winning streak?" They think the Law of Averages evens everything up, favoring heads after a series of tails. See Problem 20.

**EXAMPLE 7** Yes-no poll of $N = 2500$ voters. Is a 53%-47% outcome conclusive?

The fraction $p$ of "yes" voters in the whole population is not known. That is the reason for the poll. The deviation $\sigma = \sqrt{p(1-p)}$ is also not known, but for one voter this is never more than $\frac{1}{2}$ (when $p = \frac{1}{2}$). Therefore $\sigma/\sqrt{N}$ for 2500 voters is no larger than $\frac{1}{2}/\sqrt{2500}$, which is 1%.

The result of the poll was $\bar{X} = 53\%$. With 95% confidence, this sample is within two standard deviations (here 2%) of its mean. Therefore with 95% confidence, the unknown mean $\mu = p$ of the whole population is between 51% and 55%. This poll is conclusive.

If the true mean had been $p = 50\%$, the poll would have had only a .0013 chance of reaching 53%. The error margin on each side of a poll is amazingly simple; it is always $1/\sqrt{N}$.

**Remark 2** The New York Times has better mathematicians than the Globe. Two days after Bush defeated Dukakis, their poll of $N = 11,645$ voters was printed with the following explanation. "In theory, in 19 cases out of 20 [there is 95%] the results should differ by no more than one percentage point [there is $1/\sqrt{N}$] from what would have been obtained by seeking out all voters in the United States."

**EXAMPLE 8** Football players at Caltech (if any) have average weight $\mu = 210$ pounds and standard deviation $\sigma = 30$ pounds. Are $N = 16$ players safe on an elevator with capacity 3600 pounds? 16 times 210 is 3360.
8 Applications of the Integral

The average weight $\bar{X}$ is approximately a normal random variable with $\mu = 210$ and $\sigma = 30/\sqrt{N} = 30/4$. There is only a 2% chance that $\bar{X}$ is above $\mu + 2\sigma = 225$ (see Figure 8.12b—weights below the mean are no problem on an elevator). Since 16 times 225 is 3600, a statistician would have 98% confidence that the elevator is safe. This is an example where 98% is not good enough—I wouldn't get on.

EXAMPLE 9 (The famous Weldon Dice) Weldon threw 12 dice 26,306 times and counted the 5's and 6's. They came up in 33.77% of the 315,672 separate rolls. Thus $\bar{X} = 0.3377$ instead of the expected fraction $p = \frac{1}{6}$ of 5's and 6's. Were the dice fair?

The variance in each roll is $\sigma^2 = p(1-p) = 2/9$. The standard deviation of $\bar{X}$ is $\sigma = \sigma/\sqrt{N} = \sqrt{2/9}/\sqrt{315672} \approx 0.00084$. For fair dice, there is a 95% chance that $\bar{X}$ will differ from $\frac{1}{6}$ by less than $2\sigma$. (For Poisson probabilities that is false. Here $\bar{X}$ is normal.) But $0.3377$ differs from $0.3333$ by more than $5\sigma$. The chance of falling 5 standard deviations away from the mean is only about 1 in 10,000.†

So the dice were unfair. The faces with 5 or 6 indentations were lighter than the others, and a little more likely to come up. Modern dice are made to compensate for that, but Weldon never tried again.

Read-through questions

Discrete probability uses counting, __a__ probability uses calculus. The function $p(x)$ is the probability __b___. The chance that a random variable falls between $a$ and $b$ is __c__.

The total probability is $\int_a^b p(x) \, dx = 1$. In the discrete case $\sum p_i = 1$. The mean (or expected value) is $\mu = \int x \, p(x) \, dx$ in the continuous case and $\mu = \sum x p_i$ in the discrete case.

The Poisson distribution with mean $\lambda$ has $p_n = e^{-\lambda} \lambda^n / n!$. The sum $\sum p_n = 1$ comes from the __d__ series. The exponential distribution has $p(x) = e^{-x}$ or $2e^{-2x}$ or __e___. The standard Gaussian (or __f__ ) distribution has $\sqrt{2\pi} p(x) = e^{-x^2/2}$. Its graph is the well-known __g__ curve. The chance that the variable falls below $x$ is $F(x) = \int_0^x p(t) \, dt$. $F$ is the __h__ function. The difference $F(x + dx) - F(x)$ is about __i__ , which is the chance that $X$ is between $x$ and $x + dx$.

The variance, which measures the spread around $\mu$, is $\sigma^2 = \int (x - \mu)^2 \, p(x) \, dx$ in the continuous case and $\sigma^2 = \sum (x - \mu)^2$ in the discrete case. Its square root $\sigma$ is the __j___. The normal distribution has $p(x) = e^{-x^2}$. If $\bar{X}$ is the average of $N$ samples from any population with mean $\mu$ and variance $\sigma^2$, the Law of Averages says that $\bar{X}$ will approach __k___. The Central Limit Theorem says that the distribution for $\bar{X}$ approaches __l___. Its mean is __m__ and its variance is __n__.

8.4 EXERCISES

In a yes-no poll when the voters are 50-50, the mean for one voter is $p = 0\frac{1}{2} + 1\frac{1}{2} = \frac{1}{2}$. The variance is $(0 - p)^2 p_0 + (1 - p)^2 p_1 = \frac{1}{4}$. For a poll with $N = 100$, $\sigma$ is __o___. There is a 95% chance that $\bar{X}$ (the fraction saying yes) will be between __p__ and __q__.

1 If $p_1 = \frac{1}{4}$, $p_2 = \frac{1}{4}$, $p_3 = \frac{3}{4}$, ..., what is the probability of an outcome $X < 4$? What are the probabilities of $X = 4$ and $X > 4$?

2 With the same $p_n = \left(\frac{1}{4}\right)^n$, what is the probability that $X$ is odd? Why is $p_n = \left(\frac{1}{4}\right)^n$ an impossible set of probabilities? What multiple $c\left(\frac{1}{4}\right)^n$ is possible?

3 Why is $p(x) = e^{-2x}$ not an acceptable probability density for $x \geq 0$? Why is $p(x) = 4e^{-2x} - e^{-x}$ not acceptable?

*4 If $p_n = \left(\frac{1}{4}\right)^n$, show that the probability $P$ that $X$ is a prime number satisfies $6/16 \leq P \leq 7/16$.

5 If $p(x) = e^{-x}$ for $x \geq 0$, find the probability that $X \geq 2$ and the approximate probability that $1 \leq X \leq 1.01$.

6 If $p(x) = C/x^3$ is a probability density for $x \geq 1$, find the constant $C$ and the probability that $X \leq 2$.

7 If you choose $x$ completely at random between 0 and $\pi$, what is the density $p(x)$ and the cumulative density $F(x)$?

†Joe DiMaggio's 56-game hitting streak was much more improbable—I think it is statistically the most exceptional record in major sports.
8.4 Probability and Calculus

In 8–13 find the mean value \( \mu = \Sigma np_n \) or \( \mu = \int xp(x) \, dx \).

8. \( p_0 = 1/2, p_1 = 1/4, p_2 = 1/4 \)

9. \( p_1 = 1/7, p_2 = 1/7, \ldots, p_7 = 1/7 \)

10. \( p_x = 1/n^2 \) (\( p_0 = 1/e, p_1 = 1/e, p_2 = 1/2e, \ldots \))

11. \( p(x) = 2/\pi(1 + x^2), \quad x \geq 0 \)

12. \( p(x) = e^{-x} \) (integrate by parts)

13. \( p(x) = ae^{-ax} \) (integrate by parts)

14. Show by substitution that

\[
\int_{-\infty}^{\infty} e^{-x^2/2a^2} \, dx = \sqrt{2\pi} \int_{0}^{\infty} e^{-u^2} \, du = \sqrt{\frac{\pi}{2a}}.
\]

15. Find the cumulative probability \( F \) (the integral of \( p \)) in Problems 11, 12, 13. In terms of \( F \), what is the chance that a random sample lies between \( a \) and \( b \)?

16. Can-Do Airlines books 100 passengers when their plane only holds 98. If the average number of no-shows is 2, what is the Poisson probability that someone will be bumped?

17. The waiting time for a bus has probability density \( (1/10)e^{-x/10} \), with \( \mu = 10 \) minutes. What is the probability of waiting longer than 10 minutes?

18. You make a 3-minute telephone call. If the waiting time for the next incoming call has \( p(x) = e^{-x} \), what is the probability that your phone will be busy?

19. Supernovas are expected about every 100 years. What is the probability that you will be alive for the next one? Use a Poisson model with \( \lambda = 0.01 \) and estimate your lifetime. (Supernovas actually occurred in 1054 (Crab Nebula), 1572, 1604, and 1987. But the future distribution doesn't depend on the date of the last one.)

20. (a) A fair coin comes up heads 10 times in a row. Will heads or tails be more likely on the next toss?

(b) The fraction of heads after \( N \) tosses is \( x \). The expected fraction after \( 2N \) tosses is \( \frac{x}{2} \).

21. Show that the area between \( \mu \) and \( \mu + \sigma \) under the bell-shaped curve is a fixed number (near 1/3), by substituting \( y = \frac{x - \mu}{\sigma \sqrt{2\pi}} \):

\[
\int_{\mu}^{\mu+\sigma} \frac{1}{\sigma \sqrt{2\pi}} e^{-\frac{(x-\mu)^2}{2\sigma^2}} \, dx = \int_{0}^{1} \frac{1}{\sqrt{2\pi}} e^{-y^2/2} \, dy.
\]

What is the area between \( \mu - \sigma \) and \( \mu \)? The area outside \( (\mu - \sigma, \mu + \sigma) \)?

22. For a yes-no poll of two voters, explain why

\( p_0 = (1-p)^2, p_1 = 2p - 2p^2, p_2 = p^2. \)

Find \( \mu \) and \( \sigma^2 \). \( N \) voters give the "binomial distribution."

23. Explain the last step in this reorganization of the formula for \( \sigma^2 \):

\[
\sigma^2 = \int (x - \mu)^2 p(x) \, dx = \int (x^2 - 2\mu x + \mu^2) p(x) \, dx
= \int x^2 p(x) \, dx - 2\mu \int xp(x) \, dx + \mu^2 \int p(x) \, dx
= \int x^2 p(x) \, dx - \mu^2.
\]

24. Use \( \int (x - \mu)^2 p(x) \, dx \) and also \( \int x^2 p(x) \, dx - \mu^2 \) to find \( \sigma^2 \) for the uniform distribution: \( p(x) = 1 \) for \( 0 \leq x \leq 1 \).

25. Find \( \sigma^2 \) if \( p_0 = 1/3, p_1 = 1/3, p_2 = 1/3 \). Use \( \Sigma (n - \mu)^2 p_n \) and also \( \Sigma n^2 p_n - \mu^2 \).

26. Use Problem 23 and integration by parts (equation 7.1.10) to find \( \sigma^2 \) for the exponential distribution \( p(x) = 2e^{-2x} \) for \( x \geq 0 \), which has mean \( \frac{1}{2} \).

27. The waiting time to your next car accident has probability density \( p(x) = \frac{1}{2}e^{-x^2} \). What is \( \mu \)? What is the probability of no accident in the next four years?

28. With \( p = \frac{1}{2}, \frac{1}{4}, \frac{1}{8}, \ldots \), find the average number \( \mu \) of coin tosses by writing \( p_1 + 2p_2 + 3p_3 + \ldots \) as \( (p_1 + p_2 + p_3 + \cdots) + (p_2 + p_3 + p_4 + \cdots) + (p_3 + p_4 + p_5 + \cdots) + \cdots \).

29. In a poll of 900 Americans, 30 are in favor of war. What range can you give with 95% confidence for the percentage of peaceful Americans?

30. Sketch rough graphs of \( p(x) \) for the fraction \( x \) of heads in 4 tosses of a fair coin, and in 16 tosses. The mean value is \( \frac{1}{2} \).

31. A judge tosses a coin 2500 times. How many heads does it take to prove with 95% confidence that the coin is unfair?

32. Long-life bulbs shine an average of 2000 hours with standard deviation 150 hours. You can have 95% confidence that your bulb will fail between ______ and ______ hours.

33. Grades have a normal distribution with mean 70 and standard deviation 10. If 300 students take the test and passing is 55, how many are expected to fail? (Estimate from Figure 8.12b.) What passing grade will fail 1/10 of the class?

34. The average weight of luggage is \( \mu = 30 \) pounds with deviation \( \sigma = 8 \) pounds. What is the probability that the luggage for 64 passengers exceeds 2000 pounds? How does the answer change for 256 passengers and 8000 pounds?

35. A thousand people try independently to guess a number between 1 and 1000. This is like a lottery.

(a) What is the chance that the first person fails?

(b) What is the chance \( P_0 \) that they all fail?

(c) Explain why \( P_0 \) is approximately \( \frac{1}{e} \).

36. (a) In Problem 35, what is the chance that the first person is right and all others are wrong?

(b) Show that the probability \( P_1 \) of exactly one winner is also close to \( \frac{1}{e} \).

(c) Guess the probability \( P_n \) of \( n \) winners (fishy question).
This chapter concludes with two sections related to engineering and physics. Each
application starts with a finite number of masses or forces. Their sum is the total
mass or total force. Then comes the “continuous case,” in which the mass is spread
out instead of lumped. Its distribution is given by a density function $\rho$ (Greek rho),
and the sum changes to an integral.

The first step (hardest step?) is to get the physical quantities straight. The second
step is to move from sums to integrals (discrete to continuous, lumped to distributed).
By now we hardly stop to think about it—although this is the key idea of integral
calculus. The third step is to evaluate the integrals. For that we can use substitution
or integration by parts or tables or a computer.

Figure 8.13 shows the one-dimensional case: masses along the $x$ axis. The total
mass is the sum of the masses. The new idea is that of moments—when the mass or
force is multiplied by a distance:

$$\text{moment of mass around the } y \text{ axis} = mx = (\text{mass}) \times (\text{distance to axis}).$$

The figure has masses 1, 3, 2. The total mass is 6. The “lever arms” or “moment
arms” are the distances $x = 1, 3, 7$. The masses have moments 1 and 9 and 14 (since
$mx$ is 2 times 7). The total moment is $1 + 9 + 14 = 24$. Then the balance point is at
$x = M_x/M = 24/6 = 4$.

The total mass is the sum of the $m$'s. The total moment is the sum of $m_n x_n$
(negative on the other side of $x = 0$). If the masses are children on a seesaw, the
balance point is the center of gravity $\bar{x}$—also called the center of mass:

$$\bar{x} = \frac{\sum m_n x_n}{\sum m_n} = \frac{\text{total moment}}{\text{total mass}}. \quad (1)$$

If all masses are moved to $\bar{x}$, the total moment (6 times 4) is still 24. The moment
equals the mass $\sum m_n$ times $\bar{x}$. The masses act like a single mass at $\bar{x}$.

Also: If we move the axis to $\bar{x}$, and leave the children where they are, the seesaw
balances. The masses on the left of $\bar{x} = 4$ will offset the mass on the right. Reason:
The distances to the new axis are $x_n - \bar{x}$. The moments add to zero by equation (1):

$$\text{moment around new axis} = \sum m_n (x_n - \bar{x}) = \sum m_n x_n - \sum m_n \bar{x} = 0.$$  

Turn now to the continuous case, when mass is spread out along the line. Each
piece of length $\Delta x$ has an average density $\rho_n = (\text{mass of piece})/(\text{length of piece}) = \Delta m/\Delta x$. As the pieces get shorter, this approaches $dm/dx$—the density at the point. The limit of $(\text{small mass})/(\text{small length})$ is the density $\rho(x)$.

Integrating that derivative $\rho = dm/dx$, we recover the total mass: $\Sigma \rho_n \Delta x$ becomes

$$\text{total mass } M = \int \rho(x) \, dx. \quad (2)$$

When the mass is spread evenly, $\rho$ is constant. Then $M = \rho L = \text{density times length}$. 

---

**8.5 Masses and Moments**
8.5 Masses and Moments

The moment formula is similar. For each piece, the moment is mass $p_a \Delta x$ multiplied by distance $x$—and we add. In the continuous limit, $\rho(x) \, dx$ is multiplied by $x$ and we integrate:

$$\text{total moment around y axis} = M_y = \int \rho(x) \, dx.$$  \hfill (3)

Moment is mass times distance. Dividing by the total mass $M$ gives “average distance”:

$$\text{center of mass } \bar{x} = \frac{\text{moment}}{\text{mass}} = \frac{M_y}{M} = \int \frac{x \rho(x) \, dx}{\int \rho(x) \, dx}.$$  \hfill (4)

Remark If you studied Section 8.4 on probability, you will notice how the formulas match up. The mass $\int \rho(x) \, dx$ is like the total probability $\int \rho(x) \, dx$. The moment $\int x \rho(x) \, dx$ is like the mean $\int x \rho(x) \, dx$. The moment of inertia $\int (x - \bar{x})^2 \rho(x) \, dx$ is the variance. Mathematics keeps hammering away at the same basic ideas! The only difference is that the total probability is always 1. The mean really corresponds to the center of mass $\bar{x}$, but in probability we didn’t notice the division by $\int \rho(x) \, dx = 1$.

EXAMPLE 1 With constant density $\rho$ from 0 to $L$, the mass is $M = \rho L$. The moment is

$$M_y = \int_0^L x \rho \, dx = \frac{1}{2} \rho x^2 \bigg|_0^L = \frac{1}{2} \rho L^2.$$  

The center of mass is $\bar{x} = M_y/M = L/2$. It is halfway along.

EXAMPLE 2 With density $e^{-x}$ the mass is 1, the moment is 1, and $\bar{x}$ is 1:

$$\int_0^\infty e^{-x} \, dx = \left[-e^{-x}\right]_0^\infty = 1 \quad \text{and} \quad \int_0^\infty xe^{-x} \, dx = \left[-xe^{-x} - e^{-x}\right]_0^\infty = 1.$$

MASSES AND MOMENTS IN TWO DIMENSIONS

Instead of placing masses along the x axis, suppose $m_i$ is at the point $(x_i, y_i)$ in the plane. Similarly $m_n$ is at $(x_n, y_n)$. Now there are two moments to consider. Around the y axis $M_y = \Sigma m_n x_n$ and around the x axis $M_x = \Sigma m_n y_n$. Please notice that the x’s go into the moment $M_y$—because the x coordinate gives the distance from the y axis!

Around the x axis, the distance is y and the moment is $M_x$. The center of mass is the point $(\bar{x}, \bar{y})$ at which everything balances:

$$\bar{x} = \frac{M_y}{M} = \frac{\sum m_n x_n}{\sum m_n} \quad \text{and} \quad \bar{y} = \frac{M_x}{M} = \frac{\sum m_n y_n}{\sum m_n}.$$  \hfill (5)

In the continuous case these sums become two-dimensional integrals. The total mass is $\int \int \rho(x, y) \, dx \, dy$, when the density is $\rho = \text{mass per unit area}$. These “double integrals” are for the future (Section 14.1). Here we consider the most important case: $\rho = \text{constant}$. Think of a thin plate, made of material with constant density (say $\rho = 1$). To compute its mass and moments, the plate is cut into strips (Figure 8.14):

mass $M = \text{area of plate}$ \hfill (6)

moment $M_y = \int (\text{distance x}) \, (\text{length of vertical strip}) \, dx$ \hfill (7)

moment $M_x = \int (\text{height y}) \, (\text{length of horizontal strip}) \, dy$. \hfill (8)
Applications of the Integral

Fig. 8.14 Plates cut into strips to compute masses and moments and centroids.

The mass equals the area because \( \rho = 1 \). For moments, all points in a vertical strip are the same distance from the \( y \) axis. That distance is \( x \). The moment is \( x \) times area, or \( x \) times length times \( dx \)—and the integral accounts for all strips.

Similarly the \( x \)-moment of a horizontal strip is \( y \) times strip length times \( dy \).

EXAMPLE 3 A plate has sides \( x = 0 \) and \( y = 0 \) and \( y = 4 - 2x \). Find \( M, M_y, M_x \).

mass \( M = \text{area} = \int_0^2 \int_0^{4-2x} \rho(x, y) \, dy \, dx = \left[ \int_0^2 (4 - x^2) \, dx \right]^2 = 4. \)

The vertical strips go up to \( y = 4 - 2x \), and the horizontal strips go out to \( x = \frac{1}{2}(4 - y) \):

moment \( M_y = \int_0^2 x(4 - 2x) \, dx = \left[ 2x^2 - \frac{2}{3}x^3 \right]_0^2 = \frac{8}{3} \)

moment \( M_x = \int_0^a \frac{1}{2} (4 - y) \, dy = \left[ y^2 - \frac{1}{6}y^3 \right]_0^4 = \frac{16}{3} \).

The "center of mass" has \( \bar{x} = M_y/M = 2/3 \) and \( \bar{y} = M_x/M = 4/3 \). This is the centroid of the triangle (and also the "center of gravity"). With \( \rho = 1 \) these terms all refer to the same balance point \((\bar{x}, \bar{y})\). The plate will not tip over, if it rests on that point.

EXAMPLE 4 Find \( M_y \) and \( M_x \) for the half-circle below \( x^2 + y^2 = r^2 \).

\( M_y = 0 \) because the region is symmetric—Figure 8.14 balances on the \( y \) axis. In the \( x \)-moment we integrate \( y \) times the length of a horizontal strip (notice the factor 2):

\[ M_x = \int_0^r y \cdot 2\sqrt{r^2 - y^2} \, dy = -\frac{2}{3}(r^2 - y^2)^{3/2} \bigg|_0^r = \frac{2}{3}r^3. \]

Divide by the mass (the area \( \frac{1}{2}\pi r^2 \)) to find the height of the centroid: \( \bar{y} = M_x/M = 4r/3\pi \). This is less than \( \frac{1}{2}r \) because the bottom of the semicircle is wider than the top.

MOment of Inertia

The moment of inertia comes from multiplying each mass by the square of its distance from the axis. Around the \( y \) axis, the distance is \( x \). Around the origin, it is \( r \):

\[ I_y = \sum x^2 m_n \quad \text{and} \quad I_x = \sum y^2 m_n \quad \text{and} \quad I_o = \sum r^2 m_n. \]

Notice that \( I_x + I_y = I_o \) because \( x^2 + y^2 = r^2 \). In the continuous case we integrate.

The moment of inertia around the \( y \) axis is \( I_y = \iint x^2 \rho(x, y) \, dx \, dy \). With a constant density \( \rho = 1 \), we again keep together the points on a strip. On a vertical strip they share the same \( x \). On a horizontal strip they share \( y \):

\[ I_y = \int (x^2) (\text{vertical strip length}) \, dx \quad \text{and} \quad I_x = \int (y^2) (\text{horizontal strip length}) \, dy. \]
In engineering and physics, it is rotation that leads to the moment of inertia. Look at the energy of a mass \( m \) going around a circle of radius \( r \). It has \( I_0 = mr^2 \).

\[
\text{kinetic energy} = \frac{1}{2}mv^2 = \frac{1}{2}m(r\omega)^2 = \frac{1}{2}I_0\omega^2.
\] (9)

The angular velocity is \( \omega \) (radians per second). The speed is \( v = r\omega \) (meters per second).

An ice skater reduces \( I_0 \) by putting her arms up instead of out. She stays close to the axis of rotation \( (r \) is small). Since her rotational energy \( \frac{1}{2}I_0\omega^2 \) does not change, \( \omega \) increases as \( I_0 \) decreases. Then she spins faster.

Another example: It takes force to turn a revolving door. More correctly, it takes torque. The force is multiplied by distance from the turning axis: \( T = Fx \), so a push further out is more effective.

To see the physics, replace Newton's law \( F = ma = m \, dv/dt \) by its rotational form: \( T = I \, d\omega/dt \). Where \( F \) makes the mass move, the torque \( T \) makes it turn. Where \( m \) measures unwillingness to change speed, \( I \) measures unwillingness to change rotation.

**EXAMPLE 5** Find the moment of inertia of a rod about (a) its end and (b) its center.

The distance \( x \) from the end of the rod goes from 0 to \( L \). The distance from the center goes from \(-L/2\) to \( L/2 \). Around the center, turning is easier because \( I \) is smaller:

\[
I_{\text{end}} = \int_0^L x^2 \, dx = \frac{1}{3}L^3 \quad I_{\text{center}} = \int_{-L/2}^{L/2} x^2 \, dx = \frac{1}{12}L^3.
\] (10)

**Fig. 8.15** Moment of inertia for rod and propeller. Rolling balls beat cylinders.

**MOMENT OF INERTIA EXPERIMENT**

**Experiment**: Roll a solid cylinder (a coin), a hollow cylinder (a ring), a solid ball (a marble), and a hollow ball (not a pingpong ball) down a slope. Galileo dropped things from the Leaning Tower—this experiment requires a Leaning Table. Objects that fall together from the tower don’t roll together down the table.

**Question 1** What is the order of finish? *Record your prediction first!*

**Question 2** Does size make a difference if shape and density are the same?

**Question 3** Does density make a difference if size and shape are the same?

**Question 4** Find formulas for the velocity \( v \) and the finish time \( T \).

To compute \( v \), the key is that potential energy plus kinetic energy is practically constant. Energy loss from rolling friction is very small. If the mass is \( m \) and the vertical drop is \( h \), the energy at the top (all potential) is \( mgh \). The energy at the bottom (all kinetic) has two parts: \( \frac{1}{2}mv^2 \) from movement along the plane plus \( \frac{1}{2}I\omega^2 \) from turning. *Important fact*: \( v = r\omega \) for a rolling cylinder or ball of radius \( r \).
8 Applications of the Integral

Equate energies and set $\omega = v/r$:

$$mgh = \frac{1}{2} mv^2 + \frac{1}{2} I\omega^2 = \frac{1}{2} mv^2 \left( 1 + \frac{l}{mr^2} \right).$$

(11)

The ratio $I/mr^2$ is critical. Call it $J$ and solve (11) for $v^2$:

$$v^2 = \frac{2gh}{1 + J} \text{ (smaller } J \text{ means larger velocity).}$$

(12)

The order of $J's$, for different shapes and sizes, should decide the race. Apparently the density doesn't matter, because it is a factor in both $I$ and $m$—so it cancels in $J = I/mr^2$. A hollow cylinder has $J = 1$, which is the largest possible— all its mass is at the full distance $r$ from the axis. So the hollow cylinder should theoretically come in last. This experiment was developed by Daniel Drucker.

Problems 35–37 find the other three $J's$. Problem 40 finds the time $T$ by integration. Your experiment will show how close this comes to the measured time.

8.5 Exercises

Read-through questions

If masses $m_a$ are at distances $x_a$, the total mass is $M = \sum a$. The total moment about $x = 0$ is $M = \sum b$. The center of mass is at $x = \sum c$. In the continuous case, the mass distribution is given by the $\sum d \rho(x)$. The total mass is $M = \int e$ and the center of mass is at $x = \int f$. With $p = x$, the integrals from 0 to $L$ give $M = \int g$ and $\int h x \rho(x) \, dx = \frac{i}{j}$ and $\bar{x} = \frac{k}{l}$. The total moment is the same if the whole mass $M$ is placed at $1$.

In a plane, with masses $m_a$ at the points $(x_a, y_a)$, the moment around the $y$ axis is $M_y = n$. The center of mass has $\bar{x} = \frac{o}{p}$ and $\bar{y} = \frac{q}{r}$. For a plate with density $\rho = 1$, the mass $M$ equals the $\int s$. If the plate is divided into vertical strips of height $y(x)$, then $M = \int t y(x) \, dx$ and $M_y = \int u y(x) \, dx$. For a square plate $0 \leq x, y \leq L$, the mass is $M = \int v$, and the moment around the $y$ axis is $M_y = \int w x$. The center of mass is at $(\bar{x}, \bar{y}) = \int x$. This point is the $\sum x \text{,}$ where the plate balances.

A mass $m$ at a distance $x$ from the axis has moment of inertia $I = \int y^2 \, dx$. A rod with $\rho = 1$ from $x = a$ to $x = b$ has $I_y = \int u$. For a plate with $\rho = 1$ and strips of height $y(x)$, this becomes $I_y = \int v$. The torque $T$ is $\sum w \times$ times $x$.

Compute the mass $M$ along the $x$ axis, the moment $M_x$, around $x = 0$, and the center of mass $\bar{x} = M_x/M$.

1. $m_1 = 2$ at $x_1 = 1$, $m_2 = 4$ at $x_2 = 2$
2. $m = 3$ at $x = 0, 1, 2, 6$
3. $\rho = 1$ for $-1 \leq x \leq 3$
4. $\rho = x^2$ for $0 \leq x \leq 1$

5. $\rho = 1$ for $0 \leq x < 1$, $\rho = 2$ for $1 \leq x \leq 2$
6. $\rho = \sin x$ for $0 \leq x \leq \pi$

Find the mass $M$, the moments $M_x$ and $M_y$, and the center of mass $(\bar{x}, \bar{y})$.

7. Unit masses at $(x, y) = (1, 0), (0, 1)$, and $(1, 1)$
8. $m_1 = 1$ at $(1, 0), m_2 = 4$ at $(0, 1)$
9. $\rho = 7$ in the square $0 \leq x \leq 1, 0 \leq y \leq 1$.
10. $\rho = 3$ in the triangle with vertices $(0, 0), (a, 0), \text{ and } (0, b)$.

Find the area $M$ and the centroid $(\bar{x}, \bar{y})$ inside curves 11–16.

11. $y = \sqrt{1 - x^2}, y = 0, x = 0$ (quarter-circle)
12. $y = x, y = 2 - x, y = 0$ (triangle)
13. $y = e^{-2x}, y = 0, x = 0$ (infinite dagger)
14. $y = x^2, y = x$ (lens)
15. $x^2 + y^2 = 1, x^2 + y^2 = 4$ (ring)
16. $x^2 + y^2 = 1, x^2 + y^2 = 4, y = 0$ (half-ring).

Verify these engineering formulas for $I_x$ with $\rho = 1$:

17. Rectangle bounded by $x = 0, x = a, y = 0, y = b$: $I_x = a^3 b/12$.
18. Square bounded by $x = -\frac{1}{2} a, x = \frac{1}{2} a, y = -\frac{1}{2} a, y = \frac{1}{2} a$: $I_x = a^4 / 12$.
19. Triangle bounded by $x = 0, y = 0, x + y = a$: $I_x = a^2 / 12$.
20. Disk of radius $a$ centered at $x = y = 0$: $I_x = \pi a^4 / 4$. 

21 The moment of inertia around the point \( x = t \) of a rod with density \( \rho(x) \) is \( I = \int (x-t)^2 \rho(x) \, dx \). Expand \((x-t)^2\) and \( I \) into three terms. Show that \( dI/dt = 0 \) when \( t = \bar{x} \). The moment of inertia is smallest around the center of mass.

22 A region has \( \bar{x} = 0 \) if \( M_x = \int x \text{(height of strip)} \, dx = 0 \). The moment of inertia about any other axis \( x = c \) is \( I = \int (x-c)^2 \text{(height of strip)} \, dx \). Show that \( I = I_x + (\text{area})(c^2) \). This is the parallel axis theorem: \( I \) is smallest around the balancing axis \( c = 0 \).

23 (With thanks to Trivial Pursuit) In what state is the center of gravity of the United States—the "geographical center" or centroid?

24 Pappus (an ancient Greek) noticed that the volume is

\[ V = \int 2\pi y \text{(strip width)} \, dy = 2\pi M_x = 2\pi \bar{y} M \]

when a region of area \( M \) is revolved around the \( x \) axis. In the first step the solid was cut into ________.

25 Use this theorem of Pappus to find the volume of a torus. Revolve a disk of radius \( a \) whose center is at height \( \bar{y} = b > a \).

26 Rotate the triangle of Example 3 around the \( x \) axis and find the volume of the resulting cone—first from \( V = 2\pi \bar{y} M \), second from \( \frac{1}{3}\pi r^3 h \).

27 Find \( M_x \) and \( M_y \) for a thin wire along the semicircle \( y = \sqrt{1-x^2} \). Take \( \rho = 1 \) so \( M = \text{length} = \pi \).

28 A second theorem of Pappus gives \( A = 2\pi \bar{y} L \) as the surface area when a wire of length \( L \) is rotated around the \( x \) axis. Verify his formula for a horizontal wire along \( y = 3 \) \( (x = 0 \text{ to } x = L) \) and a vertical wire \( (y = 1 \text{ to } y = L+1) \).

29 The surface area of a sphere is \( A = 4\pi r^2 \) when \( r = 1 \). So \( A = 2\pi \bar{y} L \) leads to \( \bar{y} = ________ \) for the semicircular wire in Problem 27.

30 Rotating \( y = mx \) around the \( x \) axis between \( x = 0 \) and \( x = L \) produces the surface area \( A = ________ \).

31 Put a mass \( m \) at the point \((x, 0)\). Around the origin the torque from gravity is the force \( mg \) times the distance \( x \). This equals \( g \) times the ________ \( mx \).

32 If ten equal forces \( F \) are alternately down and up at \( x = 1, 2, \ldots, 10 \), what is their torque?

33 The solar system has nine masses \( m_i \) at distances \( r_i \) with angular velocities \( \omega_i \). What is the moment of inertia around the sun? What is the rotational energy? What is the torque provided by the sun?

34 The disk \( x^2 + y^2 \leq a^2 \) has \( I_0 = \int_0^a \pi r^4 \, dr = \frac{1}{4} \pi a^4 \). Why is this different from \( I_x \) in Problem 20? Find the radius of gyration \( r = \sqrt{I_0/M} \). (The rotational energy \( \frac{1}{2}I_0 \omega^2 \) equals \( \frac{1}{2}Mr^2 \omega^2 \)—when the whole mass is turning at radius \( r \)).

Questions 35–42 come from the moment of inertia experiment.

35 A solid cylinder of radius \( r \) is assembled from hollow cylinders of length \( l \), radius \( x \), and volume \( (2\pi x)(l)(dx) \). The solid cylinder has

\[ M = \int_0^l 2\pi x l \rho \, dx \quad \text{and} \quad I = \int_0^l x^2 2\pi x l \rho \, dx \]

With \( \rho = 7 \) find \( M \) and \( I \) and \( J = I/Mr^2 \).

36 Problem 14.4.40 finds \( J = 2/5 \) for a solid ball. It is less than \( J \) for a solid cylinder because the mass of the ball is more concentrated near ________.

37 Problem 14.4.39 finds \( J = \frac{1}{5} \int_0^\pi \sin^3 \phi \, d\phi = ________ \) for a hollow ball. The four rolling objects finish in the order ________.

38 By varying the density of the ball how could you make it roll faster than any of these shapes?

39 Answer Question 2 about the experiment.

40 For a vertical drop of \( y \), equation (12) gives the velocity along the plane: \( v^2 = 2gy/(1 + J) \). Thus \( v = cy^{1/2} \) for \( c = ________ \). The vertical velocity is \( dy/dt = v \sin \alpha \).

\[ dy/dt = cy^{1/2} \sin \alpha \quad \text{and} \quad \int y^{-1/2} \, dy = \int c \sin \alpha \, dt \]

Integrate to find \( y(t) \). Show that the bottom is reached \((y = \bar{h})\) at time \( T = 2\sqrt{\bar{h}/c \sin \alpha} \).

41 What is the theoretical ratio of the four finishing times?

42 True or false:

(a) Basketballs roll downhill faster than baseballs.
(b) The center of mass is always at the centroid.
(c) By putting your arms up you reduce \( I_x \) and \( I_y \).
(d) The center of mass of a high jumper goes over the bar (on successful jumps).
Chapter 1 introduced derivatives $df/dt$ and $df/dx$. The independent variable could be $t$ or $x$. For velocity it was natural to use the letter $t$. This section is about two important physical quantities—force and work—for which $x$ is the right choice.

The basic formula is $W = Fx$. Work equals force times distance moved (distance in the direction of $F$). With a force of 100 pounds on a car that moves 20 feet, the work is 2000 foot-pounds. If the car is rolling forward and you are pushing backward, the work is $-2000$ foot-pounds. If your force is only 80 pounds and the car doesn't move, the work is zero. In these examples the force is constant.

$W = Fx$ is completely parallel to $f = vt$. When $v$ is constant, we only need multiplication. It is a changing velocity that requires calculus. The integral $\int v(t) \, dt$ adds up small multiplications over short times. For a changing force, we add up small pieces of work $F \, dx$ over short distances:

$$W = Fx \quad \text{(constant force)} \quad W = \int F(x) \, dx \quad \text{(changing force)}.$$  

In the first case we lift a suitcase weighing $F = 30$ pounds up $x = 20$ feet of stairs. The work is $W = 600$ foot-pounds. The suitcase doesn't get heavier as we go up—it only seems that way. Actually it gets lighter (we study gravity below).

In the second case we stretch a spring, which needs more force as $x$ increases. Hooke's law says that $F(x) = kx$. The force is proportional to the stretching distance $x$. Starting from $x = 0$, the work increases with the square of $x$:

$$F = kx \quad \text{and} \quad W = \int_0^x kx \, dx = \frac{1}{2}kx^2.$$  

In metric units the force is measured in Newtons and the distance in meters. The unit of work is a Newton-meter (a joule). The 600 foot-pounds for an American suitcase would have been about 800 joules in France.

**EXAMPLE 1** Suppose a force of $F = 20$ pounds stretches a spring 1 foot.

(a) Find $k$. The elastic constant is $k = F/x = 20$ pounds per foot.

(b) Find $W$. The work is $\frac{1}{2}kx^2 = \frac{1}{2} \cdot 20 \cdot 1^2 = 10$ foot-pounds.

(c) Find $x$ when $F = -10$ pounds. This is compression not stretching: $x = -\frac{1}{2}$ foot.

Compressing the same spring through the same distance requires the same work. For compression $x$ and $F$ are negative. But the work $W = \frac{1}{2}kx^2$ is still positive. Please note that $W$ does not equal $kx$ times $x$! That is the whole point of variable force (change $Fx$ to $\int F(x) \, dx$).

May I add another important quantity from physics? It comes from looking at the situation from the viewpoint of the spring. In its natural position, the spring rests comfortably. It feels no strain and has no energy. Tension or compression gives it potential energy. More stretching or more compression means more energy. The change in energy equals the work. The potential energy of the suitcase increases by 600 foot-pounds, when it is lifted 20 feet.

Write $V(x)$ for the potential energy. Here $x$ is the height of the suitcase or the extension of the spring. In moving from $x = a$ to $x = b$, work = increase in potential:

$$W = \int_a^b F(x) \, dx = V(b) - V(a).$$  

This is absolutely beautiful. The work $W$ is the definite integral. The potential $V$ is the indefinite integral. If we carry the suitcase up the stairs and back down, our total
work is zero. We may feel tired, but the trip down should have given back our energy. (It was in the suitcase.) Starting with a spring that is compressed one foot, and ending with the spring extended one foot, again we have done no work. \( V = \frac{1}{2} k x^2 \) is the same for \( x = -1 \) and \( x = 1 \). But an extension from \( x = 1 \) to \( x = 3 \) requires work:

\[
W = \text{change in } V = \frac{1}{2} k (3)^2 - \frac{1}{2} k (1)^2.
\]

Indefinite integrals like \( V \) come with a property that we know well. They include an arbitrary constant \( C \). The correct potential is not simply \( \frac{1}{2} k x^2 \), it is \( \frac{1}{2} k x^2 + C \). To compute a change in potential, we don’t need \( C \). The constant cancels. But to determine \( V \) itself, we have to choose \( C \). By fixing \( V = 0 \) at one point, the potential is determined at all other points. A common choice is \( V = 0 \) at \( x = 0 \). Sometimes \( V = 0 \) at \( x = \infty \) (for gravity). Electric fields can be “grounded” at any point.

There is another connection between the potential \( V \) and the force \( F \). According to (2), \( V \) is the indefinite integral of \( F \). Therefore \( F(x) \) is the derivative of \( V(x) \). The fundamental theorem of calculus is also fundamental to physics:

\[
\text{force exerted on spring: } F = \frac{dV}{dx} \quad \text{(3a)}
\]

\[
\text{force exerted by spring: } F = -\frac{dV}{dx} \quad \text{(3b)}
\]

Those lines say the same thing. One is our force pulling on the spring, the other is the “restoring force” pulling back. (3a) and (3b) are a warning that the sign of \( F \) depends on the point of view. Electrical engineers and physicists use the minus sign. In mechanics the plus sign is more common. It is one of the ironies of fate that \( F = V' \), while distance and velocity have those letters reversed: \( v = f' \). Note the change to capital letters and the change to \( x \).

\[ \text{Fig. 8.16  Stretched spring; suitcase 20 feet up; moon of mass } m; \text{ oscillating spring.} \]

**EXAMPLE 2  Newton's law of gravitation** (inverse square law):

force to overcome gravity = \( GMm/x^2 \)  \quad \text{force exerted by gravity} = - \( GMm/x^2 \)

An engine pushes a rocket forward. Gravity pulls it back. The gravitational constant is \( G \) and the Earth's mass is \( M \). The mass of the rocket or satellite or suitcase is \( m \), and the potential is the indefinite integral:

\[
V(x) = \int F(x) \, dx = - \frac{GMm}{x} + C. \tag{4}
\]

Usually \( C = 0 \), which makes the potential zero at \( x = \infty \).

**Remark**  When carrying the suitcase upstairs, \( x \) changed by 20 feet. The weight was regarded as constant—which it nearly is. But an exact calculation of work uses the integral of \( F(x) \), not just the multiplication 30 times 20. The serious difference comes when the suitcase is carried to \( x = \infty \). With constant force that requires infinite work. With the correct (decreasing) force, the work equals \( V \) at infinity (which is zero) minus \( V \) at the pickup point \( x_0 \). The change in \( V \) is \( W = GMm/x_0 \).
KINETIC ENERGY

This optional paragraph carries the physics one step further. Suppose you release the spring or drop the suitcase. The external force changes to \( F = 0 \). But the internal force still acts on the spring, and gravity still acts on the suitcase. They both start moving. The potential energy of the suitcase is converted to kinetic energy, until it hits the bottom of the stairs.

Time enters the problem, either through Newton's law or Einstein's:

\[
\text{(Newton)} \quad F = ma = m \frac{dv}{dt} \quad \text{(Einstein)} \quad F = \frac{d}{dt} (mv).
\] (5)

Here we stay with Newton, and pretend the mass is constant. Exercise 21 follows Einstein; the mass increases with velocity. There \( m = m_0 \sqrt{1 - v^2/c^2} \) goes to infinity as \( v \) approaches \( c \), the speed of light. That correction comes from the theory of relativity, and is not needed for suitcases.

What happens as the suitcase falls? From \( x = a \) at the top of the stairs to \( x = b \) at the bottom, potential energy is lost. But kinetic energy \( \frac{1}{2}mv^2 \) is gained, as we see from integrating Newton's law:

\[
\text{force } F = m \frac{dv}{dt} = m \frac{dv}{dx} \frac{dx}{dt} = mv \frac{dv}{dx}
\]

\[
\text{work } \int_a^b F \, dx = \int_a^b m v \frac{dv}{dx} \, dx = \frac{1}{2} mv^2 (b) - \frac{1}{2} mv^2 (a). \tag{6}
\]

This same force \( F \) is given by \( -dV/dx \). So the work is also the change in \( V \):

\[
\text{work } \int_a^b F \, dx = \int_a^b \left( -\frac{dV}{dx} \right) \, dx = -V(b) + V(a). \tag{7}
\]

Since (6) = (7), the total energy \( \frac{1}{2}mv^2 + V \) (kinetic plus potential) is constant:

\[
\frac{1}{2}mv^2 (b) + V(b) = \frac{1}{2}mv^2 (a) + V(a). \tag{8}
\]

This is the law of conservation of energy. The total energy is conserved.

**EXAMPLE 3** Attach a mass \( m \) to the end of a stretched spring and let go. The spring's energy \( V = \frac{1}{2}kx^2 \) is gradually converted to kinetic energy of the mass. At \( x = 0 \) the change to kinetic energy is complete: the original \( \frac{1}{2}kx^2 \) has become \( \frac{1}{2}mv^2 \). Beyond \( x = 0 \) the potential energy increases, the force reverses sign and pulls back, and kinetic energy is lost. Eventually all energy is potential—when the mass reaches the other extreme. It is simple harmonic motion, exactly as in Chapter 1 (where the mass was the shadow of a circling ball). The equation of motion is the statement that the rate of change of energy is zero (and we cancel \( v = dx/dt \)):

\[
\frac{d}{dt} \left( \frac{1}{2} mv^2 + \frac{1}{2} kx^2 \right) = m \frac{dv}{dt} \frac{dx}{dt} + kx \frac{dx}{dt} = 0 \quad \text{or} \quad m \frac{d^2 x}{dt^2} + kx = 0. \tag{9}
\]

That is \( F = ma \) in disguise. For a spring, the solution \( x = \cos \sqrt{k/m} t \) will be found in this book. For more complicated structures, engineers spend a billion dollars a year computing the solution.
PRESSURE AND HYDROSTATIC FORCE

Our forces have been concentrated at a single points. That is not the case for pressure. A fluid exerts a force all over the base and sides of its container. Suppose a water tank or swimming pool has constant depth \( h \) (in meters or feet). The water has weight-density \( w \approx 9800 \text{ N/m}^3 \approx 62 \text{ lb/ft}^3 \). On the base, the pressure is \( w \) times \( h \). The force is \( wh \) times the base area \( A \):

\[
F = whA \text{ (pounds or Newtons)} \quad p = F/A = wh \text{ (lb/ft}^2 \text{ or N/m}^2). \quad (10)
\]

Thus pressure is force per unit area. Here \( p \) and \( F \) are computed by multiplication, because the depth \( h \) is constant. Pressure is proportional to depth (as divers know). Down the side wall, \( h \) varies and we need calculus.

The pressure on the side is still \( wh \)—the same in all directions. We divide the side into horizontal strips of thickness \( \Delta h \). Geometry gives the length \( l(h) \) at depth \( h \) (Figure 8.17). The area of a strip is \( l(h) \Delta h \). The pressure \( wh \) is nearly constant on the strip—the depth only changes by \( \Delta h \). The force on the strip is \( AF \). Adding those forces, and narrowing the strips so that \( \Delta h \to 0 \), the total force approaches an integral:

\[
\text{total force } F = \int wh(h) \, dh \quad (11)
\]

**EXAMPLE 4** Find the total force on the trapezoidal dam in Figure 8.17.

The side length is \( l = 60 \) when \( h = 0 \). The depth \( h \) increases from 0 to 20. The main problem is to find \( l \) at an in-between depth \( h \). With straight sides the relation is linear: \( l = 60 + ch \). We choose \( c \) to give \( l = 50 \) when \( h = 20 \). Then \( 50 = 60 + c(20) \) yields \( c = -\frac{1}{2} \).

The total force is the integral of \( whl \). So substitute \( l = 60 - \frac{1}{2}h \):

\[
F = \int_0^{20} wh(60 - \frac{1}{2}h) \, dh = \left[ 30wh^2 - \frac{1}{8}wh^3 \right]_0^{20} = 12000w - \frac{1}{8}(8000w).
\]

With distance in feet and \( w = 62 \text{ lb/ft}^3 \), \( F \) is in pounds. With distance in meters and \( w = 9800 \text{ N/m}^3 \), the force is in Newtons.

Note that \( \text{weight-density } w = \text{mass-density } \rho \times \text{gravitational acceleration } g = (1000)(9.8) \). These SI units were chosen to make the density of water at \( 0^\circ \text{C} \) exactly \( \rho = 1000 \text{ kg/m}^3 \).

**EXAMPLE 5** Find the work to pump water out of a tank. The area at depth \( h \) is \( A(h) \).

Imagine lifting out one layer of water at a time. The layer weighs \( wA(h) \Delta h \). The work to lift it to the top is its weight times the distance \( h \), or \( whA(h) \Delta h \). The work to empty the whole tank is the integral:

\[
W = \int whA(h) \, dh \quad (12)
\]
Suppose the tank is the bottom half of a sphere of radius \( R \). The cross-sectional area at depth \( h \) is \( A = \pi (R^2 - h^2) \). Then the work is the integral (12) from 0 to \( R \). It equals \( W = \pi w R^4/4 \).

Units: \( w = \text{force}/(\text{distance})^3 \times R^4 = \text{(distance)}^2 \) gives work \( W = \text{(force)} \times \text{(distance)} \).

---

### Read-through questions

Work equals \( a \) times \( b \). For a spring the force is \( F = ax \), proportional to the extension \( x \) (this is \( a \) law).

With this variable force, the work in stretching from 0 to \( x \) is \( W = \int_0^x \text{work} \). This equals the increase in the \( g \) energy \( V \). Thus \( W \) is a \( \text{area} \) integral and \( V \) is the corresponding \( \text{area} \) integral, which includes an arbitrary \( f \). The derivative \( dV/dx \) equals \( k \). The force of gravity is \( F = \_ \) and the potential is \( V = \_ m \).

In falling, \( V \) is converted to \( k \) energy \( K = \_ \). The total energy \( K + V \) is \( \_ \) (this is the law of \( g \) when there is no external force).

Pressure is force per unit \( \_ \). Water of density \( w \) in a pool of depth \( h \) and area \( A \) exerts a downward force \( F = \_ \) on the base. The pressure is \( p = \_ \). On the sides the \( \_ \) is still \( wh \) at depth \( h \), so the total force is \( \int w \text{h} \text{dh} \), where \( \_ \). In a cubic pool of side \( s \), the force on the base is \( F = \_ \), the length around the sides is \( \_ \), and the total force on the four sides is \( F = \_ \).

The work to pump the water out of the pool is \( W = \int w h A \text{dh} = \_ \).

1. (a) Find the work \( W \) when a constant force \( F = 12 \) pounds moves an object from \( x = .9 \) feet to \( x = 1.1 \) feet.
   (b) Compute \( W \) by integration when the force \( F = 12/x^3 \) varies with \( x \).

2. A 12-inch spring is stretched to 15 inches by a force of 75 pounds.
   (a) What is the spring constant \( k \) in \( \text{pounds per foot} \)?
   (b) Find the work done in stretching the spring.
   (c) Find the work to stretch it 3 more inches.

3. A shock-absorber is compressed 1 inch by a weight of 1 ton. Find its spring constant \( k \) in \( \text{pounds per foot} \). What potential energy is stored in the shock-absorber?

4. A force \( F = 20x - x^3 \) stretches a nonlinear spring by \( x \).
   (a) What work is required to stretch it from \( x = 0 \) to \( x = 2 \)?
   (b) What is its potential energy \( V \) at \( x = 2 \), if \( V(0) = 5 \)?
   (c) What is \( k = \text{dF/dx} \) for a small additional stretch at \( x = 2 \)?

5. (a) A 120-lb person makes a scale go down \( x \) inches. How much work is done?
   (b) If the same person goes \( x \) inches down the stairs, how much potential energy is lost?

6. A rocket burns its 100 kg of fuel at a steady rate to reach a height of 25 km.
   (a) Find the weight of fuel left at height \( h \).
   (b) How much work is done lifting fuel?

7. Integrate to find the work in winding up a hanging cable of length 100 feet and weight density 5 lb/ft. How much additional work is caused by a 200-pound weight hanging at the end of the cable?

8. The great pyramid (height 500' — you can see it from Cairo) has a square base 800' by 800'. Find the area \( A \) at height \( h \). If the rock weighs \( w = 100 \) lb/ft\(^3\), approximately how much work did it take to lift all the rock?

9. The force of gravity on a mass \( m \) is \( F = -G M m/x^2 \). With \( G = 6 \times 10^{-11} \) and Earth mass \( M = 6 \times 10^{24} \) and rocket mass \( m = 1000 \), compute the work to lift the rocket from \( x = 6400 \) to \( x = 6500 \). (The units are kgs and kms and Newtons, giving work in Newton-kms.)

10. The approximate work to lift a 30-pound suitcase 20 feet is 600 foot-pounds. The exact work is the change in the potential \( V = -G M m/x \). Show that \( \Delta V \) is 600 times a correction factor \( R^2/(R^2 - 10^2) \), when \( x \) changes from \( R - 10 \) to \( R + 10 \). (This factor is practically 1, when \( R \) radius of the Earth.)

11. Find the work to lift the rocket in Problem 9 from \( x = 6400 \) out to \( x = \infty \). If this work equals the original kinetic energy \( \frac{1}{2} m u^2 \), what was the original \( v \) (the escape velocity)?

12. The kinetic energy \( \frac{1}{2} m u^2 \) of a rocket is converted into potential energy \( -G M m/x \). Starting from the Earth's radius \( x = R \), what \( x \) does the rocket reach? If it reaches \( x = \infty \) show that \( v = \sqrt{2GM/R} \). This escape velocity is 25,000 miles per hour.

13. It takes 20 foot-pounds of work to stretch a spring 2 feet. How much work to stretch it one more foot?

14. A barrel full of beer is 4 feet high with a 1 foot radius and an opening at the bottom. How much potential energy is lost by the beer as it comes out of the barrel?
15 A rectangular dam is 40 feet high and 60 feet wide. Compute the total side force \( F \) on the dam when (a) the water is at the top, (b) the water level is halfway up.

16 A triangular dam has an 80-meter base at a depth of 30 meters. If water covers the triangle, find
   (a) the pressure at depth \( h \)
   (b) the length \( l \) of the dam at depth \( h \)
   (c) the total force on the dam.

17 A cylinder of depth \( H \) and cross-sectional area \( A \) stands full of water (density \( w \)). (a) Compute the work \( W = \int wAh \, dh \) to lift all the water to the top. (b) Check the units of \( W \). (c) What is the work \( W \) if the cylinder is only half full?

18 In Problem 17, compute \( W \) in both cases if \( H = 20 \) feet, \( w = 62 \) lb/ft\(^3\), and the base is a circle of radius \( r = 5 \) feet.

19 How much work is required to pump out a swimming pool, if the area of the base is 800 square feet, the water is 4 feet deep, and the top is one foot above the water level?

20 For a cone-shaped tank the cross-sectional area increases with depth: \( A = \pi r^2 h^2 / H^2 \). Show that the work to empty it is half the work for a cylinder with the same height and base. What is the ratio of volumes of water?

21 In relativity the mass is \( m = m_0 \sqrt{1 - v^2 / c^2} \). Find the correction factor in Newton's equation \( F = \dot{m}a \) to give Einstein's equation \( F = d(mv)/dt = (d(mv)/dv)(dv/dt) = \ldots m_0a \).

22 Estimate the depth of the Titanic, the pressure at that depth, and the force on a cabin door. Why doesn’t every door collapse at the bottom of the Atlantic Ocean?

23 A swimming pool is 4 meters wide, 10 meters long, and 2 meters deep. Find the weight of the water and the total force on the bottom.

24 If the pool in Problem 23 has a shallow end only one meter deep, what fraction of the water is saved? Draw a cross-section (a trapezoid) and show the direction of force on the sides and the sloping bottom.

25 In what ways is work like a definite integral and energy like an indefinite integral? Their derivative is the \ldots \ldots \ldots \ldots \ldots \ldots
# Contents

## CHAPTER 9  
**Polar Coordinates and Complex Numbers**

- **9.1 Polar Coordinates** 348  
- **9.2 Polar Equations and Graphs** 351  
- **9.3 Slope, Length, and Area for Polar Curves** 356  
- **9.4 Complex Numbers** 360

## CHAPTER 10  
**Infinite Series**

- **10.1 The Geometric Series** 368  
- **10.2 Convergence Tests: Positive Series** 374  
- **10.3 Convergence Tests: All Series** 381  
- **10.4 The Taylor Series for \( e^x, \sin x, \) and \( \cos x \)** 385  
- **10.5 Power Series** 391

## CHAPTER 11  
**Vectors and Matrices**

- **11.1 Vectors and Dot Products** 398  
- **11.2 Planes and Projections** 407  
- **11.3 Cross Products and Determinants** 416  
- **11.4 Matrices and Linear Equations** 425  
- **11.5 Linear Algebra in Three Dimensions** 435

## CHAPTER 12  
**Motion along a Curve**

- **12.1 The Position Vector** 446  
- **12.2 Plane Motion: Projectiles and Cycloids** 453  
- **12.3 Tangent Vector and Normal Vector** 459  
- **12.4 Polar Coordinates and Planetary Motion** 464

## CHAPTER 13  
**Partial Derivatives**

- **13.1 Surfaces and Level Curves** 472  
- **13.2 Partial Derivatives** 475  
- **13.3 Tangent Planes and Linear Approximations** 480  
- **13.4 Directional Derivatives and Gradients** 490  
- **13.5 The Chain Rule** 497  
- **13.6 Maxima, Minima, and Saddle Points** 504  
- **13.7 Constraints and Lagrange Multipliers** 514
Polar Coordinates and Complex Numbers

9.1 Polar Coordinates

Up to now, points have been located by their x and y coordinates. But if you were a flight controller, and a plane appeared on the screen, you would not give its position that way. Instead of x and y, you would read off the direction of the plane and its distance. The direction is given by an angle \( \theta \). The distance is given by a positive number \( r \). Those are the polar coordinates of the point, where x and y are the rectangular coordinates.

The angle \( \theta \) is measured from the horizontal. Suppose the distance is 2 and the direction is 30° or \( \pi/6 \) (degrees preferred by flight controllers, radians by mathematicians). A pilot looking along the x axis would give the plane's direction as "11 o'clock." This totally destroys our system of units, by measuring direction in hours. But the angle and the distance locate the plane.

How far to a landing strip at \( r = 1 \) and \( \theta = -\pi/2 \)? For that question polar coordinates are not good. They are perfect for distance from the origin (which equals \( r \)), but for most other distances I would switch to x and y. It is extremely simple to determine x and y from \( r \) and \( \theta \), and we will do it constantly. The most used formulas in this chapter come from Figure 9.1—where the right triangle has angle \( \theta \) and hypotenuse \( r \). The sides of that triangle are x and y:

\[
x = r \cos \theta \quad \text{and} \quad y = r \sin \theta.
\] (1)

The point at \( r = 2, \theta = \pi/6 \) has \( x = 2 \cos(\pi/6) \) and \( y = 2 \sin(\pi/6) \). The cosine of \( \pi/6 \) is \( \sqrt{3}/2 \) and the sine is \( \frac{1}{2} \). So \( x = \sqrt{3} \) and \( y = 1 \). Polar coordinates convert easily to xy coordinates—now we go the other way.

Always \( x^2 + y^2 = r^2 \). In this example \( (\sqrt{3})^2 + (1)^2 = (2)^2 \). Pythagoras produces \( r \) from x and y. The direction \( \theta \) is also available, but the formula is not so beautiful:

\[
r = \sqrt{x^2 + y^2} \quad \text{and} \quad \tan \theta = \frac{y}{x} \quad \text{and (almost)} \quad \theta = \tan^{-1} \frac{y}{x}.
\] (2)

Our point has \( y/x = 1/\sqrt{3} \). One angle with this tangent is \( \theta = \tan^{-1} (1/\sqrt{3}) = \pi/6 \).
EXAMPLE 1  Point B in Figure 9.1c is at a negative angle \( \theta = -\pi/4 \). The \( x \) coordinate \( r \cos(-\pi/4) \) is the same as \( r \cos \pi/4 \) (the cosine is even). But the \( y \) coordinate \( r \sin(-\pi/4) \) is negative. Computing \( r \) and \( \theta \) from \( x = 1 \) and \( y = 1 \), the distance is \( r = \sqrt{1+1} \) and \( \tan \theta = -1/\sqrt{2} \).

Warning  To any angle \( \theta \) we can add or subtract 2\( \pi \)—which goes a full 360° circle and keeps the same direction. Thus \( -\pi/4 \) or \( -45° \) is the same angle as \( 7\pi/4 \) or \( 315° \). So is \( 15\pi/4 \) or \( 675° \).

If we add or subtract 180°, the tangent doesn’t change. The point \((1, -1)\) is on the \(-45°\) line at \( r = \sqrt{2} \). The point \((-1, 1)\) is on the \(135°\) line also with \( r = \sqrt{2} \). Both have \( \tan \theta = -1 \). We had to write “almost” in equation (2), because a point has many \( \theta \)'s and two points have the same \( r \) and \( \tan \theta \).

Even worse, we could say that \( B = (1, -1) \) is on the \(135°\) line but at a negative distance \( r = -\sqrt{2} \). A negative \( r \) carries the point backward along the \(135°\) line, which is forward to \( B \). In giving the position of \( B \), I would always keep \( r > 0 \). But in drawing the graph of a polar equation, \( r < 0 \) is allowed. We move now to those graphs.

THE CIRCLE \( r = \cos \theta \)

The basis for Chapters 1–8 was \( y = f(x) \). The key to this chapter is \( r = F(\theta) \). That is a relation between the polar coordinates, and the points satisfying an equation like \( r = \cos \theta \) produce a polar graph.

It is not obvious why \( r = \cos \theta \) gives a circle. The equations \( r = \cos 2\theta \) and \( r = \cos^2 \theta \) produce entirely different graphs—not circles. The direct approach is to take \( \theta = 0°, 30°, 60°, \ldots \) and go out the distance \( r = \cos \theta \) on each ray. The points are marked in Figure 9.2a, and connected into a curve. It seems to be a circle of radius \( \frac{1}{2} \), with its center at the point \((\frac{1}{2}, 0)\). We have to be able to show mathematically that \( r = \cos \theta \) represents a shifted circle.

One point must be mentioned. The angles from 0 to \( \pi \) give the whole circle. The number \( r = \cos \theta \) becomes negative after \( \pi/2 \), and we go backwards along each ray.

![Fig. 9.2](image-url)  The circle \( r = \cos \theta \) and the switch to \( x \) and \( y \). The circle \( r = \sin \theta \).
At \(\theta = \pi\) (to the left of the origin) the cosine is \(-1\). Going backwards brings us to the same point as \(\theta = 0\) and \(r = 1\) — which completes the circle.

When \(\theta\) continues from \(\pi\) to \(2\pi\) we go around again. The polar equation gives the circle twice. (Or more times, when \(\theta\) continues past \(2\pi\).) If you don’t like negative \(r\)'s and multiple circles, restrict \(\theta\) to the range from \(-\pi/2\) to \(\pi/2\). We still have to see why the graph of \(r = \cos \theta\) is a circle.

**Method 1** Multiply by \(r\) and convert to rectangular coordinates \(x\) and \(y\):

\[
r = \cos \theta \implies r^2 = r \cos \theta \implies x^2 + y^2 = x.
\]

This is a circle because of \(x^2 + y^2\). From rewriting as \((x - \frac{1}{2})^2 + y^2 = (\frac{1}{2})^2\) we recognize its center and radius. Center at \(x = \frac{1}{2}\) and \(y = 0\); radius \(\frac{1}{2}\). Done.

**Method 2** Write \(x\) and \(y\) separately as functions of \(\theta\). Then \(\theta\) is a “parameter”:

\[
x = r \cos \theta = \cos^2 \theta \quad \text{and} \quad y = r \sin \theta = \sin \theta \cos \theta.
\]

These are not polar equations but parametric equations. The parameter \(\theta\) is the angle, but it could be the time — the curve would be the same. Chapter 12 studies parametric equations in detail — here we stay with the circle.

To find the circle, square \(x\) and \(y\) and add. This produces \(x^2 + y^2 = x\) in Problem 26. But here we do something new: Start with the circle and find equation (4). In case you don’t reach Chapter 12, the idea is this. Add the vectors \(OC\) to the center and \(CP\) out the radius:

The point \(P\) in Figure 9.2 has \((x, y) = OC + CP = (\frac{1}{2}, 0) + (\frac{1}{2} \cos t, \frac{1}{2} \sin t)\).

The parameter \(t\) is the angle at the center of the circle. The equations are \(x = \frac{1}{2} + \frac{1}{2} \cos t\) and \(y = \frac{1}{2} \sin t\). A trigonometric person sees a double angle and sets \(t = 2\theta\). The result is equation (4) for the circle:

\[
x = \frac{1}{2} + \frac{1}{2} \cos 2\theta = \cos^2 \theta \quad \text{and} \quad y = \frac{1}{2} \sin 2\theta = \sin \theta \cos \theta.
\]

This step rediscovers a basic theorem of geometry: The angle \(t\) at the center is twice the angle \(\theta\) at the circumference. End of quick introduction to parameters.

A second circle is \(r = \sin \theta\), drawn in Figure 9.2c. A third circle is \(r = \cos \theta + \sin \theta\), not drawn. Problem 27 asks you to find its \(xy\) equation and its radius. All calculations go back to \(x = r \cos \theta\) and \(y = r \sin \theta\) — the basic facts of polar coordinates! The last exercise shows a parametric equation with beautiful graphs, because it may be possible to draw them now. Then the next section concentrates on \(r = F(\theta)\) — and goes far beyond circles.

### 9.1 Exercises

**Read-through questions**

Polar coordinates \(r\) and \(\theta\) correspond to \(x = \_a\_\) and \(y = \_b\_.\) The points with \(r > 0\) and \(\theta = \pi\) are located \(\_c\_.\)

The points with \(r = 1\) and \(0 \leq \theta \leq \pi\) are located \(\_d\_.\) Reversing the sign of \(\theta\) moves the point \((x, y)\) to \(\_e\_.\)

Given \(x\) and \(y\), the polar distance is \(r = \_f\_.\) The tangent of \(\theta\) is \(\_g\_.\) The point \((6, 8)\) has \(r = \_h\_\) and \(\theta = \_i\_.\)

Another point with the same \(\theta\) is \(\_j\_.\) Another point with the same \(r\) is \(\_k\_.\) Another point with the same \(r\) and tan \(\theta\) is \(\_l\_.\)

The polar equation \(r = \cos \theta\) produces a shifted \(\_m\_.\) The top point is at \(\theta = \_n\_.\) which gives \(r = \_o\_.\) When \(\theta\) goes from \(0\) to \(2\pi\), we go \(\_p\_.\) times around the graph.

Rewriting as \(r^2 = r \cos \theta\) leads to the \(xy\) equation \(\_q\_.\) Substituting \(r = \cos \theta\) into \(x = r \cos \theta\) yields \(x = \_r\_.\) and similarly \(y = \_s\_.\) In this form \(x\) and \(y\) are functions of the \(\_t\_.\)

Find the polar coordinates \(r > 0\) and \(0 \leq \theta < 2\pi\) of these points.

1. \((x, y) = (0, 1)\)
2. \((x, y) = (-4, 0)\)
3 \((x, y) = (\sqrt{2}, \sqrt{2})\) 
4 \((x, y) = (-1, \sqrt{3})\) 
5 \((x, y) = (-1, -1)\) 
6 \((x, y) = (3, 4)\)

Find rectangular coordinates \((x, y)\) from polar coordinates.
7 \((r, \theta) = (2, \pi/2)\) 
8 \((r, \theta) = (1, 3\pi/2)\) 
9 \((r, \theta) = (\sqrt{2}, \pi/4)\) 
10 \((r, \theta) = (3\pi, 3\pi)\) 
11 \((r, \theta) = (2, -\pi/6)\) 
12 \((r, \theta) = (2, 5\pi/6)\)

13 What is the distance from \((x, y) = (3, 1)\) to \((1, -\pi)\)?

14 How far is the point \(r = 3, \theta = \pi/2\) from \(r = 4, \theta = \pi\)?

15 How far is \((x, y) = (r \cos \theta, r \sin \theta)\) from \((X, Y) = (R \cos \phi, R \sin \phi)\)? Simplify \((x - X)^2 + (y - Y)^2\) by using \(\cos(\theta - \phi) = \cos \theta \cos \phi + \sin \theta \sin \phi\).

16 Find a second set of polar coordinates (a different \(r\) or \(\theta\)) for the points
\((r, \theta) = (-1, \pi/2), (-1, 3\pi/4), (1, -\pi/2), (0, 0)\).

17 Using polar coordinates describe (a) the half-plane \(x > 0\); (b) the half-plane \(y < 0\); (c) the ring with \(x^2 + y^2\) between \(4 \) and \(5\); (d) the wedge \(x \geq |y|\).

18 True or false, with a reason or an example:
(a) Changing to \(-r\) and \(-\theta\) produces the same point.
(b) Each point has only one \(r\) and \(\theta\), when \(r < 0\) is not allowed.
(c) The graph of \(r = 1/\sin \theta\) is a straight line.

19 From \(x\) and \(\theta\) find \(y\) and \(r\).

20 Which other point has the same \(r\) and \(\tan \theta\) as \(x = \sqrt{3}, y = 1\) in Figure 9.1b?

21 Convert from rectangular to polar equations:
(a) \(y = x\) 
(b) \(x + y = 1\) 
(c) \(x^2 + y^2 = x + y\)

22 Show that the triangle with vertices at \((0, 0), (r_1, \theta_1),\) and \((r_2, \theta_2)\) has area \(A = \frac{1}{2} r_1 r_2 \sin(\theta_2 - \theta_1)\). Find the base and height assuming \(0 \leq \theta_1 \leq \theta_2 \leq \pi\).

Problems 23–28 are about polar equations that give circles.

23 Convert \(r = \sin \theta\) into an \(xy\) equation. Multiply first by \(r\).

24 Graph \(r = \sin \theta\) at \(\theta = 0^\circ, 30^\circ, 60^\circ, ..., 360^\circ\). These thirteen values of \(\theta\) give ______ different points on the graph. What range of \(\theta\)'s goes once around the circle?

25 Substitute \(r = \sin \theta\) into \(x = r \cos \theta\) and \(y = r \sin \theta\) to find \(x\) and \(y\) in terms of the parameter \(\theta\). Then compute \(x^2 + y^2\) to reach the \(xy\) equation.

26 From the parametric equations \(x = \cos^2 \theta\) and \(y = \sin \theta \cos \theta\) in (4), recover the \(xy\) equation. Square, add, eliminate \(\theta\).

27 (a) Multiply \(r = \cos \theta + \sin \theta\) by \(r\) to convert into an \(xy\) equation. (b) Rewrite the equation as \((x - \frac{1}{2})^2 + (y - \frac{1}{2})^2 = R^2\) to find the radius \(R\). (c) Draw the graph.

28 Find the radius of \(r = a \cos \theta + b \sin \theta\). (Multiply by \(r\).)

29 Convert \(x + y = 1\) into an \(r\theta\) equation and solve for \(r\). Then substitute this \(r\) into \(x = r \cos \theta\) and \(y = r \sin \theta\) to find parametric equations for the line.

30 The equations \(x = \cos^2 \theta\) and \(y = \sin^2 \theta\) also lead to \(x + y = 1\)—but they are different from the answer to Problem 29. Explanation: \(\theta\) is no longer the polar angle and we should have written \(t\). Find a point \(x = \cos^2 \theta, y = \sin^2 \theta\) that is not at the angle \(\theta\).

31 Convert \(r = \cos^2 \theta\) into an \(xy\) equation (of sixth degree!)

32 If you have a graphics package for parametric curves, graph some hypocycloids. The equations are \(x = (1 - b) \cos t + b \cos(1 - b)t/b, y = (1 - b) \sin t - b \sin(1 - b)t/b\). The figure shows \(b = \pi/6\) and part of \(b = .31831\).
This section studies polar curves $r = F(\theta)$. The cardioid is a sentimental favorite—maybe parabolas are more practical. The cardioid is $r = 1 + \cos \theta$, the parabola is $r = 1/(1 + \cos \theta)$. Section 12.2 adds cycloids and astroids. A graphics package can draw them and so can we.

Together with the circles $r = \text{constant}$ go the straight lines $\theta = \text{constant}$. The equation $\theta = \pi/4$ is a ray out from the origin, at that fixed angle. If we allow $r < 0$, as we do in drawing graphs, the one-directional ray changes to a full line. Important: The circles are perpendicular to the rays. We have "orthogonal coordinates"—more interesting than the $x - y$ grid of perpendicular lines. In principle $x$ could be mixed with $\theta$ (non-orthogonal), but in practice that never happens.

Other curves are attractive in polar coordinates—we look first at five examples. Sometimes we switch back to $x = r \cos \theta$ and $y = r \sin \theta$, to recognize the graph.

Example 1 The graph of $r = 1/\cos \theta$ is the straight line $x = 1$ (because $r \cos \theta = 1$).

Example 2 The graph of $r = \cos 2\theta$ is the four-petal flower in Figure 9.3.

The points at $\theta = 30^\circ$ and $-30^\circ$ and $150^\circ$ and $-150^\circ$ are marked on the flower. They all have $r = \cos 2\theta = 1/2$. There are three important symmetries—across the $x$ axis, across the $y$ axis, and through the origin. This four-petal curve has them all. So does the vertical flower $r = \sin 2\theta$—but surprisingly, the tests it passes are different.

(Across the $x$ axis: $y$ to $-y$) There are two ways to cross. First, change $\theta$ to $-\theta$. The equation $r = \cos 2\theta$ stays the same. Second, change $\theta$ to $\pi - \theta$ and also $r$ to $-r$. The equation $r = \sin 2\theta$ stays the same. Both flowers have $x$ axis symmetry.

(Across the $y$ axis: $x$ to $-x$) There are two ways to cross. First, change $\theta$ to $\pi - \theta$. The equation $r = \cos 2\theta$ stays the same. Second, change $\theta$ to $-\theta$ and $r$ to $-r$. Now $r = \sin 2\theta$ stays the same (the sine is odd). Both curves have $y$ axis symmetry.

(Through the origin) Now we change $r$ to $-r$ or $\theta$ to $\theta + \pi$. The flower equations pass the second test only: $\cos 2(\theta + \pi) = \cos 2\theta$ and $\sin 2(\theta + \pi) = \sin 2\theta$. Every equation $r^2 = F(\theta)$ passes the first test, since $(-r)^2 = r^2$.

The circle $r = \cos \theta$ has $x$ axis symmetry, but not $y$ or $r$. The spiral $r = \theta^2$ has $y$ axis symmetry, because $-r = (-\theta)^2$ is the same equation.

Question What happens if you change $r$ to $-r$ and also change $\theta$ to $\theta + \pi$?
Answer Nothing—because $(r, \theta)$ and $(-r, \theta + \pi)$ are always the same point.

Fig. 9.3 The four-petal flower $r = \cos 2\theta$ and the spiral $r = \theta$ ($r > 0$ in red).
EXAMPLE 3  The graph of \( r = \theta \) is a spiral of Archimedes—or maybe two spirals.

The spiral adds new points as \( \theta \) increases past \( 2\pi \). Our other examples are "periodic"—\( \theta = 2\pi \) gives the same point as \( \theta = 0 \). A periodic curve repeats itself. The spiral moves out by \( 2\pi \) each time it comes around. If we allow negative angles and negative \( r = \theta \), a second spiral appears.

EXAMPLE 4  The graph of \( r = 1 + \cos \theta \) is a cardioid. It is drawn in Figure 9.4c.

The cardioid has no simple \( xy \) equation. Still the curve is very attractive. It has a cusp at the origin and it is heart-shaped (hence its name). To draw it, plot \( r = 1 + \cos \theta \)

at 30° intervals and connect the points. For this curve \( r \) is never negative, since \( \cos \theta \) never goes below \(-1\).

It is a curious fact that the electrical vector in your heart almost traces out a cardioid. See Section 11.1 about electrocardiograms. If it is a perfect cardioid you are in a little trouble.

\[ r = 1 + \frac{1}{3} \cos \theta \]

no dimple

\[ r = 1 + \frac{2}{3} \cos \theta \]

dimple

\[ r = 1 + \cos \theta \]

cardioid

\[ r = 1 + \frac{2}{3} \cos \theta \]

inner loop

Fig. 9.4 Limaçons \( r = 1 + b \cos \theta \), including a cardioid and Mars seen from Earth.

EXAMPLE 5  The graph of \( r = 1 + b \cos \theta \) is a limaçon (a cardioid when \( b = 1 \)).

Limaçon (soft \( c \)) is a French word for snail—not so well known as escargot but just as indible. (I am only referring to the shell. Excusez-moi!) Figure 9.4 shows how a dimple appears as \( b \) increases. Then an inner loop appears beyond \( b = 1 \) (the cardioid at \( b = 1 \) is giving birth to a loop). For large \( b \) the curve looks more like two circles. The limiting case is a double circle, when the inner loop is the same as the outer loop. Remember that \( r = \cos \theta \) goes around the circle twice.

We could magnify the limaçon by a factor \( c \), changing to \( r = c(1 + b \cos \theta) \). We could rotate 180° to \( r = 1 - b \cos \theta \). But the real interest is whether these figures arise in applications, and Donald Saari showed me a nice example.

Mars seen from Earth  The Earth goes around the Sun and so does Mars. Roughly speaking Mars is \( 1 \frac{1}{2} \) times as far out, and completes its orbit in two Earth years.

We take the orbits as circles: \( r = 2 \) for Earth and \( r = 3 \) for Mars. Those equations tell where but not when. With time as a parameter, the coordinates of Earth and Mars are given at every instant \( t \):

\[ x_e = 2 \cos 2\pi t, \ y_e = 2 \sin 2\pi t \quad \text{and} \quad x_M = 3 \cos \pi t, \ y_M = 3 \sin \pi t. \]
At \( t = 1 \) year, the Earth completes a circle (angle = \( 2\pi \)) and Mars is halfway.

Now the key step. Subtract to find the position of Mars relative to Earth:

\[
x_{M-E} = 3 \cos \pi t - 2 \cos 2\pi t \quad \text{and} \quad y_{M-E} = 3 \sin \pi t - 2 \sin 2\pi t.
\]

Replacing \( \cos 2\pi t \) by \( 2 \cos^2 \pi t - 1 \) and \( \sin 2\pi t \) by \( 2 \sin \pi t \cos \pi t \), this is

\[
x_{M-E} = (3 - 4 \cos \pi t) \cos \pi t + 2 \quad \text{and} \quad y_{M-E} = (3 - 4 \cos \pi t) \sin \pi t.
\]

Seen from the Earth, Mars does a loop in the sky! There are two \( t \)'s for which

\[3 - 4 \cos \pi t = 0\] (or \( \cos \pi t = \frac{3}{4}\)). At both times, Mars is two units from Earth \( (x_{M-E} = 2 \) and \( y_{M-E} = 0)\). When we move the origin to that point, the 2 is subtracted away—the M - E coordinates become \( x = r \cos \pi t \) and \( y = r \sin \pi t \) with \( r = 3 - 4 \cos \pi t\). That is a limaçon with a loop, like Figure 9.4d.

**Note added in proof** I didn't realize that a 3-to-2 ratio is also responsible for heating up two spots on opposite sides of Mercury. From the newspaper of June 13, 1990:

"Astronomers today reported the first observations showing that Mercury has two extremely hot spots. That is because Mercury, the planet closest to the Sun, turns on its axis once every 59.6 days, which is a day on Mercury. It goes around the sun every 88 days, a Mercurian year. With this 3-to-2 ratio between spin and revolution, the Sun appears to stop in the sky and move backward, describing a loop over each of the hot spots."

### CONIC SECTIONS IN POLAR COORDINATES

The exercises include other polar curves, like lemniscates and 200-petal flowers. But get serious. The most important curves are the **ellipse** and **parabola** and **hyperbola**.

In Section 3.5 their equations involved \( 1, x, y, x^2, xy, y^2 \). With one focus at the origin, their polar equations are even better.

**9A** The graph of \( r = A/(1 + e \cos \theta) \) is a conic section with "eccentricity" \( e: \)

- circle if \( e = 0 \)
- ellipse if \( 0 < e < 1 \)
- parabola if \( e = 1 \)
- hyperbola if \( e > 1 \).

**EXAMPLE 6** \((e = 1)\) The graph of \( r = 1/(1 + \cos \theta) \) is a parabola. This equation is \( r + r \cos \theta = 1 \) or \( r = 1 - x \). Squaring both sides gives \( x^2 + y^2 = 1 - 2x + x^2 \). Canceling \( x^2 \) leaves \( y^2 = 1 - 2x \), the parabola in Figure 9.5b.

The amplifying factor \( A \) blows up all curves, with no change in shape.

**EXAMPLE 7** \((e = 2)\) The same steps lead from \( r(1 + 2 \cos \theta) = 1 \) to \( r = 1 - 2x \). Squaring gives \( x^2 + y^2 = 1 - 4x + 4x^2 \) and the \( x^2 \) terms do not cancel. Instead we have \( y^2 - 3x^2 = 1 - 4x \). This is the hyperbola in Figure 9.5c, with a focus at \((0, 0)\).

The hyperbola \( y^2 - 3x^2 = 1 \) (without the \(-4x\)) has its center at \((0, 0)\).

**EXAMPLE 8** \((e = \frac{3}{2})\) The same steps lead from \( r(1 + \frac{3}{2} \cos \theta) = 1 \) to \( r = 1 - \frac{3}{2}x \). Squaring gives the ellipse \( x^2 + y^2 = 1 - x + \frac{3}{4}x^2 \). Polar equations look at conics in a new way, which happens to match the sun and planets perfectly. The sun at \((0, 0)\) is not the center of the system, but a focus.

Finally \( e = 0 \) gives the circle \( r = 1 \). Center of circle = both foci = \((0, 0)\).
9.2 Polar Equations and Graphs

**Fig. 9.5** $r = 1/(1 + e \cos \theta)$ is an ellipse for $e = \frac{1}{2}$, a parabola for $e = 1$, a hyperbola for $e = 2$.

### The Directrix

The figure shows the line $d$ (the "directrix") for each curve. All points $P$ on the curve satisfy $r = |PF| = e|Pd|$. The distance to the focus is $e$ times the distance to the directrix. ($e$ is still the eccentricity, nothing to do with exponentials.) A geometer would start from this property $r = e|Pd|$ and construct the curve. We derive the property from the equation:

$$r = \frac{A}{1 + e \cos \theta} \Rightarrow r + ex = A \Rightarrow r = d\left(\frac{A}{e} - x\right). \quad (1)$$

The directrix is the line at $x = A/e$. That last equation is exactly $|PF| = e|Pd|$.

Notice how two numbers determine these curves. Here the numbers are $A$ and $e$. In Section 3.5 they were $a$ and $b$. (The ellipse was $x^2/a^2 + y^2/b^2 = 1$.) Using $A$ and $e$ we go smoothly from ellipses through parabolas (at $e = 1$) and on to hyperbolas. With three more numbers we can move the focus to any point and rotate the curve through any angle. Conics are determined by five numbers.

### Exercises

#### Read-through questions

1. In 7-14 sketch the curve and check for $x$, $y$, and $r$ symmetry.
2. Convert to $xy$ coordinates to draw and identify these curves.
3. In 17-23 find the points where the two curves meet.
4. Warning: You might set $2 \cos \theta = 1 + \cos \theta$ to find $\cos \theta = 1$.

#### 9.2 Exercises

7. $r^2 = 4 \cos 2\theta$ (lemniscate)
8. $r^2 = 4 \sin 2\theta$ (lemniscate)
9. $r = \cos 3\theta$ (three petals)
10. $r^2 = 10 + 6 \cos 4\theta$
11. $r = e^\theta$ (logarithmic spiral)
12. $r = 1/\theta$ (hyperbolic spiral)
13. $r = \tan 9$
14. $r = 1 - 2 \sin 3\theta$ (rose inside rose)
15. Convert $r = 6 \sin \theta + 8 \cos \theta$ to the $xy$ equation of a circle (what radius, what center?).
16. Squaring and adding in the Mars–Earth equation gives $x_{ME}^2 + y_{ME}^2 = 13 - 12 \cos \theta$. The graph of $r^2 = 13 - 12 \cos \theta$ is not at all like Figure 9.4d. What went wrong?

17. $r = 2 \cos \theta$ and $r = 1 + \cos \theta$

### Reading and Drawing

- **Section 3.5** they were $a$ and $b$. (The ellipse was $x^2/a^2 + y^2/b^2 = 1$.) Using $A$ and $e$ we go smoothly from ellipses through parabolas (at $e = 1$) and on to hyperbolas. With three more numbers we can move the focus to any point and rotate the curve through any angle. Conics are determined by five numbers.
18 \( r^2 = \sin 2\theta \) and \( r^2 = \cos 2\theta \)
19 \( r = 1 + \cos \theta \) and \( r = 1 - \sin \theta \)
20 \( r = 1 + \cos \theta \) and \( r = 1 - \cos \theta \)
21 \( r = 2 \) and \( r = 4 \sin 2\theta \)
22 \( r^2 = 4 \cos \theta \) and \( r = 1 - \cos \theta \)
23 \( r \sin \theta = 1 \) and \( r \cos (\theta - \pi/4) = \sqrt{2} \) (straight lines)
24 When is there a dimple in \( r = 1 + b \cos \theta \)? From \( x = (1 + b \cos \theta) \cos \theta \) find \( dx/d\theta \) and \( d^2x/d\theta^2 \) at \( \theta = \pi \). When that second derivative is negative the limaçon has a dimple.
25 How many petals for \( r = \cos 5\theta \)? For \( r = \cos \theta \) there was one, for \( r = \cos 2\theta \) there were four.
26 Explain why \( r = \cos 100 \theta \) has 200 petals but \( r = \cos 101 \theta \) only has 101. The other 101 petals are __________. What about \( r = \cos \frac{1}{3}\theta \)?
27 Find an \( xy \) equation for the cardioid \( r = 1 + \cos \theta \).
28 (a) The flower \( r = \cos 2\theta \) is symmetric across the \( x \) and \( y \) axes. Does that make it symmetric about the origin? (Do two symmetries imply the third, so \( -r = \cos 2\theta \) produces the same curve?)(b) How can \( r = 1, \theta = \pi/2 \) lie on the curve but fail to satisfy the equation?
29 Find an \( xy \) equation for the flower \( r = \cos 2\theta \).
30 Find equations for curves with these properties:
(a) Symmetric about the origin but not the \( x \) axis
(b) Symmetric across the 45° line but not symmetric in \( x \) or \( y \) or \( r \)
(c) Symmetric in \( x \) and \( y \) and \( r \) (like the flower) but changed when \( x \leftrightarrow y \) (not symmetric across the 45° line).

Problems 31–37 are about conic sections—especially ellipses.
31 Find the top point of the ellipse in Figure 9.5a, by maximizing \( y = r \sin \theta = \sin \theta/(1 + \frac{1}{4} \cos \theta) \).
32 (a) Show that all conics \( r = 1/(1 + e \cos \theta) \) go through \( x = 0, y = 1 \).
(b) Find the second focus of the ellipse and hyperbola. For the parabola \( (e = 1) \) where is the second focus?
33 The point \( Q \) in Figure 9.5c has \( y = 1 \). By symmetry find \( x \) and then \( r \) (negative!). Check that \( x^2 + y^2 = r^2 \) and \( |QF| = 2|Qd| \).
34 The equations \( r = A/(1 + e \cos \theta) \) and \( r = 1/(C + D \cos \theta) \) are the same if \( C = \) _______ and \( D = \) _______. For the mirror image across the \( y \) axis replace \( \theta \) by _______. This gives \( r = 1/(C - D \cos \theta) \) as in Figure 12.10 for a planet around the sun.
35 The ellipse \( r = A/(1 + e \cos \theta) \) has length \( 2a \) on the \( x \) axis. Add \( r \) at \( \theta = 0 \) to \( r \) at \( \theta = \pi \) to prove that \( a = A/(1 - e^2) \). The Earth’s orbit has \( a = 92,600,000 \) miles = one astronomical unit (AU).
36 The maximum height \( h \) occurs when \( y = r \sin \theta = A \sin \theta/(1 + e \cos \theta) \) has \( dy/d\theta = 0 \). Show that \( b = y_{\text{max}} = A\sqrt{1 - e^2} \).
37 Combine \( a \) and \( b \) from Problems 35–36 to find \( c = \sqrt{a^2 - b^2} = A e(1 - e^2) \). Then the eccentricity \( e \) is \( c/a \). Halley’s comet is an ellipse with \( a = 18.1 \) AU and \( b = 4.6 \) AU so \( e = \) _______.

Comets have large eccentricity, planets have much smaller \( e \): Mercury .21, Venus .01, Earth .02, Mars .09, Jupiter .05, Saturn .05, Uranus .01, Neptune .25, Pluto 1.25, Kohoutek .9999.
38 If you have a computer with software to do polar graphs, start with these:
1. Flowers \( r = A + \cos n\theta \) for \( n = 1, 3, 7, 8; A = 0, 1, 2 \)
2. Petals \( r = (\cos m\theta + 4 \cos n\theta)/\cos \theta \), \( (m, n) = (5, 3), (3, 5), (9, 1), (2, 3) \)
3. Logarithmic spiral: \( r = e^{\theta/2} \)
4. Nephroid \( r = 1 + 2 \sin \frac{1}{2} \theta \) from the bottom of a teacup
5. Dr. Fay’s butterfly \( r = e^{\cos \theta} - 2 \cos 4\theta + \sin^2(\theta/12) \)
Then create and name your own curve.

### 9.3 Slope, Length, and Area for Polar Curves

The previous sections introduced polar coordinates and polar equations and polar graphs. There was no calculus! We now tackle the problems of area (integral calculus) and slope (differential calculus), when the equation is \( r = F(\theta) \). The use of \( F \) instead of \( f \) is a reminder that the slope is \( \frac{dF}{d\theta} \) and the area is \( \frac{1}{2} \int F(\theta) d\theta \).

Start with area. The region is always divided into small pieces—what is their shape? Look between the angles \( \theta \) and \( \theta + \Delta \theta \) in Figure 9.6a. Inside the curve is a narrow wedge—almost a triangle, with \( \Delta \theta \) as its small angle. If the radius is constant...
the wedge is a sector of a circle. It is a piece of pie cut at the extremely narrow angle \(\Delta \theta\). The area of that piece is a fraction (the angle \(\Delta \theta\) divided by the whole angle \(2\pi\)) of the whole area \(\pi r^2\) of the circle:

\[
\text{area of wedge} = \frac{\Delta \theta}{2\pi} \pi r^2 = \frac{1}{2} r^2 \Delta \theta = \frac{1}{2} [F(\theta)]^2 \Delta \theta.
\]  

(1)

We admit that the exact shape is not circular. The true radius \(F(\theta)\) varies with \(\theta\)—but in a narrow angle that variation is small. When we add up the wedges and let \(\Delta \theta\) approach zero, the area becomes an integral.

**EXAMPLE 1** Find the area inside the circle \(r = \cos \theta\) of radius \(\frac{1}{2}\) (Figure 9.6).

\[
\text{area} = \int_0^{2\pi} \frac{1}{2} r^2 d\theta = \int_0^{2\pi} \frac{1}{2} [\cos \theta]^2 d\theta = \frac{2\pi}{4} = \frac{\pi}{2}.
\]

*That is wrong!* The correct area of a circle of radius \(\frac{1}{2}\) is \(\pi/4\). The mistake is that we went *twice* around the circle as \(\theta\) increased to \(2\pi\). Integrating from \(0\) to \(\pi\) gives \(\pi/4\).

**EXAMPLE 2** Find the area between the circles \(r = \cos \theta\) and \(r = \frac{1}{2}\).

The circles cross at the points where \(r = \cos \theta\) agrees with \(r = \frac{1}{2}\). Figure 9.6 shows these points at \(\pm 60^\circ\), or \(\theta = \pm \pi/3\). Those are the limits of integration, where \(\cos \theta = \frac{1}{2}\). The integral adds up the difference between two wedges, one out to \(r = \cos \theta\) and a smaller one with \(r = \frac{1}{2}\):

\[
\text{area} = \int_{-\pi/3}^{\pi/3} \frac{1}{2} [(\cos \theta)^2 - \left(\frac{1}{2}\right)^2] d\theta.
\]

(3)

*Note that chopped wedges have area \(\frac{1}{2}(F_1^2 - F_2^2)\Delta \theta\) and not \(\frac{1}{2}(F_1 - F_2)^2\Delta \theta\).*

**EXAMPLE 3** Find the area between the cardioid \(r = 1 + \cos \theta\) and the circle \(r = 1\).

\[
\text{area} = \int_{-\pi/2}^{\pi/2} \frac{1}{2} [(1 + \cos \theta)^2 - 1^2] d\theta \quad \text{limits } \theta = \pm \frac{\pi}{2} \text{ where } 1 + \cos \theta = 1
\]
SLOPE OF A POLAR CURVE

Where is the highest point on the cardioid \( r = 1 + \cos \theta \)? What is the slope at \( \theta = \pi/4 \)? Those are not the most important questions in calculus, but still we should know how to answer them. I will describe the method quickly, by switching to rectangular coordinates:

\[
\begin{align*}
\quad & x = r \cos \theta = (1 + \cos \theta) \cos \theta \\
\text{and} & \quad y = r \sin \theta = (1 + \cos \theta) \sin \theta.
\end{align*}
\]

For the highest point, maximize \( y \) by setting its derivative to zero:

\[
\frac{dy}{d\theta} = (1 + \cos \theta)(\cos \theta) + (-\sin \theta)(\sin \theta) = 0. \tag{3}
\]

Thus \( \cos \theta + \cos 2\theta = 0 \), which happens at \( 60^\circ \). The height is \( y = (1 + \frac{1}{2})(\sqrt{3}/2) \).

For the slope, use the chain rule \( \frac{dy}{dx} = \frac{dy/d\theta}{dx/d\theta} \):

\[
\frac{dy}{dx} = \frac{dy/d\theta}{dx/d\theta} = \frac{(1 + \cos \theta)(\cos \theta) + (-\sin \theta)(\sin \theta)}{(1 + \cos \theta)(-\sin \theta) + (-\sin \theta) \cos \theta}. \tag{4}
\]

Equations (3) and (4) avoid the awkward (or impossible) step of eliminating \( \theta \). Instead of trying to find \( y \) as a function of \( x \), we keep \( x \) and \( y \) as functions of \( \theta \). At \( \theta = \pi/4 \), the ratio in equation (4) yields \( \frac{dy}{dx} = -\frac{141}{\sqrt{3}/2} \).

Problem 18 finds a general formula for the slope, using \( \frac{dy}{dx} = (dy/d\theta)/(dx/d\theta) \). Problem 20 finds a more elegant formula, by looking at the question differently.

LENGTH OF A POLAR CURVE

The length integral always starts with \( ds = \sqrt{(dx)^2 + (dy)^2} \). A polar curve has \( x = r \cos \theta = F(\theta) \cos \theta \) and \( y = F(\theta) \sin \theta \). Now take derivatives by the product rule:

\[
\begin{align*}
\quad & dx = (F'(\theta) \cos \theta - F(\theta) \sin \theta) d\theta \\
\text{and} & \quad dy = (F'(\theta) \sin \theta + F(\theta) \cos \theta) d\theta.
\end{align*}
\]

Squaring and adding (note \( \cos^2 \theta + \sin^2 \theta \)) gives the element of length \( ds \):

\[
ds = \sqrt{[F'(\theta)]^2 + [F(\theta)]^2} \, d\theta. \tag{5}
\]

The figure shows \( (ds)^2 = (dr)^2 + (r d\theta)^2 \), the same formula with different letters. The total arc length is \( \int ds \).

The area of a surface of revolution is \( \int 2\pi y \, ds \) (around the \( x \) axis) or \( \int 2\pi x \, ds \) (around the \( y \) axis). Write \( x \), \( y \), and \( ds \) in terms of \( \theta \) and \( d\theta \). Then integrate.

**EXAMPLE 4** The circle \( r = \cos \theta \) has \( ds = \sqrt{1} \, d\theta \). So its length is \( \pi \) (not \( 2\pi \!! \)—don't go around twice). Revolved around the \( y \) axis the circle yields a doughnut with no hole. Since \( x = r \cos \theta = \cos^2 \theta \), the surface area of the doughnut is

\[
\int 2\pi x \, ds = \int_0^\pi 2\pi \cos^2 \theta \, d\theta = \pi^2.
\]

**EXAMPLE 5** The length of \( r = 1 + \cos \theta \) is, by symmetry, double the integral from 0 to \( \pi \):

\[
\text{length of cardioid} = 2 \int_0^\pi \sqrt{(-\sin \theta)^2 + (1 + \cos \theta)^2} \, d\theta
\]

\[
= 2 \int_0^\pi \sqrt{2 + 2 \cos \theta} \, d\theta = 4 \int_0^\pi \cos \frac{\theta}{2} \, d\theta = 8.
\]
9.3 Slope, Length, and Area for Polar Curves

We substituted $4\cos^2\frac{1}{2}\theta$ for $2 + 2\cos \theta$ in the square root. It is possible to skip symmetry and integrate from 0 to $2\pi$—but that needs the absolute value $|\cos \frac{1}{2}\theta|$ to maintain a positive square root.

**Example 6** The logarithmic spiral $r = e^{-\theta}$ has $ds = \sqrt{e^{-2\theta} + e^{-2\theta}} d\theta$. It spirals to zero as $\theta$ goes to infinity, and the total length is finite:

\[ \int ds = \int_0^\infty \sqrt{2} e^{-\theta} d\theta = -\sqrt{2} e^{-\theta} \bigg|_0^\infty = \sqrt{2}. \]

Revolve this spiral for a mathematical seashell with area $\int_0^\infty (2\pi e^{-\theta} \cos \theta) \sqrt{2} e^{-\theta} d\theta$.

**Read-through exercises**

A circular wedge with angle $\Delta \theta$ is a fraction $\frac{\Delta \theta}{2\pi}$ of a whole circle. If the radius is $r$, the wedge area is $\frac{\Delta \theta}{2\pi} \pi r^2$. Then the area inside $r = F(\theta)$ is $\int F(\theta)^2 d\theta$. The area inside $r = 0^2$ from 0 to $\pi$ is $\frac{\pi}{2}$. That spiral meets the circle $r = 1$ at $\theta = 0$. The area inside the circle and outside the spiral is $\int_0^\pi \pi \frac{1}{2} \cos \theta - \frac{1}{2} r^2 d\theta$.

In 1–6 draw the curve and find the area inside.

1. $r = 1 + \cos \theta$
2. $r = \sin \theta + \cos \theta$ from 0 to $\pi$
3. $r = 2 + \cos \theta$
4. $r = 1 + 2 \cos \theta$ (inner loop only)
5. $r = \cos 2\theta$ (one petal only)
6. $r = \cos 3\theta$ (one petal only)

Find the area between the curves in 7–12 after locating their intersections (draw them first).

7. Circle $r = \cos \theta$ and circle $r = \sin \theta$
8. Spiral $r = \theta$ and $y$ axis (first arch)
9. Outside cardioid $r = 1 + \cos \theta$ inside circle $r = 3 \cos \theta$
10. Lemniscate $r^2 = 4 \cos 2\theta$ outside $r = \sqrt{2}$
11. Circle $r = 8 \cos \theta$ beyond line $r \cos \theta = 4$
12. Circle $r = 10$ beyond line $r \cos \theta = 6$
13. Locate the mistake and find the correct area of the lemniscate $r^2 = \cos 2\theta$: area $= \int_0^{\frac{1}{2}\pi} \frac{1}{2} r^2 d\theta = \int_0^{\frac{1}{2}\pi} \frac{1}{4} \cos 2\theta d\theta = 0$.
14. Find the area between the two circles in Example 2.
15. Compute the area between the cardioid and circle in Example 3.
16. Find the complete area (carefully) between the spiral $r = e^{-\theta}$ and the origin.
17. At what $\theta$'s does the cardioid $r = 1 + \cos \theta$ have infinite slope? Which points are furthest to the left (minimum $x$)?
18. Apply the chain rule $\frac{dy}{dx} = \frac{\frac{dy}{d\theta}}{\frac{dx}{d\theta}}$ to $x = F(\theta) \cos \theta$, $y = F(\theta) \sin \theta$. Simplify to reach $\frac{dy}{dx} = \frac{F + \tan \theta}{-F + \tan \theta}$.
19. The groove in a record is nearly a spiral $r = \theta$: length $= \int \sqrt{r^2 + (dr/d\theta)^2} d\theta = \int_0^{14} \sqrt{2} r^2 + c^2 dr$.

Take $c = 0.002$ to give 636 turns between the outer radius 14 cm and the inner radius 6 cm ($14 - 6$ equals $0.002(636)2\pi$).

(a) Omit $c^2$ and just integrate $dr/c$.

(b) Compute the length integral. Tables and calculators allowed. You will never trust integrals again.

20. Show that the angle $\psi$ between the ray from the origin and the tangent line has $\tan \psi = F/dF/d\theta$.

Hint: If the tangent line is at an angle $\phi$ with the horizontal, then $\tan \phi$ is the slope $dy/dx$ in Problem 18. Therefore $\tan \psi = \tan(\phi - \theta) = \frac{\tan \phi - \tan \theta}{1 + \tan \phi \tan \theta}$.

Substitute for $\tan \phi$ and simplify like mad.

21. The circle $r = F(\theta) = 4\sin \theta$ has $\psi = 3\theta$. Draw a figure including $\theta$, $\phi$, $\psi$ and check $\tan \psi$.

22. Draw the cardioid $r = 1 - \cos \theta$, noticing the minus sign. Include the angles $\theta$, $\phi$, $\psi$ and show that $\psi = \theta/2$.

23. The first limaçon in Figure 9.4 looks like a circle centered at $(\frac{1}{2}, 0)$. Prove that it isn't.

24. Find the equation of the tangent line to the circle $r = \cos \theta$ at $\theta = \pi/6$. 
In 25–28 compute the length of the curve.

25 \( r = \theta \) (\( \theta \) from 0 to \( 2\pi \))

26 \( r = \sec \theta \) (\( \theta \) from 0 to \( \pi/4 \))

27 \( r = \sin^3(\theta/3) \) (\( \theta \) from 0 to 3\( \pi \))

28 \( r = \theta^2 \) (\( \theta \) from 0 to \( \pi \))

29 The narrow wedge in Figure 9.6 is almost a triangle. It was treated as a circular sector but triangles are more familiar. Why is the area approximately \( \frac{1}{2}r^2\theta \)?

30 In Example 4 revolve the circle around the x axis and find the surface area. We really only revolve a semicircle.

31 Compute the seashell area \( \int_0^{\pi} e^{-2\cos \theta} d\theta \) using two integrations by parts.

32 Find the surface area when the cardioid \( r = 1 + \cos \theta \) is revolved around the x axis.

33 Find the surface area when the lemniscate \( r^2 = \cos 2\theta \) is revolved around the x axis. What is \( \theta \) after one petal?

34 When \( y = f(x) \) is revolved around the x axis, the volume is \( \int \pi y^2 dx \). When the circle \( r = \cos \theta \) is revolved, switch to a 0-\( \theta \)-integral from 0 to \( \pi/2 \) and check the volume of a sphere.

35 Find the volume when the cardioid \( r = 1 + cos \theta \) is rotated around the x axis.

36 Find the surface area and volume when the graph of \( r = 1/\cos \theta \) is rotated around the y axis (0 < \( \theta \) < \( \pi/4 \)).

37 Show that the spirals \( r = \theta \) and \( r = 1/\theta \) are perpendicular when they meet at \( \theta = 1 \).

38 Draw three circles of radius 1 that touch each other and find the area of the curved triangle between them.

39 Draw the unit square 0 < \( x \) < 1, 0 < \( y \) < 1. In polar coordinates its right side is \( r = \_ \). Find the area from \( \int \frac{1}{2}r^2 d\theta \).

40 (Unravel the paradox) The area of the ellipse \( x = 4 \cos \theta, y = 3 \sin \theta \) is \( \pi \cdot 4 \cdot 3 = 12\pi \). But the integral of \( \frac{1}{2}r^2 d\theta \) is

\[
\int_0^{2\pi} \frac{1}{2} (16 \cos^2 \theta - 9 \sin^2 \theta) d\theta = 12 \frac{1}{2} \pi.
\]

9.4 Complex Numbers

Real numbers are sufficient for most of calculus. Starting from \( x^2 + 4 \), its integral \( \frac{1}{2}x^3 + 4x + C \) is also real. If we are given \( x^3 - 1 \), its derivative \( 3x^2 \) is real. But the roots (or zeros) of those polynomials are complex numbers:

\( x^2 - 4 = 0 \) and \( x^3 - 1 = 0 \) have complex solutions.

We expect two square roots of \( -4 \). There are three cube roots of \( 1 \). Complex numbers are unavoidable, in order to find \( n \) roots for each polynomial of degree \( n \).

This section explains how to work with complex numbers. You will see their relation to polar coordinates. At the end, we use them to solve differential equations.

Start with the imaginary number \( i \). Everybody knows that \( x^2 = -1 \) has no real solution. When you square a real number, the result is never negative. So the world has agreed on a solution called \( i \). (Except that electrical engineers call it \( j \).) Imaginary numbers follow the normal rules of addition, subtraction, multiplication, and division, with one difference: Whenever \( i^2 \) appears it is replaced by \( -1 \). In particular, \( -i \) times \( -i \) gives \( +i^2 = -1 \). In other words, \( -i \) is also a square root of \( -1 \). There are two solutions (\( i \) and \( -i \)) to the equation \( x^2 + 1 = 0 \).

Finding cube roots of \( 1 \) will stretch us further. We need complex numbers—real plus imaginary.

A complex number (say \( 1 + 3i \)) is the sum of a real number (1) and a pure imaginary number (3i). Addition keeps those parts separate; multiplication uses \( i^2 = -1 \):

Addition: \( (1 + 3i) + (1 + 3i) = 1 + 1 + i(3 + 3) = 2 + 6i \)

Multiplication: \( (1 + 3i)(1 + 3i) = 1 + 3i + 3i + 9i^2 = -8 + 6i \).
9.4 Complex Numbers

Adding $1 + 3i$ to $5 - i$ is easy $(6 + 2i)$. Multiplying is longer, but you see the rules:

$$(1 + 3i)(5 - i) = 5 + 15i - i - 3i^2 = 8 + 14i.$$ 

The point is this: We don't have to imagine any more new numbers. After accepting $i$, the rest is straightforward. A real number is just a complex number with no imaginary part! When $1 + 3i$ combines with its "complex conjugate" $1 - 3i$—adding or multiplying—the answer is real:

$$(1 + 3i) + (1 - 3i) = 2 \quad (\text{real})$$

$$(1 + 3i)(1 - 3i) = 1 - 3i + 3i - 9i^2 = 10. \quad (\text{real})$$

The complex conjugate offers a way to do division, by making the denominator real:

$$\frac{1}{1 + 3i} = \frac{1 - 3i}{1 + 3i(1 - 3i)} = \frac{1 - 3i}{10} \quad \text{and} \quad \frac{1}{x + iy} = \frac{x - iy}{x^2 + y^2}.$$ 

The complex number $x + iy$ has real part $x$ and imaginary part $y$. Its complex conjugate is $x - iy$. The product $(x + iy)(x - iy)$ equals $x^2 + y^2 = r^2$. The absolute value (or modulus) is $r = |x + iy| = \sqrt{x^2 + y^2}$.

**The Complex Plane**

Complex numbers correspond to points in a plane. The number $1 + 3i$ corresponds to the point $(1, 3)$. Similarly $x + iy$ is paired with $(x, y)$—which is $x$ units along the "real axis" and $y$ units up the "imaginary axis." The ordinary plane turns into the complex plane. The absolute value $r$ is the same as the polar coordinate $r$ (Figure 9.8a).

The figure shows two more copies of the complex plane. The one in the middle is for addition and subtraction. It uses rectangular coordinates. The one on the right is for multiplication and division and squaring. It uses polar coordinates. In squaring a complex number, $r$ is squared and $\theta$ is doubled—as the right figure and equation (3) both show.

Fig. 9.8 The complex plane shows $x, y, r, \theta$. Add with $x$ and $y$, multiply with $r$ and $\theta$.

Adding complex numbers is like adding vectors (Chapter 11). The real parts give $3 - 1$ and the imaginary parts give $1 + 1$. The vector sum $(2, 2)$ corresponds to the complex sum $2 + 2i$. The complex conjugate $3 - i$ is the mirror image across the real axis ($i$ reversed to $-i$). The connection to $r$ and $\theta$ is the same as before (you see it in the triangle):

$$x = r \cos \theta \quad \text{and} \quad y = r \sin \theta \quad \text{so that} \quad x + iy = r(\cos \theta + i \sin \theta). \quad (2)$$

In the third figure, $1 + i$ has $r = \sqrt{2}$ and $\theta = \pi/4$. The polar form is $\sqrt{2} \cos \pi/4 + \sqrt{2}i \sin \pi/4$. When this number is squared, its $45^\circ$ angle becomes $90^\circ$. The square is $(1 + i)^2 = 1 + 2i - 1 = 2i$. Its polar form is $2 \cos \pi/2 + 2i \sin \pi/2$. 
Multiplication adds angles, division subtracts angles, and squaring doubles angles. The absolute values are multiplied, divided, and squared:

\[(r \cos \theta + i r \sin \theta)^2 = r^2 \cos 2\theta + i r^2 \sin 2\theta.\]  (3)

For \(n\)th powers we reach \(r^n\) and \(n\theta\). For square roots, \(r\) goes to \(\sqrt{r}\) and \(\theta\) goes to \(\frac{1}{2}\theta\). The number \(-1\) is at 180°, so its square root \(i\) is at 90°.

To see why \(\theta\) is doubled in equation (3), factor out \(r^2\) and multiply as usual:

\[(\cos \theta + i \sin \theta)(\cos \theta + i \sin \theta) = \cos^2 \theta - \sin^2 \theta + 2i \sin \theta \cos \theta.\]

The right side is \(\cos 2\theta + i \sin 2\theta\). The double-angle formulas from trigonometry match the squaring of complex numbers. The cube would be \(\cos 3\theta + i \sin 3\theta\) (because \(2\theta\) and \(\theta\) add to \(3\theta\), and \(r\) is still 1). The \(n\)th power is in de Moivre's formula:

\[(\cos \theta + i \sin \theta)^n = \cos n\theta + i \sin n\theta.\]  (4)

With \(n = -1\) we get \(\cos(-\theta) + i \sin(-\theta)\)—which is \(\cos \theta - i \sin \theta\), the complex conjugate:

\[
\frac{1}{\cos \theta + i \sin \theta} = \frac{1}{\cos \theta + i \sin \theta} = \frac{\cos \theta - i \sin \theta}{\cos \theta - i \sin \theta} = \cos \theta - i \sin \theta.
\]  (5)

We are almost touching Euler's formula, the key to all numbers on the unit circle:

\[\text{Euler's formula: } \cos \theta + i \sin \theta = e^{i\theta}.\]  (6)

Squaring both sides gives \((e^{i\theta})(e^{i\theta}) = e^{2i\theta}\). That is equation (3). The \(-1\) power is \(1/e^{i\theta} = e^{-i\theta}\). That is equation (5). Multiplying any \(e^{i\theta}\) by \(e^{i\phi}\) produces \(e^{i(\theta + \phi)}\). The special case \(\phi = \theta\) gives the square, and the special case \(\phi = -\theta\) gives \(e^{i\theta}e^{-i\theta} = 1\).

Euler's formula appeared in Section 6.7, by changing \(x\) to \(i\theta\) in the series for \(e^x\):

\[e^x = 1 + x + \frac{x^2}{2} + \frac{x^3}{6} + \cdots \text{ becomes } e^{i\theta} = 1 + i\theta - \frac{\theta^2}{2} - i\frac{\theta^3}{6} + \cdots.\]

A highlight of Chapter 10 is to recognize two new series on the right. The real terms \(1 - \frac{1}{2}\theta^2 + \cdots\) add up to \(\cos \theta\). The imaginary part \(-\frac{1}{2}\theta^3 + \cdots\) adds up to \(\sin \theta\). Therefore \(e^{i\theta}\) equals \(\cos \theta + i \sin \theta\). It is fantastic that the most important periodic functions in all of mathematics come together in this graceful way.

We learn from Euler (pronounced oiler) that \(e^{2\pi i} = 1\). The cosine of \(2\pi\) is 1, the sine is zero. If you substitute \(x = 2\pi i\) into the infinite series, somehow everything cancels except the 1—which is almost a miracle. From the viewpoint of angles, \(\theta = 2\pi\) carries us around a full circle and back to \(e^{2\pi i} = 1\).

Multiplying Euler's formula by \(r\), we have a third way to write a complex number:

\[\text{Every complex number is } x + iy = r \cos \theta + ir \sin \theta = re^{i\theta}.\]  (7)

**EXAMPLE 1**  
\(2e^{i\theta}\) times \(3e^{i\theta}\) equals \(6e^{2i\theta}\). For \(\theta = \pi/2\), \(2i\) times \(3i\) is \(-6\).

**EXAMPLE 2**  
Find \(w^3\) and \(w^4\) and \(w^8\) and \(w^{25}\) when \(w = e^{i\pi/4}\).

Solution \(e^{i\pi/4}\) is \(\frac{1}{\sqrt{2}} + i\frac{1}{\sqrt{2}}\). Note that \(r^2 = \frac{1}{2} + \frac{1}{2} = 1\). Now watch angles:

\[
w^2 = e^{i\pi/2} = i \quad w^4 = e^{i\pi} = -1 \quad w^8 = 1 \quad w^{25} = w^8w^8w^8w = w.
\]

Figure 9.9 shows the eight powers of \(w\). They are the eighth roots of 1.
**9.4 Complex Numbers**

**Fig. 9.9** The eight powers of $w$ and the cube roots of 1.

**EXAMPLE 3** $(x^2 + 4 = 0)$ The square roots of $-4$ are $2i$ and $-2i$. Instead of $(i)(i) = -1$ we have $(2i)(2i) = -4$. If Euler insists, we write $2i$ and $-2i$ as $2e^{i\pi/2}$ and $2e^{3\pi/2}$.

**EXAMPLE 4** (The cube roots of 1) In rectangular coordinates we have to solve $(x + iy)^3 = 1$, which is not easy. In polar coordinates this same equation is $r^3 e^{3i\theta} = 1$. Immediately $r = 1$. The angle $\theta$ can be $2\pi/3$ or $4\pi/3$ or $6\pi/3$—the cube roots in the figure are evenly spaced:

$$(e^{2\pi i/3})^3 = e^{2\pi i} = 1 \quad (e^{4\pi i/3})^3 = e^{4\pi i} = 1 \quad (e^{6\pi i/3})^3 = e^{6\pi} = 1.$$ 

You see why the angle $8\pi/3$ gives nothing new. It completes a full circle back to $2\pi/3$.

The $n$th roots of 1 are $e^{2\pi i/n}, e^{4\pi i/n}, \ldots, 1$. There are $n$ of them. They lie at angles $2\pi/n, 4\pi/n, \ldots, 2\pi$ around the unit circle.

**SOLUTION OF DIFFERENTIAL EQUATIONS**

The algebra of complex numbers is now applied to the calculus of complex functions. The complex number is $c$, the complex function is $e^{ct}$. It will solve the equations $y'' = -4y$ and $y''' = y$, by connecting them to $c^2 = -4$ and $c^3 = 1$. Chapter 16 does the same for all linear differential equations with constant coefficients—this is an optional preview.

Please memorize the one key idea: Substitute $y = e^{ct}$ into the differential equation and solve for $c$. Each derivative brings a factor $c$, so $y' = ce^{ct}$ and $y'' = c^2 e^{ct}$:

$$d^2y/dt^2 = -4y$$

leads to $c^2 e^{ct} = -4e^{ct}$, which gives $c^2 = -4$. \(8\)

For this differential equation, $c$ must be a square root of $-4$. We know the candidates ($c = 2i$ and $c = -2i$). The equation has two “pure exponential solutions” $e^{ct}$:

$$y = e^{2it} \quad \text{and} \quad y = e^{-2it}. \quad 9$$

Their combinations $y = Ae^{2it} + Be^{-2it}$ give all solutions. In Chapter 16 we will choose the two numbers $A$ and $B$ to match two initial conditions at $t = 0$.

The solution $y = e^{2it} = \cos 2t + i \sin 2t$ is complex. The differential equation is real. For real $y$'s, take the real and imaginary parts of the complex solutions:

$$y_{\text{real}} = \cos 2t \quad \text{and} \quad y_{\text{imaginary}} = \sin 2t. \quad 10$$

These are the “pure oscillatory solutions.” When $y = e^{2it}$ travels around the unit circle, its imaginary part $\sin 2t$ moves up and down. (It is like the ball and its shadow in Section 1.4, but twice as fast because of $2t$.) The real part $\cos 2t$ goes backward and forward. By the chain rule, the second derivative of $\cos 2t$ is $-4 \cos 2t$. Thus $d^2y/dt^2 = -4y$ and we have real solutions.
EXAMPLE 5 Find three solutions and then three real solutions to \( \frac{d^3y}{dt^3} = y \).

*Key step: Substitute \( y = e^t \). The result is \( e^t e^t = e^{2t} \). Thus \( e^3 = 1 \) and \( c \) is a cube root of 1. The candidate \( c = 1 \) gives \( y = e^t \) (our first solution). The next \( c \) is complex:

\[
c = e^{2\pi i/3} = -\frac{1}{2} + i\frac{\sqrt{3}}{2} \quad \text{yields} \quad y = e^{2t} = e^{-it/2} e^{i\sqrt{3}t/2}.
\]

The real part of the exponent leads to the absolute value \( |y| = e^{-t/2} \). It decreases as \( t \) gets larger, so \( y \) moves toward zero. At the same time, the factor \( e^{i\sqrt{3}t/2} \) goes around the unit circle. Therefore \( y \) spirals in to zero (Figure 9.10). So does its complex conjugate, which is the third exponential. Changing \( i \) to \(-i\) in (11) gives the third cube root of 1 and the third solution \( e^{-it} e^{-i\sqrt{3}t/2} \).

The first real solution is \( y = e^t \). The others are the two parts of the spiral:

\[
y_{\text{real}} = e^{-t/2} \cos \frac{\sqrt{3}t}{2} \quad \text{and} \quad y_{\text{imaginary}} = e^{-t/2} \sin \frac{\sqrt{3}t}{2}.
\]

That is \( r \cos \theta \) and \( r \sin \theta \). It is the ultimate use (until Chapter 16) of polar coordinates and complex numbers. We might have discovered \( \cos 2t \) and \( \sin 2t \) without help, for \( y'' = -4y \). I don’t think these solutions to \( y'' = y \) would have been found.

EXAMPLE 6 Find four solutions to \( \frac{d^4y}{dt^4} = y \) by substituting \( y = e^t \).

Four derivatives lead to \( c^4 = 1 \). Therefore \( c \) is \( i \) or \(-1 \) or \(-i \) or \( 1 \). The solutions are \( y = e^t, e^{-t}, e^{-it}, \) and \( e^t \). If we want real solutions, \( e^t \) and \( e^{-it} \) combine into \( \cos t \) and \( \sin t \). In all cases \( y''' = y \).

![Fig. 9.10 Solutions move in the complex plane: \( y'' = -4y \) and \( y''' = y \) and \( y'''' = y \).](image)

9.4 EXERCISES

Read-through questions

The complex number \( 3 + 4i \) has real part \( a \) and imaginary part \( b \). Its absolute value is \( r = \sqrt{a^2 + b^2} \) and its complex conjugate is \( c \). Its position in the complex plane is at \( (a, b) \). Its polar form is \( r \cos \theta + ir \sin \theta = r e^{i\theta} \). Its square is \( a^2 + b^2 \) and \( c^2 = \frac{1}{2} e^{ia} \). Its n-th power is \( d \).

The sum of \( 1 + i \) and \( 1 - i \) is \( e \). The product of \( 1 + i \) and \( 1 - i \) is \( f \). In polar form this is \( \sqrt{2} e^{i\pi/4} \) times \( g \). The quotient \( (1 + i)/(1 - i) \) equals the imaginary number \( h \). The number \( (1 + i)^3 \) equals \( i \). An eighth root of \( 1 \) is \( w = o \). The other eighth roots are \( p \).

To solve \( \frac{d^8y}{dt^8} = y \), look for a solution of the form \( y = q \). Substituting and canceling \( e^t \) leads to the equation \( t \). There are \( s \) choices for \( c \), one of which is \( -1 + i/\sqrt{2} \). With that choice \( |e^t| = t \). The real solutions are \( \text{Re} e^t = u \) and \( \text{Im} e^t = v \).

In 1–6 plot each number in the complex plane.

1. \( 2 + i \) and its complex conjugate \( 2 - i \) and their sum and product
2. \( 1 + i \) and its square \( (1 + i)^2 \) and its reciprocal \( 1/(1 + i) \)
3. \(2e^{i\pi/6}\) and its reciprocal \(\frac{1}{e^{-i\pi/6}}\) and their squares.

4. The sixth roots of 1 (six of them)

5. \(\cos \frac{3\pi}{4} + i \sin \frac{3\pi}{4}\) and its square and cube.

6. \(4e^{i\pi/3}\) and its square roots.

7. For complex numbers \(c = x + iy = re^{i\theta}\) and their conjugates \(\bar{c} = x - iy = re^{-i\theta}\), find all possible locations in the complex plane of (1) \(c + \bar{c}\) (2) \(c - \bar{c}\) (3) \(c\bar{c}\) (4) \(c/\bar{c}\).

8. Find \(x\) and \(y\) for the complex numbers \(x + iy\) at angles \(\theta = 45^\circ, 90^\circ, 135^\circ\) on the unit circle. Verify directly that the square of the first is the second and the cube of the first is the third.

9. If \(c = 2 + i\) and \(d = 4 + 3i\), find \(cd\) and \(c/d\). Verify that the absolute value \(|cd|\) equals \(|c|\) times \(|d|\), and \(|c/d|\) equals \(|c|\) divided by \(|d|\).

10. Find a solution \(x\) to \(e^{ix} = i\) and a solution to \(e^{ix} = 1/e\). Then find a second solution.

Find the sum and product of the numbers in 11–14.

11. \(e^{i\theta}\) and \(e^{-i\theta}\), also \(e^{2i\theta/3}\) and \(e^{4i\theta/3}\)

12. \(e^{i\theta}\) and \(e^{3i\theta}\), also \(e^{i4\theta/3}\) and \(e^{-i\theta}\)

13. The sixth roots of 1 (add and multiply all six)

14. The two roots of \(c^2 - 4c + 5 = 0\)

15. If \(c = re^{i\theta}\) is not zero, what are \(c^6\) and \(c^{-1}\) and \(c^{-4}\)?

16. Multiply out \((\cos \theta + i \sin \theta)^3 = e^{3i\theta}\), to find the real part \(\cos 3\theta\) and the imaginary part \(\sin 3\theta\) in terms of \(\cos \theta\) and \(\sin \theta\).

17. Plot the three cube roots of a typical number \(re^{i\theta}\). Show why they add to zero. One cube root is \(r^{1/3}e^{i\theta/3}\).

18. Prove that the four fourth roots of \(re^{i\theta}\) multiply to give \(-re^{i\theta}\).

In 19–22, find all solutions of the form \(y = e^c\).

19. \(y'' + y = 0\)

20. \(y''' + y = 0\)

21. \(y'' - y' = 0\)

22. \(y'' + 6y' + 5y = 0\)

Construct two real solutions from the real and imaginary parts of \(e^c\) (first find \(c\)).

23. \(y' + 49y = 0\)

24. \(y' - 2y' + 2y = 0\)

Sketch the path of \(y = e^t\) as \(t\) increases from zero, and mark \(y = c^t\).

25. \(c = 1 - i\)

26. \(c = -1 + i\)

27. \(c = \pi i/4\)

28. What is the solution of \(dy/dt = iy\) starting from \(y_0 = 1\)? For this solution, matching real parts and imaginary parts of \(dy/dt = iy\) gives \(_______\) and \(_______\).

29. In Figure 9.10b, at what time \(t\) does the spiral cross the real axis at the far left? What does \(y\) equal at that time?

30. Show that \(\cos \theta = \frac{1}{2}(e^{i\theta} + e^{-i\theta})\) and find a similar formula for \(\sin \theta\).

31. True or false, with an example to show why:

(a) If \(c_1 + c_2\) is real, the \(c_1\)'s are complex conjugates.
(b) If \(|c_1| = 2\) and \(|c_2| = 4\) then \(c_1c_2\) has absolute value 8.
(c) If \(|c_1| = 1\) and \(|c_2| = 1\) then \(|c_1 + c_2|\) is (at least 1) (at most 2) (equal to 2).
(d) If \(e^c\) approaches zero as \(t \to \infty\), then \(c\) is negative (the real part of \(c\) is negative) (\(|c|\) is less than 1).

32. The polar form of \(re^{i\theta}\) times \(Re^{i\phi}\) is \(_______\). The rectangular form is \(_______\). Circle the terms that give \(rR\cos(\theta + \phi)\).

33. The complex number \(1/(re^{i\theta})\) has polar form \(_______\) and rectangular form \(_______\) and square roots \(_______\).

34. Show that \(\cos ix = \cosh x\) and \(\sin ix = i \sinh x\). What is the cosine of \(i\)?
CHAPTER 9   Polar Coordinates and Complex Numbers
  9.1 Polar Coordinates 348
  9.2 Polar Equations and Graphs 351
  9.3 Slope, Length, and Area for Polar Curves 356
  9.4 Complex Numbers 360

CHAPTER 10   Infinite Series
  10.1 The Geometric Series 368
  10.2 Convergence Tests: Positive Series 374
  10.3 Convergence Tests: All Series 381
  10.4 The Taylor Series for $e^x$, sin $x$, and cos $x$ 385
  10.5 Power Series 391

CHAPTER 11   Vectors and Matrices
  11.1 Vectors and Dot Products 398
  11.2 Planes and Projections 407
  11.3 Cross Products and Determinants 416
  11.4 Matrices and Linear Equations 425
  11.5 Linear Algebra in Three Dimensions 435

CHAPTER 12   Motion along a Curve
  12.1 The Position Vector 446
  12.2 Plane Motion: Projectiles and Cycloids 453
  12.3 Tangent Vector and Normal Vector 459
  12.4 Polar Coordinates and Planetary Motion 464

CHAPTER 13   Partial Derivatives
  13.1 Surfaces and Level Curves 472
  13.2 Partial Derivatives 475
  13.3 Tangent Planes and Linear Approximations 480
  13.4 Directional Derivatives and Gradients 490
  13.5 The Chain Rule 497
  13.6 Maxima, Minima, and Saddle Points 504
  13.7 Constraints and Lagrange Multipliers 514
Infinite Series

Infinite series can be a pleasure (sometimes). They throw a beautiful light on \( \sin x \) and \( \cos x \). They give famous numbers like \( \pi \) and \( e \). Usually they produce totally unknown functions—which might be good. But on the painful side is the fact that an infinite series has infinitely many terms.

It is not easy to know the sum of those terms. More than that, it is not certain that there is a sum. We need tests, to decide if the series converges. We also need ideas, to discover what the series converges to. Here are examples of convergence, divergence, and oscillation:

The first series converges. Its next term is 1/8, after that is 1/16—and every step brings us halfway to 2. The second series (the sum of 1's) obviously diverges to infinity. The oscillating example (with 1's and -1's) also fails to converge.

All those and more are special cases of one infinite series which is absolutely the most important of all:

\[
1 + \frac{1}{2} + \frac{1}{4} + \cdots = 2 \quad 1 + 1 + 1 + \cdots = \infty \quad 1 - 1 + 1 - 1 \cdots = ?
\]

The first series converges. Its next term is 1/8, after that is 1/16—and every step brings us halfway to 2. The second series (the sum of 1's) obviously diverges to infinity. The oscillating example (with 1's and -1's) also fails to converge.

All those and more are special cases of one infinite series which is absolutely the most important of all:

**The geometric series is**

\[
1 + x + x^2 + x^3 + \cdots = \frac{1}{1 - x}.
\]

This is a series of *functions*. It is a "power series." When we substitute numbers for \( x \), the series on the left may converge to the sum on the right. We need to know when it doesn't. Choose \( x = \frac{1}{2} \) and \( x = 1 \) and \( x = -1 \):

\[
1 + \frac{1}{2} + \left(\frac{1}{2}\right)^2 + \cdots \text{ is the convergent series. Its sum is } \frac{1}{1 - \frac{1}{2}} = 2
\]

\[
1 + 1 + 1 + \cdots \text{ is divergent. Its sum is } \frac{1}{1 - 1} = \frac{1}{0} = \infty
\]

\[
1 + (-1) + (-1)^2 + \cdots \text{ is the oscillating series. Its sum should be } \frac{1}{1 - (-1)} = \frac{1}{2}.
\]

The last sum bounces between one and zero, so at least its average is 1/2. At \( x = 2 \) there is no way that \( 1 + 2 + 4 + 8 + \cdots \) agrees with \( 1/(1 - 2) \).

This behavior is typical of a power series—to converge in an interval of \( x \)'s and
The Geometric Series

to diverge when \( x \) is large. The geometric series is safe for \( x \) between \(-1\) and \(1\). Outside that range it diverges.

The next example shows a repeating decimal 1.111...:

Set \( x = \frac{1}{10} \). The geometric series is \( 1 + \frac{1}{10} + \left( \frac{1}{10} \right)^2 + \left( \frac{1}{10} \right)^3 + \ldots \)

The decimal 1.111... is also the fraction \( \frac{1}{1 - \frac{1}{10}} \), which is \( \frac{10}{9} \). Every fraction leads to a repeating decimal. Every repeating decimal adds up (through the geometric series) to a fraction.

To get 3.333..., just multiply by 3. This is \( \frac{10}{3} \). To get 1.0101..., set \( x = \frac{1}{100} \). This is the fraction \( \frac{1}{1 - \frac{1}{100}} \), which is \( \frac{100}{99} \).

Here is an unusual decimal (which eventually repeats). I don't really understand it:

\[
\frac{1}{243} = .004115226337448...
\]

Most numbers are not fractions (or repeating decimals). A good example is \( \pi \):

\[
\pi = 3 + \frac{1}{10} + \frac{4}{100} + \frac{1}{1000} + \frac{5}{10000} + \ldots
\]

This is 3.1415..., a series that certainly converges. We happen to know the first billion terms (the billionth is given below). Nobody knows the 2 billionth term. Compare that series with this one, which also equals \( \pi \):

\[
\pi = 4 - \frac{4}{3} - \frac{4}{5} - \frac{4}{7} + \ldots
\]

That alternating series is really remarkable. It is typical of this chapter, because its pattern is clear. We know the first billion terms (it has a minus sign). This is not a geometric series, but in Section 10.1 it comes from a geometric series.

**Question** Does this series actually converge? What if all signs are \( + \)?

**Answer** The alternating series converges to \( \pi \) (Section 10.3). The positive series diverges to infinity (Section 10.2). The terms go to zero, but their sum is infinite.

This example begins to show what the chapter is about. Part of the subject deals with special series, adding to \( \frac{10}{9} \) or \( \pi \) or \( e^x \). The other part is about series in general, adding to numbers or functions that nobody has heard of. The situation was the same for integrals—they give famous answers like \( \ln x \) or unknown answers like \( \int x^3 \, dx \). The sum of \( 1 + \frac{1}{8} + \frac{1}{27} + \ldots \) is also unknown—although a lot of mathematicians have tried.

The chapter is not long, but it is full. The last half studies power series. We begin with a linear approximation like \( 1 + x \). Next is a quadratic approximation like \( 1 + x + x^2 \). In the end we match all the derivatives of \( f(x) \). This is the "Taylor series," a new way to create functions—not by formulas or integrals but by infinite series.

No example can be better than \( \frac{1}{1 - x} \), which dominates Section 10.1. Then we define convergence and test for it. (Most tests are really comparisons with a geometric series.) The second most important series in mathematics is the exponential series \( e^x = 1 + x + \frac{1}{2}x^2 + \frac{1}{6}x^3 + \ldots \). It includes the series for \( \sin x \) and \( \cos x \), because of the formula \( e^{ix} = \cos x + i \sin x \). Finally a whole range of new and old functions will come from Taylor series.

In the end, all the key functions of calculus appear as "infinite polynomials" (except the step function). This is the ultimate voyage from the linear function \( y = mx + b \).
10.1 The Geometric Series

We begin by looking at both sides of the geometric series:

\[ 1 + x + x^2 + x^3 + \cdots = \frac{1}{1-x} \]  

(1)

How does the series on the left produce the function on the right? How does \(1/(1-x)\) produce the series? Add up two terms of the series, then three terms, then \(n\) terms:

\[ 1 + x = \frac{1-x^2}{1-x} \]
\[ 1 + x + x^2 = \frac{1-x^3}{1-x} \]
\[ 1 + \cdots + x^{n-1} = \frac{1-x^n}{1-x} \]  

(2)

For the first, \(1 + x\) times \(1 - x\) equals \(1 - x^2\) by ordinary algebra. The second begins to make the point: \(1 + x + x^2\) times \(1 - x\) gives \(1 - x + x - x^2 + x^2 - x^3\). Between 1 at the start and \(-x^3\) at the end, everything cancels. The same happens in all cases: \(1 + \cdots + x^{n-1}\) times \(1 - x\) leaves 1 at the start and \(-x^n\) at the end. This proves equation (2)—the sum of \(n\) terms of the series.

For the whole series we will push \(n\) towards infinity. On a graph you can see what is happening. Figure 10.1 shows \(n = 1\) and \(n = 2\) and \(n = 3\) and \(n = \infty\).

**Fig. 10.1** Two terms, then three terms, then full series:

\[ 1 + x + x^2 + \cdots = \frac{1}{1-x} \]

The infinite sum gives a finite answer, provided \(x\) is between 
\(-1\) and 1. Then \(x^n\) goes to zero:

\[ \frac{1-x^n}{1-x} \rightarrow \frac{1}{1-x} \]

Now start with the function \(1/(1-x)\). How does it produce the series? One way is elementary but brutal, to do “long division” of \(1-x\) into 1 (next to the figure). Another way is to look up the binomial formula for \((1-x)^{-1}\). That is cheating—we want to discover the series, not just memorize it. The successful approach uses cal-
10.1 The Geometric Series

culus. Compute the derivatives of \( f(x) = \frac{1}{1-x} \):

\[
f'(x) = (1-x)^{-2} \quad f''(x) = 2(1-x)^{-3} \quad f'''(x) = 6(1-x)^{-4} \quad \ldots
\]

At \( x = 0 \) these derivatives are 1, 2, 6, 24, \ldots. Notice how \(-1\) from the chain rule keeps them positive. The \( n \)th derivative at \( x = 0 \) is \( n! \) factorials:

\[
f(0) = 1 \quad f'(0) = 1 \quad f''(0) = 2 \quad f'''(0) = 6 \quad \ldots \quad f^{(n)}(0) = n!
\]

Now comes the idea. To match the series with \( \frac{1}{1-x} \), match all those derivatives at \( x = 0 \). Each power \( x^n \) gets one derivative right. Its derivatives at \( x = 0 \) are zero, except the \( n \)th derivative, which is \( n! \). By adding all powers we get every derivative right—so the geometric series matches the function:

\[
1 + x + x^2 + x^3 + \ldots \text{ has the same derivatives at } x = 0 \text{ as } \frac{1}{1-x}.
\]

The linear approximation is \( 1 + x \). Then comes \( \frac{1}{2} f''(0)x^2 = x^2 \). The third derivative is supposed to be 6, and \( x^3 \) is just what we need. Through its derivatives, the function produces the series.

With that example, you have seen a part of this subject. The geometric series diverges if \(|1-x| \geq 1\). Otherwise it adds up to the function it comes from (when \(-1 < x < 1\)). To get familiar with other series, we now apply algebra or calculus—to reach the square of \( \frac{1}{1-x} \) or its derivative or its integral. The point is that these operations are applied to the series.

The best I know is to show you eight operations that produce something useful. At the end we discover series for \( \ln 2 \) and \( \pi \).

1. **Multiply the geometric series by a or ax:**

\[
a + ax + ax^2 + \ldots = \frac{a}{1-x} \quad ax + ax^2 + ax^3 + \ldots = \frac{ax}{1-x}.
\]

The first series fits the decimal 3.333\ldots. In that case \( a = 3 \). The geometric series for \( x = \frac{1}{10} \) gave 1.111\ldots = 10/9, and this series is just three times larger. Its sum is 10/3.

The second series fits other decimals that are fractions in disguise. To get 12/99, choose \( a = 12 \) and \( x = 1/100 \):

\[
.121212\ldots = \frac{12}{100} + \frac{12}{100^2} + \frac{12}{100^3} + \ldots = \frac{12/100}{1-1/100} = \frac{12}{99}.
\]

Problem 13 asks about \( .8787 \ldots \) and \( .123123\ldots \). It is usual in precalculus to write \( a + ar + ar^2 + \ldots = a/(1-r) \). But we use \( x \) instead of \( r \) to emphasize that this is a function—which we can now differentiate.

2. **The derivative of the geometric series \( 1 + x + x^2 + \ldots \) is \( 1/(1-x)^2 \):**

\[
1 + 2x + 3x^2 + 4x^3 + \ldots = \frac{d}{dx}\left( \frac{1}{1-x} \right) = \frac{1}{(1-x)^2}.
\]

At \( x = \frac{1}{10} \) the left side starts with 1.23456789. The right side is \( 1/(1-\frac{1}{10})^2 = 1/(9/10)^2 \), which is 100/81. If you have a calculator, divide 100 by 81.

The answer should also be near \( (1.11111111)^2 \), which is 1.2345678987654321.

3. **Subtract \( 1 + x + x^2 + \ldots \) from \( 1 + 2x + 3x^2 + \ldots \) as you subtract functions:**

\[
x + 2x^2 + 3x^3 + \ldots = \frac{1}{(1-x)^2} - \frac{1}{(1-x)} = \frac{x}{(1-x)^2}.
\]

Curiously, the same series comes from multiplying (5) by \( x \). It answers a question left open in Section 8.4—the average number of coin tosses until the result is heads. This
is the sum \(1(p_1) + 2(p_2) + \cdots\) from probability, with \(x = \frac{1}{2}\):
\[
1(\frac{1}{2}) + 2(\frac{1}{2})^2 + 3(\frac{1}{2})^3 + \cdots = \frac{1}{(1 - \frac{1}{2})^2} = 2. \tag{7}
\]

The probability of waiting until the \(n\)th toss is \(p_n = (\frac{1}{2})^n\). The expected value is two tosses. I suggested experiments, but now this mean value is exact.

4. **Multiply series:** the geometric series times itself is \(1/(1-x)\) squared:
\[
(1 + x + x^2 + \cdots)(1 + x + x^2 + \cdots) = 1 + 2x + 3x^2 + \cdots. \tag{8}
\]
The series on the right is not new! In equation (5) it was the derivative of \(y = 1/(1-x)\). Now it is the square of the same \(y\). The geometric series satisfies \(dy/dx = y^2\), so the function does too. We have stumbled onto a differential equation.

Notice how the series was squared. A typical term in equation (8) is \(3x^2\), coming from 1 times \(x^2\) and \(x\) times \(x^2\) and \(x^2\) times 1 on the left side. It is a lot quicker to square \(1/(1-x)\) but other series can be multiplied when we don’t know what functions they add up to.

5. **Solve \(dy/dx = y^2\) from any starting value—a new application of series:**
Suppose the starting value is \(y = 1\) at \(x = 0\). The equation \(y' = y^2\) gives \(1^2\) for the derivative. Now a key step: The derivative of the equation gives \(y'' = 2yy'\). At \(x = 0\) that is \(2 \cdot 1 \cdot 1\). Continuing upwards, the derivative of \(2yy'\) is \(2yy'' + 2(y')^2\). At \(x = 0\) that is \(y''' = 4 + 2 = 6\).

All derivatives are factorials: 1, 2, 6, 24, …. We are matching the derivatives of the geometric series \(1 + x + x^2 + x^3 + \cdots\). Term by term, we rediscover the solution to \(y' = y^2\). The solution starting from \(y(0) = 1\) is \(y = 1/(1-x)\).

A different starting value is \(-1\). Then \(y' = (-1)^2 = 1\) as before. The chain rule gives \(y'' = 2yy' = -2\) and then \(y''' = 6\). With alternating signs to match these derivatives, the solution starting from \(-1\) is
\[
y = -1 + x - x^2 + x^3 \cdots = -1/(1+x). \tag{9}
\]
It is a small challenge to recognize the function on the right from the series on the left. The series has \(-x\) in place of \(x\); then multiply by \(-1\). The sum \(y = -1/(1+x)\) also satisfies \(y' = y^2\). We can solve differential equations from all starting values by infinite series. Essentially we substitute an unknown series into the equation, and calculate one term at a time.

6. **The integrals of \(1 + x + x^2 + \cdots\) and \(1 - x + x^2 - \cdots\) are logarithms:**
\[
x + \frac{1}{2}x^2 + \frac{1}{3}x^3 + \cdots = \int_0^{\infty} dx = -\ln(1-x) \tag{10a}
\]
\[
x - \frac{1}{2}x^2 + \frac{1}{3}x^3 - \cdots = \int_0^{\infty} \frac{dx}{1+x} = +\ln(1+x) \tag{10b}
\]
The derivative of (10a) brings back the geometric series. For logarithms we find \(1/n\) not \(1/n!\) The first term \(x\) and second term \(\frac{1}{2}x^2\) give linear and quadratic approximations. Now we have the whole series. I cannot fail to substitute 1 and \(\frac{1}{2}\), to find \(\ln(1-1)\) and \(\ln(1+1)\) and \(\ln(1-\frac{1}{2})\):
\[
x = 1: \quad 1 + \frac{1}{2} + \frac{1}{3} + \frac{1}{4} + \cdots = -\ln 0 = +\infty \tag{11a}
\]
\[
x = 1: \quad 1 - \frac{1}{2} + \frac{1}{3} - \frac{1}{4} + \cdots = \ln 2 = .693 \tag{11b}
\]
\[
x = \frac{1}{2}: \quad \frac{1}{2} + \frac{1}{3} + \frac{1}{4} + \frac{1}{6} + \cdots = -\ln \frac{1}{2} = \ln 2. \tag{12}
\]
10.1 The Geometric Series

The first series diverges to infinity. This harmonic series \( 1 + \frac{1}{2} + \frac{1}{3} + \cdots \) came into the earliest discussion of limits (Section 2.6). The second series has alternating signs and converges to \( \ln 2 \). The third has plus signs and also converges to \( \ln 2 \). These will be examples for a major topic in infinite series—tests for convergence.

For the first time in this book we are able to compute a logarithm! Something remarkable is involved. The sums of numbers in (11) and (12) were discovered from the sums of functions in (10). You might think it would be easier to deal only with numbers, to compute \( \ln 2 \). But then we would never have integrated the series for \( 1/(1-x) \) and detected (10). It is better to work with \( x \), and substitute special values like \( \frac{1}{3} \) at the end.

There are two practical problems with these series. For \( \ln 2 \) they converge slowly. For \( \ln e \) they blow up. The correct answer is \( \ln e = 1 \), but the series can’t find it. Both problems are solved by adding (10a) to (10b), which cancels the even powers:

\[
2 \left( x + \frac{x^3}{3} + \frac{x^5}{5} + \cdots \right) = \ln(1 + x) - \ln(1 - x) = \ln \frac{1 + x}{1 - x}.
\]  

At \( x = \frac{1}{3} \), the right side is \( \ln \frac{4}{3} - \ln \frac{2}{3} = \ln 2 \). Powers of \( \frac{1}{3} \) are much smaller than powers of 1 or \( \frac{1}{2} \), so \( \ln 2 \) is quickly computed. All logarithms can be found from the improved series (13).

7. Change variables in the geometric series (replace \( x \) by \( x^2 \) or \(-x^2\)):

\[
1 + x^2 + x^4 + x^6 + \cdots = \frac{1}{1 - x^2}
\]

\[
1 - x^2 + x^4 - x^6 + \cdots = \frac{1}{1 + x^2}.
\]

This produces new functions (always our goal). They involve even powers of \( x \). The second series will soon be used to calculate \( \pi \). Other changes are valuable:

\[
\frac{x}{2} \text{ in place of } x: \quad 1 + \frac{x}{2} + \left(\frac{x}{2}\right)^2 + \cdots = \frac{1}{1 - (x/2)} = \frac{2}{2 - x}
\]

\[
\frac{1}{x} \text{ in place of } x: \quad 1 + \frac{1}{x} + \frac{1}{x^2} + \cdots = \frac{1}{1 - (1/x)} = \frac{x}{x - 1}.
\]

Equation (17) is a series of negative powers \( x^-n \). It converges when \( |x| \) is greater than 1. Convergence in (17) is for large \( x \). Convergence in (16) is for \( |x| < 2 \).

8. The integral of \( 1 - x^2 + x^4 - x^6 + \cdots \) yields the inverse tangent of \( x \):

\[
x - \frac{1}{3}x^3 + \frac{1}{5}x^5 - \frac{1}{7}x^7 + \cdots = \int \frac{dx}{1 + x^2} = \tan^{-1} x.
\]

We integrated (15) and got odd powers. The magical formula for \( \pi \) (discovered by Leibniz) comes when \( x = 1 \). The angle with tangent 1 is \( \pi/4 \):

\[
1 - \frac{1}{3} + \frac{1}{5} - \frac{1}{7} + \cdots = \frac{\pi}{4}.
\]

The first three terms give \( \pi \approx 3.47 \) (not very close). The 5000th term is still of size .0001, so the fourth decimal is still not settled. By changing to \( x = 1/\sqrt{3} \), the astronomer Halley and his assistant found 71 correct digits of \( \pi/6 \) (while waiting for the comet). That is one step in the long and amazing story of calculating \( \pi \). The Chudnovsky brothers recently took the latest step with a supercomputer—they have found more than one billion decimal places of \( \pi \) (see Science, June 1989). The digits look completely random, as everyone expected. But so far we have no proof that all ten digits occur \( \frac{1}{10} \) of the time.
Historical note  Archimedes located π above 3.14 and below 3\(\frac{1}{3}\). Variations of his method (polygons in circles) reached as far as 34 digits—but not for 1800 years. Then Halley found 71 digits of π/6 with equation (18). For faster convergence that series was replaced by other inverse tangents, using smaller values of \(x\):

\[
\frac{\pi}{4} = \tan^{-1} \frac{1}{2} + \tan^{-1} \frac{1}{3} = 4 \tan^{-1} \frac{1}{5} - \tan^{-1} \frac{1}{239}. \tag{20}
\]

A prodigy named Dase, who could multiply 100-digit numbers in his head in 8 hours, finally passed 200 digits of π. The climax of hand calculation came when Shanks published 607 digits. I am sorry to say that only 527 were correct. (With years of calculation he went on to 707 digits, but still only 527 were correct.) The mistake was not noticed until 1945! Then Ferguson reached 808 digits with a desk calculator.

Now comes the computer. Three days on an ENIAC (1949) gave 2000 digits. A hundred minutes on an IBM 704 (1958) gave 10,000 digits. Shanks (no relation) reached 100,000 digits. Finally a million digits were found in a day in 1973, with a CDC 7600. All these calculations used variations of equation (20).

The record after that went between Cray and Hitachi and now IBM. But the method changed. The calculations rely on an incredibly accurate algorithm, based on the “arithmetic-geometric mean iteration” of Gauss. It is also incredibly simple, all things considered:

\[
a_{n+1} = \frac{a_n + b_n}{2}, \quad b_{n+1} = \sqrt{a_n b_n}, \quad \pi_n = 2a_{n+1} \left(1 - \sum_{k=0}^{n} 2^k (a_k^2 - b_k^2)\right).
\]

The number of correct digits more than doubles at every step. By \(n = 9\) we are far beyond Shanks (the hand calculator). No end is in sight. Almost anyone can go past a billion digits, since with the Chudnovsky method we don’t have to start over again.

It is time to stop. You may think (or hope) that nothing more could possibly be done with geometric series. We have gone a long way from \(1/(1-x)\), but some functions can never be reached. One is \(e^x\) (and its relatives \(\sin x, \cos x, \sinh x, \cosh x\)). Another is \(\sqrt{1-x}\) (and its relatives \(1/\sqrt{1-x^2}, \sin^{-1} x, \sec^{-1} x, \ldots\)). The exponentials are in 10.4, with series that converge for all \(x\). The square roots are in 10.5, closer to geometric series and converging for \(|x| < 1\). Before that we have to say what convergence means.

The series came fast, but I hope you see what can be done (subtract, multiply, differentiate, integrate). Addition is easy, division is harder, all are legal. Some unexpected numbers are the sums of infinite series.

Added in proof  By e-mail I just learned that the record for π is back in Japan: 2^60 digits which is more than 1.07 billion. The elapsed time was 100 hours (75 hours of CPU time on an NEC machine). The billionth digit after the decimal point is 9.

---

10.1 Exercises

Read-through questions

The geometric series \(1 + x + x^2 + \cdots\) adds to \(\frac{1}{1-x}\). It converges provided \(|x| < 1\). The sum of \(n\) terms is \(\frac{1-x^n}{1-x}\). The derivatives of the series match the derivatives of \(1/(1-x)\) at the point \(x = 0\), where the \(n\)th derivative is \(n!\). The decimal 1.111... is the geometric series at \(x = 1/9\) and equals the fraction \(\frac{1}{9}\). The decimal .666... multiplies this by \(\frac{2}{3}\). The decimal .999... is the same as \(1\).

The derivative of the geometric series is \(\frac{x}{1-x}\). This also comes from squaring the \(1/1-x\) series. By choosing \(x = 1/9\), the decimal 0.102030405 is close to \(\frac{1}{9}\). The differential equation \(dy/dx = y^2\) is solved by the geometric series, going term by term starting from \(y(0) = 0\).
The integral of the geometric series is \( \frac{1}{1-x} \). At \( x = 1 \) this becomes the \( \frac{1}{x} \) series, which diverges. At \( x = 1 \) we have \( \ln 2 = \frac{1}{x} \). The change from \( x \) to \( -x \) produces the series \( \frac{1}{1+x} = \frac{1}{x} \) and \( \ln(1+x) = \frac{1}{x} \).

In the geometric series, changing to \( x^2 \) or \( -x^2 \) gives \( \frac{1}{1-x^2} \) and \( \frac{1}{1+x^2} \). Integrating the last one yields \( x - \frac{1}{2}x^3 + \frac{1}{3}x^5 \ldots = \frac{1}{x} \). The angle whose tangent is \( x = 1 \) is \( \tan^{-1} 1 = \frac{\pi}{4} \). Then substituting \( x = 1 \) gives the series \( \pi = \frac{\pi}{4} \).

1. The geometric series is \( 1 + x + x^2 + \ldots = G \). Another way to discover \( G \) is to multiply by \( x \). Then \( x + x^2 + x^3 + \ldots = xG \), and this can be subtracted from the original series. What does that leave, and what is \( G \)?

2. A basketball is dropped 10 feet and bounces back 6 feet. After every fall it recovers \( \frac{3}{5} \) of its height. What total distance does the ball travel, bouncing forever?

3. Find the sums of \( \frac{1}{1} + \frac{1}{2} + \frac{1}{2} + \ldots \) and \( 1 - \frac{1}{1} + \frac{1}{2} + \ldots \) and \( 10 - 1 + .1 - .01 \ldots \) and 3.040404…

4. Replace \( x \) by \( 1-x \) in the geometric series to find a series for \( 1/x \). Integrate to find a series for \( \ln x \). These are power series “around the point \( x = 1 \).” What is their sum at \( x = 0 \)?

5. What is the second derivative of the geometric series, and what is its sum at \( x = 1 \)?

6. Multiply the series \( (1 + x + x^2 + \ldots)(1 - x + x^2 - \ldots) \) and find the product by comparing with equation (14).

7. Start with the fraction \( \frac{1}{2} \). Divide 7 into 1.000… (by long division or calculator) until the numbers start repeating. Which is the first number to repeat? How do you know that the next ________ digits will be the same as the first?

Note about the fractions \( 1/q, 10/q, 100/q, \ldots \) All remainders are less than \( q \) so eventually two remainders are the same. By subtraction, \( q \) goes evenly into a power \( 10^n \) minus a smaller power \( 10^{n-a} \). Thus \( q = 10^n - 10^{n-a} \) for some \( c \) and \( 1/q \) has a repeating decimal:

\[
\frac{1}{q} = \frac{c}{10^n - 10^{n-a}} = \frac{c}{10^a} \left( \frac{1}{1 + \frac{1}{10} + \frac{1}{10^2} + \ldots} \right)
\]

Conclusion: Every fraction equals a repeating decimal.

8. Find the repeating decimal for \( \frac{1}{3} \) and read off \( c \). What is the number \( n \) of digits before it repeats?

9. From the fact that every \( q \) goes evenly into a power \( 10^n \) minus a smaller power, show that all primes except 2 or 5 go evenly into 9 or 99 or 999 or ….

10. Explain why \( .010010001 \ldots \) cannot be a fraction (the number of zeros increases).

11. Show that \( .123456789101112 \ldots \) is not a fraction.

12. From the geometric series, the repeating decimal \( 1.065065065 \ldots \) equals what fraction? Explain why every repeating decimal equals a fraction.

13. Write .878787… and .123123… as fractions and as geometric series.

14. Find the square of 1.111… as an infinite series.

Find the functions which equal the sums 15–24.

15. \( x + x^3 + x^5 + \ldots \)
16. \( 1 - 2x + 4x^2 - \ldots \)
17. \( x^3 + x^6 + x^9 + \ldots \)
18. \( \frac{1}{2}x - \frac{1}{4}x^2 + \frac{1}{8}x^3 - \ldots \)
19. \( \ln x + (\ln x)^3 + (\ln x)^5 + \ldots \)
20. \( x - 2x^2 + 3x^3 - \ldots \)
21. \( \frac{1}{x} + \frac{1}{x^3} + \frac{1}{x^5} + \ldots \)
22. \( x + \frac{x}{1 + x} + \frac{x}{1 + \sqrt{1+x^2}} + \ldots \)
23. \( \tan x - \frac{1}{3}\tan^3 x + \frac{1}{5}\tan^5 x - \ldots \)
24. \( e^x + e^{2x} + e^{3x} + \ldots \)
25. Multiply the series for \( 1/(1-x) \) and \( 1/(1+x) \) to find the coefficients of \( x, x^2, x^3 \) and \( x^4 \).

26. Compare the integral of \( 1 + x^3 + x^6 + \ldots \) to equation (13) and find \( \int dx/(1-x^2) \).

27. What fractions are close to .2468 and .987654321?

28. Find the first three terms in the series for \( 1/(1-x)^3 \).

Add up the series 29–34. Problem 34 comes from (18).

29. \( \frac{2}{3} + \frac{2}{3^2} + \frac{2}{3^3} + \ldots \)
30. \( 1 + .02 + .003 + \ldots \)
31. \( 1 + \frac{1}{2}(.01) + \frac{1}{2}(.001) + \ldots \)
32. \( 1 - \frac{1}{2}(.01) + \frac{1}{2}(.001) - \ldots \)
33. \( 1 + \frac{1}{2}(.001) + \frac{1}{3}(.00001) + \ldots \)
34. \( 1 - \frac{1}{3} + \frac{1}{5} - \frac{1}{3} + \ldots \)

35. Compute the \( n \)th derivative of \( 1 + 2x + 3x^2 + \ldots \) at \( x = 0 \). Compute also the \( n \)th derivative of \( (1-x)^{-2} \).

36. The differential equation \( dy/dx = y^2 \) starts from \( y(0) = b \). From the equation and its derivatives find \( y', y'', y''' \) at \( x = 0 \), and construct the start of a series that matches these derivatives. Can you recognize \( y(x) \)?

37. The equation \( dy/dx = y^2 \) has the differential form \( dy/y^2 = dx \). Integrate both sides and choose the integration constant so that \( y = b \) at \( x = 0 \). Solve for \( y(x) \) and compare with Problem 36.

38. In a bridge game, what is the average number \( \mu \) of deals until you get the best hand? The probability on the first deal is \( p_1 = \frac{1}{x} \). Then \( p_2 = (\frac{1}{x})^2 \) (probability of missing on the first) times (probability of winning on the second). Generally \( p_n = (\frac{1}{x})^{n-1} \). The mean value \( \mu \) is \( p_1 + 2p_2 + 3p_3 + \ldots \).

39. Show that \( (\Sigma a_n)(\Sigma b_n) = \Sigma a_n b_n \) is ridiculous.

40. Find a series for \( \ln \frac{1}{2} \) by choosing \( x \) in (10b). Find a series for \( \ln 3 \) by choosing \( x \) in (13). How is \( \ln \frac{1}{2} \) related to \( \ln 3 \), and which series converges faster?
10.2 Convergence Tests: Positive Series

This is the third time we have stopped the calculations to deal with the definitions. Chapter 2 said what a derivative is. Chapter 5 said what an integral is. Now we say what the sum of a series is—if it exists. In all three cases a limit is involved. That is the formal, careful, cautious part of mathematics, which decides if the active and progressive parts make sense.

The series \( \frac{1}{2} + \frac{1}{4} + \frac{1}{8} + \ldots \) converges to 1. The series \( 1 + \frac{1}{2} + \frac{1}{4} + \ldots \) diverges to infinity. The series \( 1 - \frac{1}{2} + \frac{1}{4} - \ldots \) converges to \( \ln 2 \). When we speak about convergence or divergence of a series, we are really speaking about convergence or divergence of its "partial sums."

**Definition 1** The partial sum \( s_n \) of the series \( a_1 + a_2 + a_3 + \ldots \) stops at \( a_n \):

\[
s_n = \text{sum of the first } n \text{ terms} = a_1 + a_2 + \ldots + a_n.
\]

Thus \( s_n \) is part of the total sum. The example \( \frac{1}{2} + \frac{1}{4} + \frac{1}{8} + \ldots \) has partial sums

\[
s_1 = \frac{1}{2}, \quad s_2 = \frac{3}{4}, \quad s_3 = \frac{7}{8}, \quad s_n = 1 - \frac{1}{2^n}.
\]

Those add up larger and larger parts of the series—what is the sum of the whole series? The answer is: The series \( \frac{1}{2} + \frac{1}{4} + \ldots \) converges to 1 because its partial sums \( s_n \) converge to 1. The series \( a_1 + a_2 + a_3 + \ldots \) converges to \( s \) when its partial sums—going further and further out—approach this limit \( s \). Add the \( a \)'s, not the \( s \)'s.

**Definition 2** The sum of a series is the limit of its partial sums \( s_n \).

We repeat: if the limit exists. The numbers \( s_n \) may have no limit. When the partial sums jump around, the whole series has no sum. Then the series does not converge. When the partial sums approach \( s \), the distant terms \( a_n \) are approaching zero. More than that, the sum of distant terms is approaching zero.

The new idea (\( \Sigma a_n = s \)) has been converted to the old idea (\( s_n \to s \)).

**Example 1** The geometric series \( \frac{1}{10} + \frac{1}{100} + \frac{1}{1000} + \ldots \) converges to \( s = \frac{1}{9} \).

The partial sums \( s_1, s_2, s_3, s_4 \) are \( .1, .11, .111, .1111 \). They are approaching \( s = \frac{1}{9} \).
Note again the difference between the series of a's and the sequence of s's. The series 
1 + 1 + 1 + \cdots diverges because the sequence of s's is 1, 2, 3, \ldots. A sharper example is 
the harmonic series: 1 + \frac{1}{2} + \frac{1}{3} + \cdots diverges because its partial sums 1, 1\frac{1}{2}, \ldots 
eventually go past every number s. We saw that in 2.6 and will see it again here.

Do not confuse \( a_n \to 0 \) with \( s_n \to s \). You cannot be sure that a series converges, just 
on the basis that \( a_n \to 0 \). The harmonic series is the best example: \( a_n = \frac{1}{n} \to 0 \) but 
still \( s_n \to \infty \). This makes infinite series into a delicate game, which mathematicians 
 enjoy. The line between divergence and convergence is hard to find and easy to cross.

A slight push will speed up \( a_n \to 0 \) and make the \( s_n \) converge. Even though \( a_n \to 0 \) 
does not by itself guarantee convergence, it is the first requirement:

**10A** If a series converges \( (s_n \to s) \) then its terms must approach zero \( (a_n \to 0) \).

Proof Suppose \( s_n \) approaches \( s \) (as required by convergence). Then also \( s_{n-1} \) 
approaches \( s \), and the difference \( s_n - s_{n-1} \) approaches zero. That difference is \( a_n \). So 
\( a_n \to 0 \).

**EXAMPLE 1** (continued) For the geometric series \( 1 + x + x^2 + \cdots \), the test \( a_n \to 0 \) is 
the same as \( x^n \to 0 \). The test is failed if \( |x| \geq 1 \), because the powers of \( x \) don't go to 
zero. Automatically the series diverges. The test is passed if \( -1 < x < 1 \). But to prove 
convergence, we cannot rely on \( a_n \to 0 \). It is the partial sums that must converge:

\[
\begin{align*}
  s_n &= 1 + x + \cdots + x^{n-1} = \frac{1 - x^n}{1 - x} \\
  s_n &\to \frac{1}{1 - x} 
\end{align*}
\]

This is \( s \).

For other series, first check that \( a_n \to 0 \) (otherwise there is no chance of convergence). The \( a_n \) will not have the special form \( x^n \)—so we need sharper tests.

The geometric series stays in our mind for this reason. Many convergence tests are 
comparisons with that series. The right comparison gives enough information:

If \( |a_1| < \frac{1}{2} \) and \( |a_2| < \frac{1}{4} \) and \( \ldots \), then \( a_1 + a_2 + \ldots \) converges faster than \( \frac{1}{2} + \frac{1}{4} + \ldots \).

More generally, the terms in \( a_1 + a_2 + a_3 + \ldots \) may be smaller than \( ax + ax^2 + ax^3 + \ldots \). Provided \( x < 1 \), the second series converges. Then \( \sum a_n \) also 
converges. We move now to convergence by comparison or divergence by comparison.

Throughout the rest of this section, all numbers \( a_n \) are assumed positive.

**COMPARISON TEST AND INTEGRAL TEST**

In practice it is rare to compute the partial sums \( s_n = a_1 + \cdots + a_n \). Usually a simple 
formula can't be found. We may never know the exact limit \( s \). But it is still possible 
to decide convergence—whether there is a sum—by comparison with another series 
that is known to converge.

**10B** (Comparison test) Suppose that \( 0 \leq a_n \leq b_n \) and \( \sum b_n \) converges. Then 
\( \sum a_n \) converges.

The smaller terms \( a_n \) add to a smaller sum: \( \sum a_n \) is below \( \sum b_n \) and must converge.

On the other hand suppose \( a_n \geq c_n \) and \( \sum c_n = \infty \). This comparison forces \( \sum a_n = \infty \). 
A series diverges if it is above another divergent series.

Note that a series of positive terms can only diverge "to infinity." It cannot oscillate, 
because each term moves it forward. Either the \( s_n \) creep up on \( s \), passing every number 
below it, or they pass all numbers and diverge. If an increasing sequence \( s_n \) is bounded 
above, it must converge. The line of real numbers is complete, and has no holes.
The harmonic series $1 + \frac{1}{2} + \frac{1}{3} + \cdots$ diverges to infinity.

A comparison series is $1 + \frac{1}{2} + \frac{1}{4} + \frac{1}{8} + \frac{1}{16} + \cdots$. The harmonic series is larger. But this comparison series is really $1 + \frac{1}{2} + \frac{1}{3} + \cdots$ because $\frac{1}{2} = \frac{2}{4} = \frac{4}{8}$.

The comparison series diverges. The harmonic series, above it, must also diverge.

To apply the comparison test, we need something to compare with. In Example 2, we thought of another series. It was convenient because of those $\frac{1}{2}$'s. But a different series will need a different comparison, and where will it come from? There is an automatic way to think of a comparison series. It comes from the integral test.

Allow me to apply the integral test to the same example. To understand the integral test, look at the areas in Figure 10.2. The test compares rectangles with curved areas.

**Example 2** (again) Compare $1 + \frac{1}{2} + \frac{1}{3} + \cdots$ with the area under the curve $y = \frac{1}{x}$.

Every term $a_n = \frac{1}{n}$ is the area of a rectangle. We are comparing it with a curved area $c_n$. Both areas are between $x = n$ and $x = n + 1$, and the rectangle is above the curve. So $a_n > c_n$:

$$\text{rectangular area } a_n = \frac{1}{n} \text{ exceeds curved area } c_n = \int_{n}^{n+1} \frac{dx}{x}.$$

Here is the point. Those $c_n$'s look complicated, but we can add them up. The sum $c_1 + \cdots + c_n$ is the whole area, from 1 to $n + 1$. It equals $\ln(n + 1)$ — we know the integral of $1/x$. We also know that the logarithm goes to infinity.

The rectangular area $1 + \frac{1}{2} + \cdots + \frac{1}{n}$ is above the curved area. By comparison of areas, the harmonic series diverges to infinity — a little faster than $\ln(n + 1)$.

**Remark** The integral of $1/x$ has another advantage over the series with $\frac{1}{2}$'s. First, the integral test was automatic. From $1/n$ in the series, we went to $1/x$ in the integral. Second, the comparison is closer. Instead of adding only $\frac{1}{2}$ when the number of terms is doubled, the true partial sums grow like $\ln n$. To prove that, put rectangles under the curve.

Rectangles below the curve give an area below the integral. Figure 10.2b omits the first rectangle, to get under the curve. Then we have the opposite to the first comparison — the sum is now smaller than the integral:

$$\frac{1}{2} + \frac{1}{3} + \cdots + \frac{1}{n} \leq \int_{1}^{n} \frac{dx}{x} = \ln n.$$

Adding 1 to both sides, $s_n$ is below $1 + \ln n$. From the previous test, $s_n$ is above $\ln(n + 1)$. That is a narrow space — we have an excellent estimate of $s_n$. The sum of $1/n$
10.2 Convergence Tests: Positive Series

and the integral of $1/x$ diverge together. Problem 43 will show that the difference between $s_n$ and $\ln n$ approaches "Euler's constant," which is $\gamma = .577 \ldots$

**Main point:** Rectangular area is $s_n$. Curved area is close. We are using integrals to help with sums (it used to be the opposite).

**Question** If a computer adds a million terms every second for a million years, how large is the partial sum of the harmonic series?

**Answer** The number of terms is $n = 60^2 \cdot 24 \cdot 365 \cdot 10^{12} < 3.2 \cdot 10^{19}$. Therefore $\ln n$ is less than $\ln 3.2 + 19 \ln 10 < 45$. By the integral test $s_n < 1 + \ln n$, the partial sum after a million years has not reached 46.

For other series, $1/x$ changes to a different function $y(x)$. At $x = n$ this function must equal $a_n$. Also $y(x)$ must be decreasing. Then a rectangle of height $a_n$ is above the graph to the right of $x = n$, and below the graph to the left of $x = n$. The series and the integral box each other in: left sum $\geq$ integral $\geq$ right sum. The reasoning is the same as it was for $a_n = 1/n$ and $y(x) = 1/x$: There is finite area in the rectangles when there is finite area under the curve.

When we can't add the $a$'s, we integrate $y(x)$ and compare areas:

| 10C (Integral test) | If $y(x)$ is decreasing and $y(n)$ agrees with $a_n$, then $a_1 + a_2 + a_3 + \ldots$ and $\int_1^\infty y(x) \, dx$ both converge or both diverge. |

**EXAMPLE 3** The "$p$-series" $\frac{1}{2^p} + \frac{1}{3^p} + \frac{1}{4^p} + \ldots$ converges if $p > 1$. Integrate $y = \frac{1}{x^p}$:

$$\frac{1}{n^p} < \int_{n-1}^n \frac{dx}{x^p}$$

so by addition

$$\sum_{n=2}^\infty \frac{1}{n^p} < \int_1^\infty \frac{dx}{x^p}$$

In Figure 10.2c, the area is finite if $p > 1$. The integral equals $\int x^{1-p} (1-p) \, dx$, which is $1/(p-1)$. **Finite area means convergent series.** If $1/1^p$ is the first term, add 1 to the curved area:

$$\frac{1}{1^p} + \frac{1}{2^p} + \frac{1}{3^p} + \ldots < 1 + \frac{1}{p-1} = \frac{p}{p-1}.$$ 

The borderline case $p = 1$ is the harmonic series (divergent). By the comparison test, every $p < 1$ also produces divergence. Thus $\Sigma 1/\sqrt{n}$ diverges by comparison with $\int dx/\sqrt{x}$ and also by comparison with $\Sigma 1/n$. Section 7.5 on improper integrals runs parallel to this section on "improper sums" (infinite series).

Notice the special cases $p = 2$ and $p = 3$. The series $1 + \frac{1}{2} + \frac{1}{3} + \ldots$ converges. Euler found $\pi^2/6$ as its sum. The series $1 + \frac{1}{2} + \frac{1}{3} + \ldots$ also converges. That is proved by comparing $\Sigma 1/n^3$ with $\Sigma 1/n^2$ or with $\int dx/x^3$. But the sum for $p = 3$ is unknown.

**Extra credit problem** The sum of the $p$-series leads to the most important problem in pure mathematics. The "zeta function" is $Z(p) = \Sigma 1/n^p$, so $Z(2) = \pi^2/6$ and $Z(3)$ is unknown. Riemann studied the complex numbers $p$ where $Z(p) = 0$ (there are infinitely many). He conjectured that the real part of those $p$ is always $1/2$. That has been tested for the first billion zeros, but never proved.

**Comparison with the Geometric Series**

We can compare any new series $a_1 + a_2 + \ldots$ with $1 + x + \ldots$. Remember that the first million terms have nothing to do with convergence. It is further out, as $n \to \infty$, that the comparison stands or falls. We still assume that $a_n > 0$. 

(Ratio test) If \( \frac{a_{n+1}}{a_n} \) approaches a limit \( L < 1 \), the series converges.

(Root test) If the \( n \)th root \( (a_n)^{1/n} \) approaches \( L < 1 \), the series converges.

Roughly speaking, these tests make \( a_n \) comparable with \( L^n \)—therefore convergent. The tests also establish divergence if \( L > 1 \). They give no decision when \( L = 1 \). Unfortunately \( L = 1 \) is the most important and the hardest case.

On the other hand, you will now see that the ratio test is fairly easy.

**EXAMPLE 4** The geometric series \( x + x^2 + \cdots \) has ratio exactly \( x \). The \( n \)th root is also exactly \( x \). So \( L = x \). There is convergence if \( x < 1 \) (known) and divergence if \( x > 1 \) (also known). The divergence of \( 1 + 1 + \cdots \) is too delicate (!) for the ratio test and root test, because \( L = 1 \).

**EXAMPLE 5** The \( p \)-series has \( a_n = 1/n^p \) and \( a_{n+1}/a_n = n^p/(n+1)^p \). The limit as \( n \to \infty \) is \( L = 1 \), for every \( p \). The ratio test does not feel the difference between \( p = 2 \) (convergence) and \( p = 1 \) (divergence) or even \( p = -1 \) (extreme divergence). Neither does the root test. So the integral test is sharper.

**EXAMPLE 6** A combination of \( p \)-series and geometric series can now be decided:

\[
\frac{x}{1^p} + \frac{x^2}{2^p} + \cdots + \frac{x^n}{n^p} + \cdots \text{ has ratio } \frac{a_{n+1}}{a_n} = \frac{x^{n+1}}{(n+1)^p} \frac{n^p}{x^n} \text{ approaching } L = x.
\]

It is \(|x| < 1\) that decides convergence, not \( p \). The powers \( x^n \) are stronger than any \( n^p \). The factorials \( n! \) will now prove stronger than any \( x^n \).

**EXAMPLE 7** The exponential series \( e^x = 1 + x + \frac{x^2}{2!} + \frac{x^3}{3!} + \cdots \) converges for all \( x \).

The terms of this series are \( x^n/n! \). The ratio between neighboring terms is

\[
\frac{x^{n+1}/(n+1)!}{x^n/n!} = \frac{x}{n+1},
\]

which approaches \( L = 0 \) as \( n \to \infty \).

With \( x = 1 \), this ratio test gives convergence of \( \sum 1/n! \). The sum is \( e \). With \( x = 4 \), the larger series \( \sum 4^n/n! \) also converges. We know this sum too—it is \( e^4 \). Also the sum of \( x^n n!/n! \) converges for any \( x \) and \( p \). Again \( L = 0 \)—the ratio test is not even close. The factorials take over, and give convergence.

Here is the proof of convergence when the ratios approach \( L < 1 \). Choose \( x \) halfway from \( L \) to 1. Then \( x < 1 \). Eventually the ratios go below \( x \) and stay below:

\[
a_{N+1}/a_N < x \quad a_{N+2}/a_{N+1} < x \quad a_{N+3}/a_{N+2} < x \quad \cdots
\]

Multiply the first two inequalities. Then multiply all three:

\[
a_{N+1}/a_N < x \quad a_{N+2}/a_N < x^2 \quad a_{N+3}/a_N < x^3 \quad \cdots
\]

Therefore \( a_{N+1} + a_{N+2} + a_{N+3} + \cdots \) is less than \( a_N(x + x^2 + x^3 + \cdots) \). Since \( x < 1 \), comparison with the geometric series gives convergence.

**EXAMPLE 8** The series \( \sum 1/n^n \) is ideal for the root test. The \( n \)th root is \( 1/n \). Its limit is \( L = 0 \). Convergence is even faster than for \( e = \sum 1/n! \). The root test is easily explained, since \( (a_n)^{1/n} < x \) yields \( a_n < x^n \) and \( x \) is close to \( L < 1 \). So we compare with the geometric series.
10.2 Convergence Tests: Positive Series

SUMMARY FOR POSITIVE SERIES

The convergence of geometric series and $p$-series and exponential series is settled. I will put these $a_n$'s in a line, going from most divergent to most convergent. The crossover to convergence is after $1/n$:

$$1 + 1 + \ldots \quad (p < 1) \quad \frac{1}{n^p} \quad \frac{1}{n} \quad \frac{1}{n^p} \quad (p > 1) \quad \frac{n}{2^n} \quad \frac{1}{2^n} \quad \frac{n}{n!} \quad \frac{1}{n!} \quad \frac{1}{n^n}$$

10A 10B and 10C 10D and 10E

You should know that this crossover is not as sharp as it looks. On the convergent side, $1/n(\ln n)^2$ comes before all those $p$-series. On the divergent side, $1/n(\ln n)$ and $1/n(\ln n)(\ln \ln n)$ belong after $1/n$. For any divergent (or convergent) series, there is another that diverges (or converges) more slowly.

Thus there is no hope of an ultimate all-purpose comparison test. But comparison is the best method available. Every series in that line can be compared with its neighbors, and other series can be placed in between. It is a topic that is understood best by examples.

EXAMPLE 9

$$\sum \frac{1}{\ln n}$$ diverges because $\sum \frac{1}{n}$ diverges. The comparison uses $\ln n < n$.

EXAMPLE 10

$$\sum \frac{1}{n(\ln n)^2} \approx \int \frac{dx}{x(\ln x)^2} < \infty \quad \sum \frac{1}{n(\ln n)} \approx \int \frac{dx}{x(\ln x)} = \infty.$$  

The indefinite integrals are $-1/(\ln x)$ and $\ln(\ln x)$. The first goes to zero as $x \to \infty$; the integral and series both converge. The second integral $\ln(\ln x)$ goes to infinity—very slowly but it gets there. So the second series diverges. These examples squeeze new series into the line, closer to the crossover.

EXAMPLE 11

$$\frac{1}{n^2 + 1} < \frac{1}{n^2} \quad \text{so} \quad \frac{1}{2} + \frac{1}{5} + \frac{1}{10} + \ldots < \frac{1}{1} + \frac{1}{4} + \frac{1}{9} + \ldots \quad (\text{convergence}).$$

The constant 1 in this denominator has no effect—and again in the next example.

EXAMPLE 12

$$\frac{1}{2n-1} > \frac{1}{2n} \quad \text{so} \quad \frac{1}{1} + \frac{1}{3} + \frac{1}{5} + \ldots > \frac{1}{2} + \frac{1}{4} + \frac{1}{6} + \ldots.$$  

$\sum 1/2n$ is $1/2$ times $\sum 1/n$, so both series diverge. Two series behave in the same way if the ratios $a_n/b_n$ approach $L > 0$. Examples 11–12 have $n^2/(n^2 + 1) \to 1$ and $2n/(2n - 1) \to 1$. That leads to our final test:

10F (Limit comparison test) If the ratio $a_n/b_n$ approaches a positive limit $L$, then $\sum a_n$ and $\sum b_n$ either both diverge or both converge.

Reason: $a_n$ is smaller than $2Lb_n$ and larger than $\frac{1}{L}b_n$, at least when $n$ is large. So the two series behave in the same way. For example $\sum \sin (7/n^p)$ converges for $p > 1$, not for $p \leq 1$. It behaves like $\sum 1/n^p$ (here $L = 7$). The tail end of a series (large $n$) controls convergence. The front end (small $n$) controls most of the sum.

There are many more series to be investigated by comparison.
10.2 EXERCISES

Read-through questions

The convergence of \( a_1 + a_2 + \cdots \) is decided by the partial sums \( s_n = \sum_{k=1}^{n} a_k \). If the \( s_n \) approach \( s \), then \( \Sigma a_n = \lim_{n \to \infty} s_n \). For the geometric series \( 1 + x + \cdots \) the partial sums are \( s_n = 1 - x^{n+1} \). In that case \( s_n \to 1/(1-x) \) if and only if \( x \leq 1 \). In all cases the limit \( s_n \to s \) requires that \( a_n \to 0 \). But the harmonic series \( a_n = 1/n \) shows that we can have \( a_n \to 0 \) and still the series \( \sum a_n \) diverges.

The comparison test says that if \( 0 < a_n < b_n \) then \( \sum b_n \) converges if and only if \( \sum a_n \) converges. By this test the \( p \)-series \( \sum 1/n^p \) converges if and only if \( p > 1 \). For the harmonic series \( (p = 1) \), \( s_n = 1 + \cdots + 1/n \) is close to the integral \( \int_1^n 1/x \, dx \).

The \( n \)-test applies when \( a_n+1/a_n \to L \). There is convergence if \( L < 1 \), divergence if \( L > 1 \), and no decision if \( L = 1 \). The same is true for the \( r \)-test, when \( (a_n)_{n=1}^\infty \) is non-negative.

There is no sharp line between \( x \) and \( y \). But if \( E_b \) converges and \( a_n/b_n \to L \), it follows from the test that \( \sum a_n \) also converges.

Here is a quick proof that a finite sum \( 1 + 1 + 1 + \cdots = s \) is impossible. Division by 2 would give \( 1/2 + 1/2 + \cdots = s/2 \). Subtraction would leave \( 1 + 1/2 + 1/4 + \cdots = s \). Those last two series cannot both add to \( s \) because \( 1/2 \) because \( n! \) grows faster than any \( x^k \) or \( n^k \).

Behind every decimal \( s = \ldots abc \ldots \) is a convergent series \( a/10 + b/100 + \cdots \). By a comparison test prove convergence.

Decide convergence or divergence in 7–10 (and give a reason).

1 Establish convergence or divergence in 11–20 by a comparison test.

21 Find the limit \( L \) in the ratio test or root test.

29 \((1 - \frac{1}{4}) + (1 - \frac{1}{3}) + (1 - \frac{1}{4})\) is "telescoping" because \( 1/4 \) and \( 1/3 \) cancel \(-\frac{1}{4}\) and \(-\frac{1}{3}\). Add the infinite telescoping series

30 Compute the sum \( s \) for other "telescoping series":

(a) \(\left(\frac{1}{1} - \frac{1}{3}\right) + \left(\frac{1}{2} - \frac{1}{5}\right) + \left(\frac{1}{3} - \frac{1}{7}\right)\)...

(b) \(\ln 1 + \ln \frac{1}{2} + \ln \frac{1}{3} + \cdots\)

31 In the integral test, what sum is larger than \( \int_1^\infty y(x) \, dx \) and what sum is smaller? Draw a figure to illustrate.

32 Comparing sums with integrals, find numbers larger and smaller than

\( s_n = 1 + \frac{1}{3} + \cdots + \frac{1}{2n-1} \) and \( s_n = 1 + \frac{1}{8} + \cdots + \frac{1}{n^3} \).

33 Which integral test shows that \( \sum a_n \) converges? What is the sum?

34 Which integral test shows that \( \sum b_n \) converges? What is the sum?
Decide for or against convergence in 35–42, based on \( \int y(x) \, dx \). 

35 \[ \sum_{n=1}^{\infty} \frac{1}{n^2 + 1} \]  
36 \[ \sum_{n=1}^{\infty} \frac{1}{3n + 5} \]  
37 \[ \sum_{n=1}^{\infty} \frac{n}{n^2 + 1} \]  
38 \[ \sum_{n=1}^{\infty} \frac{\ln n}{n} \left( \text{is } \frac{\ln x}{x} \text{ decreasing?} \right) \]  
39 \[ \sum_{n=1}^{\infty} \frac{n^2/n^e}{n} \]  
40 \[ \sum_{n=1}^{\infty} \frac{n}{n^r/n^e} \]  
41 \[ \sum_{n=1}^{\infty} \frac{n}{n+1} \]  
42 \[ \sum_{n=1}^{\infty} \frac{n}{n^2} \]  
43 (a) Explain why \( D_n = \left( 1 + \frac{1}{2} + \cdots + \frac{1}{n} \right) - \ln n \) is positive by using rectangles as in Figure 10.2.  
(b) Show that \( D_{n+1} \) is less than \( D_n \) by proving that 
\[ \frac{1}{n+1} < \int_n^{n+1} \frac{1}{x} \, dx. \]  
(c) (Calculator) The decreasing \( D_n \)’s must approach a limit. Compute them until they go below \( .6 \) and below \( .58 \) (when?). The limit of the \( D_n \) is Euler’s constant \( \gamma = .577… \).  
44 In the harmonic series, use \( s_n \approx .577 + \ln n \) to show that 
\( s_n = 1 + \frac{1}{2} + \cdots + \frac{1}{n} \) needs more than 600 terms to reach \( s_n > 7 \).  
How many terms for \( s_n > 10 \)?  
45 (a) Show that 
\[ 1 + \frac{1}{2} + \frac{1}{3} - \frac{1}{4} - \cdots - \frac{1}{2n} = \frac{1}{n+1} + \cdots + \frac{1}{2n} \]  
by adding \( 2 \left( \frac{1}{2} + \frac{1}{4} + \cdots + \frac{1}{2n} \right) \) to both sides.  
(b) Why is the right side close to \( \ln 2n - \ln n \)? Deduce that 
\[ 1 - \frac{1}{2} + \frac{1}{3} - \frac{1}{4} + \cdots \]  
approaches \( \ln 2 \).  
46 Every second a computer adds a million terms of \( \sum 1/(n \ln n) \). By comparison with \( \int dx/(x \ln x) \), estimate the partial sum after a million years (see Question in text).  
47 Estimate \[ \sum_{n=1}^{1000} \frac{1}{n^2} \] by comparison with an integral.  
48 If \( \sum a_n \) converges (all \( a_n > 0 \)) show that \( \sum a_n^2 \) converges.  
49 If \( \sum a_n \) converges (all \( a_n > 0 \)) show that \( \sum \sin a_n \) converges. How could \( \sum a_n \) converge when \( \sum a_n \) diverges?  
50 The \( n \)th prime number \( p_n \) satisfies \( p_n/n \ln n \to 1 \). Prove that 
\[ \sum \frac{1}{p_n} = \frac{1}{2} + \frac{1}{3} + \frac{1}{5} + \frac{1}{7} + \frac{1}{11} + \cdots \] diverges.  

Construct a series \( \sum a_n \) that converges faster than \( \sum a_n \), but slower than \( \sum c_n \) (meaning \( a_n/b_n \to 0 \), \( a_n/c_n \to \infty \)).  
51 \( b_n = 1/n^2 \), \( c_n = 1/n^3 \)  
52 \( b_n = n! \), \( c_n = (\frac{1}{2})^n \)  
53 \( b_n = 1/n! \), \( c_n = 1/n^e \)  
54 \( b_n = 1/n! \), \( c_n = 1/e^x \)  
In Problem 53 use Stirling’s formula \( \sqrt{2\pi n} n^n/e^n \to 1 \).  
55 For the series \( \frac{1}{2} + \frac{1}{3} + \frac{1}{6} + \frac{1}{4} + \frac{1}{6} + \cdots \), show that the ratio test fails. The roots \( (a_n)^{1/n} \) do approach a limit \( L \). Find \( L \) from the even terms \( a_{2k} = 1/2^k \). Does the series converge?  
56 (For instructors) If the ratios \( a_{n+1}/a_n \) approach a positive limit \( L \) show that the roots \( (a_n)^{1/n} \) also approach \( L \).  

Decide convergence in 57–66 and name your test.  
57 \[ \sum_{n=1}^{\infty} \frac{1}{(\ln n)^n} \]  
58 \[ \sum_{n=1}^{\infty} \frac{1}{n^n} \]  
59 \[ \sum_{n=1}^{\infty} \frac{1}{10^n} \]  
60 \[ \sum_{n=1}^{\infty} \frac{1}{\ln (10^n)} \]  
61 \[ \sum_{n=1}^{\infty} \frac{1}{n^n} \]  
62 \[ \sum_{n=1}^{\infty} \frac{1}{e^n} \]  
63 \[ \sum_{n=1}^{\infty} \frac{1}{(\ln n)^p} \] (test all \( p \))  
64 \[ \sum_{n=1}^{\infty} \frac{1}{n^{1/p}} \] (test all \( p \))  
65 \[ \sum_{n=1}^{\infty} \frac{1}{n^{1/2}} \] (test all \( p, q \))  
66 \[ \sum_{n=1}^{\infty} \frac{1}{n^{1/4}} \] (test all \( p, q \))  
67 Suppose \( a_n/b_n \to 0 \) in the limit comparison test. Prove that \( \sum a_n \) converges if \( \sum b_n \) converges.  
68 Can you invent a series whose convergence you and your instructor cannot decide?
Changing a negative number from \( a_n \) to \( |a_n| \) increases the sum. Main point: The smaller series \( \Sigma a_n \) is sure to converge if \( \Sigma |a_n| \) converges.

\[
\textbf{40G} \quad \text{If } \Sigma |a_n| \text{ converges then } \Sigma a_n \text{ converges (absolutely). But } \Sigma a_n \text{ might converge, as in the series for } \pi, \text{ even if } \Sigma |a_n| \text{ diverges to infinity.}
\]

EXAMPLE 1 Start with the positive series \( \frac{1}{2} + \frac{1}{4} + \frac{1}{6} + \cdots \). Change any signs to minus. Then the new series converges (absolutely). The right choice of signs will make it converge to any number between \(-1\) and \(1\).

EXAMPLE 2 Start with the alternating series \( 1 - \frac{1}{2} + \frac{1}{3} - \frac{1}{4} + \cdots \) which converges to \( \ln 2 \). Change to plus signs. The new series \( 1 + \frac{1}{2} + \frac{1}{3} + \cdots \) diverges to infinity. The original alternating series was not absolutely convergent. It was only "conditionally convergent." A series can converge (conditionally) by a careful choice of signs—even if \( \Sigma |a_n| = \infty \).

\[
\text{If } \Sigma |a_n| \text{ converges then } \Sigma a_n \text{ converges.} \text{ Here is a quick proof. The numbers } a_n + |a_n| \text{ are either zero (if } a_n \text{ is negative) or } 2|a_n|. \text{ By comparison with } \Sigma 2|a_n|, \text{ which converges, } \Sigma (a_n + |a_n|) \text{ must converge. Now subtract the convergent series } \Sigma |a_n|. \text{ The difference } \Sigma a_n \text{ also converges, completing the proof. All tests for positive series (integral, ratio, comparison, ... ) apply immediately to absolute convergence, because we switch to } |a_n|.
\]

EXAMPLE 3 Start with the geometric series \( \frac{1}{2} + \frac{1}{4} + \frac{1}{8} + \cdots \) which converges to \( \frac{1}{2} \). Change any of those signs to minus. Then the new series must converge (absolutely). But the sign changes cannot achieve all sums between \(-\frac{1}{2}\) and \(\frac{1}{2}\). This time the sums belong to the famous (and very thin) Cantor set of Section 3.7.

EXAMPLE 4 (looking ahead) Suppose \( \Sigma a_n x^n \) converges for a particular number \( x \). Then for every \( x \) nearer to zero, it converges absolutely. This will be proved and used in Section 10.6 on power series, where it is the most important step in the theory.

EXAMPLE 5 Since \( \Sigma 1/n^2 \) converges, so does \( \Sigma (\cos n)/n^2 \). That second series has irregular signs, but it converges absolutely by comparison with the first series (since \( |\cos n| < 1 \)). Probably \( \Sigma (\tan n)/n^2 \) does not converge, because the tangent does not stay bounded like the cosine.

### Alternating Series

The series \( 1 - \frac{1}{2} + \frac{1}{3} - \frac{1}{4} + \cdots \) converges to \( \ln 2 \). That was stated without proof. This is an example of an alternating series, in which the signs alternate between plus and minus. There is the additional property that the absolute values \( 1, \frac{1}{2}, \frac{1}{3}, \frac{1}{4}, \ldots \) decrease to zero. Those two facts—decrease to zero with alternating signs—guarantee convergence.

\[
\textbf{40H} \quad \text{An alternating series } a_1 - a_2 + a_3 - a_4 - \cdots \text{ converges (at least conditionally, maybe not absolutely) if every } a_{n+1} \leq a_n \text{ and } a_n \to 0.
\]

The best proof is in Figure 10.3. Look at \( a_1 - a_2 + a_3 \). It is below \( a_1 \), because \( a_3 \) (with plus sign) is smaller than \( a_2 \) (with minus sign). The sum of five terms is less than the
10.3 Convergence Tests: All Series

sum of three terms, because \( a_5 \) is smaller than \( a_4 \). These partial sums \( s_1, s_3, s_5, \ldots \) with an odd number of terms are decreasing.

Now look at two terms \( a_1 - a_2 \), then four terms, then six terms. Adding on \( a_3 - a_4 \) increases the sum (because \( a_3 \geq a_4 \)). Similarly \( s_6 \) is greater than \( s_5 \) (because it includes \( a_5 - a_6 \) which is positive). So the sums \( s_2, s_4, s_6, \ldots \) are increasing.

The difference between \( s_{n-1} \) and \( s_n \) is the single number \( \pm a_n \). It is required by 10H to approach zero. Therefore the decreasing sequence \( s_1, s_3, \ldots \) approaches the same limit \( s \) as the increasing sequence \( s_2, s_4, \ldots \). The series converges to \( s \), which always lies between \( s_{n-1} \) and \( s_n \).

This plus-minus pattern is special but important. The positive series \( \Sigma a_n \) may not converge. The alternating series is \( \Sigma(-1)^{n+1}a_n \).

**EXAMPLE 6** The alternating series \( 4 - \frac{1}{2} + \frac{1}{4} - \frac{1}{8} \cdots \) is conditionally convergent (to \( \pi \)). The absolute values decrease to zero. Is this series absolutely convergent? No. With plus signs, \( 4(1 + \frac{1}{2} + \frac{1}{3} + \cdots) \) diverges like the harmonic series.

**EXAMPLE 7** The alternating series \( 1 - 1 + 1 - 1 + \cdots \) is not convergent at all. Which requirement in 10H is not met? The partial sums \( s_1, s_3, s_5, \ldots \) all equal 1 and \( s_2, s_4, s_6, \ldots \) all equal 0—but they don’t approach the same limit \( s \).

**MULTIPLYING AND REARRANGING SERIES**

In Section 10.1 we added and subtracted and multiplied series. Certainly addition and subtraction are safe. If one series has partial sums \( s_n \to s \) and the other has partial sums \( t_n \to t \), then addition gives partial sums \( s_n + t_n \to s + t \). But multiplication is more dangerous, because the order of the multiplication can make a difference. More exactly, the order of terms is important when the series are conditionally convergent. For absolutely convergent series, the order makes no difference. We can rearrange their terms and multiply them in any order, and the sum and product comes out right:

**10I** Suppose \( \Sigma a_n \) converges absolutely. If \( A_1, A_2, \ldots \) is any reordering of the \( a_n \)'s, then \( \Sigma A_n \) also converges absolutely.

**10J** Suppose \( \Sigma a_n = s \) and \( \Sigma b_n = t \) converge absolutely. Then the infinitely many terms \( a_nb_n \) in their product add (in any order) to \( st \).

Rather than proving 10I and 10J, we show what happens when there is only conditional convergence. Our favorite is \( 1 - \frac{1}{2} + \frac{1}{3} - \frac{1}{4} + \cdots \), converging conditionally to
In 2. By rearranging, it will converge conditionally to anything! Suppose the desired
sum is 1000. Take positive terms \(1 + \frac{1}{5} + \cdots\) until they pass 1000. Then add negative
terms \(-\frac{1}{5} - \frac{1}{4} - \cdots\) until the subtotal drops below 1000. Then new positive terms
bring it above 1000, and so on. All terms are eventually used, since at least one new
term is needed at each step. The limit is \(s = 1000\).

We also get strange products, when series fail to converge absolutely:
\[
(1 - \frac{1}{\sqrt{2}} + \frac{1}{\sqrt{3}} \cdots)(1 - \frac{1}{\sqrt{2}} + \frac{1}{\sqrt{3}} \cdots) = 1 - \left(\frac{1}{\sqrt{2}} + \frac{1}{\sqrt{3}} \cdots\right).
\]
On the left the series converge (conditionally). The alternating terms go to zero. On
the right the series diverges. Its terms in parentheses don't even approach zero, and
the product is completely wrong.

I close by emphasizing that it is absolute convergence that matters. The most
important series are power series \(\Sigma a_n x^n\). Like the geometric series (with all \(a_n = 1\))
there is absolute convergence over an interval of \(x\)'s. They give functions of \(x\), which
is what calculus needs and wants.

We go next to the series for \(e^x\), which is absolutely convergent everywhere. From
the viewpoint of convergence tests it is too easy—the danger is gone. But from the
viewpoint of calculus and its applications, \(e^x\) is unconditionally the best.

**10.3 EXERCISES**

**Read-through questions**

The series \(\Sigma a_n\) is absolutely convergent if the series \(\Sigma |a_n|\) is
convergent. Then the original series \(\Sigma a_n\) is also \(b\). But
the series \(\Sigma a_n\) can converge without converging absolutely.
That is called \(c\) convergence, and the series \(d\) is an example.

For alternating series, the sign of each \(a_{n+1}\) is \(e\) to the
sign of \(a_n\). With the extra conditions that \(f\) and \(g\),
the series converges (at least conditionally). The partial sums
\(s_1, s_2, \ldots\) are \(h\) and the partial sums \(s_2, s_4, \ldots\) are \(i\).
The difference between \(s_n\) and \(s_{n-1}\) is \(j\). Therefore the
two series converge to the same number \(s\). An alternating
series that converges absolutely [conditionally] (not at all) is
\(k\) \([\underline{l}\]) \((m\)\). With absolute [conditional] con-
vergence a reordering (can or cannot?) change the sum.

**Do the series 1–12 converge absolutely or conditionally?**

1 \[\sum (-1)^{n+1} \frac{n}{n+3}\]  
2 \[\sum (-1)^{n-1}\sqrt{n+3}\]
3 \[\sum (-1)^{n+1} \frac{1}{n!}\]  
4 \[\sum (-1)^{n+1} \frac{3^n}{n!}\]
5 \[\sum (-1)^{n+1} 3\sqrt{n} (n+1)\]  
6 \[\sum (-1)^{n+1}\sin^2 n\]
7 \[\sum (-1)^{n+1} \frac{\ln \left(\frac{1}{n}\right)}{n}\]  
8 \[\sum (-1)^{n+1}\sin^2 n / n\]
9 \[\sum (-1)^{n+1} n^2 (1 + n^4)\]  
10 \[\sum (-1)^{n+1} 2^{1/n}\]
11 \[\sum (-1)^{n+1} n^{1/n}\]  
12 \[\sum (-1)^{n+1} (1 - n^{1/n})\]

13 Suppose \(\Sigma a_n\) converges absolutely. Explain why keeping
the positive \(a\)'s gives another convergent series.

14 Can \(\Sigma a_n\) converge absolutely if all \(a_n\) are negative?

15 Show that the alternating series \(1 - \frac{1}{3} + \frac{1}{5} - \frac{1}{7} + \cdots\)
does not converge, by computing the partial sums \(s_2, s_4, \ldots\).
Which requirement of 10H is not met?

16 Show that \(\frac{1}{5} - \frac{1}{3} + \frac{1}{7} - \frac{1}{11} + \cdots\) does not converge. Which
requirement of 10H is not met?

17 (a) For an alternating series with terms decreasing to zero,
why does the sum \(s\) always lie between \(s_{n-1}\) and \(s_n\)?
(b) Is \(s - s_n\) positive or negative if \(s_n\) stops at a positive \(a_n\)?

18 Use Problem 17 to give a bound on the difference between
\(s_n = 1 - \frac{1}{3} + \frac{1}{5} - \frac{1}{7} + \frac{1}{9}\) and the sum \(s = \ln 2\) of the infinite
series.

19 Find the sum \(1 - \frac{1}{2!} + \frac{1}{3!} - \frac{1}{4!} + \cdots\). The partial sum \(s_4\)
is (above \(s\))(below \(s\)) by less than \(\underline{2}\).

20 Give a bound on the difference between \(s_{100} = \frac{1}{1^2} - \frac{1}{2^2} + \frac{1}{3^2} - \frac{1}{100^2}\)
and \(s = \sum (-1)^{n+1} / n^2\).

21 Starting from \(\frac{1}{1^2} + \frac{1}{2^2} + \frac{1}{3^2} + \cdots = \frac{\pi^2}{6}\), with plus signs,
show that the alternating series in Problem 20 has \(s = \pi^2 / 12\).

22 Does the alternating series in 20 or the positive series in
21 give \(\pi^2\) more quickly? Compare \(1/101^2 - 1/102^2 + \cdots\) with
\(1/101^2 + 1/102^2 + \cdots\).
23 If \( \Sigma a_n \) does not converge show that \( \Sigma |a_n| \) does not converge.

24 Find conditions which guarantee that \( a_1 - a_2 + a_3 - a_4 + \ldots \) will converge (negative term follows two positive terms).

25 If the terms of \( \ln 2 = 1 - \frac{1}{2} + \frac{1}{3} - \frac{1}{4} + \ldots \) are rearranged into \( 1 - \frac{1}{4} + \frac{1}{5} + \frac{1}{6} - \frac{1}{7} + \ldots \), show that this series now adds to \( \frac{1}{4} \ln 2 \). (Combine each positive term with the following negative term.)

26 Show that the series \( 1 + \frac{1}{2} - \frac{1}{3} - \frac{1}{4} + \frac{1}{5} + \frac{1}{6} - \frac{1}{7} + \ldots \) converges to \( \frac{1}{2} \ln 2 \).

27 What is the sum of \( 1 + \frac{1}{2} - \frac{1}{3} + \frac{1}{4} + \frac{1}{5} - \frac{1}{6} + \ldots \)?

28 Find conditions which guarantee that \( a_1 + a_2 - a_3 + a_4 - a_5 + \ldots \) will converge (negative term follows two positive terms).

29 (a) Prove that this alternating series converges:
\[
1 - \int_1^\infty \frac{dx}{x^2} + 1 - \int_2^\infty \frac{dx}{x^3} + 1 - \int_3^\infty \frac{dx}{x^4} + \ldots
\]
(b) Show that its sum is Euler's constant \( \gamma \).

30 Prove that this series converges. Its sum is \( \pi/2 \):
\[
\int_0^\infty \frac{\sin x}{x} dx + \int_0^\infty \frac{\sin x}{x} \frac{dx}{x} + \ldots = \int_0^\infty \frac{\sin x}{x} dx.
\]

31 The cosine of \( \theta = 1 \) radian is \( 1 - \frac{1}{2!} + \frac{1}{4!} - \ldots \). Compute \( \cos 1 \) to five correct decimals (how many terms?).

32 The sine of \( \theta = \pi \) radians is \( \pi - \frac{\pi^3}{3!} + \frac{\pi^5}{5!} - \ldots \). Compute \( \sin \pi \) to eight correct decimals (how many terms?).

33 If \( \Sigma a_n^2 \) and \( \Sigma b_n^2 \) are convergent show that \( \Sigma a_n b_n \) is absolutely convergent.

Hint: \((a + b)^2 \geq 0 \) yields \( 2|ab| \leq a^2 + b^2 \).

34 Verify the Schwarz inequality \( (\Sigma a_n b_n)^2 \leq (\Sigma a_n^2)(\Sigma b_n^2) \) if \( a_n = \left(\frac{1}{n}\right)^p \) and \( b_n = \left(\frac{1}{n}\right)^q \).

35 Under what condition does \( \Sigma (a_{n+1} - a_n) \) converge and what is its sum?

36 For a conditionally convergent series, explain how the terms could be rearranged so that the sum is \( + \infty \). All terms must eventually be included, even negative terms.

37 Describe the terms in the product \((1 + \frac{1}{2} + \frac{1}{3} + \ldots)(1 + \frac{1}{2} + \frac{1}{4} + \ldots)\) and find their sum.

38 True or false:
(a) Every alternating series converges.
(b) \( \Sigma a_n \) converges conditionally if \( \Sigma |a_n| \) diverges.
(c) A convergent series with positive terms is absolutely convergent.
(d) If \( \Sigma a_n \) and \( \Sigma b_n \) both converge, so does \( \Sigma (a_n + b_n) \).

39 Every number \( x \) between 0 and 2 equals \( 1 + \frac{1}{2} + \frac{1}{3} + \ldots \) with suitable terms deleted. Why?

40 Every number \( s \) between \(-1 \) and \( 1 \) equals \( \pm \frac{1}{1} \pm \frac{1}{2} \pm \frac{1}{3} \pm \ldots \) with a suitable choice of signs. (Add \( 1 + \frac{1}{2} + \frac{1}{3} + \ldots \) to get Problem 39.) Which signs give \( s = -1 \) and \( s = 0 \) and \( s = \frac{1}{2} \)?

41 Show that no choice of signs will make \( 1 + \frac{1}{2} + \frac{1}{3} + \ldots \) equal to zero.

42 The sums in Problem 41 form a Cantor set centered at zero. What is the smallest positive number in the set? Choose signs to show that \( \frac{1}{3} \) is in the set.

43 Show that the tangent of \( \theta = \frac{1}{2}(\pi - 1) \) is \( \sin 1/(1 - \cos 1) \). This is the imaginary part of \( s = -\ln(1 - e^i) \). From \( s = \Sigma e^{inx} / n \) deduce the remarkable sum \( \Sigma (\sin n)/n = \frac{1}{2}(\pi - 1) \).

44 Suppose \( \Sigma a_n \) converges and \( |x| < 1 \). Show that \( \Sigma a_n x^n \) converges absolutely.

### 10.4 The Taylor Series for \( e^x \), \( \sin x \), and \( \cos x \)

This section goes back from numbers to functions. Instead of \( \Sigma a_n = s \) it deals with \( \Sigma a_n x^n = f(x) \). The sum is a function of \( x \). The geometric series has all \( a_n = 1 \) (including \( a_0 \), the constant term) and its sum is \( f(x) = 1/(1-x) \). The derivatives of \( 1 + x + x^2 + \ldots \) match the derivatives of \( f \). Now we choose the \( a_n \) differently, to match a different function.

The new function is \( e^x \). All its derivatives are \( e^x \). At \( x = 0 \), this function and its derivatives equal 1. To match these 1's, we move factorials into the denominators.
Term by term the series is
\[ e^x = 1 + \frac{x}{1!} + \frac{x^2}{2!} + \frac{x^3}{3!} + \cdots. \] (1)

\(x^n/n!\) has the correct \(n\)th derivative \((= 1)\). From the derivatives at \(x = 0\), we have built back the function! At \(x = 1\) the right side is \(1 + \frac{1}{2} + \frac{1}{3} + \cdots\) and the left side is \(e = 2.71828\ldots\). At \(x = -1\) the series gives \(1 - \frac{1}{2} - \frac{1}{3} - \cdots\), which is \(e^{-1}\).

The same term-by-term idea works for differential equations, as follows.

**EXAMPLE 1** Solve \(dy/dx = -y\) starting from \(y = 1\) at \(x = 0\).

Solution The zeroth derivative at \(x = 0\) is the function itself, \(y = 1\). Then the equation \(y' = -y\) gives \(y' = 1\) and \(y'' = -y' = 1\). The alternating derivatives \(1, -1, 1, -1, \ldots\) are matched by the alternating series for \(e^{-x}\):
\[ y = 1 - x + \frac{x^2}{2} - \frac{x^3}{6} + \cdots = e^{-x} \] (the correct solution to \(y' = -y\)).

**EXAMPLE 2** Solve \(d^2y/dx^2 = -y\) starting from \(y = 1\) and \(y' = 0\) (the answer is \(\cos x\)).

Solution The equation gives \(y'' = -1\) again at \(x = 0\). The derivative of the equation gives \(y''' = -y' = 0\). Then \(y''' = -y'' = 1\). The even derivatives are alternately +1 and -1, the odd derivatives are zero. This is matched by a series of even powers, which constructs \(\cos x\):
\[ y = 1 - \frac{1}{2!}x^2 + \frac{1}{4!}x^4 - \frac{1}{6!}x^6 + \cdots = \cos x. \]

The first terms \(1 - \frac{1}{2}x^2\) came earlier in the book. Now we have the whole alternating series. It converges absolutely for all \(x\), by comparison with the series for \(e^x\) (odd powers are dropped). The partial sums in Figure 10.4 reach further and further before they lose touch with \(\cos x\).

If we wanted plus signs instead of plus-minus, we could average \(e^x\) and \(e^{-x}\). The differential equation for \(\cosh x\) is \(d^2y/dx^2 = +y\), to give plus signs:
\[ \frac{1}{2}(e^x + e^{-x}) = 1 + \frac{1}{2!}x^2 + \frac{1}{4!}x^4 + \frac{1}{6!}x^6 + \cdots \] (which is \(\cosh x\)).

**TAYLOR SERIES**

The idea of matching derivatives by powers is becoming central to this chapter. The derivatives are given at a basepoint (say \(x = 0\)). They are numbers \(f(0), f'(0), \ldots\). The derivative \(f^{(n)}(0)\) will be the \(n\)th derivative of \(a_nx^n\), if we choose \(a_n\) to be \(f^{(n)}(0)/n!\).
Then the series \( \sum a_n x^n \) has the same derivatives at the basepoint as \( f(x) \):

The **Taylor series** that matches \( f(x) \) and all its derivatives at \( x = 0 \) is

\[
f(0) + f'(0)x + \frac{1}{2} f''(0)x^2 + \frac{1}{6} f'''(0)x^3 + \cdots = \sum_{n=0}^{\infty} \frac{f^{(n)}(0)}{n!} x^n.
\]

The first terms give the linear and quadratic approximations that we know well. The \( x^3 \) term was mentioned earlier (but not used). Now we have all the terms—an "infinite approximation" that is intended to equal \( f(x) \).

Two things are needed. First, the series must converge. Second, the function must do what the series predicts, away from \( x = 0 \). Those are true for \( e^x \) and \( \cos x \) and \( \sin x \); the series equals the function. We proceed on that basis.

The Taylor series with special basepoint \( x = 0 \) is also called the "Maclaurin series."

**EXAMPLE 3** Find the Taylor series for \( f(x) = \sin x \) around \( x = 0 \).

**Solution** The numbers \( f^{(n)}(0) \) are the values of \( f = \sin x, f' = \cos x, f'' = -\sin x, \ldots \) at \( x = 0 \). Those values are 0, 1, 0, -1, 0, 1, .... All even derivatives are zero. To find the coefficients in the Taylor series, divide by the factorials:

\[
\sin x = x - \frac{1}{6}x^3 + \frac{1}{120}x^5 - \cdots.
\]  

**EXAMPLE 4** Find the Taylor series for \( f(x) = (1 + x)^5 \) around \( x = 0 \).

**Solution** This function starts at \( f(0) = 1 \). Its derivative is \( 5(1 + x)^4 \), so \( f'(0) = 5 \). The second derivative is \( 5 \cdot 4 \cdot (1 + x)^3 \), so \( f''(0) = 5 \cdot 4 \). The next three derivatives are \( 5 \cdot 4 \cdot 3, 5 \cdot 4 \cdot 3 \cdot 2, 5 \cdot 4 \cdot 3 \cdot 2 \cdot 1 \). After that all derivatives are zero. Therefore the Taylor series stops after the \( x^5 \) term:

\[
1 + 5x + \frac{5 \cdot 4}{2!} x^2 + \frac{5 \cdot 4 \cdot 3}{3!} x^3 + \frac{5 \cdot 4 \cdot 3 \cdot 2}{4!} x^4 + \frac{5 \cdot 4 \cdot 3 \cdot 2 \cdot 1}{5!} x^5.
\]

You may recognize 1, 5, 10, 10, 5, 1. They are the binomial coefficients, which appear in Pascal's triangle (Section 2.2). By matching derivatives, we see why 0!, 1!, 2!, ... are needed in the denominators.

There is no doubt that \( x = 0 \) is the nicest basepoint. But Taylor series can be constructed around other points \( x = a \). The principle is the same—match derivatives by powers—but now the powers to use are \((x - a)^n\). The derivatives \( f^{(n)}(a) \) are computed at the new basepoint \( x = a \).

**The Taylor series begins with** \( f(a) + f'(a)(x - a) \). This is the tangent approximation at \( x = a \). The whole "infinite approximation" is centered at \( a \)—at that point it has the same derivatives as \( f(x) \).

The **Taylor series** for \( f(x) \) around the basepoint \( x = a \) is

\[
f(x) = f(a) + f'(a)(x - a) + \frac{1}{2} f''(a)(x - a)^2 + \cdots = \sum_{n=0}^{\infty} \frac{f^{(n)}(a)}{n!} (x - a)^n.
\]

**EXAMPLE 5** Find the Taylor series for \( f(x) = (1 + x)^5 \) around \( x = a = 1 \).

**Solution** At \( x = 1 \), the function is \((1 + 1)^5 = 32 \). Its first derivative \( 5(1 + x)^4 \) is \( 5 \cdot 16 = 80 \). We compute the \( n \)th derivative, divide by \( n! \), and multiply by \((x - 1)^n\):

\[
32 + 80(x - 1) + 80(x - 1)^2 + 40(x - 1)^3 + 10(x - 1)^4 + (x - 1)^5.
\]
That Taylor series (which stops at \( n = 5 \)) should agree with \((1 + x)^5\). It does. We could rewrite \(1 + x\) as \(2 + (x - 1)\), and take its fifth power directly. Then 32, 16, 8, 4, 2, 1 will multiply the usual coefficients 1, 5, 10, 10, 5, 1 to give our Taylor coefficients 32, 80, 80, 40, 10, 1. The series stops as it will stop for any polynomial—because the high derivatives are zero.

**EXAMPLE 6** Find the Taylor series for \( f(x) = e^x \) around the basepoint \( x = 1 \).

**Solution** At \( x = 1 \) the function and all its derivatives equal \( e \). Therefore the Taylor series has that constant factor (note the powers of \( x - 1 \), not \( x \)):

\[
e^x = e + e(x - 1) + \frac{e}{2!}(x - 1)^2 + \frac{e}{3!}(x - 1)^3 + \cdots. (6)
\]

### DEFINING THE FUNCTION BY ITS SERIES

Usually, we define \( \sin x \) and \( \cos x \) from the sides of a triangle. But we could start instead with the series. Define \( \sin x \) by equation (2). The logic goes backward, but it is still correct:

First, prove that the series converges.
Second, prove properties like \((\sin x)' = \cos x\).
Third, connect the definitions by series to the sides of a triangle.

We don't plan to do all this. The usual definition was good enough. But note first: There is no problem with convergence. The series for \( \sin x \) and \( \cos x \) and \( e^x \) all have terms \( \pm x^n/n! \). The factorials make the series converge for all \( x \). The general rule for \( e^x \) times \( e^y \) can be based on the series. Equation (6) is typical: \( e \) is multiplied by powers of \((x - 1)\). Those powers add to \( e^{x-1} \). So the series proves that \( e^x = ee^{x-1} \).

That is just one example of the multiplication \((e^x)(e^y) = e^{x+y}\):

\[
\left(1 + x + \frac{x^2}{2} + \ldots \right) \left(1 + y + \frac{y^2}{2} + \ldots \right) = 1 + x + y + \frac{x^2}{2} + xy + \frac{y^2}{2} + \ldots. (7)
\]

Term by term, multiplication gives the series for \( e^{x+y} \). Term by term, differentiating the series for \( e^x \) gives \( e^x \). Term by term, the derivative of \( \sin x \) is \( \cos x \):

\[
\frac{d}{dx} \left(x - \frac{x^3}{3!} + \frac{x^5}{5!} - \ldots \right) = 1 - \frac{x^2}{2!} + \frac{x^4}{4!} - \ldots. (8)
\]

We don't need the famous limit \((\sin x)/x \to 1\), by which geometry gave us the derivative. The identities of trigonometry become identities of infinite series. We could even define \( \pi \) as the first positive \( x \) at which \( x - \frac{1}{6}x^3 + \cdots \) equals zero. But it is certainly not obvious that this sine series returns to zero—much less that the point of return is near 3.14.

**The function that will be defined by infinite series is** \( e^{i\theta} \). This is the exponential of the imaginary number \( i\theta \) (a multiple of \( i = \sqrt{-1} \)). The result \( e^{i\theta} \) is a complex number, and our goal is to identify it. (We will be confirming Section 9.4.) The technique is to treat \( i\theta \) like all other numbers, real or complex, and simply put it into the series:

**DEFINITION** \( e^{i\theta} \) is the sum of \( 1 + (i\theta) + \frac{1}{2!}(i\theta)^2 + \frac{1}{3!}(i\theta)^3 + \cdots \). (9)

Now use \( i^2 = -1 \). The even powers are \( i^0 = 1, i^2 = -1, i^4 = 1, \ldots \) We are just multiplying \(-1\) by \(-1\) to get \(1\). The odd powers are \( i^1 = -i, i^3 = -i, \ldots \) There-
fore $e^{i\theta}$ splits into a real part (with no i's) and an imaginary part (multiplying i):

$$e^{i\theta} = \left(1 - \frac{1}{2!}\theta^2 + \frac{1}{4!}\theta^4 - \cdots \right) + i \left(\theta - \frac{1}{3!}\theta^3 + \frac{1}{5!}\theta^5 - \cdots \right).$$  \hfill (10)

You recognize those series. They are $\cos \theta$ and $\sin \theta$. Therefore:

**Euler's formula is** $e^{i\theta} = \cos \theta + i\sin \theta$. **Note that** $e^{i\pi} = 1$.

The real part is $x = \cos \theta$ and the imaginary part is $y = \sin \theta$. Those coordinates pick out the point $e^{i\theta}$ in the “complex plane.” Its distance from the origin $(0, 0)$ is $r = 1$, because $(\cos \theta)^2 + (\sin \theta)^2 = 1$. Its angle is $\theta$, as shown in Figure 10.5. The number $-1$ is $e^{i\pi}$, at the distance $r = 1$ and the angle $\pi$. It is on the real axis to the left of zero. If $e^{i\theta}$ is multiplied by $r = 2$ or $r = \frac{1}{2}$ or any $r \geq 0$, the result is a complex number at a distance $r$ from the origin:

**Complex numbers:** $re^{i\theta} = r(\cos \theta + i\sin \theta) = r\cos \theta + ir\sin \theta = x + iy$.

With $e^{i\theta}$, a negative number has a logarithm. **The logarithm of** $-1$ **is imaginary** (it is $i\pi$, since $e^{i\pi} = -1$). A negative number also has fractional powers. The fourth root of $-1$ is $(-1)^{\frac{1}{4}} = e^{\frac{i\pi}{4}}$. More important for calculus: **The derivative of** $x^{\frac{1}{4}}$ **is** $\frac{1}{4}x^{-\frac{3}{4}}$. That sounds old and familiar, but at $x = -1$ it was never allowed.

**Complex numbers tie up the loose ends left by the limitations of real numbers.**

The formula $e^{i\theta} = \cos \theta + i\sin \theta$ has been called “one of the greatest mysteries of undergraduate mathematicians.” Writers have used desperate methods to avoid infinite series. That proof in (10) may be the clearest (I remember sending it to a prisoner studying calculus) but here is a way to start from $d/dx(e^{ix})$:

A different proof of Euler's formula  Any complex number is $e^{ix} = r(\cos \theta + i\sin \theta)$ for some $r$ and $\theta$. Take the $x$ derivative of both sides, and substitute for $ie^{ix}$:

$$e^{ix} \frac{d}{dx}(\cos \theta + i\sin \theta) = (\cos \theta + i\sin \theta) \frac{dr}{dx} + r(-\sin \theta + i\cos \theta) \frac{d\theta}{dx} = i(r\cos \theta + i\sin \theta).$$

Comparing the real parts and also the imaginary parts, we need $dr/dx$ and $d\theta/dx$.

For the real parts on the right we have $\cos \theta + i\sin \theta = x + iy$. Therefore $r$ is always 1 and $\theta$ is $x$. Substituting into the first sentence of the proof, we have Euler's formula $e^{i\theta} = 1(\cos \theta + i\sin \theta)$.

### Read-through questions

The __a__ series is chosen to match $f(x)$ and all its __b__ at the basepoint. Around $x = 0$ the series begins with $f(0) + \_c\_ x + \_d\_ x^2$. The coefficient of $x^2$ is __e__. For $f(x) = e^x$ this series is __f__. For $f(x) = \cos x$ the series is __g___. For $f(x) = \sin x$ the series is __h___. If the signs were all positive in those series, the functions would be $cosh x$ and __i___. Addition gives $cosh x + sinh x = 1$.

In the Taylor series for $f(x)$ around $x = a$, the coefficient of $(x - a)^n$ is $b_n = \_j\_$. Then $b_n(x - a)^n$ has the same __k__ as $f$ at the basepoint. In the example $f(x) = x^2$, the Taylor coefficients are $b_0 = \_m\_$, $b_1 = \_n\_$, $b_2 = \_o\_$. The series $b_0 + b_1(x - a) + b_2(x - a)^2$ agrees with the original __p___. The series for $e^x$ around $x = a$ has $b_n = \_q\_$. Then the Taylor series reproduces the identity $e^x = \_r\_ \{ \_s\_ \}$.

**1.** Write down the series for $e^{2x}$ and compute all derivatives at $x = 0$. Give a series of numbers that adds to $e^2$.

**2.** Write down the series for $\sin 2x$ and check the third derivative at $x = 0$. Give a series of numbers that adds to $\sin 2\pi = 0$.

### Exercises

10.4 The Taylor Series for $e^x$, $\sin x$, and $\cos x$
In 3-8 find the derivatives of $f(x)$ at $x = 0$ and the Taylor series (powers of $x$) with those derivatives.

3 $f(x) = e^{ix}$  
4 $f(x) = 1/(1 + x)$  
5 $f(x) = 1/(1 - 2x)$  
6 $f(x) = \cosh x$  
7 $f(x) = \ln(1 - x)$  
8 $f(x) = \ln(1 + x)$

Problems 9-14 solve differential equations by series.

9 From the equation $dy/dx = y - 2$ find all the derivatives of $y$ at $x = 0$ starting from $y(0) = 1$. Construct the infinite series for $y$, identify it as a known function, and verify that the function satisfies $y' = y - 2$.

10 Differentiate the equation $y' = cy + s$ (c and s constant) to find all derivatives of $y$ at $x = 0$. If the starting value is $y_0 = 0$, construct the Taylor series for $y$ and identify it with the solution of $y' = cy + s$ in Section 6.3.

11 Find the infinite series that solves $y'' = -y$ starting from $y(0) = 0$ and $y'(0) = 1$ at $x = 0$.

12 Find the infinite series that solves $y' = y$ starting from $y = 1$ at $x = 3$ (use powers of $x - 3$). Identify $y$ as a known function.

13 Find the infinite series (powers of $x$) that solves $y'' - y = 0$ starting from $y(0) = 0$ and $y'(0) = 1$ at $x = 0$.

14 Solve $y'' = y$ by a series with $y' = 1$ and $y = 0$ at $x = 0$ and identify $y$ as a known function.

15 Find the Taylor series for $f(x) = (1 + x)^2$ around $x = \alpha = 0$ and around $x = a = 1$ (powers of $x - 1$). Check that both series add to $(1 + x)^2$.

16 Find all derivatives of $f(x) = x^3$ at $x = a$ and write out the Taylor series around that point. Verify that it adds to $x^3$.

17 What is the series for $(1 - x)^5$ with basepoint $a = 1$?

18 Write down the Taylor series for $f = \cos x$ around $x = 0$ and also for $f = \cos (x - 2\pi)$ around $x = 0$.

In 19-24 compute the derivatives of $f$ and its Taylor series around $x = 1$.

19 $f(x) = 1/x$  
20 $f(x) = 1/(2 - x)$  
21 $f(x) = \ln x$  
22 $f(x) = x^4$  
23 $f(x) = e^{-x}$  
24 $f(x) = e^{2x}$

In 25-33 write down the first three nonzero terms of the Taylor series around $x = 0$, from the series for $e^{x}$, $\cos x$, and $\sin x$.

25 $xe^{2x}$  
26 $\cos \sqrt{x}$  
27 $(1 - \cos x)/x^2$

28 $\sin x/x$  
29 $\int_0^x \sin x/x \, dx$  
30 $\sin x^2$

31 $e^{2x}$  
32 $e^x = e^{inx}$  
33 $e^x \cos x$

*34 For $x < 0$ the derivative of $x^n$ is still $nx^{n-1}$.

$$\frac{d}{dx} (x^n) = \frac{d}{dx} (x^n) = n(x^{n-1}) = nx^{n-1}e^{inx} \frac{d[x]}{dx}.$$  

What is $d[x]/dx$? Rewrite this answer as $nx^{n-1}$.

35 Why doesn't $f(x) = \sqrt{x}$ have a Taylor series around $x = 0$? Find the first two terms around $x = 1$.

36 Find the Taylor series for $2^x$ around $x = 0$.

In 37-44 find the first three terms of the Taylor series around $x = 0$.

37 $f(x) = \tan^{-1}x$  
38 $f(x) = \sin^{-1}x$  
39 $f(x) = \tan x$  
40 $f(x) = \ln(\cos x)$  
41 $f(x) = e^{inx}$  
42 $f(x) = \tanh^{-1}x$  
43 $f(x) = \cos^2x$  
44 $f(x) = \sec^2x$

45 From $e^{i\theta} = \cos \theta + i \sin \theta$ and $e^{-i\theta} = \cos \theta - i \sin \theta$, add and subtract to find $\cos \theta$ and $\sin \theta$.

46 Does $(e^{i\theta})^2$ equal $\cos^2 \theta + i \sin^2 \theta$ or $\cos \theta^2 + i \sin \theta^2$?

47 Find the real and imaginary parts and the 99th power of $e^{ix}$, $e^{inx}$, and $e^{-inx}$.

48 The three cube roots of 1 are 1, $e^{2\pi i/3}$, $e^{4\pi i/3}$.

(a) Find the real and imaginary parts of $e^{2\pi i/3}$.

(b) Explain why $(e^{2\pi i/3})^3 = 1$.

(c) Check this statement in rectangular coordinates.

49 The cube roots of $-1 = e^{i\pi}$ are $e^{i/3}$ and ________ and ________. Find their sum and their product.

50 Find the squares of $2e^{i\pi/3} = 1 + \sqrt{3}i$ and $4e^{2\pi i/3} = 2\sqrt{2} + i\sqrt{2}$ in both polar and rectangular coordinates.

51 Multiply $e^{ix} = \cos s + i \sin s$ times $e^{it} = \cos t + i \sin t$ to find formulas for $\cos(s + t)$ and $\sin(s + t)$.

52 Multiply $e^{ix}$ times $e^{-it}$ to find formulas for $\cos(s - t)$ and $\sin(s - t)$.

53 Find the logarithm of $i$. Then find another logarithm of $i$. (What can you add to the exponent of $e^{ix}$ without changing the result?)

54 (Proof that $e$ is irrational) If $e = p/q$ then

$$N = p\left[\frac{1}{e} - \left(1 - \frac{1}{1!} + \frac{1}{2!} - \ldots + \frac{1}{p!}\right)\right]$$

would be an integer. (Why?) The number in brackets—the distance from the alternating series to its sum $1/e$—is less than the last term which is $1/p!$. Deduce that $|N| < 1$ and reach a contradiction, which proves that $e$ cannot equal $p/q$.

55 Solve $dy/dx = y$ by infinite series starting from $y = 2$ at $x = 0$. 
This section studies the properties of a power series. When the basepoint is zero, the powers are $x^n$. The series is $\sum a_n x^n$. When the basepoint is $x = a$, the powers are $(x - a)^n$. We want to know when and where (and how quickly) the series converges to the underlying function. For $e^x$ and $\cos x$ and $\sin x$ there is convergence for all $x$—but that is certainly not true for $1/(1 - x)$. The convergence is best when the function is smooth.

First I emphasize that power series are not the only series. For many applications they are not the best choice. An alternative is a sum of sines, $f(x) = \sum b_n \sin n x$. That is a "Fourier sine series", which treats all $x$'s equally instead of picking on a basepoint. A Fourier series allows jumps and corners in the graph—it takes the rough with the smooth. By contrast a power series is terrific near its basepoint, and gets worse as you move away. The Taylor coefficients $a_n$ are totally determined at the basepoint—where all derivatives are computed. Remember the rule for Taylor series:

$$a_n = (\text{nth derivative at the basepoint})/n! = f^{(n)}(a)/n!$$

A remarkable fact is the convergence in a symmetric interval around $x = a$.

PROOF OF 10M A power series $\sum a_n x^n$ either converges for all $x$, or it converges only at the basepoint $x = 0$, or else it has a radius of convergence $r$:

$$\sum a_n x^n \text{ converges absolutely if } |x| < r \text{ and diverges if } |x| > r.$$
**EXAMPLE 2** The sine series \( \sin x = x - \frac{x^3}{3!} + \frac{x^5}{5!} - \cdots \) has \( r = \infty \) (it converges everywhere).

The ratio of \( x^{n+2}/(n+2)! \) to \( x^n/n! \) is \( x^2/(n+2)(n+1) \). This approaches \( L = 0 \).

**EXAMPLE 3** The series \( \sum (x - 5)^n/n^2 \) has radius \( r = 1 \) around its basepoint \( a = 5 \).

The ratios between terms approach \( L = x - 5 \). (The fractions \( n^2/(n+1)^2 \) go toward 1.) There is absolute convergence if \( |x - 5| < 1 \). This is the interval \( 4 < x < 6 \), symmetric around the basepoint. This series happens to converge at the endpoints 4 and 6, because of the factor \( 1/n^2 \). That factor decides the delicate question—convergence at the endpoints—but all powers of \( n \) give the same interval of convergence \( 4 < x < 6 \).

**CONVERGENCE TO THE FUNCTION: REMAINDER TERM AND RADIUS \( r \)**

Remember that a Taylor series starts with a function \( f(x) \). The derivatives at the basepoint produce the series. Suppose the series converges: Does it converge to the function? This is a question about the remainder \( R_n(x) = f(x) - s_n(x) \), which is the difference between \( f \) and the partial sum \( s_n = a_0 + \cdots + a_n(x - a)^n \). The remainder \( R_n \) is the error if we stop the series, ending with the \( n \)th derivative term \( a_n(x - a)^n \).

Suppose \( f \) has an \( (n+1) \)st derivative from the basepoint \( a \) out to \( x \). Then for some point \( c \) in between (position not known) the remainder at \( x \) equals

\[
R_n(x) = f(x) - s_n(x) = f^{(n+1)}(c)(x-a)^{n+1}/(n+1)!
\]

The error in stopping at the \( n \)th derivative is controlled by the \( (n+1) \)st derivative.

You will guess, correctly, that the unknown point \( c \) comes from the Mean Value Theorem. For \( n = 1 \) the proof is at the end of Section 3.8. That was the error \( e(x) \) in linear approximation:

\[
R_1(x) = f(x) - f(a) - f'(a)(x-a) = \frac{1}{2} f''(c)(x-a)^2.
\]

For every \( n \), the proof compares \( R_n \) with \( (x-a)^{n+1} \) \( n+1 \)st derivatives are \( f^{(n+1)}(x-a)^n \) and \( (n+1)! \). The generalized Mean Value Theorem says that the ratio of \( R_n \) to \( (x-a)^{n+1} \) equals the ratio of those derivatives, at the right point \( c \). That is equation (2). The details can stay in Section 3.8 and Problem 23, because the main point is what we want. **The error is exactly like the next term** \( a_{n+1}(x-a)^{n+1} \), except that the \( (n+1) \)st derivative is at \( c \) instead of the basepoint \( a \).

**EXAMPLE 4** When \( f \) is \( e^x \), the \( (n+1) \)st derivative is \( e^x \). Therefore the error is

\[
R_n = e^x - \left( 1 + x + \cdots + \frac{x^n}{n!} \right) = e^c \frac{x^{n+1}}{(n+1)!}.
\]

At \( x = 1 \) and \( n = 2 \), the error is \( e - (1 + 1 + \frac{1}{2}) \approx 0.218 \). The right side is \( e^c/6 \). The unknown point is \( c = \ln (0.218 \cdot 6) = 0.27 \). Thus \( c \) lies between the basepoint \( a = 0 \) and the error point \( x = 1 \), as required. The series converges to the function, because \( R_n \to 0 \).

In practice, \( n \) is the number of derivatives to be calculated. We may aim for an error \( |R_n| \) below \( 10^{-6} \). Unfortunately, the high derivative in formula (2) is awkward to estimate (except for \( e^x \)). And high derivatives in formula (1) are difficult to compute. Most real calculations use only a few terms of a Taylor series. For more accuracy we move the basepoint closer, or switch to another series.
There is a direct connection between the function and the convergence radius $r$. A hint came for $f(x) = 1/(1 - x)$. The function blows up at $x = 1$—which also ends the convergence interval for the series. Another hint comes for $f' = 1/x$, if we expand around $x = a = 1$:

$$\frac{1}{x} = \frac{1}{1 - (1 - x)} = 1 + (1 - x) + (1 - x)^2 + \cdots.$$ 

(4)

This geometric series converges for $|1 - x| < 1$. Convergence stops at the end point $x = 0$—exactly where $1/x$ blows up. The failure of the function stops the convergence of the series. But note that $1/(1 + x^2)$, which never seems to fail, also has convergence radius $r = 1$:

$$\frac{1}{1 + x^2} = 1 - x^2 + x^4 - x^6 + \cdots$$ converges only for $|x| < 1$.

When you see the reason, you will know why $r$ is a "radius." There is a circle, and the function fails at the edge of the circle. The circle contains complex numbers as well as real numbers. The imaginary points $i$ and $-i$ are at the edge of the circle. The function fails at those points because $1/(1 + i^2) = \infty$.

Complex numbers are pulling the strings, out of sight. The circle of convergence reaches out to the nearest "singularity" of $f(x)$, real or imaginary or complex. For $1/(1 + x^2)$, the singularities at $i$ and $-i$ make $r = 1$. If we expand around $a = 3$, the distance to $i$ and $-i$ is $r = \sqrt{10}$. If we change to $\ln(1 + x)$, which blows up at $x = -1$, the radius of convergence of $x - \frac{1}{2}x^2 + \frac{1}{3}x^3 - \cdots$ is $r = 1$.

The function fails at those points because $1/(1 + i^2) = \infty$.

Complex numbers are pulling the strings, out of sight. The circle of convergence reaches out to the nearest "singularity" of $f(x)$, real or imaginary or complex. For $1/(1 + x^2)$, the singularities at $i$ and $-i$ make $r = 1$. If we expand around $a = 3$, the distance to $i$ and $-i$ is $r = \sqrt{10}$. If we change to $\ln(1 + x)$, which blows up at $x = -1$, the radius of convergence of $x - \frac{1}{2}x^2 + \frac{1}{3}x^3 - \cdots$ is $r = 1$.

\[ \frac{1}{1 + i^2} = \infty \text{ also at } -i \]

**Fig. 10.6** Convergence radius $r$ is distance from basepoint $a$ to nearest singularity.

**THE BINOMIAL SERIES**

We close this chapter with one more series. It is the Taylor series for $(1 + x)^p$, around the basepoint $x = 0$. A typical power is $p = \frac{1}{2}$, where we want the terms in

$$\sqrt{1 + x} = 1 + \frac{1}{2}x + a_2x^2 + \cdots.$$ 

The slow way is to square both sides, which gives $1 + x + (2a_2 + \frac{1}{2})x^2$ on the right. Since $1 + x$ is on the left, $a_2 = -\frac{1}{2}$ is needed to remove the $x^2$ term. Eventually $a_3$ can be found. The fast way is to match the derivatives of $f = (1 + x)^{1/2}$:

$$f' = \frac{1}{2}(1 + x)^{-1/2} \quad f'' = \left(\frac{1}{2}\right)(-\frac{1}{2})(1 + x)^{-3/2} \quad f''' = \left(\frac{1}{2}\right)(-\frac{1}{2})(-\frac{1}{2})(1 + x)^{-5/2}.$$
At $x = 0$ those derivatives are $\frac{1}{2}, -\frac{1}{8}, \frac{1}{16}$. Dividing by $1!, 2!, 3!$ gives

$$a_0 = \frac{1}{2}, \quad a_1 = -\frac{1}{8}, \quad a_2 = \frac{1}{16}, \quad a_n = \frac{1}{n!} \left( \frac{1}{2} \right) \left( \frac{1}{2} - 1 \right) \cdots \left( \frac{1}{2} - n + 1 \right).$$

These are the **binomial coefficients** when the power is $p = \frac{1}{2}$.

Notice the difference from the binomials in Chapter 2. For those, the power $p$ was a positive integer. The series $(1 + x)^p = 1 + 2x + x^2$ stopped at $x^2$. The coefficients for $p = 2$ were 1, 2, 0, 0, 0, ... For fractional $p$ or negative $p$ those later coefficients are not zero, and we find them from the derivatives of $(1 + x)^p$:

$$(1 + x)^p \quad p(1 + x)^{p-1} \quad (p-1)(1 + x)^{p-2} \quad f^{(n)}(0) = p(p-1)\cdots(p-n+1)(1+x)^{p-n}.$$

Dividing by $0!, 1!, 2!, \ldots, n!$ at $x = 0$, the binomial coefficients are

$$\binom{p}{0} = 1, \quad \binom{p}{1}, \quad \binom{p}{2}, \quad \cdots, \quad \frac{f^{(n)}(0)}{n!} = \frac{p(p-1)\cdots(p-n+1)}{n!}.$$

For $p = n$ that last binomial coefficient is $n!/n! = 1$. It gives the final $x^n$ at the end of $(1 + x)^p$. For other values of $p$, the binomial series never stops. It converges for $|x| < 1$:

$$(1 + x)^p = 1 + px + \binom{p}{2} x^2 + \cdots = \sum_{n=0}^{\infty} \frac{p(p-1)\cdots(p-n+1)}{n!} x^n.$$

When $p = 1, 2, 3, \ldots$ the binomial coefficient $p!/n!(n-p)!$ counts the number of ways to select a group of $n$ friends out of a group of $p$ friends. If you have 20 friends, you can choose 2 of them in $20(19)/2 = 190$ ways.

Suppose $p$ is not a positive integer. What goes wrong with $(1 + x)^p$, to stop the convergence at $|x| = 1$? The failure is at $x = -1$. If $p$ is negative, $(1 + x)^p$ blows up. If $p$ is positive, as in $\sqrt{1 + x}$, the higher derivatives blow up. Only for a positive integer $p = n$ does the convergence radius move out to $r = \infty$. In that case the series for $(1 + x)^n$ stops at $x^n$, and $f$ never fails.

A power series is a function in a new form. It is not a simple form, but sometimes it is the only form. To compute $f$ we have to sum the series. To square $f$ we have to multiply series. But the operations of calculus—derivative and integral—are easier. That explains why power series help to solve differential equations, which are a rich source of new functions. (Numerically the series are not always so good.) I should have said that the derivative and integral are easy for each separate term $a_n x^n$—and fortunately the convergence radius of the whole series is not changed.

If $f(x) = \sum a_n x^n$ has convergence radius $r$, so do its derivative and its integral:

$$\frac{df}{dx} = \sum na_n x^{n-1} \quad \text{and} \quad \int f(x) dx = \sum a_n x^{n+1}/(n+1) \text{ also converge for } |x| < r.$$

**EXAMPLE 5** The series for $1/(1-x)$ and its derivative $1/(1-x)^2$ and its integral $-\ln(1-x)$ all have $r = 1$ (because they all have trouble at $x = 1$). The series are $\sum x^n$ and $\sum nx^{n-1}$ and $\sum x^{n+1}/(n+1)$.

**EXAMPLE 6** We can integrate $e^{x^2}$ (previously impossible) by integrating every term in its series:

$$\int e^{x^2} \, dx = \int \left( 1 + x^2 + \frac{1}{2!} x^4 + \cdots \right) \, dx = x + \frac{x^3}{3} + \frac{1}{2!} \left( \frac{x^5}{5} \right) + \frac{1}{3!} \left( \frac{x^7}{7} \right) + \cdots.$$

This always converges ($r = \infty$). The derivative of $e^{x^2}$ was never a problem.
Read-through questions

If $|x| < |X|$ and $\Sigma a_n X^n$ converges, then the series $\Sigma a_n x^n$ also converges in a $\underline{a}$ interval around the $\underline{e}$. For $\Sigma (2x)^n$ the convergence radius is $r = \underline{d}$. For $\Sigma x^n/n!$ the radius is $r = \underline{e}$. For $\Sigma (x - 3)^n$ there is convergence for $|x - 3| < \underline{f}$. Then $x$ is between $\underline{g}$ and $\underline{h}$.

Starting with $f(x)$, its Taylor series $\Sigma a_n x^n$ has $a_n = \underline{i}$. With basepoint $a$, the coefficient of $(x - a)^n$ is $\underline{j}$. The error after the $x^n$ term is called the $\underline{k}$ remainder $R_n(x)$. It is equal to $\underline{l}$ where the unknown point $c$ is between $\underline{m}$. Thus the error is controlled by the $\underline{n}$ derivative.

The circle of convergence reaches out to the first point where $f(x)$ fails. For $f = 4/(2 - x)$, that point is $x = \underline{o}$. Around the basepoint $a = 5$, the convergence radius would be $r = \underline{p}$. For $\sin x$ and $\cos x$ the radius is $r = \underline{q}$.

The series for $\sqrt{1 + x}$ is the $\underline{r}$ series with $p = \frac{1}{2}$. Its coefficients are $a_n = \underline{s}$. Its convergence radius is $\underline{t}$. Its square is the very short series $1 + x$.

In 1--6 find the Taylor series for $f(x)$ around $x = 0$ and its radius of convergence $r$. At what point does $f(x)$ blow up?

1. $f(x) = 1/(1 - 4x)$
2. $f(x) = 1/(1 - 4x^2)$
3. $f(x) = e^{1-x}$
4. $f(x) = \tan x$ (through $x^3$)
5. $f(x) = \ln(e + x)$
6. $f(x) = 1/(1 + 4x^2)$

Find the interval of convergence and the function in 7--10.

7. $f(x) = \sum_{n=0}^{\infty} \frac{(x - 1)^n}{2^n}$
8. $f(x) = \sum_{n=0}^{\infty} \frac{n(x - a)^{n-1}}{n+1}$
9. $f(x) = \sum_{n=0}^{\infty} \frac{1}{n+1} (x - a)^{n+1}$
10. $f(x) = (x - 2n) - \frac{(x - 2n)^3}{3!} + \cdots$

11. Write down the Taylor series for $(e^x - 1)/x$, based on the series for $e^x$. At $x = 0$ the function is 0/0. Evaluate the series at $x = 0$. Check by l'Hôpital's Rule on $(e^x - 1)/x$.

12. Write down the Taylor series for $xe^x$ around $x = 0$. Integrate and substitute $x = 1$ to find the sum of $1/(n(n + 2))$.

13. If $f(x)$ is an even function, so $f(-x) = f(x)$, what can you say about its Taylor coefficients in $f = \Sigma a_n x^n$?

14. Puzzle out the sums of the following series:
   a. $x + x^2 + x^3 + x^4 + x^5 + \cdots$
   b. $1 + \frac{x^2}{4!} + \frac{x^4}{8!} + \cdots$
   c. $(x - 1) - \frac{1}{3}(x - 1)^2 + \frac{1}{5}(x - 1)^3 - \cdots$

15. From the series for $(1 - \cos x)/x^2$ find the limit as $x \to 0$ faster than l'Hôpital's rule.

16. Construct a power series that converges for $0 < x < 2\pi$.

17--24 are about remainders and 25--36 are about binomials.

17. If the cosine series stops before $x^9/8!$ show from (2) that the remainder $R_n$ is less than $x^9/8!$ Does this also follow because the series is alternating?

18. If the sine series around $x = \pi$ stops after the terms in problem 10, estimate the remainder from equation (2).

19. Estimate by (2) the remainder $R_n = x^{n-1} + x^{n+2} + \cdots$ in the geometric series. Then compute $R_n$ exactly and find the unknown point $c$ for $n = 2$ and $x = \frac{1}{2}$.

20. For $-\ln(1-x) = x + \frac{x^2}{2} + \frac{x^3}{3} + R_3$, use equation (2) to show that $R_3 \leq \frac{1}{6}$ at $x = \frac{1}{2}$.

21. Find $R_n$ in Problem 20 and show that the series converges to the function at $x = \frac{1}{2}$ (prove that $R_n \to 0$).

22. By estimating $R_n$ prove that the Taylor series for $e^x$ around $x = 1$ converges to $e^x$ as $n \to \infty$.

23. (Proof of the remainder formula when $n = 2$)
   a. At $x = a$ find $R_2, R_3, R_5, R_6^2$.
   b. At $x = a$ evaluate $g(x) = (x - a)^3$ and $g'(x), g''(x), g'''(x)$.
   c. What rule gives $R_2(x) - R_2(a) = \frac{R_3(c_1)}{g'(c_1)}$?
   d. $\ln \frac{R_2(c_2) - R_2(a)}{g'(c_2)} = \frac{R_3(c_2)}{g'(c_2)}$ and
   e. $\frac{R_3(c_2) - R_2(a)}{g''(c_2) - g'(a)} = \frac{R_3(c_2)}{g''(c_2)}$ where are $c_1$ and $c_2$ and $c$?
   f. Combine (a-b-c-d) into the remainder formula (2).

24. All derivatives of $f(x) = e^{-1/x^2}$ are zero at $x = 0$, including $f(0) = 0$. What is $f'(1)$? What is the Taylor series around $x = 0$? What is the radius of convergence? Where does the series converge to $f(x)$? For $x = 1$ and $n = 1$ what is the remainder estimate in (2)?

25. (a) Find the first three terms in the binomial series for $1/\sqrt{1 - x^2}$.
   (b) Integrate to find the first three terms in the Taylor series for $\sin^{-1}x$.

26. Show that the binomial coefficients in $1/\sqrt{1 - x} = \Sigma a_n x^n$ are $a_n = 1 \cdot 3 \cdot 5 \cdots (2n - 1)/2n!$.

27. For $p = -1$ and $p = -2$ find nice formulas for the binomial coefficients.

28. Change the dummy variable and add lower limits to make $\Sigma^n p x^{n-1} = \Sigma^n (n+1)x^n$. 

10.5 EXERCISES
29 In \((1 - x)^{-1} = \sum x^n\) the coefficient of \(x^n\) is the number of groups of \(n\) friends that can be formed from 1 friend (not binomial repetition is allowed!). The coefficient is 1 and there is only one group—the same friend \(n\) times.

(a) Describe all groups of \(n\) friends that can be formed from 2 friends. (There are \(n + 1\) groups.)

(b) How many groups of 5 friends can be formed from 3 friends?

30 (a) What is the coefficient of \(x^n\) when \(1 + x + x^2 + \ldots\) multiplies \(1 + x + x^2 + \ldots\)? Write the first three terms.

(b) What is the coefficient of \(x^5\) in \((\sum x^n)^3\)?

31 Show that the binomial series for \(\sqrt{1 + 4x}\) has integer coefficients. (Note that \(x^n\) changes to \((4x)^n\). These coefficients are important in counting trees, paths, parentheses...)

32 In the series for \(1/\sqrt{1 - 4x}\), show that the coefficient of \(x^n\) is \((2n)!/(n!)^2\).

Use the binomial series to compute 33–36 with error less than \(1/1000\).

33 \((15)^{1/4}\)

34 \((100)^{1/13}\)

35 \((1.1)^{1.1}\)

36 \(e^{1/1000}\)

37 From \(\sec x = 1/\sqrt{1 - \cos x}\) find the Taylor series of \(\sec x\) up to \(x^6\). What is the radius of convergence \(r\) (distance to blowup point)?

38 From \(\sec^2 x = 1/\sqrt{1 - \sin^2 x}\) find the Taylor series up to \(x^2\). Check by squaring the secant series in Problem 37. Check by differentiating the tangent series in Problem 39.

39 (Division of series) Find \(\tan x\) by long division of \(\sin x/\cos x\):

\[
\frac{\frac{x - x^3}{6} + \frac{x^5}{120} \ldots}{\left(1 - \frac{x^2}{2} + \frac{x^4}{24} \ldots\right)} = x + \frac{x^3}{3} + \frac{2x^5}{15} + \ldots.
\]

40 (Composition of series) If \(f = a_0 + a_1 x + a_2 x^2 + \cdots\) and \(g = b_1 x + b_2 x^2 + \cdots\) find the \(1, x, x^2\) coefficients of \(f(g(x))\). Test on \(f = 1/(1 + x), g = x/(1 - x)\), with \(f(g(x)) = 1 - x\).

41 (Multiplication of series) From the series for \(\sin x\) and \(1/(1 - x)\) find the first four terms for \(f = \sin x/(1 - x)\).

42 (Inversion of series) If \(f = a_1 x + a_2 x^2 + \cdots\) find coefficients \(b_1, b_2\) in \(g = b_1 x + b_2 x^2 + \cdots\) so that \(f(g(x)) = x\). Compute \(b_1, b_2\) for \(f = e^x - 1, g = f^{-1} = \ln(1 + x)\).

43 From the multiplication \((\sin x)(\sin x)\) or the derivatives of \(f(x) = \sin^2 x\) find the first three terms of the series. Find the first four terms for \(\cos^2 x\) by an easy trick.

44 Somehow find the first six nonzero terms for \(f = (1 - x)/(-x^2)\).

45 Find four terms of the series for \(1/\sqrt{1 - x}\). Then square the series to reach a geometric series.

46 Compute \(\int e^{-x^2} dx\) to 3 decimals by integrating the power series.

47 Compute \(\int \sin^2 t dt\) to 4 decimals by power series.

48 Show that \(\sum x^n/n\) converges at \(x = -1\), even though its derivative \(\sum x^{n-1}\) diverges. How can they have the same convergence radius?

49 Compute \(\lim_{x \to 0} (\sin x - \tan x)/x^3\) from the series.

50 If the \(n\)th root of \(a_n\) approaches \(L > 0\), explain why \(\sum a_n x^n\) has convergence radius \(r = 1/L\).

51 Find the convergence radius \(r\) around basepoints \(a = 0\) and \(a = 1\) from the blowup points of \((1 + \tan x)/(1 + x^2)\).
CHAPTER 9  Polar Coordinates and Complex Numbers
  9.1  Polar Coordinates  348
  9.2  Polar Equations and Graphs  351
  9.3  Slope, Length, and Area for Polar Curves  356
  9.4  Complex Numbers  360

CHAPTER 10  Infinite Series
  10.1  The Geometric Series  368
  10.2  Convergence Tests: Positive Series  374
  10.3  Convergence Tests: All Series  381
  10.4  The Taylor Series for \(e^x, \sin x, \text{ and } \cos x\)  385
  10.5  Power Series  391

CHAPTER 11  Vectors and Matrices
  11.1  Vectors and Dot Products  398
  11.2  Planes and Projections  407
  11.3  Cross Products and Determinants  416
  11.4  Matrices and Linear Equations  425
  11.5  Linear Algebra in Three Dimensions  435

CHAPTER 12  Motion along a Curve
  12.1  The Position Vector  446
  12.2  Plane Motion: Projectiles and Cycloids  453
  12.3  Tangent Vector and Normal Vector  459
  12.4  Polar Coordinates and Planetary Motion  464

CHAPTER 13  Partial Derivatives
  13.1  Surfaces and Level Curves  472
  13.2  Partial Derivatives  475
  13.3  Tangent Planes and Linear Approximations  480
  13.4  Directional Derivatives and Gradients  490
  13.5  The Chain Rule  497
  13.6  Maxima, Minima, and Saddle Points  504
  13.7  Constraints and Lagrange Multipliers  514
CHAPTER 11

Vectors and Matrices

This chapter opens up a new part of calculus. It is multidimensional calculus, because the subject moves into more dimensions. In the first ten chapters, all functions depended on time $t$ or position $x$—but not both. We had $f(t)$ or $y(x)$. The graphs were curves in a plane. There was one independent variable ($x$ or $t$) and one dependent variable ($y$ or $f$). Now we meet functions $f(x, t)$ that depend on both $x$ and $t$. Their graphs are surfaces instead of curves. This brings us to the calculus of several variables.

Start with the surface that represents the function $f(x, t)$ or $f(x, y)$ or $f(x, y, t)$. I emphasize functions, because that is what calculus is about.

**EXAMPLE 1** $f(x, t) = \cos(x - t)$ is a traveling wave (cosine curve in motion).

At $t = 0$ the curve is $f = \cos x$. At a later time, the curve moves to the right (Figure 11.1). At each $t$ we get a cross-section of the whole $x$-$t$ surface. For a wave traveling along a string, the height depends on position as well as time.

A similar function gives a wave going around a stadium. Each person stands up and sits down. Somehow the wave travels.

**EXAMPLE 2** $f(x, y) = 3x + y + 1$ is a sloping roof (fixed in time).

The surface is two-dimensional—you can walk around on it. It is flat because $3x + y + 1$ is a linear function. In the $y$ direction the surface goes up at $45^\circ$. If $y$ increases by 1, so does $f$. That slope is 1. In the $x$ direction the roof is steeper (slope 3). There is a direction in between where the roof is steepest (slope $\sqrt{10}$).

**EXAMPLE 3** $f(x, y, t) = \cos(x - y - t)$ is an ocean surface with traveling waves.

This surface moves. At each time $t$ we have a new $x$-$y$ surface. There are three variables, $x$ and $y$ for position and $t$ for time. I can't draw the function, it needs four dimensions! The base coordinates are $x$, $y$, $t$ and the height is $f$. The alternative is a movie that shows the $x$-$y$ surface changing with $t$.

At time $t = 0$ the ocean surface is given by $\cos(x - y)$. The waves are in straight lines. The line $x - y = 0$ follows a crest because $\cos 0 = 1$. The top of the next wave is on the parallel line $x - y = 2\pi$, because $\cos 2\pi = 1$. Figure 11.1 shows the ocean surface at a fixed time.

The line $x - y = t$ gives the crest at time $t$. The water goes up and down (like people in a stadium). *The wave goes to shore, but the water stays in the ocean.*
11 Vectors and Matrices

Fig. 11.1 Moving cosine with a small optical illusion—the darker bands seem to go from top to bottom as you turn.

Fig. 11.2 Linear functions give planes.

Of course multidimensional calculus is not only for waves. In business, demand is a function of price and date. In engineering, the velocity and temperature depend on position $x$ and time $t$. Biology deals with many variables at once (and statistics is always looking for linear relations like $z = x + 2y$). A serious job lies ahead, to carry derivatives and integrals into more dimensions.

11.1 Vectors and Dot Products

In a plane, every point is described by two numbers. We measure across by $x$ and up by $y$. Starting from the origin we reach the point with coordinates $(x, y)$. I want to describe this movement by a vector—the straight line that starts at $(0,0)$ and ends at $(x, y)$. This vector $v$ has a direction, which goes from $(0,0)$ to $(x, y)$ and not the other way.

In a picture, the vector is shown by an arrow. In algebra, $v$ is given by its two components. For a column vector, write $x$ above $y$:

$$v = \begin{bmatrix} x \\ y \end{bmatrix}$$

(x and $y$ are the components of $v$).

Note that $v$ is printed in boldface; its components $x$ and $y$ are in lightface. The vector $-v$ in the opposite direction changes signs. Adding $v$ to $-v$ gives the zero vector (different from the zero number and also in boldface):

$$-v = \begin{bmatrix} -x \\ -y \end{bmatrix} \quad \text{and} \quad v - v = \begin{bmatrix} x - x \\ y - y \end{bmatrix} = \begin{bmatrix} 0 \\ 0 \end{bmatrix} = 0. \quad (2)$$

Notice how vector addition or subtraction is done separately on the $x$'s and $y$'s:

$$v + w = \begin{bmatrix} 3 \\ 1 \end{bmatrix} + \begin{bmatrix} -1 \\ 2 \end{bmatrix} = \begin{bmatrix} 2 \\ 3 \end{bmatrix}. \quad (3)$$

Another way to indicate a vector is $\vec{v}$. You will recognize vectors without needing arrows.
11.1 Vectors and Dot Products

The vector $v$ has components $v_1 = 3$ and $v_2 = 1$. (I write $v_1$ for the first component and $v_2$ for the second component.) The vector $w$ has $w_1 = -1$ and $w_2 = 2$. To add the vectors, add the components. To draw this addition, place the start of $w$ at the end of $v$. Figure 11.3 shows how $w$ starts where $v$ ends.

VECTORS WITHOUT COORDINATES

In that head-to-tail addition of $v + w$, we did something new. The vector $w$ was moved away from the origin. Its length and direction were not changed! The new arrow is parallel to the old arrow—only the starting point is different. The vector is the same as before.

A vector can be defined without an origin and without $x$ and $y$ axes. The purpose of axes is to give the components—the separate distances $x$ and $y$. Those numbers are necessary for calculations. But $x$ and $y$ coordinates are not necessary for head-to-tail addition $v + w$, or for stretching to $2v$, or for linear combinations $2v + 3w$. Some applications depend on coordinates, others don't.

 Generally speaking, physics works without axes—it is "coordinate-free." A velocity has direction and magnitude, but it is not tied to a point. A force also has direction and magnitude, but it can act anywhere—not only at the origin. In contrast, a vector that gives the prices of five stocks is not floating in space. Each component has a meaning—there are five axes, and we know when prices are zero. After examples from geometry and physics (no axes), we return to vectors with coordinates.

EXAMPLE 1 (Geometry) Take any four-sided figure in space. Connect the midpoints of the four straight sides. Remarkable fact: Those four midpoints lie in the same plane.

More than that, they form a parallelogram.

Frankly, this is amazing. Figure 11.4a cannot do justice to the problem, because it is printed on a flat page. Imagine the vectors $A$ and $D$ coming upward. $B$ and $C$ go down at different angles. Notice how easily we indicate the four sides as vectors, not caring about axes or origin.

I will prove that $V = W$. That shows that the midpoints form a parallelogram.

What is $V$? It starts halfway along $A$ and ends halfway along $B$. The small triangle at the bottom shows $V = \frac{1}{2}A + \frac{1}{2}B$. This is vector addition—the tail of $\frac{1}{2}B$ is at the head of $\frac{1}{2}A$. Together they equal the shortcut $V$. For the same reason $W = \frac{1}{2}C + \frac{1}{2}D$.

The heart of the proof is to see these relationships.

One step is left. Why is $\frac{1}{2}A + \frac{1}{2}B$ equal to $\frac{1}{2}C + \frac{1}{2}D$? In other words, why is $A + B$ equal to $C + D$? (I multiplied by 2.) When the right question is asked, the answer jumps out. A head-to-tail addition $A + B$ brings us to the point $R$. Also $C + D$ brings us to $R$. The proof comes down to one line:

$$A + B = PR = C + D.$$ 

Then $V = \frac{1}{2}A + \frac{1}{2}B$ equals $W = \frac{1}{2}C + \frac{1}{2}D$. 

![Figure 11.3](image-url) Parallelogram for $v + w$, stretching for $2v$, signs reversed for $-v$. 

The vector $v$ has components $v_1 = 3$ and $v_2 = 1$. (I write $v_1$ for the first component and $v_2$ for the second component.) The vector $w$ has $w_1 = -1$ and $w_2 = 2$. To add the vectors, add the components. To draw this addition, place the start of $w$ at the end of $v$. Figure 11.3 shows how $w$ starts where $v$ ends.
**EXAMPLE 2** (Also geometry) In any triangle, draw lines from the corners to the midpoints of the opposite sides. To prove by vectors: *Those three lines meet at a point.* Problem 38 finds the meeting point in Figure 11.4. Problem 37 says that *the three vectors add to zero.*

**EXAMPLE 3** (Medicine) An electrocardiogram shows the sum of many small vectors, the voltages in the wall of the heart. What happens to this sum—the heart vector \( V \)—in two cases that a cardiologist is watching for?

**Case 1.** Part of the heart is dead (infarction).

**Case 2.** Part of the heart is abnormally thick (hypertrophy).

A heart attack kills part of the muscle. A defective valve, or hypertension, overworks it. In case 1 the cells die from the cutoff of blood (loss of oxygen). In case 2 the heart wall can triple in size, from excess pressure. The causes can be chemical or mechanical. The effect we see is electrical.

*The machine is adding small vectors and “projecting” them in twelve directions.* The leads on the arms, left leg, and chest give twelve directions in the body. Each graph shows the component of \( V \) in one of those directions. Three of the projections in the vertical plane, plus lead 2 for front-back—produce the “mean QRS vector” in Figure 11.5. That is the sum \( V \) when the ventricles start to contract. The left ventricle is larger, so the heart vector normally points down and to the left.

**Fig. 11.4** Four midpoints form a parallelogram (\( V = W \)). Three medians meet at \( P \).

**Fig. 11.5** \( V \) is a sum of small voltage vectors, at the moment of depolarization.
11.1 Vectors and Dot Products

Fig. 11.6 Changes in $V$ show dead muscle and overworked muscle.

We come soon to projections, but here the question is about $V$ itself. How does the ECG identify the problem?

**Case 1: Heart attack** The dead cells make no contribution to the electrical potential. Some small vectors are missing. Therefore the sum $V$ turns away from the infarcted part.

**Case 2: Hypertrophy** The overwork increases the contribution to the potential. Some vectors are larger than normal. Therefore $V$ turns toward the thickened part.

When $V$ points in an abnormal direction, the ECG graphs locate the problem. The $P, Q, R, S, T$ waves on separate graphs can all indicate hypertrophy, in different regions of the heart. Infarctions generally occur in the left ventricle, which needs the greatest blood supply. When the supply of oxygen is cut back, that ventricle feels it first. The result can be a heart attack (= myocardial infarction = coronary occlusion). Section 11.2 shows how the projections on the ECG point to the location.

First come the basic facts about vectors—components, lengths, and dot products.

**COORDINATE VECTORS AND LENGTH**

To compute with vectors we need axes and coordinates. The picture of the heart is "coordinate-free," but calculations require numbers. A vector is known by its components. The unit vectors along the axes are $i$ and $j$ in the plane and $i, j, k$ in space:

$$
i = \begin{bmatrix} 1 \\ 0 \end{bmatrix}, \quad j = \begin{bmatrix} 0 \\ 1 \end{bmatrix} \quad \text{in 2D}; \quad i = \begin{bmatrix} 1 \\ 0 \end{bmatrix}, \quad j = \begin{bmatrix} 0 \\ 1 \end{bmatrix}, \quad k = \begin{bmatrix} 0 \\ 0 \end{bmatrix} \quad \text{in 3D}.$$  

Notice how easily we moved into three dimensions! The only change is that vectors have three components. The combinations of $i$ and $j$ (or $i, j, k$) produce all vectors $v$ in the plane (and all vectors $V$ in space):

$$v = 3i + j = \begin{bmatrix} 3 \\ 1 \end{bmatrix}, \quad V = i + 2j - 2k = \begin{bmatrix} 1 \\ 2 \\ -2 \end{bmatrix}.$$
Those vectors are also written \( \mathbf{v} = (3, 1) \) and \( \mathbf{V} = (1, 2, -2) \). The components of the vector are also the coordinates of a point. (The vector goes from the origin to the point.) This relation between point and vector is so close that we allow them the same notation: \( P = (x, y, z) \) and \( \mathbf{v} = (x, y, z) = xi + yj + zk \).

The sum \( \mathbf{v} + \mathbf{V} \) is totally meaningless. Those vectors live in different dimensions.

From the components we find the length. The length of \((3, 1)\) is \( \sqrt{3^2 + 1^2} = \sqrt{10} \). This comes directly from a right triangle. In three dimensions, \( \mathbf{V} \) has a third component to be squared and added. The length of \( \mathbf{V} = (x, y, z) \) is \( |\mathbf{V}| = \sqrt{x^2 + y^2 + z^2} \).

Vertical bars indicate length, which takes the place of absolute value. The length of \( \mathbf{v} = 3i + j \) is the distance from the point \((0, 0)\) to the point \((3, 1)\):

\[
|\mathbf{v}| = \sqrt{3^2 + 1^2} = \sqrt{10} \quad |\mathbf{V}| = \sqrt{1^2 + 2^2 + (-2)^2} = 3.
\]

A unit vector is a vector of length one. Dividing \( \mathbf{v} \) and \( \mathbf{V} \) by their lengths produces unit vectors in the same directions:

\[
\frac{\mathbf{v}}{|\mathbf{v}|} = \left[ \begin{array}{c} 3/\sqrt{10} \\ 1/\sqrt{10} \end{array} \right] \quad \text{and} \quad \frac{\mathbf{V}}{|\mathbf{V}|} = \left[ \begin{array}{c} 1/3 \\ 2/3 \\ -2/3 \end{array} \right]
\]

Each nonzero vector has a positive length \(|\mathbf{v}|\). The direction of \( \mathbf{v} \) is given by a unit vector \( \mathbf{u} = \mathbf{v}/|\mathbf{v}| \). Then length times direction equals \( \mathbf{v} \).

A unit vector in the plane is determined by its angle \( \theta \) with the \( x \) axis:

\[
\mathbf{u} = \left[ \begin{array}{c} \cos \theta \\ \sin \theta \end{array} \right] = (\cos \theta)\mathbf{i} + (\sin \theta)\mathbf{j} \quad \text{is a unit vector: } |\mathbf{u}|^2 = \cos^2 \theta + \sin^2 \theta = 1.
\]

In 3-space the components of a unit vector are its "direction cosines":

\[
\mathbf{U} = (\cos \alpha)\mathbf{i} + (\cos \beta)\mathbf{j} + (\cos \gamma)\mathbf{k} : \quad \alpha, \beta, \gamma = \text{angles with } x, y, z \text{ axes.}
\]

Then \( \cos^2 \alpha + \cos^2 \beta + \cos^2 \gamma = 1 \). We are doing algebra with numbers while we are doing geometry with vectors. It was the great contribution of Descartes to see how to study algebra and geometry at the same time.

![Coordinate vectors i, j, k. Perpendicular vectors v \cdot w = (6)(1) + (-2)(3) = 0.](image)

**The Dot Product of Two Vectors**

There are two basic operations on vectors. First, vectors are added \((\mathbf{v} + \mathbf{w})\). Second, a vector is multiplied by a scalar \((7\mathbf{v} \text{ or } -2\mathbf{w})\). That leaves a natural question—how do you multiply two vectors? The main part of the answer is—you don't. But there
is an extremely important operation that begins with two vectors and produces a number. It is usually indicated by a dot between the vectors, as in \( \mathbf{v} \cdot \mathbf{w} \), so it is called the *dot product*.

**Definition 1** The dot product multiplies the lengths \(|\mathbf{v}|\) times \(|\mathbf{w}|\) times a cosine:

\[
\mathbf{v} \cdot \mathbf{w} = |\mathbf{v}| |\mathbf{w}| \cos \theta, \quad \theta = \text{angle between } \mathbf{v} \text{ and } \mathbf{w}.
\]

**Example**

\[
\begin{bmatrix} 3 \\ 0 \end{bmatrix} \text{ has length } 3, \quad \begin{bmatrix} 2 \\ 2 \end{bmatrix} \text{ has length } \sqrt{8}, \text{ the angle is } 45^\circ.
\]

The dot product is \(|\mathbf{v}| |\mathbf{w}| \cos \theta = (3)(\sqrt{8})(1/\sqrt{2})\), which simplifies to 6. The square roots in the lengths are "canceled" by square roots in the cosine. For computing \( \mathbf{v} \cdot \mathbf{w} \), a second and much simpler way involves no square roots in the first place.

**Definition 2** The dot product \( \mathbf{v} \cdot \mathbf{w} \) multiplies component by component and adds:

\[
\mathbf{v} \cdot \mathbf{w} = v_1 w_1 + v_2 w_2
\]

**Example**

\[
\begin{bmatrix} 3 \\ 0 \end{bmatrix} \cdot \begin{bmatrix} 2 \\ 2 \end{bmatrix} = (3)(2) + (0)(2) = 6.
\]

The first form \(|\mathbf{v}| |\mathbf{w}| \cos \theta \) is coordinate-free. The second form \( v_1 w_1 + v_2 w_2 \) computes with coordinates. Remark 4 explains why these two forms are equal.

**Remark 1** When \( \mathbf{V} = \mathbf{W} \) the angle is zero but not the cosine! In this case \( \cos \theta = 1 \) and \( \mathbf{V} \cdot \mathbf{V} = |\mathbf{V}|^2 \). The dot product of \( \mathbf{V} \) with itself is the length squared:

\[
\mathbf{V} \cdot \mathbf{V} = |\mathbf{V}|^2 = (V_1^2 + V_2^2 + V_3^2).
\]

**Remark 2** The dot product of \( \mathbf{i} = (1, 0, 0) \) with \( \mathbf{j} = (0, 1, 0) \) is \( \mathbf{i} \cdot \mathbf{j} = 0 \). The axes are perpendicular. Similarly \( \mathbf{i} \cdot \mathbf{k} = 0 \) and \( \mathbf{j} \cdot \mathbf{k} = 0 \). Those are unit vectors: \( \mathbf{i} \cdot \mathbf{i} = \mathbf{j} = \mathbf{j} = \mathbf{k} \cdot \mathbf{k} = 1 \).

**Remark 3** The dot product has three properties that keep the algebra simple:

1. \( \mathbf{V} \cdot \mathbf{W} = \mathbf{W} \cdot \mathbf{V} \)
2. \( (c \mathbf{V}) \cdot \mathbf{W} = c(\mathbf{V} \cdot \mathbf{W}) \)
3. \( (\mathbf{U} + \mathbf{V}) \cdot \mathbf{W} = \mathbf{U} \cdot \mathbf{W} + \mathbf{V} \cdot \mathbf{W} \)

When \( \mathbf{V} \) is doubled \( (c = 2) \) the dot product is doubled. When \( \mathbf{V} \) is split into \( \mathbf{i}, \mathbf{j}, \mathbf{k} \) components, the dot product splits in three pieces. The same applies to \( \mathbf{W} \), since
11 Vectors and Matrices

Fig. 11.8 Length squared = \((V - W) \cdot (V - W)\), from coordinates and the cosine law.

\[ V \cdot W = W \cdot V. \]  The nine dot products of \(i, j, k\) are zeros and ones, and a giant splitting of both \(V\) and \(W\) gives back the correct \(V \cdot W:\)

\[ V \cdot W = V_1 i \cdot W_1 i + V_2 j \cdot W_2 j + V_3 k \cdot W_3 k + six\ zeros = V_1 W_1 + V_2 W_2 + V_3 W_3. \]

**Remark 4**  The two forms of the dot product are equal. This comes from computing \(|V - W|^2\) by coordinates and also by the “law of cosines”:

- with coordinates: \(|V - W|^2 = (V_1 - W_1)^2 + (V_2 - W_2)^2 + (V_3 - W_3)^2\)
- from cosine law: \(|V - W|^2 = |V|^2 + |W|^2 - 2|V||W| \cos \theta.\)

Compare those two lines. Line 1 contains \(V_1^2\) and \(V_2^2\) and \(V_3^2\). Their sum matches \(|V|^2\) in the cosine law. Also \(W_1^2 + W_2^2 + W_3^2\) matches \(|W|^2\). Therefore the terms containing \(-2\) are the same (you can mentally cancel the \(-2\). The definitions agree:

\[-2(V_1 W_1 + V_2 W_2 + V_3 W_3) = -2|V||W| \cos \theta = -2V \cdot W.\]

The cosine law is coordinate-free. It applies to all triangles (even in \(n\) dimensions). Its vector form in Figure 11.8 is \(|V - W|^2 = |V|^2 - 2V \cdot W + |W|^2\). This application to \(V \cdot W\) is its brief moment of glory.

**Remark 5**  The dot product is the best way to compute the cosine of \(\theta:\)

\[ \cos \theta = \frac{V \cdot W}{|V||W|}. \]  (6)

Here are examples of \(V\) and \(W\) with a range of angles from \(0\) to \(\pi\):

- \(i\) and \(3i\) have the same direction \(\cos \theta = 1\)  \(\theta = 0\)
- \(i \cdot (i + j) = 1\) is positive \(\cos \theta = \frac{1}{\sqrt{2}}\)  \(\theta = \pi/4\)
- \(i\) and \(j\) are perpendicular: \(i \cdot j = 0\) \(\cos \theta = 0\)  \(\theta = \pi/2\)
- \(i \cdot (-i + j) = -1\) is negative \(\cos \theta = -\frac{1}{\sqrt{2}}\)  \(\theta = 3\pi/4\)
- \(i\) and \(-3i\) have opposite directions \(\cos \theta = -1\)  \(\theta = \pi\)

**Remark 6**  The Cauchy-Schwarz inequality \(|V \cdot W| \leq |V||W|\) comes from \(|\cos \theta| \leq 1\).

The left side is \(|V||W||\cos \theta|\). It never exceeds the right side \(|V||W|\). This is a key inequality in mathematics, from which so many others follow:

- **Geometric mean** \(\sqrt{xy} \leq \text{arithmetic mean} \ \frac{1}{2}(x + y)\) (true for any \(x \geq 0\) and \(y \geq 0\)).
- **Triangle inequality** \(|V + W| \leq |V| + |W|\) (\(|V|, |W|, |V + W|\) are lengths of sides).

These and other examples are in Problems 39 to 44. The Schwarz inequality \(|V \cdot W| \leq |V||W|\) becomes an equality when \(|\cos \theta| = 1\) and the vectors are \(\underline{\text{\ldots}}\).
11.4 Exercises

### True or false in three dimensions:

1. If both U and V make a 30° angle with W, so does U + V.
2. If they make a 90° angle with W, so does U + V.
3. If they make a 90° angle with W they are perpendicular:
   \[ u \cdot v = 0. \]

4. From \( W = (1, 2, 3) \) subtract a multiple of \( V = (1, 1, 1) \) so that \( W - cV \) is perpendicular to \( V \).
5. Draw \( W \) and \( W - eV \).

13 (a) What is the sum \( V \) of the twelve vectors from the center of a clock to the hours?
(b) If the 4 o'clock vector is removed, find \( V \) for the other eleven vectors.
(c) If the vectors to 1, 2, 3 are cut in half, find \( V \) for the twelve vectors.

14 (a) By removing one or more of the twelve clock vectors, make the length \( |V| \) as large as possible.
(b) Suppose the vectors start from the top instead of the center (the origin is moved to 12 o'clock, so \( v_{12} = \mathbf{0} \)). What is the new sum \( V^* \)?

15 Find the angle \( POQ \) by vector methods if \( P = (1, 1, 0), O = (0, 0, 0), Q = (1, 2, -2). \)

16 (a) Draw the unit vectors \( u_1 = (\cos \theta, \sin \theta) \) and \( u_2 = (\cos \phi, \sin \phi) \). By dot products find the formula for \( \cos(\theta - \phi) \).
(b) Draw the unit vector \( u_3 \) from a 90° rotation of \( u_2 \). By dot products find the formula for \( \sin(\theta + \phi) \).

17 Describe all points \((x, y)\) such that \( v = xi + yj \) satisfies
(a) \( |v| = 2 \)
(b) \( v - i = 2 \)
(c) \( v \cdot i = 2 \)
(d) \( v \cdot i = |v| \)

18 (Important) If \( A \) and \( B \) are non-parallel vectors from the origin, describe
(a) the endpoints of \( tB \) for all numbers \( t \)
(b) the endpoints of \( A + tB \) for all \( t \)
(c) the endpoints of \( sA + tB \) for all \( s \) and \( t \)
(d) the vectors \( v \) that satisfy \( v \cdot A = v \cdot B \)

19 (a) If \( v + 2w = i \) and \( 2v + 3w = j \) find \( v \) and \( w \).
(b) If \( v = i + j \) and \( w = 3i + 4j \) then \( i = \quad \quad v + \quad \quad w \).

20 If \( P = (0, 0) \) and \( R = (0, 1) \) choose \( Q \) so the angle \( PQR \) is 90°. All possible \( Q \)'s lie in a ________.

21 (a) Choose \( d \) so that \( A = 2i + 3j \) is perpendicular to \( B = 9i + dj \).
(b) Find a vector \( C \) perpendicular to \( A = i + j + k \) and \( B = i - k \).
22 If a boat has velocity $V$ with respect to the water and the water has velocity $W$ with respect to the land, then _______. The speed of the boat is not $|V| + |W|$ but _______.

23 Find the angle between the diagonal of a cube and (a) an edge (b) the diagonal of a face (c) another diagonal of the cube. Choose lines that meet.

24 Draw the triangle $PQR$ in Example 1 (the four-sided figure in space). By geometry not vectors, show that $PR$ is twice as long as $V$. Similarly $|PR| = 2|W|$. Also $V$ is parallel to $W$ because both are parallel to _______. So $V = W$ as before.

25 (a) If $A$ and $B$ are unit vectors, show that they make equal angles with $A + B$.

(b) If $A$, $B$, $C$ are unit vectors with $A + B + C = 0$, they form a _______ triangle and the angle between any two is _______.

26 (a) Find perpendicular unit vectors $I$ and $J$ in the plane that are different from $i$ and $j$.

(b) Find perpendicular unit vectors $I$, $J$, $K$ different from $i$, $j$, $k$.

27 If $I$ and $J$ are perpendicular, take their dot products with $A = ai + bj$ to find $a$ and $b$.

28 Suppose $I = (i + j)/\sqrt{2}$ and $J = (i - j)/\sqrt{2}$. Check $I \cdot J = 0$ and write $A = 2i + 3j$ as a combination $ai + bj$. (Best method: use $a$ and $b$ from Problem 27. Alternative: Find $i$ and $j$ from $I$ and $J$ and substitute into $A$.)

29 (a) Find the position vector $OP$ and the velocity vector $PQ$ when the point $P$ moves around the unit circle (see figure) with speed 1. (b) Change to speed 2.

30 The sum $(A \cdot i)^2 + (A \cdot j)^2 + (A \cdot k)^2$ equals _______.

31 In the semicircle find $C$ and $D$ in terms of $A$ and $B$. Prove that $C \cdot D = 0$ (they meet at right angles).

32 The diagonal $PR$ has $|PR|^2 = (A + B) \cdot (A + B) = A \cdot A + A \cdot B + B \cdot A + B \cdot B$. Add $|QS|^2$ from the other diagonal to prove the parallelogram law: $|PR|^2 + |QS|^2 = \text{sum of squares of the four side lengths.}$

33 If $(1, 2, 3)$, $(3, 4, 7)$, and $(2, 1, 2)$ are corners of a parallelogram, find all possible fourth corners.

34 The diagonals of the parallelogram are $A + B$ and _______. If they have the same length, prove that $A \cdot B = 0$ and the region is a _______.

35 The vector from the earth’s center to Seattle is $ai + bj + ck$.

(a) Along the circle at the latitude of Seattle, what two functions of $a$, $b$, $c$ stay constant? $k$ goes to the North Pole.

(b) On the circle at the longitude of Seattle—the meridian—what two functions of $a$, $b$, $c$ stay constant?

(c) Extra credit: Estimate $a$, $b$, $c$ in your present position. The $0^\circ$ meridian through Greenwich has $b = 0$.

36 If $|A + B|^2 = |A|^2 + |B|^2$, prove that $A$ is perpendicular to $B$.

37 In Figure 11.4, the medians go from the corners to the midpoints of the opposite sides. Express $M_1$, $M_2$, $M_3$ in terms of $A$, $B$, $C$. Prove that $M_1 + M_2 + M_3 = 0$. What relation holds between $A$, $B$, $C$?

38 The point $\frac{1}{3}$ of the way along is the same for all three medians. This means that $A + \frac{1}{3}M_3 = \frac{1}{3}M_2 = \frac{1}{3}M_3$. Prove that those three vectors are equal.

39 (a) Verify the Schwarz inequality $|V \cdot W| \leq |V||W|$ for $V = i + 2j + 2k$ and $W = 2i + 2j + k$.

(b) What does the inequality become when $V = (\sqrt{x}, \sqrt{y})$ and $W = (\sqrt{y}, \sqrt{x})$?

40 By choosing the right vector $W$ in the Schwarz inequality, show that $(V + W)^2 < 3(V + W)$.

41 The Schwarz inequality for $ai + bj$ and $ci + dj$ says that $(a^2 + b^2)(c^2 + d^2) \geq (ac + bd)^2$. Multiply out to show that the difference is $\geq 0$.

42 The vectors $A$, $B$, $C$ form a triangle if $A + B + C = 0$. The triangle inequality $|A + B| \leq |A| + |B|$ says that any one side length is less than _______. The proof comes from Schwarz:

$$|A + B|^2 = A \cdot A + 2A \cdot B + B \cdot B \leq |A|^2 + |B|^2 = (|A| + |B|)^2.$$
11.2 Planes and Projections

The most important "curves" are straight lines. The most important functions are linear. Those sentences take us back to the beginning of the book—the graph of \( mx + b \) is a line. The goal now is to move into three dimensions, where graphs are surfaces. Eventually the surfaces will be curved. But calculus starts with the flat surfaces that correspond to straight lines:

What are the most important surfaces? **Planes.**

What are the most important functions? **Still linear.**

The geometrical idea of a plane is turned into algebra, by finding the equation of a plane. Not just a general formula, but the particular equation of a particular plane.

A line is determined by one point \((x_0, y_0)\) and the slope \(m\). The point-slope equation is \( y - y_0 = m(x - x_0) \). That is a linear equation, it is satisfied when \( y = y_0 \) and \( x = x_0 \), and \( dy/dx = m \). For a plane, we start again with a particular point—which is now \((x_0, y_0, z_0)\). But the slope of a plane is not so simple. Many planes climb at a 45° angle—with "slope 1"—and more information is needed.

The direction of a plane is described by a vector \( \mathbf{N} \). The vector is not in the plane, but perpendicular to the plane. In the plane, there are many directions. Perpendicular to the plane, there is only one direction. A vector in that perpendicular direction is a normal vector.

The normal vector \( \mathbf{N} \) can point "up" or "down". The length of \( \mathbf{N} \) is not crucial (we often make it a unit vector and call it \( \mathbf{n} \)). Knowing \( \mathbf{N} \) and the point \( P_0 = (x_0, y_0, z_0) \), we know the plane (Figure 11.9). For its equation we switch to algebra and use the dot product—which is the key to perpendicularity.

\( \mathbf{N} \) is described by its components \((a, b, c)\). In other words \( \mathbf{N} = ai + bj + ck \). This vector is perpendicular to every direction in the plane. A typical direction goes from...
11 Vectors and Matrices

To another point $P = (x, y, z)$ in the plane. The vector from $P_0$ to $P$ has components $(x - x_0, y - y_0, z - z_0)$. This vector lies in the plane, so its dot product with $N$ is zero:

$$\mathbf{N} \cdot (x - x_0, y - y_0, z - z_0) = 0$$

The plane through $P_0$ perpendicular to $N = (a, b, c)$ has the equation

$$(a, b, c) \cdot (x - x_0, y - y_0, z - z_0) = 0$$

or

$$a(x - x_0) + b(y - y_0) + c(z - z_0) = 0. \quad (1)$$

The point $P$ lies on the plane when its coordinates $x, y, z$ satisfy this equation.

**EXAMPLE 1** The plane through $P_0 = (1, 2, 3)$ perpendicular to $N = (1, 1, 1)$ has the equation

$$(x - 1) + (y - 2) + (z - 3) = 0.$$  That can be rewritten as $x + y + z = 6$. Notice three things. First, $P_0$ lies on the plane because $1 + 2 + 3 = 6$. Second, $N = (1, 1, 1)$ can be recognized from the $x, y, z$ coefficients in $x + y + z = 6$. Third, we could change $N$ to $(2, 2, 2)$ and we could change $P_0$ to $(8, 2, -4)$—because $N$ is still perpendicular and $P_0$ is still in the plane: $8 + 2 - 4 = 6$.

The new normal vector $N = (2, 2, 2)$ produces $2(x - 1) + 2(y - 2) + 2(z - 3) = 0$. That can be rewritten as $2x + 2y + 2z = 12$. Same normal direction, same plane.

The new point $P_0 = (8, 2, -4)$ produces $(x - 8) + (y - 2) + (z + 4) = 0$. That is another form of $x + y + z = 6$. All we require is a perpendicular $N$ and a point $P_0$ in the plane.

**EXAMPLE 2** The plane through $(1, 2, 4)$ with the same $N = (1, 1, 1)$ has a different equation: $(x - 1) + (y - 2) + (z - 4) = 0$. This is $x + y + z = 7$ (instead of 6). These planes with 7 and 6 are parallel.

Starting from $a(x - x_0) + b(y - y_0) + c(z - z_0) = 0$, we often move $ax_0 + by_0 + cz_0$ to the right hand side—and call this constant $d$:

$$a(x - x_0) + b(y - y_0) + c(z - z_0) = ax_0 + by_0 + cz_0 = d. \quad (2)$$

A different $d$ gives a parallel plane; $d = 0$ gives a plane through the origin.

**EXAMPLE 3** The plane $x - y + 3z = 0$ goes through the origin $(0, 0, 0)$. The normal vector is read directly from the equation: $N = (1, -1, 3)$. The equation is satisfied by $P_0 = (1, 1, 0)$ and $P = (1, 4, 1)$. Subtraction gives a vector $V = (0, 3, 1) that is in the plane, and $N \cdot V = 0$.

The parallel planes $x - y + 3z = d$ have the same $N$ but different $d$'s. These planes miss the origin because $d$ is not zero ($x = 0, y = 0, z = 0$ on the left side needs $d = 0$ on the right side). Note that $3x - 3y + 9z = -15$ is parallel to both planes. $N$ is changed to $3N$ in Figure 11.9, but its direction is not changed.

**EXAMPLE 4** The angle between two planes is the angle between their normal vectors. The planes $x - y + 3z = 0$ and $3y + z = 0$ are perpendicular, because $(1, -1, 3) \cdot (0, 3, 1) = 0$. The planes $z = 0$ and $y = 0$ are also perpendicular, because $(0, 0, 1) \cdot (0, 1, 0) = 0$. (Those are the xy plane and the xz plane.) The planes $x + y = 0$ and $x + z = 0$ make a 60° angle, because $\cos 60° = (1, 1, 0) \cdot (1, 0, 1)/\sqrt{2}/\sqrt{2} = \frac{1}{2}$.

The cosine of the angle between two planes is $|N_1 \cdot N_2|/|N_1| |N_2|$. See Figure 11.10.
11.2 Planes and Projections

Fig. 11.10 Angle between planes = angle between normals. Parallel and perpendicular to a line. A line in space through $P_0$ and $Q$.

**Remark 1** We gave the “point-slope” equation of a line (using $m$), and the “point-normal” equation of a plane (using $N$). What is the normal vector $N$ to a line?

The vector $V = (1, m)$ is parallel to the line $y = mx + b$. The line goes across by 1 and up by $m$. The perpendicular vector is $N = (-m, 1)$. The dot product $N \cdot V$ is $-m + m = 0$. Then the point-normal equation matches the point-slope equation:

$$-m(x - x_0) + 1(y - y_0) = 0$$

is the same as $y - y_0 = m(x - x_0)$. (3)

**Remark 2** What is the point-slope equation for a plane? The difficulty is that a plane has different slopes in the $x$ and $y$ directions. The function $f(x, y) = m(x - x_0) + M(y - y_0)$ has two derivatives $m$ and $M$.

This remark has to stop. In Chapter 13, “slopes” become “partial derivatives.”

**A LINE IN SPACE**

In three dimensions, a line is not as simple as a plane. A line in space needs two equations. Each equation gives a plane, and the line is the intersection of two planes.

The equations $x + y + z = 3$ and $2x + 3y + z = 6$ determine a line.

Two points on that line are $P_0 = (1, 1, 1)$ and $Q = (3, 0, 0)$. They satisfy both equations so they lie on both planes. Therefore they are on the line of intersection. The direction of that line, subtracting coordinates of $P_0$ from $Q$, is along the vector $V = 2i - j - k$.

The line goes through $P_0 = (1, 1, 1)$ in the direction of $V = 2i - j - k$.

Starting from $(x_0, y_0, z_0) = (1, 1, 1)$, add on any multiple $tV$. Then $x = 1 + 2t$ and $y = 1 - t$ and $z = 1 - t$. Those are the components of the vector equation $P = P_0 + tV$—which produces the line.

Here is the problem. The line needs two equations—or a vector equation with a parameter $t$. Neither form is as simple as $ax + by + cz = d$. Some books push ahead anyway, to give full details about both forms. After trying this approach, I believe that those details should wait. Equations with parameters are the subject of Chapter 12, and a line in space is the first example. Vectors and planes give plenty to do here—especially when a vector is projected onto another vector or a plane.

**PROJECTION OF A VECTOR**

What is the projection of a vector $B$ onto another vector $A$? One part of $B$ goes along $A$—that is the projection. The other part of $B$ is perpendicular to $A$. We now compute these two parts, which are $P$ and $B - P$. 
In geometry, projections involve \( \cos \theta \). In algebra, we use the dot product (which is closely tied to \( \cos \theta \)). In applications, the vector \( \mathbf{B} \) might be a velocity \( \mathbf{V} \) or a force \( \mathbf{F} \):

An airplane flies northeast, and a 100-mile per hour wind blows due east. What is the projection of \( \mathbf{V} = (100, 0) \) in the flight direction \( \mathbf{A} \)?

Gravity makes a ball roll down the surface \( 2x + 2y + z = 0 \). What are the projections of \( \mathbf{F} = (0, 0, -mg) \) in the plane and perpendicular to the plane?

The component of \( \mathbf{V} \) along \( \mathbf{A} \) is the push from the wind (tail wind). The other component of \( \mathbf{V} \) pushes sideways (crosswind). Similarly the force parallel to the surface makes the ball move. Adding the two components brings back \( \mathbf{V} \) or \( \mathbf{F} \).

\[ \text{We now compute the projection of } \mathbf{B} \text{ onto } \mathbf{A}. \]

Call this projection \( \mathbf{P} \). Since its direction is known—\( \mathbf{P} \) is along \( \mathbf{A} \)—we can describe \( \mathbf{P} \) in two ways:

1) \textit{Give the length of } \mathbf{P} \textit{ along } \mathbf{A}.
2) \textit{Give the vector } \mathbf{P} \textit{ as a multiple of } \mathbf{A}.

Figure 11.11b shows the projection \( \mathbf{P} \) and its length. The hypotenuse is \( |\mathbf{B}| \). The length is \(|\mathbf{P}| = |\mathbf{B}| \cos \theta \). The perpendicular component \( \mathbf{B} - \mathbf{P} \) has length \( |\mathbf{B}| \sin \theta \). The cosine is positive for angles less than 90°. The cosine (and \( \mathbf{P} \)) are zero when \( \mathbf{A} \) and \( \mathbf{B} \) are perpendicular. \( |\mathbf{B}| \cos \theta \) is negative for angles greater than 90°, and the projection points along \(-\mathbf{A}\) (the length is \( |\mathbf{B}| \cos \theta \)). Unless the angle is 0° or 30° or 45° or 60° or 90°, we don't want to compute cosines—and we don't have to. The dot product does it automatically:

\[ |\mathbf{A}||\mathbf{B}| \cos \theta = \mathbf{A} \cdot \mathbf{B} \text{ so the length of } \mathbf{P} \text{ along } \mathbf{A} \text{ is } |\mathbf{B}| \cos \theta = \frac{\mathbf{A} \cdot \mathbf{B}}{|\mathbf{A}|}. \quad (4) \]

Notice that the length of \( \mathbf{A} \) cancels out at the end of (4). If \( \mathbf{A} \) is doubled, \( \mathbf{P} \) is unchanged. But if \( \mathbf{B} \) is doubled, the projection is doubled.

What is the vector \( \mathbf{P} \)? Its length along \( \mathbf{A} \) is \( \mathbf{A} \cdot \mathbf{B} / |\mathbf{A}| \). If \( \mathbf{A} \) is a unit vector, then \( |\mathbf{A}| = 1 \) and the projection is \( \mathbf{A} \cdot \mathbf{B} \) times \( \mathbf{A} \). Generally \( \mathbf{A} \) is not a unit vector, until we divide by \( |\mathbf{A}| \). \textit{Here is the projection } \mathbf{P} \textit{ of } \mathbf{B} \textit{ along } \mathbf{A}:

\[ \mathbf{P} = \text{(length of } \mathbf{P} \text{)(unit vector)} = \left( \frac{\mathbf{A} \cdot \mathbf{B}}{|\mathbf{A}|} \right) \left( \frac{\mathbf{A}}{|\mathbf{A}|} \right) = \frac{\mathbf{A} \cdot \mathbf{B}}{|\mathbf{A}|^2} \mathbf{A}. \quad (5) \]
11.2 Planes and Projections

EXAMPLE 5  For the wind velocity \( V = (100, 0) \) and flying direction \( A = (1, 1) \), find \( P \).
Here \( V \) points east, \( A \) points northeast. The projection of \( V \) onto \( A \) is \( P \):

\[
\text{length } |P| = \frac{A \cdot V}{|A|} = \frac{100}{\sqrt{2}} \quad \text{vector } P = \frac{A \cdot V}{|A|^2} A = \frac{100}{2} (1, 1) = (50, 50).
\]

EXAMPLE 6  Project \( F = (0, 0, -mg) \) onto the plane with normal \( N = (2, 2, 1) \).

The projection of \( F \) along \( N \) is not the answer. But compute that first:

\[
\frac{F \cdot N}{|N|} = -\frac{mg}{3} \quad P = \frac{F \cdot N}{|N|^2} N = -\frac{mg}{9} (2, 2, 1).
\]

\( P \) is the component of \( F \) perpendicular to the plane. It does not move the ball. The in-plane component is the difference \( F - P \). Any vector \( B \) has two projections, along \( A \) and perpendicular:

\[
The \text{ projection } P = \frac{A \cdot B}{|A|^2} A \text{ is perpendicular to the remaining component } B - P.
\]

EXAMPLE 7  Express \( B = i - j \) as the sum of a vector \( P \) parallel to \( A = 3i + j \) and a vector \( B - P \) perpendicular to \( A \). Note \( A \cdot B = 2 \).

\[
\text{Solution } P = \frac{A \cdot B}{|A|^2} A = \frac{2}{10} A = \frac{2}{10} i + \frac{6}{10} j. \quad \text{Then } B - P = \frac{4}{10} i - \frac{12}{10} j.
\]

Check: \( P \cdot (B - P) = \left(\frac{6}{10}\right)\left(\frac{2}{10}\right) - \left(\frac{2}{10}\right)\left(\frac{6}{10}\right) = 0 \). These projections of \( B \) are perpendicular.

\[
\text{Pythagoras: } |P|^2 + |B - P|^2 \text{ equals } |B|^2. \text{ Check that too: } 0.4 + 1.6 = 2.0.
\]

Question  When is \( P = 0 \)? Answer  When \( A \) and \( B \) are perpendicular.

EXAMPLE 8  Find the nearest point to the origin on the plane \( x + 2y + 2z = 5 \).

The shortest distance from the origin is along the normal vector \( N \). The vector \( P \) to the nearest point (Figure 11.12) is \( t \) times \( N \), for some unknown number \( t \). We find \( t \) by requiring \( P = tN \) to lie on the plane.

The plane \( x + 2y + 2z = 5 \) has normal vector \( N = (1, 2, 2) \). Therefore \( P = tN = (t, 2t, 2t) \). To lie on the plane, this must satisfy \( x + 2y + 2z = 5 \):

\[
t + 2(2t) + 2(2t) = 5 \quad \text{or } 9t = 5 \quad \text{or } t = \frac{5}{9}.
\]

Then \( P = \frac{5}{9}N = (\frac{5}{9}, \frac{10}{9}, \frac{10}{9}) \). That locates the nearest point. The distance is \( \frac{5}{9}|N| = \frac{5}{9} \).

This example is important enough to memorize, with letters not numbers:

| 11E | On the plane \( ax + by + cz = d \), the nearest point to \((0, 0, 0)\) is \( P = \left(\frac{da}{a^2 + b^2 + c^2}, \frac{db}{a^2 + b^2 + c^2}, \frac{dc}{a^2 + b^2 + c^2}\right) \). The distance is \( \frac{|d|}{\sqrt{a^2 + b^2 + c^2}} \). |

The steps are the same. \( N \) has components \( a, b, c \). The nearest point on the plane is a multiple \( (ta, tb, tc) \). It lies on the plane if \( a(ta) + b(tb) + c(tc) = d \).

Thus \( t = \frac{d}{a^2 + b^2 + c^2} \). The point \((ta, tb, tc) = tN \) is in equation (7). The distance to the plane is \( |tN| = |d|/|N| \).
Vectors and Matrices

11.12 Vector to the nearest point $P$ is a multiple $tN$. The distance is in (7) and (9).

**Question** How far is the plane from an arbitrary point $Q = (x_l, y_l, z_l)$?

**Answer** The vector from $Q$ to $P$ is our multiple $tN$. In vector form $P = Q + tN$. This reaches the plane if $P \cdot N = d$, and again we find $t$:

$$ (Q + tN) \cdot N = d \quad \text{yields} \quad t = (d - Q \cdot N)/||N||^2. \quad (8) $$

This new term $Q \cdot N$ enters the distance from $Q$ to the plane:

$$ \text{distance} = ||tN|| = |d - Q \cdot N|/||N|| = |d - ax_1 - by_1 - cz_1|/\sqrt{a^2 + b^2 + c^2}. \quad (9) $$

When the point is on the plane, that distance is zero—because $ax_1 + by_1 + cz_1 = d$. When $Q$ is $i + 3j + 2k$, the figure shows $Q \cdot N = 11$ and distance $= 2$.

**PROJECTIONS OF THE HEART VECTOR**

An electrocardiogram has leads to your right arm–left arm–left leg. You produce the voltage. The machine amplifies and records the readings. There are also six chest leads, to add a front-back dimension that is monitored across the heart. We will concentrate on the big “Einthoven triangle,” named after the inventor of the ECG.

The graphs show voltage variations plotted against time. The first graph plots the voltage difference between the arms. Lead II connects the left leg to the right arm. Lead III completes the triangle, which has roughly equal sides (especially if you are a little lopsided). So the projections are based on 60° and 120° angles.

The heart vector $V$ is the sum of many small vectors—all moved to the same origin. $V$ is the net effect of action potentials from the cells—small dipoles adding to a single dipole. The pacemaker ($S$–$A$ node) starts the impulse. The atria depolarize to give the P wave on the graphs. This is actually a P loop of the heart vector—the
graphs only show its projections. The impulse reaches the AV node, pauses, and moves quickly through the ventricles. This produces the QRS complex—the large sharp movement on the graph.

The total QRS interval should not exceed $1/10$ second ($2\frac{1}{2}$ spaces on the printout). 

V points first toward the right shoulder. This direction is opposite to the leads, so the tracings go slightly down. That is the Q wave, small and negative. Then the heart vector sweeps toward the left leg. In positions 3 and 4, its projection on lead I (between the arms) is strongly positive. The R wave is this first upward deflection in each lead. Closing the loop, the S wave is negative (best seen in leads I and aVR).

**Question 1** How many graphs from the arms and leg are really independent?

**Answer** Only two! In a plane, the heart vector $V$ has two components. If we know two projections, we can compute the others. (The ECG does that for us.) Different vectors show better in different projections. A mathematician would use $90^\circ$ angles, with an electrode at your throat.

**Question 2** How are the voltages related? What is the aVR lead?

**Answer** Project the heart vector $V$ onto the sides of the triangle:

The lead vectors have $L_1 - L_II + L_{III} = 0$—they form a triangle.

The projections have $V_L - V_{II} + V_{III} = V \cdot L_1 - V \cdot L_II + V \cdot L_{III} = 0$.

The aVR lead is $-\frac{1}{2}L_1 - \frac{3}{2}L_{III}$. It is pure algebra (no wire). By vector addition it points toward the electrode on the right arm. Its length is $\sqrt{3}$ if the other lengths are 2.

Including aVL and aVF to the left arm and foot, there are six leads intersecting at equal angles. Visualize them going out from a single point (the origin in the chest).

**Question 3** If the heart vector is $V = 2i - j$, what voltage differences are recorded?

**Answer** The leads around the triangle have length 2. The machine projects $V$:

- Lead I is the horizontal vector $2i$. So $V \cdot L_1 = 4$.
- Lead II is the $-60^\circ$ vector $i - \sqrt{3}j$. So $V \cdot L_{II} = 2 + \sqrt{3}$.
- Lead III is the $-120^\circ$ vector $-i - \sqrt{3}j$. So $V \cdot L_{III} = -2 + \sqrt{3}$.

The first and third add to the second. The largest R waves are in leads I and II. In aVR the projection of $V$ will be negative (Problem 46), and will be labeled an S wave.
Question 4  What about the potential (not just its differences). Is it zero at the center? 

Answer: It is zero if we say so. The potential contains an arbitrary constant C. (It is like an indefinite integral. Its differences are like definite integrals.) Cardiologists define a “central terminal” where the potential is zero.

The average of V over a loop is the mean heart vector H. This average requires \( \int V \, dt \), by Chapter 5. With no time to integrate, the doctor looks for a lead where the area under the QRS complex is zero. Then the direction of H (the axis) is perpendicular to that lead. There is so much to say about calculus in medicine.

---

11.2 Exercises

Read-through questions

A plane in space is determined by a point \( P_0 = (x_0, y_0, z_0) \) and a \( \vec{a} \) vector \( \vec{N} \) with components \( (a, b, c) \). The point \( P = (x, y, z) \) is on the plane if the dot product of \( \vec{N} \) with \( \vec{a} \) is zero. (That answer was not P!) The equation of this plane is \( a(x - x_0) + b(y - y_0) + c(z - z_0) = 0 \). The equation is also written as \( ax + by + cz = d \), where \( d \) equals \( \vec{N} \cdot \vec{a} \). A parallel plane has the same \( \vec{a} \) and a different \( \vec{N} \). A plane through the origin has \( d = 0 \).

The equation of the plane through \( P_0 = (2, 1, 0) \) perpendicular to \( \vec{N} = (3, 4, 5) \) is \( 3x + 4y + 5z = \) _1_. A second point in the plane is \( P = (0, 0, _k_) \). The vector from \( P_0 \) to \( P \) is \( \vec{a} = \) _1_. and it is \( \vec{N} \). (Check by dot product.) The plane through \( P_0 = (2, 1, 0) \) perpendicular to the \( z \) axis has \( \vec{N} = \) _1_ and equation _1_.

The component of \( \vec{B} \) in the direction of \( \vec{A} \) is _1_, where \( \theta \) is the angle between the vectors. This is \( \vec{A} \cdot \vec{B} \) divided by _1_. The projection vector \( \vec{P} \) is \( |\vec{B}| \cos \theta \times \vec{A} \) times a _1_ vector in the direction of \( \vec{A} \). Then \( \vec{P} = (|\vec{B}| \cos \theta) (\vec{A}/|\vec{A}|) \) simplifies to _1_. When \( \vec{B} \) is doubled, \( \vec{P} \) is _1_. When \( \vec{A} \) is doubled, \( \vec{P} \) is _1_. If \( \vec{B} \) reverses direction then \( \vec{P} \) _1_. If \( \vec{A} \) reverses direction then \( \vec{P} \) _1_.

When \( \vec{B} \) is a velocity vector, \( \vec{P} \) represents the _1_. When \( \vec{B} \) is a force vector, \( \vec{P} \) is _1_. The component of \( \vec{B} \) perpendicular to \( \vec{A} \) equals _1_. The shortest distance from \( (0, 0, 0) \) to the plane \( ax + by + cz = d \) is along the _1_ vector. The distance is _1_ and the closest point on the plane is \( \vec{P} = \) _1_. The distance from \( Q = (x_1, y_1, z_1) \) to the plane is _1_.

Find two points \( P \) and \( P_0 \) on the planes 1–6 and a normal vector \( \vec{N} \). Verify that \( \vec{N} \cdot (\vec{P} - \vec{P}_0) = 0 \).

1 \( x + 2y + 3z = 0 \) 2 \( x + 2y + 3z = 6 \) 3 the \( yz \) plane  
4 the plane through \( (0, 0, 0) \) perpendicular to \( i + j - k \)  
5 the plane through \( (1, 1, 1) \) perpendicular to \( i + j - k \)  
6 the plane through \( (0, 0, 0) \) and \( (1, 0, 0) \) and \( (0, 1, 1) \).

Find an \( x - y - z \) equation for planes 7–10.

7 The plane through \( P_0 = (1, 2, -1) \) perpendicular to \( \vec{N} = i + j \)  
8 The plane through \( P_0 = (1, 2, -1) \) perpendicular to \( \vec{N} = i + 2j - k \)  
9 The plane through \( (1, 0, 1) \) parallel to \( x + 2y + z = 0 \)  
10 The plane through \( (x_0, y_0, z_0) \) parallel to \( x + y + z = 1 \).

11 When is a plane with normal vector \( \vec{N} \) parallel to the vector \( \vec{V} \)? When is it perpendicular to \( \vec{V} \)?

12 (a) If two planes are perpendicular (front wall and side wall), is every line in one plane perpendicular to every line in the other?

(b) If a third plane is perpendicular to the first, it might be (parallel) (perpendicular) (at a 45° angle) to the second.

13 Explain why a plane cannot

(a) contain \( (1, 2, 3) \) and \( (2, 3, 4) \) and be perpendicular to \( \vec{N} = i + j \)  
(b) be perpendicular to \( \vec{N} = i + j \) and parallel to \( \vec{V} = i + k \)  
(c) contain \( (1, 0, 0), (0, 1, 0), (0, 0, 1) \), and \( (1, 1, 1) \)  
(d) contain \( (1, 1, -1) \) if it has \( \vec{N} = i + j - k \) (maybe it can)

(e) go through the origin and have the equation \( ax + by + cz = 1 \).

14 The equation \( 3x + 4y + 7z - t = 0 \) yields a hyperplane in four dimensions. Find its normal vector \( \vec{N} \) and two points \( P, Q \) on the hyperplane. Check \( (\vec{P} - \vec{Q}) \cdot \vec{N} = 0 \).

15 The plane through \( (x, y, z) \) perpendicular to \( ai + bj + ck \) goes through \( (0, 0, 0) \) if _1_. The plane goes through \( (x_0, y_0, z_0) \) if _1_.

16 A curve in three dimensions is the intersection of _1_ surfaces. A line in four dimensions is the intersection of _1_ hyperplanes.

17 (angle between planes) Find the cosine of the angle between \( x + 2y + 2z = 0 \) and \( (a) x + 2z = 0 \) \( (b) x + 2z = 0 \) \( (c) x = 0 \).
18 \( \mathbf{N} \) is perpendicular to a plane and \( V \) is along a line. Draw the angle \( \theta \) between the plane and the line, and explain why \( V \cdot \mathbf{N}/|V||\mathbf{N}| \) is \( \sin \theta \) not \( \cos \theta \). Find the angle between the \( xy \) plane and \( V = i + j + \sqrt{2}k \).

In Problems 19–26 find the projection \( P \) of \( B \) along \( A \). Also find \( |P| \).

19 \( A = (4, 2, 4), B = (1, -1, 0) \)

20 \( A = (1, -1, 0), B = (4, 2, 4) \)

21 \( B \) is unit vector at \( 60^\circ \) angle with \( A \)

22 \( B = \text{vector of length 2 at } 60^\circ \text{ angle with } A \)

23 \( B = -A \quad 24 \ A = i + j, B = i + k \)

25 \( A \) is perpendicular to \( x - y + z = 0 \), \( B = i + j \).

26 \( A \) is perpendicular to \( x - y + z = 5 \), \( B = i + j + 5k \).

27 The force \( F = 3i - 4k \) acts at the point \((1, 2, 2)\). How much force pulls toward the origin? How much force pulls vertically down? Which direction does a mass move under the force \( F \)?

28 The projection of \( B \) along \( A \) is \( P = \ldots \). The projection of \( B \) perpendicular to \( A \) is \( \ldots \). Check the dot product of the two projections.

29 \( P = (x, y, z) \) is on the plane \( ax + by + cz = 5 \) if \( P \cdot \mathbf{N} = |P||\mathbf{N}| \cos \theta = 5 \). Since the largest value of \( \cos \theta \) is 1, the smallest value of \( |P| \) is \( \ldots \). This is the distance between \( \ldots \).

30 If the air speed of a jet is 500 and the wind speed is 50, what information do you need to compute the jet’s speed over land? What is that speed?

31 How far is the plane \( x + y - z = 1 \) from \((0, 0, 0)\) and also from \((1, 1, -1)\)? Find the nearest points.

32 Describe all points at a distance 1 from the plane \( x + 2y + 2z = 3 \).

33 The shortest distance from \( Q = (2, 1, 1) \) to the plane \( x + y + z = 0 \) is along the vector \( \ldots \). The point \( P = Q + t\mathbf{N} = (2 + t, 1 + t, 1 + t) \) lies on the plane if \( t = \ldots \). Then \( P = \ldots \) and the shortest distance is \( \ldots \). (This distance is not \( |P| \).)

34 The plane through \((1, 1, 1)\) perpendicular to \( \mathbf{N} = i + 2j + 2k \) is a distance \( \ldots \) from \((0, 0, 0)\).

35 \( \text{Distance between planes} \) \( 2x - 2y + z = 1 \) is parallel to \( 2x - 2y + z = 3 \) because \( \ldots \). Choose a vector \( Q \) on the first plane and find \( t \) so that \( Q + t\mathbf{N} \) lies on the second plane. The distance is \( |\mathbf{N}| = \ldots \).

36 The distance between the planes \( x + y + 5z = 7 \) and \( 3x + 2y + z = 1 \) is zero because \( \ldots \).

In Problems 37–41 all points and vectors are in the \( xy \) plane.

37 The line \( 3x + 4y = 10 \) is perpendicular to the vector \( \mathbf{N} = \ldots \). On the line, the closest point to the origin is \( P = \ldots \). Find \( t \) and \( P \) and \( |P| \).

38 Draw the line \( x + 2y = 4 \) and the vector \( \mathbf{N} = 1 + 2j \). The closest point to \( Q = (3, 3) \) is \( P = Q + t\mathbf{N} \). Find \( t \). Find \( P \).

39 A new way to find \( P \) in Problem 37: minimize \( x^2 + y^2 = x^2 + (1/4 - 1/4x)^2 \). By calculus find the best \( x \) and \( y \).

40 To catch a drug runner going from \((0, 0)\) to \((4, 0)\) at 8 meters per second, you must travel from \((0, 3)\) to \((4, 0)\) at \( \ldots \) meters per second. The projection of your velocity vector onto his velocity vector will have length \( \ldots \).

41 Show by vectors that the distance from \((x_1, y_1)\) to the line \( ax + by = d \) is \( |d - ax_1 - by_1|/\sqrt{a^2 + b^2} \).

42 It takes three points to determine a plane. So why does \( ax + by + cz = d \) contain four numbers \( a, b, c, d \)? When does \( ex + fy + gz = 1 \) represent the same plane?

43 \( \text{(projections by calculus)} \) The dot product of \( B - t\mathbf{A} \) with itself is \( |B|^2 - 2\mathbf{A} \cdot \mathbf{B} + t^2|\mathbf{A}|^2 \). (a) This has a minimum at \( t = \ldots \). (b) Then \( t\mathbf{A} \) is the projection of \( B - t\mathbf{A} \).

44 A figure showing \( B, t\mathbf{A}, \) and \( B - t\mathbf{A} \) is worth 1000 words.

45 From their equations, how can you tell if two planes are (a) parallel (b) perpendicular (c) at a 45° angle?

Problems 45–48 are about the ECG and heart vector.

45 The \( \text{aVR lead} = -\mathbf{H}_1 - \mathbf{H}_N \). Find the \( \text{aVL} \) and \( \text{aVF} \) leads toward the left arm and foot. Show that \( \text{aVR} + \text{aVL} + \text{aVF} = 0 \). They go out from the center at 120° angles.

46 Find the projection on the \( \text{aVR lead} \) of \( \mathbf{V} = 2i - j \) in Question 3.

47 If the potentials are \( \phi_{RA} = 1 \) (right arm) and \( \phi_{LA} = 2 \) and \( \phi_{UL} = -3 \), find the heart vector \( \mathbf{V} \). The differences in potential are the projections of \( \mathbf{V} \).

48 If \( \mathbf{V} \) is perpendicular to a lead \( L \), the reading on that lead is \( \ldots \). If \( \int \mathbf{V}(t)dt \) is perpendicular to lead \( L \), why is the area under the reading zero?
After saying that vectors are not multiplied, we offered the dot product. Now we contradict ourselves further, by defining the cross product. Where $\mathbf{A} \cdot \mathbf{B}$ was a number, the cross product $\mathbf{A} \times \mathbf{B}$ is a vector. It has length and direction:

The length is $|\mathbf{A}||\mathbf{B}| \sin \theta$. The direction is perpendicular to $\mathbf{A}$ and $\mathbf{B}$.

The cross product (also called vector product) is defined in three dimensions only. $\mathbf{A}$ and $\mathbf{B}$ lie on a plane through the origin. $\mathbf{A} \times \mathbf{B}$ is along the normal vector $\mathbf{N}$, perpendicular to that plane. We still have to say whether it points “up” or “down” along $\mathbf{N}$.

The length of $\mathbf{A} \times \mathbf{B}$ depends on $\sin \theta$, where $\mathbf{A} \cdot \mathbf{B}$ involved $\cos \theta$. The dot product rewards vectors for being parallel ($\cos 0 = 1$). The cross product is largest when $\mathbf{A}$ is perpendicular to $\mathbf{B}$ ($\sin \pi/2 = 1$). At every angle

$$|\mathbf{A} \cdot \mathbf{B}|^2 + |\mathbf{A} \times \mathbf{B}|^2 = |\mathbf{A}|^2|\mathbf{B}|^2 \cos^2 \theta + |\mathbf{A}|^2|\mathbf{B}|^2 \sin^2 \theta = |\mathbf{A}|^2|\mathbf{B}|^2.$$  

That will be a bridge from geometry to algebra. This section goes from definition to formula to volume to determinant. Equations (6) and (14) are the key formulas for $\mathbf{A} \times \mathbf{B}$.

Notice that $\mathbf{A} \times \mathbf{A} = \mathbf{0}$. (This is the zero vector, not the zero number.) When $\mathbf{B}$ is parallel to $\mathbf{A}$, the angle is zero and the sine is zero. Parallel vectors have $\mathbf{A} \times \mathbf{B} = \mathbf{0}$. Perpendicular vectors have $\sin \theta = 1$ and $|\mathbf{A} \times \mathbf{B}| = |\mathbf{A}||\mathbf{B}| = \text{area of rectangle with sides } \mathbf{A} \text{ and } \mathbf{B}$.

Here are four examples that lead to the cross product $\mathbf{A} \times \mathbf{B}$.

**EXAMPLE 1** (From geometry) Find the area of a parallelogram and a triangle.

Vectors $\mathbf{A}$ and $\mathbf{B}$, going out from the origin, form two sides of a triangle. They produce the parallelogram in Figure 11.13, which is twice as large as the triangle.

The area of a parallelogram is base times height (perpendicular height not sloping height). The base is $|\mathbf{A}|$. The height is $|\mathbf{B}| \sin \theta$. We take absolute values because height and area are not negative. Then the area is the length of the cross product:

$$\text{area of parallelogram} = |\mathbf{A}| |\mathbf{B}| \sin \theta = |\mathbf{A} \times \mathbf{B}|.$$  

**EXAMPLE 2** (From physics) The torque vector $\mathbf{T} = \mathbf{R} \times \mathbf{F}$ produces rotation.

The force $\mathbf{F}$ acts at the point $(x, y, z)$. When $\mathbf{F}$ is parallel to the position vector $\mathbf{R} = xi + yj + zk$, the force pushes outward (no turning). When $\mathbf{F}$ is perpendicular to $\mathbf{R}$, the force creates rotation. For in-between angles there is an outward force $|\mathbf{F}| \cos \theta$ and a turning force $|\mathbf{F}| \sin \theta$. The turning force times the distance $|\mathbf{R}|$ is the **moment** $|\mathbf{R}||\mathbf{F}| \sin \theta$. 

Fig. 11.13 Area $|\mathbf{A} \times \mathbf{B}|$ and moment $|\mathbf{R} \times \mathbf{F}|$. Cross products are perpendicular to the page.
11.3 Cross Products and Determinants

The moment gives the magnitude and sign of the **torque vector** \( \mathbf{T} = \mathbf{R} \times \mathbf{F} \). The direction of \( \mathbf{T} \) is along the axis of rotation, at right angles to \( \mathbf{R} \) and \( \mathbf{F} \).

**EXAMPLE 3** Does the cross product go up or down? *Use the right-hand rule.*

Forces and torques are probably just fine for physicists. Those who are not natural physicists want to see something turn.† We can visualize a record or compact disc rotating around its axis—which comes up through the center.

At a point on the disc, you give a push. When the push is outward (hard to do), nothing turns. Rotation comes from force “around” the axis. The disc can turn either way—depending on the angle between force and position. A sign convention is necessary, and it is the **right-hand rule**:

\[ \mathbf{A} \times \mathbf{B} \text{ points along your right thumb when the fingers curl from } \mathbf{A} \text{ toward } \mathbf{B}. \]

This rule is simplest for the vectors \( \mathbf{i}, \mathbf{j}, \mathbf{k} \) in Figure 11.14—which is all we need.

Suppose the fingers curl from \( \mathbf{i} \) to \( \mathbf{j} \). The thumb points along \( \mathbf{k} \). *The \( x-y-z \) axes form a "right-handed triple."* Since \( |\mathbf{i}| = 1 \) and \( |\mathbf{j}| = 1 \) and \( \sin \pi/2 = 1 \), the length of \( \mathbf{i} \times \mathbf{j} \) is 1. *The cross product is* \( \mathbf{i} \times \mathbf{j} = \mathbf{k} \). The disc turns counterclockwise—it's angular velocity is up—when the force acts at \( \mathbf{i} \) in the direction \( \mathbf{j} \).

Figure 11.14b reverses \( \mathbf{i} \) and \( \mathbf{j} \). The force acts at \( \mathbf{j} \) and its direction is \( \mathbf{i} \). The disc turns clockwise (the way records and compact discs actually turn). When the fingers curl from \( \mathbf{j} \) to \( \mathbf{i} \), the thumb points down. Thus \( \mathbf{j} \times \mathbf{i} = -\mathbf{k} \). This is a special case of an amazing rule:

*The cross product is anticommutative:* \[ \mathbf{B} \times \mathbf{A} = - (\mathbf{A} \times \mathbf{B}). \] (3)

That is quite remarkable. Its discovery by Hamilton produced an intellectual revolution in 19th century algebra, which had been totally accustomed to \( \mathbf{A} \mathbf{B} = \mathbf{B} \mathbf{A} \). This commutative law is old and boring for numbers (it is new and boring for dot products). Here we see its opposite for vector products \( \mathbf{A} \times \mathbf{B} \). Neither law holds for matrix products.

**EXAMPLE 4** A screw goes into a wall or out, following the right-hand rule.

The disc was in the \( xy \) plane. So was the force. (We are not breaking records here.) The axis was up and down. To see the cross product more completely we need to turn a screw into a wall.

Figure 11.14b shows the \( xz \) plane as the wall. The screw is in the \( y \) direction. By turning from \( x \) toward \( z \) we drive the screw into the wall—which is the negative \( y \) direction. In other words \( \mathbf{i} \times \mathbf{k} \) equals minus \( \mathbf{j} \). We turn the screw clockwise to make it go in. To take out the screw, twist from \( \mathbf{k} \) toward \( \mathbf{i} \). Then \( \mathbf{k} \times \mathbf{i} \) equals plus \( \mathbf{j} \).

†Everybody is a natural mathematician. That is the axiom behind this book.
To summarize: \( k \times i = j \) and \( j \times k = i \) have plus signs because \( kij \) and \( jki \) are in the same "cyclic order" as \( ijk \). (Anticyclic is minus.) The \( z-x-y \) and \( y-z-x \) axes form right-handed triples like \( x-y-z \).

**The Formula for the Cross Product**

We begin the algebra of \( A \times B \). It is essential for computation, and it comes out beautifully. The square roots in \( |A||B| \sin \theta \) will disappear in formula (6) for \( A \times B \). (The square roots also disappeared in \( A \cdot B \), which is \( |A||B| \cos \theta \). But \( |A||B| \tan \theta \) would be terrible.) Since \( A \times B \) is a vector we need to find three components.

Start with the two-dimensional case. The vectors \( a_1 i + a_2 j \) and \( b_1 i + b_2 j \) are in the \( xy \) plane. Their cross product must go in the \( z \) direction. Therefore \( A \times B = \_? \_k \) and there is only one nonzero component. It must be \( |A||B| \sin \theta \) (with the correct sign), but we want a better formula. There are two clean ways to compute \( A \times B \), either by algebra \((a)\) or by a bridge \((b)\) to the dot product and geometry:

(a) \( (a_1 i + a_2 j) \times (b_1 i + b_2 j) = a_1 b_2 i \times i + a_1 b_1 j \times j + a_2 b_1 i \times j + a_2 b_2 j \times i \). \((4)\)

On the right are 0, \(-a_1 b_1 k\), and 0. The cross product is \( (a_1 b_2 - a_2 b_1)k \).

(b) Rotate \( B = b_1 i + b_2 j \) clockwise through 90° into \( B^* = b_2 i - b_1 j \). Its length is unchanged (and \( B \cdot B^* = 0 \)). Then \( |A||B| \sin \theta \) equals \( |A||B^*| \cos \theta \), which is \( A \cdot B^* \):

\[
|A||B| \sin \theta = A \cdot B^* = \begin{bmatrix} a_1 \\ a_2 \end{bmatrix} \cdot \begin{bmatrix} b_2 \\ -b_1 \end{bmatrix} = a_1 b_2 - a_2 b_1. \quad (5)
\]

**Example 5** For \( A = i + 2j \) and \( B = 4i + 5j \) the cross product is \( (1 \cdot 5 - 2 \cdot 4)k = -3k \). Area of parallelogram = 3, area of triangle = 3/2. The minus sign in \( A \times B = -3k \) is absent in the areas.

**Note** Splitting \( A \times B \) into four separate cross products is correct, but it does not follow easily from \( |A||B| \sin \theta \). Method \((a)\) is not justified until Remark 1 below. An algebraist would change the definition of \( A \times B \) to start with the distributive law (splitting rule) and the anticommutative law:

\[
A \times (B + C) = (A \times B) + (A \times C) \quad \text{and} \quad A \times B = -(B \times A).
\]

**The Cross Product Formula (3 Components)**

We move to three dimensions. The goal is to compute all three components of \( A \times B \) (not just the length). Method \((a)\) splits each vector into its \( i, j, k \) components, making nine separate cross products:

\[
(a_1 i + a_2 j + a_3 k) \times (b_1 i + b_2 j + b_3 k) = a_1 b_2 (i \times i) + a_1 b_3 (i \times j) + a_2 b_3 (j \times i) + \text{seven more terms}.
\]

Remember \( i \times i = j \times j = k \times k = 0 \). Those three terms disappear. The other six terms come in pairs, and please notice the cyclic pattern:

**Formula** \( A \times B = (a_2 b_3 - a_3 b_2)i + (a_3 b_1 - a_1 b_3)j + (a_1 b_2 - a_2 b_1)k. \) \((6)\)

The \( k \) component is the \( 2 \times 2 \) answer, when \( a_3 = b_3 = 0 \). The \( i \) component involves indices 2 and 3, \( j \) involves 3 and 1, \( k \) involves 1 and 2. The cross product formula is
written as a "determinant" in equation (14) below—many people use that form to compute \( \mathbf{A} \times \mathbf{B} \).

**EXAMPLE 6** \((i+2j+3k) \times (4i+5j+6k) = (2 \cdot 6 - 3 \cdot 5)i + (3 \cdot 4 - 1 \cdot 6)j + (1 \cdot 5 - 2 \cdot 4)k\). The \( i, j, k \) components give \( \mathbf{A} \times \mathbf{B} = -3i + 6j - 3k \). Never add the -3, 6, and -3.

**Remark 1** The three-dimensional formula (6) is still to be matched with \( \mathbf{A} \times \mathbf{B} \) from geometry. One way is to rotate \( \mathbf{B} \) into \( \mathbf{B}^* \) as before, staying in the plane of \( \mathbf{A} \) and \( \mathbf{B} \). Fortunately there is an easier test. The vector in equation (6) satisfies all four geometric requirements on \( \mathbf{A} \times \mathbf{B} \): perpendicular to \( \mathbf{A} \), perpendicular to \( \mathbf{B} \), correct length, right-hand rule. The length is checked in Problem 16—here is the zero dot product with \( \mathbf{A} \):

\[
\mathbf{A} \cdot (\mathbf{A} \times \mathbf{B}) = a_1(a_2b_3 - a_3b_2) + a_2(a_3b_1 - a_1b_3) + a_3(a_1b_2 - a_2b_1) = 0. 
\] (7)

**Remark 2** (Optional) There is a wonderful extension of the Pythagoras formula \( a^2 + b^2 = c^2 \). Instead of sides of a triangle, we go to areas of projections on the \( yz, xz, \) and \( xy \) planes. \( 3^2 + 6^2 + 3^2 \) is the square of the parallelogram area in Example 6.

For triangles these areas are cut in half. Figure 11.15a shows three projected triangles of area \( \frac{1}{2} \). Its Pythagoras formula is \( (\frac{1}{2})^2 + (\frac{1}{2})^2 + (\frac{1}{2})^2 = \text{(area of } \mathbf{PQR})^2 \).

**EXAMPLE 7** \( P = (1, 0, 0), Q = (0, 1, 0), R = (0, 0, 1) \) lie in a plane. Find its equation.

**Idea for any \( P, Q, R \):** Find vectors \( \mathbf{A} \) and \( \mathbf{B} \) in the plane. Compute the normal \( \mathbf{N} = \mathbf{A} \times \mathbf{B} \).

**Solution** The vector from \( P \) to \( Q \) has components -1, 1, 0. It is \( \mathbf{A} = j - i \) (subtract to go from \( P \) to \( Q \)). Similarly the vector from \( P \) to \( R \) is \( \mathbf{B} = k - i \). Since \( \mathbf{A} \) and \( \mathbf{B} \) are in the plane of Figure 11.15, \( \mathbf{N} = \mathbf{A} \times \mathbf{B} \) is perpendicular:

\[
(j - i) \times (k - i) = (j \times k) - (i \times k) - (j \times i) + (i \times i) = i + j + k. 
\] (8)

*The normal vector is \( \mathbf{N} = i + j + k \). The equation of the plane is \( 1x + 1y + 1z = d \).*

With the right choice \( d = 1 \), this plane contains \( P, Q, R \). The equation is \( x + y + z = 1 \).

**EXAMPLE 8** What is the area of this same triangle \( \mathbf{PQR} \)?

**Solution** The area is half of the cross-product length \( |\mathbf{A} \times \mathbf{B}| = |i + j + k| = \sqrt{3} \).

---

**DETERMINANTS AND VOLUMES**

We are close to good algebra. The two plane vectors \( a_1i + a_2j \) and \( b_1i + b_2j \) are the sides of a parallelogram. Its area is \( a_1b_2 - a_2b_1 \), possibly with a sign change. There
is a special way to write these four numbers—in a "square matrix." There is also a
name for the combination that leads to area. It is the "determinant of the matrix":

\[ \begin{vmatrix} a_1 & a_2 \\ b_1 & b_2 \end{vmatrix}, \text{ its determinant is } \begin{vmatrix} a_1 & a_2 \\ b_1 & b_2 \end{vmatrix} = a_1 b_2 - a_2 b_1. \]

This is a 2 by 2 matrix (notice brackets) and a 2 by 2 determinant (notice vertical bars). The matrix is an array of four numbers and the determinant is one number:

Examples of determinants: 
\[ \begin{vmatrix} 2 & 1 \\ 4 & 3 \end{vmatrix} = 6 - 4 = 2, \quad \begin{vmatrix} 2 & 1 \\ 2 & 1 \end{vmatrix} = 0, \quad \begin{vmatrix} 1 & 0 \\ 0 & 1 \end{vmatrix} = 1. \]

The second has no area because \( \mathbf{A} = \mathbf{B} \). The third is a unit square (\( \mathbf{A} = \mathbf{i}, \mathbf{B} = \mathbf{j} \)).

Now move to three dimensions, where determinants are most useful. The parallelogram becomes a parallelepiped. The word "box" is much shorter, and we will use it, but remember that the box is squashed. (Like a rectangle squashed to a parallelogram, the angles are generally not 90°.) The three edges from the origin are \( \mathbf{A} = (a_1, a_2, a_3), \mathbf{B} = (b_1, b_2, b_3), \mathbf{C} = (c_1, c_2, c_3) \). Those edges are at right angles only when \( \mathbf{A} \cdot \mathbf{B} = \mathbf{A} \cdot \mathbf{C} = \mathbf{B} \cdot \mathbf{C} = 0 \).

**Question:** What is the volume of the box? The right-angle case is easy—it is length times width times height. The volume is \( |\mathbf{A}| \times |\mathbf{B}| \times |\mathbf{C}| \), when the angles are 90°. For a squashed box (Figure 11.15) we need the perpendicular height, not the sloping height.

There is a beautiful formula for volume. \( \mathbf{B} \) and \( \mathbf{C} \) give a parallelogram in the base, and \( \mathbf{B} \times \mathbf{C} \) is the base area. This cross product points straight up. The third vector \( \mathbf{A} \) points up at an angle—its perpendicular height is \( |\mathbf{A}| \cos \theta \). Thus the volume is \( |\mathbf{B} \times \mathbf{C}| \times |\mathbf{A}| \times \cos \theta \). **The volume is the dot product of \( \mathbf{A} \) with \( \mathbf{B} \times \mathbf{C} \).**

**11G The triple scalar product is** \( \mathbf{A} \cdot (\mathbf{B} \times \mathbf{C}) \). **Volume of box** \( = |\mathbf{A} \cdot (\mathbf{B} \times \mathbf{C})| \).

**Important:** \( \mathbf{A} \cdot (\mathbf{B} \times \mathbf{C}) \) is a number, not a vector. This volume is zero when \( \mathbf{A} \) is in the same plane as \( \mathbf{B} \) and \( \mathbf{C} \) (the box is totally flattened). Then \( \mathbf{B} \times \mathbf{C} \) is perpendicular to \( \mathbf{A} \) and their dot product is zero.

**Useful facts:** \( \mathbf{A} \cdot (\mathbf{B} \times \mathbf{C}) = (\mathbf{A} \times \mathbf{B}) \cdot \mathbf{C} = \mathbf{C} \cdot (\mathbf{A} \times \mathbf{B}) = \mathbf{B} \cdot (\mathbf{C} \times \mathbf{A}) \).

All those come from the same box, with different sides chosen as base—but no change in volume. Figure 11.15 has \( \mathbf{B} \) and \( \mathbf{C} \) in the base but it can be \( \mathbf{A} \) and \( \mathbf{B} \) or \( \mathbf{A} \) and \( \mathbf{C} \). **The triple product** \( \mathbf{A} \cdot (\mathbf{C} \times \mathbf{B}) \) **has opposite sign**, since \( \mathbf{C} \times \mathbf{B} = - (\mathbf{B} \times \mathbf{C}) \). This order \( \mathbf{ACB} \) is not cyclic like \( \mathbf{ABC} \) and \( \mathbf{CAB} \) and \( \mathbf{BCA} \).

To compute this triple product \( \mathbf{A} \cdot (\mathbf{B} \times \mathbf{C}) \), we take \( \mathbf{B} \times \mathbf{C} \) from equation (6):

\[ \mathbf{A} \cdot (\mathbf{B} \times \mathbf{C}) = a_1 (b_2 c_3 - b_3 c_2) + a_2 (b_3 c_1 - b_1 c_3) + a_3 (b_1 c_2 - b_2 c_1). \] (9)

The numbers \( a_1, a_2, a_3 \) multiply 2 by 2 determinants to give a 3 by 3 determinant! There are three terms with plus signs (like \( a_1 b_2 c_3 \)). The other three have minus signs (like \( -a_1 b_3 c_2 \)). The plus terms have indices 123, 231, 312 in cyclic order. The minus terms have anticyclic indices 132, 213, 321. Again there is a special way to write the nine components of \( \mathbf{A}, \mathbf{B}, \mathbf{C} \)—as a "3 by 3 matrix." The combination in (9), which
1.3 Cross Products and Determinants

gives volume, is a "3 by 3 determinant:

\[
\begin{vmatrix}
  a_1 & a_2 & a_3 \\
  b_1 & b_2 & b_3 \\
  c_1 & c_2 & c_3
\end{vmatrix}
= \begin{vmatrix}
  a_1 & a_2 & a_3 \\
  b_1 & b_2 & b_3 \\
  c_1 & c_2 & c_3
\end{vmatrix}
\]

A single number is produced out of nine numbers, by formula (9). The nine numbers are multiplied three at a time, as in \( a_1 b_1 c_2 - \) except this product is not allowed. **Each row and column must be represented once.** This gives the six terms in the determinant:

\[
\begin{vmatrix}
  a_1 & a_2 & a_3 \\
  b_1 & b_2 & b_3 \\
  c_1 & c_2 & c_3
\end{vmatrix} = a_1 b_2 c_3 + a_2 b_3 c_1 + a_3 b_1 c_2 - a_1 b_3 c_2 - a_2 b_1 c_3 - a_3 b_2 c_1
\]

(10)

The trick is in the \( \pm \) signs. Products down to the right are "plus":

\[
\begin{vmatrix}
  2 & 1 & 1 \\
  1 & 2 & 1 \\
  1 & 1 & 2
\end{vmatrix}
= 2 \cdot 2 \cdot 2 + 1 \cdot 1 \cdot 1 + 1 \cdot 1 \cdot 1 = 8 + 1 + 1 = 4.
\]

With practice the six products like 2·2·2 are done in your head. Write down only 8 + 1 + 2 - 2 - 2 = 4. This is the determinant and the volume.

Note the special case when the vectors are \( i, j, k \). The box is a unit cube:

\[
\text{volume of cube} = \begin{vmatrix}
  1 & 0 & 0 \\
  0 & 1 & 0 \\
  0 & 0 & 1
\end{vmatrix} = 1 + 0 + 0 = 1.
\]

If \( A, B, C \) lie in the same plane, the volume is zero. A zero determinant is the test to see whether three vectors lie in a plane. Here row \( A = \text{row } B - \text{row } C \):

\[
\begin{vmatrix}
  0 & 1 & -1 \\
  -1 & 1 & 0 \\
  -1 & 0 & 1
\end{vmatrix} = 0 \cdot 1 \cdot 1 + 1 \cdot 0 \cdot (-1) + (-1) \cdot (-1) \cdot 0 = 0.
\]

(11)

Zeros in the matrix simplify the calculation. All three products with plus signs—down to the right—are zero. The only two nonzero products cancel each other.

If the three \( -1 \)'s are changed to \( +1 \)'s, the determinant is \( -2 \). The determinant can be negative when all nine entries are positive! A negative determinant only means that the rows \( A, B, C \) form a "left-handed triple." This extra information from the sign—right-handed vs. left-handed—is free and useful, but the volume is the absolute value.

The determinant yields the volume also in higher dimensions. In physics, four dimensions give space-time. Ten dimensions give superstrings. Mathematics uses all dimensions. The 64 numbers in an 8 by 8 matrix give the volume of an eight-dimensional box—with \( 8! = 40,320 \) terms instead of \( 3! = 6 \). Under pressure from my class I omit the formula.
Question: When is the point \((x, y, z)\) on the plane through the origin containing \(B\) and \(C\)? For the vector \(A = xi + yj + zk\) to lie in that plane, the volume \(A \cdot (B \times C)\) must be zero. The equation of the plane is \textit{determinant} = zero.

Follow this example for \(B = j - i\) and \(C = k - i\) to find the plane parallel to \(B\) and \(C\):

\[
\begin{vmatrix}
  x & y & z \\
  -1 & 1 & 0 \\
  -1 & 0 & 1 \\
\end{vmatrix}
= x \cdot 1 \cdot 1 + y \cdot 0 \cdot (-1) + z \cdot 0 \cdot (-1) = 0.
\]

This equation is \(x + y + z = 0\). The normal vector \(N = B \times C\) has components 1, 1, 1.

\[\text{THE CROSS PRODUCT AS A DETERMINANT}\]

There is a connection between 3 by 3 and 2 by 2 determinants that you have to see. The numbers in the top row multiply determinants from the other rows:

\[
\begin{vmatrix}
  a_1 & a_2 & a_3 \\
  b_1 & b_2 & b_3 \\
  c_1 & c_2 & c_3 \\
\end{vmatrix}
= a_1 \begin{vmatrix}
  b_2 & b_3 \\
  c_2 & c_3 \\
\end{vmatrix}
- a_2 \begin{vmatrix}
  b_1 & b_3 \\
  c_1 & c_3 \\
\end{vmatrix}
+ a_3 \begin{vmatrix}
  b_1 & b_2 \\
  c_1 & c_2 \\
\end{vmatrix}.
\]

The highlighted product \(a_1(b_2c_3 - b_3c_2)\) gives two of the six terms. \textit{All six products contain an} \(a\) \textit{and} \(b\) \textit{and} \(c\) \textit{from different columns}. There are \(3! = 6\) different orderings of columns 1, 2, 3. Note how \(a_2\) multiplies a determinant from columns 1 and 2.

Equation (13) is identical with equations (9) and (10). We are meeting the same six terms in different ways. The new feature is the minus sign in front of \(a_2\)—and the common mistake is to forget that sign. In a 4 by 4 determinant, \(a_1, -a_2, a_3, -a_4\) would multiply 3 by 3 determinants.

Now comes a key step. We write \(A \times B\) as a determinant. The vectors \(i, j, k\) go in the top row, the components of \(A\) and \(B\) go in the other rows. The "determinant" is exactly \(A \times B\):

\[
A \times B = \begin{vmatrix}
  i & j & k \\
  a_1 & a_2 & a_3 \\
  b_1 & b_2 & b_3 \\
\end{vmatrix}
= i \begin{vmatrix}
  a_2 & a_3 \\
  b_2 & b_3 \\
\end{vmatrix}
- j \begin{vmatrix}
  a_1 & a_3 \\
  b_1 & b_3 \\
\end{vmatrix}
+ k \begin{vmatrix}
  a_1 & a_2 \\
  b_1 & b_2 \\
\end{vmatrix}.
\]

This time we highlighted the \(j\) component with its minus sign. There is no great mathematics in formula (14)—it is probably illegal to mix \(i, j, k\) with six numbers but it works. This is the good way to remember and compute \(A \times B\). In the example \((j - i) \times (k - i)\) from equation (8), those two vectors go into the last two rows:

\[
\begin{vmatrix}
  i & j & k \\
  -1 & 1 & 0 \\
  -1 & 0 & 1 \\
\end{vmatrix}
= i \begin{vmatrix}
  1 & 0 \\
  0 & 1 \\
\end{vmatrix}
- j \begin{vmatrix}
  1 & 0 \\
  0 & 1 \\
\end{vmatrix}
+ k \begin{vmatrix}
  -1 & 1 \\
  -1 & 0 \\
\end{vmatrix}
= i + j + k.
\]

The \(k\) component is highlighted, to see \(a_1b_2 - a_2b_1\) again. Note the change from equation (11), which had 0, 1, −1 in the top row. That triple product was a number (zero). This cross product is a vector \(i + j + k\).
11.3 Cross Products and Determinants

Review question 1 With the i, j, k row changed to 3, 4, 5, what is the determinant?
Answer $3 \cdot 1 \cdot 4 + 1 \cdot 5 + 1 \cdot 1 = 12$. That triple product is the volume of a box.

Review question 2 When is $A \times B = 0$ and when is $A \cdot (B \times C) = 0$? Zero vector, zero number.
Answer When A and B are on the same line. When A, B, C are in the same plane.

Review question 3 Does the parallelogram area $|A \times B|$ equal a 2 by 2 determinant?
Answer If A and B lie in the xy plane, yes. Generally no.

Review question 4 What are the vector triple products $(A \times B) \times C$ and $A \times (B \times C)$?
Answer Not computed yet. These are two new vectors in Problem 47.

Review question 5 Find the plane through the origin containing $A = i + j + 2k$ and $B = i + k$. Find the cross product of those same vectors $A$ and $B$.
Answer The position vector $P = xi + yj + zk$ is perpendicular to $N = A \times B$:

\[
\begin{vmatrix}
  x & y & z \\
  1 & 1 & 2 \\
  0 & 1 & 0
\end{vmatrix} = x + y - z = 0. \\
A \times B = \begin{vmatrix}
  i & j & k \\
  1 & 1 & 2 \\
  1 & 0 & 1
\end{vmatrix} = i + j - k.
\]

Read-through questions

The cross product $A \times B$ is a vector whose length is to $A$ and $B$. That length is the area of a parallelogram whose base is $|A|$ and whose height is the height. When $A = a_1 i + a_2 j + a_3 k$ and $B = b_1 i + b_2 j$, the area is $f$. This equals a 2 by 2 determinant. In general $|A \cdot B|^2 + |A \times B|^2 = h$.

The rules for cross products are $A \times A = 0$ and $A \times B = -(B \times A)$ and $A \times (B + C) = A \times B + A \times C$. In particular $A \times B$ needs the right-hand rule to decide its direction. If the fingers curl from A towards B (not more than 180°), then $m$ points $n$. By this rule $i \times j = o$ and $i \times k = p$ and $j \times k = q$.

The vectors $a_1 i + a_2 j + a_3 k$ and $b_1 i + b_2 j + b_3 k$ have cross product $i + j + k$. The vectors $A = i + j + k$ and $B = i + j$ have $A \times B = u$. (This is also the 3 by 3 determinant.) Perpendicular to the plane containing $(0, 0, 0), (1, 1, 0), (1, 1, 1)$ is the normal vector $N = w$. The area of the triangle with those three vertices is $v$, which is half the area of the parallelogram with fourth vertex at $z$. Vectors $A, B, C$ from the origin determine a vector. Its volume $|A \cdot (\vec{A} \times \vec{B})|$ comes from a 3 by 3 $B$. There are six terms, $C$, with a plus sign and $D$, with minus. In every term each row and $E$ is represented once. The rows $(1, 0, 0), (0, 0, 1),$ and $(0, 1, 0)$ have determinant $f$. That box is a $G$, but its sides form a $H$-handed triple in the order given.

Exercises

1. Compute the cross products 1–8 from formula (6) or the determinant (14). Do one example both ways.
   1. $(i \times j) \times k$
   2. $(i \times j) \times i$
   3. $(2i + 3j) \times (i + k)$
   4. $(2i + 3j + k) \times (2i + 3j - k)$
   5. $(2i + 3j + k) \times (i - j - k)$
   6. $(i + j - k) \times (i + j + k)$
   7. $(i + 2j + 3k) \times (4i - 9j)$
   8. $(i \cos \theta + j \sin \theta) \times (i \sin \theta - j \cos \theta)$
   9. When are $|A \times B| = |A||B|$ and $|A \times (B \times C)| = |A||B||C|$?
   10. True or false:
       (a) $A \times B$ never equals $A \cdot B$.
       (b) If $A \times B = 0$ and $A \cdot B = 0$, then either $A = 0$ or $B = 0$.
       (c) If $A \times B = A \times C$ and $A \neq 0$, then $B = C$. 
In 11–16 find \(|A \times B|\) by equation (1) and then by computing \(A \times B\) and its length.

11. \(A = i + j + k, B = i\)

12. \(A = i + j, B = i - j\)

13. \(A = -B\)

14. \(A = i + j, B = j + k\)

15. \(A = a_1i + a_2j, B = b_1i + b_2j\)

16. \(A = (a_1, a_2, a_3), B = (b_1, b_2, b_3)\)

In Problem 16 (the general case), equation (1) proves that the length from equation (6) is correct.

17. True or false, by testing on \(A = i, B = j, C = k:\)
   
   (a) \(A \times (A \times B) = 0\)
   
   (b) \(A \cdot (B \times C) = (A \times B) \cdot C\)
   
   (c) \(A \cdot (B \times C) = C \cdot (B \times A)\)
   
   (d) \((A - B) \times (A + B) = 2(A \times B)\)

18. (a) From \(A \times B = -B \times A\) deduce that \(A \times A = 0\).
   
   (b) Split \((A + B) \times (A + B)\) into four terms, to deduce that \((A \times B) = -(B \times A)\).

33. When \(B = 3i + j\) is rotated 90° clockwise in the xy plane it becomes \(B^* = \ldots\). When rotated 90° counterclockwise it is \(\ldots\). When rotated 180° it is \(\ldots\).

34. From formula (6) verify that \(B \cdot (A \times B) = 0\).

35. Compute

\[
\begin{bmatrix}
1 & 2 & 3 \\
2 & 3 & 4 \\
3 & 4 & 6
\end{bmatrix}
\begin{bmatrix}
1 & 0 & 2 \\
1 & 2 & 1 \\
0 & 1 & 2
\end{bmatrix}
\]

36. Which of the following are equal to \(A \times B\)?
   
   \(A + B) \times B, \ (-B) \times (-A), \ |A||B| \sin \theta, \ (A + C) \times (B - C), \ \frac{1}{2}(A - B) \times (A + B)\).

37. Compare the six terms on both sides to prove that

\[
\begin{bmatrix}
a_1 & b_1 & c_1 \\
a_2 & b_2 & c_2 \\
a_3 & b_3 & c_3
\end{bmatrix}
\]

The matrix is "transposed"—same determinant.

38. Compare the six terms to prove that

\[
\begin{bmatrix}
a_1 & a_2 & a_3 \\
b_1 & b_2 & b_3
\end{bmatrix}
\begin{bmatrix}
c_1 & c_2 & c_3
\end{bmatrix}
\]

This is an "expansion on row 2." Note minus signs.

39. Choose the signs and 2 by 2 determinants in

\[
\begin{bmatrix}
a_1 & a_2 & a_3 \\
b_1 & b_2 & b_3
\end{bmatrix}
\begin{bmatrix}
c_1 & c_2 & c_3
\end{bmatrix}
\]

40. Show that \((A \times B) + (B \times C) + (C \times A)\) is perpendicular to \(B - A\) and \(C - B\) and \(A - C\).

Problems 41–44 compute the areas of triangles.

41. The triangle \(PQR\) in Example 7 has squared area \((\sqrt{3}/2)^2 = (i)^2 + (i)^2 + (i)^2\), from the 3D version of Pythagoras in Remark 2. Find the area of \(PQR\) when \(P = (a, 0, 0), Q = (0, b, 0),\) and \(R = (0, 0, c)\). Check with \(\frac{1}{2}|A \times B|\).

42. A triangle in the xy plane has corners at \((a_1, b_1), (a_2, b_2)\) and \((a_3, b_3)\). Its area \(A\) is half the area of a parallelogram. Find two sides of the parallelogram and explain why

\[A = \frac{1}{2}|(a_2 - a_1) (b_3 - b_1) - (a_3 - a_1) (b_2 - b_1)|\]

43. By Problem 42 find the area \(A\) of the triangle with corners \((2, 1)\) and \((4, 2)\) and \((1, 2)\). Where is a fourth corner to make a parallelogram?
Lifting the triangle of Problem 42 up to the plane \( z = 1 \) gives corners \((a_1, b_1, 1), (a_2, b_2, 1), (a_3, b_3, 1)\). The area of the triangle times \( \frac{1}{3} \) is the volume of the upside-down pyramid from \((0, 0, 0)\) to these corners. This pyramid volume is \( \frac{1}{3} \) the box volume, so \( \frac{1}{3} \) (area of triangle) = \( \frac{1}{4} \) (volume of box):

\[
\text{area of triangle} = \frac{1}{2} \begin{vmatrix} a_1 & b_1 & 1 \\ a_2 & b_2 & 1 \\ a_3 & b_3 & 1 \end{vmatrix}.
\]

Find the area \( A \) in Problem 43 from this determinant.

(a) The projections of \( A = a_1 i + a_2 j + a_3 k \) and \( B = b_1 i + b_2 j + b_3 k \) onto the \( xy \) plane are \( \boxed{\text{_________}} \).

(b) The parallelogram with sides \( A \) and \( B \) projects to a parallelogram with area \( \boxed{\text{_________}} \).

(c) General fact: The projection onto the plane normal to the unit vector \( n \) has area \( (A \times B) \cdot n \). Verify for \( n = k \).

(a) For \( A = i + j + 4k \) and \( B = -i + j \), compute \( (A \times B) \cdot i \) and \( (A \times B) \cdot j \) and \( (A \times B) \cdot k \). By Problem 45 those are the areas of projections onto the \( xz \) and \( xy \) and \( xy \) planes.

(b) Square and add those areas to find \( |A \times B|^2 \). This is the Pythagoras formula in space (Remark 2).

(a) The triple cross product \( (A \times B) \times C \) is in the plane of \( A \) and \( B \), because it is perpendicular to the cross product \\

(b) Compute \( (A \times B) \times C \) when \( A = a_1 i + a_2 j + a_3 k, \; B = b_1 i + b_2 j + b_3 k, \; C = i \).

(c) Compute \( (A \cdot C)B - (B \cdot C)A \) when \( C = i \). The answers in (b) and (c) should agree. This is also true if \( C = j \) or \( C = k \) or \( C = c_1 i + c_2 j + c_3 k \).

That proves the tricky formula

\[
(A \times B) \times C = (A \cdot C)B - (B \cdot C)A.
\]

Take the dot product of equation \((*)\) with \( D \) to prove

\[
(A \times B) \cdot (C \times D) = (A \cdot C)(B \cdot D) - (B \cdot C)(A \cdot D).
\]

The plane containing \( P = (0, 1, 1) \) and \( Q = (1, 0, 1) \) and \( R = (1, 1, 0) \) is perpendicular to the cross product \( N \cdot \boxed{\text{_________}} \). Find the equation of the plane and the area of triangle \( PQR \).

Let \( P = (1, 0, -1), \; Q = (1, 1, 1), \; R = (2, 2, 1). \) Choose \( S \) so that \( PQRS \) is a parallelogram and compute its area. Choose \( T, \; U, \; V \) so that \( OPQRSTUV \) is a box (parallelepiped) and compute its volume.

---

11.4 Matrices and Linear Equations

We are moving from geometry to algebra. Eventually we get back to calculus, where functions are nonlinear—but linear equations come first. In Chapter 1, \( y = mx + b \) produced a line. Two equations produce two lines. If they cross, the intersection point solves both equations—and we want to find it.

Three equations in three variables \( x, y, z \) produce three planes. Again they go through one point (usually). Again the problem is to find that intersection point—which solves the three equations.

The ultimate problem is to solve \( n \) equations in \( n \) unknowns. There are \( n \) hyperplanes in \( n \)-dimensional space, which meet at the solution. We need a test to be sure they meet. We also want the solution. These are the objectives of linear algebra, which joins with calculus at the center of pure and applied mathematics.

Like every subject, linear algebra requires a good notation. To state the equations and solve them, we introduce a "matrix." The problem will be \( Au = d \). The solution will be \( u = A^{-1}d \). It remains to understand where the equations come from, where the answer comes from, and what the matrices \( A \) and \( A^{-1} \) stand for.

---

**TWO EQUATIONS IN TWO unknowns**

Linear algebra has no reason to choose one variable as special. The equation \( y - y_0 = m(x - x_0) \) separates \( y \) from \( x \). A better equation for a line is \( ax + by = d. \) (A vertical
line like \( x = 5 \) appears when \( b = 0 \). The first form did not allow slope \( m = \infty \).) This section studies two lines:

\[
\begin{align*}
  a_1 x + b_1 y &= d_1 \\
  a_2 x + b_2 y &= d_2.
\end{align*}
\]

(1)

By solving both equations at once, we are asking \((x, y)\) to lie on both lines. The practical question is: Where do the lines cross? The mathematician's question is: Does a solution exist and is it unique?

To understand everything is not possible. There are parts of life where you never know what is going on (until too late). But two equations in two unknowns can have no mysteries. There are three ways to write the system—by rows, by columns, and by matrices. Please look at all three, since setting up a problem is generally harder and more important than solving it. After that comes the concession to the real world: we compute \( x \) and \( y \).

**EXAMPLE 1** How do you invest \$5000 to earn \$400 a year interest, if a money market account pays 5% and a deposit account pays 10%?

**Set up equations by rows:** With \( x \) dollars at 5% the interest is .05\( x \). With \( y \) dollars at 10% the interest is .10\( y \). One row for principal, another row for interest:

\[
\begin{align*}
  x + y &= 5000 \\
  .05x + .10y &= 400.
\end{align*}
\]

(2)

**Same equations by columns:** The left side of (2) contains \( x \) times one vector plus \( y \) times another vector. The right side is a third vector. The equation by columns is

\[
\begin{bmatrix}
  x \\
  .05x + .10y
\end{bmatrix}
= 
\begin{bmatrix}
  5000 \\
  400
\end{bmatrix}.
\]

(3)

**Same equations by matrices:** Look again at the left side. There are two unknowns \( x \) and \( y \), which go into a vector \( u \). They are multiplied by the four numbers 1, .05, 1, and .10, which go into a two by two matrix \( A \). The left side becomes a matrix times a vector:

\[
Au = 
\begin{bmatrix}
  1 & 1 \\
  .05 & .10
\end{bmatrix}
\begin{bmatrix}
  x \\
  y
\end{bmatrix}
= 
\begin{bmatrix}
  5000 \\
  400
\end{bmatrix}.
\]

(4)

Now you see where the "rows" and "columns" came from. They are the rows and columns of a matrix. The rows entered the separate equations (2). The columns entered the vector equation (3). The matrix-vector multiplication \( Au \) is defined so that all these equations are the same:

\[
\begin{align*}
  Au & \text{ by rows:} \\
  \begin{bmatrix}
    a_1 & b_1 \\
    a_2 & b_2
  \end{bmatrix}
  \begin{bmatrix}
    x \\
    y
  \end{bmatrix}
  &= 
  \begin{bmatrix}
    a_1 x + b_1 y \\
    a_2 x + b_2 y
  \end{bmatrix} \quad \text{(each row is a dot product)}
  \\
  Au & \text{ by columns:} \\
  \begin{bmatrix}
    a_1 & b_1 \\
    a_2 & b_2
  \end{bmatrix}
  \begin{bmatrix}
    x \\
    y
  \end{bmatrix}
  &= 
  x \begin{bmatrix}
    a_1 \\
    a_2
  \end{bmatrix}
  + 
  y \begin{bmatrix}
    b_1 \\
    b_2
  \end{bmatrix} \quad \text{(combination of column vectors)}
\end{align*}
\]

\(A\) is the coefficient matrix. The unknown vector is \( u \). The known vector on the right side, with components 5000 and 400, is \( d \). The matrix equation is \( Au = d \).
This notation \( Au = d \) continues to apply when there are more equations and more unknowns. The matrix \( A \) has a row for each equation (usually \( m \) rows). It has a column for each unknown (usually \( n \) columns). For 2 equations in 3 unknowns it is a 2 by 3 matrix (therefore rectangular). For 6 equations in 6 unknowns the matrix is 6 by 6 (therefore square). The best way to get familiar with matrices is to work with them. Note also the pronunciation: "matrisees" and never "matrixes."

Answer to the practical question  The solution is \( x = 2000, y = 3000 \). That is the intersection point in the row picture (Figure 11.16). It is also the correct combination in the column picture. The matrix equation checks both at once, because matrices are multiplied by rows or by columns. The product either way is \( d \):

\[
\begin{bmatrix}
1 & 1 & 2000 \\
.05 & .10 & 3000
\end{bmatrix}
\begin{bmatrix}
2000 \\
3000
\end{bmatrix}
= \begin{bmatrix}
2000 + 3000 \\
(.05)2000 + (.10)3000
\end{bmatrix}
= \begin{bmatrix}
5000 \\
400
\end{bmatrix}
= d.
\]

Singular case  In the row picture, the lines cross at the solution. But there is a case that gives trouble. When the lines are parallel, they never cross and there is no solution. When the lines are the same, there is an infinity of solutions:

\[
\begin{align*}
2x + y &= 0 \\
2x + y &= 1
\end{align*}
\]

\[
\begin{align*}
4x + 2y &= 0
\end{align*}
\]

This trouble also appears in the column picture. The columns are vectors \( a \) and \( b \). The equation \( Au = d \) is the same as \( xa + yb = d \). We are asked to find the combination of \( a \) and \( b \) (with coefficients \( x \) and \( y \)) that produces \( d \). In the singular case \( a \) and \( b \) lie along the same line (Figure 11.17). No combination can produce \( d \), unless it happens to lie on this line.

![Fig. 11.17  Row and column pictures: singular (no solution) and nonsingular (\( x = y = 1 \)).](image-url)
The investment problem is nonsingular, and 2000 \(a + 3000 \, b\) equals \(d\). We also drew

**Example 2:** The matrix \(A\) multiplies \(u = (1, 1)\) to solve \(x + 2y = 3\) and \(x - y = 0\):

\[
Au = \begin{bmatrix} 1 & 2 \\ 1 & -1 \end{bmatrix} \begin{bmatrix} 1 \\ 1 \end{bmatrix} = \begin{bmatrix} 1+2 \\ 1-1 \end{bmatrix} = \begin{bmatrix} 3 \\ 0 \end{bmatrix}.
\]

By columns \[
\begin{bmatrix} 1 \\ 1 \end{bmatrix} + \begin{bmatrix} 2 \\ -1 \end{bmatrix} = \begin{bmatrix} 3 \\ 0 \end{bmatrix}.
\]

The crossing point is \((1, 1)\) in the row picture. The solution is \(x = 1, \, y = 1\) in the column picture (Figure 11.17b). Then \(1\) times \(a\) plus \(1\) times \(b\) equals the right side \(d\).

**SOLUTION BY DETERMINANTS**

Up to now we just wrote down the answer. The real problem is to find \(x\) and \(y\) when they are unknown. We solve two equations with letters not numbers:

\[
a_1x + b_1y = d_1
\]

\[
a_2x + b_2y = d_2.
\]

The key is to eliminate \(x\). Multiply the first equation by \(a_2\) and the second equation by \(a_1\). Subtract the first from the second and the \(x\)'s disappear:

\[
(a_1b_2 - a_2b_1)y = (a_1d_2 - a_2d_1).
\]

(6)

To eliminate \(y\), subtract \(b_1\) times the second equation from \(b_2\) times the first:

\[
(b_2a_1 - b_1a_2)x = (b_2d_1 - b_1d_2).
\]

(7)

What you see in those parentheses are 2 by 2 determinants! Remember from Section 11.3:

The determinant of

\[
\begin{bmatrix} a_1 & b_1 \\ a_2 & b_2 \end{bmatrix}
\]

is the number

\[
\begin{vmatrix} a_1 & b_1 \\ a_2 & b_2 \end{vmatrix} = a_1b_2 - a_2b_1.
\]

This number appears on the left side of (6) and (7). The right side of (7) is also a determinant—but it has \(d\)'s in place of \(a\)'s. The right side of (6) has \(d\)'s in place of \(b\)'s. So \(x\) and \(y\) are ratios of determinants, given by Cramer's Rule:

\[
11H \quad Cramer's \ Rule \quad \frac{\begin{vmatrix} a_1 & b_1 \\ a_2 & b_2 \end{vmatrix}}{\begin{vmatrix} a_1 & d_1 \\ a_2 & d_2 \end{vmatrix}}, \quad y = \frac{\begin{vmatrix} d_1 & b_1 \\ d_2 & b_2 \end{vmatrix}}{\begin{vmatrix} a_1 & d_1 \\ a_2 & d_2 \end{vmatrix}}.
\]

The investment example is solved by three determinants from the three columns:

\[
\begin{vmatrix} 1 & 1 \\ .05 & .10 \end{vmatrix} = .05, \quad \begin{vmatrix} 5000 & 1 \\ 400 & .10 \end{vmatrix} = 100, \quad \begin{vmatrix} 1 & 5000 \\ .05 & 400 \end{vmatrix} = 150.
\]

Cramer's Rule has \(x = 100/.05 = 2000\) and \(y = 150/.05 = 3000\). This is the solution.

The singular case is when the determinant of \(A\) is zero—and we can't divide by it.

\[
11I \quad Cramer's \ Rule \ breaks \ down \ when \ \text{det} \ \! A = 0—\text{which \ is \ the \ singular \ case.} \quad \text{Then \ the \ lines \ in \ the \ row \ picture \ are \ parallel, \ and \ one \ column \ is \ a \ multiple \ of \ the \ other \ column.}
\]
EXAMPLE 3  The lines $2x + y = 0, 2x + y = 1$ are parallel. The determinant is zero:

$$\begin{bmatrix} 2 & 1 \\ 2 & 1 \end{bmatrix} \begin{bmatrix} x \\ y \end{bmatrix} = \begin{bmatrix} 0 \\ 1 \end{bmatrix} \text{ has } \det A = \begin{vmatrix} 2 & 1 \\ 2 & 1 \end{vmatrix} = 0.$$

The lines in Figure 11.17a don't meet. Notice the columns: $\begin{bmatrix} 1 \\ 1 \end{bmatrix}$ is a multiple of $\begin{bmatrix} 1 \\ 1 \end{bmatrix}$.

One final comment on 2 by 2 systems. They are small enough so that all solution methods apply. Cramer's Rule uses determinants. Larger systems use elimination (3 by 3 matrices are on the borderline). A third solution (the same solution!) comes from the inverse matrix $A^{-1}$, to be described next. But the inverse is more a symbol for the answer than a new way of computing it, because to find $A^{-1}$ we still use determinants or elimination.

THE INVERSE OF A MATRIX

The symbol $A^{-1}$ is pronounced "$A$ inverse." It stands for a matrix—the one that solves $Au = d$. I think of $A$ as a matrix that takes $u$ to $d$. Then $A^{-1}$ is a matrix that takes $d$ back to $u$. If $Au = d$ then $u = A^{-1}d$ (provided the inverse exists). This is exactly like functions and inverse functions: $g(x) = y$ and $x = g^{-1}(y)$. Our goal is to find $A^{-1}$ when we know $A$.

The first approach will be very direct. Cramer's Rule gave formulas for $x$ and $y$, the components of $u$. From that rule we can read off $A^{-1}$, assuming that $D = a_1b_2 - a_2b_1$ is not zero. $D$ is det $A$ and we divide by it:

$$Cramer: \quad u = \frac{1}{D} \begin{bmatrix} b_2d_1 - b_1d_2 \\ -a_2d_1 + a_1d_2 \end{bmatrix} \quad \text{This is } A^{-1}d = \frac{1}{D} \begin{bmatrix} b_2 & -b_1 \\ -a_2 & a_1 \end{bmatrix} \begin{bmatrix} d_1 \\ d_2 \end{bmatrix} \quad (8)$$

The matrix on the right (including $1/D$ in all four entries) is $A^{-1}$. Notice the sign pattern and the subscript pattern. The inverse exists if $D$ is not zero—this is important. Then the solution comes from a matrix-vector multiplication, $A^{-1}$ times $d$. We repeat the rules for that multiplication:

DEFINITION  A matrix $M$ times a vector $v$ equals a vector of dot products:

$$Mv = \begin{bmatrix} \text{row 1} \\ \text{row 2} \end{bmatrix} \begin{bmatrix} v \end{bmatrix} = (\text{row 1}) \cdot v = (\text{row 2}) \cdot v \quad (9)$$

Equation (8) follows this rule with $M = A^{-1}$ and $v = d$. Look at Example 1:

$$A = \begin{bmatrix} 1 & 1 \\ .05 & .10 \end{bmatrix}, \quad \det A = .05, \quad A^{-1} = \frac{1}{.05} \begin{bmatrix} .10 & -1 \\ -.05 & 1 \end{bmatrix} = \begin{bmatrix} 2 & -20 \\ -1 & 20 \end{bmatrix}.$$

There stands the inverse matrix. It multiplies $d$ to give the solution $u$:

$$A^{-1}d = \begin{bmatrix} 2 & -20 \\ -1 & 20 \end{bmatrix} \begin{bmatrix} 5000 \\ 400 \end{bmatrix} = \begin{bmatrix} (2)(5000) - (20)(400) \\ (-1)(5000) + (20)(400) \end{bmatrix} = \begin{bmatrix} 2000 \\ 3000 \end{bmatrix}.$$

The formulas work perfectly, but you have to see a direct way to reach $A^{-1}d$. Multiply both sides of $Au = d$ by $A^{-1}$. The multiplication "cancels" $A$ on the left side, and leaves $u = A^{-1}d$. This approach comes next.
To understand the power of matrices, we must multiply them. The product of $A^{-1}$ with $Au$ is a matrix times a vector. But that multiplication can be done another way. First $A^{-1}$ multiplies $A$, a matrix times a matrix. The product $A^{-1}A$ is another matrix (a very special matrix). Then this new matrix multiplies $u$.

The matrix-matrix rule comes directly from the matrix-vector rule. Effectively, a vector $v$ is a matrix $V$ with only one column. When there are more columns, $M$ times $V$ splits into separate matrix-vector multiplications, side by side:

**DEFINITION** A matrix $M$ times a matrix $V$ equals a matrix of dot products:

$$MV = \begin{bmatrix} \text{row 1} \\ \text{row 2} \end{bmatrix} \begin{bmatrix} v_1 \\ v_2 \end{bmatrix} = \begin{bmatrix} (\text{row 1}) \cdot v_1 \\ (\text{row 1}) \cdot v_2 \\ (\text{row 2}) \cdot v_1 \\ (\text{row 2}) \cdot v_2 \end{bmatrix}.$$ (10)

**EXAMPLE 4**

$$\begin{bmatrix} 1 & 2 \\ 3 & 4 \end{bmatrix} \begin{bmatrix} 5 & 6 \\ 7 & 8 \end{bmatrix} = \begin{bmatrix} 1 \cdot 5 + 2 \cdot 7 \\ 1 \cdot 6 + 2 \cdot 8 \end{bmatrix} = \begin{bmatrix} 19 \\ 22 \end{bmatrix}.$$  

**EXAMPLE 5**

Multiplying $A^{-1}$ times $A$ produces the "identity matrix" $I$,

$$A^{-1}A = \begin{bmatrix} b_2 & -b_1 \\ -a_2 & a_1 \end{bmatrix} \begin{bmatrix} a_1 & b_1 \\ a_2 & b_2 \end{bmatrix} = \begin{bmatrix} a_1b_2 - a_2b_1 & 0 \\ 0 & -a_2b_1 + a_1b_2 \end{bmatrix} = \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix}.$$ (11)

This identity matrix is denoted by $I$. It has 1's on the diagonal and 0's off the diagonal. It acts like the number 1. **Every vector satisfies** $Iu = u$.

**44J (Inverse matrix and identity matrix)** $AA^{-1} = I$ and $A^{-1}A = I$ and $Iu = u$:

$$A = \begin{bmatrix} a & b \\ c & d \end{bmatrix} \quad A^{-1} = \frac{1}{D} \begin{bmatrix} d & -b \\ -c & a \end{bmatrix} \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix} \begin{bmatrix} x \\ y \end{bmatrix} = \begin{bmatrix} x \\ y \end{bmatrix}. \quad (12)$$

Note the placement of $a, b, c, d$. With these letters $D$ is $ad - bc$.

The next section moves to three equations. The algebra gets more complicated (and 4 by 4 is worse). It is not easy to write out $A^{-1}$. So we stay longer with the 2 by 2 formulas, where each step can be checked. Multiplying $Au = d$ by the inverse matrix gives $A^{-1}Au = A^{-1}d$—and the left side is $Iu = u$.

**Fig. 11.18** Rotate $v$ forward into $Av$. Rotate $d$ backward into $A^{-1}d$. 
11.4 Matrices and Linear Equations

**Example 6**  
\[ A = \begin{bmatrix} \cos \theta & -\sin \theta \\ \sin \theta & \cos \theta \end{bmatrix} \]  
rotates every \( \mathbf{v} \) to \( A\mathbf{v} \), through the angle \( \theta \).

**Question 1**  
Where is the vector \( \mathbf{v} = \begin{bmatrix} 1 \\ 0 \end{bmatrix} \) rotated to?

**Question 2**  
What is \( A^{-1} \)?

**Question 3**  
Which vector \( \mathbf{u} \) is rotated into \( \mathbf{d} = \begin{bmatrix} 0 \\ 1 \end{bmatrix} \)?

**Solution 1**  
\( \mathbf{v} \) rotates into \( A\mathbf{v} = \begin{bmatrix} \cos \theta & -\sin \theta \\ \sin \theta & \cos \theta \end{bmatrix} \begin{bmatrix} 1 \\ 0 \end{bmatrix} = \begin{bmatrix} \cos \theta \\ \sin \theta \end{bmatrix} \).

**Solution 2**  
det \( A = 1 \) so \( A^{-1} = \begin{bmatrix} \cos \theta & \sin \theta \\ -\sin \theta & \cos \theta \end{bmatrix} \) = rotation through \(-\theta\).

**Solution 3**  
If \( A\mathbf{u} = \mathbf{d} \) then \( \mathbf{u} = A^{-1}\mathbf{d} = \begin{bmatrix} \cos \theta & \sin \theta \\ -\sin \theta & \cos \theta \end{bmatrix} \begin{bmatrix} 0 \\ 1 \end{bmatrix} = \begin{bmatrix} \sin \theta \\ \cos \theta \end{bmatrix} \).

**Historical note**  
I was amazed to learn that it was Leibniz (again!) who proposed the notation we use for matrices. The entry in row \( i \) and column \( j \) is \( a_{ij} \). The identity matrix has \( a_{11} = a_{22} = 1 \) and \( a_{12} = a_{21} = 0 \). This is in a linear algebra book by Charles Dodgson—better known to the world as Lewis Carroll, the author of *Alice in Wonderland*. I regret to say that he preferred his own notation \( ij \) instead of \( a_{ij} \). “I have turned the symbol toward the left, to avoid all chance of confusion with \( \ell \).” It drove his typesetter mad.

**Projection onto a Plane = Least Squares Fitting by a Line**

We close with a genuine application. It starts with three-dimensional vectors \( \mathbf{a}, \mathbf{b}, \mathbf{d} \) and leads to a 2 by 2 system. One good feature: \( \mathbf{a}, \mathbf{b}, \mathbf{d} \) can be \( n \)-dimensional with no change in the algebra. In practice that happens. Second good feature: There is a calculus problem in the background. The example is to fit points by a straight line.

There are three ways to state the problem, and they look different:

1. Solve \( x\mathbf{a} + y\mathbf{b} = \mathbf{d} \) as well as possible (three equations, two unknowns \( x \) and \( y \)).
2. Project the vector \( \mathbf{d} \) onto the plane of the vectors \( \mathbf{a} \) and \( \mathbf{b} \).
3. Find the closest straight line ("least squares") to three given points.

Figure 11.19 shows a three-dimensional vector \( \mathbf{d} \) above the plane of \( \mathbf{a} \) and \( \mathbf{b} \). Its projection onto the plane is \( \mathbf{p} = x\mathbf{a} + y\mathbf{b} \). The numbers \( x \) and \( y \) are unknown, and our goal is to find them. The calculation will use the dot product, which is always the key to right angles.

The difference \( \mathbf{d} - \mathbf{p} \) is the "error." There has to be an error, because no combination of \( \mathbf{a} \) and \( \mathbf{b} \) can produce \( \mathbf{d} \) exactly. (Otherwise \( \mathbf{d} \) is in the plane.) The projection \( \mathbf{p} \) is the closest point to \( \mathbf{d} \), and it is governed by one fundamental law: The error is perpendicular to the plane. That makes the error perpendicular to both vectors \( \mathbf{a} \) and \( \mathbf{b} \):

\[ \mathbf{a} \cdot (x\mathbf{a} + y\mathbf{b} - \mathbf{d}) = 0 \quad \text{and} \quad \mathbf{b} \cdot (x\mathbf{a} + y\mathbf{b} - \mathbf{d}) = 0. \]  \[ (13) \]
Rewrite those as two equations for the two unknown numbers x and y:

\[(a - a)x + (a - b)y = a \cdot d \]
\[(b - a)x + (b - b)y = b \cdot d. \tag{14}\]

These are the famous normal equations in statistics, to compute x and y and p.

**EXAMPLE 7**  For \(a = (1, 1, 1)\) and \(b = (1, 2, 3)\) and \(d = (0, 5, 4)\), solve equation (14):

\[
\begin{align*}
3x + 6y &= 9 \\
6x + 14y &= 22
\end{align*}
\]

Notice the three equations that we are not solving (we can't): \(xa + yb = d\) is
\[
\begin{align*}
x + y &= 0 \\
x + 2y &= 5 \quad \text{with the 3 by 2 matrix } A = \begin{bmatrix} 1 & 1 \\ 1 & 2 \\ 1 & 3 \end{bmatrix}, \tag{15}\n\end{align*}
\]

For \(d = (0, 5, 4)\) there is no solution; \(d\) is not in the plane of \(a\) and \(b\). For \(p = (1, 3, 5)\) there is a solution, \(x = -1\) and \(y = 2\). The vector \(p\) is in the plane. The error \(d - p\) is \((-1, 2, -1)\). This error is perpendicular to the columns \((1, 1, 1)\) and \((1, 2, 3)\), so it is perpendicular to their plane.

**SAME EXAMPLE**  (written as a line-fitting problem) Fit the points \((1, 0)\) and \((2, 5)\) and \((3, 4)\) as closely as possible ("least squares") by a straight line.

Two points determine a line. The example asks the line \(f = x + yt\) to go through three points. That gives the three equations in (15), which can't be solved with two unknowns. We have to settle for the closest line, drawn in Figure 11.19b. This line is computed again below, by calculus.

Notice that the closest line has heights 1, 3, 5 where the data points have heights 0, 5, 4. Those are the numbers in \(p\) and \(d\)! The heights 1, 3, 5 fit onto a line; the heights 0, 5, 4 do not. In the first figure, \(p = (1, 3, 5)\) is in the plane and \(d = (0, 5, 4)\) is not. Vectors in the plane lead to heights that lie on a line.

Notice another coincidence. The coefficients \(x = -1\) and \(y = 2\) give the projection \(-a + 2b\). They also give the closest line \(f = -1 + 2t\). All numbers appear in both figures.

![Fig. 11.19 Projection onto plane is (1, 3, 5) with coefficients -1, 2. Closest line has heights 1, 3, 5 with coefficients -1, 2. Error in both pictures is -1, 2, -1.](image-url)
Remark  Finding the closest line is a calculus problem: Minimize a sum of squares. The numbers \( x \) and \( y \) that minimize \( E \) give the least squares solution:

\[
E(x, y) = (x + y - 0)^2 + (x + 2y - 5)^2 + (x + 3y - 4)^2.
\]  

(16)

Those are the three errors in equation (15), squared and added. They are also the three errors in the straight line fit, between the line and the data points. The projection minimizes the error (by geometry), the normal equations (14) minimize the error (by algebra), and now calculus minimizes the error by setting the derivatives of \( E \) to zero.

The new feature is this: \( E \) depends on two variables \( x \) and \( y \). Therefore \( E \) has two derivatives. They both have to be zero at the minimum. That gives two equations for \( x \) and \( y \):

\[
\text{\( x \) derivative of \( E \) is zero:} \quad 2(x + y) + 2(x + 2y - 5) + 2(x + 3y - 4) = 0
\]

\[
\text{\( y \) derivative of \( E \) is zero:} \quad 2(x + y) + 2(x + 2y - 5)(2) + 2(x + 3y - 4)(3) = 0
\]

When we divide by 2, those are the normal equations

\[
3x + 6y = 9 \quad \text{and} \quad 6x + 14y = 22.
\]

The minimizing \( x \) and \( y \) from calculus are the same numbers \(-1\) and \( 2 \).

The \( x \) derivative treats \( y \) as a constant. The \( y \) derivative treats \( x \) as a constant. These are partial derivatives. This calculus approach to least squares is in Chapter 13, as an important application of partial derivatives.

We now summarize the least squares problem—to find the closest line to \( n \) data points. In practice \( n \) may be 1000 instead of 3. The points have horizontal coordinates \( b_1, b_2, ..., b_n \). These vectors \( \mathbf{b} \) and \( \mathbf{d} \), together with \( \mathbf{a} = (1, 1, ..., 1) \), determine a projection—the combination \( \mathbf{p} = x\mathbf{a} + y\mathbf{b} \) that is closest to \( \mathbf{d} \). This problem is the same in \( n \) dimensions—the error \( \mathbf{d} - \mathbf{p} \) is perpendicular to \( \mathbf{a} \) and \( \mathbf{b} \). That is still tested by dot products, \( \mathbf{p} \cdot \mathbf{a} = \mathbf{d} \cdot \mathbf{a} \) and \( \mathbf{p} \cdot \mathbf{b} = \mathbf{d} \cdot \mathbf{b} \), which give the normal equations for \( x \) and \( y \):

\[
(a \cdot a)x + (a \cdot b)y = a \cdot d \quad \text{or} \quad (n) \quad x + (\Sigma b_i)y = \Sigma d_i
\]

\[
(b \cdot a)x + (b \cdot b)y = b \cdot d \quad \text{or} \quad (\Sigma b_i)x + (\Sigma b_i^2)y = \Sigma b_i d_i.
\]  

(17)

11K The least squares problem projects \( \mathbf{d} \) onto the plane of \( \mathbf{a} \) and \( \mathbf{b} \). The projection is \( \mathbf{p} = x\mathbf{a} + y\mathbf{b} \), in \( n \) dimensions. The closest line is \( f = x + yt \), in two dimensions. The normal equations (17) give the best \( x \) and \( y \).
A matrix–matrix multiplication $MV$ yields a matrix of dot products, from the rows of $M$ and the columns of $V$:

$$
\begin{bmatrix}
A \\
B
\end{bmatrix}
\begin{bmatrix}
v_1 \\
v_2
\end{bmatrix} =
\begin{bmatrix}
3 & 1 \\
1 & 1
\end{bmatrix}
\begin{bmatrix}
1/2 \\
-1/2
\end{bmatrix} =
\begin{bmatrix}
1/2 \\
-1/2
\end{bmatrix}.
\]

The last line contains the $u$ matrix, denoted by $I$. It has the property that $IA = AI = v$ for every matrix $A$, and $Iu = u$ for every vector $u$. The inverse matrix satisfies $A^{-1}A = I$. Then $Au = d$ is solved by multiplying both sides by $v$, to give $u = w$. There is no inverse matrix when $A$ is singular.

The combination $ra + sb$ is the projection of $d$ when the error $b$ is perpendicular to $c$. If $a = (1, 1, 1), b = (1, 2, 3)$, and $d = (0, 8, 4)$, the equations for $x$ and $y$ are $e = f = 0$. Solving them also gives the closest $t$ to the data points $(1, 0), (2, 0), (3, 4)$. The solution is $x = 0, y = 2$, which means the best line is $H$. The projection is $0a + 2b = -1$. The three error components are $0, 1, 1$. Check perpendicularity: $0 = 0$ and $1 = 1 = 1$. Applying calculus to this problem, $x$ and $y$ minimize the sum of squares $E = M$.

In 1–8 find the point $(x, y)$ where the two lines intersect (if they do). Also show how the right side is a combination of the columns on the left side (if it is). Also find the determinant $D$.

1. $x + y = 7$
   $x - y = 3$
2. $2x + y = 11$
   $x + y = 6$
3. $3x - y = 8$
   $x - 3y = 0$
4. $2x + y = 3$
   $2x + y = 7$
5. $2x - 4y = 0$
   $x - 2y = 0$
6. $10x + y = 1$
   $x + y = 1$
7. $ax + by = 0$
   $2ax + by = 2$
8. $ax + by = 1$
   $cx + dy = 1$

9 Solve Problem 3 by Cramer's Rule.

10 Try to solve Problem 4 by Cramer's Rule.

11 What are the ratios for Cramer's Rule in Problem 5?

12 If $A = I$ show how Cramer's Rule solves $Au = d$.

13 Draw the row picture and column picture for Problem 1.

14 Draw the row and column pictures for Problem 6.

15 Find $A^{-1}$ in Problem 1.

16 Find $A^{-1}$ in Problem 8 if $ad - bc = 1$.

17 A 2 by 2 system is singular when the two lines in the row picture ________. This system is still solvable if one equation is a ________ of the other equation. In that case the two lines are ________ and the number of solutions is ________.

18 Try Cramer's Rule when there is no solution or infinitely many:

$$
\begin{align*}
3x + y &= 0 \\
3x + y &= 1
\end{align*}
$$

or

$$
\begin{align*}
6x + 2y &= 2 \\
6x + 2y &= 2.
\end{align*}
$$

19 $Au = d$ is singular when the columns of $A$ are ________. A solution exists if the right side $d$ is ________. In this solvable case the number of solutions is ________.

20 The equations $x - y = d_1$ and $9x - 9y = d_2$ can be solved if ________.

21 Suppose $x = \frac{1}{2}$ billion people live in the U.S. and $y = 5$ billion live outside. If 4 per cent of those inside move out and 2 per cent of those outside move in, find the populations $d_1$ inside and $d_2$ outside after the move. Express this as a matrix multiplication $Au = d$ (and find the matrix).

22 In Problem 21 what is special about $a_1 + a_2$ and $b_1 + b_2$ (the sums down the columns of $A$)? Explain why $d_1 + d_2$ equals $x + y$.

23 With the same percentages moving, suppose $d_1 = 0.58$ billion are inside and $d_2 = 4.92$ billion are outside at the end. Set up and solve two equations for the original populations $x$ and $y$.

24 What is the determinant of $A$ in Problems 21–23? What is $A^{-1}$. Check that $A^{-1}A = I$.

25 The equations $ax + y = 0$, $x + ay = 0$ have the solution $x = y = 0$. For which two values of $a$ are there other solutions (and what are the other solutions)?

26 The equations $ax + by = 0$, $cx + dy = 0$ have the solution $x = y = 0$. There are other solutions if the two lines are ________. This happens if $a, b, c, d$ satisfy ________.

27 Find the determinant and inverse of $A = \begin{bmatrix} 3 & 4 \\ 2 & -1 \end{bmatrix}$. Do the same for $2A, A^{-1}$, $-A$, and $I$.

28 Show that the determinant of $A^{-1}$ is $1/\det A$:

$$
A^{-1} = \begin{bmatrix}
d(ad - bc) & -b(ad - bc) \\
-c(ad - bc) & a(ad - bc)
\end{bmatrix}.
$$

29 Compute $AB$ and $BA$ and also $BC$ and $CB$:

$$
A = \begin{bmatrix}
1 & 4 \\
2 & -1
\end{bmatrix},
B = \begin{bmatrix}
3 & 1 \\
1 & 1
\end{bmatrix},
C = \begin{bmatrix}
0 & 1 \\
2 & 0
\end{bmatrix}.
$$

Verify the associative law: $AB$ times $C$ equals $A$ times $BC$.

30 (a) Find the determinants of $A, B, AB$, and $BA$ above.
   (b) Propose a law for the determinant of $BC$ and test it.

31 For $A = \begin{bmatrix} a & b \\ c & d \end{bmatrix}$ and $B = \begin{bmatrix} e & f \\ g & h \end{bmatrix}$ write out $AB$ and factor its determinant into $(ad - bc)(eh - fg)$. Therefore $\det(AB) = (\det A)(\det B)$. 

32 Usually $\det(A + B)$ does not equal $\det A + \det B$. Find examples of inequality and equality.

33 Find the inverses, and check $A^{-1}A = I$ and $BB^{-1} = I$, for

$$A = \begin{bmatrix} 1 & 4 \\ 0 & 2 \end{bmatrix} \text{ and } B = \begin{bmatrix} 2 & 2 \\ 0 & 1 \end{bmatrix}.$$  

34 In Problem 33 compute $AB$ and the inverse of $AB$. Check that this inverse equals $B^{-1}$ times $A^{-1}$.

35 The matrix product $AB^{-1}A^{-1}$ equals the _______ matrix. Therefore the inverse of $AB$ is _______. Important: The associative law in Problem 29 allows you to multiply $BB^{-1}$ first.

36 The matrix multiplication $C^{-1}B^{-1}A^{-1}ABC$ yields the _______ matrix. Therefore the inverse of $ABC$ is _______.

37 The equations $x + 2y + 3z$ and $4x + 5y + cz = 0$ always have a nonzero solution. The vector $u = (x, y, z)$ is required to be _______ to $v = (1, 2, 3)$ and $w = (4, 5, c)$. So choose $u = _______$. 

38 Find the combination $p = xa + yb$ of the vectors $a = (1, 1, 1)$ and $b = (-1, 0, 1)$ that comes closest to $d = (2, 6, 4)$. (a) Solve the normal equations (14) for $x$ and $y$. (b) Check that the error $d - p$ is perpendicular to $a$ and $b$.

39 Plot the three data points $(-1, 2), (0, 6), (1, 4)$ in a plane. Draw the straight line $x + yt$ with the same $x$ and $y$ as in Problem 38. Locate the three errors up or down from the data points and compare with Problem 38.

40 Solve equation (14) to find the combination $xa + yb$ of $a = (1, 1, 1)$ and $b = (-1, 1, 2)$ that is closest to $d = (1, 1, 3)$. Draw the corresponding straight line for the data points $(-1, 1), (1, 1)$, and $(2, 3)$. What is the vector of three errors and what is it perpendicular to?

41 Under what condition on $d_1, d_2, d_3$ do the three points $(0, d_1), (1, d_2), (2, d_3)$ lie on a line?

42 Find the matrices that reverse $x$ and $y$ and project:

$$M \begin{bmatrix} x \\ y \end{bmatrix} = \begin{bmatrix} y \\ x \end{bmatrix} \text{ and } P \begin{bmatrix} x \\ y \end{bmatrix} = \begin{bmatrix} x^7 \\ 0 \end{bmatrix}.$$  

43 Multiplying by $P = \begin{bmatrix} .5 & .5 \\ .5 & .5 \end{bmatrix}$ projects $u$ onto the 45° line.

(a) Find the projection $Pu$ of $u = [1]$. 
(b) Why does $P$ times $P$ equal $P$? 
(c) Does $P^{-1}$ exist? What vectors give $Pu = 0$?

44 Suppose $u$ is not the zero vector but $Au = 0$. Then $A^{-1}$ can't exist: It would multiply _______ and produce $u$.

### 11.5 Linear Algebra

This section moves from two to three dimensions. There are three unknowns $x, y, z$ and also three equations. This is at the crossover point between formulas and algorithms—it is real linear algebra. The formulas give a direct solution using determinants. The algorithms use elimination and the numbers $x, y, z$ appear at the end. In practice that end result comes quickly. Computers solve linear equations by elimination.

The situation for a nonlinear equation is similar. Quadratic equations $ax^2 + bx + c = 0$ are solved by a formula. Cubic equations are solved by Newton's method (even though a formula exists). For equations involving $x^4$ or $x^6$, algorithms take over completely.

Since we are at the crossover point, we look both ways. This section has a lot to do, in mixing geometry, determinants, and 3 by 3 matrices:

1. The row picture: three planes intersect at the solution  
2. The column picture: a vector equation combines the columns  
3. The formulas: determinants and Cramer's Rule  
4. Matrix multiplication and $A^{-1}$  
5. The algorithm: Gaussian elimination.

Part of our goal is three-dimensional calculus. Another part is $n$-dimensional algebra. And a third possibility is that you may not take mathematics next year. If that
11 Vectors and Matrices

happens, I hope you will use mathematics. Linear equations are so basic and important, in such a variety of applications, that the effort in this section is worth making.

An example is needed. It is convenient and realistic if the matrix contains zeros. Most equations in practice are fairly simple—a thousand equations each with 990 zeros would be very reasonable. Here are three equations in three unknowns:

\[
\begin{align*}
  x + y &= 1 \\
  x + 2z &= 0 \\
  -2y + 2z &= -4.
\end{align*}
\]  

In matrix-vector form, the unknown \( u \) has components \( x, y, z \). The right sides \( 1, 0, -4 \) go into \( d \). The nine coefficients, including three zeros, enter the matrix \( A \):

\[
\begin{bmatrix}
  1 & 1 & 0 \\
  1 & 0 & 2 \\
  0 & -2 & 2
\end{bmatrix}
\begin{bmatrix}
  x \\
  y \\
  z
\end{bmatrix}
= 
\begin{bmatrix}
  1 \\
  0 \\
  -4
\end{bmatrix}
\text{ or } Au = d. 
\]  

The goal is to understand that system geometrically, and then solve it.

**THE ROW PICTURE: INTERSECTING PLANES**

Start with the first equation \( x + y = 1 \). In the \( xy \) plane that produces a line. In three dimensions it is a plane. It has the usual form \( ax + by + cz = d \), except that \( c \) happens to be zero. The plane is easy to visualize (Figure 11.20a), because it cuts straight down through the line. The equation \( x + y = 1 \) allows \( z \) to have any value, so the graph includes all points above and below the line.

The second equation \( x + 2z = 0 \) gives a second plane, which goes through the origin. When the right side is zero, the point \( (0, 0, 0) \) satisfies the equation. This time \( y \) is absent from the equation, so the plane contains the whole \( y \) axis. All points \( (0, y, 0) \) meet the requirement \( x + 2z = 0 \). The normal vector to the plane is \( N = i + 2k \). The plane cuts across, rather than down, in 11.20b.

Before the third equation we combine the first two. The intersection of two planes *is a line*. In three-dimensional space, two equations (not one) describe a line. The points on the line have to satisfy \( x + y = 1 \) and also \( x + 2z = 0 \). A convenient point is \( P = (0, 1, 0) \). Another point is \( Q = (-1, 2, \frac{1}{2}) \). The line through \( P \) and \( Q \) extends out in both directions.

The solution is on that line. The third plane decides where.

Fig. 11.20 First plane, second plane, intersection line meets third plane at solution.
11.5 Linear Algebra

The third equation \(-2y + 2z = -4\) gives the third plane—which misses the origin because the right side is not zero. What is important is the point where the three planes meet. The intersection line of the first two planes crosses the third plane. We used determinants (but elimination is better) to find \(x = -2, y = 3, z = 1\). This solution satisfies the three equations and lies on the three planes.

A brief comment on 4 by 4 systems. The first equation might be \(x + y + z - t = 0\). It represents a three-dimensional "hyperplane" in four-dimensional space. (In physics this is space-time.) The second equation gives a second hyperplane, and its intersection with the first one is two-dimensional. The third equation (third hyperplane) reduces the intersection to a line. The fourth hyperplane meets that line at a point, which is the solution. It satisfies the four equations and lies on the four hyperplanes. In this course three dimensions are enough.

COLUMN PICTURE: COMBINATION OF COLUMN VECTORS

There is an extremely important way to rewrite our three equations. In (1) they were separate, in (2) they went into a matrix. Now they become a vector equation:

\[
\begin{bmatrix}
1 \\
1 \\
0
\end{bmatrix} x + \begin{bmatrix}
1 \\
0 \\
-2
\end{bmatrix} y + \begin{bmatrix}
2 \\
2 \\
-4
\end{bmatrix} z = \begin{bmatrix}
1 \\
0 \\
-4
\end{bmatrix}
\]

(3)

The columns of the matrix are multiplied by \(x, y, z\). That is a special way to see matrix-vector multiplication: \(Au\) is a combination of the columns of \(A\). We are looking for the numbers \(x, y, z\) so that the combination produces the right side \(d\).

The column vectors \(a, b, c\) are shown in Figure 11.21a. The vector equation is \(xa + yb + xc = d\). The combination that solves this equation must again be \(x = -2, y = 3, z = 1\). That agrees with the intersection point of the three planes in the row picture.

\[c = \begin{bmatrix}
0 \\
2 \\
2
\end{bmatrix}
\]

\[b = \begin{bmatrix}
1 \\
0 \\
-2
\end{bmatrix}
\]

\[a = \begin{bmatrix}
1 \\
1 \\
0
\end{bmatrix}
\]

\[d = \begin{bmatrix}
1 \\
0 \\
-4
\end{bmatrix}
\]

\[d = 1c + 3b - 2a
\]

The determinant and the inverse matrix

For a 3 by 3 determinant, the section on cross products gave two formulas. One was the triple product \(a \cdot (b \times c)\). The other wrote out the six terms:

\[\det A = a_1(b_2c_3 - b_3c_2) + a_2(b_3c_1 - b_1c_3) + a_3(b_1c_2 - b_2c_1).\]
Geometrically this is the volume of a box. The columns a, b, c are the edges going out from the origin. In our example the determinant and volume are 2:

\[
\begin{vmatrix}
  a_1 & b_1 & c_1 \\
  a_2 & b_2 & c_2 \\
  a_3 & b_3 & c_3
\end{vmatrix} = 1 \cdot 1 \cdot 0 - 1 \cdot 2 \cdot 2 + 1 \cdot 2 \cdot 0 = 2.
\]

A slight dishonesty is present in that calculation, and will be admitted now. In Section 11.3 the vectors A, B, C were rows. In this section a, b, c are columns. It doesn't matter, because the determinant is the same either way. Any matrix can be "transposed"—exchanging rows for columns—without altering the determinant. The six terms \((a_1b_2c_3)\) is the first) may come in a different order, but they are the same six terms. Here four of those terms are zero, because of the zeros in the matrix. The sum of all six terms is \(D = \det A = 2\).

Since \(D\) is not zero, the equations can be solved. The three planes meet at a point. The column vectors a, b, c produce a genuine box, and are not flattened into the same plane (with zero volume). The solution involves dividing by \(D\)—which is only possible if \(D = \det A\) is not zero.

When the determinant \(D\) is not zero, \(A\) has an inverse: \(AA^{-1} = A^{-1}A = I\). Then the equations \(Au = d\) have one and only one solution \(u = A^{-1}d\).

The 3 by 3 identity matrix \(I\) is at the end of equation (5). Always \(Iu = u\).

We now compute \(A^{-1}\), first with letters and then with numbers. The neatest formula uses cross products of the columns of \(A\)—it is special for 3 by 3 matrices.

\[\text{Every entry is divided by } D:\quad \text{The inverse matrix is } A^{-1} = \frac{1}{D} \begin{bmatrix} b \times c \\ c \times a \\ a \times b \end{bmatrix}. \quad (4)\]

To test this formula, multiply by \(A\). Matrix multiplication produces a matrix of dot products—from the rows of the first matrix and the columns of the second, \(A^{-1}A = I\):

\[
\begin{bmatrix}
  b \times c \\
  c \times a \\
  a \times b
\end{bmatrix}
\begin{bmatrix}
  a & b & c
\end{bmatrix} = \frac{1}{D}
\begin{bmatrix}
  a \cdot (b \times c) & b \cdot (b \times c) & c \cdot (b \times c) \\
  a \cdot (c \times a) & b \cdot (c \times a) & c \cdot (c \times a) \\
  a \cdot (a \times b) & b \cdot (a \times b) & c \cdot (a \times b)
\end{bmatrix} = \begin{bmatrix} 1 & 0 & 0 \\ 0 & 1 & 0 \\ 0 & 0 & 1 \end{bmatrix}. \quad (5)
\]

On the right side, six of the triple products are zero. They are the off-diagonals like \(b \cdot (b \times c)\), which contain the same vector twice. Since \(b \times c\) is perpendicular to \(b\), this triple product is zero. The same is true of the others, like \(a \cdot (a \times b) = 0\). That is the volume of a box with two identical sides. The six off-diagonal zeros are the volumes of completely flattened boxes.

On the main diagonal the triple products equal \(D\). The order of vectors can be abc or bca or cab, and the volume of the box stays the same. Dividing by this number \(D\), which is placed outside for that purpose, gives the 1's in the identity matrix \(I\).

Now we change to numbers. The goal is to find \(A^{-1}\) and to test it.

**Example 1** The inverse of \(A =\)

\[
\begin{bmatrix}
  1 & 1 & 0 \\
  1 & 0 & 2 \\
  0 & -2 & 2
\end{bmatrix}
\]
is \(A^{-1} =\)

\[
\begin{bmatrix}
  4 & -2 & 2 \\
 -2 & 2 & -2 \\
 -2 & 2 & -1
\end{bmatrix}
\].
That comes from the formula, and it absolutely has to be checked. Do not fail to multiply $A^{-1}$ times $A$ (or $A$ times $A^{-1}$). Matrix multiplication is much easier than the formula for $A^{-1}$. We highlight row 3 times column 1, with dot product zero:

$$
\begin{array}{ccc}
1 & 1 & 0 \\
2 & 0 & 2 \\
1 & 2 & 2 \\
\end{array}
\begin{array}{ccc}
4 & -2 & 2 \\
-2 & 2 & -2 \\
-2 & 2 & -1 \\
\end{array}
= \frac{1}{2}
\begin{array}{ccc}
4-2 & 4-4 & -4+4 \\
2-2+2 & 2-2+4 & 4-4 \\
2-2+2 & 2-2+4 & 4-2 \\
\end{array}
= \begin{array}{ccc}
1 & 0 & 0 \\
0 & 1 & 0 \\
0 & 0 & 1 \\
\end{array}.
$$

**Remark on $A^{-1}$** Inverting a matrix requires $D \neq 0$. We divide by $D = \det A$. The cross products $b \times c$ and $c \times a$ and $a \times b$ give $A^{-1}$ in a neat form, but errors are easy. We prefer to avoid writing $i, j, k$. There are nine $2 \times 2$ determinants to be calculated, and here is $A^{-1}$ in full—containing the nine "cofactors" divided by $D$:

$$
A^{-1} = \frac{1}{D}
\begin{vmatrix}
 b_2c_3 - b_3c_2 & b_3c_1 - b_1c_3 & b_1c_2 - b_2c_1 \\
 c_2a_3 - c_3a_2 & c_3a_1 - c_1a_3 & c_1a_2 - c_2a_1 \\
 a_2b_3 - a_3b_2 & a_3b_1 - a_1b_3 & a_1b_2 - a_2b_1 \\
\end{vmatrix}.
$$

**Important:** The first row of $A^{-1}$ does not use the first column of $A$, except in $1/D$. In other words, $b \times c$ does not involve $a$. Here are the $2 \times 2$ determinants that produce $4, -2, 2$—which is divided by $D = 2$ in the top row of $A^{-1}$:

$$
\begin{array}{ccc}
1 & 1 & 0 \\
1 & 0 & 2 \\
0 & -2 & 2 \\
\end{array}
\begin{array}{ccc}
1 & 1 & 0 \\
1 & 0 & 2 \\
0 & -2 & 2 \\
\end{array}
= \begin{array}{ccc}
+ & - & + \\
- & + & - \\
+ & - & + \\
\end{array}.
$$

The second highlighted determinant looks like +2 not −2. But the *sign matrix* on the right assigns a minus to that position in $A^{-1}$. We reverse the sign of $b_1c_3 - b_2c_1$, to find the cofactor $b_3c_1 - b_1c_3$ in the top row of (6).

To repeat: For a row of $A^{-1}$, cross out the corresponding column of $A$. Find the three $2 \times 2$ determinants, use the sign matrix, and divide by $D$.

**Example 2** $B =
\begin{bmatrix}
1 & 1 \\
0 & 1 \\
0 & 0 \\
\end{bmatrix}$ has $D = 1$ and $B^{-1} =
\begin{bmatrix}
1 & -1 & 0 \\
0 & 1 & -1 \\
0 & 0 & 1 \\
\end{bmatrix}$.

The multiplication $BB^{-1} = I$ checks the arithmetic. Notice how $\frac{1}{1}$ in $B$ leads to a zero in the top row of $B^{-1}$. To find row 1, column 3 of $B^{-1}$ we ignore column 1 and row 3 of $B$. (Also: the inverse of a triangular matrix is triangular.) The minus signs come from the sign matrix.

**The solution $u = A^{-1}d$**

The purpose of $A^{-1}$ is to solve the equation $Au = d$. Multiplying by $A^{-1}$ produces $Iu = A^{-1}d$. The matrix becomes the identity, $Iu$ equals $u$, and the solution is immediate:

$$
u = A^{-1}d = \frac{1}{D}
\begin{bmatrix}
 b \times c \\
 c \times a \\
 a \times b \\
\end{bmatrix}
\begin{bmatrix}
d \cdot (b \times c) \\
 d \cdot (c \times a) \\
 d \cdot (a \times b) \\
\end{bmatrix}.
$$

$$
\begin{bmatrix}
 b \times c \\
 c \times a \\
 a \times b \\
\end{bmatrix}
\begin{bmatrix}
d \\
 d \cdot (b \times c) \\
 d \cdot (c \times a) \\
 d \cdot (a \times b) \\
\end{bmatrix}.
$$

$$
\begin{bmatrix}
 b \times c \\
 c \times a \\
 a \times b \\
\end{bmatrix}
\begin{bmatrix}
d \cdot (b \times c) \\
 d \cdot (c \times a) \\
 d \cdot (a \times b) \\
\end{bmatrix}.
$$
By writing those components \( x, y, z \) as ratios of determinants, we have Cramer's Rule:

\[
\begin{align*}
11 \text{(Cramer's Rule)} \\
\text{The solution is } & x = \frac{|d \ b \ c|}{|a \ b \ c|}, \quad y = \frac{|a \ d \ c|}{|a \ b \ c|}, \quad z = \frac{|a \ b \ d|}{|a \ b \ c|}. \quad (10)
\end{align*}
\]

The right side \( d \) replaces, in turn, columns \( a \) and \( b \) and \( c \). All denominators are \( D = a \cdot (b \times c) \). The numerator of \( x \) is the determinant \( d \cdot (b \times c) \) in (9). The second numerator agrees with the second component \( d \cdot (c \times a) \), because the cyclic order is correct. The third determinant with columns \( abd \) equals the triple product \( d \cdot (a \times b) \) in \( A^{-1}u \). Thus (10) is the same as (9).

**EXAMPLE A:** Multiply by \( A^{-1} \) to find the known solution \( x = -2, y = 3, z = 1 \):

\[
\begin{align*}
A^{-1}d = & \frac{1}{2} \begin{bmatrix}
4 & -2 & 2 \\
-2 & 2 & -2 \\
-2 & 2 & -1
\end{bmatrix} \\
& = \frac{1}{2} \begin{bmatrix}
4 - 8 \\
-2 + 8 \\
-2 + 4
\end{bmatrix} = \begin{bmatrix}
-2 \\
3 \\
1
\end{bmatrix}.
\end{align*}
\]

**EXAMPLE B:** Multiply by \( B^{-1} \) to solve \( Bu = d \) when \( d \) is the column \( (6, 5, 4) \):

\[
\begin{align*}
u = B^{-1}d &= \begin{bmatrix}
1 & -1 & 0 \\
0 & 1 & -1 \\
0 & 0 & 1
\end{bmatrix} \begin{bmatrix}
6 \\
5 \\
4
\end{bmatrix} = \begin{bmatrix}
1 & 1 & 1 \\
0 & 1 & 1 \\
0 & 0 & 1
\end{bmatrix} = \begin{bmatrix}
6 \\
5 \\
4
\end{bmatrix}.
\end{align*}
\]

**EXAMPLE C:** Put \( d = (6, 5, 4) \) in each column of \( B \). Cramer's Rule gives \( u = (1, 1, 4) \):

\[
\begin{align*}
\begin{vmatrix}
6 & 1 & 1 \\
5 & 1 & 1 \\
4 & 0 & 1
\end{vmatrix} = & 1 \begin{vmatrix}
1 & 1 & 1 \\
0 & 1 & 1 \\
0 & 0 & 1
\end{vmatrix} = 1 \begin{vmatrix}
1 & 1 & 1 \\
0 & 1 & 5 \\
0 & 0 & 4
\end{vmatrix} = 4 \quad \text{all divided by } D = 1.
\end{align*}
\]

This rule fills the page with determinants. Those are good ones to check by eye, without writing down the six terms (three + and three -).

The formulas for \( A^{-1} \) are honored chiefly in their absence. They are not used by the computer, even though the algebra is in some ways beautiful. In big calculations, the computer never finds \( A^{-1} \)—just the solution.

We now look at the singular case \( D = 0 \). Geometry-algebra-algorithm must all break down. After that is the algorithm: Gaussian elimination.

**THE SINGULAR CASE**

Changing one entry of a matrix can make the determinant zero. The triple product \( a \cdot (b \times c) \), which is also the volume, becomes \( D = 0 \). The box is flattened and the matrix is singular. That happens in our example when the lower right entry is changed from 2 to 4:

\[
\begin{bmatrix}
1 & 1 & 0 \\
1 & 0 & 2 \\
0 & -2 & 4
\end{bmatrix}
\]

has determinant \( D = 0 \).
11.5 Linear Algebra

This does more than change the inverse. It destroys the inverse. We can no longer divide by \( D \). There is no \( S^{-1} \).

What happens to the row picture and column picture? For 2 by 2 systems, the singular case had two parallel lines. Now the row picture has three planes, which need not be parallel. Here the planes are not parallel. Their normal vectors are the rows of \( S \), which go in different directions. But somehow the planes fail to go through a common point.

What happens is more subtle. The intersection line from two planes misses the third plane. The line is parallel to the plane and stays above it (Figure 11.22a). When all three planes are drawn, they form an open tunnel. The picture tells more than the numbers, about how three planes can fail to meet. The third figure shows an end view, where the planes go directly into the page. Each pair meets in a line, but those lines don’t meet in a point.

![Fig. 11.22](image)

The row picture in the singular case: no intersection point, no solution.

When two planes are parallel, the determinant is again zero. One row of the matrix is a multiple of another row. The extreme case has all three planes parallel—as in a matrix with nine 1’s.

The column picture must also break down. In the 2 by 2 failure (previous section), the columns were on the same line. Now the three columns are in the same plane. The combinations of those columns produce \( d \) only if it happens to lie in that particular plane. Most vectors \( d \) will be outside the plane, so most singular systems have no solution.

When the determinant is zero, \( Au = d \) has no solution or infinitely many.

THE ELIMINATION ALGORITHM

Go back to the 3 by 3 example \( Au = d \). If you were given those equations, you would never think of determinants. You would—quite correctly—start with the first equation. It gives \( x = 1 - y \), which goes into the next equation to eliminate \( x \):

\[
\begin{align*}
x + y &= 1 \\
x + 2z &= 0 \\
-2y + 2z &= -4
\end{align*}
\]

Stop there for a minute. On the right is a 2 by 2 system for \( y \) and \( z \). The first equation and first unknown are eliminated—exactly what we want. But that step was not organized in the best way, because a “1” ended up on the left side. Constants should stay on the right side—the pattern should be preserved. It is better to take the same
step by subtracting the first equation from the second:

\[
\begin{align*}
    x + y &= 1 \\
    x + 2z &= 0 \quad \rightarrow \quad -y + 2z &= -1 \\
    -2y + 2z &= -4 \quad \rightarrow \quad -2y + 2z &= -4.
\end{align*}
\] (11)

Same equations, better organization. Now look at the corner term \(-y\). Its coefficient \(-1\) is the second pivot. (The first pivot was \(+1\), the coefficient of \(x\) in the first corner.)

We are ready for the next elimination step:

**Plan:** Subtract a multiple of the "pivot equation" from the equation below it.

**Goal:** To produce a zero below the pivot, so \(y\) is eliminated.

**Method:** Subtract 2 times the pivot equation to cancel \(-2y\).

\[
\begin{align*}
    -y + 2z &= -1 \\
    -2y + 2z &= -4 \quad \rightarrow \quad -2z &= -2.
\end{align*}
\] (12)

The answer comes by back substitution. Equation (12) gives \(z = 1\). Then equation (11) gives \(y = 3\). Then the first equation gives \(x = -2\). This is much quicker than determinants. You may ask: Why use Cramer's Rule? Good question.

With numbers elimination is better. It is faster and also safer. (To check against error, substitute \(-2, 3, 1\) into the original equations.) The algorithm reaches the answer without the determinant and without the inverse. Calculations with letters use \(\text{det } A\) and \(A^{-1}\).

Here are the steps in a definite order (top to bottom):

1. Subtract a multiple of equation 1 to produce \(0x\) in equation 2
2. Subtract a multiple of equation 1 to produce \(0x\) in equation 3
3. Subtract a multiple of equation 2 (new) to produce \(0y\) in equation 3.

**EXAMPLE** (notice the zeros appearing under the pivots):

\[
\begin{align*}
    x + y + z &= 1 \\
    2x + 5y + 3z &= 7 \quad \rightarrow \quad 3y + z &= 5 \quad \rightarrow \quad 3y + z &= 5 \\
    4x + 7y + 6z &= 11 \quad \rightarrow \quad 3y + 2z &= 7 \quad \rightarrow \quad z &= 2.
\end{align*}
\]

Elimination leads to a triangular system. The coefficients below the diagonal are zero. First \(z = 2\), then \(y = 1\), then \(x = -2\). Back substitution solves triangular systems (fast).

As a final example, try the singular case \(Su = d\) when the corner entry is changed from 2 to 4. With \(D = 0\), there is no inverse matrix \(S^{-1}\). Elimination also fails, by reaching an impossible equation \(0 = -2\):

\[
\begin{align*}
    x + y &= 1 \\
    x + 2z &= 0 \quad \rightarrow \quad -y + 2z &= -1 \quad \rightarrow \quad -y + 2z &= -1 \\
    -2y + 4z &= -4 \quad \rightarrow \quad -2y + 4z &= -4 \quad \rightarrow \quad 0 &= -2
\end{align*}
\]

The three planes do not meet at a point—a fact that was not obvious at the start. Algebra discovers this fact from \(D = 0\). Elimination discovers it from \(0 = -2\). The chapter is ending at the point where my linear algebra book begins.
One final comment. In actual computing, you will use a code written by professionals. The steps will be the same as above. A multiple of equation 1 is subtracted from each equation below it, to eliminate the first unknown $x$. With one fewer unknown and equation, elimination starts again. (A parallel computer executes many steps at once.) Extra instructions are included to reduce roundoff error. You only see the result! But it is more satisfying to know what the computer is doing.

In the end, solving linear equations is the key step in solving nonlinear equations. The central idea of differential calculus is to linearize near a point.

11.5 Exercises

Read-through questions

Three equations in three unknowns can be written as $Au = d$. The vector $u$ has components $x$, $y$, $z$ and $A$ is a $b$. The row picture has $c$ for each equation. The first two planes intersect in a $d$, and all three planes intersect in a $e$, which is $f$. The column picture starts with vectors $a$, $b$, $c$ from the columns of $g$ and combines them to produce $h$. The vector equation is $i = d$.

The determinant of $A$ is the triple product $j$. This is the volume of a box, whose edges from the origin are $k$. If $det A = l$, then the system is $m$. Otherwise there is an $n$ matrix such that $A^{-1}A = o$ (the $p$ matrix). In this case the solution to $Au = d$ is $u = q$.

The rows of $A^{-1}$ are the cross products $b \times c$, $r$, $s$, divided by $D$. The entries of $A^{-1}$ are $2$ by $2$ $t$, divided by $D$. The upper left entry equals $u$. The $2$ by $2$ determinants needed for a row of $A^{-1}$ do not use the corresponding $v$ of $A$.

The solution is $u = A^{-1}d$. Its first component $x$ is a ratio of determinants, $|d| b c| / |w$. Cramer's Rule breaks down when $det A = x$. Then the columns $a$, $b$, $c$ lie in the same $y$. There is no solution to $xa + yb + zc = d$, if $d$ is not on that $z$. In a singular row picture, the intersection of planes 1 and 2 is $A$ to the third plane.

In practice $u$ is computed by $b$. The algorithm starts by subtracting a multiple of row 1 to eliminate $x$ from $c$. If the first two equations are $x - y = 1$ and $3x + z = 7$, this elimination step leaves $d$. Similarly $x$ is eliminated from the third equation, and then $e$ is eliminated. The equations are solved by back $f$. When the system has no solution, we reach an impossible equation like $g$. The example $x - y = 1$, $3x + z = 7$ has no solution if the third equation is $h$.

Rewrite 1-4 as matrix equations $Au = d$ (do not solve).

1 $d = (0, 0, 8)$ is a combination of $a = (1, 2, 0)$ and $b = (2, 3, 2)$ and $c = (2, 5, 2)$.

2 The planes $x + y = 0$, $x + y + z = 1$, and $y + z = 0$ intersect at $u = (x, y, z)$.

3 The point $u = (x, y, z)$ is on the planes $x = y$, $y = z$, $x - z = 1$.

4 A combination of $a = (1, 0, 0)$ and $b = (0, 2, 0)$ and $c = (0, 0, 3)$ equals $d = (5, 2, 0)$.

5 Show that Problem 3 has no solution in two ways: find the determinant of $A$, and combine the equations to produce $1 = 0$.

6 Solve Problem 2 in two ways: by inspiration and Cramer's Rule.

7 Solve Problem 4 in two ways: by inspection and by computing the determinant and inverse of the diagonal matrix

$$
A = \begin{bmatrix}
1 & 0 & 0 \\
0 & 2 & 0 \\
0 & 0 & 3
\end{bmatrix}
$$

8 Solve the three equations of Problem 1 by elimination.

9 The vectors $b$ and $c$ lie in a plane which is perpendicular to the vector $d$. In case the vector $a$ also lies in that plane, it is also perpendicular and $a \cdot e = 0$. The $f$ of the matrix with columns in a plane is $g$.

10 The plane $a_1, x + b_1y + c_1z = d_1$ is perpendicular to its normal vector $N_1 = h$. The plane $a_2, x + b_2y + c_2z = d_2$ is perpendicular to $N_2 = i$. The planes meet in a line that is perpendicular to both vectors, so the line is parallel to their $j$ product. If this line is also parallel to the third plane and perpendicular to $N_3$, the system is $k$. The matrix has no $l$, which happens when $(N_1 \times N_2) \cdot N_3 = 0$.

Problems 11-24 use the matrices $A, B, C$.

$$
A = \begin{bmatrix}
1 & 4 & 0 \\
0 & 2 & 6 \\
0 & 0 & 3
\end{bmatrix},
B = \begin{bmatrix}
0 & 0 & 1 \\
2 & 1 & 0 \\
6 & 4 & 0
\end{bmatrix},
C = \begin{bmatrix}
1 & -1 & -3 \\
-1 & 2 & 0 \\
0 & -1 & 3
\end{bmatrix}
$$
11 Find the determinants $|A|, |B|, |C|$. Since $A$ is triangular, its determinant is the product ________.

12 Compute the cross products of each pair of columns in $B$ (three cross products).

13 Compute the inverses of $A$ and $B$ above. Check that $A^{-1}A = I$ and $B^{-1}B = I$.

14 Solve $Au = \begin{bmatrix} 1 \\ 0 \\ 0 \end{bmatrix}$ and $Bu = \begin{bmatrix} 0 \\ 0 \\ 0 \end{bmatrix}$. With this right side $d$, why is $u$ the first column of the inverse?

15 Suppose all three columns of a matrix add to zero, as in $C$ above. The dot product of each column with $v = (1, 1, 1)$ is ________. All three columns lie in the same _________. The determinant of $C$ must be ________.

16 Find a nonzero solution to $Cu = 0$. Find all solutions to $Cu = 0$.

17 Choose any right side $d$ that is perpendicular to $v = (1, 1, 1)$ and solve $Cu = d$. Then find a second solution.

18 Choose any right side $d$ that is not perpendicular to $v = (1, 1, 1)$. Show by elimination (reach an impossible equation) that $Cu = d$ has no solution.

19 Compute the matrix product $AB$ and then its determinant. How is $\det AB$ related to $\det A$ and $\det B$?

20 Compute the matrix products $BC$ and $CB$. All columns of $CB$ add to ________, and its determinant is ________.

21 Add $A$ and $C$ by adding each entry of $A$ to the corresponding entry of $C$. Check whether the determinant of $A + C$ equals $\det A + \det C$.

22 Compute $2A$ by multiplying each entry of $A$ by 2. The determinant of $2A$ equals ________ times the determinant of $A$.

23 Which four entries of $A$ give the upper left corner entry $p$ of $A^{-1}$, after dividing by $D = \det A$? Which four entries of $A$ give the entry $q$ in row 1, column 2 of $A^{-1}$? Find $p$ and $q$.

24 The 2 by 2 determinants from the first two rows of $B$ are $-1$ (from columns 2, 3) and $-2$ (from columns 1, 3) and ________ (from columns 1, 2). These numbers go into the third ________ of $B^{-1}$, after dividing by ________ and changing the sign of ________.

25 Why does every inverse matrix $A^{-1}$ have an inverse?

26 From the multiplication $AB^{-1}A^{-1} = I$ it follows that the inverse of $AB$ is ________. The separate inverses come in ________ order. If you put on socks and then shoes, the inverse begins by taking off ________.

27 Find the determinants of these four permutation matrices:

$$P = \begin{bmatrix} 0 & 1 & 0 \\ 1 & 0 & 0 \\ 0 & 0 & 1 \end{bmatrix}, \quad Q = \begin{bmatrix} 0 & 0 & 1 \\ 0 & 1 & 0 \\ 1 & 0 & 0 \end{bmatrix}, \quad PQ = \begin{bmatrix} 0 & 1 & 0 \\ 0 & 0 & 1 \\ 1 & 0 & 0 \end{bmatrix}$$

and $QP = ________$. Multiply $u = (x, y, z)$ by each permutation to find $Pu, Qu, PQu,$ and $QPu$.

28 Find all six of the 3 by 3 permutation matrices (including $I$), with a single 1 in each row and column. Which of them are "even" (determinant 1) and which are "odd" (determinant -1)?

29 How many 2 by 2 permutation matrices are there, including $I$? How many 4 by 4?

30 Multiply any matrix $A$ by the permutation matrix $P$ and explain how $PA$ is related to $A$. In the opposite order explain how $AP$ is related to $A$.

31 Eliminate $x$ from the last two equations by subtracting the first equation. Then eliminate $y$ from the new third equation by using the new second equation:

$$x + y + z = 2 \quad \text{(a)} \quad x + y = 1$$
$$x + 3y + 3z = 0 \quad \text{(b)} \quad x + z = 3$$
$$x + 3y + 7z = 2 \quad y + z = 5.$$ 

After elimination solve for $z, y, x$ (back substitution).

32 By elimination and back substitution solve:

$$x + 2y + 2z = 0 \quad x - y = 1 \quad \text{(a)} \quad 2x + 3y + 5z = 0 \quad 2x + 3y + 5z = 0 \quad \text{(b)} \quad x - 3z = 4$$
$$2y + 2z = 8 \quad y - z = 7.$$ 

33 Eliminate $x$ from equation 2 by using equation 1:

$$x + 2y + 2z = 0 \quad 2x + 4y + 5z = 0 \quad \text{(a)} \quad 2x + 2y + 2z = 0 \quad 2y + z = 8.$$ 

Why can't the new second equation eliminate $y$ from the third equation? Is there a solution or is the system singular?

Note: If elimination creates a zero in the "pivot position," try to exchange that pivot equation with an equation below it. Elimination succeeds when there is a full set of pivots.

34 The pivots in Problem 32a are 1, -1, and 4. Circle those as they appear along the diagonal in elimination. Check that the product of the pivots equals the determinant. (This is how determinants are computed.)

35 Find the pivots and determinants in Problem 31.

36 Find the inverse of $A = \begin{bmatrix} 1 & 1 & 0 \\ 0 & 1 & 1 \\ 0 & 0 & 1 \end{bmatrix}$ and also of $B = A^2$. 
11.5 Linear Algebra

37. The symbol $a_{ij}$ stands for the entry in row $i$, column $j$. Find $a_{12}$ and $a_{21}$ in Problem 36. The formula $\sum a_{ij}b_{jk}$ gives the entry in which row and column of the matrix product $AB$?

38. Write down a 3 by 3 singular matrix $S$ in which no two rows are parallel. Find a combination of rows 1 and 2 that is parallel to row 3. Find a combination of columns 1 and 2 that is parallel to column 3. Find a nonzero solution to $Su = 0$.

39. Compute these determinants. The 2 by 2 matrix is invertible if _______. The 3 by 3 matrix (is/is not) invertible.
CHAPTER 9  Polar Coordinates and Complex Numbers
  9.1 Polar Coordinates 348
  9.2 Polar Equations and Graphs 351
  9.3 Slope, Length, and Area for Polar Curves 356
  9.4 Complex Numbers 360

CHAPTER 10  Infinite Series
  10.1 The Geometric Series 368
  10.2 Convergence Tests: Positive Series 374
  10.3 Convergence Tests: All Series 381
  10.4 The Taylor Series for $e^x$, sin $x$, and cos $x$ 385
  10.5 Power Series 391

CHAPTER 11  Vectors and Matrices
  11.1 Vectors and Dot Products 398
  11.2 Planes and Projections 407
  11.3 Cross Products and Determinants 416
  11.4 Matrices and Linear Equations 425
  11.5 Linear Algebra in Three Dimensions 435

CHAPTER 12  Motion along a Curve
  12.1 The Position Vector 446
  12.2 Plane Motion: Projectiles and Cycloids 453
  12.3 Tangent Vector and Normal Vector 459
  12.4 Polar Coordinates and Planetary Motion 464

CHAPTER 13  Partial Derivatives
  13.1 Surfaces and Level Curves 472
  13.2 Partial Derivatives 475
  13.3 Tangent Planes and Linear Approximations 480
  13.4 Directional Derivatives and Gradients 490
  13.5 The Chain Rule 497
  13.6 Maxima, Minima, and Saddle Points 504
  13.7 Constraints and Lagrange Multipliers 514
This chapter is about "vector functions." The vector $2i + 4j + 8k$ is constant. The vector $R(t) = t \hat{i} + t^2 \hat{j} + t^3 \hat{k}$ is moving. It is a function of the parameter $t$, which often represents time. At each time $t$, the position vector $R(t)$ locates the moving body:

$$\text{position vector} = R(t) = x(t) \hat{i} + y(t) \hat{j} + z(t) \hat{k}. \quad (1)$$

Our example has $x = t$, $y = t^2$, $z = t^3$. As $t$ varies, these points trace out a curve in space. The parameter $t$ tells when the body passes each point on the curve. The constant vector $2i + 4j + 8k$ is the position vector $R(2)$ at the instant $t = 2$.

What are the questions to be asked? Every student of calculus knows the first question: Find the derivative. If something moves, the Navy salutes it and we differentiate it. At each instant, the body moving along the curve has a speed and a direction. This information is contained in another vector function—the velocity vector $v(t)$ which is the derivative of $R(t)$:

$$v(t) = \frac{dR}{dt} = \frac{dx}{dt} \hat{i} + \frac{dy}{dt} \hat{j} + \frac{dz}{dt} \hat{k}. \quad (2)$$

Since $\hat{i}$, $\hat{j}$, $\hat{k}$ are fixed vectors, their derivatives are zero. In polar coordinates $\hat{i}$ and $\hat{j}$ are replaced by moving vectors. Then the velocity $v$ has more terms from the product rule (Section 12.4).

Two important cases are uniform motion along a line and around a circle. We study those motions in detail ($v$ = constant on line, $v$ = tangent to circle). This section also finds the speed and distance and acceleration for any motion $R(t)$.

Equation (2) is the computing rule for the velocity $dR/dt$. It is not the definition of $dR/dt$, which goes back to basics and does not depend on coordinates:

$$\frac{dR}{dt} = \lim_{\Delta t \to 0} \frac{\Delta R}{\Delta t} = \lim_{\Delta t \to 0} \frac{R(t + \Delta t) - R(t)}{\Delta t}.$$ 

We repeat: $R$ is a vector so $\Delta R$ is a vector so $dR/dt$ is a vector. All three vectors are in Figure 12.1 ($t$ is not a vector!). This figure reveals the key fact about the geometry: the velocity $v = dR/dt$ is tangent to the curve.
12.1 The Position Vector

The vector \( \Delta R \) goes from one point on the curve to a nearby point. Dividing by \( \Delta t \) changes its length, not its direction. That direction lines up with the tangent to the curve, as the points come closer.

**EXAMPLE 1**

\[ \mathbf{R}(t) = t\mathbf{i} + t^2\mathbf{j} + t^3\mathbf{k} \quad \mathbf{v}(t) = \mathbf{i} + 2t\mathbf{j} + 3t^2\mathbf{k} \]

This curve swings upward as \( t \) increases. When \( t = 0 \) the velocity is \( \mathbf{v} = \mathbf{i} \). The tangent is along the \( x \) axis, since the \( j \) and \( k \) components are zero. When \( t = 1 \) the velocity is \( \mathbf{i} + 2\mathbf{j} + 3\mathbf{k} \), and the curve is climbing.

For the shadow on the \( xy \) plane, drop the \( k \) component. Position on the shadow is \( t\mathbf{i} + t^2\mathbf{j} \). Velocity along the shadow is \( \mathbf{i} + 2t\mathbf{j} \). The shadow is a plane curve.

**EXAMPLE 2** Uniform motion in a straight line: the velocity vector \( \mathbf{v} \) is constant.

The speed and direction don't change. The position vector moves with \( d\mathbf{R}/dt = \mathbf{v} \):

\[ \mathbf{R}(t) = \mathbf{R}_0 + t\mathbf{v} \quad (\mathbf{R}_0 \text{ fixed, } \mathbf{v} \text{ fixed, } t \text{ varying}) \quad (3) \]

That is the *equation of a line* in vector form. Certainly \( d\mathbf{R}/dt = \mathbf{v} \). The starting point \( \mathbf{R}_0 = x_0\mathbf{i} + y_0\mathbf{j} + z_0\mathbf{k} \) is given. The velocity \( \mathbf{v} = v_1\mathbf{i} + v_2\mathbf{j} + v_3\mathbf{k} \) is also given. Separating the \( x, y \) and \( z \) components, equation (3) for a line is

**line with parameter:**

\[ x = x_0 + tv_1, \quad y = y_0 + tv_2, \quad z = z_0 + tv_3. \]

The speed along the line is \( |\mathbf{v}| = \sqrt{v_1^2 + v_2^2 + v_3^2} \). The direction of the line is the unit vector \( \mathbf{v}/|\mathbf{v}| \). We have three equations for \( x, y, z \), and eliminating \( t \) leaves two equations. The parameter \( t \) equals \( (x - x_0)/v_1 \) from equation (4). It also equals \( (y - y_0)/v_2 \) and \( (z - z_0)/v_3 \). So these ratios equal each other, and \( t \) is gone:

**line without parameter:**

\[ \frac{x - x_0}{v_1} = \frac{y - y_0}{v_2} = \frac{z - z_0}{v_3}. \]

**An example** is \( x = y/2 = z/3 \). In this case \( (x_0, y_0, z_0) = (0, 0, 0) \)—the line goes through the origin. Another point on the line is \( (x, y, z) = (2, 4, 6) \). Because \( t \) is gone, we cannot say when we reach that point and how fast we are going. The equations \( x/4 = y/8 = z/12 \) give the same line. Without \( t \) we can't know the velocity \( \mathbf{v} = d\mathbf{R}/dt \).

**EXAMPLE 3** Find an equation for the line through \( P = (0, 2, 1) \) and \( Q = (1, 3, 3) \).

**Solution** We have choices! \( \mathbf{R}_0 \) can go to *any point* on the line. The velocity \( \mathbf{v} \) can be *any multiple* of the vector from \( P \) to \( Q \). The decision on \( \mathbf{R}_0 \) controls where we start, and \( \mathbf{v} \) controls our speed.

The vector from \( P \) to \( Q \) is \( \mathbf{i} + \mathbf{j} + 2\mathbf{k} \). Those numbers 1, 1, 2 come from subtracting 0, 2, 1 from 1, 3, 3. We choose this vector \( \mathbf{i} + \mathbf{j} + 2\mathbf{k} \) as a first \( \mathbf{v} \), and double it for a
second v. We choose the vector \( \mathbf{R}_0 = \mathbf{P} \) as a first start and \( \mathbf{R}_0 = \mathbf{Q} \) as a second start.

Here are two different expressions for the same line—they are \( \mathbf{P} + tv \) and \( \mathbf{Q} + t(2v) \):

\[
\mathbf{R}(t) = (2i + k) + t(i + j + 2k) \quad \mathbf{R}^*(t) = (i + 3j + 3k) + t(2i + 2j + 4k).
\]

The vector \( \mathbf{R}(t) \) gives \( x = t, y = 2 + t, z = 1 + 2t \). The vector \( \mathbf{R}^* \) is at a different point on the same line at the same time: \( x^* = 1 + 2t, y^* = 3 + 2t, z^* = 3 + 4t \).

If I pick \( t = 1 \) in \( \mathbf{R} \) and \( t = 0 \) in \( \mathbf{R}^* \), the point is \( (1, 3, 3) \). We arrive there at different times. You are seeing how parameters work, to tell "where" and also "when." If \( t \) goes from \( -\infty \) to \( +\infty \), all points on one line are also on the other line. The path is the same, but the "twins" are going at different speeds.

**Question 1** When do these twins meet? When does \( \mathbf{R}(t) = \mathbf{R}^*(t) \)?

**Answer** They meet at \( t = -1 \), when \( \mathbf{R} = \mathbf{R}^* = -i + j - k \).

**Question 2** What is an equation for the segment between \( P \) and \( Q \) (not beyond)?

**Answer** In the equation for \( \mathbf{R}(t) \), let \( t \) go from 0 to 1 (not beyond):

\[
x = t \quad y = 2 + t \quad z = 1 + 2t \quad \text{[0} \leq t \leq 1 \text{ for segment]}.
\]

At \( t = 0 \) we start from \( P = (0, 2, 1) \). At \( t = 1 \) we reach \( Q = (1, 3, 3) \).

**Question 3** What is an equation for the line without the parameter \( t \)?

**Answer** Solve equations (6) for \( t \) or use (5): \( x/t = (y - 2)/1 = (z - 1)/2 \).

**Question 4** Which point on the line is closest to the origin?

**Answer** The derivative of \( x^2 + y^2 + z^2 = t^2 + (2 + t)^2 + (1 + 2t)^2 \) is \( 8 + 8t \). This derivative is zero at \( t = -1 \). So the closest point is \( (-1, 1, -1) \).

**Question 5** Where does the line meet the plane \( x + y + z = 11 \)?

**Answer** Equation (6) gives \( x + y + z = 3 + 4t = 11 \). So \( t = 2 \). The meeting point is \( x = t = 2, y = t + 2 = 4, z = 1 + 2t = 5 \).

**Question 6** What line goes through \( (3, 1, 1) \) perpendicular to the plane \( x - y - z = 1 \)?

**Answer** The normal vector to the plane is \( \mathbf{N} = i - j - k \). That is \( v \). The position vector to \( (3, 1, 1) \) is \( \mathbf{R}_0 = 3i + j + k \). Then \( \mathbf{R} = \mathbf{R}_0 + tv \).

**COMPARING LINES AND PLANES**

A line has one parameter or two equations. We give the starting point and velocity: \( (x, y, z) = (x_0, y_0, z_0) + t(v_1, v_2, v_3) \). That tells directly which points are on the line. Or we eliminate \( t \) to find the two equations in (5).

A plane has one equation or two parameters! The equation is \( ax + by + cz = d \). That tells us indirectly which points are on the plane. (Instead of knowing \( x, y, z \), we know the equation they satisfy. Instead of directions \( v \) and \( w \) in the plane, we are told the perpendicular direction \( \mathbf{N} = (a, b, c) \).) With parameters, the line contains \( \mathbf{R}_0 + tv \) and the plane contains \( \mathbf{R}_0 + tv + sw \). A plane looks worse with parameters \( (t \text{ and } s) \), a line looks better.

Questions 5 and 6 connected lines to planes. Here are two more. See Problems 41–44:

**Question 7** When is the line \( \mathbf{R}_0 + tv \) parallel to the plane? When is it perpendicular?

**Answer** The test is \( v \cdot \mathbf{N} = 0 \). The test is \( v \times \mathbf{N} = 0 \).

**EXAMPLE 4** Find the plane containing \( P_0 = (1, 2, 1) \) and the line of points \( (1, 0, 0) + t(2, 0, -1) \). That vector \( v \) will be in the plane.
The Position Vector

Solution The vector \( \mathbf{v} = 2\mathbf{i} - \mathbf{k} \) goes along the line. The vector \( \mathbf{w} = 2\mathbf{j} + \mathbf{k} \) goes from \((1, 0, 0)\) to \((1, 2, 1)\). Their cross product is

\[
\mathbf{N} = \mathbf{v} \times \mathbf{w} = \begin{vmatrix} \mathbf{i} & \mathbf{j} & \mathbf{k} \\ 2 & 0 & -1 \\ 0 & 2 & 1 \end{vmatrix} = 2\mathbf{i} - 2\mathbf{j} + 4\mathbf{k}.
\]

The plane \( 2x - 2y + 4z = 2 \) has this normal \( \mathbf{N} \) and contains the point \((1, 2, 1)\).

SPEED, DIRECTION, DISTANCE, ACCELERATION

We go back to the curve traced out by \( \mathbf{R}(t) \). The derivative \( \mathbf{v}(t) = d\mathbf{R}/dt \) is the velocity vector along that curve. The speed is the magnitude of \( \mathbf{v} \):

\[
\text{speed} = |\mathbf{v}| = \sqrt{(dx/dt)^2 + (dy/dt)^2 + (dz/dt)^2}.
\]  

(7)

The direction of the velocity vector is \( \mathbf{v}/|\mathbf{v}| \). This is a unit vector, since \( \mathbf{v} \) is divided by its length. The unit tangent vector \( \mathbf{v}/|\mathbf{v}| \) is denoted by \( \mathbf{T} \).

The tangent vector is constant for lines. It changes direction for curves.

EXAMPLE 5 (important) Find \( \mathbf{v} \) and \( |\mathbf{v}| \) and \( \mathbf{T} \) for steady motion around a circle:

\[
x = r \cos \omega t, \quad y = r \sin \omega t, \quad z = 0.
\]

Solution The position vector is \( \mathbf{R} = r \cos \omega t \mathbf{i} + r \sin \omega t \mathbf{j} \). The velocity is

\[
\mathbf{v} = d\mathbf{R}/dt = -\omega r \sin \omega t \mathbf{i} + \omega r \cos \omega t \mathbf{j} \quad \text{(tangent, not unit tangent)}
\]

The speed is the radius \( r \) times the angular velocity \( \omega \):

\[
|\mathbf{v}| = \sqrt{(-\omega r \sin \omega t)^2 + (\omega r \cos \omega t)^2} = \omega r.
\]

The unit tangent vector is \( \mathbf{v} \) divided by \( |\mathbf{v}| \):

\[
\mathbf{T} = -\sin \omega t \mathbf{i} + \cos \omega t \mathbf{j} \quad \text{(length 1 since } \sin^2 \omega t + \cos^2 \omega t = 1). \]

Think next about the distance traveled. Distance along a curve is always denoted by \( s \) (called arc length). I don't know why we use \( s \)—certainly not as the initial for speed. In fact speed is distance divided by time. The ratio \( s/t \) gives average speed; \( ds/dt \) is instantaneous speed. We are back to Chapter 1 and Section 8.3, the relation of speed to distance:

\[
\text{speed }|\mathbf{v}| = ds/dt \quad \text{distance } s = \int (ds/dt) \, dt = \int |\mathbf{v}(t)| \, dt.
\]

Notice that \( |\mathbf{v}| \) and \( s \) and \( t \) are scalars. The direction vector is \( \mathbf{T} \):

\[
\mathbf{T} = \frac{\mathbf{v}}{|\mathbf{v}|} = \frac{d\mathbf{R}/dt}{ds/dt} = \frac{d\mathbf{R}}{ds} = \text{unit tangent vector}.
\]  

(8)

In Figure 12.3, the chord length (straight) is \( |\Delta \mathbf{R}| \). The arc length (curved) is \( \Delta s \). As \( \Delta \mathbf{R} \) and \( \Delta s \) approach zero, the ratio \( |\Delta \mathbf{R}|/\Delta s \) approaches \( |\mathbf{T}| = 1 \).

Think finally about the acceleration vector \( \mathbf{a}(t) \). It is the rate of change of velocity (not the rate of change of speed):

\[
\mathbf{a} = \frac{d\mathbf{v}}{dt} = \frac{d^2 \mathbf{R}}{dt^2} = \frac{d^2 x}{dt^2} \mathbf{i} + \frac{d^2 y}{dt^2} \mathbf{j} + \frac{d^2 z}{dt^2} \mathbf{k}.
\]  

(9)
For steady motion along a line, as in \( x = t, y = 2 + t, z = 1 + 2t \), there is no acceleration. The second derivatives are all zero. For steady motion around a circle, there is acceleration. In driving a car, you accelerate with the gas pedal or the brake. You also accelerate by turning the wheel. It is the velocity vector that changes, not the speed.

**EXAMPLE 6** Find the distance \( s(t) \) and acceleration \( a(t) \) for circular motion.

**Solution** The speed in Example 5 is \( ds/dt = \omega r \). After integrating, the distance is \( s = \omega rt \). At time \( t \) we have gone through an angle of \( \omega t \). The radius is \( r \), so the distance traveled agrees with \( \omega t \) times \( r \). Note that the dimension of \( \omega \) is \( 1/\text{time} \). (Angles are dimensionless.) At time \( t = 2\pi/\omega \) we have gone once around the circle—to \( s = 2\pi r \) not back to \( s = 0 \).

The acceleration is \( a = d^2\mathbf{R}/dt^2 \). Remember \( \mathbf{R} = r \cos \omega t \mathbf{i} + r \sin \omega t \mathbf{j} \):

\[
a(t) = -\omega^2 r \cos \omega t \mathbf{i} - \omega^2 r \sin \omega t \mathbf{j}. \quad (10)
\]

That direction is opposite to \( \mathbf{R} \). This is a special motion, with no action on the gas pedal or the brake. All the acceleration is from turning. The magnitude is \( |a| = \omega^2 r \), with the correct dimension of distance/(time)^2.

**EXAMPLE 7** Find \( \mathbf{v} \) and \( s \) and \( a \) around the helix \( \mathbf{R} = \cos t \mathbf{i} + \sin t \mathbf{j} + t \mathbf{k} \).

**Solution** The velocity is \( \mathbf{v} = -\sin t \mathbf{i} + \cos t \mathbf{j} + \mathbf{k} \). The speed is

\[
\frac{ds}{dt} = |\mathbf{v}| = \sqrt{\sin^2 t + \cos^2 t + 1} = \sqrt{2} \quad \text{(constant)}.
\]

Then distance is \( s = \sqrt{2} t \). At time \( t = \pi \), a half turn is complete. The distance along the shadow is \( \pi \) (a half circle). The distance along the helix is \( \sqrt{2} \pi \), because of its 45° slope.

The unit tangent vector is velocity/speed, and the acceleration is \( d\mathbf{v}/dt \):

\[
\mathbf{T} = (-\sin t \mathbf{i} + \cos t \mathbf{j} + \mathbf{k})/\sqrt{2} \quad \mathbf{a} = -\cos t \mathbf{i} - \sin t \mathbf{j}.
\]

**EXAMPLE 8** Find \( \mathbf{v} \) and \( s \) and \( a \) around the ellipse \( x = \cos t, y = 2 \sin t, z = 0 \).

**Solution** Take derivatives: \( \mathbf{v} = -\sin t \mathbf{i} + 2 \cos t \mathbf{j} \) and \( |\mathbf{v}| = \sqrt{\sin^2 t + 4 \cos^2 t} \). This is the speed \( ds/dt \). For the distance \( s \), something bad happens (or something normal). The speed is not simplified by \( \sin^2 t + \cos^2 t = 1 \). We cannot integrate \( ds/dt \) to find a formula for \( s \). The square root defeats us.

The acceleration \( -\cos t \mathbf{i} - 2 \sin t \mathbf{j} \) still points to the center. This is not the Earth.
12.1 The Position Vector

going around the sun. The path is an ellipse but the speed is wrong. See Section 12.4 (the pound note) for a terrible error in the position of the sun.

12A The basic formulas for motion along a curve are

\[
\begin{align*}
\mathbf{v} &= \frac{d\mathbf{R}}{dt} \\
\mathbf{a} &= \frac{d\mathbf{v}}{dt} \\
|\mathbf{v}| &= \frac{ds}{dt} \\
\mathbf{T} &= \frac{\mathbf{v}}{|\mathbf{v}|} = \frac{d\mathbf{R}/dt}{ds/dt} = \frac{d\mathbf{R}}{ds}.
\end{align*}
\]

Suppose we know the acceleration \( \mathbf{a}(t) \) and the initial velocity \( \mathbf{v}_0 \) and position \( \mathbf{R}_0 \). Then \( \mathbf{v}(t) \) and \( \mathbf{R}(t) \) are also known. We integrate each component:

\[
\begin{align*}
\mathbf{a}(t) &= \text{constant} \Rightarrow \mathbf{v}(t) &= \mathbf{v}_0 + \mathbf{a}t \\
\mathbf{a}(t) &= \cos \, t \, \mathbf{k} \Rightarrow \mathbf{v}(t) &= \mathbf{v}_0 + \sin \, t \, \mathbf{k} \\
\mathbf{a}(t) &= \cos \, t \, \mathbf{k} \Rightarrow \mathbf{R}(t) &= \mathbf{R}_0 + \mathbf{v}_0 t - \cos \, t \, \mathbf{k}.
\end{align*}
\]

THE CURVE OF A BASEBALL

There is a nice discussion of curve balls in the calculus book by Edwards and Penney. We summarize it here (optionally). The ball leaves the pitcher's hand five feet off the ground: \( \mathbf{R}_0 = 0\mathbf{i} + 0\mathbf{j} + 5\mathbf{k} \). The initial velocity is \( \mathbf{v}_0 = 120\mathbf{i} - 2\mathbf{j} + 2\mathbf{k} \) (120 ft/sec is more than 80 miles per hour). The acceleration is \( -32\mathbf{k} \) from gravity, plus a new term from spin. If the spin is around the \( z \) axis, and the ball goes along the \( x \) axis, then this acceleration is in the \( y \) direction. (It comes from the cross product \( \mathbf{k} \times \mathbf{i} \)—there is a pressure difference on the sides of the ball.) A good pitcher can achieve \( \mathbf{a} = 16\mathbf{j} - 32\mathbf{k} \).

The batter integrates as fast as he can:

\[
\begin{align*}
\mathbf{v}(t) &= \mathbf{v}_0 + \mathbf{a}t = 120\mathbf{i} + (-2 + 16t)\mathbf{j} + (2 - 32t)\mathbf{k} \\
\mathbf{R}(t) &= \mathbf{R}_0 + \mathbf{v}_0 t + \frac{1}{2}\mathbf{a}t^2 = 120t \mathbf{i} + (-2t + 8t^2)\mathbf{j} + (5 + 2t - 16t^2)\mathbf{k}.
\end{align*}
\]

Notice the \( t^2 \). The effect of spin is small at first, then suddenly bigger (as every batter knows). So is the effect of gravity—the ball starts to dive. At \( t = \frac{1}{2} \), the \( \mathbf{i} \) component is 60 feet and the ball reaches the batter. The \( \mathbf{j} \) component is 1 foot and the \( \mathbf{k} \) component is 2 feet—the curve goes low over the outside corner.

At \( t = \frac{1}{4} \), when the batter saw the ball halfway, the \( \mathbf{j} \) component was zero. It looked as if it was coming right over the plate.

Fig. 12.4 A curve ball approaches home plate. Halfway it is on line.

12.1 EXERCISES

**Read-through questions**

The position vector \( \mathbf{a} \) along the curve changes with the parameter \( t \). The velocity is \( \mathbf{b} \). The acceleration is \( \mathbf{c} \). If the position is \( \mathbf{i} + \mathbf{j} + t^2\mathbf{k} \), then \( \mathbf{v} = \mathbf{d} \) and \( \mathbf{a} = \mathbf{e} \). In that example the speed is \( |\mathbf{v}| = \mathbf{f} \). This equals \( ds/dt \), where \( s \) measures the \( \mathbf{g} \). Then \( s = \int \mathbf{h} \). The tangent vector is in the same direction as the \( \mathbf{i} \) vector. In general \( \mathbf{T} = \mathbf{k} \) and in the example \( \mathbf{T} = \mathbf{j} \).

Steady motion along a line has \( \mathbf{a} = \mathbf{m} \). If the line is \( x = y = z \), the unit tangent vector is \( \mathbf{T} = \mathbf{n} \). If the speed is
12 Motion Along a Curve

|t| = \sqrt{3}, the velocity vector is \( \mathbf{v} = -q \). If the initial position is \((1, 0, 0)\), the position vector is \( \mathbf{R}(t) = -p \). The general equation of a line is \( x = x_0 + \alpha t \), \( y = y_0 + \beta t \), \( z = t \). In vector notation this is \( \mathbf{R}(t) = \mathbf{p} + t \mathbf{i} + t \mathbf{j} + \mathbf{k} \). Eliminating \( t \) leaves the equations \( x - x_0 = (y - y_0)/\alpha = (z - z_0)/\beta = \mathbf{v}/v \). A line in space needs \( \mathbf{v} \) equations where a plane needs \( \mathbf{v} \). A line has one parameter where a plane has \( \mathbf{v} \). The line from \( \mathbf{R}_0 = (1, 0, 0) \) to \( \mathbf{R}_1 = (2, 2, 2) \) with \( |\mathbf{v}| = 3 \) is \( \mathbf{R}(t) = \mathbf{x} \).

Steady motion around a circle (radius \( r \), angular velocity \( \omega \)) has \( x = r \cos t \), \( y = r \sin t \), \( z = 0 \). The velocity is \( \mathbf{v} = \omega \mathbf{R} \).

The speed is \( |\mathbf{v}| = \omega \). The acceleration is \( \mathbf{a} = -\omega^2 \mathbf{R} \), which has magnitude \( \omega^2 r \) and direction \( -\mathbf{a} \). Combining upward motion \( \mathbf{R} = tk \) with this circular motion produces motion around a \( \mathbf{f} \). Then \( \mathbf{v} = \mathbf{g} \) and \( |\mathbf{v}| = \mathbf{h} \).

1 Sketch the curve with parametric equations \( x = t \), \( y = t^2 \). Find the velocity vector and the speed at \( t = 1 \).
2 Sketch the path with parametric equations \( x = 1 + t \), \( y = 1 - t \). Find the xy equation of the path and the speed along it.
3 On the circle \( x = \cos t \), \( y = \sin t \) explain by the chain rule and then by geometry why \( dy/dx = -
\frac{\sin t}{\cos t} \).
4 Locate the highest point on the curve \( x = 6t \), \( y = 6t^2 - t^3 \). This curve is a _______. What is the acceleration \( \mathbf{a} \)?
5 Find the velocity vector and the \( xy \) equation of the tangent line to \( x = e^t \), \( y = e^{-t} \) at \( t = 0 \). What is the \( xy \) equation of the curve?
6 Describe the shapes of these curves: (a) \( x = 2t \), \( y = 4t^3 \); (b) \( x = 4t^3 \), \( y = 8t \); (c) \( x = 4t \), \( y = 4t^2 \).

Note: To find “parametric equations” is to find \( x(t) \), \( y(t) \), and possibly \( z(t) \).

7 Find parametric equations for the line through \( \mathbf{P} = (1, 2, 4) \) and \( \mathbf{Q} = (3, 4, 9) \). Probably your speed is 5; change the equations so the speed is 10. Probably your \( \mathbf{R}_3 \) is \( \mathbf{P} \); change the start to \( \mathbf{Q} \).
8 Find an equation for any one plane that is perpendicular to the line in Problem 7. Also find equations for any one line that is perpendicular.
9 On a straight line from \( (2, 3, 4) \) with velocity \( \mathbf{v} = i - k \), the position vector is \( \mathbf{R}(t) = \mathbf{p} \). If the velocity vector is changed to \( \mathbf{v} = i - 3k \), then \( \mathbf{R}(t) = \mathbf{p} \). The path is still _______.
10 Find parametric equations for steady motion from \( \mathbf{P} = (3, 1, -2) \) at \( t = 0 \) on a line to \( \mathbf{Q} = (0, 0, 0) \) at \( t = 3 \). What is the speed? Change parameters so the speed is \( e^t \).
11 The equations \( x = t + y - 2 = 3z - 2 \) describe a _______. The same path is given parametrically by \( x = t + y \), \( y = \ldots \), \( z = \ldots \). The same path is also given by \( x = 1 + 2t \), \( y = \ldots \), \( z = \ldots \).
12 Find parametric equations to go around the unit circle with speed \( e^t \) starting from \( x = 1 \), \( y = 0 \). When is the circle completed?
13 The path \( x = 2y - 3z = 6t \) is a _______ traveled with speed _______. If \( t \) is restricted by \( t \geq 1 \) the path starts at _______. If \( t \) is restricted by \( 0 \leq t \leq 1 \) the path is a _______.
14 Find the closest point to the origin on the line \( x = 1 + t \), \( y = 2 - t \). When and where does it cross the 45° line through the origin? Find the equation of a line it never crosses.
15 (a) How far apart are the two parallel lines \( x = y \) and \( x = y + 1 ? \) (b) How far is the point \( x = 1 \), \( y = t \) from the point \( x = 1 \), \( y = t + 1 ? \) (c) What is the closest distance if their speeds are different: \( x = 1 \), \( y = t \) and \( x = 2t \), \( y = 2t + 1 \)?
16 Which vectors follow the same path as \( \mathbf{R} = t \mathbf{i} + t^2 \mathbf{j} ? \) The speed along the path may be different.
17 Find a parametric form for the straight line \( y = mx + b \).
18 The line \( x = 1 + t \), \( y = 2 + 2t \) passes through the origin provided \( v_1 + v_2 = 0 \). This line crosses the 45° line \( y = x \) unless \( v_1 + v_2 = 0 \).
19 Find the velocity \( \mathbf{v} \) and speed \( |\mathbf{v}| \) and tangent vector \( \mathbf{T} \) for these motions: (a) \( \mathbf{R} = [t + t^{-1}] \mathbf{i} \) (b) \( \mathbf{R} = t \cos t \mathbf{i} + t \sin t \mathbf{j} \) (c) \( \mathbf{R} = (t + 1) \mathbf{i} + (2t + 1) \mathbf{j} + (2t + 2) \mathbf{k} \).
20 If the velocity \( dx/dt \mathbf{i} + dy/dt \mathbf{j} \) is always perpendicular to the position vector \( \mathbf{x} = x \mathbf{i} + y \mathbf{j} \), show from their dot product that \( x^2 + y^2 \) is constant. The point stays on a circle.
21 Find two paths \( \mathbf{R}(t) \) with the same \( \mathbf{v} = \cos t \mathbf{i} + \sin t \mathbf{j} \). Find a third path with a different \( \mathbf{v} \) but the same acceleration.
22 If the acceleration is a constant vector, the path must be _______. If the path is a straight line, the acceleration vector must be _______.
23 Find the minimum and maximum speed if \( x = t \cos t \), \( y = t - \sin t \). Show that \( |\mathbf{a}| \) is constant but not \( \mathbf{a} \). The point is going around a circle while the center is moving on what line?
24 Find \( x(t) \), \( y(t) \) so that the point goes around the circle \( (x - 1)^2 + (y - 3)^2 = 4 \) with speed 1.
25 A ball that is circling with \( x = \cos 2t \), \( y = \sin 2t \) flies off on a tangent at \( t = \pi/8 \). Find its departure point and its position at a later time \( t \) (linear motion; compute its constant velocity \( \mathbf{v} \)).
26 Why is \( |\mathbf{a}| \) generally different from \( ds/dt \)? Give an example of the difference, and an example where they are equal.
27 Change \( t \) so that the speed along the helix \( \mathbf{R} = \cos t \mathbf{i} + \sin t \mathbf{j} + t \mathbf{k} \) is 1 instead of \( \sqrt{2} \). Call the new parameter \( s \).
28 Find the speed \( ds/dt \) on the line \( x = 1 + 6t \), \( y = 2 + 3t \), \( z = 2t \). Integrate to find the length \( s \) from \( (1, 2, 0) \) to \( (13, 8, 4) \). Check by using \( 12^2 + 6^2 + 4^2 \).
29 Find $v$ and $|v|$ and $a$ for the curve $x = \tan t$, $y = \sec t$. What is this curve? At what time does it go to infinity, and along what line?

30 Construct parametric equations for travel on a helix with speed $r$.

31 Suppose the unit tangent vector $T(t)$ is the derivative of $R(t)$. What does that say about the speed? Give a noncircular example.

32 For travel on the path $y = f(x)$, with no parameter, it is impossible to find the _______ but still possible to find the _______ at each point of the path.

Find $x(t)$ and $y(t)$ for paths 33–36.

33 Around the square bounded by $x = 0$, $x = 1$, $y = 0$, $y = 1$, with speed 2. The formulas have four parts.

34 Around the unit circle with speed $e^{-t}$. Do you get all the way around?

35 Around a circle of radius 4 with acceleration $|a| = 1$.

36 Up and down the $y$ axis with constant acceleration $-j$, returning to $(0, 0)$ at $t = 10$.

37 True (with reason) or false (with example):
   (a) If $|R| = 1$ for all $t$ then $|v| =$ constant.
   (b) If $a = 0$ then $R =$ constant.
   (c) If $v \cdot v =$ constant then $v \cdot a = 0$.
   (d) If $v \cdot R = 0$ then $R \cdot R =$ constant.
   (e) There is no path with $v = a$.

38 Find the position vector to the shadow of $t i + t^2 j + t^3 k$ on the $xz$ plane. Is the curve ever parallel to the line $x = y = z$?

39 On the ellipse $x = a \cos t$, $y = b \sin t$, the angle $\theta$ from the center is not the same as $t$ because _______.

40 Two particles are racing from $(1, 0)$ to $(0, 1)$. One follows $x = \cos t$, $y = \sin t$, the other follows $x = 1 + v_1 t$, $y = v_2 t$. Choose $v_1$ and $v_2$ so that the second particle goes slower but wins.

41 Two lines in space are given by $R(t) = P + tv$ and $R(t) = Q + tw$. Four possibilities: The lines are parallel or the same or intersecting or skew. Decide which is which based on the vectors $v$ and $w$ and $u = Q - P$ (which goes between the lines):
   (a) The lines are parallel if _______ are parallel.
   (b) The lines are the same if _______ are parallel.
   (c) The lines intersect if _______ are not parallel but _______ lie in the same plane.
   (d) The lines are skew if the triple product $u \cdot (v \times w)$ is _______.

42 If the lines are skew (not in the same plane), find a formula based on $u$, $v$, $w$ for the distance between them. The vector $u$ may not be perpendicular to the two lines, so project it onto a vector that is.

43 The distance from $Q$ to the line $P + tv$ is the projection of $u = Q - P$ perpendicular to $v$. How far is $Q = (9, 4, 5)$ from the line $x = 1 + t$, $y = 1 + 2t$, $z = 3 + 2t$?

44 Solve Problem 43 by calculus: substitute for $x$, $y$, $z$ in $(x - 9)^2 + (y - 4)^2 + (z - 5)^2$ and minimize. Which $(x, y, z)$ on the line is closest to $(9, 4, 5)$?

45 Practice with parameters, starting from $x = F(t)$, $y = G(t)$.
   (a) The mirror image across the $45^\circ$ line is $x =$ _______, $y =$ _______.
   (b) Write the curve $x = t^2$, $y = t^3$ as $y = f(x)$.
   (c) Why can't $x = t^2$, $y = t^3$ be written as $y = f(x)$?
   (d) If $F$ is invertible then $t = F^{-1}(x)$ and $y =$ _______.(x).

46 From 12:00 to 1:00 a snail crawls steadily out the minute hand (one meter in one hour). Find its position at time $t$ starting from $(0, 0)$.
The usual initial position is \((0,0)\). Some flights start higher, at \((0,h)\). The initial velocity is \((v_0 \cos \alpha, v_0 \sin \alpha)\), where \(v_0\) is the speed and \(\alpha\) is the angle with the horizontal. The acceleration from gravity is purely vertical: \(\frac{d^2y}{dt^2} = -g\). So the horizontal velocity stays at its initial value. The upward velocity decreases by \(-gt\):

\[
dx/dt = v_0 \cos \alpha, \quad dy/dt = v_0 \sin \alpha - gt.
\]

The horizontal distance \(x(t)\) is steadily increasing. The height \(y(t)\) increases and then decreases. To find the position, integrate the velocities (for a high start add \(h\) to \(y\)):

**The projectile path** is \(x(t) = (v_0 \cos \alpha)t, \quad y(t) = (v_0 \sin \alpha)t - \frac{1}{2}gt^2\). \(1\)

This path is a parabola. But it is not written as \(y = ax^2 + bx + c\). It could be, if we eliminated \(t\). Then we would lose track of time. The parabola is \(y(x)\), with no parameter, where we have \(x(t)\) and \(y(t)\).

Basic question: **Where does the projectile hit the ground?** For the parabola, we solve \(y(x) = 0\). That gives the position \(x\). For the projectile we solve \(y(t) = 0\). That gives the time it hits the ground, not the place. If that time is \(T\), then \(x(T)\) gives the place.

The information is there. It takes two steps instead of one, but we learn more.

**EXAMPLE 1** Water leaves the hose at 10 meters/second (this is \(v_0\)). It starts up at the angle \(\alpha\). Find the time \(T\) when \(y\) is zero again, and find where the projectile lands.

**Solution** The flight ends when \(y = (10 \sin \alpha)T - \frac{1}{2}g T^2 = 0\). The flight time is \(T = \frac{(20 \sin \alpha)}{g}\). At that time, the horizontal distance is

\[
x(T) = (10 \cos \alpha)T = \frac{(200 \cos \alpha \sin \alpha)}{g}.
\]

This is the range \(R\).

The projectile (or water from the hose) hits the ground at \(x = R\). To simplify, replace \(200 \cos \alpha \sin \alpha\) by \(100 \sin 2\alpha\). Since \(g = 9.8\) meters/sec\(^2\), we can't reach the car.

The range \(R = \frac{(100 \sin 2\alpha)}{9.8}\) is at most \(100/9.8\). This is less than 12.

The range is greatest when \(\sin 2\alpha = 1\) (\(\alpha\) is \(45^\circ\)). To reach 12 meters we could stand on a ladder (Problem 14). To hit a baseball against air resistance, the best angle is nearer to \(35^\circ\). Figure 12.5 shows symmetric parabolas (no air resistance) and unsymmetric flight paths that drop more steeply.

**12B** The flight time \(T\) and the horizontal range \(R = x(T)\) are reached when \(y = 0\), which means \((v_0 \sin \alpha)T = \frac{1}{2}g T^2\):

\[
T = \frac{(2v_0 \sin \alpha)}{g} \quad \text{and} \quad R = (v_0 \cos \alpha)T = \frac{(v_0^2 \sin 2\alpha)}{g}.
\]

![Fig. 12.5 Equal range \(R\), different times \(T\). Baseballs hit at 35\(^\circ\) with increasing \(v_0\). The dots are at half-seconds (from *The Physics of Baseball* by Robert Adair. Harper and Row 1990).](image-url)
EXAMPLE 2 What are the correct angles $\alpha$ for a given range $R$ and given $v_0$?

Solution The range is $R = (v_0^2 \sin 2\alpha)/g$. This determines the sine of $2\alpha$—but two angles can have the same sine. We might find $2\alpha = 60^\circ$ or $120^\circ$. The starting angles $\alpha = 30^\circ$ and $\alpha = 60^\circ$ in Figure 12.5 give the same $\sin 2\alpha$ and the same range $R$. The flight times contain $\sin \alpha$ and are different.

By calculus, the maximum height occurs when $dy/dt = 0$. Then $v_0 \sin \alpha = gt$, which means that $t = (v_0 \sin \alpha)/g$. This is half of the total flight time $T$—the time going up equals the time coming down. The value of $y$ at this halfway time $t = \frac{1}{2}T$ is

$$y_{\text{max}} = (v_0 \sin \alpha)(v_0 \sin \alpha)/g = (v_0 \sin \alpha)^2/2g. \quad (2)$$

EXAMPLE 3 If a ski jumper goes 90 meters down a $30^\circ$ slope, after taking off at 28 meters/second, find equations for the flight time and the ramp angle $\alpha$.

Solution The jumper lands at the point $x = 90 \cos 30^\circ$, $y = -90 \sin 30^\circ$ (minus sign for obvious reasons). The basic equation (2) is $x = (28 \cos \alpha)t$, $y = (28 \sin \alpha)t - \frac{1}{2}gt^2$. Those are two equations for $\alpha$ and $t$. Note that $t$ is not $T$, the flight time to $y = 0$.

Conclusion The position of a projectile involves three parameters $v_0$, $\alpha$, and $t$. Three pieces of information determine the flight (almost). The reason for the word almost is the presence of $\sin \alpha$ and $\cos \alpha$. Some flight requirements cannot be met (reaching a car at 12 meters). Other requirements can be met in two ways (when the car is close). The equation $\sin \alpha = c$ is more likely to have no solution or two solutions than exactly one solution.

Watch for the three pieces of information in each problem. When a football starts at $v_0 = 20$ meters/second and hits the ground at $x = 40$ meters, the third fact is _______. This is like a lawyer who is asked the fee and says $\$1000 for three questions. "Isn't that steep?" says the client. "Yes," says the lawyer, "now what's your last question?"

CYCLOIDS

A projectile's path is a parabola. To compute it, eliminate $t$ from the equations for $x$ and $y$. Problem 5 finds $y = ax^2 + bx$, a parabola through the origin. The path of a point on a wheel seems equally simple, but eliminating $t$ is virtually impossible. The cycloid is a curve that really needs and uses a parameter.

To trace out a cycloid, roll a circle of radius $a$ along the $x$ axis. Watch the point that starts at the bottom of the circle. It comes back to the bottom at $x = 2\pi a$, after a complete turn of the circle. The path in between is shown in Figure 12.6. After a century of looking for the $xy$ equation, a series of great scientists (Galileo, Christopher Wren, Huygens, Bernoulli, even Newton and l'Hôpital) found the right way to study a cycloid—by introducing a parameter. We will call it $\theta$; it could also be $t$.

Fig. 12.6 Path of $P$ on a rolling circle is a cycloid. Fastest slide to $Q$. 
The parameter is the angle \( \theta \) through which the circle turns. (This angle is not at the origin, like \( \theta \) in polar coordinates.) The circle rolls a distance \( a\theta \), radius times times angle, along the x axis. So the center of the circle is at \( x = a\theta \), \( y = a \). To account for the segment \( CP \), subtract \( a \sin \theta \) from \( x \) and \( a \cos \theta \) from \( y \):

**The point \( P \) has** \( x = a(\theta - \sin \theta) \) and \( y = a(1 - \cos \theta) \).  

(3)

At \( \theta = 0 \) the position is \((0, 0)\). At \( \theta = 2\pi \) the position is \((2\pi a, 0)\). In between, the slope of the cycloid comes from the chain rule:

\[
\frac{dy}{dx} = \frac{dy/d\theta}{dx/d\theta} = \frac{a \sin \theta}{a(1 - \cos \theta)}.
\]

(4)

This is infinite at \( \theta = 0 \). The point on the circle starts straight upward and the cycloid has a cusp. Note how all calculations use the parameter \( \theta \). We go quickly:

**Question 1** Find the area under one arch of the cycloid (\( \theta = 0 \) to \( \theta = 2\pi \)).

**Answer** The area is \( \int y \, dx = \int_{0}^{2\pi} a(1 - \cos \theta)(1 - \cos \theta) \, d\theta \). This equals \( 3\pi a^2 \).

**Question 2** Find the length of the arch, using \( ds = \sqrt{(dx/d\theta)^2 + (dy/d\theta)^2} \, d\theta \).

**Answer** \( \int ds = \int_{0}^{2\pi} a \sqrt{(1 - \cos \theta)^2 + (\sin \theta)^2} \, d\theta = \int_{0}^{2\pi} a \sqrt{2 - 2 \cos \theta} \, d\theta \).

Now substitute \( 1 - \cos \theta = 2 \sin^2 \frac{\theta}{2} \). The square root is \( 2 \sin \frac{\theta}{2} \). The length is \( 8a \).

**Question 3** If the cycloid is turned over (\( y \) is downward), find the time to slide to the bottom. The slider starts with \( v = 0 \) at \( y = 0 \).

**Answer** Kinetic plus potential energy is \( \frac{1}{2}mv^2 - mgy = 0 \) (it starts from zero and can't change). So the speed is \( v = \sqrt{2gy} \). This is \( ds/dt \) and we know \( ds \):

\[
\text{sliding time} = \int ds = \int_{0}^{\pi/2} \sqrt{2gy} = \int_{0}^{\pi/2} \frac{a \sqrt{2 - 2 \cos \theta}}{\sqrt{2a(1 - \cos \theta)}} \, d\theta = \pi \sqrt{a/g}.
\]

Check dimensions: \( a = \text{distance}, \ g = \text{distance}/(\text{time})^2, \ \pi \sqrt{a/g} = \text{time}. \) That is the shortest sliding time for any curve. The cycloid solves the "brachistochrone problem," which minimizes the time down curves from \( O \) to \( Q \) (Figure 12.6). You might think a straight path would be quicker—it is certainly shorter. A straight line has the equation \( x = \pi y/2 \), so the sliding time is

\[
\int dt = \int ds/\sqrt{2gy} = \int_{0}^{\pi/2} \sqrt{(\pi/2)^2 + 1} \, dy/\sqrt{2gy} = \sqrt{\pi^2 + 4} \sqrt{a/g}.
\]

(5)

This is larger than the cycloid time \( \pi \sqrt{a/g} \). It is better to start out vertically and pick up speed early, even if the path is longer.

Instead of publishing his solution, John Bernoulli turned this problem into an international challenge: Prove that the cycloid gives the fastest slide. Most mathematicians couldn't do it. The problem reached Isaac Newton (this was later in his life). As you would expect, Newton solved it. For some reason he sent back his proof with no name. But when Bernoulli received the answer, he was not fooled for a moment: "I recognize the lion by his claws."

What is also amazing is a further property of the cycloid: The time to \( Q \) is the same if you begin anywhere along the path. Starting from rest at \( P \) instead of \( O \), the bottom is reached at the same time. This time Bernoulli got carried away: "You will be petrified with astonishment when I say...".

There are other beautiful curves, closely related to the cycloid. For an epicycloid, the circle rolls around the outside of another circle. For a hypocycloid, the rolling circle is inside the fixed circle. The astroid is the special case with radii in the ratio 1 to 4. It is the curved star in Problem 34, where \( x = a \cos^3 \theta \) and \( y = a \sin^3 \theta \).
12.2 Plane Motion: Projectiles and Cycloids

The cycloid even solves the old puzzle: What point moves backward when a train starts forward? The train wheels have a flange that extends below the track, and \( dx/dt < 0 \) at the bottom of the flange.

---

**Read-through questions**

A projectile starts with speed \( v_0 \) and angle \( \alpha \). At time \( t \) its velocity is \( dx/dt = a \), \( dy/dt = b \), \( v_0 = c \), \( v = d \). Knowing \( v_0 \) and the position of the target, we (can) (cannot) solve for \( \alpha \). Knowing \( \alpha \) and the position of the target, we (can) (cannot) solve for \( v_0 \).

A \( \square \) is traced out by a point on a rolling circle. If the radius is \( a \) and the turning angle is \( \theta \), the center of the circle is at \( x = c \), \( y = d \). The point is at \( x = m \), \( y = n \), starting from \( (0, 0) \). It travels a distance \( o \) in a full turn of the circle. The curve has a \( \square \) at the end of every turn. An upside-down cycloid gives the \( \square \) slide between two points.

**Problems 1–18 and 41 are about projectiles**

1. Find the time of flight \( T \), the range \( R \), and the maximum height \( Y \) of a projectile with \( v_0 = 16 \) ft/sec and
   (a) \( \alpha = 30^\circ \)
   (b) \( \alpha = 60^\circ \)
   (c) \( \alpha = 90^\circ \).

2. If \( v_0 = 32 \) ft/sec and the projectile returns to the ground at \( T = 1 \), find the angle \( \alpha \) and the range \( R \).

3. A ball is thrown at \( 60^\circ \) with \( v_0 = 20 \) meters/sec to clear a wall 2 meters high. How far away is the wall?

4. If \( v(0) = 3i + 3j \) find \( v(t), v(1), v(2) \) and \( \mathbf{R}(t), \mathbf{R}(1), \mathbf{R}(2) \).

5. (a) Eliminate \( t \) from \( x = t, y = -\frac{1}{2}t^2 \) to find the \( xy \) equation of the path. At what \( x \) is \( y = 0 \)?
   (b) Do the same for any \( v_0 \) and \( \alpha \).

6. Find the angle \( \alpha \) for a ball kicked at 30 meters/second if it clears 6 meters traveling horizontally.

7. How far out does a stone hit the water \( h \) feet below, starting with velocity \( v_0 \) at angle \( \alpha = 0^\circ \)?

8. How far out does the same stone go, starting at angle \( \alpha \)? Find an equation for the angle that maximizes the range.

9. A ball starting from \( (0, 0) \) passes through \( (5, 2) \) after 2 seconds. Find \( v_0 \) and \( \alpha \). (The units are meters.)

*10 With \( x \) and \( y \) from equation (1), show that
   \[ \frac{v_0^2}{g} = \frac{(gx/v_0)^2}{2gy}. \]

---

**Problems 19–34 are about cycloids and related curves**

19. Find the unit tangent vector \( \mathbf{T} \) to the cycloid. Also find the speed at \( \theta = 0 \) and \( \theta = \pi \), if the wheel turns at \( d\theta/dt = 1 \).

20. The slope of the cycloid is infinite at \( \theta = 0 \):
   \[ \frac{dy}{dx} = \frac{dy/d\theta}{dx/d\theta} = \frac{\sin \theta}{1 - \cos \theta}. \]
   By whose rule? Estimate the slope at \( \theta = \frac{\pi}{2} \) and \( \theta = -\frac{\pi}{2} \). Where does the slope equal one?

21. Show that the tangent to the cycloid at \( P \) in Figure 12.6a goes through \( x = ab, y = 2a \). Where is this point on the rolling circle?

22. For a trochoid, the point \( P \) is a distance \( d \) from the center
12 Motion Along a Curve

of the rolling circle. Redraw Figure 12.6b to find \( x = a\theta - d\sin \theta \) and \( y = \) _______.

23 If a circle of radius \( a \) rolls inside a circle of radius \( 2a \), show that one point on the small circle goes across on a straight line.

24 Find \( d^2y/dx^2 \) for the cycloid, which is concave _______.

25 If \( d\theta/dx = c \), find the velocities \( dx/dt \) and \( dy/dt \) along the cycloid. Where is \( dx/dt \) greatest and where is \( dy/dt \) greatest?

26 Experiment with graphs of \( x = a\cos \theta + b\sin \theta \), \( y = c\cos \theta + d\sin \theta \) using a computer. What kind of curves are they? Why are they closed?

27 A stone in a bicycle tire goes along a cycloid. Find equations for the stone’s path if it flies off at the top (a projectile).

28 Draw curves on a computer with \( x = a\cos \theta + b\cos 3\theta \) and \( y = c\sin \theta + d\sin 3\theta \). Is there a limit to the number of loops?

29 When a penny rolls completely around another penny, the head makes _______ turns. When it rolls inside a circle four times larger (for the astroid), the head makes _______ turns.

30 Display the cycloid family with computer graphics:
   (a) cycloid
   (b) epicycloid \( x = C\cos \theta - \cos C\theta \), \( y = C\sin \theta + \sin C\theta \)
   (c) hypocycloid \( x = e\cos \theta + \cos e\theta \), \( y = e\sin \theta - \sin e\theta \)
   (d) astroid \((c = 3)\)
   (e) deltoid \((c = 2)\).

31 If one arch of the cycloid is revolved around the \( x \) axis, find the surface area and volume.

32 For a hypocycloid the fixed circle has radius \( c + 1 \) and the circle rolling inside has radius 1. There are \( c + 1 \) cusps if \( c \) is an integer. How many cusps (use computer graphics if possible) for \( c = 1/2? \), \( c = 3/2? \), \( c = \sqrt{2} \)? What curve for \( c = 1 \)?

33 When a string is unwound from a circle find \( x(\theta) \) and \( y(\theta) \) for point \( P. \) Its path is the “involute” of the circle.

34 For the point \( P \) on the astroid, explain why \( x = 3\cos \theta + 3\sin 3\theta \) and \( y = 3\sin \theta - \sin 3\theta \). The angle in the figure is \( 3\theta \) because both circular arcs have length _______. Convert to \( x = 4\cos^3 \theta \), \( y = 4\sin^3 \theta \) by triple-angle formulas.

35 Find the area inside the astroid.

36 Explain why \( x = 2a\cos \theta \) and \( y = 2a\sin^2 \theta \) for the point \( P \) on the witch of Agnesi. Eliminate \( \theta \) to find the \( xy \) equation. Note: Maria Agnesi wrote the first three-semester calculus text (l'Hôpital didn’t do integral calculus). The word “witch” is a total mistranslation, nothing to do with her or the curve.

37 For a cardioid the radius \( C - 1 \) of the fixed circle equals the radius 1 of the circle rolling outside (epicycloid with \( C = 2 \)).
   (a) The coordinates of \( P \) are \( x = -1 + 2\cos \theta - \cos 2\theta \), \( y = \) _______.
   (b) The double-angle formulas yield \( x = 2\cos \theta(1 - \cos \theta) \), \( y = \) _______.
   (c) \( x^2 + y^2 = \) _______ so its square root is \( r = \) _______.

38 Explain the last two steps in equation (5) for the sliding time down a straight path.

39 On an upside-down cycloid the slider takes the same time \( T \) to reach bottom wherever it starts. Starting at \( \theta = x \), write \( 1 - \cos \theta = 2\sin^2 \theta/2 \) and \( 1 - \cos \theta = 2\sin^2 \theta/2 \) to show that

\[
T = \int_0^\pi \sqrt{2a^2(1 - \cos \theta)} \, d\theta = \pi \sqrt{a}. \]

40 Suppose a heavy weight is attached to the top of the rolling circle. What is the path of the weight?

41 The wall in Fenway Park is 37 feet high and 315 feet from home plate. A baseball hit 3 feet above the ground at \( \alpha = 22.5^\circ \) will just go over if \( v_0 = \) _______. The time to reach the wall is _______.
A driver produces acceleration three ways—by the gas pedal, the brake, and steering wheel. The first two change the speed. Turning the wheel changes the direction. All three change the velocity (they give acceleration). For steady motion around a circle, the change is from steering—the acceleration $dv/dt$ points to the center. We now look at motion along other curves, to separate change in the speed $|v|$ from change in the direction $T$.

The direction of motion is $T = v/|v|$. It depends on the path but not the speed (because we divide by $|v|$). For turning we measure two things:

1. How fast $T$ turns: this will be the curvature $\kappa$ (kappa).
2. Which direction $T$ turns: this will be the normal vector $N$.

$\kappa$ and $N$ depend, like $s$ and $T$, only on the shape of the curve. Replacing $t$ by $2t$ or $t^2$ leaves them unchanged. For a circle we give the answers in advance. The normal vector $N$ points to the center. The curvature $\kappa$ is $1/radius$.

A smaller turning circle means a larger curvature $\kappa$: more bending.

The curvature $\kappa$ is change in direction $|dT|$ divided by change in position $|ds|$. There are three formulas for $\kappa$—a direct one for graphs $y(x)$, a brutal but valuable one for any parametric curve $(x(t), y(t))$, and a neat formula that uses the vectors $v$ and $a$. We begin with the definition and the neat formula.

**DEFINITION**

$\kappa = |dT|/|ds|

**FORMULA**

$\kappa = |v \times a|/|v|^3 \quad (1)$

The definition does not involve the parameter $t$—but the calculations do. The position vector $R(t)$ yields $v = dR/dt$ and $a = dv/dt$. If $t$ is changed to $2t$, the velocity $v$ is doubled and $a$ is multiplied by 4. Then $|v \times a|$ and $|v|^3$ are multiplied by 8, and their ratio $\kappa$ is unchanged.

**Proof of formula (1)** Start from $v = |v|T$ and compute its derivative $a$:

$$a = \frac{dv}{dt} T + |v| \frac{dT}{dt}$$

by the product rule.

Now take the cross product with $v = |v|T$. Remember that $T \times T = 0$:

$$v \times a = |v|T \times |v| \frac{dT}{dt} = |v\times a| \frac{dT}{dt} \quad (2)$$

We know that $|T| = 1$. Equation (4) will show that $T$ is perpendicular to $dT/dt$. So $|v \times a|$ is the first length $|v|$ times the second length $|v| |dT/dt|$. The factor $\sin \theta$ in the length of a cross product is 1 from the $90^\circ$ angle. In other words

$$\frac{dT}{dt} = \frac{|v \times a|}{|v|^2}$$

and

$$\kappa = \frac{|dT/dt|}{ds/dt} = \frac{|v \times a|}{|v|^3} \quad (3)$$

The chain rule brings the extra $|ds/dt| = |v|$ into the denominator.

Before any examples, we show that $dT/dt$ is perpendicular to $T$. The reason is that $T$ is a unit vector. Differentiate both sides of $T \cdot T = 1$:

$$\frac{dT}{dt} \cdot T + T \cdot \frac{dT}{dt} = 0 \quad \text{or} \quad 2T \cdot \frac{dT}{dt} = 0. \quad (4)$$
12 Motion Along a Curve

That proof used the product rule \( U' \cdot V + U \cdot V' \) for the derivative of \( U \cdot V \) (Problem 23, with \( U = V = T \)). Think of the vector \( T \) moving around the unit sphere. To keep a constant length \((T + dT) \cdot (T + dT) = 1\), we need \( 2T \cdot dT = 0 \). Movement \( dT \) is perpendicular to radius vector \( T \).

Our first examples will be plane curves. The position vector \( R(t) \) has components \( x(t) \) and \( y(t) \) but no \( z(t) \). Look at the components of \( v \) and \( a \) and \( v \times a \) (\( x' \) means \( dx/dt \)):

\[
\begin{align*}
R & \quad x(t) \quad y(t) \quad 0 \\
v & \quad x'(t) \quad y'(t) \quad 0 \quad |v| = \sqrt{|x'|^2 + |y'|^2} \\
a & \quad x''(t) \quad y''(t) \quad 0 \quad \kappa = \frac{|x'y'' - y'x''|}{\sqrt{(x')^2 + (y')^2}} \\
v \times a & \quad 0 \quad 0 \quad 0 \quad x'y'' - y'x''
\end{align*}
\]

Equation (5) is the brutal but valuable formula for \( \kappa \). Apply it to movement around a circle. We should find \( \kappa = 1/\text{radius} \).

**EXAMPLE 1** When \( x = a \cos \omega t \) and \( y = a \sin \omega t \) we substitute \( x', y', x'', y'' \) into (5):

\[
\kappa = \frac{(-\omega a \sin \omega t)(-\omega^2 a \sin \omega t) - (\omega a \cos \omega t)(-\omega^2 a \cos \omega t)}{[(\omega a \sin \omega t)^2 + (\omega a \cos \omega t)^2]^{3/2}} = \frac{\omega^3 a^2}{[\omega^2 a^2]^{3/2}}.
\]

This is \( \omega^3 a^2/\omega^3 a^3 \) and \( \omega \) cancels. The speed makes no difference to \( \kappa = 1/a \).

The third formula for \( \kappa \) applies to an ordinary plane curve given by \( y(x) \). The parameter \( t \) is \( x \). You see the square root in the speed \( |v| = ds/dx \):

\[
\begin{align*}
\begin{array}{lll}
R & \quad x \quad y(x) \quad 0 \\
v & \quad 1 \quad dy/dx \quad 0 \quad |v| = \sqrt{1 + (dy/dx)^2} \\
a & \quad 0 \quad d^2y/dx^2 \quad 0 \quad \kappa = \frac{|d^2y/dx^2|}{(1 + (dy/dx)^2)^{3/2}}
\end{array}
\end{align*}
\]

In practice this is the most popular formula for \( \kappa \). The most popular approximation is \( |d^2y/dx^2| \). (The denominator is omitted.) For the bending of a beam, the nonlinear equation uses \( \kappa \) and the linear equation uses \( d^2y/dx^2 \). We can see the difference for a parabola:

**EXAMPLE 2** The curvature of \( y = \frac{1}{2}x^2 \) is \( \kappa = |y''|/(1 + (y')^2)^{3/2} = 1/(1 + x^2)^{3/2} \).

\[\text{Fig. 12.7} \quad \text{Normal} \ N \ \text{divided by curvature} \ \kappa \ \text{for circle and parabola and unit helix.}\]
The approximation is \( y'' = 1 \). This agrees with \( \kappa \) at \( x = 0 \), where the parabola turns the corner. But for large \( x \), the curvature approaches zero. Far out on the parabola, we go a long way for a small change in direction.

The parabola \( y = -\frac{1}{2}x^2 \), opening down, has the same \( \kappa \). Now try a space curve.

**EXAMPLE 3**  Find the curvature of the unit helix \( \mathbf{R} = \cos t \mathbf{i} + \sin t \mathbf{j} + \mathbf{t} \).

Take the cross product of \( \mathbf{v} = -\sin t \mathbf{i} + \cos t \mathbf{j} + \mathbf{k} \) and \( \mathbf{a} = -\cos t \mathbf{i} - \sin t \mathbf{j} \):

\[
\begin{vmatrix}
\mathbf{i} & \mathbf{j} & \mathbf{k} \\
-\sin t & \cos t & 1 \\
-\cos t & -\sin t & 0
\end{vmatrix}
= \sin t \mathbf{i} - \cos t \mathbf{j} + \mathbf{k}.
\]

This cross product has length \( \sqrt{2} \). Also the speed is \( |\mathbf{v}| = \sqrt{\sin^2 t + \cos^2 t + 1} = \sqrt{2} \):

\[
\kappa = \frac{|\mathbf{v} \times \mathbf{a}|}{|\mathbf{v}|^3} = \sqrt{2}/(\sqrt{2})^3 = \frac{1}{2}.
\]

Compare with a unit circle. Without the climbing term \( \mathbf{t} \), the curvature would be 1. Because of climbing, each turn of the helix is longer and \( \kappa = \frac{1}{2} \).

That makes one think: Is the helix twice as long as the circle? No. The length of a turn is only increased by \( |\mathbf{v}| = \sqrt{2} \). The other \( \sqrt{2} \) is because the tangent \( \mathbf{T} \) slopes upward. The shadow in the base turns a full 360°, but \( \mathbf{T} \) turns less.

**THE NORMAL VECTOR \( \mathbf{N} \)**

The discussion is bringing us to an important vector. Where \( \kappa \) measures the rate of turning, the unit vector \( \mathbf{N} \) gives the direction of turning. \( \mathbf{N} \) is perpendicular to \( \mathbf{T} \), and in the plane that leaves practically no choice. Turn left or right. For a space curve, follow \( d\mathbf{T}/dt \). Remember equation (4), which makes \( d\mathbf{T} \) perpendicular to \( \mathbf{T} \).

*The normal vector \( \mathbf{N} \) is a unit vector along \( d\mathbf{T}/dt \). It is perpendicular to \( \mathbf{T} \):*

\[
\text{DEFINITION} \quad \mathbf{N} = \frac{d\mathbf{T}/ds}{|d\mathbf{T}/ds|} = \frac{1}{\kappa} \frac{d\mathbf{T}}{ds} \quad \text{FORMULA} \quad \mathbf{N} = \frac{d\mathbf{V}/dt}{|d\mathbf{V}/dt|}.
\]

**EXAMPLE 4**  Find the normal vector \( \mathbf{N} \) for the same helix \( \mathbf{R} = \cos t \mathbf{i} + \sin t \mathbf{j} + \mathbf{t} \).

Solution  Copy \( \mathbf{v} \) from Example 3, divide by \( |\mathbf{v}| \), and compute \( d\mathbf{T}/dt \):

\[
\mathbf{T} = \mathbf{v}/|\mathbf{v}| = (-\sin t \mathbf{i} + \cos t \mathbf{j} + \mathbf{k})/\sqrt{2} \quad \text{and} \quad \frac{d\mathbf{T}}{dt} = (-\cos t \mathbf{i} - \sin t \mathbf{j})/\sqrt{2}.
\]

To change \( d\mathbf{T}/dt \) into a unit vector, cancel the \( \sqrt{2} \). The normal vector is \( \mathbf{N} = -\cos t \mathbf{i} - \sin t \mathbf{j} \). It is perpendicular to \( \mathbf{T} \). Since the \( \mathbf{k} \) component is zero, \( \mathbf{N} \) is horizontal. The tangent \( \mathbf{T} \) slopes up at 45°—it goes around the circle at that latitude. The normal \( \mathbf{N} \) is tangent to this circle (\( \mathbf{N} \) is tangent to the path of the tangent!). So \( \mathbf{N} \) stays horizontal as the helix climbs.

There is also a third direction, perpendicular to \( \mathbf{T} \) and \( \mathbf{N} \). It is the binormal vector \( \mathbf{B} = \mathbf{T} \times \mathbf{N} \), computed in Problems 25–30. The unit vectors \( \mathbf{T} \), \( \mathbf{N} \), \( \mathbf{B} \) provide the natural coordinate system for the path—along the curve, in the plane of the curve, and out of that plane. The theory is beautiful but the computations are not often done—we stop here.
12 Motion Along a Curve

TANGENTIAL AND NORMAL COMPONENTS OF ACCELERATION

May I return a last time to the gas pedal and brake and steering wheel? The first two give acceleration along \( T \). Turning gives acceleration along \( N \). The rate of turning (curvature \( \kappa \)) and the direction \( N \) are established. We now ask about the force required. Newton's Law is \( F = ma \), so we need the acceleration \( a \)—especially its component along \( T \) and its component along \( N \).

\[
\text{The acceleration is } a = \frac{d^2 s}{dt^2} T + \kappa \left[ \frac{ds}{dt} \right]^2 N. \tag{8}
\]

For a straight path, \( d^2 s/dt^2 \) is the only acceleration—the ordinary second derivative. The term \( \kappa (ds/dt)^2 \) is the acceleration in turning. Both have the dimension of length/(time)^2.

The force to steer around a corner depends on curvature and speed—as all drivers know. Acceleration is the derivative of \( v = |v|T = (ds/dt)T \):

\[
a = \frac{d^2 s}{dt^2} T + \frac{ds}{dt} \frac{dT}{dt} = \frac{d^2 s}{dt^2} T + \frac{ds}{dt} \frac{dT}{ds} ds. \tag{9}
\]

That last term is \( \kappa (ds/dt)^2 N \), since \( dT/ds = \kappa N \) by formula (7). So (8) is proved.

**EXAMPLE 5** A fixed speed \( ds/dt = 1 \) gives \( d^2 s/dt^2 = 0 \). The only acceleration is \( \kappa N \).

**EXAMPLE 6** Find the components of \( a \) for circular speed-up \( R(t) = \cos t^2 \hat{i} + \sin t^2 \hat{j} \).

Without stopping to think, compute \( dR/dt = v \) and \( ds/dt = |v| \) and \( v/|v| = T \):

\[
v = -2t \sin t^2 \hat{i} + 2t \cos t^2 \hat{j}, \quad |v| = 2t, \quad T = -\sin t^2 \hat{i} + \cos t^2 \hat{j}.
\]

The derivative of \( ds/dt = |v| \) is \( d^2 s/dt^2 = 2 \). The derivative of \( v \) is \( a \):

\[
a = -2 \sin t^2 \hat{i} + 2 \cos t^2 \hat{j} - 4t^2 \cos t^2 \hat{i} - 4t^2 \sin t^2 \hat{j}.
\]

In the first terms of \( a \) we see \( 2T \). In the last terms we must be seeing \( \kappa |v|^3 N \). Certainly \( |v|^2 = 4t^2 \) and \( \kappa = 1 \), because the circle has radius 1. Thus \( a = 2T + 4t^2 N \) has the tangential component 2 and normal component \( 4t^2 \)—acceleration along the circle and in to the center.

**Table of Formulas**

<table>
<thead>
<tr>
<th>Formula</th>
<th>Description</th>
</tr>
</thead>
<tbody>
<tr>
<td>( v = dR/dt )</td>
<td>( a = d\mathbf{v}/dt )</td>
</tr>
<tr>
<td>(</td>
<td>v</td>
</tr>
<tr>
<td>Curvature ( \kappa =</td>
<td>dT/ds</td>
</tr>
<tr>
<td>Normal vector ( \mathbf{N} = -\frac{dT}{\kappa ds} = \frac{dT/dt}{</td>
<td>dT/dt</td>
</tr>
</tbody>
</table>

![Fig. 12.8 Components of \( a \) as car turns corner.](image)
The curvature tells how fast the curve \( a \). For a circle of radius \( a \), the direction changes by \( 2\pi \) in a distance \( b \), so \( \kappa = \frac{a}{b} \). For a plane curve \( y = f(x) \) the formula is \( \kappa = \left| y'' \right| / a \). The curvature of \( y = \sin x \) is \( b \). At a point where \( y'' = 0 \) (an \( c \) point) the curve is momentarily straight and \( \kappa = \frac{d}{b} \). For a space curve \( \kappa = r \times a \). The normal vector \( N \) is perpendicular to \( dN/ds \). At unit speed on any curve, the time \( d \) is constant. True or false: The path is a circle.

**Compute the curvature \( \kappa \) in Problems 1–8.**

1. \( y = e^x \)
2. \( y = \ln x \) (where is \( \kappa \) largest?)
3. \( x = 2 \cos t \), \( y = 2 \sin t \)
4. \( x = \cos t^2 \), \( y = \sin t \)
5. \( x = 1 + t^2 \), \( y = 3t^2 \) (the path is a ________).
6. \( x = \cos^3 t \), \( y = \sin^3 t \)
7. \( r = \theta + t \) (so \( x = \cos t \), \( y = \))
8. \( x = t \), \( y = \ln \cos t \)
9. Find \( T \) and \( N \) in Problem 4.
10. Show that \( N = \sin t \) \( i + \cos t \) \( j \) in Problem 6.
11. Compute \( T \) and \( N \) in Problem 8.
12. Find the speed \( \|v\| \) and curvature \( \kappa \) of a projectile:
   \[ x = (v_0 \cos \alpha) t, \ y = (v_0 \sin \alpha) t - \frac{1}{2} gt^2. \]
13. Find \( T \) and \( \|v\| \) and \( \kappa \) for the helix \( R = 3 \cos t \) \( i + 3 \sin t \) \( j + 4t \) \( k \). How much longer is a turn of the helix than the corresponding circle? What is the upward slope of \( T \)?
14. When \( \kappa = 0 \) the path is a _______. This happens when \( v \) and \( a \) are _______. Then \( v \times a = _______ \).
15. Find the curvature of a cycloid \( x = a(t - \sin t), \ y = a(1 - \cos t) \).
16. If all points of a curve are moved twice as far from the origin \( (x \rightarrow 2x, \ y \rightarrow 2y) \), what happens to \( \kappa \)? What happens to \( N \)?

**17** Find \( \kappa \) and \( N \) at \( \theta = \pi \) for the hypocycloid \( x = 4 \cos \theta + \cos 4\theta, \ y = 4 \sin \theta - \sin 4\theta \).

**18** From \( v = ||T|| \) and \( a \) in equation (8), derive \( \kappa = ||v \times a|| ||v||^3 \).

**19** From a point on the curve, go along the vector \( N/\kappa \) to find the center of curvature. Locate this center for the point \( (1, 0) \) on the circle \( x = \cos t \), \( y = \sin t \) and the ellipse \( x = \cos t \), \( y = 2 \sin t \) and the parabola \( y = \frac{1}{2}(x^2 - 1) \). The path of the center of curvature is the “evolute” of the curve.

**20** Which of these depend only on the shape of the curve, and which depend also on the speed? \( v, T, ||v||, \kappa, a, N, B \).

**21** A plane curve through \( (0, 0) \) and \( (2, 0) \) with constant curvature \( \kappa \) is the circular arc _______. For which \( \kappa \) is there no such curve?

**22** Sketch a smooth curve going through \( (0, 0), (1, -1), \) and \( (2, 0) \). Somewhere \( d^2y/dx^2 \) is at least _______. Somewhere the curvature is at least _______. (Proof is for instructors only.)

**23** For plane vectors, the ordinary product rule applied to \( U, V \), \( U = V' + U' \) shows that \( (U \cdot V)' = U' \cdot V + U \cdot V' \).

**24** If \( v \) is perpendicular to \( a \), prove that the speed is constant. True or false: The path is a circle.

**Problems 25–30 work with the T-N-B system—along the curve, in the plane of the curve, perpendicular to that plane.**

**25** Compute \( B = T \times N \) for the helix \( R = \cos t \) \( i + \sin t \) \( j + tk \) in Examples 3–4.

**26** Using Problem 23, differentiate \( B \cdot T = 0 \) and \( B \cdot B = 1 \) to show that \( B \) is perpendicular to \( T \) and \( B \). So \( dB/ds = -\tau N \) for some number \( \tau \) called the torsion.

**27** Compute the torsion \( \tau = |dB/ds| \) for the helix in Problem 25.

**28** Find \( B = T \times N \) for the curve \( x = t, y = t, z = t^2 \).

**29** A circle lies in the xy plane. Its normal \( N \) lies _______ and \( B = _______ \) and \( \tau = |dB/ds| = _______ \).

**30** The Serret-Frenet formulas are \( dT/ds = \kappa N, dN/ds = -\kappa T + \tau B, dB/ds = -\tau N \). We know the first and third. Differentiate \( N = -T \times B \) to find the second.

**31** The angle \( \theta \) from the x axis to the tangent line is \( \theta = \tan^{-1}(dy/dx) \), when \( dy/dx \) is the slope of the curve.
   (a) Compute \( d\theta/dx \).
   (b) Divide by \( ds/dx = (1 - (dy/dx)^2)^{1/2} \) to show that \( |d\theta/ds| \) is \( \kappa \) in equation (5). Curvature is change in direction \( |d\theta| \) divided by change in position \( |ds| \).

**32** If the tangent direction is at angle \( \theta \) then \( T = \cos \theta \) \( i + \sin \theta \) \( j \). In Problem 31 \( |d\theta/ds| \) agreed with \( \kappa = |dT/ds| \) because \( |dT/d\theta| = _______ \).
12 Motion Along a Curve

In 33–37 find the $T$ and $N$ components of acceleration.

33 $x = 5 \cos \omega t$, $y = 5 \sin \omega t$, $z = 0$ (circle)
34 $x = 1 + t$, $y = 1 + 2t$, $z = 1 + 3t$ (line)
35 $x = t \cos t$, $y = t \sin t$, $z = 0$ (spiral)
36 $x = e^t \cos t$, $y = e^t \sin t$, $z = 0$ (spiral)
37 $x = 1$, $y = t$, $z = t^2$.
38 For the spiral in 36, show that the angle between $\mathbf{R}$ and $\mathbf{a}$ (position and acceleration) is constant. Find the angle.
39 Find the curvature of a polar curve $r = F(\theta)$.

12.4 Polar Coordinates and Planetary Motion

This section has a general purpose—to do vector calculus in polar coordinates. It also has a specific purpose—to study central forces and the motion of planets. The main gravitational force on a planet is from the sun. It is a central force, because it comes from the sun at the center. Polar coordinates are natural, so the two purposes go together.

You may feel that the planets are too old for this course. But Kepler’s laws are more than theorems, they are something special in the history of mankind—“the greatest scientific discovery of all time.” If we can recapture that glory we should do it. Part of the greatness is in the difficulty—Kepler was working sixty years before Newton discovered calculus. From pages of observations, and some terrific guesses, a theory was born. We will try to preserve the greatness without the difficulty, and show how elliptic orbits come from calculus. The first conclusion is quick.

**Motion in a central force field always stays in a plane.**

$F$ is a multiple of the vector $\mathbf{R}$ from the origin (central force). $F$ also equals $m\mathbf{a}$ (Newton’s Law). Therefore $\mathbf{R}$ and $\mathbf{a}$ are in the same direction and $\mathbf{R} \times \mathbf{a} = 0$. Then $\mathbf{R} \times \mathbf{v}$ has zero derivative and is constant:

$$
\frac{d}{dt} (\mathbf{R} \times \mathbf{v}) = \mathbf{v} \times \mathbf{v} + \mathbf{R} \times \mathbf{a} = \mathbf{0} + \mathbf{0}.
$$

(1)

$\mathbf{R} \times \mathbf{v}$ is a constant vector $\mathbf{H}$. So $\mathbf{R}$ stays in the plane perpendicular to $\mathbf{H}$.

How does a planet move in that plane? We turn to polar coordinates. At each point except the origin (where the sun is), $\mathbf{u}_r$ is *the unit vector pointing outward*. It is the position vector $\mathbf{R}$ divided by its length $r$ (which is $\sqrt{x^2 + y^2}$):

$$
\mathbf{u}_r = \frac{\mathbf{R}}{r} = (xi + yj)/r = \cos \theta \mathbf{i} + \sin \theta \mathbf{j}.
$$

(2)

That is a unit vector because $\cos^2 \theta + \sin^2 \theta = 1$. It goes out from the center. Figure 12.9 shows $\mathbf{u}_r$ and the second unit vector $\mathbf{u}_\theta$ at a 90° angle:

$$
\mathbf{u}_\theta = -\sin \theta \mathbf{i} + \cos \theta \mathbf{j}.
$$

(3)

The dot product is $\mathbf{u}_r \cdot \mathbf{u}_\theta = 0$. The subscripts $r$ and $\theta$ indicate direction (not derivative).

**Question 1:** How do $\mathbf{u}_r$ and $\mathbf{u}_\theta$ change as $r$ changes (out a ray)? They don’t.

**Question 2:** How do $\mathbf{u}_r$ and $\mathbf{u}_\theta$ change as $\theta$ changes? Take the derivative:

$$
\frac{d\mathbf{u}_r}{d\theta} = -\sin \theta \mathbf{i} + \cos \theta \mathbf{j} = \mathbf{u}_\theta,
$$

$$
\frac{d\mathbf{u}_\theta}{d\theta} = -\cos \theta \mathbf{i} - \sin \theta \mathbf{j} = -\mathbf{u}_r.
$$

(4)
12.4 Polar Coordinates and Planetary Motion

Since \( u_r = \frac{\mathbf{R}}{r} \), one formula is simple: The position vector is \( \mathbf{R} = ru_r \). For its derivative \( v = \frac{d\mathbf{R}}{dt} \), use the chain rule \( du_r/dt = (du_r/d\theta)(d\theta/dt) = (d\theta/dt)u_\theta \):

\[
\text{The velocity is } v = \frac{d}{dt} (ru_r) = \frac{dr}{dt} u_r + r \frac{d\theta}{dt} u_\theta.
\]

The outward speed is \( dr/dt \). The circular speed is \( r \frac{d\theta}{dt} \). The sum of squares is \( |v|^2 \).

Return one more time to steady motion around a circle, say \( r = 3 \) and \( \theta = 2t \). The velocity is \( v = 6u_\theta \), all circular. The acceleration is \(-12u_r\), all inward. For circles \( u_\theta \) is the tangent vector \( \mathbf{T} \). But the unit vector \( u_r \) points outward and \( \mathbf{N} \) points inward—the way the curve turns.

Now we tackle acceleration for any motion in polar coordinates. There can be speedup in \( r \) and speedup in \( \theta \) (also change of direction). Differentiate \( v \) in (5) by the product rule:

\[
\frac{dv}{dt} = \frac{d^2r}{dt^2} u_r + \frac{dr}{dt} \frac{du_r}{dt} + \frac{d\theta}{dt} u_\theta + r \frac{d^2\theta}{dt^2} u_\theta + r \frac{d\theta}{dt} \frac{du_r}{dt}.
\]

For \( du_r/dt \) and \( du_\theta/dt \), multiply equation (4) by \( d\theta/dt \). Then all terms contain \( u_r \) or \( u_\theta \).

The formula for \( \mathbf{a} \) is famous but not popular (except it got us to the moon):

\[
\mathbf{a} = \frac{dv}{dt} = \left( \frac{d^2r}{dt^2} - r \left( \frac{d\theta}{dt} \right)^2 \right) u_r + \left( \frac{r d^2\theta}{dt^2} + 2 \frac{dr}{dt} \frac{d\theta}{dt} \right) u_\theta.
\]

In the steady motion with \( r = 3 \) and \( \theta = 2t \), only one acceleration term is nonzero: \( \mathbf{a} = -12u_r \). Formula (6) can be memorized (maybe). Problem 14 gives a new way to reach it, using \( re^\theta \).

**EXAMPLE 1** Find \( \mathbf{R} \) and \( v \) and \( \mathbf{a} \) for speedup \( \theta = t^2 \) around the circle \( r = 1 \).

Solution The position vector is \( \mathbf{R} = u_r \). Then \( v \) and \( \mathbf{a} \) come from (5–6):

\[
v = (r \frac{d\theta}{dt})u_\theta = 2ru_\theta \quad \text{and} \quad \mathbf{a} = -(2t)^2u_r + 2u_\theta.
\]

This question and answer were also in Example 6 of the previous section. The acceleration was \( 2\mathbf{T} + 4t^2\mathbf{N} \). Notice again that \( \mathbf{T} = u_\theta \) and \( \mathbf{N} = -u_r \), going round the circle.

**EXAMPLE 2** Find \( \mathbf{R} \) and \( v \) and \( |v| \) and \( \mathbf{a} \) for the spiral motion \( r = 3t \), \( \theta = 2t \).

Solution The position vector is \( \mathbf{R} = 3ru_r \). Equation (5) gives velocity and speed:

\[
v = 3u_r + 6ru_\theta \quad \text{and} \quad |v| = \sqrt{(3)^2 + (6t)^2}.
\]
12 Motion Along a Curve

The motion goes out and also around. From (6) the acceleration is $-12t \mathbf{u} + 12\mathbf{u}_\theta$. The same answers would come more slowly from $\mathbf{R} = 3t \cos 2t \mathbf{i} + 3t \sin 2t \mathbf{j}$.

This example uses polar coordinates, but the motion is not circular. One of Kepler's inspirations, after many struggles, was to get away from circles.

**KEPLER'S LAWS**

You may know that before Newton and Leibniz and calculus and polar coordinates, Johannes Kepler discovered three laws of planetary motion. He was the court mathematician to the Holy Roman Emperor, who mostly wanted predictions of wars. Kepler also determined the date of every Easter—no small problem. His triumph was to discover patterns in the observations made by astronomers (especially by Tycho Brahe). Galileo and Copernicus expected circles, but Kepler found ellipses.

**Law 1:** Each planet travels in an ellipse with one focus at the sun.

**Law 2:** The vector from sun to planet sweeps out area at a steady rate: $\frac{dA}{dt} = \text{constant}$.

**Law 3:** The length of the planet's year is $T = ka^{3/2}$, where $a =$ maximum distance from the center (not the sun) and $k = 2\pi/\sqrt{GM}$ is the same for all planets.

With calculus the proof of these laws is a thousand times quicker. But Law 2 is the only easy one. The sun exerts a central force. Equation (I) gave $\mathbf{R} \times \mathbf{v} = \mathbf{H} = \text{constant}$ for central forces. Replace $\mathbf{R}$ by $ru$, and replace $\mathbf{v}$ by equation (5):

$$\mathbf{H} = ru \times \left(\frac{dr}{dt} \mathbf{u} + r \frac{d\theta}{dt} \mathbf{u}_\theta\right) = r^2 \frac{d\theta}{dt} (u_r \times u_\theta).$$

(7)

This vector $\mathbf{H}$ is constant, so its length $h = r^2 d\theta/dt$ is constant. In polar coordinates, the area is $dA = \frac{1}{2} r^2 d\theta$. This area $dA$ is swept out by the planet (Figure 12.10), and we have proved Law 2:

$$\frac{dA}{dt} = \frac{1}{2} r^2 \frac{d\theta}{dt} = \frac{1}{2} h = \text{constant}. \quad (8)$$

Near the sun $r$ is small. So $d\theta/dt$ is big and planets go around faster.

**Fig. 12.10** The planet is on an ellipse with the sun at a focus. Note $a$, $b$, $c$, $q$.

Now for Law 1, about ellipses. We are aiming for $1/r = C - D \cos \theta$, which is the polar coordinate equation of an ellipse. It is easier to write $q$ than $1/r$, and find an equation for $q$. The equation we will reach is $d^2 q/d\theta^2 + q = C$. The desired $q = C - D \cos \theta$ solves that equation (check this), and gives us Kepler's ellipse.
The first step is to connect $dr/dt$ to $dq/d\theta$ by the chain rule:

$$\frac{dr}{dt} = \frac{d}{dt} \left( \frac{1}{r} \right) = -\frac{1}{r^2} \frac{dq}{dt} = -\frac{1}{q^2} \frac{dq}{d\theta} \frac{d\theta}{dt} = -\frac{h}{q^2} \frac{dq}{d\theta} \frac{d\theta}{dt} = -\frac{h}{q^2} \frac{dq}{d\theta}.$$ \hfill (9)

Notice especially $d\theta/dt = h/r^2 = hq^2$. What we really want are second derivatives:

$$\frac{d^2r}{dt^2} = -\frac{h}{r^2} \frac{d}{dt} \left( \frac{dq}{d\theta} \right) = -\frac{h}{q^2} \frac{d}{d\theta} \left( \frac{dq}{d\theta} \right) \frac{d\theta}{dt} = -\frac{h^2q^2}{d\theta^2}.$$ \hfill (10)

After this trick of introducing $q$, we are ready for physics. The planet obeys Newton's Law $F = ma$, and the central force $F$ is the sun's gravity:

$$\frac{F}{m} = a = -\frac{GM}{r^2} = \frac{d^2r}{dt^2} = -r \left( \frac{d\theta}{dt} \right)^2.$$ \hfill (11)

That right side is the $u$, component of $a$ in (6). Change $r$ to $1/q$ and change $d\theta/dt$ to $hq^2$. The preparation in (10) allows us to rewrite $d^2r/dt^2$ in equation (11). That equation becomes

$$-GM q^2 = -h^2q^2 \frac{d^2q}{d\theta^2} - \frac{1}{q^2} (hq^2)^2.$$ 

Dividing by $-h^2q^2$ gives what we hoped for—the simple equation for $q$:

$$\frac{d^2q}{d\theta^2} + q = GM/h^2 = C \quad (a \text{ constant}).$$ \hfill (12)

**The solution is** $q = C - D \cos \theta$. Section 9.3 gave this polar equation for an ellipse or parabola or hyperbola. To be sure it is an ellipse, an astronomer computes $C$ and $D$ from the sun's mass $M$ and the constant $G$ and the earth's position and velocity. The main point is that $C > D$. Then $q$ is never zero and $r$ is never infinite. Hyperbolas and parabolas are ruled out, and the orbit in Figure 12.10 must be an ellipse.†

Astronomy is really impressive. You should visit the Greenwich Observatory in London, to see how Halley watched his comet. He amazed the world by predicting the day it would return. Also the discovery of Neptune was pure mathematics—the path of Uranus was not accounted for by the sun and known planets. LeVerrier computed a point in the sky and asked a Berlin astronomer to look. Sure enough Neptune was there.

Recently one more problem was solved—to explain the gap in the asteroids around Jupiter. The reason is "chaos"—the three-body problem goes unstable and an asteroid won't stay in that orbit. We have come a long way from circles.

**Department of Royal Mistakes** The last pound note issued by the Royal Mint showed Newton looking up from his great book *Principia Mathematica*. He is not smiling and we can see why. The artist put the sun at the center! Newton has just proved it is at the focus. True, the focus is marked $S$ and the planet is $P$. But those rays at the center brought untold headaches to the Mint—the note is out of circulation. I gave an antique dealer three pounds for it (in coins).

**Kepler's third law** gives the time $T$ to go around the ellipse—the planet's year. What is special in the formula is $a^{3/2}$—and for Kepler himself, the 15th of May 1618 was unforgettable: "the right ratio outfought the darkness of my mind, by the great proof afforded by my labor of seventeen years on Brahe's observations." The second

---

†An amateur sees the planet come around again, and votes for an ellipse.
law $\frac{dA}{dt} = \frac{1}{2} h$ is the key, plus two facts about an ellipse—its area $\pi ab$ and the height $b^2/a$ above the sun:

1. The area $A = \int_0^T \frac{dA}{dt} dt = \frac{1}{2} hT$ must equal $\pi ab$, so $T = \frac{2\pi ab}{h}$

2. The distance $r = 1/C$ at $\theta = \pi/2$ must equal $b^2/a$, so $b = \sqrt{a/C}$.

The height $b^2/a$ is in Figure 12.10 and Problems 25–26. The constant $C = GM/h^2$ is in equation (12). Put them together to find the period:

$$T = \frac{2\pi ab}{h} = \frac{2\pi}{h} \sqrt{\frac{a}{C}} = \frac{2\pi}{\sqrt{GM}} a^{3/2}.$$  \hspace{1cm} (13)

To think of Kepler guessing $a^{3/2}$ is amazing. To think of Newton proving Kepler's laws by calculus is also wonderful—because we can do it too.

**EXAMPLE 3** When a satellite goes around in a circle, find the time $T$.

Let $r$ be the radius and $\omega$ be the angular velocity. The time for a complete circle (angle $2\pi$) is $T = 2\pi/\omega$. The acceleration is $GM/r^2$ from gravity, and it is also $r\omega^2$ for circular motion. Therefore Kepler is proved right:

$$r\omega^2 = GM/r^2 \Rightarrow \omega = \sqrt{GM/r^3} \Rightarrow T = 2\pi/\omega = 2\pi r^{3/2}/\sqrt{GM}.$$ 

### 12.4 EXERCISES

**Read-through questions**

A central force points toward __a__. Then $\mathbf{R} \times d^2\mathbf{R}/dt^2 = 0$ because __b__. Therefore $\mathbf{R} \times d\mathbf{R}/dt$ is a __c__ (called H).

In polar coordinates, the outward unit vector is $\mathbf{u}_r = \cos \theta \mathbf{i} + \sin \theta \mathbf{j}$. Rotated by 90° this becomes $\mathbf{u}_r = \mathbf{e}_\theta$. The position vector $\mathbf{R}$ is the distance $r$ times __d__. The velocity $\mathbf{v} = d\mathbf{R}/dt$ is __e__ $\mathbf{u}_r + \mathbf{f}$ $\mathbf{u}_\theta$. For steady motion around the circle $r = 5$ with $\theta = 4t$, $\mathbf{v}$ is __g__ and $|\mathbf{v}|$ is __h__ and a is __i__.

For motion under a circular force, $r^2$ times __j__ is constant. Dividing by 2 gives Kepler's second law $dA/dt = __k__$. The first law says that the orbit is an __l__ with the sun at __m__. The polar equation for a conic section is __n__ $= C - D \cos \theta$. Using $\mathbf{F} = m\mathbf{a}$ we found $q_{00} + __o__ = C$. So the path is a conic section; it must be an ellipse because __p__.

The properties of an ellipse lead to the period $T = __q__$, which is Kepler's third law.
12.4 Polar Coordinates and Planetary Motion

1. Find the unit vectors \( \mathbf{u}_r \) and \( \mathbf{u}_\theta \) at the point \((0, 2)\). The \( \mathbf{u}_r \) and \( \mathbf{u}_\theta \) components of \( \mathbf{v} = i + j \) at that point are ________.

2. Find \( \mathbf{u}_r \) and \( \mathbf{u}_\theta \) at \((3, 3)\). If \( \mathbf{v} = i + j \) then \( \mathbf{v} = \) ________ \( \mathbf{u}_r \) and \( \mathbf{u}_\theta \) components of \( \mathbf{v} = i + j \) with the complex number \( re^{i\theta} \) and take its derivatives: \( d\mathbf{v}/dt = \) ________ and \( d\theta/dt = \) ________.

3. At the point \((1, 2)\), velocities in the direction ________ will give \( dr/dt = 0 \). Velocities in the direction ________ will give \( d\theta/dt = 0 \).

4. Traveling on the cardioid \( r = 1 - \cos \theta \) with \( d\theta/dt = 2 \), what is \( v \)? How long to go around the cardioid (no integration involved)?

5. If \( r = e^t \) and \( \theta = 2t \), find \( v \) and \( \mathbf{a} \) when \( t = 1 \).

6. If \( r = 1 \) and \( \theta = \sin t \), describe the path and find \( v \) and \( \mathbf{a} \) from equations (5-6). Where is the velocity zero?

7. (Important) \( \mathbf{R} = 4 \cos 5t \mathbf{i} + 4 \sin 5t \mathbf{j} \) travels on a circle of radius 4 with \( \theta = 5t \) and speed 20. Find the components of \( v \) and \( \mathbf{a} \) in three systems: \( \mathbf{i} \) and \( \mathbf{j} \), \( \mathbf{T} \) and \( \mathbf{N} \), \( v \) and \( \mathbf{a} \) at the return to the starting point \((4, 0)\).

8. When is the circle \( r = 4 \) completed, if the speed is \( 8t \)? Find \( v \) and \( \mathbf{a} \) at the return to the starting point \((4, 0)\).

9. The \( \mathbf{u}_r \) component of acceleration is ________ = 0 for a central force, which is in the direction of ________.

10. If \( r^2 d\theta/dt = 2 \) for travel up the line \( x = 1 \), draw a triangle to show that \( r = \sec \theta \) and integrate to find the time to reach \((1, 1)\).

11. A satellite is \( r = 10,000 \) km from the center of the Earth, traveling perpendicular to the radius vector at 4 km/sec. Find \( dr/dt \) and \( \mathbf{a} \).

12. From \( |\mathbf{v}| = 1 \), it follows that \( du_r/dr \) and \( du_\theta/d\theta \) are ________ to \( \mathbf{u}_r \) (Section 12.3). In fact \( du_r/dr \) is ________ and \( du_\theta/d\theta \) is ________.

13. Momentum is \( m\mathbf{v} \) and its derivative is \( m\mathbf{a} \) = force. Angular momentum is \( m\mathbf{R} \times \mathbf{v} \) and its derivative is ________ = torque. Angular momentum is constant under a central force because the ________ is zero.

14. To find (and remember) \( \mathbf{v} \) and \( \mathbf{a} \) in polar coordinates, start with the complex number \( re^{i\theta} \) and take its derivatives:

\[
\frac{d\mathbf{R}}{dt} = \frac{dr}{dt} e^{i\theta} + ir \frac{d\theta}{dt} e^{i\theta}
\]

\[
\frac{d^2\mathbf{R}}{dt^2} = \quad + \quad + \quad + \quad + \quad .
\]

Key idea: The coefficients of \( e^{i\theta} \) and \( ie^{i\theta} \) are the \( \mathbf{u}_r \) and \( \mathbf{u}_\theta \) components of \( \mathbf{R} \), \( \mathbf{v} \), \( \mathbf{a} \):

\[
\mathbf{R} = ru_r + 0u_\theta \quad \mathbf{v} = \frac{dr}{dt} \mathbf{u}_r + r\frac{d\theta}{dt} \mathbf{u}_\theta \quad \mathbf{a} = \quad .
\]

(a) Fill in the five terms from the derivative of \( d\mathbf{R}/dt \)
(b) Convert \( e^{i\theta} \) to \( \mathbf{u}_r \) and \( ie^{i\theta} \) to \( \mathbf{u}_\theta \) to find \( \mathbf{a} \)
(c) Compare \( \mathbf{R} \), \( \mathbf{v} \), \( \mathbf{a} \) with formulas (5-6)
(d) (for instructors only) Why does this method work?

Note: how \( e^{i\theta} = \cos \theta + i \sin \theta \) corresponds to \( \mathbf{u}_r = \cos \theta \mathbf{i} + \sin \theta \mathbf{j} \). This is one place where electrical engineers are allowed to write \( j \) instead of \( i \) for \( \sqrt{-1} \).

15. If the period is \( T \) find from (13) a formula for the distance \( a \).

16. To stay above New York what should be the period of a satellite? What should be its distance \( a \) from the center of the Earth?

17. From \( T \) and \( a \) find a formula for the mass \( M \).

18. If the moon has a period of 28 days at an average distance of \( a = 380,000 \) km, estimate the mass of the ________.

19. The Earth takes 365.2 days to go around the sun at a distance \( a \approx 93 \) million miles \( \approx 150 \) million kilometers. Find the mass of the sun.

20. True or false:

(a) The paths of all comets are ellipses.
(b) A planet in a circular orbit has constant speed.
(c) Orbits in central force fields are conic sections.

21. \( \sqrt{GM} \approx 2 \times 10^7 \) in what units, based on the Earth's mass \( M = 6 \times 10^{24} \) kg and the constant \( G = 6.67 \times 10^{-11} \) Nm²/kg²? A force of one kg meter/sec² is a Newton N.

22. If a satellite circles the Earth at 9000 km from the center, estimate its period \( T \) in seconds.

23. The Viking 2 orbiter around Mars had a period of about 10,000 seconds. If the mass of Mars is \( M = 6.4 \times 10^{23} \) kg, what was the value of \( a \)?

24. Convert 1/r = \( C - D \cos \theta \), or 1 = \( C \mathbf{r} - D \mathbf{x} \), into the \( xy \) equation of an ellipse.

25. The distances \( a \) and \( c \) on the ellipse give the constants in \( r = 1/(C - D \cos \theta) \). Substitute \( \theta = 0 \) and \( \theta = \pi \) as in Figure 12.10 to find \( D = c/(a^2 - c^2) \) and \( C = a/(a^2 - c^2) = a/b^2 \).

26. Show that \( x = -c \), \( y = b^2/a \) lies on the ellipse \( x^2/a^2 + y^2/b^2 = 1 \). Thus \( y \) is the height \( 1/C \) above the sun in Figure 12.10. The distance from the sun to the center has \( c^2 = a^2 - b^2 \).

27. The point \( x = a \cos 2\pi t/T \), \( y = b \sin 2\pi t/T \) travels around an ellipse centered at \((0, 0)\) and returns at time \( T \). By symmetry it sweeps out area at the same rate at both ends of the major axis. Why does this break Kepler's second law?

28. If a central force is \( F = -ma(r)\mathbf{u}_r \), explain why \( d^2r/dt^2 - r(d\theta/dt)^2 = -a(r) \). What is \( a(r) \) for gravity? Equation (12) for \( q = 1/r \) leads to \( q_{\mathbf{u}_r} + q = r^2a(r) \).

29. When \( F = 0 \) the body should travel in a straight ________.
The equation $q_{pp} + q = 0$ allows $q = \cos \theta$, in which case the path $1/r = \cos \theta$ is _________. Extra credit: Mark off equal distances on a line, connect them to the sun, and explain why the triangles have equal area. So $dA/dt$ is still constant.

30. The strong nuclear force increases with distance, $a(r) = r$. It binds quarks so tightly that up to now no top quarks have been seen (reliably). Problem 28 gives $q_{pp} + q = 1/q^3$.
   (a) Multiply by $q$, and integrate to find $\frac{1}{2}q^2 - \frac{1}{2}q^2 = _____ + C$.

   *(b) Integrate again (with tables) after setting $u = q^2, u = 2qq^2$.

31. The path of a quark in 30(b) can be written as $r^2(A + B \cos 2\theta) = 1$. Show that this is the same as the ellipse $(A + B)x^2 + (A - B)y^2 = 1$ with the origin at the center. The nucleus is not at a focus, and the pound note is correct for Newton watching quarks. (Quantum mechanics not accounted for.)

32. When will Halley’s comet appear again? It disappeared in 1986 and its mean distance to the sun (average of $a + c$ and $a - c$) is $a = 1.6 \times 10^8$ kilometers.

33. You are walking at 2 feet/second toward the center of a merry-go-round that turns once every ten seconds. Starting from $r = 20, \theta = 0$ find $r(t), \theta(t), v(t), a(t)$ and the length of your path to the center.

34. From Kepler’s laws $r = 1/(C - D \cos \theta)$ and $r^2 d\theta/dt = h$, show that:

1. $dr/dt = -Dh \sin \theta$
2. $d^2r/dt^2 = \left(\frac{1}{r} - C\right)h^2/r^2$
3. $d^2r/dt^2 - r(d\theta/dt)^2 = -Ch^2/r^2$.

When Newton reached 3, he knew that Kepler’s laws required a central force of $Ch^2/r^2$. This is his inverse square law. Then he went backwards, in our equations (8–12), to show that this force yields Kepler’s laws.

35. How long is our year? The Earth’s orbit has $a = 149.57 \times 10^6$ kilometers.
### Contents

**CHAPTER 9  Polar Coordinates and Complex Numbers**

- 9.1 Polar Coordinates 348
- 9.2 Polar Equations and Graphs 351
- 9.3 Slope, Length, and Area for Polar Curves 356
- 9.4 Complex Numbers 360

**CHAPTER 10  Infinite Series**

- 10.1 The Geometric Series 368
- 10.2 Convergence Tests: Positive Series 374
- 10.3 Convergence Tests: All Series 381
- 10.4 The Taylor Series for $e^x$, $\sin x$, and $\cos x$ 385
- 10.5 Power Series 391

**CHAPTER 11  Vectors and Matrices**

- 11.1 Vectors and Dot Products 398
- 11.2 Planes and Projections 407
- 11.3 Cross Products and Determinants 416
- 11.4 Matrices and Linear Equations 425
- 11.5 Linear Algebra in Three Dimensions 435

**CHAPTER 12  Motion along a Curve**

- 12.1 The Position Vector 446
- 12.2 Plane Motion: Projectiles and Cycloids 453
- 12.3 Tangent Vector and Normal Vector 459
- 12.4 Polar Coordinates and Planetary Motion 464

**CHAPTER 13  Partial Derivatives**

- 13.1 Surfaces and Level Curves 472
- 13.2 Partial Derivatives 475
- 13.3 Tangent Planes and Linear Approximations 480
- 13.4 Directional Derivatives and Gradients 490
- 13.5 The Chain Rule 497
- 13.6 Maxima, Minima, and Saddle Points 504
- 13.7 Constraints and Lagrange Multipliers 514
CHAPTER 13

Partial Derivatives

This chapter is at the center of multidimensional calculus. Other chapters and other topics may be optional; this chapter and these topics are required. We are back to the basic idea of calculus—*the derivative*. There is a function \( f \), the variables move a little bit, and \( f \) moves. The question is how much \( f \) moves and how fast. Chapters 1–4 answered this question for \( f(x) \), a function of one variable. Now we have \( f(x, y) \) or \( f(x, y, z) \)—with two or three or more variables that move independently. As \( x \) and \( y \) change, \( f \) changes. The fundamental problem of differential calculus is to connect \( \Delta x \) and \( \Delta y \) to \( \Delta f \).

Calculus solves that problem in the limit. It connects \( dx \) and \( dy \) to \( df \). In using this language I am building on the work already done. You know that \( \frac{df}{dx} \) is the limit of \( \frac{Af}{\Delta x} \). Calculus computes the rate of change—which is the slope of the tangent line. The goal is to extend those ideas to

\[
\begin{align*}
f(x, y) &= x^2 - y^2 & \text{or} & \quad f(x, y) &= \sqrt{x^2 + y^2} & \text{or} & \quad f(x, y, z) &= 2x + 3y + 4z.
\end{align*}
\]

These functions have graphs, they have derivatives, and they must have tangents.

The heart of this chapter is summarized in six lines. The subject is *differential* calculus—small changes in a short time. Still to come is *integral* calculus—adding up those small changes. We give the words and symbols for \( f(x, y) \), matched with the words and symbols for \( f(x) \). Please use this summary as a guide, to know where calculus is going.

*Curve* \( y = f(x) \) *vs.* *Surface* \( z = f(x, y) \)

\[
\begin{align*}
\frac{df}{dx} & \quad \text{becomes two partial derivatives} & \frac{\partial f}{\partial x} \quad \text{and} \quad \frac{\partial f}{\partial y} \\
\frac{d^2f}{dx^2} & \quad \text{becomes four second derivatives} & \frac{\partial^2 f}{\partial x^2}, \frac{\partial^2 f}{\partial x \partial y}, \frac{\partial^2 f}{\partial y \partial x}, \frac{\partial^2 f}{\partial y^2} \\
\Delta f \approx \frac{df}{dx} \Delta x & \quad \text{becomes the linear approximation} & \Delta f \approx \frac{\partial f}{\partial x} \Delta x + \frac{\partial f}{\partial y} \Delta y \\
tangent \ line & \quad \text{becomes the tangent plane} & z - z_o = \frac{\partial f}{\partial x}(x - x_o) + \frac{\partial f}{\partial y}(y - y_o) \\
\frac{dy}{dt} = \frac{dy}{dx} \frac{dx}{dt} & \quad \text{becomes the chain rule} & \frac{dz}{dt} = \frac{\partial z}{\partial x} \frac{dx}{dt} + \frac{\partial z}{\partial y} \frac{dy}{dt} \\
\frac{df}{dx} = 0 & \quad \text{becomes two maximum-minimum equations} & \frac{\partial f}{\partial x} = 0 \quad \text{and} \quad \frac{\partial f}{\partial y} = 0.
\end{align*}
\]
13 Partial Derivatives

13.1 Surfaces and Level Curves

The graph of \( y = f(x) \) is a curve in the \( xy \) plane. There are two variables—\( x \) is independent and free, \( y \) is dependent on \( x \). Above \( x \) on the base line is the point \((x, y)\) on the curve. The curve can be displayed on a two-dimensional printed page.

The graph of \( z = f(x, y) \) is a surface in \( xyz \) space. There are three variables—\( x \) and \( y \) are independent, \( z \) is dependent. Above \((x, y)\) in the base plane is the point \((x, y, z)\) on the surface (Figure 13.1). Since the printed page remains two-dimensional, we shade or color or project the surface. The eyes are extremely good at converting two-dimensional images into three-dimensional understanding—they get a lot of practice. The mathematical part of our brain also has something new to work on—two partial derivatives.

This section uses examples and figures to illustrate surfaces and their level curves. The next section is also short. Then the work begins.

**EXAMPLE 1** Describe the surface and the level curves for \( z = f(x, y) = \sqrt{x^2 + y^2} \).

*The surface is a cone.* Reason: \( \sqrt{x^2 + y^2} \) is the distance in the base plane from \((0, 0)\) to \((x, y)\). When we go out a distance 5 in the base plane, we go up the same distance 5 to the surface. The cone climbs with slope 1. The distance out to \((x, y)\) equals the distance up to \(z\) (this is a \(45^\circ\) cone).

*The level curves are circles.* At height 5, the cone contains a circle of points—all at the same “level” on the surface. The plane \( z = 5 \) meets the surface \( z = \sqrt{x^2 + y^2} \) at those points (Figure 13.1b). The circle below them (in the base plane) is the level curve.

**DEFINITION** A level curve or contour line of \( z = f(x, y) \) contains all points \((x, y)\) that share the same value \( f(x, y) = c \). Above those points, the surface is at the height \( z = c \).

There are different level curves for different \( c \). To see the curve for \( c = 2 \), cut through the surface with the horizontal plane \( z = 2 \). The plane meets the surface above the points where \( f(x, y) = 2 \). The level curve in the base plane has the equation \( f(x, y) = 2 \). Above it are all the points at “level 2” or “level \( c \)” on the surface.

Every curve \( f(x, y) = c \) is labeled by its constant \( c \). This produces a contour map (the base plane is full of curves). For the cone, the level curves are given by \( \sqrt{x^2 + y^2} = c \), and the contour map consists of circles of radius \( c \).

**Question** What are the level curves of \( z = f(x, y) = x^2 + y^2 \)?

*Answer* Still circles. But the surface is not a cone (it bends up like a parabola). The circle of radius 3 is the level curve \( x^2 + y^2 = 9 \). On the surface above, the height is 9.

---

---

![Fig. 13.1](image-url) The surface for \( z = f(x, y) = \sqrt{x^2 + y^2} \) is a cone. The level curves are circles.
13.1 Surfaces and Level Curves

EXAMPLE 2 For the linear function \( f(x, y) = 2x + y \), the surface is a plane. Its level curves are straight lines. The surface \( z = 2x + y \) meets the plane \( z = c \) in the line \( 2x + y = c \). That line is above the base plane when \( c \) is positive, and below when \( c \) is negative. The contour lines are in the base plane. Figure 13.2b labels these parallel lines according to their height in the surface.

Question If the level curves are all straight lines, must they be parallel?

Answer No. The surface \( z = y/x \) has level curves \( y/x = c \). Those lines \( y = cx \) swing around the origin, as the surface climbs like a spiral playground slide.

Fig. 13.2 A plane has parallel level lines. The spiral slide \( z = y/x \) has lines \( y/x = c \).

EXAMPLE 3 The weather map shows contour lines of the temperature function. Each level curve connects points at a constant temperature. One line runs from Seattle to Omaha to Cincinnati to Washington. In winter it is painful even to think about the line through L.A. and Texas and Florida. USA Today separates the contours by color, which is better. We had never seen a map of universities.

Fig. 13.3 The temperature at many U.S. and Canadian universities. Mt. Monadnock in New Hampshire is said to be the most climbed mountain (except Fuji?) at 125,000/year. Contour lines every 6 meters.
Question From a contour map, how do you find the highest point?

Answer The level curves form loops around the maximum point. As $c$ increases the loops become tighter. Similarly the curves squeeze to the lowest point as $c$ decreases.

Example 4 A contour map of a mountain may be the best example of all. Normally the level curves are separated by 100 feet in height. On a steep trail those curves are bunched together—the trail climbs quickly. In a flat region the contour lines are far apart. Water runs perpendicular to the level curves. On my map of New Hampshire that is true of creeks but looks doubtful for rivers.

Question Which direction in the base plane is uphill on the surface?

Answer The steepest direction is perpendicular to the level curves. This is important. Proof to come.

Example 5 In economics $x^2y$ is a utility function and $x^2y = c$ is an indifference curve.

The utility function $x^2y$ gives the value of $x$ hours awake and $y$ hours asleep. Two hours awake and fifteen minutes asleep have the value $f = (2^2)(1)$. This is the same as one hour of each: $f = (1^2)(1)$. Those lie on the same level curve in Figure 13.4a. We are indifferent, and willing to exchange any two points on a level curve.

The indifference curve is “convex.” We prefer the average of any two points. The line between two points is up on higher level curves.

Figure 13.4b shows an extreme case. The level curves are straight lines $4x + y = c$. Four quarters are freely substituted for one dollar. The value is $f = 4x + y$ dollars.

Figure 13.4c shows the other extreme. Extra left shoes or extra right shoes are useless. The value (or utility) is the smaller of $x$ and $y$. That counts pairs of shoes.

Fig. 13.4 Utility functions $x^2y$, $4x + y$, $\min(x, y)$. Convex, straight substitution, complements.

13.1 Exercises

Read-through questions

The graph of $z = f(x, y)$ is a ___ in ___-dimensional space. The ___ curve $f(x, y) = 7$ lies down in the base plane. Above this level curve are all points at height ___ in the surface. The ___ $z = 7$ cuts through the surface at those points. The level curves $f(x, y) = ___$ are drawn in the $xy$ plane and labeled by ___. The family of labeled curves is a ___ map.

For $z = f(x, y) = x^2 - y^2$, the equation for a level curve is ____. This curve is a ___. For $z = x - y$ the curves are _____. Level curves never cross because ____. They crowd together when the surface is ___. The curves tighten to a point when ___. The steepest direction on a mountain is ____ to the ____.
1. Draw the surface \( z = f(x, y) \) for these four functions:
   \[ f_1 = \sqrt{4 - x^2 - y^2} \quad f_2 = 2 - \sqrt{x^2 + y^2} \]
   \[ f_3 = 2 - \frac{1}{2}(x^2 + y^2) \quad f_4 = 1 + e^{-x^2 - y^2} \]

2. The level curves of all four functions are ________. They enclose the maximum at _________. Draw the four curves \( f(x, y) = c \) and rank them by increasing radius.

3. Set \( y = 0 \) and compute the \( x \) derivative of each function at \( x = 2 \). Which mountain is flattest and which is steepest at that point?

4. Set \( y = 1 \) and compute the \( x \) derivative of each function at \( x = 1 \).

For \( f_5 \) to \( f_{10} \) draw the level curves \( f = 0, 1, 2 \). Also \( f = -4 \).

5. \( f_5 = x - y \)  \( f_6 = (x + y)^2 \)
6. \( f_7 = xe^{-y} \)  \( f_8 = \sin(x - y) \)
7. \( f_9 = y - x^2 \)  \( f_{10} = y/x^2 \)

11. Suppose the level curves are parallel straight lines. Does the surface have to be a plane?

12. Construct a function whose level curve \( f = 0 \) is in two separate pieces.

13. Construct a function for which \( f = 0 \) is a circle and \( f = 1 \) is not.

14. Find a function for which \( f = 0 \) has infinitely many pieces.

15. Draw the contour map for \( f = xy \) with level curves \( f = -2, -1, 0, 1, 2 \). Describe the surface.

16. Find a function \( f(x, y) \) whose level curve \( f = 0 \) consists of a circle and all points inside it.

Draw two level curves in 17–20. Are they ellipses, parabolas, or hyperbolas? Write \( \sqrt{-2x + c} \) as \( \sqrt{c + 2x} \) before squaring both sides.

17. \( f = \sqrt{4x^2 + y^2} \)
18. \( f = \sqrt{4x^2 + y^2 - 2x} \)
19. \( f = \sqrt{5x^2 + y^2 - 2x} \)
20. \( f = \sqrt{3x^2 + y^2 - 2x} \)

21. The level curves of \( f = (y - 2)/(x - 1) \) are ________ through the point \((1, 2)\) except that this point is not ________.

22. Sketch a map of the US with lines of constant temperature (isotherms) based on today's paper.

23. (a) The contour lines of \( z = x^2 + y^2 - 2x - 2y \) are circles around the point ________, where \( z \) is a minimum.
   (b) The contour lines of \( f = ________ \) are the circles \( x^2 + y^2 = c + 1 \) on which \( f = c \).

24. Draw a contour map of any state or country (lines of constant height above sea level). Florida may be too flat.

25. The graph of \( w = F(x, y, z) \) is a ________-dimensional surface in \( xyzw \) space. Its level sets \( F(x, y, z) = c \) are ________-dimensional surfaces in \( xyz \) space. For \( w = x - 2y + z \) those level sets are ________. For \( w = x^2 + y^2 + z^2 \) those level sets are ________.

26. The surface \( x^2 + y^2 - z^2 = -1 \) is in Figure 13.8. There is empty space when \( z^2 \) is smaller than 1 because ________.

27. The level sets of \( F = x^2 + y^2 + qz^2 \) look like footballs when \( q \) is ________, like basketballs when \( q \) is ________, and like frisbees when \( q \) is ________.

28. Let \( T(x, y) \) be the driving time from your home at \((0, 0)\) to nearby towns at \((x, y)\). Draw the level curves.

29. (a) The level curves of \( f(x, y) = \sin(x - y) \) are ________.
   (b) The level curves of \( g(x, y) = \sin(x^2 - y^2) \) are ________.
   (c) The level curves of \( h(x, y) = \sin(x - y) \) are ________.

30. Prove that if \( x_1 y_1 = 1 \) and \( x_2 y_2 = 1 \) then their average \( x = \frac{1}{2}(x_1 + x_2), y = \frac{1}{2}(y_1 + y_2) \) has \( xy \geq 1 \). The function \( f = xy \) has convex level curves (hyperbolas).

31. The hours in a day are limited by \( x + y = 24 \). Write \( x^2 y \) as \( x^2(24 - x) \) and maximize to find the optimal number of hours to stay awake.

32. Near \( x = 16 \) draw the level curve \( x^2 y = 2048 \) and the line \( x + y = 24 \). Show that the curve is convex and the line is tangent.

33. The surface \( z = 4x + y \) is a ________. The surface \( z = \min(x, y) \) is formed from two ________. We are willing to exchange 6 left and 2 right shoes for 2 left and 4 right shoes but better is the average ________.

34. Draw a contour map of the top of your shoe.

13.2 Partial Derivatives

The central idea of differential calculus is the derivative. A change in \( x \) produces a change in \( f \). The ratio \( \Delta f/\Delta x \) approaches the derivative, or slope, or rate of change.

What to do if \( f \) depends on both \( x \) and \( y \)?

The new idea is to vary \( x \) and \( y \) one at a time. First, only \( x \) moves. If the function is \( x + xy \), then \( \Delta f \) is \( \Delta x + y \Delta x \). The ratio \( \Delta f/\Delta x \) is \( 1 + y \). The "x derivative" of \( x + xy \)
is 1 + y. For all functions the method is the same: Keep y constant, change x, take the limit of \( \Delta f/\Delta x \):

\[
\frac{\partial f}{\partial x}(x, y) = \lim_{\Delta x \to 0} \frac{\Delta f}{\Delta x} = \lim_{\Delta x \to 0} \frac{f(x + \Delta x, y) - f(x, y)}{\Delta x}.
\]

On the left is a new symbol \( \partial f/\partial x \). It signals that only x is allowed to vary—\( \partial f/\partial x \) is a partial derivative. The different form \( \partial \) of the same letter (still say "d") is a reminder that x is not the only variable. Another variable y is present but not moving.

**EXAMPLE 1** \( f(x, y) = x^2y^2 + xy + y \) \( \frac{\partial f}{\partial x} (x, y) = 2xy^2 + y + 0 \).

Do not treat y as zero! Treat it as a constant, like 6. Its x derivative is zero. If \( f(x) = \sin 6x \) then \( \partial f/\partial x = 6 \cos 6x \). If \( f(x, y) = \sin xy \) then \( \partial f/\partial x = y \cos xy \).

Spoken aloud, \( \partial f/\partial x \) is still "\( \partial f/\partial x \). It is a function of x and y. When more is needed, call it "the partial of \( f \) with respect to x." The symbol \( f' \) is no longer available, since it gives no special indication about x. Its replacement \( f_x \) is pronounced "\( f \) sub x," which is shorter than \( \partial f/\partial x \) and means the same thing.

We may also want to indicate the point \( (x_0, y_0) \) where the derivative is computed:

\[
\frac{\partial f}{\partial x}(x_0, y_0) \quad \text{or} \quad f_x(x_0, y_0) \quad \text{or} \quad \frac{\partial f}{\partial x}\bigg|_{(x_0, y_0)} \quad \text{or just} \quad \left( \frac{\partial f}{\partial x} \right)_0.
\]

**EXAMPLE 2** \( f(x, y) = \sin 2x \cos y \) \( f_x = 2 \cos 2x \cos y \) (cos y is constant for \( \partial /\partial x \))

The particular point \( (x_0, y_0) \) is \( (0, 0) \). The height of the surface is \( f(0, 0) = 0 \). The slope in the x direction is \( f_x = 2 \). At a different point \( x_0 = \pi, y_0 = \pi \) we find \( f_x(\pi, \pi) = -2 \).

Now keep x constant and vary y. The ratio \( \Delta f/\Delta y \) approaches \( \partial f/\partial y \):

\[
f_y(x, y) = \lim_{\Delta y \to 0} \frac{\Delta f}{\Delta y} = \lim_{\Delta y \to 0} \frac{f(x, y + \Delta y) - f(x, y)}{\Delta y}.
\]

This is the slope in the y direction. Please realize that a surface can go up in the x direction and down in the y direction. The plane \( f(x, y) = 3x - 4y \) has \( f_x = 3 \) (up) and \( f_y = -4 \) (down). We will soon ask what happens in the 45° direction.

**EXAMPLE 3** \( f(x, y) = \sqrt{x^2 + y^2} \) \( \frac{\partial f}{\partial x} = \frac{x}{\sqrt{x^2 + y^2}} = \frac{\partial f}{\partial y} = \frac{y}{\sqrt{x^2 + y^2}} \)

The x derivative of \( \sqrt{x^2 + y^2} \) is really one-variable calculus, because y is constant. The exponent drops from \( \frac{1}{2} \) to \( \frac{-1}{2} \), and there is \( 2x \) from the chain rule. This distance function has the curious derivative \( \partial f/\partial x = x/f \).

The graph is a cone. Above the point \( (0, 2) \) the height is \( \sqrt{0^2 + 2^2} = 2 \). The partial derivatives are \( f_x = 0/2 \) and \( f_y = 2/2 \). At that point, Figure 13.5 climbs in the y direction. It is level in the x direction. An actual step \( \Delta x \) will increase \( O^2 + 2^2 \) to \( (\Delta x)^2 + 2^2 \). But this change is of order \( (\Delta x)^2 \) and the x derivative is zero.

Figure 13.5 is rather important. It shows how \( \partial f/\partial x \) and \( \partial f/\partial y \) are the ordinary derivatives of \( f(x, y_0) \) and \( f(x_0, y) \). It is natural to call these partial functions. The first has y fixed at \( y_0 \) while x varies. The second has x fixed at \( x_0 \) while y varies. Their graphs are cross sections down the surface—cut out by the vertical planes \( y = y_0 \) and \( x = x_0 \). Remember that the level curve is cut out by the horizontal plane \( z = c \).
Partial Derivatives

The limits of \( \Delta f/\Delta x \) and \( \Delta f/\Delta y \) are computed as always. With partial functions we are back to a single variable. **The partial derivative is the ordinary derivative of a partial function** (constant \( y \) or constant \( x \)). For the cone, \( \partial f/\partial y \) exists at all points except \((0, 0)\). The figure shows how the cross section down the middle of the cone produces the absolute value function: \( f(0, y) = |y| \). It has one-sided derivatives but not a two-sided derivative.

Similarly \( \partial f/\partial x \) will not exist at the sharp point of the cone. We develop the idea of a continuous function \( f(x, y) \) as needed (the definition is in the exercises). Each partial derivative involves one direction, but limits and continuity involve all directions. The distance function is continuous at \((0, 0)\), where it is not differentiable.

**EXAMPLE 4**

\[
f(x, y) = y^2 - x^2 \quad \partial f/\partial x = -2x \quad \partial f/\partial y = 2y
\]

Move in the \( x \) direction from \((1, 3)\). Then \( y^2 - x^2 \) has the partial function \( 9 - x^2 \). With \( y \) fixed at 3, a parabola opens downward. In the \( y \) direction (along \( x = 1 \)) the partial function \( y^2 - 1 \) opens upward. The surface in Figure 13.6 is called a *hyperbolic paraboloid*, because the level curves \( y^2 - x^2 = c \) are hyperbolas. Most people call it a saddle, and the special point at the origin is a saddle point.

The origin is special for \( y^2 - x^2 \) because both derivatives are zero. **The bottom of the \( y \) parabola at \((0, 0)\) is the top of the \( x \) parabola.** The surface is momentarily flat in all directions. It is the top of a hill and the bottom of a mountain range at the same time.
time. A saddle point is neither a maximum nor a minimum, although both derivatives are zero.

Note Do not think that \( f(x, y) \) must contain \( y^2 \) and \( x^2 \) to have a saddle point. The function \( 2xy \) does just as well. The level curves \( 2xy = c \) are still hyperbolas. The partial functions \( 2xy \) and \( 2x^2y \) now give straight lines—which is remarkable. Along the \( 45^\circ \) line \( x = y \), the function is \( 2x^2 \) and climbing. Along the \( -45^\circ \) line \( x = -y \), the function is \( -2x^2 \) and falling. The graph of \( 2xy \) is Figure 13.6 rotated by \( 45^\circ \).

**EXAMPLES 5–6** \( f(x, y, z) = x^2 + y^2 + z^2 \) \( P(T, V) = nRT/V \)

Example 5 shows more variables. Example 6 shows that the variables may not be named \( x \) and \( y \). Also, the function may not be named \( f \)! Pressure and temperature and volume are \( P \) and \( T \) and \( V \). The letters change but nothing else:

\[
\frac{\partial P}{\partial T} = nR/V \quad \frac{\partial P}{\partial V} = -nRT/V^2 \quad \text{(note the derivative of } 1/V)\]

There is no \( \partial P/\partial R \) because \( R \) is a constant from chemistry—not a variable.

Physics produces six variables for a moving body—the coordinates \( x, y, z \) and the momenta \( p_x, p_y, p_z \). Economics and the social sciences do better than that. If there are 26 products there are 26 variables—sometimes 52, to show prices as well as amounts. The profit can be a complicated function of these variables. The partial derivatives are the marginal profits, as one of the 52 variables is changed. A spreadsheet shows the 52 values and the effect of a change. An infinitesimal spreadsheet shows the derivative.

**SECOND DERIVATIVE**

Genius is not essential, to move to second derivatives. The only difficulty is that two first derivatives \( f_x \) and \( f_y \) lead to four second derivatives \( f_{xx} \) and \( f_{xy} \) and \( f_{yx} \) and \( f_{yy} \). (Two subscripts: \( f_{xx} \) is the \( x \) derivative of the \( x \) derivative. Other notations are \( \partial^2 f/\partial x^2 \) and \( \partial^2 f/\partial x \partial y \) and \( \partial^2 f/\partial y \partial x \) and \( \partial^2 f/\partial y^2 \).) Fortunately \( f_{xy} \) equals \( f_{yx} \), as we see first by example.

**EXAMPLE 7** \( f = x/y \) has \( f_x = 1/y \), which has \( f_{xx} = 0 \) and \( f_{xy} = -1/y^2 \).

The function \( x/y \) is linear in \( x \) (which explains \( f_{xx} = 0 \)). Its \( y \) derivative is \( f_y = -x/y^2 \).

This has the \( x \) derivative \( f_{yx} = -1/y^2 \). The mixed derivatives \( f_{xy} \) and \( f_{yx} \) are equal.

In the pure \( y \) direction, the second derivative is \( f_{yy} = 2x/y^3 \). One-variable calculus is sufficient for all these derivatives, because only one variable is moving.

**EXAMPLE 8** \( f = 4x^2 + 3xy + y^2 \) has \( f_x = 8x + 3y \) and \( f_y = 3x + 2y \).

Both “cross derivatives” \( f_{xy} \) and \( f_{yx} \) equal 3. The second derivative in the \( x \) direction is \( \partial^2 f/\partial x^2 = 8 \) or \( f_{xx} = 8 \). Thus “\( f \, x \, x \)” is “\( d \, f \, d \, x \, squared. \)” Similarly \( \partial^2 f/\partial y^2 = 2 \). The only change is from \( d \) to \( \partial \).

If \( f(x, y) \) has continuous second derivatives then \( f_{xy} = f_{yx} \). Problem 43 sketches a proof based on the Mean Value Theorem. For third derivatives almost any example shows that \( f_{xxx} = f_{xxy} = f_{xyx} \) is different from \( f_{yyy} = f_{yyx} = f_{xyy} \).

**Question** How do you plot a space curve \( x(t), y(t), z(t) \) in a plane? One way is to look parallel to the direction (1, 1, 1). On your XY screen, plot \( X = (y - x)/\sqrt{2} \) and \( Y = (2z - x - y)/\sqrt{6} \). The line \( x = y = z \) goes to the point \( (0, 0) \)!
13.2 Partial Derivatives

How do you graph a surface \( z = f(x, y) \)? Use the same \( X \) and \( Y \). Fix \( x \) and let \( y \) vary, for curves one way in the surface. Then fix \( y \) and vary \( x \), for the other partial function. For a parametric surface like \( x = (2 + v \sin \frac{1}{2}u) \cos u, y = (2 + v \sin \frac{1}{2}u) \sin u, z = v \cos \frac{1}{2}u \), vary \( u \) and then \( v \). Dick Williamson showed how this draws a one-sided “Möbius strip.”

Read-through questions

The derivative \( \frac{\partial f}{\partial y} \) comes from fixing \( b \) and moving \( c \). It is the limit of \( \frac{f(c) - f(b)}{c - b} \). If \( f = e^{2x} \sin y \) then \( \frac{\partial f}{\partial x} = \frac{\partial f}{\partial y} \) and \( \frac{\partial f}{\partial y} = \frac{\partial f}{\partial x} \). If \( f = \sqrt{x^2 + y^2} \) then \( f_x = \frac{y}{x} \) and \( f_y = \frac{x}{y} \). At \((x_0, y_0)\) the partial derivative \( f_x \) is the ordinary derivative of the \( f(x, y) \) function. Similarly \( f_y \) comes from \( f(y, x) \). Those functions are cut out by vertical planes \( x = x_0 \) and \( y = y_0 \), while the level curves are cut out by \( z = f(x, y) \) planes.

The four second derivatives are \( f_{xx}, f_{yy}, f_{xy}, f_{yx} \). For \( f = xy \) they are \( f_{xx} = -2 \), \( f_{yy} = -2 \), \( f_{xy} = 1 \), \( f_{yx} = 1 \). In those examples the derivatives \( f_{xx} \) and \( f_{yy} \) are the same. That is always true when the second derivatives are the same. At the origin, \( \cos 2x \cos 3y \) is curving \( u \) in the \( x \) and \( y \) directions, while \( \sin xy \) goes \( v \) in the \( 45\degree \) direction and \( w \) in the \(-45\degree \) direction.

Find \( \frac{\partial f}{\partial x} \) and \( \frac{\partial f}{\partial y} \) for the functions in 1–12.

1. \( 3x - y + x^2y^2 \)
2. \( \sin(3x - y) + y \)
3. \( x^3y^2 - x^2 - e^x \)
4. \( xe^{x+y} \)
5. \( (x+y)/(x-y) \)
6. \( 1/\sqrt{x^2 + y^2} \)
7. \( (x^2 + y^2)^{-1} \)
8. \( \ln(x+2y) \)
9. \( \ln(x^2 + y^2) \)
10. \( y^2 \)
11. \( \tan^{-1}(y/x) \)
12. \( \ln(xy) \)

Compute \( f_{xx}, f_{yy}, f_{xy}, f_{yx} \), for the functions in 13–20.

13. \( x^2 + 3xy + 2y^2 \)
14. \( x + 3y^2 \)
15. \( (x + iy)^3 \)
16. \( e^{x+iy} \)
17. \( 1/\sqrt{x^2 + y^2} \)
18. \( (x + y)^2 \)
19. \( \cos ax \cos by \)
20. \( 1/(x + iy) \)

Find the domain and range (all inputs and outputs) for the functions 21–26. Then compute \( f_x, f_y, f_x, f_y \).

21. \( 1/(x - y)^2 \)
22. \( \sqrt{x^2 + y^2 - t^2} \)
23. \( (y - x)/(x - t) \)
24. \( \ln(x + t) \)

Exercises

25. \( x^2 + y^2 \) Why does this equal \( e^{x+y} \)?
26. \( \cos x \cos^{-1}y \)

27. Verify \( f_{xx} = f_{yy} \) for \( f = x^2y^2 \). If \( f_{xx} = 0 \) then \( f_x \) does not depend on \( x \), and it is independent of \( y \). The function must have the form \( f(x, y) = G(x) + \) constant.

28. In terms of \( u \), compute \( f_x \) and \( f_y \) for \( f(x, y) = \int u \, dt \). First vary \( x \). Then vary \( y \).

29. Compute \( \frac{\partial f}{\partial x} \) for \( f(x, y) = \int e^t \, dt \). Keep \( y \) constant.

30. What is \( f(x, y) = \int e^t \, dt \) and what are \( f_x \) and \( f_y \)?

31. Calculate all eight third derivatives \( f_{xxx}, f_{xxy}, f_{xx}, f_{xy}, f_{yy}, f_{yy} \), \( f_{yy}, f_{yy}, f_{yy} \), \( f_{yy}, f_{yy}, f_{yy} \), \( f_{yy}, f_{yy}, f_{yy} \). How many are different?

32–35, equation.

32. \( f_x + f_y = 0 \)
33. \( f_x = f_y \)
34. \( f_x = f_y \)
35. \( f_{xx} = 4f_{yy} \)

36. Show that \( t^{-1/2} e^{-x^2/4t} \) satisfies the heat equation \( f_t = f_{xx} \). This \( f(x, t) \) is the temperature at position \( x \) and time \( t \) due to a point source of heat at \( x = 0, t = 0 \).

37. The equation for heat flow in the \( xy \) plane is \( f_t = f_{xx} + f_{yy} \). Show that \( f(x, y, t) = e^{-2t} \sin x \sin y \) is a solution. What exponent in \( f = e^{-2t} \sin 2x \sin 3y \) gives a solution?

38. Find solutions \( f(x, y) = e^{-2t} \sin mx \cos ny \) of the heat equation \( f_t = f_{xx} + f_{yy} \). Show that \( f(x, y, t) = e^{-2t} \sin 2x \sin 3y \) is also a solution.

39. The basic wave equation \( f_{tt} = f_{xx} \). Verify that \( f(x, t) = \sin(x + t) \) and \( f(x, t) = \sin(x - t) \) are solutions. Draw both graphs at \( t = \pi/4 \). Which wave moved to the left and which moved to the right?

40. Continuing 39, the peaks of the waves moved a distance \( \Delta x \) in the time step \( \Delta t = \pi/4 \). The wave velocity is \( \Delta x/\Delta t = \) constant.

41. Which of these satisfy the wave equation \( f_{tt} = e^{3f_{xx}}? \)

\( \sin(x - ct), \cos(x + ct), e^{x-t}, e^x - e^{-x}, e^x \cos ct \).

42. Suppose \( \partial f/\partial t = \partial f/\partial x \). Show that \( \partial^2 f/\partial t^2 = \partial^2 f/\partial x^2 \).
43 The proof of $f_{xy} = f_{yx}$ studies $f(x, y)$ in a small rectangle. The top-bottom difference is $g(x) = f(x, B) - f(x, A)$. The difference at the corners 1, 2, 3, 4 is:

$$Q = [f_4 - f_3] - [f_2 - f_1] = g(b) - g(a) \quad \text{(definition of g)}$$

$$= (b - a)g'(c) \quad \text{(Mean Value Theorem)}$$

$$= (b - a)[f_4(c, B) - f_4(c, A)] \quad \text{(compute g(x)}$$

$$= (b - a)(B - A) f_{xy}(c, C) \quad \text{(MVT again)}$$

(a) The right-left difference is $h(y) = f(b, y) - f(a, y)$. The same $Q$ is $h(B) - h(A)$. Change the steps to reach $Q = (B - A)h'(c^*, C^*)$.

(b) The two forms of $Q$ make $f_{xy}$ at $(c, C)$ equal to $f_{yx}$ at $(c^*, C^*)$. Shrink the rectangle toward $(a, A)$. What assumption yields $f_{xy}$ at $(c, C)$ equal to $f_{yx}$ at that typical point?

44 Find $\frac{\partial f}{\partial x}$ and $\frac{\partial f}{\partial y}$ where they exist, based on equations (1) and (2).

(a) $f = xy$  \hspace{1cm} (b) $f = x^2 + y^2$ if $x \neq 0, f = 0$ if $x = 0$

Questions 45–52 are about limits in two dimensions.

45 Complete these four correct definitions of limit: 1 The points $(x_n, y_n)$ approach the point $(a, b)$ if $x_n$ converges to $a$ and ________. 2 For any circle around $(a, b)$, the points $(x_n, y_n)$ eventually go ________ the circle and stay ________.

46 Find $(x_2, y_3)$ and $(x_4, y_4)$ and the limit $(a, b)$ if it exists. Start from $(x_0, y_0) = (1, 0)$.

(a) $(x_n, y_n) = \left(\frac{1}{n^2}, \frac{1}{n^2}\right)$  \hspace{1cm} (b) $(x_n, y_n) = \left(\frac{1}{n}, \frac{1}{n}\right)$

47 $(\text{Limit of } f(x, y))$ 1 Informal definition: the numbers $f(x_n, y_n)$ approach $L$ when the points $(x_n, y_n)$ approach $(a, b)$.

2 Epsilon-delta definition: For each $\varepsilon > 0$ there is a $\delta > 0$ such that $|f(x, y) - L|$ is less than ________ when the distance from $(x, y)$ to $(a, b)$ is ________.

The value of $f$ at $(a, b)$ is not involved.

48 Write down the limit $L$ as $(x, y) \to (a, b)$. At which points $(a, b)$ does $f(x, y)$ have no limit?

(a) $f(x, y) = \sqrt{x^2 + y^2}$  \hspace{1cm} (b) $f(x, y) = x/y$

(c) $f(x, y) = 1/(x + y)$  \hspace{1cm} (d) $f(x, y) = xy/(x^2 + y^2)$

In (d) find the limit at $(0, 0)$ along the line $y = mx$. The limit changes with $m$, so $L$ does not exist at $(0, 0)$. Same for $x/y$.

49 Definition of continuity. $f(x, y)$ is continuous at $(a, b)$ if $f(a, b)$ is defined and $f(x, y)$ approaches the limit ________ as $(x, y)$ approaches $(a, b)$. Construct a function that is not continuous at $(1, 2)$.

50 Show that $x^2y/(x^4 + y^2) \to 0$ along every straight line $y = mx$ to the origin. But traveling down the parabola $y = x^2$, the ratio equals ________.

51 Can you define $f(0, 0)$ so that $f(x, y)$ is continuous at $(0, 0)$?

(a) $f = |x + y| - 1$  \hspace{1cm} (b) $f = (1 + x)^y$  \hspace{1cm} (c) $f = x^1 + y^1$.

52 Which functions approach zero as $(x, y) \to (0, 0)$ and why?

(a) $\frac{xy^2}{x^2 + y^2}$  \hspace{1cm} (b) $\frac{x^2y^2}{x^2 + y^2}$  \hspace{1cm} (c) $\frac{x^my^n}{x^n + y^n}$

13.3 Tangent Planes and Linear Approximations

Over a short range, a smooth curve $y = f(x)$ is almost straight. The curve changes direction, but the tangent line $y - y_0 = f'(x_0)(x - x_0)$ keeps the same slope forever. The tangent line immediately gives the linear approximation to $y = f(x)$:

$$y \approx y_0 + f'(x_0)(x - x_0).$$

What happens with two variables? The function is $z = f(x, y)$, and its graph is a surface. We are at a point on that surface, and we are near-sighted. We don’t see far away. The surface may curve out of sight at the horizon, or it may be a bowl or a saddle. To our myopic vision, the surface looks flat. We believe we are on a plane (not necessarily horizontal), and we want the equation of this tangent plane.
13.3 Tangent Planes and Linear Approximations

**Notation** The basepoint has coordinates \( x_0 \) and \( y_0 \). The height on the surface is \( z_0 = f(x_0, y_0) \). Other letters are possible: the point can be \((a, b)\) with height \( w \). The subscript \(_0\) indicates the value of \( x \) or \( y \) or \( z \) or \( \partial f/\partial x \) or \( \partial f/\partial y \) at the point.

With one variable the tangent line has slope \( df/dx \). With two variables there are two derivatives \( \partial f/\partial x \) and \( \partial f/\partial y \). At the particular point, they are \((\partial f/\partial x)_0\) and \((\partial f/\partial y)_0\). **Those are the slopes of the tangent plane.** Its equation is the key to this chapter:

**EXAMPLE 1** Find the tangent plane to \( z = 14 - x^2 - y^2 \) at \((x_0, y_0, z_0) = (1, 2, 9)\).

**Solution** The derivatives are \( \partial f/\partial x = -2x \) and \( \partial f/\partial y = -2y \). When \( x = 1 \) and \( y = 2 \) those are \((\partial f/\partial x)_0 = -2 \) and \((\partial f/\partial y)_0 = -4 \). The equation of the tangent plane is

\[
z - 9 = -2(x - 1) - 4(y - 2) \quad \text{or} \quad z + 2x + 4y = 19. \tag{2}
\]

This \( z(x, y) \) has derivatives \(-2\) and \(-4\), just like the surface. So the plane is tangent.

The normal vector \( N \) has components \(-2, -4, -1\). **The equation of the normal line is** \((x, y, z) = (1, 2, 9) + t(-2, -4, -1)\). Starting from \((1, 2, 9)\) the line goes out along \( N \)—perpendicular to the plane and the surface.

Figure 13.7 shows more detail about the tangent plane. The dotted lines are the \( x \) and \( y \) tangent lines. They lie in the plane. All tangent lines lie in the tangent plane! These particular lines are tangent to the "partial functions"—where \( y \) is fixed at \( y_0 = 2 \) or \( x \) is fixed at \( x_0 = 1 \). The plane is balancing on the surface and touching at the tangent point.

More is true. In the surface, **every curve through the point is tangent to the plane.** Geometrically, the curve goes up to the point and "kisses" the plane.† The tangent \( T \) to the curve and the normal \( N \) to the surface are perpendicular: \( T \cdot N = 0 \).

†A safer word is "osculate." At saddle points the plane is kissed from both sides.
EXAMPLE 2  Find the tangent plane to the sphere \( z^2 = 14 - x^2 - y^2 \) at (1, 2, 3).

Solution  Instead of \( z = 14 - x^2 - y^2 \) we have \( z = \sqrt{14 - x^2 - y^2} \). At \( x_0 = 1, \ y_0 = 2 \) the height is now \( z_0 = 3 \). The surface is a sphere with radius \( \sqrt{14} \). The only trouble from the square root is its derivatives:

\[
\frac{\partial z}{\partial x} = \frac{\partial}{\partial x} \sqrt{14 - x^2 - y^2} = \frac{-\frac{2x}{2\sqrt{14 - x^2 - y^2}}}{\sqrt{14 - x^2 - y^2}} \quad \text{and} \quad \frac{\partial z}{\partial y} = \frac{-\frac{2y}{2\sqrt{14 - x^2 - y^2}}}{\sqrt{14 - x^2 - y^2}} \tag{3}
\]

At \((1, 2)\) those slopes are \(-\frac{1}{4}\) and \(-\frac{1}{2}\). The equation of the tangent plane is linear:

\[z - 3 = -\frac{1}{4}(x - 1) - \frac{1}{2}(y - 2).\]

I cannot resist improving the equation, by multiplying through by \(3\) and moving all terms to the left side:

\[
tangent \plane \ to \ sphere: \quad 1(x - 1) + 2(y - 2) + 3(z - 3) = 0. \tag{4}
\]

If mathematics is the "science of patterns," equation (4) is a prime candidate for study. The numbers 1, 2, 3 appear twice. The coordinates are \((x_0, y_0, z_0) = (1, 2, 3)\). The normal vector is \(1i + 2j + 3k\). The tangent equation is \(1x + 2y + 3z = 14\). None of this can be an accident, but the square root of \(14 - x^2 - y^2\) made a simple pattern look complicated.

This square root is not necessary. Calculus offers a direct way to find \(dz/dx\)—implicit differentiation. Just differentiate every term as it stands:

\[
x^2 + y^2 + z^2 = 14 \quad \text{leads to} \quad 2x + 2z \frac{\partial z}{\partial x} = 0 \quad \text{and} \quad 2y + 2z \frac{\partial z}{\partial y} = 0. \tag{5}
\]

Canceling the 2's, the derivatives on a sphere are \(-x/z\) and \(-y/z\). Those are the same as in (3). The equation for the tangent plane has an extremely symmetric form:

\[
z - z_0 = -\frac{x_0}{z_0}(x - x_0) - \frac{y_0}{z_0}(y - y_0) \quad \text{or} \quad x_0(x - x_0) + y_0(y - y_0) + z_0(z - z_0) = 0. \tag{6}
\]

Reading off \(\mathbf{N} = x_0\mathbf{i} + y_0\mathbf{j} + z_0\mathbf{k}\) from the last equation, calculus proves something we already knew: The normal vector to a sphere points outward along the radius.

**THE TANGENT PLANE TO** \(F(x, y, z) = c\)

The sphere suggests a question that is important for other surfaces. Suppose the equation is \(F(x, y, z) = c\) instead of \(z = f(x, y)\). Can the partial derivatives and tangent plane be found directly from \(F\)?

The answer is yes. It is not necessary to solve first for \(z\). The derivatives of \(F\),
13.3 Tangent Planes and Linear Approximations

computed at \((x_0, y_0, z_0)\), give a second formula for the tangent plane and normal vector.

\[ \frac{\partial F}{\partial x} (x - x_0) + \frac{\partial F}{\partial y} (y - y_0) + \frac{\partial F}{\partial z} (z - z_0) = 0. \]  

(7)

The normal vector is \( \mathbf{N} = \frac{\partial F}{\partial x} \mathbf{i} + \frac{\partial F}{\partial y} \mathbf{j} + \frac{\partial F}{\partial z} \mathbf{k} \).

Notice how this includes the original case \( z = f(x, y) \). The function \( F \) becomes \( f(x, y) - z \). Its partial derivatives are \( \frac{\partial f}{\partial x} \) and \( -1 \). (The \(-1\) is from the derivative of \(-z\).) Then equation (7) is the same as our original tangent equation (1).

**EXAMPLE 3** The surface \( F = x^2 + y^2 - z^2 = c \) is a hyperboloid. Find its tangent plane.

**Solution** The partial derivatives are \( F_x = 2x, F_y = 2y, F_z = -2z \). Equation (7) is

\[ 2x_0(x - x_0) + 2y_0(y - y_0) - 2z_0(z - z_0) = 0. \]  

(8)

We can cancel the 2's. The normal vector is \( \mathbf{N} = x_0 \mathbf{i} + y_0 \mathbf{j} - z_0 \mathbf{k} \). For \( c > 0 \) this hyperboloid has one sheet (Figure 13.8). For \( c = 0 \) it is a cone and for \( c < 0 \) it breaks into two sheets (Problem 13.1.26).

**DIFFERENTIALS**

Come back to the linear equation \( z - z_0 = \left( \frac{\partial z}{\partial x} \right)_0 (x - x_0) + \left( \frac{\partial z}{\partial y} \right)_0 (y - y_0) \) for the tangent plane. That may be the most important formula in this chapter. Move along the tangent plane instead of the curved surface. Movements in the plane are \( dx \) and \( dy \) and \( dz \)—while \( \Delta x \) and \( \Delta y \) and \( \Delta z \) are movements in the surface. The \( \Delta \)'s are governed by the tangent equation—the \( \Delta \)'s are governed by \( z = f(x, y) \). In Chapter 2 the \( \Delta \)'s were **differentials** along the tangent line:

\[ dy = \left( \frac{dy}{dx} \right) dx \quad \text{(straight line)} \quad \text{and} \quad \Delta y \approx \left( \frac{dy}{dx} \right) \Delta x \quad \text{(on the curve)}. \]  

(9)

Now \( y \) is independent like \( x \). The dependent variable is \( z \). The idea is the same. The distances \( x - x_0 \) and \( y - y_0 \) and \( z - z_0 \) (on the tangent plane) are \( dx \) and \( dy \) and \( dz \). The equation of the plane is

\[ dz = \left( \frac{\partial z}{\partial x} \right)_0 dx + \left( \frac{\partial z}{\partial y} \right)_0 dy \quad \text{or} \quad df = f_x dx + f_y dy. \]  

(10)

This is the **total differential**. All letters \( dz \) and \( df \) and \( dw \) can be used, but \( \partial z \) and \( \tilde{f} \) are not used. Differentials suggest small movements in \( x \) and \( y \); then \( dz \) is the resulting movement in \( z \). On the tangent plane, equation (10) holds exactly.

A “centering transform” has put \( x_0, y_0, z_0 \) at the center of coordinates. Then the “zoom transform” stretches the surface into its tangent plane.

**EXAMPLE 4** The area of a triangle is \( A = \frac{1}{2}ab \sin \theta \). Find the total differential \( dA \).

**Solution** The base has length \( b \) and the sloping side has length \( a \). The angle between them is \( \theta \). You may prefer \( A = \frac{1}{2}bh \), where \( h \) is the perpendicular height \( a \sin \theta \). Either way we need the partial derivatives. If \( A = \frac{1}{2}ab \sin \theta \), then

\[ \frac{\partial A}{\partial a} = \frac{1}{2} b \sin \theta \quad \frac{\partial A}{\partial b} = \frac{1}{2} a \sin \theta \quad \frac{\partial A}{\partial \theta} = \frac{1}{2} ab \cos \theta. \]  

(11)
These lead immediately to the total differential \(dA\) (like a product rule):
\[
dA = \left(\frac{\partial A}{\partial a}\right) da + \left(\frac{\partial A}{\partial b}\right) db + \left(\frac{\partial A}{\partial \theta}\right) d\theta = \frac{1}{2}b \sin \theta \, da + \frac{1}{2}a \sin \theta \, db + \frac{1}{2}ab \cos \theta \, d\theta.
\]

**EXAMPLE 5** The volume of a cylinder is \(V = \pi r^2 h\). Decide whether \(V\) is more sensitive to a change from \(r = 1.0\) to \(r = 1.1\) or from \(h = 1.0\) to \(h = 1.1\).

Solution The partial derivatives are \(\partial V/\partial r = 2\pi rh\) and \(\partial V/\partial h = \pi r^2\). They measure the sensitivity to change. Physically, they are the side area and base area of the cylinder. The volume differential \(dV\) comes from a shell around the side plus a layer on top:
\[
dV = \text{shell} + \text{layer} = 2\pi rh \, dr + \pi r^2 dh.
\]

Starting from \(r = h = 1\), that differential is \(dV = 2\pi dh + \pi dr\). With \(dr = dh = .1\), the shell volume is \(.2\pi\) and the layer volume is only \(.1\pi\). So \(V\) is sensitive to \(dr\).

For a short cylinder like a penny, the layer has greater volume. \(V\) is more sensitive to \(dh\). In our case \(V = \pi r^2 h\) increases from \(\pi(1)^3\) to \(\pi(1.1)^3\). Compare \(\Delta V\) to \(dV\):
\[
\Delta V = \pi(1.1)^3 - \pi(1)^3 = .331\pi \quad \text{and} \quad dV = 2\pi(1.1) + \pi(1) = .300\pi.
\]

The difference is \(\Delta V - dV = .031\pi\). The shell and layer missed a small volume in Figure 13.9, just above the shell and around the layer. The mistake is of order \((dr)^2 + (dh)^2\). For \(V = \pi r^2 h\), the differential \(dV = 2\pi rh \, dr + \pi r^2 dh\) is a linear approximation to the true change \(\Delta V\). We now explain that properly.

**LINEAR APPROXIMATION**

Tangents lead immediately to linear approximations. That is true of tangent planes as it was of tangent lines. The plane stays close to the surface, as the line stayed close to the curve. Linear functions are simpler than \(f(x)\) or \(f(x, y)\) or \(F(x, y, z)\). All we need are first derivatives at the point. Then the approximation is good near the point.

This key idea of calculus is already present in differentials. On the plane, \(df\) equals \(f_d dx + f_y dy\). On the curved surface that is a linear approximation to \(\Delta f\):

\[
13C \quad \text{The linear approximation to } f(x, y) \text{ near the point } (x_0, y_0) \text{ is:}
\]
\[
f(x, y) \approx f(x_0, y_0) + \left(\frac{\partial f}{\partial x}\right)_0 (x-x_0) + \left(\frac{\partial f}{\partial y}\right)_0 (y-y_0).
\]

In other words \(\Delta f \approx f_d \Delta x + f_y \Delta y\), as proved in Problem 24. The right side of (13) is a linear function \(f_L(x, y)\). At \((x_0, y_0)\), the functions \(f\) and \(f_L\) have the same slopes. Then \(f(x, y)\) curves away from \(f_L\) with an error of \(\text{"second order"}
\]
\[
|f(x, y) - f_L(x, y)| \leq M[(x-x_0)^2 + (y-y_0)^2].
\]

This assumes that \(f_{xx}, f_{xy}, \text{ and } f_{yy}\) are continuous and bounded by \(M\) along the line from \((x_0, y_0)\) to \((x, y)\). Example 3 of Section 13.5 shows that \(|f_d| \leq 2M\) along that line. A factor \(\frac{1}{2}\) comes from equation 3.8.12, for the error \(f - f_L\) with one variable.

For the volume of a cylinder, \(r\) and \(h\) went from 1.0 to 1.1. The second derivatives of \(V = \pi r^2 h\) are \(V_{rr} = 2\pi h\) and \(V_{rh} = 2\pi r\) and \(V_{hh} = 0\). They are below \(M = 2.2\pi\). Then \(14\) gives the error bound \(2.2\pi(1^2 + 1^2) = .044\pi\), not far above the actual error \(.031\pi\). The main point is that the error in linear approximation comes from the quadratic terms—those are the first terms to be ignored by \(f_L\).
13.3 Tangent Planes and Linear Approximations

EXAMPLE 6  Find a linear approximation to the distance function \( r = \sqrt{x^2 + y^2} \).

Solution  The partial derivatives are \( x/r \) and \( y/r \). Then \( \Delta r \approx (x/r)\Delta x + (y/r)\Delta y \).

For \((x, y, r)\) near \((1, 2, \sqrt{5})\): \( \sqrt{x^2 + y^2} \approx \sqrt{1^2 + 2^2} + (x - 1)/\sqrt{5} + 2(y - 2)/\sqrt{5} \).

If \( y \) is fixed at 2, this is a one-variable approximation to \( \sqrt{x^2 + 2^2} \). If \( x \) is fixed at 1, it is a linear approximation in \( y \). Moving both variables, you might think \( dr \) would involve \( dx \) and \( dy \) in a square root. It doesn't. Distance involves \( x \) and \( y \) in a square root, but change of distance is linear in \( \Delta x \) and \( \Delta y \)—to a first approximation.

There is a rough point at \( x = 0, y = 0 \). Any movement from \((0, 0, 0)\) gives \( \Delta r = \sqrt{(\Delta x)^2 + (\Delta y)^2} \). The square root has returned. The reason is that the partial derivatives \( x/r \) and \( y/r \) are not continuous at \((0, 0)\). The cone has a sharp point with no tangent plane. Linear approximation breaks down.

The next example shows how to approximate \( \Delta z \) from \( \Delta x \) and \( \Delta y \), when the equation is \( F(x, y, z) = c \). We use the implicit derivatives in (7) instead of the explicit derivatives in (1). The idea is the same: Look at the tangent equation as a way to find \( \Delta z \), instead of an equation for \( z \). Here is Example 6 with new letters.

EXAMPLE 7  From \( F = -x^2 - y^2 + z^2 = 0 \) find a linear approximation to \( \Delta z \).

Solution  (implicit derivatives) Use the derivatives of \( F: -2x\Delta x - 2y\Delta y + 2z\Delta z = 0 \). Then solve for \( \Delta z \), which gives \( \Delta z \approx (x/z)\Delta x + (y/z)\Delta y \)—the same as Example 6.

EXAMPLE 8  How does the equilibrium price change when the supply curve changes?

The equilibrium price is at the intersection of the supply and demand curves (supply = demand). As the price \( p \) rises, the demand \( q \) drops (the slope is \(-0.2\)).

\[
\text{demand line } DD: \quad p = -0.2q + 40. \tag{15}
\]

The supply (also \( q \)) goes up with the price. The slope \( s \) is positive (here \( s = 0.4 \)).

\[
\text{supply line } SS: \quad p = sq + t = 0.4q + 10.
\]

Those lines are in Figure 13.10. They meet at the equilibrium price \( P = 30 \). The quantity \( Q = 50 \) is available at \( P \) (on \( SS \)) and demanded at \( P \) (on \( DD \)). So it is sold.

Where do partial derivatives come in? The reality is that those lines \( DD \) and \( SS \) are not fixed for all time. Technology changes, and competition changes, and the value of money changes. Therefore the lines move. Therefore the crossing point \((Q, P)\) also moves. Please recognize that derivatives are hiding in those sentences.
Main point: The equilibrium price $P$ is a function of $s$ and $t$. Reducing $s$ by better technology lowers the supply line to $p = .3q + 10$. The demand line has not changed. The customer is as eager or stingy as ever. But the price $P$ and quantity $Q$ are different. The new equilibrium is at $Q = 60$ and $P = \$28$, where the new line $XX$ crosses $DD$.

If the technology is expensive, the supplier will raise $t$ when reducing $s$. Line $YY$ is $p = .3q + 20$. That gives a higher equilibrium $P = \$32$ at a lower quantity $Q = 40$—the demand was too weak for the technology.

Calculus question: Find $\partial P/\partial s$ and $\partial P/\partial t$. The difficulty is that $P$ is not given as a function of $s$ and $t$. So take implicit derivatives of the supply = demand equations:

\[
\text{supply = demand: } P = -.2Q + 40 = sQ + t
\]

$s$ derivative:

\[
P_s = -.2Q_s = sQ_s + Q \quad \text{(note } t_s = 0)\]

$t$ derivative:

\[
P_t = -.2Q_t = sQ_t + 1 \quad \text{(note } t_t = 1)\]

Now substitute $s = .4$, $t = 10$, $P = 30$, $Q = 50$. That is the starting point, around which we are finding a linear approximation. The last two equations give $P_s = 50/3$ and $P_t = 1/3$ (Problem 25). The linear approximation is

\[
P = 30 + 50(s - .4)/3 + (t - 10)/3.
\]

Comment: This example turned out to be subtle (so is economics). I hesitated before including it. The equations are linear and their derivatives are easy, but something in the problem is hard—there is no explicit formula for $P$. The function $P(s, t)$ is not known. Instead of a point on a surface, we are following the intersection of two lines. The solution changes as the equation changes. The derivative of the solution comes from the derivative of the equation.

Summary: The foundation of this section is equation (1) for the tangent plane. Everything builds on that—total differential, linear approximation, sensitivity to small change. Later sections go on to the chain rule and “directional derivatives” and “gradients.” The central idea of differential calculus is $\Delta f \approx f_s \Delta x + f_t \Delta y$.

Newton’s Method for Two Equations

Linear approximation is used to solve equations. To find out where a function is zero, look first to see where its approximation is zero. To find out where a graph crosses the $xy$ plane, look to see where its tangent plane crosses.

Remember Newton’s method for $f(x) = 0$. The current guess is $x_n$. Around that point, $f(x)$ is close to $f(x_n) + (x - x_n)f'(x_n)$. This is zero at the next guess $x_{n+1} = x_n - f(x_n)/f'(x_n)$. That is where the tangent line crosses the $x$ axis.

With two variables the idea is the same—but two unknowns $x$ and $y$ require two equations. We solve $g(x, y) = 0$ and $h(x, y) = 0$. Both functions have linear approximations that start from the current point $(x_n, y_n)$—where derivatives are computed:

\[
g(x, y) \approx g(x_n, y_n) + (\partial g/\partial x)(x - x_n) + (\partial g/\partial y)(y - y_n)
\]

\[
h(x, y) \approx h(x_n, y_n) + (\partial h/\partial x)(x - x_n) + (\partial h/\partial y)(y - y_n).
\]

The natural idea is to set these approximations to zero. That gives linear equations for $x - x_n$ and $y - y_n$. Those are the steps $\Delta x$ and $\Delta y$ that take us to the next guess.
in Newton's method:

\[ \begin{align*}
13D & \quad \text{Newton's method to solve } g(x, y) = 0 \text{ and } h(x, y) = 0 \text{ has linear equations for the steps } \Delta x \text{ and } \Delta y \text{ that go from } (x_n, y_n) \text{ to } (x_{n+1}, y_{n+1}): \\
& \quad \left( \frac{\partial g}{\partial x} \right) \Delta x + \left( \frac{\partial g}{\partial y} \right) \Delta y = -g(x_n, y_n) \quad \text{and} \quad \left( \frac{\partial h}{\partial x} \right) \Delta x + \left( \frac{\partial h}{\partial y} \right) \Delta y = -h(x_n, y_n). \quad (19)
\end{align*} \]

**EXAMPLE 9** \[ g = x^3 - y = 0 \text{ and } h = y^3 - x = 0 \] have 3 solutions \((1, 1), (0, 0), (-1, -1)\). I will start at different points \((x_0, y_0)\). The next guess is \[ x_1 = x_0 + \Delta x, \quad y_1 = y_0 + \Delta y. \]

It is of extreme interest to know which solution Newton's method will choose—if it converges at all. I made three small experiments.

1. Suppose \((x_0, y_0) = (2, 1)\). At that point \[ g = 2^3 - 1 = 7 \text{ and } h = 1^3 - 2 = -1. \] The derivatives are \[ g_x = 3x^2 = 12, \quad g_y = -1, \quad h_x = -1, \quad h_y = 3y^2 = 3. \] The steps \(\Delta x\) and \(\Delta y\) come from solving (19):

\[ \begin{align*}
12\Delta x - \Delta y &= -7 \quad \Rightarrow \quad \Delta x = -4/7 \quad \Rightarrow \quad x_1 = x_0 + \Delta x = 10/7 \\
-\Delta x + 3\Delta y &= +1 \quad \Rightarrow \quad \Delta y = +1/7 \quad \Rightarrow \quad y_1 = y_0 + \Delta y = 8/7.
\end{align*} \]

This new point \((10/7, 8/7)\) is closer to the solution at \((1, 1)\). The next point is \((1.1, 1.05)\) and convergence is clear. Soon convergence is fast.

2. Start at \((x_0, y_0) = (1/2, 0)\). There we find \[ g = 1/8 \text{ and } h = -1/2. \]

\[ \begin{align*}
(3/4)\Delta x - \Delta y &= -1/8 \quad \Rightarrow \quad \Delta x = -1/2 \quad \Rightarrow \quad x_1 = x_0 + \Delta x = 0 \\
-\Delta x + 0\Delta y &= +1/2 \quad \Rightarrow \quad \Delta y = +1/4 \quad \Rightarrow \quad y_1 = y_0 + \Delta y = -1/4.
\end{align*} \]

Newton has jumped from \((1/2, 0)\) on the \(x\) axis to \((0, -1/2)\) on the \(y\) axis. The next step goes to \((1/32, 0)\), back on the \(x\) axis. We are in the "basin of attraction" of \((0, 0)\).

3. Now start further out the axis at \((1, 0)\), where \[ g = 1 \text{ and } h = -1. \]

\[ \begin{align*}
3\Delta x - \Delta y &= -1 \quad \Rightarrow \quad \Delta x = -1 \quad \Rightarrow \quad x_1 = x_0 + \Delta x = 0 \\
-\Delta x + 0\Delta y &= +1 \quad \Rightarrow \quad \Delta y = -2 \quad \Rightarrow \quad y_1 = y_0 + \Delta y = -2.
\end{align*} \]

Newton moves from \((1, 0)\) to \((0, -2)\) to \((16, 0)\). Convergence breaks down—the method blows up. This danger is ever-present, when we start far from a solution.

Please recognize that even a small computer will uncover amazing patterns. It can start from hundreds of points \((x_0, y_0)\), and follow Newton's method. Each solution has a basin of attraction, containing all \((x_0, y_0)\) leading to that solution. There is also a basin leading to infinity. The basins in Figure 13.11 are completely mixed together—a color figure shows them as fractals. The most extreme behavior is on the borderline between basins, when Newton can't decide which way to go. Frequently we see chaos.

Chaos is irregular movement that follows a definite rule. Newton's method determines an iteration from each point \((x_n, y_n)\) to the next. In scientific problems it normally converges to the solution we want. (We start close enough.) But the computer makes it possible to study iterations from faraway points. This has created a new part of mathematics—so new that any experiments you do are likely to be original.
Section 3.7 found chaos when trying to solve $x^2 + 1 = 0$. But don’t think Newton’s method is a failure. On the contrary, it is the best method to solve nonlinear equations. The error is squared as the algorithm converges, because linear approximations have errors of order $(\Delta x)^2 + (\Delta y)^2$. Each step doubles the number of correct digits, near the solution. The example shows why it is important to be near.

Fig. 13.11 The basins of attraction to $(1, 1)$, $(0, 0)$, $(-1, -1)$, and infinity.

13.3 EXERCISES

Read-through questions

The tangent line to $y = f(x)$ is $y - y_0 = \frac{\partial f(x)}{\partial x} (x - x_0)$. The tangent plane to $w = f(x, y)$ is $w - w_0 = \frac{\partial f(x, y)}{\partial x} (x - x_0) + \frac{\partial f(x, y)}{\partial y} (y - y_0)$. The normal vector is $N = \nabla f$. For $w = x^3 + y^3$ the tangent equation at $(1, 1, 2)$ is $2$. The normal vector is $N = \nabla f$. For a sphere, the direction of $N$ is $\hat{u}$. The surface given implicitly by $F(x, y, z) = c$ has tangent equation $\nabla F \cdot (x - x_0) + \nabla F \cdot (y - y_0) + \nabla F \cdot (z - z_0) = 0$. For $xyz = 6$ at $(1, 2, 3)$ the tangent plane is $2$. On that plane the differentials satisfy $I dx + J dy + K dz = 0$. The differential of $z = f(x, y)$ is $dz = \frac{\partial f(x, y)}{\partial x} dx + \frac{\partial f(x, y)}{\partial y} dy$. This holds exactly on the tangent plane, while $\Delta z \approx \nabla f \cdot (x - x_0)$ holds approximately on the tangent plane. The height $z = 3x + 7y$ is more sensitive to a change in $\hat{u}$ than in $x$, because the partial derivative $\frac{\partial z}{\partial y}$ is larger than $\frac{\partial z}{\partial x}$.

The linear approximation to $f(x, y)$ is $f(x_0, y_0) + \frac{\partial f(x, y)}{\partial x} (x - x_0) + \frac{\partial f(x, y)}{\partial y} (y - y_0)$. This is the same as $\Delta f \approx \nabla f \cdot \Delta x + \frac{\partial f(x, y)}{\partial y} \Delta y$. The error is of order $\Delta x^2 + \Delta y^2$. For $f = \sin xy$ the linear approximation around $(0, 0)$ is $f(x, y) = xy$. We are moving along the $y$ instead of the $x$. When the equation is given as $F(x, y, z) = c$, the linear approximation is $\nabla F \cdot (x - x_0) + \nabla F \cdot (y - y_0) + \nabla F \cdot (z - z_0) = 0$.

Newton’s method solves $g(x, y) = 0$ and $h(x, y) = 0$ by a $\frac{\partial g(x, y)}{\partial x}$ approximation. Starting from $x_0, y_0$, the equations are replaced by $C$ and $\frac{\partial h(x, y)}{\partial y}$. The steps $\Delta x$ and $\Delta y$ go to the next point $E$. Each solution has a basin of $E$. Those basins are likely to be $G$.

In 1–8 find the tangent plane and the normal vector at $P$.

1. $z = \sqrt{x^2 + y^2}$, $P = (0, 1, 1)$
2. $x + y + z = 17$, $P = (3, 4, 10)$
3. $z = x/y$, $P = (6, 3, 2)$
4. $z = e^{x^2}$, $P = (0, 0, 1)$
5. $x^2 + y^2 + z^2 = 6$, $P = (1, 2, 1)$
6. $x^2 + y^2 + z^2 = 7$, $P = (1, 2, 1)$
7. $z = x^2$, $P = (1, 1, 1)$
8. $V = \pi r^2 h$, $P = (2, 2, 8\pi)$.

9. Show that the tangent plane to $z^2 - x^2 - y^2 = 0$ goes through the origin and makes a $45^\circ$ angle with the $z$ axis.

10. The planes $z = x + 4y$ and $z = 2x + 3y$ meet at $(1, 1, 5)$. The whole line of intersection is $(x, y, z) = (1, 1, 5) + t \vec{v}$. Find $\vec{v} = N_1 \times N_2$.

11. If $z = 3x - 2y$ find $dz$ from $dx$ and $dy$. If $z = x^3/y^2$ find $dz$ from $dx$ and $dy$ at $x_0 = 1, y_0 = 1$. If $x$ moves to 1.02 and $y$ moves to 1.03, find the approximate $dz$ and exact $\Delta z$ for both functions. The first surface is the $\boxed{\text{to the second surface.}}$
12 The surfaces \( z = x^2 + 4y \) and \( z = 2x + 3y^3 \) meet at \((1, 1, 5)\). Find the normals \( N_1 \) and \( N_2 \) and also \( v = N_1 \times N_2 \). The line in this direction \( v \) is tangent to what curve?

13 The normal \( N \) to the surface \( F(x, y, z) = 0 \) has components \( F_x, F_y, F_z \). The normal line has \( x = x_0 + F_x \), \( y = y_0 + F_y \), \( z = z_0 + F_z \). For the surface \( xyz - 3 = 0 \), find the tangent plane and normal line at \((2, 3, 1)\).

14 For the surface \( x^2 + y^2 - z = 0 \), the normal line at \((1, 2, 3)\) has \( x = __________, y = __________, z = __________. \)

15 For the sphere \( x^2 + y^2 + z^2 = 9 \), find the equation of the tangent plane through \((2, 1, 2)\). Also find the equation of the normal line and show that it goes through \((0, 0, 0)\).

16 If the normal line at every point on \( F(x, y, z) = 0 \) goes through \((0, 0, 0)\), show that \( F_x = cx \), \( F_y = cy \), \( F_z = cz \). The normal line of any point on \( F(x, y, z) = 0 \) has \( x = __________, y = __________, z = __________. \)

17 For \( w = xy \) near \((x_0, y_0)\), the linear approximation is \( dw = \). This looks like the _____ rule for derivatives. The difference between \( \Delta w = xy - x_0y_0 \) and this approximation is __________.

18 If \( f = xyz \) (3 independent variables) what is \( df? \)

19 You invest \( P = $4000 \) at \( R = 8\% \) to make \( I = $320 \) per year. If the numbers change by \( dp \) and \( dR \) what is \( dI? \) If the rate drops by \( dR = .002 \) (to 7.8%) what change \( dP \) keeps \( dI = 0? \) Find the exact interest \( I \) after those changes in \( R \) and \( P \).

20 Resistances \( R_1 \) and \( R_2 \) have parallel resistance \( R \), where \( 1/R = 1/R_1 + 1/R_2 \). Is \( R \) more sensitive to \( \Delta R_1 \), or \( \Delta R_2 \)? If \( R_1 = 1 \) and \( R_2 = 2? \)

21 (a) If your batting average is \( A = (25 \text{ hits})/(100 \text{ at bats}) = .250 \), compute the increase \((to 26/101)\) with a hit and the decrease \((to 25/101)\) with an out. (b) If \( A = x/y \) then \( dA = \frac{ydx - xdy}{y^2} \) gives \( dA = (1 - A)dy/y \). A hit \((dx = dy = 1)\) gives \( dA = -A \), y. An out \((dy = 1)\) gives \( dA = -A \), y. So at \( A = .250 \) a hit has \( times \) the effect of an out.

22 (a) A hit and 3 outs \((dx = 2, dy = 5)\) will raise your average \((dA > 0)\) provided \( A \) is less than \( \). (b) A player batting \( A = .500 \) with \( y = 400 \) at bats needs \( dx = \) hits to raise his average to \( .505 \).

23 If \( x \) and \( y \) change by \( \Delta x \) and \( \Delta y \), find the approximate change \( \Delta A \) in the angle \( \theta = \tan^{-1}(y/x) \).

24 The Fundamental Lemma behind equation (13) writes \( \Delta f = a \Delta x + b \Delta y \). The Lemma says that \( a \rightarrow f_x(x_0, y_0) \) and \( b \rightarrow f_y(x_0, y_0) \) when \( \Delta x \rightarrow 0 \) and \( \Delta y \rightarrow 0 \). The proof takes \( \Delta x \) first and then \( \Delta y \): (1) \( f(x_0 + \Delta x, y_0) = f(x_0, y_0) = \Delta x f_x(c, y_0) \) where \( c \) is between __________ and __________ (by which theorem? (2) \( f(x_0 + \Delta x, y_0 + \Delta y) = f(x_0, y_0) = \Delta y f_y(x_0 + \Delta x, C) \) where \( C \) is between __________ and __________.

25 If the supplier reduces \( s \), Figure 13.10 shows that \( P \) decreases and \( Q \). (a) Find \( P_1 = 50/3 \) and \( P_2 = 1/3 \) in the economics equation (17) by solving the equations above it for \( Q_2 \) and \( Q_1 \). (b) What is the linear approximation to \( Q \) around \( s = 4 \), \( t = 10 \), \( P = 30 \), \( Q = 50 \)?

26 Solve the equations \( P = -2Q + 40 \) and \( P = sQ + t \) for \( P \) and \( Q \). Then find \( \partial P/\partial s \) and \( \partial P/\partial t \) explicitly. At the same \( s, t, P, Q \) check \( 50/3 \) and \( 1/3 \).

27 If the supply = demand equation (16) changes to \( P = xQ + t = -Q + 50 \), find \( P_1, P_2 \) at \( s = 1, t = 10 \).

28 To find out how the roots of \( x^2 + bx + c = 0 \) vary with \( b \), take partial derivatives of the equation with respect to _______. Compare \( \partial x/\partial b \) with \( \partial x/\partial c \) to show that a root at \( x = 2 \) is more sensitive to \( b \).

29 Find the tangent planes to \( z = xy \) and \( z = x^2 - y^2 \) at \( x = 2, y = 1 \). Find the Newton point where those planes meet the \( xy \) plane \((z = 0 \text{ in the tangent equations})\).

30 (a) To solve \( g(x, y) = 0 \) and \( h(x, y) = 0 \) is to find the meeting point of three surfaces: \( g(x, y) = 0 \) and \( h(x, y) = 0 \).

(b) Newton finds the meeting point of three planes: the tangent plane to the graph of \( g, \), and ________.

Problems 31–36 go further with Newton's method for \( g = x^3 - y \) and \( h = y^3 - x \). This is Example 9 with solutions \((1, 1), (0, 0), (-1, -1)\).

31 Start from \( x_0 = 1, y_0 = 1 \) and find \( \Delta x \) and \( \Delta y \). Where are \( x_1 \) and \( y_1 \), and what line is Newton's method moving on?

32 Start from \((1, 1)\) and find the next point. This is in the basin of attraction of which solution?

33 Starting from \((a, -a)\) find \( \Delta y \) which is also \( -\Delta x \). Newton goes toward \((0, 0)\). But can you find the sharp point in Figure 13.11 where the lemon meets the spade?

34 Starting from \((a, 0)\) show that Newton's method goes to \((0, -2a^3)\) and find the next point \((x_2, y_2)\). Which numbers \( a \) lead to convergence? Which special number \( a \) leads to a cycle, in which \((x_2, y_2)\) is the same as the starting point \((a, 0)\)?

35 Show that \( x^3 = y, y^3 = x \) has exactly three solutions.

36 Locate a point from which Newton's method diverges.

37 Apply Newton's method to a linear problem: \( g = x + 2y - 5 = 0, h = 3x - 3 = 0 \). From any starting point show that \((x_1, y_1)\) is the exact solution (convergence in one step).
38. The complex equation \((x + iy)^3 = 1\) contains two real equations, \(x^3 - 3xy^2 = 1\) from the real part and \(3x^2y - y^3 = 0\) from the imaginary part. Search by computer for the basins of attraction of the three solutions \((1, 0), (-1/2, \sqrt{3}/2),\) and \((-1/2, -\sqrt{3}/2)\)—which give the cube roots of 1.

39. In Newton's method the new guess comes from \((x_n, y_n)\) by an iteration: \(x_{n+1} = G(x_n, y_n)\) and \(y_{n+1} = H(x_n, y_n)\). What are \(G\) and \(H\) for \(g = x^2 - y = 0, h = x - y = 0?\) First find \(\Delta x\) and \(\Delta y;\) then \(x_n + \Delta x\) gives \(G\) and \(y_n + \Delta y\) gives \(H\).

40. In Problem 39 find the basins of attraction of the solution \((0, 0)\) and \((1, 1)\).

41. The matrix in Newton's method is the Jacobian:
\[
J = \begin{bmatrix}
\frac{\partial g}{\partial x} & \frac{\partial g}{\partial y} \\
\frac{\partial h}{\partial x} & \frac{\partial h}{\partial y}
\end{bmatrix}
\quad \text{and} \quad
J \begin{bmatrix} \Delta x \\ \Delta y \end{bmatrix} = \begin{bmatrix} -B_n \\ -h_n \end{bmatrix}.
\]

Find \(J\) and \(\Delta x\) and \(\Delta y\) for \(g = e^x - 1, h = e^x + x\).

42. Find the Jacobian matrix at \((1, 1)\) when \(g = x^2 + y^2\) and \(h = xy\). This matrix is _______ and Newton's method fails. The graphs of \(g\) and \(h\) have _______ tangent planes.

43. Solve \(g = x^2 - y^2 + 1 = 0\) and \(h = 2xy = 0\) by Newton's method from three starting points: \((0, 2)\) and \((-1, 1)\) and \((2, 0)\). Take ten steps by computer or one by hand. The solution \((0, 1)\) attracts when \(y_0 > 0\). If \(y_0 = 0\) you should find the chaos iteration \(x_{n+1} = \frac{3}{2}(x_n - x_n^{-1})\).

---

### 13.4 Directional Derivatives and Gradients

As \(x\) changes, we know how \(f(x, y)\) changes. The partial derivative \(\frac{\partial f}{\partial x}\) treats \(y\) as constant. Similarly \(\frac{\partial f}{\partial y}\) keeps \(x\) constant, and gives the slope in the \(y\) direction. But east-west and north-south are not the only directions to move. We could go along a 45° line, where \(\Delta x = \Delta y\). In principle, before we draw axes, no direction is preferred. The graph is a surface with slopes in all directions.

On that surface, calculus looks for the rate of change (or the slope). There is a *directional derivative*, whatever the direction. In the 45° case we are inclined to divide \(\Delta f\) by \(\Delta x\), but we would be wrong.

Let me state the problem. We are given \(f(x, y)\) around a point \(P = (x_0, y_0)\). We are also given a direction \(u\) (a unit vector). There must be a natural definition of \(D_u f\)—the *derivative of \(f\) in the direction \(u\)*. To compute this slope at \(P\), we need a formula. Preferably the formula is based on \(\frac{\partial f}{\partial x}\) and \(\frac{\partial f}{\partial y}\), which we already know.

Note that the 45° direction has \(u = \frac{1}{\sqrt{2}} + j\frac{1}{\sqrt{2}}\). The square root of 2 is going to enter the derivative. This shows that dividing \(\Delta f\) by \(\Delta x\) is wrong. We should divide by the step length \(\Delta s\).

**EXAMPLE 1** Stay on the surface \(z = xy\). When \((x, y)\) moves a distance \(\Delta s\) in the 45° direction from \((1, 1)\), what is \(\Delta z/\Delta s\)?

**Solution** The step is \(\Delta s\) times the unit vector \(u\). Starting from \(x = y = 1\) the step ends at \((1 + \Delta s/\sqrt{2}, 1 + \Delta s/\sqrt{2})\). Then \(z = xy\) is

\[
z = (1 + \Delta s/\sqrt{2})^2 = 1 + \sqrt{2}\Delta s + \frac{1}{4}(\Delta s)^2,
\]

which means \(\Delta z = \sqrt{2}\Delta s + \frac{1}{4}(\Delta s)^2\).

The ratio \(\Delta z/\Delta s\) approaches \(\sqrt{2}\) as \(\Delta s \to 0\). That is the slope in the 45° direction.

**DEFINITION** The derivative of \(f\) in the direction \(u\) at the point \(P\) is \(D_u f(P)\):

\[
D_u f(P) = \lim_{\Delta s \to 0} \frac{\Delta f}{\Delta s} = \lim_{\Delta s \to 0} \frac{f(P + \Delta s u) - f(P)}{\Delta s} \quad (1)
\]

The step from \(P = (x_0, y_0)\) has length \(\Delta s\). It takes us to \((x_0 + \Delta s u_1, y_0 + \Delta s u_2)\). We compute the change \(\Delta f\) and divide by \(\Delta s\). But formula (2) below saves time.
13.4 Directional Derivatives and Gradients

The $x$ direction is $u = (1, 0)$. Then $u\Delta s$ is $(\Delta s, 0)$ and we recover $\frac{\partial f}{\partial x}$:

$$\frac{\Delta f}{\Delta s} = \frac{f(x_0 + \Delta s, y_0) - f(x_0, y_0)}{\Delta s}$$

approaches $D_{(1,0)}f = \frac{\partial f}{\partial x}$. 

Similarly $D_{u}f = \frac{\partial f}{\partial y}$, when $u = (0, 1)$ is in the $y$ direction. What is $D_{u}f$ when $u = (0, -1)$? That is the negative $y$ direction, so $D_{u}f = -\frac{\partial f}{\partial y}$.

CALCULATING THE DIRECTIONAL DERIVATIVE

$D_{u}f$ is the slope of the surface $z = f(x, y)$ in the direction $u$. How do you compute it? From $\frac{\partial f}{\partial x}$ and $\frac{\partial f}{\partial y}$, in two special directions, there is a quick way to find $D_{u}f$ in all directions. Remember that $u$ is a unit vector.

13E The directional derivative $D_{u}f$ in the direction $u = (u_1, u_2)$ equals

$$D_{u}f = \frac{\partial f}{\partial x}u_1 + \frac{\partial f}{\partial y}u_2.$$  

(2)

The reasoning goes back to the linear approximation of $\Delta f$:

$$\Delta f \approx \frac{\partial f}{\partial x}\Delta x + \frac{\partial f}{\partial y}\Delta y = \frac{\partial f}{\partial x}u_1\Delta s + \frac{\partial f}{\partial y}u_2\Delta s.$$  

Divide by $\Delta s$ and let $\Delta s$ approach zero. Formula (2) is the limit of $\Delta f/\Delta s$, as the approximation becomes exact. A more careful argument guarantees this limit provided $f_x$ and $f_y$ are continuous at the basepoint $(x_0, y_0)$.

Main point: Slopes in all directions are known from slopes in two directions.

EXAMPLE 1 (repeated) $f = xy$ and $P = (1, 1)$ and $u = (1/\sqrt{2}, 1/\sqrt{2})$. Find $D_{u}f(P)$.

The derivatives $f_x = y$ and $f_y = x$ equal 1 at $P$. The $45^\circ$ derivative is

$$D_{u}f(P) = f_xu_1 + f_yu_2 = 1(1/\sqrt{2}) + 1(1/\sqrt{2}) = \sqrt{2}$$

as before.

EXAMPLE 2 The linear function $f = 3x + y + 1$ has slope $D_{u}f = 3u_1 + u_2$.

The $x$ direction is $u = (1, 0)$, and $D_{u}f = 3$. That is $\frac{\partial f}{\partial x}$. In the $y$ direction $D_{u}f = 1$.

Two other directions are special—along the level lines of $f(x, y)$ and perpendicular:

Level direction: $D_{u}f$ is zero because $f$ is constant
Steepest direction: $D_{u}f$ is as large as possible (with $u_1^2 + u_2^2 = 1$).

To find those directions, look at $D_{u}f = 3u_1 + u_2$. The level direction has $3u_1 + u_2 = 0$. Then $u$ is proportional to $(1, -3)$. Changing $x$ by 1 and $y$ by $-3$ produces no change in $f = 3x + y + 1$.

In the steepest direction $u$ is proportional to $(3, 1)$. Note the partial derivatives $f_x = 3$ and $f_y = 1$. The dot product of $(3, 1)$ and $(1, -3)$ is zero—steepest direction is perpendicular to level direction. To make $(3, 1)$ a unit vector, divide by $\sqrt{10}$.

Steepest climb: $D_{u}f = 3(3/\sqrt{10}) + 1(1/\sqrt{10}) = 10/\sqrt{10} = \sqrt{10}$
Steepest descent: Reverse to $u = (-3/\sqrt{10}, -1/\sqrt{10})$ and $D_{u}f = -\sqrt{10}$.

The contour lines around a mountain follow $D_{u}f = 0$. The creeks are perpendicular. On a plane like $f = 3x + y + 1$, those directions stay the same at all points (Figure 13.12). On a mountain the steepest direction changes as the slopes change.
13 Partial Derivatives

THE GRADIENT VECTOR

Look again at \( f_u u_1 + f_u u_2 \), which is the directional derivative \( D_u f \). This is the dot product of two vectors. One vector is \( u = (u_1, u_2) \), which sets the direction. The other vector is \( (f_x, f_y) \), which comes from the function. This second vector is the gradient.

**DEFINITION**

The **gradient** of \( f(x, y) \) is the vector whose components are \( \frac{\partial f}{\partial x} \) and \( \frac{\partial f}{\partial y} \):

\[
\text{grad} \ f = \nabla f = \frac{\partial f}{\partial x} \mathbf{i} + \frac{\partial f}{\partial y} \mathbf{j}
\]

(add \( \frac{\partial f}{\partial z} \mathbf{k} \) in three dimensions).

The space-saving symbol \( \nabla \) is read as “grad.” In Chapter 15 it becomes “del.”

For the linear function \( 3x + y + 1 \), the gradient is the constant vector \((3, 1)\). It is the way to climb the plane. For the nonlinear function \( x^2 + xy \), the gradient is the non-constant vector \((2x + y, x)\). Notice that grad \( f \) shares the two derivatives in \( \mathbf{N} = (f_x, f_y, -1) \). But the gradient is not the normal vector. \( \mathbf{N} \) is in three dimensions, pointing away from the surface \( z = f(x, y) \). The **gradient vector** is in the \( xy \) plane! The gradient tells which way on the surface is up, but it does that from down in the base.

The level curve is also in the \( xy \) plane, perpendicular to the gradient. The contour map is a projection on the base plane of what the hiker sees on the mountain. The vector grad \( f \) tells the direction of climb, and its length \( |\text{grad} \ f| \) gives the steepness.

The example \( f = 3x + y + 1 \) had \( \text{grad} \ f = (3, 1) \). Its steepest slope was in the direction \( u = (3, 1)/\sqrt{10} \). The maximum slope was \( \sqrt{10} \). That is \( |\text{grad} \ f| = \sqrt{9 + 1} \).

Important point: **The maximum of** \( (\text{grad} \ f) \cdot \mathbf{u} \) **is the length** \( |\text{grad} \ f| \). In nonlinear examples, the gradient and steepest direction and slope will vary. But look at one particular point in Figure 13.13. Near that point, and near any point, the linear picture takes over.

On the graph of \( f \), the special vectors are the level direction \( \mathbf{L} = (f_x, -f_x, 0) \) and the uphill direction \( \mathbf{U} = (f_x, f_y, f_x^2 + f_y^2) \) and the normal \( \mathbf{N} = (f_x, f_y, -1) \). Problem 18 checks that those are perpendicular.
EXAMPLE 3  The gradient of \( f(x, y) = (14 - x^2 - y^2)/3 \) is \( \nabla f = (-2x/3, -2y/3) \).

On the surface, the normal vector is \( N = (-2x/3, -2y/3, -1) \). At the point \((1, 2, 3)\), this perpendicular is \( N = (-2/3, -4/3, -1) \). At the point \((1, 2)\) down in the base, the gradient is \((-2/3, -4/3)\). The length of \( \text{grad} \ f \) is the slope \( \sqrt{20}/3 \).

Probably a hiker does not go straight up. A “grade” of \( \sqrt{20}/3 \) is fairly steep (almost 150%). To estimate the slope in other directions, measure the distance along the path between two contour lines. If \( \Delta f = 1 \) in a distance \( \Delta s = 3 \) the slope is about \( 1/\sqrt{3} \). This calculation is not exact until the limit of \( \Delta f/\Delta s \), which is \( D_u f \).

![Fig. 13.13  N perpendicular to surface and \( \text{grad} \ f \) perpendicular to level line (in the base).](image)

EXAMPLE 4  The gradient of \( f(x, y, z) = xy + yz + xz \) has three components.

The pattern extends from \( f(x, y) \) to \( f(x, y, z) \). The gradient is now the three-dimensional vector \((f_x, f_y, f_z)\). For this function \( \text{grad} \ f \) is \((y + z, x + z, x + y)\). To draw the graph of \( w = f(x, y, z) \) would require a four-dimensional picture, with axes in the \( xywz \) directions.

Notice the dimensions. The graph is a 3-dimensional “surface” in 4-dimensional space. The gradient is down below in the 3-dimensional base. The level sets of \( f \) come from \( xy + yz + xz = c \)—they are 2-dimensional. The gradient is perpendicular to that level set (still down in 3 dimensions). The gradient is not \( N \! \! \! \! \! \! \! \i \)). The normal vector is \((f_x, f_y, f_z, -1)\), perpendicular to the surface up in 4-dimensional space.

EXAMPLE 5  Find \( \text{grad} \ z \) when \( z(x, y) \) is given implicitly: \( F(x, y, z) = x^2 + y^2 - z^2 = 0 \).

In this case we find \( z = \pm \sqrt{x^2 + y^2} \). The derivatives are \( \pm x/\sqrt{x^2 + y^2} \) and \( \pm y/\sqrt{x^2 + y^2} \), which go into \( \text{grad} \ z \). But the point is this: To find that gradient faster, differentiate \( F(x, y, z) \) as it stands. Then divide by \( F_z \):

\[
F_x \, dx + F_y \, dy + F_z \, dz = 0 \quad \text{or} \quad dz = (-F_x \, dx - F_y \, dy)/F_z. \tag{3}
\]

The gradient is \((-F_x/F_z, -F_y/F_z)\). Those derivatives are evaluated at \((x_0, y_0)\). The computation does not need the explicit function \( z = f(x, y) \):

\[
F = x^2 + y^2 - z^2 \Rightarrow F_x = 2x, F_y = 2y, F_z = -2z \Rightarrow \text{grad} \ z = (x/z, y/z).
\]

To go uphill on the cone, move in the direction \((x/z, y/z)\). That gradient direction goes radially outward. The steepness of the cone is the length of the gradient vector:

\[
|\text{grad} \ z| = \sqrt{(x/z)^2 + (y/z)^2} = 1 \quad \text{because} \quad z^2 = x^2 + y^2 \quad \text{on the cone}.
\]
Partial Derivatives

13 Partial Derivatives

Derivatives along curved paths

On a straight path the derivative of \(f\) is \(D_uf = (\text{grad } f) \cdot u\). What is the derivative on a curved path? The path direction \(u\) is the tangent vector \(T\). So replace \(u\) by \(T\), which gives the “direction” of the curve.

The path is given by the position vector \(R(t) = x(t)i + y(t)j\). The velocity is \(v = (dx/dt)i + (dy/dt)j\). The tangent vector is \(T = v/|v|\). Notice the choice—to move at any speed (with \(v\)) or to go at unit speed (with \(T\)). There is the same choice for the derivative of \(f(x, y)\) along this curve:

\[
\begin{align*}
\text{rate of change} & \quad \frac{df}{dt} = (\text{grad } f) \cdot v = \frac{\partial f}{\partial x} \frac{dx}{dt} + \frac{\partial f}{\partial y} \frac{dy}{dt} \\
\text{slope} & \quad \frac{df}{ds} = (\text{grad } f) \cdot T = \frac{\partial f}{\partial x} \frac{dx}{ds} + \frac{\partial f}{\partial y} \frac{dy}{ds}
\end{align*}
\]

Equations (4) and (5) look like chain rules. They are chain rules. The next section extends \(df/dt = (df/dx)(dx/dt)\) to more variables, proving (4) and (5). Here we focus on the meaning: \(df/ds\) is the derivative of \(f\) in the direction \(u = T\) along the curve.

Example 7

Find \(df/dt\) and \(df/ds\) for \(f = r\). The curve is \(x = t^2, y = t\) in Figure 13.14a.

Solution

The velocity along the curve is \(v = 2ti + j\). At the typical point \(t = 1\) it is \(v = 2i + j\). The unit tangent is \(T = v/\sqrt{5}\). The gradient is a unit vector \(1/\sqrt{2}i + j/\sqrt{2}\) pointing outward, when \(f(x, y)\) is the distance \(r\) from the center. The dot product with \(v\) is \(df/dt = 3/\sqrt{5}\). The dot product with \(T\) is \(df/ds = 3/\sqrt{10}\).

When we slow down to speed 1 (with \(T\)), the changes in \(f(x, y)\) slow down too.

Example 8

Find \(df/ds\) for \(f = xy\) along the circular path \(x = \cos t, y = \sin t\).

First take a direct approach. On the circle, \(xy\) equals \((\cos t)(\sin t)\). Its derivative comes from the product rule: \(df/dt = \cos^2t - \sin^2t\). Normally this is different from \(df/ds\), because the time \(t\) need not equal the arc length \(s\). There is a speed factor \(ds/dt\) to divide by—but here the speed is 1. (A circle of length \(s = 2\pi\) is completed at \(t = 2\pi\).) Thus the slope \(df/ds\) along the roller-coaster in Figure 13.14 is \(\cos^2t - \sin^2t\).

**Fig. 13.14** The distance \(f = r\) changes along the curve. The slope of the roller-coaster is \((\text{grad } f) \cdot T\). The distance \(D\) from \((x_0, y_0)\) has \(\text{grad } D = \text{unit vector}\).
The second approach uses the vectors grad \( f \) and \( T \). The gradient of \( f = xy \) is 
\( \nabla f = (y, x) \). The unit tangent vector to the path is 
\( T = (-\sin t, \cos t) \). Their dot product is the same 
\[ \frac{df}{ds} = (\nabla f) \cdot T = -\sin^2 t + \cos^2 t. \]

**Gradients without Coordinates**

This section ends with a little "philosophy." What is the coordinate-free definition of the gradient? Up to now, \( \nabla f = (f_x, f_y) \) depended totally on the choice of \( x \) and \( y \) axes. But the steepness of a surface is independent of the axes. Those are added later, to help us compute.

The steepness \( df/ds \) involves only \( f \) and the direction, nothing else. The gradient should be a "tensor"—its meaning does not depend on the coordinate system. The gradient has different formulas in different systems (\( xy \) or \( rt \) or \( \ldots \)), but the direction and length of \( \nabla f \) are determined by \( df/ds \)—without any axes:

The direction of \( \nabla f \) is the one in which \( df/ds \) is largest.

The length \( |\nabla f| \) is that largest slope.

The key equation is 
\[ \text{change in } f = (\nabla f) \cdot \text{change in position}. \]

That is another way to write 
\[ \Delta f \approx f_x \Delta x + f_y \Delta y. \]

It is the multivariable form—we used two variables—of the basic linear approximation 
\[ \Delta y \approx (dy/dx) \Delta x. \]

**Example 9**

\( D(x, y) = \text{distance from } (x, y) \text{ to } (x_0, y_0) \). Without derivatives prove 
\[ |\nabla D| = 1. \]

The graph of \( D(x, y) \) is a cone with slope 1 and sharp point \((x_0, y_0)\).

**First question** In which direction does the distance \( D(x, y) \) increase fastest?

**Answer** Going directly away from \((x_0, y_0)\). Therefore this is the direction of \( \nabla D \).

**Second question** How quickly does \( D \) increase in that steepest direction?

**Answer** A step of length \( \Delta s \) increases \( D \) by \( \Delta s \). Therefore 
\[ |\nabla D| = \Delta s/\Delta s = 1. \]

**Conclusion** \( \nabla D \) is a unit vector. The derivatives of \( D \) in Problem 48 are 
\[ (x - x_0)^2 + (y - y_0)^2 \] equals \( D^2 \).

**Exercises**

**Read-through questions**

\( D_x f \) gives the rate of change of \( u \) in the direction \( v \). It can be computed from the two derivatives \( e_u \) in the special directions \( d \). In terms of \( u_1, u_2 \) the formula is 
\[ D_x f = \langle u_1, u_2 \rangle \cdot \nabla f. \]

This is a \( f \) product of \( u \) with the vector \( \nabla f \), which is called the \( n \). For the linear function \( f = ax + by \), the gradient is \( \nabla f = \langle a, b \rangle \). The directional derivative is 
\[ D_x f = \langle a, b \rangle \cdot \hat{u}. \]

The gradient \( \nabla f = (f_x, f_y) \) is not a vector in \( n \) dimensions, it is a vector in the \( m \). It is perpendicular to the \( n \) lines. It points in the direction of \( n \) climb. Its magnitude \( |\nabla f| \) is \( p \). For \( f = x^2 + y^2 \) the gradient points \( q \) and the slope in that steepest direction is \( r \).

The gradient of \( f(x, y, z) \) is \( s \). This is different from the gradient on the surface \( F(x, y, z) = 0 \), which is 
\[ -(F_x/F_z)i + (F_y/F_z)j + (F_z/F_z)k. \]

Traveling with velocity \( v \) on a curved path, the rate of change of \( f \) is 
\[ df/dt = \langle a, b, c \rangle \cdot \hat{v}. \]

When the tangent direction is \( t \), the slope of \( f \) is \( df/ds = \langle a, b, c \rangle \cdot \hat{t} \). In a straight direction \( u \), 
\( df/ds \) is the same as \( w \).

Compute \( \nabla f \), then \( D_x f = (\nabla f) \cdot u \), then \( D_x f \) at \( P \).

1. \( f(x, y) = x^2 - y^2 \)
   \[ u = (\sqrt{3}/2, 1/2) \]
   \[ P = (1, 0) \]
2. \( f(x, y) = 3x + 4y + 7 \)
   \[ u = (3/5, 4/5) \]
   \[ P = (0, \pi/2) \]
3. \( f(x, y) = e^x \cos y \)
   \[ u = (0, 1) \]
   \[ P = (0, \pi/2) \]
4. \( f(x, y) = y^{10} \)
   \[ u = (0, -1) \]
   \[ P = (1, -1) \]
5 \( f(x, y) = \text{distance to (0, 3)} \) \( u = (1, 0) \) \( P = (1, 1) \)

Find \( \nabla f = (f_x, f_y, f_z) \) for the functions 6–8 from physics.

6 \( \sqrt{x^2 + y^2} \) (point source at the origin)

7 \( \ln(x^2 + y^2) \) (line source along z axis)

8 \( \sqrt{(x-1)^2 + y^2 + z^2} - 1 - \sqrt{(x+1)^2 + y^2 + z^2} \) (dipole)

9 \( f = 3x^2 + 2y^2 \) find the steepest direction and the level direction at (1, 2). Compute \( D_n f \) in those directions.

10 Example 2 claimed that \( f = 3x + y + 1 \) has steepest slope \( \text{Maximize } D_u f = 3u_1 + u_2 = 3u_1 + \sqrt{1 - u_1^2} \).

11 True or false, when \( f(x, y) \) is any smooth function:

(a) There is a direction \( u \) at \( P \) in which \( D_u f = 0 \).

(b) There is a direction \( u \) in which \( D_u f = \text{grad } f \).

(c) There is a direction \( u \) in which \( D_u f = 1 \).

(d) The gradient of \( f \) is \( g \) grad \( f + f \) grad \( g \).

12 What is the gradient of \( f(x, y) \)? (One component only.) What are the two possible directions \( u \) and the derivatives \( D_u f \)?

What is the normal vector \( N \) to the curve \( y = f(x) \)? (Two components.)

In 13–16 find the direction \( u \) in which \( f \) increases fastest at \( P = (1, 2) \). How fast?

13 \( f(x, y) = ax + by \)

14 \( f(x, y) = \text{smaller of } 2x \text{ and } y \)

15 \( f(x, y) = e^{x-y} \)

16 \( f(x, y) = \sqrt{5 - x^2 - y^2} \) (careful)

17 (Looking ahead) At a point where \( f(x, y) \) is a maximum, what is grad \( f \)? Describe the level curve containing the maximum point (\( x, y \)).

18 (a) Check by dot products that the normal and uphill and level directions on the graph are perpendicular: \( N = (f_x, f_y, -1), U = (f_x, f_y, f_x^2 + f_y^2), L = (f_x, -f_y, 0) \).

(b) \( N \) is \( \boxed{\text{to the tangent plane.}} \) \( U \) and \( L \) are \( \boxed{\text{to the tangent plane}} \).

(c) The gradient is the xy projection of \( \boxed{\text{grad } f} \) and also of \( \boxed{\text{grad } f} \). The projection of \( L \) points along the \( \boxed{\text{grad } f} \).

19 Compute the \( N, U, L \) vectors for \( f = 1 - x + y \) and draw them at a point on the flat surface.

20 Compute \( N, U, L \) for \( x^2 + y^2 - z^2 = 0 \) and draw them at a typical point on the cone.

With gravity in the negative z direction, in what direction \( -U \) will water flow down the roofs 21–24?

21 \( z = 2x \) (flat roof)

22 \( z = 4x - 3y \) (flat roof)

23 \( z = \sqrt{1 - x^2 - y^2} \) (sphere)

24 \( z = -\sqrt{x^2 + y^2} \) (cone)

25 Choose two functions \( f(x, y) \) that depend only on \( x + 2y \). Their gradients at \( (1, 1) \) are in the direction \( \boxed{\text{.}} \). Their level curves are \( \boxed{\text{.}} \).

26 The level curve of \( f = y/x \) through \( (1, 1) \) is \( \boxed{\text{.}} \). The direction of the gradient must be \( \boxed{\text{.}} \). Check grad \( f \).

27 Grad \( f \) is perpendicular to \( 2i + j \) with length 1, and grad \( g \) is parallel to \( 2i + j \) with length 5. Find grad \( f, g, f, g \).

28 True or false:

(a) If we know \( \text{grad } f \), we know \( f \).

(b) The line \( x = y = z \) is perpendicular to the plane \( z = x + y \).

(c) The gradient of \( z = x + y \) lies along that line.

29 Write down the level direction \( u \) for \( z = \tan^{-1}(y/x) \) at the point \( (3, 4) \). Then compute \( \theta \) and check \( D_u \theta = 0 \).

30 On a circle around the origin, distance is \( \Delta s = \pi \Delta \theta \). Then \( d\theta/ds = 1/r \). Verify by computing \( \theta \) and \( \Delta \theta \) and \( (\text{grad } \theta) \cdot T \).

31 At the point \( (2, 1, 6) \) on the mountain \( z = 9 - x - y^2 \), which way is up? On the roof \( z = x + 2y + 2 \), which way is down? The roof is \( \boxed{\text{.}} \) to the mountain.

32 Around the point \( (1, -2) \) the temperature \( T = e^{-x^2-y^2} \) has \( \Delta T \approx \boxed{\Delta x + \boxed{\Delta y}} \). In what direction \( u \) does it get hot fastest?

33 Figure A shows level curves of \( z = f(x, y) \).

(a) Estimate the direction and length of \( \text{grad } f \) at \( P, Q, R \).

(b) Locate two points where \( \text{grad } f \) is parallel to \( i + j \).

(c) Where is \( |\text{grad } f| \) largest? Where is it smallest?

(d) What is your estimate of \( z_{\text{max}} \) on this figure?

(e) On the straight line from \( P \) to \( R \), describe \( z \) and estimate its maximum.
A quadratic function \( ax^2 + by^2 + cx + dy \) has the gradients shown in Figure B. Estimate \( a, b, c, d \) and sketch two level curves.

The level curves of \( f(x, y) \) are circles around \((1, 1)\). The curve \( f = c \) has radius 2c. What is \( f \)? What is \( \nabla f \)?

Suppose \( \nabla f \) is tangent to the hyperbolas \( xy = \text{constant} \) in Figure C. Draw three level curves of \( f(x, y) \). Is \( |\nabla f| \) larger at \( P \) or \( Q \)? Is \( |\nabla f| \) constant along the hyperbolas? Choose a function that could be \( f \): \( x^2 + y^2 \), \( x^2 - y^2 \), \( xy \), \( x^2y^2 \).

The distance \( D \) from \((x, y)\) to \((1, 2)\) has \( D^2 = (x - 1)^2 + (y - 2)^2 \). Show that \( \partial D/\partial x = \left( x - 1 \right)/D \) and \( \partial D/\partial y = \left( y - 2 \right)/D \) and \( |\nabla D| = 1 \). The graph of \( D(x, y) \) is a ________ with its vertex at ________.

If \( f = 1 \) and \( \nabla f = (2, 3) \) at the point \((4, 5)\), find the tangent plane at \((4, 5)\). If \( f \) is a linear function, find \( f(x, y) \).

The slope of \( f \) along a level curve is \( df/ds = \_\_\_\_\_\_\_ = 0 \). This says that \( \nabla f \) is perpendicular to the vector ________ in the level direction.

### 13.5 The Chain Rule

Calculus goes back and forth between solving problems and getting ready for harder problems. The first is "application," the second looks like "theory." If we minimize \( f \) to save time or money or energy, that is an application. If we don't take derivatives to find the minimum—maybe because \( f \) is a function of other functions, and we don't have a chain rule—then it is time for more theory. The chain rule is a fundamental working tool, because \( f(g(x)) \) appears all the time in applications. So do \( f(x(t), y(t)) \) and \( f(x(t), y(t)) \) and worse. We have to know their derivatives. Otherwise calculus can't continue with the applications.

You may instinctively say: Don't bother with the theory, just teach me the formulas. That is not possible. You now regard the derivative of \( \sin 2x \) as a trivial problem, unworthy of an answer. That was not always so. Before the chain rule, the slopes of \( \sin 2x \) and \( \sin x^2 \) and \( \sin^2 x^2 \) were hard to compute from \( df/dx \). We are now at the same point for \( f(x, y) \). We know the meaning of \( \partial f/\partial x \), but if \( f = r \tan \theta \) and \( x = r \cos \theta \) and \( y = r \sin \theta \), we need a way to compute \( \partial f/\partial x \). A little theory is unavoidable, if the problem-solving part of calculus is to keep going.

To repeat: The chain rule applies to a function of a function. In one variable that was \( f(g(x)) \). With two variables there are more possibilities:

1. \( f(x) \) with \( z = g(x, y) \)  
   Find \( \partial f/\partial x \) and \( \partial f/\partial y \)
2. \( f(x, y) \) with \( x = x(t), y = y(t) \)  
   Find \( df/dt \)
3. \( f(x, y) \) with \( x = x(t, u), y = y(t, u) \)  
   Find \( df/dt \) and \( df/du \)
All derivatives are assumed continuous. More exactly, the input derivatives like $\partial g/\partial x$ and $dx/dt$ and $\partial x/\partial u$ are continuous. Then the output derivatives like $\partial f/\partial x$ and $df/dt$ and $\partial f/\partial u$ will be continuous from the chain rule. We avoid points like $r = 0$ in polar coordinates—where $\partial r/\partial x = x/r$ has a division by zero.

A Typical Problem Start with a function of $x$ and $y$, for example $x$ times $y$. Thus $f(x, y) = xy$. Change $x$ to $r \cos \theta$ and $y$ to $r \sin \theta$. The function becomes $(r \cos \theta)(r \sin \theta)$. We want its derivatives with respect to $r$ and $\theta$. First we have to decide on its name.

To be correct, we should not reuse the letter $f$. The new function can be $F$:

$$f(x, y) = xy \quad f(r \cos \theta, r \sin \theta) = (r \cos \theta)(r \sin \theta) = F(r, \theta).$$

Why not call it $f(r, \theta)$? Because strictly speaking that is $r$ times $\theta$! If we follow the rules, then $f(x, y)$ is $xy$ and $f(r, \theta)$ should be $r\theta$. The new function $F$ does the right thing—it multiplies $(r \cos \theta)(r \sin \theta)$. But in many cases, the rules get bent and the letter $F$ is changed back to $f$.

This crime has already occurred. The end of the last page ought to say $\partial F/\partial t$. Instead the printer put $df/dt$. The purpose of the chain rule is to find derivatives in the new variables $t$ and $u$ (or $r$ and $\theta$). In our example we want the derivative of $F$ with respect to $r$. Here is the chain rule:

$$\frac{\partial F}{\partial r} = \frac{\partial f}{\partial x} \frac{\partial x}{\partial r} + \frac{\partial f}{\partial y} \frac{\partial y}{\partial r} = (y)(\cos \theta) + (x)(\sin \theta) = 2r \sin \theta \cos \theta.$$

I substituted $r \sin \theta$ and $r \cos \theta$ for $y$ and $x$. You immediately check the answer: $F(r, \theta) = r^2 \cos \theta \sin \theta$ does lead to $\partial F/\partial r = 2r \cos \theta \sin \theta$. The derivative is correct. The only incorrect thing—but we do it anyway—is to write $f$ instead of $F$.

Question What is $\frac{\partial f}{\partial \theta}$? Answer It is $\frac{\partial f}{\partial x} \frac{\partial x}{\partial \theta} + \frac{\partial f}{\partial y} \frac{\partial y}{\partial \theta}$.

THE DERIVATIVES OF $f(g(x, y))$

Here $g$ depends on $x$ and $y$, and $f$ depends on $g$. Suppose $x$ moves by $dx$, while $y$ stays constant. Then $g$ moves by $dg = (\partial g/\partial x)dx$. When $g$ changes, $f$ also changes: $df = (df/dg)dg$. Now substitute for $dg$ to make the chain: $df = (df/dg)(\partial g/\partial x)dx$. This is the first rule:

$$13G \quad \text{Chain rule for } f(g(x, y)): \quad \frac{\partial f}{\partial x} = \frac{df}{dg} \frac{\partial g}{\partial x} \quad \text{and} \quad \frac{\partial f}{\partial y} = \frac{df}{dg} \frac{\partial g}{\partial y}. \quad (1)$$

EXAMPLE 1 Every $f(x + cy)$ satisfies the 1-way wave equation $\partial f/\partial y = c \partial f/\partial x$.

The inside function is $g = x + cy$. The outside function can be anything, $g^2$ or $\sin g$ or $e^g$. The composite function is $(x + cy)^2$ or $\sin(x + cy)$ or $e^{x+cy}$. In each separate case we could check that $\partial f/\partial y = c \partial f/\partial x$. The chain rule produces this equation in all cases at once, from $\partial g/\partial x = 1$ and $\partial g/\partial y = c$:

$$\frac{\partial f}{\partial x} = \frac{df}{dg} \frac{\partial g}{\partial x} = 1 \frac{df}{dg} \quad \text{and} \quad \frac{\partial f}{\partial y} = \frac{df}{dg} \frac{\partial g}{\partial y} = c \frac{df}{dg} \quad \text{so} \quad \frac{\partial f}{\partial y} = c \frac{\partial f}{\partial x}. \quad (2)$$

This is important: $\partial f/\partial y = c \partial f/\partial x$ is our first example of a partial differential equation. The unknown $f(x, y)$ has two variables. Two partial derivatives enter the equation.
Up to now we have worked with \( \frac{dy}{dt} \) and ordinary differential equations. The independent variable was time or space (and only one dimension in space). For partial differential equations the variables are time and space (possibly several dimensions in space). The great equations of mathematical physics—heat equation, wave equation, Laplace's equation—are partial differential equations.

Notice how the chain rule applies to \( f = \sin xy \). Its \( x \) derivative is \( y \cos xy \). A patient reader would check that \( f \) is \( \sin g \) and \( g \) is \( xy \) and \( f_x = f_g g_x \). Probably you are not so patient—you know the derivative of \( \sin xy \). Therefore we pass quickly to the next chain rule. Its outside function depends on two inside functions, and each of those depends on \( t \). We want \( \frac{df}{dt} \).

THE DERIVATIVE OF \( f(x(t), y(t)) \)

Before the formula, here is the idea. Suppose \( t \) changes by \( \Delta t \). That affects \( x \) and \( y \); they change by \( \Delta x \) and \( \Delta y \). There is a domino effect on \( f \); it changes by \( \Delta f \). Tracing backwards,

\[
\Delta f \approx \frac{\partial f}{\partial x} \Delta x + \frac{\partial f}{\partial y} \Delta y \quad \text{and} \quad \Delta x \approx \frac{dx}{dt} \Delta t \quad \text{and} \quad \Delta y \approx \frac{dy}{dt} \Delta t.
\]

Substitute the last two into the first, connecting \( \Delta f \) to \( \Delta t \). Then let \( \Delta t \to 0 \):

\[
\frac{df}{dt} = \frac{\partial f}{\partial x} \frac{dx}{dt} + \frac{\partial f}{\partial y} \frac{dy}{dt}.
\]

This is close to the one-variable rule \( \frac{dz}{dx} = (\frac{dz}{dy})(\frac{dy}{dx}) \). There we could “cancel” \( dy \). (We actually canceled \( \Delta y \) in \( (\Delta z/\Delta y)(\Delta y/\Delta x) \), and then approached the limit.) Now \( \Delta t \) affects \( \Delta f \) in two ways, through \( x \) and through \( y \). The chain rule has two terms. If we cancel in \( (\frac{\partial f}{\partial x})(dx/dt) \) we only get one of the terms.

We mention again that the true name for \( f(x(t), y(t)) \) is \( F(t) \) not \( f(t) \). For \( f(x, y, z) \) the rule has three terms: \( f_x x_x + f_y y_t + f_z z_t \) is \( f_t \) (or better \( df/dt \)).

EXAMPLE 2 How quickly does the temperature change when you drive to Florida?

Suppose the Midwest is at 30°F and Florida is at 80°F. Going 1000 miles south increases the temperature \( f(x, y) \) by 50°, or \( .05 \) degrees per mile. Driving straight south at 70 miles per hour, the rate of increase is \( (.05)(70) = 3.5 \) degrees per hour. Note how \( \frac{\text{degrees/mile}}{\text{miles/hour}} = \frac{\text{degrees/hour}}{} \). That is the ordinary chain rule \( \frac{df}{dx}(dx/dt) = \frac{df}{dt} \)—there is no \( y \) variable going south.

If the road goes southeast, the temperature is \( f = 30 + .05x + .01y \). Now \( x(t) \) is distance south and \( y(t) \) is distance east. What is \( \frac{df}{dt} \) if the speed is still 70?

Solution

\[
\frac{df}{dt} = \frac{\partial f}{\partial x} \frac{dx}{dt} + \frac{\partial f}{\partial y} \frac{dy}{dt} = .05 \frac{70}{\sqrt{2}} + .01 \frac{70}{\sqrt{2}} \approx 3 \text{ degrees/hour}.
\]

In reality there is another term. The temperature also depends directly on \( t \), because of night and day. The factor \( \cos(2\pi t/24) \) has a period of 24 hours, and it brings an extra term into the chain rule:

\[
\text{For } f(x, y, t) \text{ the chain rule is } \frac{df}{dt} = \frac{\partial f}{\partial x} \frac{dx}{dt} + \frac{\partial f}{\partial y} \frac{dy}{dt} + \frac{\partial f}{\partial t}.
\]

This is the \textit{total} derivative \( df/dt \), from all causes. Changes in \( x, y, t \) all affect \( f \). The partial derivative \( df/dt \) is only one part of \( df/dt \). (Note that \( dt/dt = 1 \).) If night and
day add 12 cos(2πt/24) to f, the extra term is $\frac{\partial f}{\partial t} = -\pi \sin(2\pi t/24)$. At nightfall that is $-\pi$ degrees per hour. You have to drive faster than 70 mph to get warm.

**SECOND DERIVATIVES**

What is $d^2 f/dt^2$? We need the derivative of (4), which is painful. It is like acceleration in Chapter 12, with many terms. So start with movement in a straight line.

Suppose $x = x_0 + t \cos \theta$ and $y = y_0 + t \sin \theta$. We are moving at the fixed angle $\theta$, with speed 1. The derivatives are $x_t = \cos \theta$ and $y_t = \sin \theta$ and $\cos^2 \theta + \sin^2 \theta = 1$. Then $df/dt$ is immediate from the chain rule:

$$f_t = f_x x_t + f_y y_t = f_x \cos \theta + f_y \sin \theta. \quad (5)$$

For the second derivative $f_{tt}$, apply this rule to $f_t$. Then $f_{tt}$ is

$$(f_t)_x \cos \theta + (f_t)_y \sin \theta = (f_{xx} \cos \theta + f_{xy} \sin \theta) \cos \theta + (f_{yx} \cos \theta + f_{yy} \sin \theta) \sin \theta.$$

Collect terms:

$$f_{tt} = f_{xx} \cos^2 \theta + 2f_{xy} \cos \theta \sin \theta + f_{yy} \sin^2 \theta. \quad (6)$$

In polar coordinates change $t$ to $r$. When we move in the $r$ direction, $\theta$ is fixed. Equation (6) gives $f_{rr}$ from $f_{xx}, f_{xy}, f_{yy}$. Second derivatives on curved paths (with new terms from the curving) are saved for the exercises.

**EXAMPLE 3**

If $f_{xx}, f_{xy}, f_{yy}$ are all continuous and bounded by $M$, find a bound on $f_{tt}$. This is the second derivative along any line.

**Solution**

Equation (6) gives $|f_{tt}| \leq M \cos^2 \theta + M \sin 2\theta + M \sin^2 \theta \leq 2M$. This upper bound $2M$ was needed in equation 13.3.14, for the error in linear approximation.

**THE DERIVATIVES OF $f(x(t, u), y(t, u))$**

Suppose there are two inside functions $x$ and $y$, each depending on $t$ and $u$. When $t$ moves, $x$ and $y$ both move: $dx = x_t dt$ and $dy = y_t dt$. Then $dx$ and $dy$ force a change in $f$: $df = f_x dx + f_y dy$. The chain rule for $\partial f/\partial t$ is no surprise:

\[
\begin{aligned}
\frac{\partial f}{\partial t} &= \frac{\partial f}{\partial x} \frac{\partial x}{\partial t} + \frac{\partial f}{\partial y} \frac{\partial y}{\partial t}. \\
\end{aligned}
\]

This rule has $\partial f/\partial t$ instead of $d/dt$, because of the extra variable $u$. The symbols remind us that $u$ is constant. Similarly $t$ is constant while $u$ moves, and there is a second chain rule for $\partial f/\partial u$: $f_u = f_x x_u + f_y y_u$.

**EXAMPLE 4**

In polar coordinates find $f_\theta$ and $f_{\theta\theta}$. Start from $f(x, y) = f(r \cos \theta, r \sin \theta)$.

The chain rule uses the $\theta$ derivatives of $x$ and $y$:

\[
\frac{\partial f}{\partial \theta} = \frac{\partial f}{\partial x} \frac{\partial x}{\partial \theta} + \frac{\partial f}{\partial y} \frac{\partial y}{\partial \theta} = \left( \frac{\partial f}{\partial x} \right) (-r \sin \theta) + \left( \frac{\partial f}{\partial y} \right) (r \cos \theta). \quad (8)
\]

The second $\theta$ derivative is harder, because (8) has four terms that depend on $\theta$. Apply the chain rule to the first term $\partial f/\partial x$. It is a function of $x$ and $y$, and $x$ and $y$ are functions of $\theta$:

\[
\frac{\partial}{\partial \theta} \left( \frac{\partial f}{\partial x} \right) = \frac{\partial}{\partial \theta} \left( \frac{\partial f}{\partial x} \right) \frac{\partial x}{\partial \theta} + \frac{\partial}{\partial \theta} \left( \frac{\partial f}{\partial x} \right) \frac{\partial y}{\partial \theta} = f_{xx}(-r \sin \theta) + f_{xy}(r \cos \theta). 
\]
The \( \theta \) derivative of \( \partial f / \partial y \) is similar. So apply the product rule to (8):
\[
f_{\theta \theta} = [f_{xx}(- r \sin \theta) + f_{x}(r \cos \theta)](- r \sin \theta) + f_{x}(- r \cos \theta)
+ [f_{xx}(- r \sin \theta) + f_{x}(r \cos \theta)](r \cos \theta) + f_{x}(r \sin \theta).
\] (9)

This formula is not attractive. In mathematics, a messy formula is almost always a signal of asking the wrong question. In fact the combination \( f_{xx} + f_{yy} \) is much more special than the separate derivatives. We might hope the same for \( f_{rr} + f_{\theta \theta} \), but dimensionally that is impossible—since \( r \) is a length and \( \theta \) is an angle. The dimensions of \( f_{xx} \) and \( f_{yy} \) are matched by \( f_{rr} \) and \( f_{\theta \theta} / r \) and \( f_{\theta \theta} / r^2 \). We could even hope that
\[
f_{xx} + f_{yy} = f_{rr} + f_{\theta \theta} + \frac{1}{r} f_r + \frac{1}{r^2} f_{\theta \theta}.
\] (10)

This equation is true. Add (5) + (6) + (9) with \( t \) changed to \( r \). *Laplace's equation* \( f_{xx} + f_{yy} = 0 \) is now expressed in polar coordinates: \( f_{rr} + f_r/r + f_{\theta \theta}/r^2 = 0 \).

**A PARADOX**

Before leaving polar coordinates there is one more question. It goes back to \( \partial r / \partial x \), which was practically the first example of partial derivatives:
\[
\frac{\partial r}{\partial x} = \frac{\partial}{\partial x} \sqrt{x^2 + y^2} = x / \sqrt{x^2 + y^2} = x / r.
\] (11)

My problem is this. We know that \( x = r \cos \theta \). So \( x / r \) on the right side is \( \cos \theta \). On the other hand \( r \) is \( x / \cos \theta \). So \( \partial r / \partial x \) is also \( 1 / \cos \theta \). How can \( \partial r / \partial x \) lead to \( \cos \theta \) one way and \( 1 / \cos \theta \) the other way?

I will admit that this cost me a sleepless night. There must be an explanation—we cannot end with \( \cos \theta = 1 / \cos \theta \). This paradox brings a new respect for partial derivatives. May I tell you what I finally noticed? You could cover up the next paragraph and think about the puzzle first.

The key to partial derivatives is to ask: Which variable is held constant? In equation (11), \( y \) is constant. But when \( r = x / \cos \theta \) gave \( \partial r / \partial x = 1 / \cos \theta \), \( \theta \) was constant. In both cases we change \( x \) and look at the effect on \( r \). The movement is on a horizontal line (constant \( y \)) or on a radial line (constant \( \theta \)). Figure 13.15 shows the difference.

**Remark** This example shows that \( \partial r / \partial x \) is different from \( 1 / (\partial x / \partial r) \). The neat formula \( (\partial r / \partial x)(\partial x / \partial r) = 1 \) is not generally true. May I tell you what takes its place? We have to include \( (\partial r / \partial y)(\partial y / \partial r) \). With two variables \( xy \) and two variables \( r \theta \), we need 2 by 2 matrices! Section 14.4 gives the details:
\[
\begin{bmatrix}
\partial r / \partial x & \partial r / \partial y \\
\partial \theta / \partial x & \partial \theta / \partial y
\end{bmatrix}
\begin{bmatrix}
\partial x / \partial r & \partial x / \partial \theta \\
\partial y / \partial r & \partial y / \partial \theta
\end{bmatrix}
= \begin{bmatrix}
1 & 0 \\
0 & 1
\end{bmatrix}.
\]

![Fig. 13.15](image)

\( dr = dx \cos \theta \) when \( y \) is constant, \( dr = dx / \cos \theta \) when \( \theta \) is constant.
This paradox points to a serious problem. In computing partial derivatives of \( f(x, y, z) \), we assumed that \( x, y, z \) were independent. Up to now, \( x \) could move while \( y \) and \( z \) were fixed. In physics and chemistry and economics that may not be possible. If there is a relation between \( x, y, z \), then \( x \) can't move by itself.

**EXAMPLE 5**

The gas law \( PV = nRT \) relates pressure to volume and temperature. \( P, V, T \) are not independent. What is the meaning of \( \partial V/\partial P \)? Does it equal \( 1/(\partial P/\partial V) \)?

Those questions have no answers, until we say what is held constant. In the paradox, \( \partial r/\partial x \) had one meaning for fixed \( y \) and another meaning for fixed \( \theta \). To indicate what is held constant, use an extra subscript (not denoting a derivative):

\[
(\partial r/\partial x)_y = \cos \theta \quad (\partial r/\partial x)_\theta = 1/\cos \theta.
\]

(\( \partial f/\partial P \))\( _\nu \) has constant volume and (\( \partial f/\partial P \))\( _r \) has constant temperature. The usual \( \partial f/\partial P \) has both \( V \) and \( T \) constant. But then the gas law won't let us change \( P \).

**EXAMPLE 6**

Let \( f = 3x + 2y + z \). Compute \( \partial f/\partial x \) on the plane \( z = 4x + y \).

**Solution 1**

Think of \( x \) and \( y \) as independent. Replace \( z \) by \( 4x + y \):

\[
f = 3x + 2y + (4x + y) \quad \text{so} \quad (\partial f/\partial x)_y = 7.
\]

**Solution 2**

Keep \( x \) and \( y \) independent. Deal with \( z \) by the chain rule:

\[
(\partial f/\partial x)_y = \partial f/\partial x + (\partial f/\partial z)(\partial z/\partial x) = 3 + (1)(4) = 7.
\]

**Solution 3**

(different) Make \( x \) and \( z \) independent. Then \( y = z - 4x \):

\[
(\partial f/\partial x)_z = \partial f/\partial x + (\partial f/\partial y)(\partial y/\partial x) = 3 + (2)(-4) = -5.
\]

Without a subscript, \( \partial f/\partial x \) means: Take the \( x \) derivative the usual way. The answer is \( \partial f/\partial x = 3 \), when \( y \) and \( z \) don't move. But on the plane \( z = 4x + y \), one of them must move! \( 3 \) is only part of the total answer, which is \( (\partial f/\partial x)_y = 7 \) or \( (\partial f/\partial x)_z = -5 \).

Here is the geometrical meaning. We are on the plane \( z = 4x + y \). The derivative \( (\partial f/\partial x)_y \) moves \( x \) but not \( y \). To stay on the plane, \( dz \) is \( 4dx \). The change in \( f = 3x + 2y + z \) is \( df = 3dx + 0 + dz = 7dx \).

**EXAMPLE 7**

On the world line \( x^2 + y^2 + z^2 = t^2 \) find \( (\partial f/\partial y)_{x,z} \) for \( f = xyzt \).

The subscripts \( x, z \) mean that \( x \) and \( z \) are fixed. The chain rule skips \( \partial f/\partial x \) and \( \partial f/\partial z \):

\[
(\partial f/\partial y)_{x,z} = \partial f/\partial y + (\partial f/\partial t)(\partial t/\partial y) = xzt + (xyz)(y/t).
\]

Why \( y/t \)?

**EXAMPLE 8**

From the law \( PV = T \), compute the product \( (\partial P/\partial V)_T(\partial V/\partial T)_p(\partial T/\partial P)_V \).

Any intelligent person cancels \( \partial V \)'s, \( \partial T \)'s, and \( \partial P \)'s to get 1. The right answer is \( -1 \):

\[
(\partial P/\partial V)_T = -T/V^2 \quad (\partial V/\partial T)_p = 1/P \quad (\partial T/\partial P)_V = V.
\]

The product is \( -T/PV \). This is \( -1 \) not \( +1 \! \) The chain rule is tricky (Problem 42).

**EXAMPLE 9**

Implicit differentiation was used in Chapter 4. The chain rule explains it:

\[
\text{If } F(x, y) = 0 \text{ then } F_x + F_y y_x = 0 \text{ so } dy/dx = -F_x/F_y. \quad (\text{13})
\]
1. Check that \( f \) and its real part and imaginary part moving and how fast?

Along a path, the derivative of \( f(x(t), y(t)) \) is \( df/dt = \ldots \). The derivative of \( f(x(t), y(t), z(t)) \) is \( df/dt = \ldots \). If \( x = x(t) \) and \( y = y(t) \), the path is \( \ldots \). The chain rule for \( f(x, y) \) gives \( df/dt = \ldots \). That is the same as the ___ rule! When \( x = u(t) \) and \( y = v(t) \), the path is \( \ldots \). The chain rule for \( f(x, y) \) gives \( df/dt = \ldots \). That is the ____ derivative \( D_u f \).

The chain rule for \( f(x, u), y(u, t) \) is \( df/dt = \ldots \). We don't write \( df/dt \) because ___.

With a relation like \( x y z = 1 \), the three variables are ___. independent. The derivatives (\( df/dx \)) and (\( df/dy \)), mean ___, and ___, and ___. For \( f = \sqrt{x^2 + y^2} \) with \( x y z = 1 \), we compute (\( df/dx \)) from the chain rule ___. In that rule \( \partial z/\partial x = \ldots \) from the relation \( x y z = 1 \).

Find \( f_x \) and \( f_y \) in Problems 1–4. What equation connects them?

1. \( f(x, y) = \sin(x + cy) \)
2. \( f(x, y) = (ax + by)^{10} \)
3. \( f(x, y) = e^{x/y} \)
4. \( f(x, y) = \ln(x + y) \)

5. Find both terms in the \( r \) derivative of \( (g(x(t), y(t)))^3 \).

6. If \( f(x, y) = xy \) and \( x = u(t) \) and \( y = v(t) \), what is \( df/dt \)? Probably all other rules for derivatives follow from the chain rule.

7. The step function \( f(x) \) is zero for \( x < 0 \) and one for \( x > 0 \). Graph \( f(x) \) and \( g(x) = f(x + 2) \) and \( h(x) = f(x + 4) \). If \( f(x + 2t) \) represents a wall of water (a tidal wave), which way is it moving and how fast?

8. The wave equation is \( f_{tt} = c^2 f_{xx} \). (a) Show that \( (x + ct)^n \) is a solution. (b) Find \( C \) different from \( c \) so that \( (x + Ct)^n \) is also a solution.

9. If \( f = \sin(x - t) \), draw two lines in the \( xt \) plane along which \( f = 0 \). Between those lines sketch a sine wave. Skiing on top of the sine wave, what is your speed \( dx/dt \)?

10. If you float at \( x = 0 \) in Problem 9, do you go up first or down first? At time \( t = 4 \) what is your height and upward velocity?

11. Laplace's equation is \( f_{xx} + f_{yy} = 0 \). Show from the chain rule that any function \( f(x + iy) \) satisfies this equation if \( i^2 = -1 \). Check that \( f = (x + iy)^2 \) and its real part ___ and its imaginary part ___ all satisfy Laplace's equation.

12. Equation (10) gave the polar form \( f_r + f_i/r + f_\theta/r^2 = 0 \) of Laplace's equation. (a) Check that \( f = r^2 e^{2\theta} \) and its real part \( r^2 \cos \theta \) and its imaginary part \( r^2 \sin \theta \) all satisfy Laplace's equation. (b) Show from the chain rule that any function \( f(re^{i\theta}) \) satisfies this equation if \( i^2 = -1 \).

In Problems 13–18 find \( df/dt \) from the chain rule (3).

13. \( f = x^2 + y^2, x = t, y = t^2 \)
14. \( f = \sqrt{x^2 + y^2}, x = t, y = t^2 \)
15. \( f = xy, x = 1 - \sqrt{t}, y = 1 + \sqrt{t} \)
16. \( f = x/y, x = e^t, y = 2e^t \)
17. \( f = \ln(x + y), x = e^t, y = e^t \)
18. \( f = x^4, x = t, y = t \)

19. If a cone grows in height by \( dh/dt = 1 \) in and radius by \( dr/dt = 2t \), starting from zero, how fast is its volume growing at \( t = 3 \)?

20. If a rocket has speed \( dx/dt = 6 \) down range and \( dy/dt = 2t \) upward, how fast is it moving away from the launch point at \( t = 0, \theta \)? How fast is the angle \( \theta \) changing, if \( \tan \theta = y/x \)?

21. If a train approaches a crossing at \( 60 \) mph and a car approaches (at right angles) at \( 45 \) mph, how fast are they coming together? (a) Assume they are both going at \( 90 \) miles from the crossing. (b) Assume they are going to hit.

22. In Example 2 does the temperature increase faster if you drive due south at \( 70 \) mph or southeast at \( 80 \) mph?

23. On the line \( x = u_1 t, y = u_2 t, z = u_3 t \), what combination of \( f_x, f_y, f_z \) gives \( df/dt \)? This is the directional derivative in 3D.

24. On the same line \( x = u_1 t, y = u_2 t, z = u_3 t \), find a formula for \( d^2 f/dt^2 \). Apply it to \( f = x y z \).

25. For \( f(x, y, z) = x + y + z \) find \( df/dt \) and \( df/dt \) when \( x = 2t \) and \( y = 3t \). Explain the difference.

26. If \( z = (x + y)^2 \) then \( x = \sqrt{z - y} \). Does \( \partial z/\partial x (\partial x/\partial z) \) = \( 1 \)?

27. Suppose \( x_1 = t \) and \( y_1 = 2t \), not constant as in (5–6). For \( f(x, y) \), find \( f_1 \), and \( f_2 \). The answer involves \( f_x, f_y, f_{xx}, f_{xy}, f_{yy} \).

28. Suppose \( x_1 = t \) and \( y_1 = t^2 \). For \( f = (x + y)^3 \) find \( f_1 \) and then \( f_{yy} \) from the chain rule.

29. Derive \( \partial f/\partial \theta = (\partial f/\partial x) \cos \theta + (\partial f/\partial y) \sin \theta \) from the chain rule. Why do we take \( \partial x/\partial \theta \) as \( \cos \theta \) and not \( 1/\cos \theta \)?

30. Compute \( f_{xx} \) for \( f(x, y) = (ax + by)^j \). If \( x = t \) and \( y = t \) compute \( f_{xx} \). True or false: \( \partial f/\partial x (\partial x/\partial z) = \partial f/\partial t \).

31. Show that \( \partial^2 r/\partial x^2 = y^2/r^3 \) in two ways:

(1) Find the \( x \) derivative of \( \partial r/\partial x = x/\sqrt{x^2 + y^2} \)

(2) Find the \( x \) derivative of \( \partial r/\partial x = x/r \) by the chain rule.
32. Reversing \( x \) and \( y \) in Problem 31 gives \( r_{xy} = x^3/r^3 \). But show that \( r_{xy} = -xy/r^3 \).

33. If \( \sin z = x + y \) find \( (\partial z/\partial x) \), in two ways:
   
   (1) Write \( z = \sin^{-1}(x + y) \) and compute its derivative.
   
   (2) Take \( x \) derivatives of \( \sin z = x + y \). Verify that these answers, explicit and implicit, are equal.

34. By direct computation find \( f_x \) and \( f_{xx} \) and \( f_{xy} \) for \( f = \sqrt{x^2 + y^2} \).

35. Find a formula for \( \partial^2 f/\partial \theta \partial \phi \) in terms of the \( x \) and \( y \) derivatives of \( f(x, y) \).

36. Suppose \( z = f(x, y) \) is solved for \( x \) to give \( x = g(y, z) \). Is it true that \( \partial z/\partial x = 1/(\partial x/\partial z) \)? Test on examples.

37. Suppose \( z = e^x \) and therefore \( x = (\ln z)/y \). Is it true or not that \( (\partial z/\partial x) = 1/(\partial x/\partial z) \)?

38. If \( x = x(t, u, v) \) and \( y = y(t, u, v) \) and \( z = z(t, u, v) \), find the \( t \) derivative of \( f(x(t, u, v), y(t, u, v)) \) in equation (7). What is \( f_t \)?

39. The \( t \) derivative of \( f(x(t, u, v), y(t, u, v)) \) is in equation (7). What is \( f_u \) ?

40. (a) For \( f = x^2 + y^2 + z^2 \) compute \( \partial f/\partial x \) (no subscript, \( x, y, z \) all independent).
   
   (b) When there is a further relation \( z = x^2 + y^2 \), use it to remove \( z \) and compute \( (\partial f/\partial x)_y \).
   
   (c) Compute \( (\partial f/\partial x)_y \) using the chain rule \( (\partial f/\partial x) + (\partial f/\partial z)(\partial z/\partial x) \).
   
   (d) Why doesn’t that chain rule contain \( (\partial f/\partial y)(\partial y/\partial x) \)?

41. For \( f = ax + by \) on the plane \( z = 3x + 5y \), find \( (\partial f/\partial x)_y \) and \( (\partial f/\partial z)_y \).

42. The gas law in physics is \( PV = nRT \) or a more general relation \( F(P, V, T) = 0 \). Show that the three derivatives in Example 8 still multiply to give \(-1\). First find \( \partial F/\partial V \) from \( \partial F/\partial V + (\partial F/\partial P)(\partial P/\partial V) = 0 \).

43. If Problem 42 changes to four variables related by \( F(x, y, z, t) = 0 \), what is the corresponding product of four derivatives?

44. Suppose \( x = t + u \) and \( y = tu \). Find the \( t \) and \( u \) derivatives of \( f(x, y) \). Check when \( f(x, y) = x^2 - 2y \).

45. (a) For \( f = r^2 \sin^2 \theta \) find \( f_x \) and \( f_y \).
   
   (b) For \( f = x^2 + y^2 \) find \( f_x \) and \( f_y \).

46. On the curve \( \sin x + \sin y = 0 \), find \( dy/dx \) and \( d^2 y/dx^2 \) by implicit differentiation.

47. (horrible) Suppose \( f_{xx} + f_{yy} = 0 \). If \( x = u + v \) and \( y = u - v \) and \( f(x, y) = g(u, v) \), find \( g_u \) and \( g_v \). Show that \( g_u + g_v = 0 \).

48. A function has constant returns to scale if \( f(cx, cy) = cf(x, y) \). When \( x \) and \( y \) are doubled so are \( f = \sqrt{x^2 + y^2} \) and \( f = \sqrt{xy} \). In economics, input/output is constant. In mathematics \( f \) is homogeneous of degree one.

   Prove that \( x \partial f/\partial x + y \partial f/\partial y = f(x, y) \), by computing the \( c \) derivative at \( c = 1 \). Test this equation on the two examples and find a third example.

49. True or false: The directional derivative of \( f(r, \theta) \) in the direction of \( u_0 \) is \( \partial f/\partial \theta \).

#### 13.6 Maxima, Minima, and Saddle Points

The outstanding equation of differential calculus is also the simplest: \( df/dx = 0 \). The slope is zero and the tangent line is horizontal. Most likely we are at the top or bottom of the graph—a maximum or a minimum. This is the point that the engineer or manager or scientist or investor is looking for—maximum stress or production or velocity or profit. With more variables in \( f(x, y) \) and \( f(x, y, z) \), the problem becomes more realistic. The question still is: How to locate the maximum and minimum?

The answer is in the partial derivatives. When the graph is level, they are zero. Deriving the equations \( f_x = 0 \) and \( f_y = 0 \) is pure mathematics and pure pleasure. Applying them is the serious part. We watch out for saddle points, and also for a minimum at a boundary point—this section takes extra time. Remember the steps for \( f(x) \) in one-variable calculus:

1. The leading candidates are stationary points (where \( df/dx = 0 \)).
2. The other candidates are rough points (no derivative) and endpoints \((a \text{ or } b)\).
3. Maximum vs. minimum is decided by the sign of the second derivative.

In two dimensions, a stationary point requires \( df/dx = 0 \) and \( df/dy = 0 \). The tangent line becomes a tangent plane. The endpoints \( a \) and \( b \) are replaced by a boundary curve. In practice boundaries contain about 40% of the minima and 80% of the work.
Finally there are three second derivatives \( f_{xx}, f_{xy}, \) and \( f_{yy} \). They tell how the graph bends away from the tangent plane—up at a \textit{minimum}, down at a \textit{maximum}, both ways at a \textit{saddle point}. This will be determined by comparing \((f_{xx})(f_{yy})\) with \((f_{xy})^2\).

### STATIONARY POINT → HORIZONTAL TANGENT → ZERO DERIVATIVES

Suppose \( f \) has a minimum at the point \((x_0, y_0)\). This may be an \textit{absolute minimum} or only a \textit{local minimum}. In both cases \( f(x_0, y_0) \leq f(x, y) \) near the point. For an absolute minimum, this inequality holds wherever \( f \) is defined. For a local minimum, the inequality can fail far away from \((x_0, y_0)\). The bottom of your foot is an absolute minimum, the end of your finger is a local minimum.

We assume for now that \((x_0, y_0)\) is an \textit{interior point} of the domain of \( f \). At a boundary point, we cannot expect a horizontal tangent and zero derivatives.

Main conclusion: At a minimum or maximum (absolute or local) a nonzero derivative is impossible. The tangent plane would tilt. In some direction \( f \) would decrease. Note that the minimum \textit{point} is \((x_0, y_0)\), and the minimum \textit{value} is \( f(x_0, y_0) \).

### EXAMPLE 1

Minimize the quadratic \( f(x, y) = x^2 + xy + y^2 - x - y + 1 \).

To locate the minimum (or maximum), set \( f_x = 0 \) and \( f_y = 0 \):

\[
\begin{align*}
f_x &= 2x + y - 1 = 0 \\
f_y &= x + 2y - 1 = 0.
\end{align*}
\]
Notice what's important: There are two equations for two unknowns \( x \) and \( y \). Since \( f \) is quadratic, the equations are linear. Their solution is \( x_0 = \frac{1}{3}, y_0 = \frac{1}{3} \) (the stationary point). This is actually a minimum, but to prove that you need to read further.

The constant 1 affects the minimum value \( f = \frac{1}{3} \)—but not the minimum point. The graph shifts up by 1. The linear terms \(-x - y\) affect \( f_x \) and \( f_y \). They move the minimum away from \((0,0)\). The quadratic part \( x^2 + xy + y^2 \) makes the surface curve upwards. Without that curving part, a plane has its minimum and maximum at boundary points.

**EXAMPLE 2 (Steiner's problem)** Find the point that is nearest to three given points.

This example is worth your attention. We are locating an airport close to three cities. Or we are choosing a house close to three jobs. The problem is to get as near as possible to the corners of a triangle. The best point depends on the meaning of "near."

The distance to the first corner \((x_1, y_1)\) is \( d_1 = \sqrt{(x-x_1)^2 + (y-y_1)^2} \). The distances to the other corners \((x_2, y_2)\) and \((x_3, y_3)\) are \( d_2 \) and \( d_3 \). Depending on whether cost equals \((\text{distance})\) or \((\text{distance})^2\) or \((\text{distance})^n\), our problem will be:

\[
\text{Minimize } d_1 + d_2 + d_3 \quad \text{or} \quad d_1^2 + d_2^2 + d_3^2 \quad \text{or even} \quad d_1^n + d_2^n + d_3^n.
\]

The second problem is the easiest, when \( d_1^2 \) and \( d_2^2 \) and \( d_3^2 \) are quadratics:

\[
f(x, y) = (x - x_1)^2 + (y - y_1)^2 + (x - x_2)^2 + (y - y_2)^2 + (x - x_3)^2 + (y - y_3)^2
\]

\[
\partial f/\partial x = 2[x - x_1 + x - x_2 + x - x_3] = 0 \quad \partial f/\partial y = 2[y - y_1 + y - y_2 + y - y_3] = 0.
\]

Solving \( \partial f/\partial x = 0 \) gives \( x = \frac{1}{3}(x_1 + x_2 + x_3) \). Then \( \partial f/\partial y = 0 \) gives \( y = \frac{1}{3}(y_1 + y_2 + y_3) \). The best point is the centroid of the triangle (Figure 13.17a). It is the nearest point to the corners when the cost is \((\text{distance})^2\). Note how squaring makes the derivatives linear. Least squares dominates an enormous part of applied mathematics.

![Fig. 13.17 The centroid minimizes \( d_1^2 + d_2^2 + d_3^2 \). The Steiner point minimizes \( d_1 + d_2 + d_3 \).](image)

The real "Steiner problem" is to minimize \( f(x, y) = d_1 + d_2 + d_3 \). We are laying down roads from the corners, with cost proportional to length. The equations \( f_x = 0 \) and \( f_y = 0 \) look complicated because of square roots. But the nearest point in Figure 13.17b has a remarkable property, which you will appreciate.

Calculus takes derivatives of \( d_1^2 = (x-x_1)^2 + (y-y_1)^2 \). The \( x \) derivative leaves \( 2d_1(\partial d_1/\partial x) = 2(x-x_1) \). Divide both sides by \( 2d_1 \):

\[
\frac{\partial d_1}{\partial x} = \frac{x-x_1}{d_1} \quad \text{and} \quad \frac{\partial d_1}{\partial y} = \frac{y-y_1}{d_1} \quad \text{so grad } d_1 = \left( \frac{x-x_1}{d_1}, \frac{y-y_1}{d_1} \right). \tag{3}
\]

This gradient is a unit vector. The sum of \((x-x_1)^2/d_1^2\) and \((y-y_1)^2/d_1^2\) is \( d_1^2/d_1^2 = 1 \). This was already in Section 13.4: Distance increases with slope 1 away from the center. The gradient of \( d_1 \) (call it \( u_1 \)) is a unit vector from the center point \((x_1, y_1)\).
Similarly the gradients of \( d_2 \) and \( d_3 \) are unit vectors \( u_2 \) and \( u_3 \). They point directly away from the other corners of the triangle. The total cost is \( f(x, y) = d_1 + d_2 + d_3 \), so we add the gradients. The equations \( f_x = 0 \) and \( f_y = 0 \) combine into the vector equation

\[
\text{grad } f = u_1 + u_2 + u_3 = 0 \text{ at the minimum.}
\]

The three unit vectors add to zero! Moving away from one corner brings us closer to another. The nearest point to the three corners is where those movements cancel. This is the meaning of “\( \text{grad } f = 0 \) at the minimum.”

It is unusual for three unit vectors to add to zero—this can only happen in one way. **The three directions must form angles of 120°.** The best point has this property, which is repeated in Figure 13.18a. The unit vectors cancel each other. At the “Steiner point,” the roads to the corners make 120° angles. This optimal point solves the problem, except for one more possibility.

Fig. 13.18 Gradients \( u_1 + u_2 + u_3 = 0 \) for 120° angles. Corner wins at wide angle. **Four corners.** In this case two branchpoints are better—still 120°.

The other possibility is a minimum at a rough point. The graph of the distance function \( d_1(x, y) \) is a cone. It has a sharp point at the center \((x_1, y_1)\). All three corners of the triangle are rough points for \( d_1 + d_2 + d_3 \), so all of them are possible minimizers.

**Suppose the angle at a corner exceeds 120°.** Then there is no Steiner point. Inside the triangle, the angle would become even wider. The best point must be a rough point—one of the corners. The winner is the corner with the wide angle. In the figure that means \( d_1 = 0 \). Then the sum \( d_2 + d_3 \) comes from the two shortest edges.

**Summary** The solution is at a 120° point or a wide-angle corner. That is the theory. The real problem is to compute the Steiner point—which I hope you will do.

**Remark 1** Steiner’s problem for four points is surprising. We don’t minimize \( d_1 + d_2 + d_3 + d_4 \)—there is a better problem. Connect the four points with roads, minimizing their total length, and allow the roads to branch. A typical solution is in Figure 13.18c. The angles at the branch points are 120°. There are at most two branch points (two less than the number of corners).

**Remark 2** For other powers \( p \), the cost is \((d_1)^p + (d_2)^p + (d_3)^p\). The x derivative is

\[
\frac{\partial f}{\partial x} = p(d_1)^{p-2}(x - x_1) + p(d_2)^{p-2}(x - x_2) + p(d_3)^{p-2}(x - x_3).
\]

The key equations are still \( \frac{\partial f}{\partial x} = 0 \) and \( \frac{\partial f}{\partial y} = 0 \). Solving them requires a computer and an algorithm. To share the work fairly, I will supply the algorithm (Newton’s method) if you supply the computer. Seriously, this is a terrific example. It is typical of real problems—we know \( \frac{\partial f}{\partial x} \) and \( \frac{\partial f}{\partial y} \) but not the point where they are zero. You can calculate that nearest point, which changes as \( p \) changes. You can also discover new mathematics, about how that point moves. I will collect all replies I receive to Problems 38 and 39.
Steiner's problem had no boundaries. The roads could go anywhere. But most applications have restrictions on $x$ and $y$, like $x \geq 0$ or $y \leq 0$ or $x^2 + y^2 \geq 1$. The minimum with these restrictions is probably higher than the absolute minimum. There are three possibilities:

1) **stationary point** $f_x = 0, f_y = 0$
2) **rough point**
3) **boundary point**

That third possibility requires us to maximize or minimize $f(x, y)$ along the boundary.

**EXAMPLE 3** Minimize $f(x, y) = x^2 + xy + y^2 - x - y + 1$ in the half-plane $x \geq 0$.

The minimum in Example 1 was $\frac{3}{4}$. It occurred at $x_0 = \frac{1}{3}, y_0 = \frac{1}{3}$. This point is still allowed. It satisfies the restriction $x \geq 0$. So the minimum is not moved.

**EXAMPLE 4** Minimize the same $f(x, y)$ restricted to the lower half-plane $y \leq 0$.

Now the absolute minimum at $(\frac{1}{3}, \frac{1}{3})$ is not allowed. There are no rough points. We look for a minimum on the boundary line $y = 0$ in Figure 13.19a. Set $y = 0$, so $f$ depends only on $x$. Then choose the best $x$:

$$f(x, 0) = x^2 - x - 1 = 0$$
and

$$f_x = 2x - 1 = 0.$$  

The minimum is at $x = \frac{1}{2}$ and $y = 0$, where $f = \frac{3}{4}$. This is up from $\frac{1}{2}$.

**Fig. 13.19** The boundaries $y = 0$ and $x^2 + y^2 = 1$ contain the minimum points.

**EXAMPLE 5** Minimize the same $f(x, y)$ on or outside the circle $x^2 + y^2 = 1$.

One possibility is $f_x = 0$ and $f_y = 0$. But this is at $(\frac{1}{3}, \frac{1}{3})$, inside the circle. The other possibility is a minimum at a boundary point, **on the circle**.

To follow this boundary we can set $y = \sqrt{1 - x^2}$. The function $f$ gets complicated, and $df/dx$ is worse. There is a way to avoid square roots. Set $x = \cos t$ and $y = \sin t$. Then $f = x^2 + xy + y^2 - x - y + 1$ is a function of the angle $t$:

$$f(t) = 1 + \cos t \sin t - \cos t - \sin t + 1$$

$$df/dt = \cos^2 t - \sin^2 t + \sin t - \cos t = (\cos t - \sin t)(\cos t + \sin t - 1).$$

Now $df/dt = 0$ locates a minimum or maximum along the boundary. The first factor $(\cos t - \sin t)$ is zero when $x = y$. The second factor is zero when $\cos t + \sin t = 1$, or $x + y = 1$. Those points on the circle are the candidates. Problem 24 sorts them out, and Section 13.7 finds the minimum in a new way, using “Lagrange multipliers.”
13.6  \textbf{Maxima, Minima, and Saddle Points}

Minimization on a boundary is a serious problem—it gets difficult quickly—and multipliers are ultimately the best solution.

\textbf{MAXIMUM VS. MINIMUM VS. SADDLE POINT}

How to separate the maximum from the minimum? When possible, try all candidates and decide. Compute \( f \) at every stationary point and other critical point (maybe also out at infinity), and compare. Calculus offers another approach, based on \textit{second derivatives}.

With one variable the second derivative test was simple: \( f'' > 0 \) at a minimum, \( f'' = 0 \) at an inflection point, \( f'' < 0 \) at a maximum. This is a local test, which may not give a global answer. But it decides whether the slope is increasing (bottom of the graph) or decreasing (top of the graph). We now find a similar test for \( f(x, y) \).

The new test involves all three second derivatives. It applies where \( f_x = 0 \) and \( f_y = 0 \). The tangent plane is horizontal. \textit{We ask whether the graph of \( f \) goes above or below that plane.} The tests \( f_{xx} > 0 \) and \( f_{yy} > 0 \) guarantee a minimum in the \( x \) and \( y \) directions, but there are other directions.

\textbf{EXAMPLE 6}  \( f(x, y) = x^2 + 10xy + y^2 \) has \( f_{xx} = 2, f_{xy} = 10, f_{yy} = 2 \) (minimum or not?)

All second derivatives are positive—but wait and see. The stationary point is \( (0, 0) \), where \( \partial f/\partial x \) and \( \partial f/\partial y \) are both zero. Our function is the sum of \( x^2 + y^2 \), which goes upward, and \( 10xy \) which has a saddle. The second derivatives must decide whether \( x^2 + y^2 \) or \( 10xy \) is stronger.

Along the \( x \) axis, where \( y = 0 \) and \( f = x^2 \), our point is at the bottom. The minimum in the \( x \) direction is at \( (0, 0) \). Similarly for the \( y \) direction. But \( (0, 0) \) is \textit{not a minimum point} for the whole function, because of \( 10xy \).

Try \( x = 1, y = -1 \). Then \( f = 1 - 10 + 1 \), which is negative. The graph goes below the \( xy \) plane in that direction. The stationary point at \( x = y = 0 \) is a \textit{saddle point}.

\textbf{EXAMPLE 7}  \( f(x, y) = x^2 + xy + y^2 \) has \( f_{xx} = 2, f_{xy} = 1, f_{yy} = 2 \) (minimum or not?)

The second derivatives 2, 1, 2 are again positive. The graph curves up in the \( x \) and \( y \) directions. But there is a big difference from Example 6: \( f_{xy} \) is reduced from 10 to 1. \textit{It is the size of} \( f_{xy} \) \textit{(not its sign!) that makes the difference.} The extra terms \(-x - y + 4\) in Example 1 moved the stationary point to \( \left(\frac{1}{2}, \frac{1}{2}\right) \). The second derivatives are still 2, 1, 2, and they pass the test for a minimum:

\begin{center}
\begin{tabular}{lll}
\textbf{13K} & At (0, 0) the quadratic function \( f(x, y) = ax^2 + 2bxy + cy^2 \) has a & \\
& \textit{minimum} if & \textit{maximum} if & \textit{saddle point} if \\
& \begin{align*}
& a > 0 \\
& ac > b^2
\end{align*} & \begin{align*}
& a < 0 \\
& ac > b^2
\end{align*} & \begin{align*}
& ac < b^2
\end{align*}
\end{tabular}
\end{center}
For a direct proof, split \( f(x, y) \) into two parts by "completing the square:"

\[
ax^2 + 2bxy + cy^2 = a\left(x + \frac{b}{a}y\right)^2 + \frac{ac - b^2}{a}y^2.
\]

That algebra can be checked (notice the \( 2b \)). It is the conclusion that's important:

- if \( a > 0 \) and \( ac > b^2 \), both parts are positive: \textit{minimum} at \((0, 0)\)
- if \( a < 0 \) and \( ac > b^2 \), both parts are negative: \textit{maximum} at \((0, 0)\)
- if \( ac < b^2 \), the parts have opposite signs: \textit{saddle point} at \((0, 0)\).

Since the test involves the square of \( b \), its sign has no importance. Example 6 had \( b = 5 \) and a saddle point. Example 7 had \( b = \frac{1}{2} \) and a minimum. Reversing to \(-x^2 - xy - y^2\) yields a maximum. So does \(-x^2 + xy - y^2\).

Now comes the final step, from \( ax^2 + 2bxy + cy^2 \) to a general function \( f(x, y) \). For all functions, quadratics or not, it is the \textit{second order terms} that we test.

**Example 8** \( f(x, y) = e^x - x - \cos y \) has a stationary point at \( x = 0, y = 0 \).

The first derivatives are \( e^x - 1 \) and \( \sin y \), both zero. The second derivatives are \( f_{xx} = e^x = 1 \) and \( f_{yy} = \cos y = 1 \) and \( f_{xy} = 0 \). We only use the derivatives at the stationary point. The first derivatives are zero, so the second order terms come to the front in the series for \( e^x - x - \cos y \):

\[
(1 + x + \frac{1}{2}x^2 + \cdots) - x - (1 - \frac{1}{2}y^2 + \cdots) = \frac{1}{2}x^2 + \frac{1}{2}y^2 + \text{higher order terms}. \quad (7)
\]

There is a \textit{minimum} at the origin. The quadratic part \( \frac{1}{2}x^2 + \frac{1}{2}y^2 \) goes upward. The \( x^3 \) and \( y^4 \) terms are too small to protest. Eventually those terms get large, but near a stationary point it is the quadratic that counts. We didn't need the whole series, because from \( f_{xx} = f_{yy} = 1 \) and \( f_{xy} = 0 \) we knew it would start with \( \frac{1}{2}x^2 + \frac{1}{2}y^2 \).

**Example 9** \( f(x, y) = e^{xy} \) has \( f_x = ye^{xy} \) and \( f_y = xe^{xy} \). The stationary point is \((0, 0)\).

The second derivatives at that point are \( a = f_{xx} = 0 \), \( b = f_{xy} = 1 \), and \( c = f_{yy} = 0 \). The test \( b^2 > ac \) makes this a saddle point. Look at the infinite series:

\[
e^{xy} = 1 + xy + \frac{1}{2}x^2y^2 + \cdots.
\]

No linear term because \( f_x = f_y = 0 \): The origin is a \textit{stationary point}. No \( x^2 \) or \( y^2 \) term (only \( xy \)): The stationary point is a \textit{saddle point}.

At \( x = 2, \ y = -2 \) we find \( f_{xx}f_{yy} > (f_{xy})^2 \). The graph is concave up at that point—but it's not a minimum since the first derivatives are not zero.

The series begins with the constant term—not important. Then come the linear terms—extremely important. Those terms are decided by \textit{first} derivatives, and they give the tangent plane. It is only at stationary points—when the linear part disappears and the tangent plane is horizontal—that second derivatives take over. Around any basepoint, \textit{these constant-linear-quadratic terms are the start of the Taylor series}. 

\[13L\] The test in 13K applies to the second derivatives \( a = f_{xx}, \ b = f_{xy}, \ c = f_{yy} \) of any \( f(x, y) \) at any stationary point. At all points the test decides whether the graph is concave up, concave down, or "indefinite."
We now put together the whole infinite series. It is a "Taylor series"—which means it is a power series that matches all derivatives of \( f \) (at the basepoint). For one variable, the powers were \( x^n \) when the basepoint was 0. For two variables, the powers are \( x^n y^m \) when the basepoint is \((0, 0)\). Chapter 10 multiplied the \( n \)th derivative \( d^n f/dx^n \) by \( x^n/n! \) Now every mixed derivative \( (\partial/\partial x)^n(\partial/\partial y)^m f(x, y) \) is computed at the basepoint (subscript \( o \)).

We multiply those numbers by \( x^n y^m/n!m! \) to match each derivative of \( f(x, y) \):

The first three terms are the linear approximation to \( f(x, y) \). They give the tangent plane at the basepoint. The \( x^2 \) term has \( n = 2 \) and \( m = 0 \), so \( n!m! = 2 \). The \( xy \) term has \( n = m = 1 \), and \( n!m! = 1 \). The quadratic part \( \frac{1}{2}(ax^2 + 2bxy + cy^2) \) is in control when the linear part is zero.

**EXAMPLE 10** All derivatives of \( e^{x+y} \) equal one at the origin. The Taylor series is

\[
e^{x+y} = 1 + x + y + \frac{x^2}{2} + xy + \frac{y^2}{2} + \cdots = \sum x^n y^m/n!m!
\]

This happens to have \( ac = b^2 \), the special case that was omitted in 13M and 13N. It is the two-dimensional version of an inflection point. The second derivatives fail to decide the concavity. When \( f_{xx} f_{yy} = (f_{xy})^2 \), the decision is passed up to the higher derivatives. But in ordinary practice, the Taylor series is stopped after the quadratics.

If the basepoint moves to \((x_0, y_0)\), the powers become \((x - x_0)^n(y - y_0)^m\)—and all derivatives are computed at this new basepoint.

**Final question:** How would you compute a minimum numerically? One good way is to solve \( f_x = 0 \) and \( f_y = 0 \). These are the functions \( g \) and \( h \) of Newton’s method (Section 13.3). At the current point \((x_n, y_n)\), the derivatives of \( g = f_x \) and \( h = f_y \) give linear equations for the steps \( \Delta x \) and \( \Delta y \). Then the next point \( x_{n+1} = x_n + \Delta x, \ y_{n+1} = y_n + \Delta y \) comes from those steps. The input is \((x_n, y_n)\), the output is the new point, and the linear equations are

\[
(g_x)\Delta x + (g_y)\Delta y = -g(x_n, y_n) \quad \text{or} \quad (f_{xx})\Delta x + (f_{xy})\Delta y = -f_x(x_n, y_n)
\]

\[
(h_x)\Delta x + (h_y)\Delta y = -h(x_n, y_n) \quad \text{or} \quad (f_{xy})\Delta x + (f_{yy})\Delta y = -f_y(x_n, y_n).
\]

When the second derivatives of \( f \) are available, use Newton’s method.

When the problem is too complicated to go beyond first derivatives, here is an alternative—steepest descent. The goal is to move down the graph of \( f(x, y) \), like a boulder rolling down a mountain. The steepest direction at any point is given by the gradient, with a minus sign to go down instead of up. So move in the direction \( \Delta x = -s \partial f/\partial x \) and \( \Delta y = -s \partial f/\partial y \).
The question is: How far to move? Like a boulder, a steep start may not aim directly toward the minimum. The stepsize is monitored, to end the step when the function f starts upward again (Problem 54). At the end of each step, compute first derivatives and start again in the new steepest direction.

13.6 Exercises

Read-through questions

A minimum occurs at a ______ point (where \( f_1 = f_2 = 0 \)) or a ______ point (no derivative) or a ______ point. Since \( f = x^2 - xy + 2y \) has \( f_1 = \frac{d}{dx} \) and \( f_2 = \frac{d}{dy} \), the stationary point is \( x = 1 \), \( y = 2 \). This is not a minimum, because \( f \) decreases when ______.

The minimum of \( d^2 = (x - x_1)^2 + (y - y_1)^2 \) occurs at the rough point ______. The graph of \( d \) is a ______ and grad \( d \) is a ______ vector that points ______. The graph of \( f = |xy| \) touches bottom along the lines ______. Those are "rough lines" because the derivative ______. The maximum of \( d \) and \( f \) must occur on the ______ of the allowed region because it doesn't occur ______.

When the boundary curve is \( x = x(t), y = y(t) \), the derivative of \( f(x, y) \) along the boundary is ______ (chain rule). If \( f = x^2 + 2y^2 \) and the boundary is \( x = \cos t, y = \sin t \), then \( df/dt = ______ \). It is zero at the points ______. The maximum is at ______ and the minimum is at ______. Inside the circle \( f \) has an absolute minimum at ______.

To separate maximum from minimum from ______, compute the ______ derivatives at a ______ point. The tests for a minimum are ______. The tests for a maximum are ______. In case \( ac < 0 \) or \( f_{xx} f_{yy} < 0 \), we have a ______. At all points these tests decide between concave up and ______ and "indefinite." For \( f = 8x^2 - 6xy + y^3 \), the origin is a ______. The signs of \( f \) at \((1, 0)\) and \((1, 3)\) are ______.

The Taylor series for \( f(x, y) \) begins with the six terms ______. The coefficient of \( x^ny^m \) is ______. To find a stationary point numerically, use ______ or ______.

Find all stationary points \((f_1 = f_2 = 0)\) in 1-16. Separate minimum from maximum from saddle point. Test 13K applies to \( a = f_{xx}, b = f_{xy}, c = f_{yy} \).

1. \( x^2 + 2xy + 3y^2 \)
2. \( xy - x + y \)
3. \( x^2 + 4xy + 3y^2 - 6x - 12y \)
4. \( x^2 - y^2 + 4y \)
5. \( x^2y^2 - x \)
6. \( xy^3 - e^x \)
7. \( x^2 + 2xy - 3y^2 \)
8. \( (x + y)^2 + (x + 2y - 6)^2 \)
9. \( x^2 + y^2 + z^2 - 4z \)
10. \( (x + y)(x + 2y - 6) \)
11. \( (x - y)^2 \)
12. \( (1 + x^3)(1 + y^3) \)
13. \( (x + y)^2 - (x + 2y)^2 \)
14. \( \sin x - \cos y \)
15. \( x^3 + y^3 - 3x^2 + 3y^2 \)
16. \( 8xy - x^4 - y^4 \)
17. A rectangle has sides on the \( x \) and \( y \) axes and a corner on the line \( x + 3y = 12 \). Find its maximum area.
18. A box has a corner at \((0, 0, 0)\) and all edges parallel to the axes. If the opposite corner \((x, y, z)\) is on the plane \( 3x + 2y + z = 1 \), what position gives maximum volume? Show first that the problem maximizes \( xy - 3x^2y - 2xy \).
19. (Straight line fit, Section 11.4) Find \( x \) and \( y \) to minimize the error
\[ E = (x + y)^2 + (x + 2y - 5)^2 + (x + 3y - 4)^2. \]
Show that this gives a minimum not a saddle point.
20. (Least squares) What numbers \( x, y \) come closest to satisfying the three equations \( x - y = 1, 2x + y = -1, x + 2y = 1 \)? Square and add the errors, \( (x - y - 1)^2 + (2x + y + 1)^2 + (x + 2y - 1)^2 \). Then minimize.
21. Minimize \( f = x^2 + xy + y^2 - x - y \) restricted by
   \( a) x \leq 0 \)
   \( b) y \geq 0 \)
   \( c) x \leq 0 \) and \( y \geq 0 \).
22. Minimize \( f = x^2 + y^2 + 2x + 4y \) in the regions
   \( a) \) all \( x, y \)
   \( b) y \geq 0 \)
   \( c) \) \( x \geq 0 \), \( y \geq 0 \).
23. Maximize and minimize \( f = x + \sqrt{3y} \) on the circle \( x = \cos t, y = \sin t \).
24. Example 5 followed \( f = x^2 + xy + y^2 - x - y + 1 \) around the circle \( x^2 + y^2 = 1 \). The four stationary points have \( x = y \) or \( x + y = 1 \). Compute \( f \) at those points and locate the minimum.
25. \( a) \) Maximize \( f = ax + by \) on the circle \( x^2 + y^2 = 1 \).
   \( b) \) Minimize \( f = x^2 + y^2 \) on the line \( ax + by = 1 \).
26. For \( f(x, y) = \frac{1}{2}x^4 + xy + \frac{1}{2}y^4 \), what are the equations \( f_1 = 0 \) and \( f_2 = 0 \)? What are their solutions? What is \( f_{xx} \)?
27. Choose \( c > 0 \) so that \( f = x^2 + xy + cy^2 \) has a saddle point at \((0, 0)\). Note that \( f > 0 \) on the lines \( x = 0 \) and \( y = 0 \) and \( y = x \) and \( y = -x \), so checking four directions does not confirm a minimum.

Problems 28-42 minimize the Steiner distance \( f = d_1 + d_2 + d_3 \) and related functions. A computer is needed for 33 and 36-39.

28. Draw the triangle with corners at \((0, 0), (1, 1), \) and \((1, -1)\).
By symmetry the Steiner point will be on the \( x \) axis. Write down the distances \( d_1, d_2, d_3 \) to \((x, 0)\) and find the \( x \) that minimizes \( d_1 + d_2 + d_3 \). Check the 120° angles.
29 Suppose three unit vectors add to zero. Prove that the angles between them must be 120°.

30 In three dimensions, Steiner minimizes the total distance \( f(x, y, z) = d_1 + d_2 + d_3 + d_4 \) from four points. Show that grad \( d_1 \) is still a unit vector (in which direction?) At what angles do four unit vectors add to zero?

31 With four points in a plane, the Steiner problem allows branches (Figure 13.18c). Find the shortest network connecting the corners of a rectangle, if the side lengths are (a) 1 and 2 (b) 1 and 1 (two solutions for a square) (c) 1 and 0.1.

32 Show that a Steiner point (120° angles) can never be outside the triangle.

33 Write a program to minimize \( f(x, y) = d_1 + d_2 + d_3 \) by Newton's method in equation (5). Fix two corners at (0, 0), (3, 0), vary the third from (1, 1) to (2, 1) to (3, 1) to (4, 1), and compute Steiner points.

34 Suppose one side of the triangle goes from \((-1, 0)\) to \((1, 0)\). Above that side are points from which the lines to \((-1, 0)\) and \((1, 0)\) meet at a 120° angle. Those points lie on a circular arc—draw it and find its center and its radius.

35 Continuing Problem 34, there are circular arcs for all three sides of the triangle. On the arcs, every point sees one side of the triangle at a 120° angle. Where is the Steiner point? (Sketch three sides with their arcs.)

36 Invent an algorithm to converge to the Steiner point based on Problem 35. Test it on the triangles of Problem 33.

37 Write a code to minimize \( f = d_1^2 + d_2^2 + d_3^2 \) by solving \( f_x = 0 \) and \( f_y = 0 \). Use Newton's method in equation (5).

38 Extend the code to allow all powers \( p \geq 1 \), not only \( p = 4 \). Follow the minimizing point from the centroid at \( p = 2 \) to the Steiner point at \( p = 1 \) (try \( p = 1.8, 1.6, 1.4, 1.2 \)).

39 Follow the minimizing point with your code as \( p \) increases: \( p = 2, p = 4, p = 8, p = 16 \). Guess the limit at \( p = \infty \) and test whether it is equally distant from the three corners.

40 At \( p = \infty \) we are making the largest of the distances \( d_1, d_2, d_3 \) as small as possible. The best point for a 1, 1, \( \sqrt{2} \) right triangle is ________.

41 Suppose the road from corner 1 is wider than the others, and the total cost is \( f(x, y) = \sqrt{2} d_1 + d_2 + d_3 \). Find the gradient of \( f \) and the angles at which the best roads meet.

42 Solve Steiner's problem for two points. Where is \( d_1 + d_2 \) a minimum? Solve also for three points if only the three corners are allowed.

Find all derivatives at (0, 0). Construct the Taylor series:

43 \( f(x, y) = (x + y)^3 \) \hspace{1cm} 44 \( f(x, y) = x^p \)

45 \( f(x, y) = \ln(1 - xy) \)

Find \( f_x, f_y, f_{xx}, f_{xy}, f_{yy} \) at the basepoint. Write the quadratic approximation to \( f(x, y) \)—the Taylor series through second-order terms:

46 \( f = e^{x+y} \) at \((0, 0)\) \hspace{1cm} 47 \( f = e^{x+y} \) at \((1, 1)\)

48 \( f = \sin x \cos y \) at \((0, 0)\) \hspace{1cm} 49 \( f = x^2 + y^2 \) at \((1, -1)\)

50 The Taylor series around \((x, y)\) is also written with steps \( h \) and \( k \): \( f(x + h, y + k) = f(x, y) + h \frac{\partial f}{\partial x} + k \frac{\partial f}{\partial y} + \frac{1}{2} h^2 \frac{\partial^2 f}{\partial x^2} + \frac{1}{2} k^2 \frac{\partial^2 f}{\partial y^2} + \frac{1}{2} hk \frac{\partial^2 f}{\partial x \partial y} + \cdots \). Fill in those four blanks.

51 Find lines along which \( f(x, y) \) is constant (these functions have \( f_x, f_y = f_{xy} = f_{yy} \) or \( ac = b^2 \)): \( a) f = x^2 - 4xy + 4y^2 \) \hspace{1cm} \( b) f = e^{x+y} \)

52 For \( f(x, y, z) \) the first three terms after \( f(0, 0, 0) \) in the Taylor series are ________. The next six terms are ________.

53 (a) For the error \( f - f_L \) in linear approximation, the Taylor series at \((0, 0)\) starts with the quadratic terms ________.

(b) The graph of \( f \) goes up from its tangent plane (and \( f \geq f_L \)) if ________. Then \( f \) is concave upward.

(c) For \((0, 0)\) to be a minimum we also need ________.

54 The gradient of \( x^2 + 2y^2 \) at the point \((1, 1)\) is \((2, 4)\). Steepest descent is along the line \( x = 1 - 2s, y = 1 - 4s \) (minus sign to go downward). Minimize \( x^2 + 2y^2 \) with respect to the stepsize \( s \). That locates the next point ________, where steepest descent begins again.

55 Newton's method minimizes \( x^2 + 2y^2 \) in one step. Starting at \((x_0, y_0) = (1, 1)\), find \( \Delta x \) and \( \Delta y \) from equation (5).

56 If \( f_{xx} + f_{yy} = 0 \), show that \( f(x, y) \) cannot have an interior maximum or minimum (only saddle points).

57 The value of \( x \) theorems and \( y \) exercises is \( f = x^2 y \) (maybe). The most that a student or author can deal with is 4x \( + y = 12 \). Substitute \( y = 12 - 4x \) and maximize \( f \). Show that the line \( 4x + y = 12 \) is tangent to the level curve \( x^2 y = f_{max} \).

58 The desirability of \( x \) houses and \( y \) yachts is \( f(x, y) \). The constraint \( px + qy = k \) limits the money available. The cost of a house is ________, the cost of a yacht is ________. Substitute \( y = (k - px)/q \) into \( f(x, y) = F(x) \) and use the chain rule for \( dF/dx \). Show that the slope \( -f_x/df_y \) at the best \( x \) is \(-p/q\).

59 At the farthest point in a baseball field, explain why the fence is perpendicular to the line from home plate. Assume it is not a rough point (corner) or endpoint (foul line).
This section faces up to a practical problem. We often minimize one function \( f(x, y) \) while another function \( g(x, y) \) is fixed. There is a constraint on \( x \) and \( y \), given by \( g(x, y) = k \). This restricts the material available or the funds available or the energy available. With this constraint, the problem is to do the best possible (\( f_{\text{max}} \) or \( f_{\text{min}} \)).

At the absolute minimum of \( f(x, y) \), the requirement \( g(x, y) = k \) is probably violated. In that case the minimum point is not allowed. We cannot use \( f_x = 0 \) and \( f_y = 0 \)—those equations don't account for \( g \).

**Step 1** Find equations for the constrained minimum or constrained maximum. They will involve \( f_x \) and \( f_y \) and also \( g_x \) and \( g_y \), which give local information about \( f \) and \( g \). To see the equations, look at two examples.

**EXAMPLE 1** Minimize \( f = x^2 + y^2 \) subject to the constraint \( g = 2x + y = k \).

*Trial runs* The constraint allows \( x = 0, y = k \), where \( f = k^2 \). Also \((\frac{k}{2}, 0)\) satisfies the constraint, and \( f = \frac{1}{2}k^2 \) is smaller. Also \( x = y = \frac{1}{2}k \) gives \( f = \frac{1}{8}k^2 \) (best so far).

*Idea of solution* Look at the level curves of \( f(x, y) \) in Figure 13.21. They are circles \( x^2 + y^2 = c \). When \( c \) is small, the circles do not touch the line \( 2x + y = k \). There are no points that satisfy the constraint, when \( c \) is too small. *Now increase \( c \).* Eventually the growing circles \( x^2 + y^2 = c \) will just touch the line \( 2x + y = k \). The point where they touch is the winner. It gives the smallest value of \( c \) that can be achieved on the line. The touching point is \((x_{\text{min}}, y_{\text{min}})\), and the value of \( c \) is \( f_{\text{min}} \).

What equation describes that point? When the circle touches the line, they are tangent. They have the same slope. The perpendiculars to the circle and the line go in the same direction. That is the key fact, which you see in Figure 13.21a. The direction perpendicular to \( f = c \) is given by \( \text{grad } f = (f_x, f_y) \). The direction perpendicular to \( g = k \) is given by \( \text{grad } g = (g_x, g_y) \). The key equation says that those two vectors are parallel. One gradient vector is a multiple of the other gradient vector, with a multiplier \( \lambda \) (called lambda) that is unknown:

\[
\text{grad } f = \lambda \text{ grad } g \quad \text{so} \quad \frac{\partial f}{\partial x} = \lambda \frac{\partial g}{\partial x} \quad \text{and} \quad \frac{\partial f}{\partial y} = \lambda \frac{\partial g}{\partial y} \quad \text{for} \quad \lambda \neq 0. \tag{1}
\]

**Step 2** There are now three unknowns \( x, y, \lambda \). There are also *three equations*:

\[
\frac{\partial f}{\partial x} = \lambda \frac{\partial g}{\partial x} \quad \text{is} \quad 2x = 2\lambda \tag{2}
\]
\[
\frac{\partial f}{\partial y} = \lambda \frac{\partial g}{\partial y} \quad \text{is} \quad 2y = \lambda \tag{2}
\]
\[
g(x, y) = k \quad \text{is} \quad 2x + y = k.
\]

In the third equation, substitute \( 2\lambda \) for \( 2x \) and \( \frac{1}{2}\lambda \) for \( y \). Then \( 2x + y \) equals \( \frac{1}{2}\lambda \) equals \( k \). Knowing \( \lambda = \frac{2}{5}k \), go back to the first two equations for \( x, y, \) and \( f_{\text{min}} \):

\[
x = \lambda = \frac{2}{5}k, \quad y = \frac{1}{2}\lambda = \frac{1}{5}k, \quad f_{\text{min}} = \left(\frac{2}{5}k\right)^2 + \left(\frac{1}{5}k\right)^2 = \frac{2}{25}k^2 = \frac{1}{5}k^2.
\]

The winning point \((x_{\text{min}}, y_{\text{min}})\) is \((\frac{2}{5}k, \frac{1}{5}k)\). It minimizes the "distance squared," \( f = x^2 + y^2 = \frac{1}{5}k^2 \), from the origin to the line.
Question  What is the meaning of the Lagrange multiplier λ?

Mysterious answer  The derivative of \( \frac{1}{2}k^2 \) is \( \frac{1}{2}k \), which equals λ. The multiplier λ is the derivative of \( f_{\text{min}} \) with respect to \( k \). Move the line by \( \Delta k \), and \( f_{\text{min}} \) changes by about \( \lambda \Delta k \). Thus the Lagrange multiplier measures the sensitivity to \( k \).

Pronounce his name “Lagrange” or better “Lagrongh” as if you are French.

**EXAMPLE 2** Maximize and minimize \( f = x^2 + y^2 \) on the ellipse \( g = (x-1)^2 + 4y^2 = 4 \).

Idea and equations  The circles \( x^2 + y^2 = c \) grow until they touch the ellipse. The touching point is \((x_{\text{min}}, y_{\text{min}})\) and that smallest value of \( c \) is \( f_{\text{min}} \). As the circles grow they cut through the ellipse. Finally there is a point \((x_{\text{max}}, y_{\text{max}})\) where the last circle touches. That largest value of \( c \) is \( f_{\text{max}} \).

The minimum and maximum are described by the same rule: the circle is tangent to the ellipse (Figure 13.21b). The **perpendiculars go in the same direction**. Therefore \((f_x, f_y)\) is a multiple of \((g_x, g_y)\), and the unknown multiplier is \( \lambda \):

\[
\begin{align*}
    f_x &= \lambda g_x; \\
    f_y &= \lambda g_y; \\
    g &= k; \\
    &\quad \text{(x-1)}^2 + 4y^2 = 4. \tag{3}
\end{align*}
\]

Solution  The second equation allows two possibilities: \( y = 0 \) or \( \lambda = \frac{1}{4} \). Following up \( y = 0 \), the last equation gives \( (x-1)^2 = 4 \). Thus \( x = 3 \) or \( x = -1 \). Then the first equation gives \( \lambda = 3/2 \) or \( \lambda = 1/2 \). The values of \( f \) are \( x^2 + y^2 = 3^2 + 0^2 = 9 \) and \( x^2 + y^2 = (-1)^2 + 0^2 = 1 \).

Now follow \( \lambda = 1/4 \). The first equation yields \( x = -1/3 \). Then the last equation requires \( y^2 = 5/9 \). Since \( x^2 = 1/9 \) we find \( x^2 + y^2 = 6/9 = 2/3 \). This is \( f_{\text{min}} \).

Conclusion  The equations (3) have four solutions, at which the circle and ellipse are tangent. The four points are \((3, 0), (-1, 0), (-1/3, \sqrt{5}/3), \) and \((-1/3, -\sqrt{5}/3) \). The four values of \( f \) are \( 9, 1, 3, 3/2 \).

**Summary**  The three equations are \( f_x = \lambda g_x \) and \( f_y = \lambda g_y \) and \( g = k \). The unknowns are \( x, y, \) and \( \lambda \). There is no absolute system for solving the equations (unless they are linear; then use elimination or Cramer’s Rule). Often the first two equations yield \( x \) and \( y \) in terms of \( \lambda \), and substituting into \( g = k \) gives an equation for \( \lambda \).

At the minimum, the level curve \( f(x, y) = c \) is tangent to the constraint curve \( g(x, y) = k \). If that constraint curve is given parametrically by \( x(t) \) and \( y(t) \), then
minimizing \( f(x(t), y(t)) \) uses the chain rule:
\[
\frac{df}{dt} = \frac{\partial f}{\partial x} \frac{dx}{dt} + \frac{\partial f}{\partial y} \frac{dy}{dt} = 0 \quad \text{or} \quad (\text{grad } f) \cdot (\text{tangent to curve}) = 0.
\]

This is the calculus proof that grad \( f \) is perpendicular to the curve. Thus grad \( f \) is parallel to grad \( g \). This means \((f_x, f_y) = \lambda(g_x, g_y)\).

We have lost \( f_x = 0 \) and \( f_y = 0 \). But a new function \( L \) has three zero derivatives:

\[
\frac{\partial L}{\partial x} = \frac{\partial f}{\partial x} - \lambda \frac{\partial g}{\partial x} = 0 \quad \frac{\partial L}{\partial y} = \frac{\partial f}{\partial y} - \lambda \frac{\partial g}{\partial y} = 0 \quad \frac{\partial L}{\partial \lambda} = -g + k = 0. \tag{4}
\]

Note that \( \frac{\partial L}{\partial \lambda} = 0 \) automatically produces \( g = k \). The constraint is "built in" to \( L \). Lagrange has included a term \( \lambda(g - k) \), which is destined to be zero—but its derivatives are absolutely needed in the equations! At the solution, \( g = k \) and \( L = f \) and \( \partial L / \partial k = \lambda \).

What is important is \( f_x = \lambda g_x \) and \( f_y = \lambda g_y \), coming from \( L_x = L_y = 0 \). In words: The constraint \( g = k \) forces \( dg = g_x dx + g_y dy = 0 \). This restricts the movements \( dx \) and \( dy \). They must keep to the curve. The equations say that \( df = f_x dx + f_y dy \) is equal to \( \lambda dg \). Thus \( df \) is zero in the allowed direction—which is the key point.

\[\text{MAXIMUM AND MINIMUM WITH TWO CONSTRAINTS}\]

The whole subject of min(max)imization is called optimization. Its applications to business decisions make up operations research. The special case of linear functions is always important—in this part of mathematics it is called linear programming. A book about those subjects won't fit inside a calculus book, but we can take one more step—to allow a second constraint.

The function to minimize or maximize is now \( f(x, y, z) \). The constraints are \( g(x, y, z) = k_1 \) and \( h(x, y, z) = k_2 \). The multipliers are \( \lambda_1 \) and \( \lambda_2 \). We need at least three variables \( x, y, z \) because two constraints would completely determine \( x \) and \( y \).

\[
\text{To minimize } f(x, y, z) \text{ subject to } g(x, y, z) = k_1 \text{ and } h(x, y, z) = k_2, \text{ solve five equations for } x, y, z, \lambda_1, \lambda_2. \text{ Combine } g = k_1 \text{ and } h = k_2 \text{ with}
\]
\[
\begin{align*}
\frac{\partial f}{\partial x} &= \lambda_1 \frac{\partial g}{\partial x} + \lambda_2 \frac{\partial h}{\partial x} \quad \frac{\partial f}{\partial y} &= \lambda_1 \frac{\partial g}{\partial y} + \lambda_2 \frac{\partial h}{\partial y} \quad \frac{\partial f}{\partial z} &= \lambda_1 \frac{\partial g}{\partial z} + \lambda_2 \frac{\partial h}{\partial z}.
\end{align*}
\tag{5}
\]

Figure 13.22a shows the geometry behind these equations. For convenience \( f \) is \( x^2 + y^2 + z^2 \), so we are minimizing distance (squared). The constraints \( g = x + y + z = 9 \) and \( h = x + 2y + 3z = 20 \) are linear— their graphs are planes. The constraints keep \( (x, y, z) \) on both planes—and therefore on the line where they meet. We are finding the squared distance from \((0, 0, 0)\) to a line.

What equation do we solve? The level surfaces \( x^2 + y^2 + z^2 = c \) are spheres. They grow as \( c \) increases. The first sphere to touch the line is tangent to it. That touching point gives the solution (the smallest \( c \)). All three vectors grad \( f \), grad \( g \), grad \( h \) are perpendicular to the line:

- line tangent to sphere \( \Rightarrow \) grad \( f \) perpendicular to line
- line in both planes \( \Rightarrow \) grad \( g \) and grad \( h \) perpendicular to line.
Thus \( \nabla f, \nabla g, \nabla h \) are in the same plane—perpendicular to the line. With three vectors in a plane, \( \nabla f \) is a combination of \( \nabla g \) and \( \nabla h \):

\[
(f_x, f_y, f_z) = \lambda_1(g_x, g_y, g_z) + \lambda_2(h_x, h_y, h_z).
\] (6)

This is the key equation (5). It applies to curved surfaces as well as planes.

**EXAMPLE 3** Minimize \( x^2 + y^2 + z^2 \) when \( x + y + z = 9 \) and \( x + 2y + 3z = 20 \).

In Figure 13.22b, the normals to those planes are \( \nabla g = (1, 1, 1) \) and \( \nabla h = (1, 2, 3) \). The gradient of \( f = x^2 + y^2 + z^2 \) is \( (2x, 2y, 2z) \). The equations (5)-(6) are

\[
2x = \lambda_1 + \lambda_2, \quad 2y = 2\lambda_1 + 2\lambda_2, \quad 2z = 3\lambda_1 + 3\lambda_2.
\]

Substitute these \( x, y, z \) into the other two equations \( g = x + y + z = 9 \) and \( h = 20 \):

\[
\frac{\lambda_1 + \lambda_2}{2} + \frac{\lambda_1 + 2\lambda_2}{2} + \frac{\lambda_1 + 3\lambda_2}{2} = 9 \quad \text{and} \quad \frac{\lambda_1 + \lambda_2}{2} + 2\frac{\lambda_1 + 2\lambda_2}{2} + 3\frac{\lambda_1 + 3\lambda_2}{2} = 20.
\]

After multiplying by 2, these simplify to \( 3\lambda_1 + 6\lambda_2 = 18 \) and \( 6\lambda_1 + 14\lambda_2 = 40 \). The solutions are \( \lambda_1 = 2 \) and \( \lambda_2 = 2 \). Now the previous equations give \( (x, y, z) = (2, 3, 4) \). The Lagrange function with two constraints is \( L(x, y, z, \lambda_1, \lambda_2) = f - \lambda_1(g - k_1) - \lambda_2(h - k_2) \). Its five derivatives are zero—those are our five equations. Lagrange has increased the number of unknowns from 3 to 5, by adding \( \lambda_1 \) and \( \lambda_2 \). The best point \( (2, 3, 4) \) gives \( f_{\text{min}} = 29 \). The \( \lambda \)'s give \( \partial f / \partial k \)—the sensitivity to changes in 9 and 20.

**INEQUALITY CONSTRAINTS**

In practice, applications involve *inequalities* as well as equations. The constraints might be \( g \leq k \) and \( h \geq 0 \). The first means: It is not required to use the whole resource \( k \), but you cannot use more. The second means: \( h \) measures a quantity that cannot be negative. *At the minimum point, the multipliers must satisfy the same inequalities:* \( \lambda_1 \leq 0 \) and \( \lambda_2 \geq 0 \). There are inequalities on the \( \lambda \)'s when there are inequalities in the constraints.

Brief reasoning: With \( g \leq k \) the minimum can be on or inside the constraint curve. Inside the curve, where \( g < k \), we are free to move in all directions. The constraint is not really constraining. This brings back \( f_x = 0 \) and \( f_y = 0 \) and \( \lambda = 0 \)—an ordinary minimum. On the curve, where \( g = k \) constrains the minimum from going lower, we have \( \lambda < 0 \). We don't know in advance which to expect.
For 100 constraints \( g_i \leq k_i \), there are 100 \( \lambda \)'s. Some \( \lambda \)'s are zero (when \( g_i < k_i \)) and some are nonzero (when \( g_i = k_i \)). It is those 2\(^{100}\) possibilities that make optimization interesting. In \textit{linear programming} with two variables, the constraints are \( x \geq 0, y \geq 0 \):

**Example 4** Minimize \( f = 5x + 6y \) with \( g = x + y = 4 \) and \( h = x \geq 0 \) and \( H = y \geq 0 \).

The constraint \( g = 4 \) is an equation, \( h \) and \( H \) yield inequalities. Each has its own Lagrange multiplier—and the inequalities require \( \lambda_2 \geq 0 \) and \( \lambda_3 \geq 0 \). The derivatives of \( f, g, h, H \) are no problem to compute:

\[
\begin{align*}
\frac{\partial f}{\partial x} &= \lambda_1 \frac{\partial g}{\partial x} + \lambda_2 \frac{\partial h}{\partial x} + \lambda_3 \frac{\partial H}{\partial x} \quad \text{yields} \quad 5 = \lambda_1 + \lambda_2 \\
\frac{\partial f}{\partial y} &= \lambda_1 \frac{\partial g}{\partial y} + \lambda_2 \frac{\partial h}{\partial y} + \lambda_3 \frac{\partial H}{\partial y} \quad \text{yields} \quad 6 = \lambda_1 + \lambda_3.
\end{align*}
\]

Those equations make \( \lambda_3 \) larger than \( \lambda_2 \). Therefore \( \lambda_3 > 0 \), which means that the constraint on \( H \) must be an equation. (Inequality for the multiplier means equality for the constraint.) In other words \( H = y = 0 \). Then \( x + y = 4 \) leads to \( x = 4 \). The solution is at \( (x_{\min}, y_{\min}) = (4, 0) \), where \( f_{\min} = 20 \).

At this minimum, \( h = x = 4 \) is above zero. The multiplier for the constraint \( h \geq 0 \) must be \( \lambda_2 = 0 \). Then the first equation gives \( \lambda_1 = 5 \). As always, the multiplier measures sensitivity. When \( g = 4 \) is increased by \( \Delta k \), the cost \( f_{\min} = 20 \) is increased by \( 5\Delta k \). In economics \( \lambda_1 = 5 \) is called a \textit{shadow price}—it is the cost of increasing the constraint.

Behind this example is a nice problem in geometry. The constraint curve \( x + y = 4 \) is a line. The inequalities \( x \geq 0 \) and \( y \geq 0 \) leave a piece of that line—from \( P \) to \( Q \) in Figure 13.23. The level curves \( f = 5x + 6y = c \) move out as \( c \) increases, until they touch the line. The first touching point is \( Q = (4, 0) \), which is the solution. It is always an \textit{endpoint}—or a corner of the triangle \( PQR \). It gives the smallest cost \( f_{\min} \), which is \( c = 20 \).

**Fig. 13.23** Linear programming: \( f \) and \( g \) are linear, inequalities cut off \( x \) and \( y \).
13.7 Constraints and Lagrange Multipliers

To minimize \( f = x^2 - y \) subject to \( g = x - y = 0 \), the three equations for \( x, y, \lambda \) are \( \frac{\partial f}{\partial x} = \frac{\partial f}{\partial y} = \frac{\partial f}{\partial \lambda} = 0 \). The solution is \( (x, y, \lambda) \). In this example the curve \( f(x, y) = f_{\text{min}} \) is \( \frac{\partial f}{\partial y} = 0 \) which is \( \frac{\partial f}{\partial x} \) to the level \( g = 0 \) at \((x_{\text{min}}, y_{\text{min}})\).

With two constraints \( g(x, y, z) = k_1 \) and \( h(x, y, z) = k_2 \) there are \( \lambda \) multipliers. The five unknowns are \( x, y, \lambda, k_1, \text{and} k_2 \). The five equations are \( \lambda \)!. The level surface \( f = f_{\text{min}} \) is \( \nabla_f \) to the curve where \( g = k_1 \) and \( h = k_2 \). Then grad \( f \) is \( \nabla_f \) to this curve, and so are grad \( g \) and \( \nabla_h \). Thus \( x \) is a combination of grad \( g \) and \( \nabla_h \). With nine variables and six constraints, there will be \( \frac{9}{3} \) multipliers and eventually \( \frac{15}{3} \) equations. If a constraint is an \( \frac{3}{3} \), \( g \leq k \), then its multiplier must satisfy \( \lambda \leq 0 \) at a minimum.

1 Example 1 minimized \( f = x^2 + y^2 \) subject to \( 2x + y = k \). Solve the constraint equation for \( y = k - 2x \), substitute into \( f \), and minimize this function of \( x \). The minimum is at \((x, y) = \ldots \), where \( f = \ldots \).

Note: This direct approach reduces to one unknown \( x \). Lagrange increases to \( x, y, \lambda \). But Lagrange is better when the first step of solving for \( y \) is difficult or impossible.

Minimize and maximize \( f(x, y) \) in 2–6. Find \( x, y, \) and \( \lambda \).

2 \( f = x^2 y \) with \( g = x^2 + y^2 = 1 \)

3 \( f = x + y \) with \( g = \frac{1}{x} + \frac{1}{y} = 1 \)

4 \( f = 3x + y \) with \( g = x^2 + 9y^2 = 1 \)

5 \( f = x^2 + y^2 \) with \( g = x^6 + y^6 = 2 \).

6 \( f = x + y \) with \( g = x^{1/3} y^{2/3} = k \). With \( x = \text{capital} \) and \( y = \text{labor} \), \( g \) is a Cobb-Douglas function in economics. Draw two of its level curves.

7 Find the point on the circle \( x^2 + y^2 = 13 \) where \( f = 2x - 3y \) is a maximum. Explain the answer.

8 Maximize \( ax + by + cz \) subject to \( x^2 + y^2 + z^2 = k^2 \). Write your answer as the Schwarz inequality for dot products: \( (a, b, c) \cdot (x, y, z) \leq \ldots k \).

9 Find the plane \( z = ax + by + c \) that best fits the points \((x, y, z) = (0, 0, 1, 1, 0, 1, 1, 2, 0, 1, 2) \). The answer \( a, b, c \) minimizes the sum of \((z - ax - by - c)^2 \) at the four points.

10 The base of a triangle is the top of a rectangle (5 sides, combined area = 1). What dimensions minimize the distance around?

11 Draw the hyperbola \( xy = -1 \) touching the circle \( g = x^2 + y^2 = 2 \). The minimum of \( f = xy \) on the circle is reached at the points \ldots \). The equations \( f_x = \lambda g_x \) and \( f_y = \lambda g_y \) are satisfied at those points with \( \lambda = \ldots \).

12 Find the maximum of \( f = xy \) on the circle \( g = x^2 + y^2 = 2 \) by solving \( f_x = \lambda g_x \) and \( f_y = \lambda g_y \) and substituting \( x \) and \( y \) into \( f \). Draw the level curve \( f = f_{\text{max}} \) that touches the circle.

13 Draw the level curves of \( f = x^2 + y^2 \) with a closed curve \( C \) across them to represent \( g(x, y) = \). Mark a point where \( C \) crosses a level curve. Why is that point not a minimum of \( f \) on \( C \)? Mark a point where \( C \) is tangent to a level curve. Is that the minimum of \( f \) on \( C \)?

14 On the circle \( g = x^2 + y^2 = 1 \), Example 5 of 13.6 minimized \( f = xy - x - y \). Example 5 of 13.6 minimized \( f = x^2 + y^2 = 1 \). Example 5 of 13.6 minimized \( f = xy - x - y \).

(a) Set up the three Lagrange equations for \( x, y, \lambda \). (b) The first two equations give \( x = y = \ldots \).

(c) There is another solution for the special value \( \lambda = -\frac{1}{k} \), when the equations become \( \ldots \). This is easy to miss but it gives \( f_{\text{min}} = -1 \) at the point \( \ldots \).

Problems 15–18 develop the theory of Lagrange multipliers.

15 (Sensitivity) Certainly \( L = f - \lambda (g - k) \) has \( \frac{\partial L}{\partial \lambda} = \lambda \). Since \( L = f_{\text{min}} \) and \( g = k \) at the minimum point, this seems to prove the key formula \( \frac{\partial f_{\text{min}}}{\partial \lambda} = \lambda \). But \( f_{\text{min}}, f_{\text{max}}, \lambda, \text{and} f_{\text{min}} \) all change with \( k \). We need the total derivative of \( L(x, y, \lambda, k) \):

\[
\frac{dL}{dk} = \frac{\partial L}{\partial k} + \frac{\partial L}{\partial x} \frac{dx}{dk} + \frac{\partial L}{\partial y} \frac{dy}{dk} + \frac{\partial L}{\partial \lambda} \frac{d\lambda}{dk} + \frac{\partial L}{d\lambda} \frac{d\lambda}{dk}.
\]

Equation (1) at the minimum point should now yield the sensitivity formula \( \frac{\partial f_{\text{min}}}{\partial \lambda} = \lambda \).

16 (Theory behind \( \lambda \)) When \( g(x, y) = k \) is solved for \( y \), it gives a curve \( y = R(x) \). Then minimizing \( f(x, y) \) along this curve yields

\[
\frac{\partial f}{\partial x} + \frac{\partial f}{\partial y} \frac{dy}{dx} = \frac{\partial f}{\partial x} + \frac{\partial f}{\partial y} \frac{dR}{dx} = 0.
\]

Those come from the \( \frac{\partial f}{\partial y} \frac{dy}{dx} = 0 \) at the minimum and \( \frac{\partial g}{\partial x} \frac{dx}{dx} = 0 \) along the curve because \( g = \ldots \).

Multiplying the second equation by \( \lambda = \frac{\partial f}{\partial y} \frac{dR}{dx} \) and subtracting from the first gives \( \frac{\partial f}{\partial x} = 0 \). Also \( \frac{\partial f}{\partial y} = \lambda \frac{\partial g}{\partial y} \). These are the equations (1) for \( x, y, \lambda \).

17 (Example of failure) \( \lambda = f_y / g_y \) breaks down if \( g_y = 0 \) at the minimum point.

(a) \( g = x^2 - y^2 = 0 \) does not allow negative \( y \) because \( \ldots \).

(b) When \( g = 0 \) the minimum of \( f = x^2 + y \) is at the point \( \ldots \).

(c) At that point \( f_y = \lambda g_y \) becomes \( \ldots \) which is impossible.

(d) Draw the pointed curve \( g = 0 \) to see why it is not tangent to a level curve of \( f \).

18 (No maximum) Find a point on the line \( g = x + y = 1 \) where \( f(x, y) = 2x + y \) is greater than 100 (or 1000). Write out \( \frac{\partial f}{\partial x} = \lambda \frac{\partial g}{\partial x} \) to see that there is no solution.

19 Find the minimum of \( f = x^2 + 2y^2 + z^2 \) if \( (x, y, z) \) is restricted to the planes \( g = x + y + z = 0 \) and \( h = x - z = 1 \).

20 (a) Find by Lagrange multipliers the volume \( V = xyz \) of the largest box with sides adding up to \( x + y + z = k \). (b) Check that \( \lambda = dV_{\text{max}} / dk \). (c) United Airlines accepts baggage with \( x + y + z = 108^\circ \). If it changes to \( 111^\circ \), approximately how much (by \( \lambda \)) and exactly how much does \( V_{\text{max}} \) increase?
21 The planes \( x = 0 \) and \( y = 0 \) intersect in the line \( x = y = 0 \), which is the \( z \) axis. Write down a vector perpendicular to the plane \( x = 0 \) and a vector perpendicular to the plane \( y = 0 \). Find \( \lambda_1 \) times the first vector plus \( \lambda_2 \) times the second. This combination is perpendicular to the line ________.

22 Minimize \( f = x^2 + y^2 + z^2 \) on the plane \( ax + by + cz = d \) — one constraint and one multiplier. Compare \( f_{\text{min}} \) with the distance formula \( |d|/\sqrt{a^2 + b^2 + c^2} \) in Section 11.2.

23 At the absolute minimum of \( f(x, y) \), the derivatives ________ are zero. If this point happens to fall on the curve \( g(x, y) = k \) then the equations \( f_x = \lambda g_x \) and \( f_y = \lambda g_y \) hold with \( \lambda = ________ \).

Problems 24–33 allow inequality constraints, optional but good.

24 Find the minimum of \( f = 3x + 5y \) with the constraints \( g = x + 2y = 4 \) and \( h = x \geq 0 \) and \( H = y \geq 0 \), using equations like (7). Which multiplier is zero?

25 Figure 13.23 shows the constraint plane \( g = x + y + z = 1 \) chopped off by the inequalities \( x \geq 0, y \geq 0, z \geq 0 \). What are the three "endpoints" of this triangle? Find the minimum and maximum of \( f = 4x - 2y + 5z \) on the triangle, by testing \( f \) at the endpoints.

26 With an inequality constraint \( g \leq k \), the multiplier at the minimum satisfies \( \lambda \leq 0 \). If \( k \) is increased, \( f_{\text{min}} \) goes down (since \( \lambda = df_{\text{min}}/dk \)). Explain the reasoning: By increasing \( k \), (more) (fewer) points satisfy the constraints. Therefore (more) (fewer) points are available to minimize \( f \). Therefore \( f_{\text{min}} \) goes (up) (down).

27 With an inequality constraint \( g \leq k \), the multiplier at a maximum point satisfies \( \lambda > 0 \). Change the reasoning in 26.

28 When the constraint \( h \geq k \) is a strict inequality \( h > k \) at the minimum, the multiplier is \( \lambda = 0 \). Explain the reasoning: For a small increase in \( k \), the same minimizer is still available (since \( h > k \) leaves room to move). Therefore \( f_{\text{min}} \) is (changed)(unchanged), and \( \lambda = df_{\text{min}}/dk \) is ________.

29 Minimize \( f = x^2 + y^2 \) subject to the inequality constraint \( x + y \leq 4 \). The minimum is obviously at ________, where \( f_x \) and \( f_y \) are zero. The multiplier is \( \lambda = ________ \). A small change from 4 will leave \( f_{\text{min}} = ________ \) so the sensitivity \( df_{\text{min}}/dk \) still equals \( \lambda \).

30 Minimize \( f = x^2 + y^2 \) subject to the inequality constraint \( x + y \geq 4 \). Now the minimum is at ________, where \( f_x \) and \( f_y \) are zero. The multiplier is \( \lambda = ________ \) and \( f_{\text{min}} = ________ \). A small change to \( 4 + dk \) changes \( f_{\text{min}} \), what multiple of \( dk \)?

31 Minimize \( f = 5x + 6y \) with \( g = x + y = 4 \) and \( h = x \geq 0 \) and \( H = y \leq 0 \). Now \( \lambda_3 \leq 0 \) and the sign change destroys Example 4. Show that equation (7) has no solution, and choose \( x, y \) to make \( 5x + 6y < -1000 \).

32 Minimize \( f = 2x + 3y + 4z \) subject to \( g = x + y + z = 1 \) and \( x, y, z \geq 0 \). These constraints have multipliers \( \lambda_2 \geq 0, \lambda_3 \geq 0, \lambda_4 \geq 0 \). The equations are \( 2 = \lambda_1 + \lambda_2 \), ________, and \( 4 = \lambda_1 + \lambda_4 \). Explain why \( \lambda_2 > 0 \) and \( \lambda_4 > 0 \) and \( f_{\text{min}} = 2 \).

33 A wire 40" long is used to enclose one or two squares (side \( x \) and side \( y \)). Maximize the total area \( x^2 + y^2 \) subject to \( x \geq 0, y \geq 0 \), \( 4x + 4y = 40 \).
CHAPTER 14  

Multiple Integrals

14.1 Double Integrals  
14.2 Changing to Better Coordinates  
14.3 Triple Integrals  
14.4 Cylindrical and Spherical Coordinates

CHAPTER 15  

Vector Calculus

15.1 Vector Fields  
15.2 Line Integrals  
15.3 Green’s Theorem  
15.4 Surface Integrals  
15.5 The Divergence Theorem  
15.6 Stokes’ Theorem and the Curl of F

CHAPTER 16  

Mathematics after Calculus

16.1 Linear Algebra  
16.2 Differential Equations  
16.3 Discrete Mathematics  

Study Guide For Chapter 1  
Answers to Odd-Numbered Problems  
Index  
Table of Integrals
CHAPTER 14

Multiple Integrals

14.1 Double Integrals

This chapter shows how to integrate functions of two or more variables. First, a double integral is defined as the limit of sums. Second, we find a fast way to compute it. The key idea is to replace a double integral by two ordinary “single” integrals.

The double integral \( \iint f(x, y) \, dy \, dx \) starts with \( \int f(x, y) \, dy \). For each fixed \( x \) we integrate with respect to \( y \). The answer depends on \( x \). Now integrate again, this time with respect to \( x \). The limits of integration need care and attention! Frequently those limits on \( y \) and \( x \) are the hardest part.

Why bother with sums and limits in the first place? Two reasons. There has to be a definition and a computation to fall back on, when the single integrals are difficult or impossible. And also — this we emphasize — multiple integrals represent more than area and volume. Those words and the pictures that go with them are the easiest to understand. You can almost see the volume as a “sum of slices” or a “double sum of thin sticks.” The true applications are mostly to other things, but the central idea is always the same: Add up small pieces and take limits.

We begin with the area of \( R \) and the volume of \( V \) by double integrals.

A LIMIT OF SUMS

The graph of \( z = f(x, y) \) is a curved surface above the \( xy \) plane. At the point \((x, y)\) in the plane, the height of the surface is \( z \). (The surface is above the \( xy \) plane only when \( z \) is positive. Volumes below the plane come with minus signs, like areas below the \( x \) axis.) We begin by choosing a positive function — for example \( z = 1 + x^2 + y^2 \).

The base of our solid is a region \( R \) in the \( xy \) plane. That region will be chopped into small rectangles (sides \( \Delta x \) and \( \Delta y \)). When \( R \) itself is the rectangle \( 0 \leq x \leq 1, \, 0 \leq y \leq 2 \), the small pieces fit perfectly. For a triangle or a circle, the rectangles miss part of \( R \). But they do fit in the limit, and any region with a piecewise smooth boundary will be acceptable.

Question: What is the volume above \( R \) and below the graph of \( z = f(x, y) \)?

Answer: It is a double integral — the integral of \( f(x, y) \) over \( R \). To reach it we begin with a sum, as suggested by Figure 14.1.
For single integrals, the interval \([a, b]\) is divided into short pieces of length \(\Delta x\). For double integrals, \(R\) is divided into small rectangles of area \(\Delta A = (\Delta x)(\Delta y)\). Above the \(i\)th rectangle is a "thin stick" with small volume. That volume is the base area \(\Delta A\) times the height above it—except that this height \(z = f(x, y)\) varies from point to point. Therefore we select a point \((x_i, y_i)\) in the \(i\)th rectangle, and compute the volume from the height above that point:

\[
\text{volume of one stick} = f(x_i, y_i)\Delta A \quad \text{volume of all sticks} = \sum f(x_i, y_i)\Delta A.
\]

This is the crucial step for any integral—to see it as a sum of small pieces.

Now take limits: \(\Delta x \to 0\) and \(\Delta y \to 0\). The height \(z = f(x, y)\) is nearly constant over each rectangle. (We assume that \(f\) is a continuous function.) The sum approaches a limit, which depends only on the base \(R\) and the surface above it. The limit is the volume of the solid, and it is the double integral of \(f(x, y)\) over \(R\):

\[
\iint_R f(x, y) \, dA = \lim_{\Delta x \to 0} \lim_{\Delta y \to 0} \sum f(x_i, y_i)\Delta A.
\]

To repeat: The limit is the same for all choices of the rectangles and the points \((x_i, y_i)\). The rectangles will not fit exactly into \(R\), if that base area is curved. The heights are not exact, if the surface \(z = f(x, y)\) is also curved. But the errors on the sides and top, where the pieces don’t fit and the heights are wrong, approach zero. Those errors are the volume of the "icing" around the solid, which gets thinner as \(\Delta x \to 0\) and \(\Delta y \to 0\). A careful proof takes more space than we are willing to give. But the properties of the integral need and deserve attention:

1. **Linearity:** \(\iint (f + g) \, dA = \iint f \, dA + \iint g \, dA\)
2. **Constant comes outside:** \(\iint cf(x, y) \, dA = c \iint f(x, y) \, dA\)
3. **\(R\) splits into \(S\) and \(T\) (not overlapping):** \(\iint_S f \, dA = \iint_S f \, dA + \iint_T f \, dA\).

In 1 the volume under \(f + g\) has two parts. The "thin sticks" of height \(f + g\) split into thin sticks under \(f\) and under \(g\). In 2 the whole volume is stretched upward by \(c\). In 3 the volumes are side by side. As with single integrals, these properties help in computations.

By writing \(dA\), we allow shapes other than rectangles. Polar coordinates have an extra factor \(r\) in \(dA = r \, dr \, d\theta\). By writing \(dx \, dy\), we choose rectangular coordinates and prepare for the splitting that comes now.
14.1 Double Integrals

SPLITTING A DOUBLE INTEGRAL INTO TWO SINGLE INTEGRALS

The double integral \( \iint f(x, y) \, dy \, dx \) will now be reduced to single integrals in \( y \) and then \( x \). (Or vice versa. Our first integral could equally well be \( \int f(x, y) \, dx \).) Chapter 8 described the same idea for solids of revolution. First came the area of a slice, which is a single integral. Then came a second integral to add up the slices. For solids formed by revolving a curve, all slices are circular disks—now we expect other shapes.

Figure 14.2 shows a slice of area \( A(x) \). It cuts through the solid at a fixed value of \( x \). The cut starts at \( y = c \) on one side of \( R \), and ends at \( y = d \) on the other side. This particular example goes from \( y = 0 \) to \( y = 2 \) (\( R \) is a rectangle). The area of a slice is the \( y \) integral of \( f(x, y) \). Remember that \( x \) is fixed and \( y \) goes from \( c \) to \( d \):

\[
A(x) = \text{area of slice} = \int_c^d f(x, y) \, dy \quad \text{(the answer is a function of } x)\.
\]

**EXAMPLE 1**

\[
A = \int_{y=0}^{y=2} (1 + x^2 + y^2) \, dy = \left[ y + x^2y + \frac{y^3}{3} \right]_{y=0}^{y=2} = 2 + 2x^2 + \frac{8}{3}.
\]

This is the reverse of a partial derivative! The integral of \( x^2 \, dy \), with \( x \) constant, is \( x^2y \). This "partial integral" is actually called an *inner integral*. After substituting the limits \( y = 2 \) and \( y = 0 \) and subtracting, we have the area \( A(x) = 2 + 2x^2 + \frac{8}{3} \). Now the outer integral adds slices to find the volume \( \int A(x) \, dx \). The answer is a number:

\[
\text{volume} = \int_{x=0}^{x=1} \left( 2 + 2x^2 + \frac{8}{3} \right) \, dx = \left[ 2x + \frac{2x^3}{3} + \frac{8}{3}x \right]_0^1 = 2 + \frac{2}{3} + \frac{8}{3} = \frac{16}{3}.
\]

![Fig. 14.2](image)

A slice of \( V \) at a fixed \( x \) has area \( A(x) = \int f(x, y) \, dy \).

To complete this example, check the volume when the \( x \) integral comes first:

inner integral = \[
\int_{x=0}^{x=1} (1 + x^2 + y^2) \, dx = \left[ x + \frac{1}{3}x^3 + y^2x \right]_{x=0}^{x=1} = \frac{4}{3} + y^2
\]

outer integral = \[
\int_{y=0}^{y=2} \left( \frac{4}{3} + y^2 \right) \, dy = \left[ \frac{4}{3}y + \frac{1}{3}y^3 \right]_{y=0}^{y=2} = \frac{8}{3} + \frac{8}{3} = \frac{16}{3}.
\]

The fact that double integrals can be split into single integrals is *Fubini's Theorem*.

14A If \( f(x, y) \) is continuous on the rectangle \( R \), then

\[
\iint_R f(x, y) \, dA = \int_a^b \left[ \int_c^d f(x, y) \, dy \right] \, dx = \int_c^d \left[ \int_a^b f(x, y) \, dx \right] \, dy. \quad (2)
\]
The inner integrals are the cross-sectional areas $A(x)$ and $a(y)$ of the slices. The outer integrals add up the volumes $A(x)dx$ and $a(y)dy$. Notice the reversing of limits.

Normally the brackets in (2) are omitted. When the $y$ integral is first, $dy$ is written inside $dx$. The limits on $y$ are inside too. I strongly recommend that you compute the inner integral on one line and the outer integral on a separate line.

**EXAMPLE 2** Find the volume below the plane $z = x - 2y$ and above the base triangle $R$.

The triangle $R$ has sides on the $x$ and $y$ axes and the line $x + y = 1$. The strips in the $y$ direction have varying lengths. (So do the strips in the $x$ direction.) This is the main point of the example—the base is not a rectangle. The upper limit on the inner integral changes as $x$ changes. The top of the triangle is at $y = 1 - x$.

Figure 14.3 shows the strips. The region should always be drawn (except for rectangles). Without a figure the limits are hard to find. A sketch of $R$ makes it easy:

$y$ goes from $c = 0$ to $d = 1 - x$. Then $x$ goes from $a = 0$ to $b = 1$.

The inner integral has variable limits and the outer integral has constant limits:

inner: $\int_{y=0}^{y=1-x} (x - 2y)dy = \left[ xy - y^2 \right]_{y=0}^{y=1-x} = x(1-x) - (1-x)^2 = -1 + 3x - 2x^2$

outer: $\int_{x=0}^{x=1} (-1 + 3x - 2x^2)dx = \left[ -x + \frac{3}{2}x^2 - \frac{2}{3}x^3 \right]_{0}^{1} = -1 + \frac{3}{2} - \frac{2}{3} = -\frac{1}{6}$

The volume is negative. Most of the solid is below the $xy$ plane. To check the answer $-\frac{1}{6}$, do the $x$ integral first: $x$ goes from 0 to $1 - y$. Then $y$ goes from 0 to 1.

inner: $\int_{x=0}^{x=1} (x - 2y)dx = \left[ \frac{1}{2}x^2 - 2xy \right]_{0}^{1} = \frac{1}{2}(1-y)^2 - 2(1-y)y = \frac{1}{2} - 3y + \frac{5}{2}y^2$

outer: $\int_{y=0}^{y=1} \left( \frac{1}{2} - 3y + \frac{5}{2}y^2 \right)dy = \left[ \frac{1}{2}y - \frac{3}{2}y^2 + \frac{5}{6}y^3 \right]_{0}^{1} = \frac{1}{2} - \frac{3}{2} + \frac{5}{6} = -\frac{1}{6}$

Same answer, very probably right. The next example computes $\int \int 1 \ dy \ dx = \text{area of } R$.

**EXAMPLE 3** The area of $R$ is $\int_{x=0}^{x=1} \int_{y=0}^{y=1-x} dy \ dx$ and also $\int_{y=0}^{y=1} \int_{x=0}^{x=1-y} dx \ dy$.

The first has vertical strips. The inner integral equals $1 - x$. Then the outer integral (of $1 - x$) has limits 0 and 1, and the area is $\frac{1}{2}$. It is like an indefinite integral inside a definite integral.

![Fig. 14.3 Thin sticks above and below (Example 2). Reversed order (Examples 3 and 4).](image-url)
EXAMPLE 4  Reverse the order of integration in \( \int_{x=0}^{2} \int_{y=x^2}^{2x} x^3 \, dy \, dx \).

Solution  Draw a figure! The inner integral goes from the parabola \( y = x^2 \) up to the straight line \( y = 2x \). This gives vertical strips. The strips sit side by side between \( x = 0 \) and \( x = 2 \). They stop where \( 2x \) equals \( x^2 \), and the line meets the parabola.

The problem is to put the \( x \) integral first. It goes along horizontal strips. On each line \( y = \) constant, we need the entry value of \( x \) and the exit value of \( x \). From the figure, \( x \) goes from \( \frac{1}{2}y \) to \( \sqrt{y} \). Those are the inner limits. Pay attention also to the outer limits, because they now apply to \( y \). The region starts at \( y = 0 \) and ends at \( y = 4 \).

No change in the integrand \( x^3 \)—that is the height of the solid:

\[
\int_{x=0}^{2} \int_{y=x^2}^{2x} x^3 \, dy \, dx \quad \text{is reversed to} \quad 
\int_{y=0}^{4} \int_{x=\frac{1}{2}y}^{\sqrt{y}} x^3 \, dx \, dy. \tag{3}
\]

EXAMPLE 5  Find the volume bounded by the planes \( x = 0, y = 0, z = 0, \) and \( 2x + y + z = 4 \).

Solution  The solid is a tetrahedron (four sides). It goes from \( z = 0 \) (the \( xy \) plane) up to the plane \( 2x + y + z = 4 \). On that plane \( z = 4 - 2x - y \). This is the height function \( f(x, y) \) to be integrated.

Figure 14.4 shows the base \( R \). To find its sides, set \( z = 0 \). The sides of \( R \) are the lines \( x = 0 \) and \( y = 0 \) and \( 2x + y = 4 \). Taking vertical strips, \( dy \) is inner:

\[
\text{inner: } \int_{y=0}^{4-2x} (4 - 2x - y) \, dy = \left[ (4 - 2x)y - \frac{1}{2}y^2 \right]_0^{4-2x} = \frac{1}{2} (4 - 2x)^2
\]

\[
\text{outer: } \int_{x=0}^{2} \frac{1}{2} (4 - 2x)^2 \, dx = \left[ \frac{(4 - 2x)^3}{2 \cdot 3 \cdot 2} \right]_0^2 = \frac{4^3}{2 \cdot 3 \cdot 2} = \frac{16}{3}
\]

Question  What is the meaning of the inner integral \( \frac{1}{2} (4 - 2x)^2 \) \( (\text{and also } \frac{16}{3}) \)?

Answer  The first is \( A(x) \), the area of the slice. \( \frac{16}{3} \) is the solid volume.

Question  What if the inner integral \( \int f(x, y) \, dy \) has limits that depend on \( y \)?

Answer  It can't. Those limits must be wrong. Find them again.

![Diagram of a tetrahedron, semicircle, and triangle](Fig. 14.4 Tetrahedron in Example 5, semicircle in Example 6, triangle in Example 7.)

EXAMPLE 6  Find the mass in a semicircle \( 0 \leq y \leq \sqrt{1-x^2} \) if the density is \( \rho = y \).

This is a new application of double integrals. The total mass is a sum of small masses \( (\rho \text{ times } \Delta A) \) in rectangles of area \( \Delta A \). The rectangles don't fit perfectly inside the semicircle \( R \), and the density is not constant in each rectangle—but those problems
14 Multiple Integrals

disappear in the limit. We are left with a double integral:

\[ \text{total mass } M = \iint_R \rho \, dA = \int \int_R \rho(x, y) \, dx \, dy. \]  

Set \( \rho = y \). Figure 14.4 shows the limits on \( x \) and \( y \) (try both \( dy \, dx \) and \( dx \, dy \)):

\[
\begin{align*}
\text{mass } M &= \int_{x=1}^1 \int_{y=0}^{1-x^2} y \, dy \, dx \quad \text{and also} \quad M = \int_{y=0}^1 \int_{x=1-y^2}^{1-x^2} y \, dx \, dy.
\end{align*}
\]

The first inner integral is \( \frac{1}{2}y^2 \). Substituting the limits gives \( \frac{1}{2}(1-x^2) \). The outer integral of \( (1-x^2) \) yields the total mass \( M = 3 \).

The second inner integral is \( xy \). Substituting the limits on \( x \) gives \( \text{ } \). Then the outer integral is \( -(1-y^2)^{3/2} \). Substituting \( y = 1 \) and \( y = 0 \) yields \( M = \text{ } \).

Remark This same calculation also produces the moment around the \( x \) axis, when the density is \( \rho = 1 \). The factor \( y \) is the distance to the \( x \) axis. The moment is \( M_x = \text{ } \).

Dividing by the area of the semicircle (which is \( \pi/2 \)) locating the centroid:

\[
\bar{x} = 0 \text{ by symmetry and } \bar{y} = \text{height of centroid} = \frac{\text{moment}}{\text{area}} = \frac{2/3}{\pi/2} = \frac{4}{3\pi}.
\]  

This is the “average height” of points inside the semicircle, found earlier in 8.5.

EXAMPLE 7 Integrate \( \int_{y=0}^1 \int_{x=0}^{x^2} \cos x^2 \, dx \, dy \) avoiding the impossible \( \int \cos x^2 \, dx \).

This is a famous example where reversing the order makes the calculation possible. The base \( R \) is the triangle in Figure 14.4 (note that \( x \) goes from \( y \) to 1). In the opposite order \( y \) goes from 0 to \( x \). Then \( \int \cos x^2 \, dy = x \cos x^2 \) contains the factor \( x \) that we need:

outer integral: \( \int_0^1 x \cos x^2 \, dx = \left[ \frac{1}{2} \sin x^2 \right]_0^1 = \frac{1}{2} \sin 1 \).

14.1 EXERCISES

Read-through questions

The double integral \( \iint_R f(x, y) \, dA \) gives the volume between \( R \) and \( \text{ } \). The base is first cut into small \( \text{ } \) of area \( \Delta A \). The volume above the \( i \)th piece is approximately \( \text{ } \). The limit of the sum \( \text{ } \) is the volume integral. Three properties of double integrals are \( \text{ } \) (linearity) and \( \text{ } \) and \( \text{ } \). If \( R \) is the rectangle \( 0 \leq x \leq 4, 4 \leq y \leq 6 \), the integral \( \iint_R x \, dA \) can be computed two ways. One is \( \int_y^1 \int_x^0 x \, dy \, dx \), when the inner integral is \( \text{ } \). The outer integral gives \( \frac{1}{2} \). When the \( x \) integral comes first it equals \( \int x \, dx = \frac{1}{2} \). Then the \( y \) integral equals \( \text{ } \). This is the volume between \( \text{ } \) (describe \( V \)).

The area of \( R \) is \( \iint \text{ } \, dy \, dx \). When \( R \) is the triangle between \( x = 0, y = 2x \), and \( y = 1 \), the inner limits on \( y \) are \( \text{ } \). This is the length of a \( \text{ } \) strip. The (outer) limits on \( x \) are \( \text{ } \). The area is \( \text{ } \). In the opposite order, the (inner) limits on \( x \) are \( \text{ } \). Now the strip is \( \text{ } \) and the outer integral is \( \text{ } \). When the density is \( \rho(x, y) \), the total mass in the region \( R \) is \( \iint \text{ } \). The moments are \( M_y = \text{ } \) and \( M_x = \text{ } \). The centroid has \( \bar{x} = M_y / M \).

Compute the double integrals 1–4 by two integrations.

1. \( \int_{y=0}^1 \int_{x=0}^2 x^2 \, dx \, dy \) and \( \int_{x=0}^1 \int_{y=0}^2 y^2 \, dx \, dy \)

2. \( \int_{x=0}^2 \int_{y=0}^c 2xy \, dx \, dy \) and \( \int_{y=0}^e \int_{x=1}^{2-e} \, dx \, dy \)

3. \( \int_{x=0}^2 \int_{y=0}^n \sin(x + y) \, dx \, dy \) and \( \int_{y=0}^3 \int_{x=2}^{2-x} \, dy \, dx \)

4. \( \int_{y=0}^1 \int_{x=0}^2 ye^x \, dx \, dy \) and \( \int_{y=0}^1 \int_{x=0}^1 \sqrt{3 + 2x + y} \, dx \, dy \)
14.2 Change to Better Coordinates

You don't go far with double integrals before wanting to change variables. Many regions simply do not fit with the x and y axes. Two examples are in Figure 14.5, a tilted square and a ring. Those are excellent shapes—in the right coordinates.

In 5–10, draw the region and compute the area.

\[
5 \int_{x=1}^{2} \int_{y=1}^{2x} dy \, dx \\
6 \int_{0}^{1} \int_{x=1}^{x} dy \, dx
\]

\[
7 \int_{0}^{1} \int_{y=2x}^{e^{-x}} dy \, dx \\
8 \int_{-1}^{1} \int_{x=-2x}^{1-x^2} dy \, dx
\]

\[
9 \int_{-1}^{1} \int_{y^2}^{1} dx \, dy \\
10 \int_{-1}^{1} \int_{x=-y}^{1} dy \, dx
\]

In 11–16 reverse the order of integration (and find the new limits) in 5–10 respectively.

In 17–24 find the limits on \( \int \int dy \, dx \) and \( \int \int dx \, dy \). Draw \( R \) and compute its area.

17 \( R \) = triangle inside the lines \( x = 0, y = 1, y = 2x \).
18 \( R \) = triangle inside the lines \( x = -1, y = 0, x + y = 0 \).
19 \( R \) = triangle inside the lines \( y = x, y = -x, y = 3 \).
20 \( R \) = triangle inside the lines \( y = x, y = 2x, y = 4 \).
21 \( R \) = triangle with vertices \( (0, 0), (4, 4), (4, 8) \).
22 \( R \) = triangle with vertices \( (0, 0), (-2, -1), (1, -2) \).
23 \( R \) = triangle with vertices \( (0, 0), (2, 0), (1, b) \). Here \( b > 0 \).
24 \( R \) = triangle with vertices \( (0, 0), (a, b), (c, d) \). The sides are \( y = bx/a, y = dx/c, \) and \( y = b + (x-a)(d-b)/(c-a) \). Find \( A = \int \int dy \, dx \) when \( 0 < a < c, 0 < d < b \).

25 Evaluate \( \int_{0}^{b} \int_{0}^{a} \partial f/\partial x \, dx \, dy \).
26 Evaluate \( \int_{0}^{b} \int_{0}^{a} \partial f/\partial y \, dx \, dy \).

In 27–28, divide the unit square \( R \) into triangles \( S \) and \( T \) and verify \( \int_{R} f \, dA = \int_{S} f \, dA + \int_{T} f \, dA \).

27 \( f(x, y) = 2x - 3y + 1 \)
28 \( f(x, y) = xe^y - ye^x \)

29 The area under \( y = f(x) \) is a single integral from \( a \) to \( b \) or a double integral (find the limits):

\[
\int_{a}^{b} f(x) \, dx = \int \int 1 \, dy \, dx
\]

30 Find the limits and the area under \( y = 1 - x^2 \):

\[
\int (1-x^2) \, dx \text{ and } \int \int 1 \, dy \, dx \text{ (reversed from 29).}
\]

31 A city inside the circle \( x^2 + y^2 = 100 \) has population density \( \rho(x, y) = 10(100 - x^2 - y^2) \). Integrate to find its population.

32 Find the volume bounded by the planes \( x = 0, y = 0, z = 0 \), and \( ax + by + cz = 1 \).

In 33–34 the rectangle with corners \( (1, 1), (1, 3), (2, 1), (2, 3) \) has density \( \rho(x, y) = x^2 \). The moments are \( M_x = \int \int x \rho \, dA \) and \( M_y = \int \int y \rho \, dA \).

33 Find the mass.
34 Find the center of mass.

In 35–36 the region is a circular wedge of radius 1 between the lines \( y = x \) and \( y = -x \).

35 Find the area.
36 Find the centroid \( (\bar{x}, \bar{y}) \).

37 Write a program to compute \( \int_{0}^{1} \int_{0}^{1} f(x, y) \, dx \, dy \) by the midpoint rule (midpoints of \( n \times n \) small squares). Which \( f(x, y) \) are integrated exactly by your program?

38 Apply the midpoint code to integrate \( x^2 \) and \( xy \) and \( y^2 \). The errors decrease like what power of \( \Delta x = \Delta y = 1/n? \)

Use the program to compute the volume under \( f(x, y) \) in 39–42. Check by integrating exactly or doubling \( n \).

39 \( f(x, y) = 3x + 4y + 5 \)
40 \( f(x, y) = 1/\sqrt{x^2 + y^2} \)
41 \( f(x, y) = x^y \)
42 \( f(x, y) = e^x \sin \pi y \)

43 In which order is \( \int \int x^2 \, dx \, dy = \int \int y^2 \, dx \, dy \) easier to integrate over the square \( 0 \leq x \leq 1, 0 \leq y \leq 1 \)? By reversing order, integrate \( (x-1)/\ln x \) from 0 to 1—if its antiderivative is unknown.

44 Explain in your own words the definition of the double integral of \( f(x, y) \) over the region \( R \).

45 \( \sum_{y} y \Delta A \) might not approach \( \int \int y \, dA \) if we only know that \( \Delta A \to 0 \). In the square \( 0 \leq x, y \leq 1 \), take rectangles of sides \( \Delta x \) and \( 1 \) (not \( \Delta x \) and \( \Delta y \)). If \( (x_1, y_1) \) is a point in the rectangle where \( y_1 = 1 \), then \( \sum_{y_1} \Delta A = \) _______. But \( \int \int y \, dA = \) _______. 
We have to be able to answer basic questions like these:

Find the area \( \iint dA \) and moment \( \iint x \, dA \) and moment of inertia \( \iint x^2 \, dA \).

The problem is: What is \( dA \)? We are leaving the \( xy \) variables where \( dA = dx \, dy \).

The reason for changing is this: The limits of integration in the \( y \) direction are miserable. I don't know them and I don't want to know them. For every \( x \) we would need the entry point \( P \) of the line \( x = \) constant, and the exit point \( Q \). The heights of \( P \) and \( Q \) are the limits on \( \int dy \), the inner integral. The geometry of the square and ring are totally missed, if we stick rigidly to \( x \) and \( y \).

Which coordinates are better? Any sensible person agrees that the area of the tilted square is 1. "Just turn it and the area is obvious." But that sensible person may not know the moment or the center of gravity or the moment of inertia. So we actually have to do the turning.

The new coordinates \( u \) and \( v \) are in Figure 14.6a. The limits of integration on \( v \) are 0 and 1. So are the limits on \( u \). But when you change variables, you don't just change limits. Two other changes come with new variables:

1. The small area \( dA = dx \, dy \) becomes \( dA = ______ du \, dv \).
2. The integral of \( x \) becomes the integral of ______.

Substituting \( u = \sqrt{x} \) in a single integral, we make the same changes. Limits \( x = 0 \) and \( x = 4 \) become \( u = 0 \) and \( u = 2 \). Since \( x \) is \( u^2 \), \( dx \) is \( 2u \, du \). The purpose of the change is to find an antiderivative. For double integrals, the usual purpose is to improve the limits—but we have to accept the whole package.

To turn the square, there are formulas connecting \( x \) and \( y \) to \( u \) and \( v \). The geometry is clear—rotate axes by \( \alpha \)—but it has to be converted into algebra:

\[
\begin{align*}
    u &= x \cos \alpha + y \sin \alpha \\
    v &= -x \sin \alpha + y \cos \alpha
\end{align*}
\]

and in reverse

\[
\begin{align*}
    x &= u \cos \alpha - v \sin \alpha \\
    y &= u \sin \alpha + v \cos \alpha
\end{align*}
\]

Figure 14.6 shows the rotation. As points move, the whole square turns. A good way to remember equation (1) is to follow the corners as they become \((1,0)\) and \((0,1)\).

The change from \( \iint x \, dA \) to \( \iint ______\, du \, dv \) is partly decided by equation (1). It gives \( x \) as a function of \( u \) and \( v \). We also need \( dA \). For a pure rotation the first guess is correct: The area \( dx \, dy \) equals the area \( du \, dv \). For most changes of variable this is false. The general formula for \( dA \) comes after the examples.
14.2 Change to Better Coordinates

Example 1 Find \( \iint dA \) and \( \iint x \, dA \) and \( x \) and also \( \iint x^2 \, dA \) for the tilted square.

Solution The area of the square is \( \int_0^1 \int_0^1 du \, dv = 1 \). Notice the good limits. Then

\[
\iint x \, dA = \int_0^1 \int_0^1 (u \cos x - v \sin x) \, du \, dv = \frac{1}{2} \cos x - \frac{1}{3} \sin x. \tag{2}
\]

This is the moment around the \( y \) axis. The factors \( \frac{1}{2} \) come from \( \frac{1}{2} u^2 \) and \( \frac{1}{2} v^2 \). The \( x \) coordinate of the center of gravity is

\[
\bar{x} = \frac{1}{\iint x \, dA} \int \int x \, dA = \left( \frac{1}{2} \cos x - \frac{1}{3} \sin x \right) \cdot 1.
\]

Similarly the integral of \( y \) leads to \( \bar{y} \). The answer is no mystery—the point \((\bar{x}, \bar{y})\) is at the center of the square! Substituting \( x = u \cos x - v \sin x \) made \( x \, dA \) look worse, but the limits 0 and 1 are much better.

The moment of inertia \( I_x \), around the \( y \) axis is also simplified:

\[
\iint x^2 \, dA = \int_0^1 \int_0^1 (u \cos x - v \sin x)^2 \, du \, dv = \frac{\cos^2 x}{3} - \frac{\cos x \sin x}{2} + \frac{\sin^2 x}{3}. \tag{3}
\]

You know this next fact but I will write it anyway: The answers don't contain \( u \) or \( v \). Those are dummy variables like \( x \) and \( y \). The answers do contain \( x \), because the square has turned. (The area is fixed at 1.) The moment of inertia \( I_x = \iint y^2 \, dA \) is the same as equation (3) but with all plus signs.

Question The sum \( I_x + I_y \) simplifies to \( \frac{5}{3} \) (a constant). Why no dependence on \( x \)?

Answer \( I_x + I_y \) equals \( I_0 \). This moment of inertia around \((0,0)\) is unchanged by rotation. We are turning the square around one of its corners.

Change to Polar Coordinates

The next change is to \( r \) and \( \theta \). A small area becomes \( dA = r \, dr \, d\theta \) (definitely not \( dr \, d\theta \)). Area always comes from multiplying two lengths, and \( d\theta \) is not a length. Figure 14.7 shows the crucial region—a "polar rectangle" cut out by rays and circles. Its area \( \Delta A \) is found in two ways, both leading to \( r \, dr \, d\theta \):

(Approximate) The straight sides have length \( \Delta r \). The circular arcs are close to \( r \Delta \theta \). The angles are 90°. So \( \Delta A \) is close to \((\Delta r)(r \Delta \theta)\).

(Exact) A wedge has area \( \frac{1}{2} r^2 \Delta \theta \). The difference between wedges is \( \Delta A \):

\[
\Delta A = \frac{1}{2} \left( r + \frac{\Delta r}{2} \right)^2 \Delta \theta - r^2 \left( \frac{\Delta r}{2} \right) \Delta \theta = r \Delta r \Delta \theta.
\]
The exact method places \( r \) dead center (see figure). The approximation says: Forget the change in \( r \Delta \theta \) as you move outward. Keep only the first-order terms.

A third method is coming, which requires no picture and no geometry. Calculus always has a third method! The change of variables \( x = r \cos \theta \), \( y = r \sin \theta \) will go into a general formula for \( dA \), and out will come the area \( r \, dr \, d\theta \).

**Fig. 14.7** Ring and polar rectangle in \( xy \) and \( r\theta \), with stretching factor \( r = 4.5 \).

### EXAMPLE 2
Find the area and center of gravity of the ring. Also find \( \iint x^2 \, dA \).

**Solution** The limits on \( r \) are 4 and 5. The limits on \( \theta \) are 0 and \( 2\pi \). *Polar coordinates are perfect for a ring.* Compared with limits like \( x = \sqrt{25-y^2} \), the change to \( r \, dr \, d\theta \) is a small price to pay:

\[
\text{area} = \int_0^{2\pi} \int_4^5 r \, dr \, d\theta = 2\pi \left[ \frac{1}{2} r^2 \right]_4^5 = \pi 5^2 - \pi 4^2 = 9\pi.
\]

The \( \theta \) integral is \( 2\pi \) (full circle). Actually the ring is a giant polar rectangle. We could have used the exact formula \( r \, dA \), with \( \Delta \theta = 2\pi \) and \( \Delta r = 5 - 4 \). When the radius \( r \) is centered at 4.5, the product \( r \, dA \, \Delta \theta \) is \((4.5)(1)(2\pi) = 9\pi \) as above.

Since the ring is symmetric around \((0, 0)\), the integral of \( x \, dA \) must be zero:

\[
\iint x \, dA = 2\pi \int_4^5 (r \cos \theta) r \, dr \, d\theta = \left[ \frac{1}{4} r^4 \right]_4^5 \left[ \sin \theta \right]_0^{2\pi} = 0.
\]

Notice \( r \cos \theta \) from \( x \)—the other \( r \) is from \( dA \). The moment of inertia is

\[
\iint x^2 \, dA = 2\pi \int_4^5 r^2 \cos^2 \theta \, r \, dr \, d\theta = \left[ \frac{1}{4} r^4 \right]_4^5 \int_0^{2\pi} \cos^2 \theta \, d\theta = \frac{1}{4} (5^4 - 4^4) \pi.
\]

This \( \theta \) integral is \( \pi \) not \( 2\pi \), because the average of \( \cos^2 \theta \) is \( \frac{1}{2} \) not 1.

For reference here are the moments of inertia when the density is \( \rho(x, y) \):

\[
I_y = \iint x^2 \rho \, dA \quad I_x = \iint y^2 \rho \, dA \quad I_0 = \iint r^2 \rho \, dA = \text{polar moment} = I_x + I_y. \tag{4}
\]

### EXAMPLE 3
Find masses and moments for semicircular plates: \( \rho = 1 \) and \( \rho = 1 - r \).

**Solution** The semicircles in Figure 14.8 have \( r = 1 \). The angle goes from 0 to \( \pi \) (the upper half-circle). Polar coordinates are best. *The mass is the integral of the density \( \rho \):*

\[
M = \int_0^\pi \int_0^1 r \, dr \, d\theta = \left( \frac{1}{2} \right)(\pi) \quad \text{and} \quad M = \int_0^\pi \int_0^1 (1-r) r \, dr \, d\theta = \left( \frac{1}{2} \right)(\pi).
\]
The first mass \( \pi/2 \) equals the area (because \( \rho = 1 \)). The second mass \( \pi/6 \) is smaller (because \( \rho < 1 \)). Integrating \( \rho = 1 \) is the same as finding a volume when the height is \( z = 1 \) (part of a cylinder). Integrating \( \rho = 1 - r \) is the same as finding a volume when the height is \( z = 1 - r \) (part of a cone). Volumes of cones have the extra factor \( \frac{1}{2} \).

The center of gravity involves the moment \( M_x = \iiint y \rho \, dA \). The distance from the \( x \) axis is \( y \), the mass of a small piece is \( \rho \, dA \), integrate to add mass times distance. Polar coordinates are still best, with \( y = r \sin \theta \). Again \( \rho = 1 \) and \( \rho = 1 - r \):

\[
\iiint y \, dA = \int_0^\pi \int_0^1 r \sin \theta \, dr \, d\theta = \frac{2}{3} \quad \iiint (1 - r) \, dA = \int_0^\pi \int_0^1 r \sin (1 - r) \, dr \, d\theta = \frac{1}{6}.
\]

The height of the center of gravity is \( \bar{y} = M_x/M = \text{moment divided by mass} \):

\[
\bar{y} = \frac{2/3}{\pi/2} = \frac{4}{3\pi} \quad \text{when } \rho = 1 \quad \bar{y} = \frac{1/6}{\pi/6} = \frac{1}{\pi} \quad \text{when } \rho = 1 - r.
\]

**Fig. 14.8** Semicircles with density piled above them.

**Fig. 14.9** Bell-shaped curve.

**Question** Compare \( \bar{y} \) for \( \rho = 1 \) and \( \rho = 1 - r \), other positive constants and \( \rho = 1 - r \).

**Answer** Any constant \( \rho \) gives \( \bar{y} = 4/3\pi \). Since \( 1 - r \) is dense at \( r = 0 \), \( \bar{y} \) drops to \( 1/\pi \).

**Question** How is \( \bar{y} = 4/3\pi \) related to the "average" of \( y \) in the semicircle?

**Answer** They are identical. This is the point of \( \bar{y} \). Divide the integral by the area:

\[
\text{The average value of a function is } \bar{y} = \frac{\iint f(x, y) \, dA}{\iint dA}.
\]

The integral of \( f \) is divided by the integral of 1 (the area). In one dimension \( \int_a^b v(x) \, dx \) was divided by \( \int_a^b 1 \, dx \) (the length \( b - a \)). That gave the average value of \( v(x) \) in Section 5.6. Equation (5) is the same idea for \( f(x, y) \).

**EXAMPLE 4** Compute \( A = \int_{-\infty}^{\infty} e^{-x^2} \, dx = \sqrt{\pi} \) from \( A^2 = \int_{-\infty}^{\infty} e^{-x^2} \, dx \int_{-\infty}^{\infty} e^{-y^2} \, dy = \pi \).

\( A \) is the area under a "bell-shaped curve"—see Figure 14.9. This is the most important definite integral in the study of probability. It is difficult because a factor \( 2x \) is not present. Integrating \( 2xe^{-x^2} \) gives \( -e^{-x^2} \), but integrating \( e^{-x^2} \) is impossible—except approximately by a computer. How can we hope to show that \( A \) is exactly \( \sqrt{\pi} \)?

The trick is to go from an area integral \( A \) to a volume integral \( A^2 \). This is unusual (and hard to like), but the end justifies the means:

\[
A^2 = \int_{x=-\infty}^{\infty} \int_{y=-\infty}^{\infty} e^{-x^2} e^{-y^2} \, dy \, dx = \int_0^{2\pi} \int_0^{\infty} e^{-r^2} r \, dr \, d\theta.
\]

The double integrals cover the whole plane. The \( r^2 \) comes from \( x^2 + y^2 \), and the key factor \( r \) appears in polar coordinates. It is now possible to substitute \( u = r^2 \). The \( r \) integral is \( \frac{1}{2} \int_0^{\infty} e^{-u} \, du = \frac{1}{2} \). The \( \theta \) integral is \( 2\pi \). The double integral is \( \frac{1}{2}(2\pi) \). Therefore \( A^2 = \pi \) and the single integral is \( A = \sqrt{\pi} \).
EXAMPLE 5  Apply Example 4 to the "normal distribution" $p(x) = e^{-x^2/2}/\sqrt{2\pi}$.

Section 8.4 discussed probability. It emphasized the importance of this particular $p(x)$. At that time we could not verify that $\int p(x)\,dx = 1$. Now we can:

$$x = \sqrt{2}y \quad \text{yields} \quad \frac{1}{\sqrt{2\pi}} \int_{-\infty}^{\infty} e^{-x^2/2}\,dx = \frac{1}{\sqrt{\pi}} \int_{-\infty}^{\infty} e^{-y^2}\,dy = 1.$$  \hspace{1cm} (7)

**Question** Why include the 2's in $p(x)$? The integral of $e^{-x^2}/\sqrt{\pi}$ also equals 1.

**Answer** With the 2's the "variance" is $\int x^2p(x)\,dx = 1$. This is a convenient number.

**CHANGE TO OTHER COORDINATES**

A third method was promised, to find $r\,dr\,d\theta$ without a picture and without geometry. The method works directly from $x = r\cos \theta$ and $y = r\sin \theta$. It also finds the 1 in $du\,dv$, after a rotation of axes. Most important, this new method finds the factor $J$ in the area $dA = J\,du\,dv$, for any change of variables. The change is from $xy$ to $uv$.

For single integrals, the "stretching factor" $J$ between the original $dx$ and the new $du$ is (not surprisingly) the ratio $dx/du$. Where we have $dx$, we write $(dx/du)du$. Where we have $(du/dx)dx$, we write $du$. That was the idea of substitutions—the main way to simplify integrals.

For double integrals the stretching factor appears in the area: $dx\,dy$ becomes $|J|\,du\,dv$. The old and new variables are related by $x = x(u, v)$ and $y = y(u, v)$. The point with coordinates $u$ and $v$ comes from the point with coordinates $x$ and $y$. A whole region $S$, full of points in the $uv$ plane, comes from the region $R$ full of corresponding points in the $xy$ plane. A small piece with area $|J|\,du\,dv$ comes from a small piece with area $dx\,dy$. The formula for $J$ is a two-dimensional version of $dx/du$.

14B  The stretching factor for area is the 2 by 2 Jacobian determinant $J(u, v)$:

$$J = \begin{vmatrix} \frac{\partial x}{\partial u} & \frac{\partial x}{\partial v} \\ \frac{\partial y}{\partial u} & \frac{\partial y}{\partial v} \end{vmatrix} = \frac{\partial x}{\partial u} \frac{\partial y}{\partial v} - \frac{\partial x}{\partial v} \frac{\partial y}{\partial u}. \hspace{1cm} (8)$$

An integral over $R$ in the $xy$ plane becomes an integral over $S$ in the $uv$ plane:

$$\int \int_R f(x, y)\,dx\,dy = \int \int_S f(x(u, v), y(u, v)) |J|\,du\,dv. \hspace{1cm} (9)$$

The determinant $J$ is often written $\partial(x, y)/\partial(u, v)$, as a reminder that this stretching factor is like $dx/du$. We require $J \neq 0$. That keeps the stretching and shrinking under control.

You naturally ask: Why take the absolute value $|J|$ in equation (9)? Good question—it wasn't done for single integrals. The reason is in the limits of integration. The single integral $\int_0^1 x\,dx$ is $\int_0^{-1} (-du)$ after changing $x$ to $-u$. We keep the minus sign and allow single integrals to run backward. Double integrals could too, but normally they go left to right and down to up. We use the absolute value $|J|$ and run forward.

**EXAMPLE 6**  Polar coordinates have $x = u\cos v = r\cos \theta$ and $y = u\sin v = r\sin \theta$.

With no geometry:

$$J = \begin{vmatrix} \frac{\partial x}{\partial r} & \frac{\partial x}{\partial \theta} \\ \frac{\partial y}{\partial r} & \frac{\partial y}{\partial \theta} \end{vmatrix} = \begin{vmatrix} \cos \theta & -r \sin \theta \\ \sin \theta & r \cos \theta \end{vmatrix} = r. \hspace{1cm} (10)$$
EXAMPLE 7  Find $J$ for the linear change to $x = au + bv$ and $y = cu + dv$.

Ordinary determinant:  $J = \begin{vmatrix} \frac{\partial x}{\partial u} & \frac{\partial x}{\partial v} \\ \frac{\partial y}{\partial u} & \frac{\partial y}{\partial v} \end{vmatrix} = \begin{vmatrix} a & b \\ c & d \end{vmatrix} = ad - bc. \tag{11}$

Why make this simple change, in which $a, b, c, d$ are all constant? It straightens parallelograms into squares (and rotates those squares). Figure 14.10 is typical.

Common sense indicated $J = 1$ for pure rotation—no change in area. Now $J = 1$ comes from equations (1) and (11), because $ad - bc$ is $\cos^2 \alpha + \sin^2 \alpha$.

In practice, $xy$ rectangles generally go into $uv$ rectangles. The sides can be curved (as in polar rectangles) but the angles are often 90°. The change is "orthogonal." The next example has angles that are not 90°, and $J$ still gives the answer.

\begin{figure}
\centering
\includegraphics{fig14.10}
\caption{Change from $xy$ to $uv$ has $J = \frac{1}{3}$.}
\end{figure}

EXAMPLE 8  Find the area of $R$ in Figure 14.10. Also compute $\iint_R e^{x+y} \, dx \, dy$.

Solution  The figure shows $x = \frac{3}{2}u + \frac{1}{2}v$ and $y = \frac{1}{2}u + \frac{3}{2}v$. The determinant is

$$J = \begin{vmatrix} \frac{\partial x}{\partial u} & \frac{\partial x}{\partial v} \\ \frac{\partial y}{\partial u} & \frac{\partial y}{\partial v} \end{vmatrix} = \begin{vmatrix} 2/3 & 1/3 \\ 1/3 & 2/3 \end{vmatrix} = \frac{4}{9} - \frac{1}{9} = \frac{1}{3}.$$  

The area of the $xy$ parallelogram becomes an integral over the $uv$ square:

$$\iint_R dx \, dy = \iint_S |J| \, du \, dv = \int_0^3 \frac{3}{3} \, du \, dv = \frac{1}{3} \cdot 3 = 3.$$  

The square has area 9, the parallelogram has area 3. I don’t know if $J = \frac{1}{3}$ is a stretching factor or a shrinking factor. The other integral $\iint e^{x+y} \, dx \, dy$ is

$$\int_0^3 \int_0^3 e^{2u^3 - v^3} \, \frac{1}{3} \, du \, dv = \left[ \int_0^3 e^{2u^3} \right]_0^3 \int_0^3 e^{0^3} \, \frac{1}{3} = \frac{3}{2} (e^2 - 1)(e - 1).$$

Main point: The change to $u$ and $v$ makes the limits easy (just 0 and 3).

Why is the stretching factor $J$ a determinant? With straight sides, this goes back to Section 11.3 on vectors. The area of a parallelogram is a determinant. Here the sides are curved, but that only produces $(du)^2$ and $(dv)^2$, which we ignore.

A change $du$ gives one side of Figure 14.11—it is $(\partial x/\partial u \, i + \partial y/\partial u \, j) \, du$. Side 2 is $(\partial x/\partial v \, i + \partial y/\partial v \, j) \, dv$. The curving comes from second derivatives. The area (the cross product of the sides) is $|J| \, du \, dv$.  

\begin{figure}
\centering
\includegraphics{fig14.11}
\caption{Curved areas are also $dA = |J| \, du \, dv$.}
\end{figure}
14 Multiple Integrals

Final remark I can't resist looking at the change in the reverse direction. Now the rectangle is in xy and the parallelogram is in uv. In all formulas, exchange x for u and y for v:

\[ \text{new } J = \begin{vmatrix} \frac{\partial u}{\partial x} & \frac{\partial u}{\partial y} \\ \frac{\partial v}{\partial x} & \frac{\partial v}{\partial y} \end{vmatrix} = \frac{\partial (u, v)}{\partial (x, y)} = \frac{1}{\text{old } J}. \] (12)

This is exactly like \( du/dx = 1/(dx/du) \). It is the derivative of the inverse function. The product of slopes is 1—stretch out, shrink back. From xy to uv we have 2 by 2 matrices, and the identity matrix \( I \) takes the place of 1:

\[ \frac{dx}{du} \cdot \frac{du}{dx} = 1. \]

The first row times the first column is \( (\partial x/\partial u)(\partial u/\partial x) + (\partial x/\partial v)(\partial v/\partial x) = \partial x/\partial x = 1 \).

The first row times the second column is \( (\partial x/\partial u)(\partial u/\partial y) + (\partial x/\partial v)(\partial v/\partial y) = \partial x/\partial y = 0 \).

The matrices are inverses of each other. The determinants of a matrix and its inverse obey our rule: old \( J \) times new \( J = 1 \). Those \( J \)'s cannot be zero, just as \( dx/du \) and \( du/dx \) were not zero. (Inverse functions increase steadily or decrease steadily.)

In two dimensions, an area \( dx \, dy \) goes to \( J \, du \, dv \) and comes back to \( dx \, dy \).
10 If the centroid of \( R \) is \((0, 0)\), the centroid of \( R^* \) is _______. The centroid of the circle with radius 3 and center \((1, 2)\) is _______. The centroid of the upper half of that circle is _______.

11 The moments of inertia \( I_x, I_y, I_0 \) of the original wedge \( R \) are _______.

12 The moments of inertia \( I_x, I_y, I_0 \) of the shifted wedge \( R^* \) are _______.

Problems 13–16 change four-sided regions to squares.

13 \( R \) has straight sides \( y = 2x, x = 1, y = 1 + 2x, x = 0 \). Locate its four corners and draw \( R \). Find its area by geometry.

14 Choose \( a, b, c, d \) so that the change \( x = au + bv, y = cu + dv \) takes the previous \( R \) into \( S \), the unit square \( 0 \leq u \leq 1, 0 \leq v \leq 1 \). From the stretching factor \( J = ad - bc \), find the area of \( R \).

15 The region \( R \) has straight sides \( x = 0, x = 1, y = 0, y = 2x + 3 \). Choose \( a, b, c \) so that \( x = au + bv \) and \( y = au + bv + cu \) change \( R \) to the unit square \( S \).

16 A nonlinear term \( uv \) was needed in Problem 15. Which regions \( R \) could change to the square \( S \) with a linear \( x = au + bv, y = cu + dv \)?

Draw the \( xy \) region \( R \) that corresponds in 17–22 to the \( uv \) square \( S \) with corners \((0, 0), (1, 0), (0, 1), (1, 1) \). Locate the corners of \( R \) and then its sides (like a jigsaw puzzle).

17 \( x = 2u + v, y = u + 2v \)

18 \( x = 3u + 2v, y = u + v \)

19 \( x = e^{2u}, y = e^u + 2v \)

20 \( x = uv, y = v^2 - u^2 \)

21 \( x = u, y = u(1 + u) \)

22 \( x = u \cos v, y = u \sin v \) (only three corners)

23 In Problems 17 and 19, compute \( J \) from equation (8). Then find the area of \( R \) from \( \int\int_S J \, du \, dv \).

24 In 18 and 20, find \( J = \partial(x, y)/\partial(u, v) \) and the area of \( R \).

25 If \( R \) lies between \( x = 0 \) and \( x = 1 \) under the graph of \( y = f(x) > 0 \), then \( x = u, y = ef(u) \) takes \( R \) to the unit square \( S \). Locate the corners of \( R \) and the point corresponding to \( u = 1/2, v = 1 \). Compute \( J \) to prove what we know:

area of \( R \) = \( \int_0^1 f(x) \, dx \) = \( \int_0^1 J \, du \, dv \).

26 From \( r = \sqrt{x^2 + y^2} \) and \( \theta = \tan^{-1}(y/x) \), compute \( \partial r/\partial x, \partial r/\partial y, \partial \theta/\partial x, \partial \theta/\partial y \), and the determinant \( J = \partial(r, \theta)/\partial(x, y) \). How is this \( J \) related to the factor \( r \, dr \, d\theta \) that enters \( r \, dr \, d\theta \)?

27 Example 4 integrated \( e^{-r^2} \) from 0 to \( \infty \) (answer \( \sqrt{\pi} \)). Also \( B = \int_0^\infty e^{-x^2} \, dx \) leads to \( B^2 = \int_0^\infty e^{-x^2} \, dx \int_0^\infty e^{-y^2} \, dy \). Change this double integral over the unit square to \( r \) and \( \theta \) — and find the limits on \( r \) that make exact integration impossible.

28 Integrate by parts to prove that the standard normal distribution \( p(x) = e^{-x^2}/(\sqrt{2\pi}) \) has \( \sigma^2 = \int_{-\infty}^{\infty} x^2 p(x) \, dx = 1 \).

29 Find the average distance from a point on a circle to the points inside. Suggestion: Let \((0, 0)\) be the point and let \( 0 \leq r \leq 2a \cos \theta, 0 \leq \theta \leq \pi \) be the circle (radius \( a \)). The distance is \( r \), so the average distance is \( \bar{r} = \int_{\theta=0}^{\theta=\pi} \int_{r=0}^{r=2a \cos \theta} r \, dr \, d\theta \).

30 Draw the region \( R: 0 \leq x \leq 1, 0 \leq y < \infty \) and describe it with polar coordinates (limits on \( r \) and \( \theta \)). Integrate \( \int_S (x^2 + y^2)^{-3/2} \, dx \, dy \) in polar coordinates.

31 Using polar coordinates, find the volume under \( z = x^2 + y^2 \) above the unit disk \( x^2 + y^2 \leq 1 \).

32 The end of Example 1 stated the moment of inertia \( \int \int_S y^2 \, dA \). Check that integration.

33 In the square \(-1 \leq x \leq 2, -2 \leq y \leq 1\), where could you distribute a unit mass (with \( \int \int p \, dxdy = 1 \)) to maximize

(a) \( \int \int x^2 \rho \, dA \)  \hspace{1em} (b) \( \int \int y^2 \rho \, dA \)  \hspace{1em} (c) \( \int \int r^2 \rho \, dA \)

34 True or false, with a reason:

(a) If the \( uv \) region \( S \) corresponds to the \( xy \) region \( R \), then the area of \( S = \text{area of } R \).

(b) \( \int \int_S x \, dA \leq \int \int x^2 \, dA \)

(c) The average value of \( f(x, y) \) is \( \int \int f(x, y) \, dA \)

(d) \( \int_{-\infty}^{\infty} x e^{-x^2} \, dx = 0 \)

(e) A polar rectangle has the same area as a straight-sided region with the same corners.

35 Find the moment of the tilted square in Example 1 if the density is \( \rho = xy \).

36 Find the mass of the ring in Example 2 if the density is \( \rho = x^2 + y^2 \). This is the same as which moment of inertia with which density?

37 Find the polar moment of inertia \( I_0 \) of the ring in Example 2 if the density is \( \rho = x^2 + y^2 \).

38 Give the following statement an appropriate name:

\( \int \int_S f(x, y) \, dA = f(P) \times \text{area of } R \), where \( P \) is a point in \( R \). Which point \( P \) makes this correct for \( f = x \) and \( f = y \)?

39 Find the \( xy \) coordinates of the top point in Figure 14.6a and check that it goes to \((u, v) = (1, 1) \).
At this point in the book, I feel I can speak to you directly. You can guess what triple integrals are like. Instead of a small interval or a small rectangle, there is a small box. Instead of length \( dx \) or area \( dx \, dy \), the box has volume \( dV = dx \, dy \, dz \). That is length times width times height. The goal is to put small boxes together (by integration). The main problem will be to discover the correct limits on \( x, y, z \).

We could dream up more and more complicated regions in three-dimensional space. But I don't think you can see the method clearly without seeing the region clearly. In practice six shapes are the most important:

- box
- prism
- cylinder
- cone
- tetrahedron
- sphere.

The box is easiest and the sphere may be the hardest (but no problem in spherical coordinates). Circular cylinders and cones fall in the middle, where \( xyz \) coordinates are possible but \( r\theta z \) are the best. I start with the box and prism and \( xyz \).

**EXAMPLE 1**  By triple integrals find the volume of a box and a prism (Figure 14.12).

\[
\iiint_{\text{box}} dV = \int_{z=0}^{2} \int_{y=0}^{3} \int_{x=0}^{2} dx \, dy \, dz \quad \text{and} \quad \iiint_{\text{prism}} dV = \int_{z=0}^{3} \int_{y=0}^{3-3z} \int_{x=0}^{2} dx \, dy \, dz
\]

The inner integral for both is \( \int dx = 2 \). Lines in the \( x \) direction have length 2, cutting through the box and the prism. The middle integrals show the limits on \( y \) (since \( dy \) comes second):

\[
\int_{y=0}^{3} 2 \, dy = 6 \quad \text{and} \quad \int_{y=0}^{3-3z} 2 \, dy = 6 - 6z.
\]

After two integrations these are areas. The first area 6 is for a plane section through the box. The second area 6 - 6z is cut through the prism. The shaded rectangle goes from \( y = 0 \) to \( y = 3 - 3z \)—we needed and used the equation \( y + 3z = 3 \) for the boundary of the prism. At this point \( z \) is still constant! But the area depends on \( z \), because the prism gets thinner going upwards. The base area is \( 6 - 6z = 6 \), the top area is \( 6 - 6z = 0 \).

The outer integral multiplies those areas by \( dz \), to give the volume of slices. They are horizontal slices because \( z \) came last. Integration adds up the slices to find the total volume:

- box volume \( = \int_{z=0}^{1} 6 \, dz = 6 \)
- prism volume \( = \int_{z=0}^{1} (6 - 6z)dz = \left[ 6z - 3z^2 \right]_0^1 = 3 \).

The box volume \( 2 \cdot 3 \cdot 1 \) didn't need calculus. The prism is half of the box, so its volume was sure to be 3—but it is satisfying to see how \( 6z - 3z^2 \) gives the answer. Our purpose is to see how a triple integral works.

**Fig. 14.12**  Box with sides 2, 3, 1. The prism is half of the box: volume \( \int (6 - 6z)dz \) or \( \frac{3}{2} dx \).
14.3 Triple Integrals

**Question** Find the prism volume in the order $dz\ dy\ dx$ (six orders are possible).

**Answer** \[ \int_0^2 \int_0^3 \int_0^{(3-y)^3} dz\ dy\ dx = \int_0^2 \int_0^3 \left(\frac{3-y}{3}\right) dy\ dx = \int_0^2 \frac{3}{2} dx = 3. \]

To find those limits on the $z$ integral, follow a line in the $z$ direction. It enters the prism at $z = 0$ and exits at the sloping face $y + 3z = 3$. That gives the upper limit $z = (3 - y)/3$. It is the height of a thin stick as in Section 14.1. This section writes out $\int dz$ for the height, but a quicker solution starts at the double integral.

What is the number $\frac{1}{3}$ in the last integral? It is the area of a vertical slice, cut by a plane $x = \text{constant}$. The outer integral adds up slices.

\[ \iiint f(x, y, z)\ dV \text{ is computed from three single integrals} \int \left[ \int \left( \int f\ dx \right) dy \right] dz. \]

That step cannot be taken in silence—some basic calculus is involved. The triple integral is the limit of $\sum f_i \Delta V$, a sum over small boxes of volume $\Delta V$. Here $f_i$ is any value of $f(x, y, z)$ in the $i$th box. (In the limit, the boxes fit a curved region.) Now take those boxes in a certain order. Put them into lines in the $x$ direction and put the lines of boxes into planes. The lines lead to the inner $x$ integral, whose answer depends on $y$ and $z$. The $y$ integral combines the lines into planes. Finally the outer integral accounts for all planes and all boxes.

Example 2 is important because it displays more possibilities than a box or prism.

**EXAMPLE 2** Find the volume of a tetrahedron (4-sided pyramid). Locate $(x, y, z)$.

**Solution** A tetrahedron has four flat faces, all triangles. The fourth face in Figure 14.13 is on the plane $x + y + z = 1$. A line in the $x$ direction enters at $x = 0$ and exits at $x = 1 - y - z$. (The length depends on $y$ and $z$. The equation of the boundary plane gives $x$.) Then those lines are put into plane slices by the $y$ integral:

\[ \int_{y=0}^{1-z} \int_{x=0}^{1-y-z} dx\ dy = \int_{y=0}^{1-z} (1 - y - z) dy = \left[ y - \frac{1}{2} y^2 - zy \right]_{y=0}^{1-z} = \frac{1}{2}(1 - z)^2. \]

What is this number $\frac{1}{2}(1 - z)^2$? It is the area at height $z$. The plane at that height slices out a right triangle, whose legs have length $1 - z$. The area is correct, but look at the limits of integration. If $x$ goes to $1 - y - z$, why does $y$ go to $1 - z$? Reason: We are assembling lines, not points. The figure shows a line at every $y$ up to $1 - z$.

![Fig. 14.13](image-url) Lines end at plane $x + y + z = 1$. Triangles end at edge $y + z = 1$. The average height is $\bar{z} = \frac{1}{2} \int z\ dV/\iiint dV$.
Adding the slices gives the volume: \( \int_0^1 \frac{1}{2}(1-z)^2 \, dz = \frac{1}{6}(z-1)^3 \bigg|_0^1 = \frac{1}{3} \). This agrees with \( \frac{1}{3} \) (base times height), the volume of a pyramid.

The height \( \bar{z} \) of the centroid is \( \bar{z} \) _average_. We compute \( \iiint z \, dV \) and divide by the volume. Each horizontal slice is multiplied by its height \( z \), and the limits of integration don't change:

\[
\iiint z \, dV = \int_0^1 \int_0^{1-x} \int_0^{1-y-z} z \, dx \, dy \, dz = \int_0^1 \frac{z(1-z)^2}{2} \, dz = \frac{1}{24}.
\]

This is quick because \( z \) is constant in the \( x \) and \( y \) integrals. Each triangular slice contributes \( z \) times its area \( \frac{1}{2}(1-z)^2 \) times \( dz \). Then the \( z \) integral gives the moment 1/24. To find the average height, divide 1/24 by the volume:

\[
\bar{z} = \text{height of centroid} = \frac{\iiint z \, dV}{\iiint dV} = \frac{1/24}{1/6} = \frac{1}{4}.
\]

By symmetry \( \bar{x} = \frac{1}{3} \) and \( \bar{y} = \frac{1}{3} \). The centroid is the point \( \left( \frac{1}{3}, \frac{1}{3}, \frac{1}{3} \right) \). Compare that with \( \left( \frac{1}{3}, \frac{1}{3}, \frac{1}{3} \right) \), the centroid of the standard right triangle. Compare also with \( \frac{1}{3} \), the center of the unit interval. There must be a five-sided region in four dimensions centered at \( \left( \frac{1}{3}, \frac{1}{3}, \frac{1}{3} \right) \).

For area and volume we meet another pattern. Length of standard interval is 1, area of standard triangle is \( \frac{1}{2} \), volume of standard tetrahedron is \( \frac{1}{2} \), hypervolume in four dimensions must be \( \frac{1}{24} \). The interval reaches the point \( x = 1 \), the triangle reaches the line \( x + y = 1 \), the tetrahedron reaches the plane \( x + y + z = 1 \). The four-dimensional region stops at the hyperplane \( x + y + z = 1 \).

**EXAMPLE 3** Find the volume \( \iiint dx \, dy \, dz \) inside the unit sphere \( x^2 + y^2 + z^2 = 1 \).

First question: What are the limits on \( x \)? If a needle goes through the sphere in the \( x \) direction, where does it enter and leave? Moving in the \( x \) direction, the numbers \( y \) and \( z \) stay constant. The inner integral deals only with \( x \). The smallest and largest \( x \) are at the boundary where \( x^2 + y^2 + z^2 = 1 \). This equation does the work—we solve it for \( x \). Look at the limits on the \( x \) integral:

\[
\text{volume of sphere} = \iiint dx \, dy \, dz = \iiint \sqrt{1-y^2-z^2} \, dx \, dy \, dz = \int \int \int 2\sqrt{1-y^2-z^2} \, dy \, dz. \tag{1}
\]

The limits on \( y \) are \( -\sqrt{1-z^2} \) and \( +\sqrt{1-z^2} \). You can use algebra on the boundary equation \( x^2 + y^2 + z^2 = 1 \). But notice that \( x \) is gone! We want the smallest and largest \( y \), for each \( z \). It helps very much to draw the plane at height \( z \), slicing through the sphere in Figure 14.14. The slice is a circle of radius \( r = \sqrt{1-z^2} \). So the area is \( \pi r^2 \), which must come from the \( y \) integral:

\[
\int 2\sqrt{1-y^2-z^2} \, dy = \text{area of slice} = \pi(1-z^2). \tag{2}
\]

I admit that I didn't integrate. Is it cheating to use the formula \( \pi r^2 \)? I don't think so. Mathematics is hard enough, and we don't have to work blindfolded. The goal is understanding, and if you know the area then use it. Of course the integral of \( \sqrt{1-y^2-z^2} \) can be done if necessary—use Section 7.2.

The triple integral is down to a single integral. We went from one needle to a circle of needles and now to a sphere of needles. The volume is a sum of slices of area \( \pi(1-z^2) \). The South Pole is at \( z = -1 \), the North Pole is at \( z = +1 \), and the integral
is the volume $4\pi/3$ inside the unit sphere:

$$
\int_{-1}^{1} \pi(1-z^2)dz = \pi \left( \frac{1}{3} \pi \right) \left. \left( \frac{z}{3} - \frac{1}{3}z^3 \right) \right|_{-1}^{1} = \frac{2}{3} \pi - \left( \frac{2}{3} \pi \right) = \frac{4}{3} \pi.
$$

**Question 1** A cone also has circular slices. How is the last integral changed?  
**Answer** The slices of a cone have radius $1-z$. Integrate $(1-z)^2$ not $\sqrt{1-z^2}$.

**Question 2** How does this compare with a circular cylinder (height $1$, radius $1$)?  
**Answer** Now all slices have radius 1. Above $z=0$, a cylinder has volume $\pi$ and a half-sphere has volume $\frac{1}{2}\pi$ and a cone has volume $\frac{1}{4}\pi$.  
For solids with equal surface area, the sphere has largest volume.

**Question 3** What is the average height $\bar{z}$ in the cone and half-sphere and cylinder?  
**Answer**

\[
\bar{z} = \frac{\int z \text{(slice area)} \, dz}{\int \text{(slice area)} \, dz} = \frac{1}{4} \text{ and } \frac{3}{8} \text{ and } \frac{1}{2}.
\]

![Fig. 14.14](image)

**Example 4** Find the volume $\iiint dx \, dy \, dz$ inside the ellipsoid $x^2/a^2 + y^2/b^2 + z^2/c^2 = 1$.

The limits on $x$ are now $\pm \sqrt{1-y^2/b^2 - z^2/c^2}$. The algebra looks terrible. The geometry is better—all slices are ellipses. A **change of variable** is absolutely the best.

Introduce $u = x/a$ and $v = y/b$ and $w = z/c$. Then the outer boundary becomes $u^2 + v^2 + w^2 = 1$. In these new variables the shape is a sphere. The triple integral for a sphere is $\iiint du \, dv \, dw = 4\pi/3$. But what volume $dV$ in $xyz$ space corresponds to a small box with sides $du$ and $dv$ and $dw$?

Every $uvw$ box comes from an $xyz$ box. The box is stretched with no bending or twisting. Since $u = x/a$, the length $dx$ is $a \, du$. Similarly $dy = b \, dv$ and $dz = c \, dw$. The volume of the $xyz$ box (Figure 14.14) is $dx \, dy \, dz = (abc) \, du \, dv \, dw$. The **strecthing factor** $J = abc$ is a constant, and the volume of the ellipsoid is

\[
\iiint_{\text{ellipsoid}} dx \, dy \, dz = \iiint_{\text{sphere}} (abc) \, du \, dv \, dw = \frac{4\pi}{3} abc.
\]

You realize that this is special—other volumes are much more complicated. The sphere and ellipsoid are curved, but the small $xyz$ boxes are straight. The next section introduces spherical coordinates, and we can finally write "**good limits.**" But then we need a different $J$. 


14.3 Exercises

Read-through questions

Six important solid shapes are ____, ____, ____. The integral \( \iiint dx
dy
dz \) adds the volume of small ____. For computation it becomes ____, single integrals. The inner integral \( \int dx \) is the ____, line through the solid. The variables ____, and ____ are held constant. The double integral \( \iint dx
dy \) is the ____, slice, with ____ held constant. Then the z integral adds up the volumes of _____.

To find the total mass in \( V \) if the density is \( \rho = e^z \), we compute the integral ____. In the new variables \( u = x, v = 2y, w = 3z \), the equation of the outer boundary is ____. The volume of the tetrahedron in \( uvw \) space is ____. From \( dx = du \) and \( dy = dy/2 \) and \( dz = e^z \), the volume of an \( xyz \) box is \( dx
dy
dz = e^z
du
dv
dw \). So the volume of \( V \) is ___.

The average height \( z \) in \( V \) we compute ____. To find the total mass in \( V \) if the density is \( \rho = e^z \) we compute the integral ____. To find the average density we compute ____. In the order \( \iint dx
dy
dz \) the limits on the inner integral can depend on ____. The limits on the middle integral can depend on ____. The limits on the outer integral can depend on ____. The outer limits for the ellipsoid \( x^2 + 2y^2 + 3z^2 = 8 \) are ___.

1. For the solid region \( 0 \leq x \leq y \leq z \leq 1 \), find the limits in \( \iiint dx
dy
dz \) and compute the volume.

2. Reverse the order in Problem 1 to \( \iiint dy
dx
dz \) and find the limits of integration. The four faces of this tetrahedron are the planes \( x = 0 \) and \( y = x \) and ________.

3. This tetrahedron and five others like it fill the unit cube. Change the inequalities in Problem 1 to describe the other five.

4. Find the centroid \( (\bar{x}, \bar{y}, \bar{z}) \) in Problem 1.

Find the limits of integration in \( \iiint dx
dy
dz \) and the volume of solids 5–16. Draw a very rough picture.

5. A cube with sides of length 2, centered at \((0, 0, 0)\).

6. Half of that cube, the box above the \( xy \) plane.

7. Part of the same cube, the prism above the plane \( z = y \).

8. Part of the same cube, above \( z = y \) and \( z = 0 \).

9. Part of the same cube, above \( x = z \) and below \( z = y \).

10. Part of the same cube, where \( x \leq y \leq z \). What shape is this?

11. The tetrahedron bounded by planes \( x = 0, y = 0, z = 0 \), and \( x + y + 2z = 2 \).

12. The tetrahedron with corners \((0, 0, 0), (2, 0, 0), (0, 4, 0), (0, 0, 4)\). First find the plane through the last three corners.

13. The part of the tetrahedron in Problem 11 below \( z = 1/2 \).

14. The tetrahedron in Problem 12 with its top sliced off by the plane \( z = 1 \).

15. The volume above \( z = 0 \) below the cone \( \sqrt{x^2 + y^2} = 1 - z \).

16. The tetrahedron in Problem 12, after it falls across the \( x \) axis onto the \( xy \) plane.

In 17–20 find the limits in \( \iiint dx
dy
dz \) or \( \iiint dz
dy
dx \). Compute the volume.

17. A circular cylinder with height 6 and base \( x^2 + y^2 \leq 1 \).

18. The part of that cylinder below the plane \( z = x \). Watch the base. Draw a picture.

19. The volume shared by the cube (Problem 5) and cylinder.

20. The same cylinder lying along the \( x \) axis.

21. A cube is inscribed in a sphere: radius 1, both centers at \((0, 0, 0)\). What is the volume of the cube?

22. Find the volume and the centroid of the region bounded by \( x = 0, y = 0, z = 0 \), and \( x/a + y/b + z/c = 1 \).

23. Find the volume and centroid of the solid \( 0 \leq z \leq 4 - x^2 - y^2 \).

24. Based on the text, what is the volume inside \( x^2 + 4y^2 + 9z^2 = 16 \)? What is the “hypervolume” of the 4-dimensional pyramid that stops at \( x + y + z + w = 1 \)?

25. Find the partial derivatives \( \partial I/\partial x, \partial I/\partial y, \partial I/\partial z \) of

\[ I = \iiint f(x, y, z) \, dx \, dy \, dz \]

26. Define the average value of \( f(x, y, z) \) in a solid \( V \).

27. Find the moment of inertia \( \iiint r^2 \, dV \) of the cube \( |x| \leq 1, |y| \leq 1, |z| \leq 1 \) when \( r \) is the distance to (a) the \( x \) axis (b) the edge \( y = z = 1 \) (c) the diagonal \( x = y = z \).

28. Add upper limits to produce the volume of a unit cube from small cubes: \( V = \sum_{i=1}^{3\Delta x} \sum_{j=1}^{3\Delta y} \sum_{k=1}^{3\Delta z} (\Delta x)^3 = 1 \).

29. Find the limit as \( \Delta x \to 0 \) of \( \sum_{i=1}^{3\Delta x} \sum_{j=1}^{3\Delta y} \sum_{k=1}^{3\Delta z} (\Delta x)^3 \).

30. The midpoint rule for an integral over the unit cube chooses the center value \( f(\frac{1}{2}, \frac{1}{2}, \frac{1}{2}) \). Which functions \( f = x^m y^n z^p \) are integrated correctly?

31. The trapezoidal rule estimates \( \int_0^1 f(x, y, z) \, dx \, dy \, dz \) as \( \frac{1}{2} \) times the sum of \( f(x, y, z) \) at 8 corners. This correctly integrates \( x^m y^n z^p \) for which \( m, n, p \)?

32. Propose a 27-point “Simpson's Rule” for integration over a cube. If many small cubes fill a large box, why are there only 8 new points per cube?
Cylindrical coordinates are good for describing solids that are symmetric around an axis. The solid is three-dimensional, so there are three coordinates \( r, \theta, z \):

\[
\begin{align*}
  r: & \text{ out from the axis} \\
  \theta: & \text{ around the axis} \\
  z: & \text{ along the axis}.
\end{align*}
\]

This is a mixture of polar coordinates \( r\theta \) in a plane, plus \( z \) upward. You will not find \( r\theta z \) difficult to work with. Start with a cylinder centered on the \( z \) axis:

- **solid cylinder**: \( 0 \leq r \leq 1 \) flat bottom and top: \( 0 \leq z \leq 3 \) half-cylinder: \( 0 \leq \theta \leq \pi \)

Integration over this half-cylinder is \( \int_0^3 \int_0^1 \int_0^r ? \ dr \ d\theta \ dz \). These limits on \( r, \theta, z \) are especially simple. Two other axially symmetric solids are almost as convenient:

- **cone**: integrate to \( r + z = 1 \)  
- **sphere**: integrate to \( r^2 + z^2 = R^2 \)

I would not use cylindrical coordinates for a box. Or a tetrahedron.

The integral needs one thing more—the volume \( dV \). The movements \( dr \) and \( d\theta \) and \( dz \) give a “curved box” in \( xyz \) space, drawn in Figure 14.15c. The base is a polar rectangle, with area \( r \ dr \ d\theta \). The new part is the height \( dz \). The volume of the curved box is \( r \ dr \ d\theta \ dz \). Then \( r \) goes in the blank space in the triple integral—the stretching factor is \( J = r \). There are six orders of integration (we give two):

\[
\begin{align*}
\text{volume} = & \int \int \int r \ dr \ d\theta \ dz = \int \int \int r \ dr \ dz \ d\theta. \quad (1)
\end{align*}
\]

**EXAMPLE 1** (Volume of the half-cylinder). The integral of \( r \ dr \) from 0 to 1 is \( \frac{1}{2} \). The \( \theta \) integral is \( \pi \) and the \( z \) integral is 3. The volume is \( 3\pi/2 \).

**EXAMPLE 2** The surface \( r = 1 - z \) encloses the cone in Figure 14.16. Find its volume.

First solution Since \( r \) goes out to \( 1 - z \), the integral of \( r \ dr \) is \( \frac{1}{2}(1 - z)^2 \). The \( \theta \) integral is \( 2\pi \) (a full rotation). Stop there for a moment.

We have reached \( \int \int r \ dr \ d\theta = \frac{1}{4}(1 - z)^2 \). This is the area of a slice at height \( z \). The slice is a circle, its radius is \( 1 - z \), its area is \( \pi(1 - z)^2 \). The \( z \) integral adds those slices to give \( \pi/3 \). That is correct, but it is not the only way to compute the volume.
Second solution  Do the $z$ and $\theta$ integrals first. Since $z$ goes up to $1-r$, and $\theta$ goes around to $2\pi$, those integrals produce $\int \int r \, dz \, d\theta = r(1-r)2\pi$. Stop again—this must be the area of something.

After the $z$ and $\theta$ integrals we have a shell at radius $r$. The height is $1-r$ (the outer shells are shorter). This height times $2\pi r$ gives the area around the shell. The choice between shells and slices is exactly as in Chapter 8. Different orders of integration give different ways to cut up the solid.

The volume of the shell is area times thickness $dr$. The volume of the complete cone is the integral of shell volumes: $\int_0^1 r(1-r)2\pi \, dr = \frac{\pi}{3}$.

Third solution  Do the $r$ and $z$ integrals first: $\int \int r \, dr \, dz = \frac{1}{3}$. Then the $\theta$ integral is $\int \frac{1}{2} \, d\theta$, which gives $\frac{1}{2}$ times $2\pi$. This is the volume $\frac{\pi}{3}$—but what is $\frac{1}{2} \, d\theta$?

The third cone is cut into wedges. The volume of a wedge is $\frac{1}{2} \, d\theta$. It is quite common to do the $\theta$ integral last, especially when it just multiplies by $2\pi$. It is not so common to think of wedges.

**Question** Is the volume $\frac{1}{2} \, d\theta$ equal to an area $\frac{1}{2}$ times a thickness $d\theta$?

**Answer** No! The triangle in the third cone has area $\frac{1}{2}$ not $\frac{1}{6}$. *Thickness is never $d\theta$.*

---

![Fig. 14.16 A cone cut three ways: slice at height $z$, shell at radius $r$, wedge at angle $\theta$.](image)

This cone is typical of a **solid of revolution**. The axis is in the $z$ direction. The $\theta$ integral yields $2\pi$, whether it comes first, second, or third. The $r$ integral goes out to a radius $f(z)$, which is 1 for the cylinder and $1-z$ for the cone. The integral $\int \int \int r \, dr \, d\theta$ is $\pi f(z)^2 = \text{area of circular slice}$. This leaves the $z$ integral $\int \pi f(z)^2 \, dz$. That is our old volume formula $\int \pi f(z)^2 \, dz$ from Chapter 8, where the slices were cut through the $x$ axis.

**Example 3** The moment of inertia around the $z$ axis is $\int \int \int r^2 \, dr \, d\theta \, dz$. The extra $r^2$ is (distance to axis)$^2$. For the cone this triple integral is $\frac{\pi}{10}$.

**Example 4** The moment around the $z$ axis is $\int \int \int r^2 \, dr \, d\theta \, dz$. For the cone this is $\frac{\pi}{6}$. The average distance $\bar{r}$ is (moment)/(volume) = $\frac{\pi}{6}/(\pi/3) = \frac{1}{2}$.

**Example 5** A sphere of radius $R$ has the boundary $r^2 + z^2 = R^2$, in cylindrical coordinates. The outer limit on the $r$ integral is $\sqrt{R^2 - z^2}$. That is not acceptable in difficult problems. To avoid it we now change to coordinates that are natural for a sphere.
14.4 Cylindrical and Spherical Coordinates

SPHERICAL COORDINATES

The Earth is a solid sphere (or near enough). On its surface we use two coordinates—latitude and longitude. To dig inward or fly outward, there is a third coordinate—the distance \( \rho \) from the center. This Greek letter \( \rho \) replaces \( r \) to avoid confusion with cylindrical coordinates. Where \( r \) is measured from the \( z \) axis, \( \rho \) is measured from the origin. Thus \( r^2 = x^2 + y^2 \) and \( \rho^2 = x^2 + y^2 + z^2 \).

The angle \( \theta \) is the same as before. It goes from 0 to \( 2\pi \). It is the longitude, which increases as you travel east around the Equator.

The angle \( \phi \) is new. It equals 0 at the North Pole and \( \pi \) (not \( 2\pi \)) at the South Pole. It is the polar angle, measured down from the \( z \) axis. The Equator has a latitude of 0 but a polar angle of \( \pi/2 \) (halfway down). Here are some typical shapes:

- solid sphere (or ball): \( 0 \leq \rho \leq R \)
- surface of sphere: \( \rho = R \)
- upper half-sphere: \( 0 < \phi < \pi/2 \)
- eastern half-sphere: \( 0 < \theta < \pi \)

The volume \( dV \) of a spherical box is explained by Figure 14.17. That shows a "spherical box" with right angles and curved edges. Two edges are \( d\rho \) and \( \rho d\phi \). The third edge is horizontal. The usual \( rd\theta \) becomes \( \rho \sin \phi \, d\theta \). Multiplying those lengths gives \( dV \).

The volume integral is explained by Figure 14.17c. That shows a "spherical box" with right angles and curved edges. Two edges are \( d\rho \) and \( \rho d\phi \). The third edge is horizontal. The usual \( rd\theta \) becomes \( \rho \sin \phi \, d\theta \). Multiplying those lengths gives \( dV \).

The volume of the box is \( dV = \rho^2 \sin \phi \, d\rho \, d\phi \, d\theta \). This is a distance cubed, from \( \rho^2 \, d\rho \).
EXAMPLE 6  A solid ball of radius $R$ has known volume $V = \frac{4}{3}\pi R^3$. Notice the limits:

$$\int_0^{2\pi} \int_0^\pi \int_0^R \rho^2 \sin \phi \, d\rho \, d\phi \, d\theta = \left[ \frac{\rho^3}{3} \right]_0^R \left[ -\cos \phi \right]_0^\pi \left[ \theta \right]_0^{2\pi} = \frac{1}{2} R^3(2\pi) = (\frac{1}{2} R^3)(\frac{\pi}{2})(2\pi).$$

**Question**  What is the volume above the cone in Figure 14.17?

**Answer**  The $\phi$ integral stops at $[\cos \phi]_0^\pi = \frac{1}{2}$. The volume is $(\frac{1}{2} R^3)(\frac{\pi}{4})(2\pi)$.

EXAMPLE 7  The surface area of a sphere is $A = 4\pi R^2$. Forget the $\rho$ integral:

$$A = \int_0^{2\pi} \int_0^\pi R^2 \sin \phi \, d\phi \, d\theta = R^2 \left[ -\cos \phi \right]_0^\pi \left[ \theta \right]_0^{2\pi} = R^2(2\pi).$$

*After those examples from geometry, here is the real thing from science.* I want to compute one of the most important triple integrals in physics—"the gravitational attraction of a solid sphere." For some reason Isaac Newton had trouble with this integral. He refused to publish his masterpiece on astronomy until he had solved it. I think he didn’t use spherical coordinates—and the integral is not easy even now.

The answer that Newton finally found is beautiful. The sphere acts as if all its mass were concentrated at the center. At an outside point $(0, 0, D)$, the force of gravity is proportional to $1/D^2$. The force from a uniform solid sphere equals the force from a point mass, at every outside point $P$. That is exactly what Newton wanted and needed, to explain the solar system and to prove Kepler’s laws.

Here is the difficulty. Some parts of the sphere are closer than $D$, some parts are farther away. The actual distance $q$, from the outside point $P$ to a typical inside point, is shown in Figure 14.18. The average distance $\bar{q}$ to all points in the sphere is not $D$. But what Newton needed was a different average, and by good luck or some divine calculus it works perfectly: The average of $1/q$ is $1/D$. This gives the potential energy:

$$\text{potential at point } P = \iiint_{\text{sphere}} \frac{1}{q} \, dV = \frac{\text{volume of sphere}}{D}. \quad (2)$$

A small volume $dV$ at the distance $q$ contributes $dV/q$ to the potential (Section 8.6, with density 1). The integral adds the contributions from the whole sphere. Equation (2) says that the potential at $r = D$ is not changed when the sphere is squeezed to the center. The potential equals the whole volume divided by the single distance $D$.

Important point: The average of $1/q$ is $1/D$ and not $1/\bar{q}$. The average of $\frac{1}{2}$ and $\frac{1}{4}$ is not $\frac{1}{3}$. Smaller point: I wrote "sphere" where I should have written "ball." The sphere is solid: $0 \leq \rho \leq R$, $0 \leq \phi \leq \pi$, $0 \leq \theta \leq 2\pi$.

What about the force? For the small volume it is proportional to $dV/q^2$ (this is the inverse square law). But force is a vector, pulling the outside point toward $dV$—not toward the center of the sphere. The figure shows the geometry and the symmetry. **We want the $z$ component of the force.** (By symmetry the overall $x$ and $y$ components are zero.) The angle between the force vector and the $z$ axis is $\alpha$, so for the $z$ component we multiply by $\cos \alpha$. The total force comes from the integral that Newton discovered:

$$\text{force at point } P = \iiint_{\text{sphere}} \frac{\cos \alpha}{q^2} \, dV = \frac{\text{volume of sphere}}{D^2}. \quad (3)$$

I will compute the integral (2) and leave you the privilege of solving (3). I mean that word seriously. If you have come this far, you deserve the pleasure of doing what at
one time only Isaac Newton could do. Problem 26 offers a suggestion (just the law
of cosines) but the integral is yours.

The law of cosines also helps with (2). For the triangle in the figure it gives
\[ q^2 = D^2 - 2D\rho \cos \phi + \rho^2. \]
Call this whole quantity \( u \). We do the surface integral first
\((d\phi \text{ and } d\theta \text{ with } \rho \text{ fixed}). \) Then \( q^2 = u \) and \( q = \sqrt{u} \) and \( du = 2\rho D \sin \phi \ d\phi \):

\[
\int_0^{2\pi} \int_0^\phi \frac{\rho^2 \sin \phi \ d\phi \ d\theta}{q} = \int_0^{2\pi} \int_0^{2\rho D} \frac{du}{2\rho D \sqrt{u}} = \left[ \frac{2\rho \sqrt{u}}{D} \right]_0^{\phi = \pi}.
\]

\( 2\pi \) came from the \( \theta \) integral. The integral of \( du/\sqrt{u} \) is \( 2\sqrt{u} \). Since \( \cos \phi = -1 \) at the upper limit, \( u \) is \( D^2 + 2\rho D + \rho^2 \). The square root of \( u \) is \( D + \rho \). At the lower limit \( \cos \phi = +1 \) and \( u = D^2 - 2\rho D + \rho^2 \). This is another perfect square—its square root
is \( D - \rho \). The surface integral (4) with fixed \( \rho \) is

\[
\int \int \frac{dA}{q} = \frac{2\pi \rho}{D} \left[ (D + \rho) - (D - \rho) \right] = \frac{4\pi \rho^2}{D}.
\]

Last comes the \( \rho \) integral: \( \int_0^\rho 4\pi \rho^2 \ d\rho / D = \frac{4}{3}\pi R^3 / D \). This proves formula (2):

**potential equals volume of the sphere divided by \( D \).**

**Note 1** Physicists are also happy about equation (5). The average of \( 1/q \) is \( 1/D \) not
only over the solid sphere but over each spherical shell of area \( 4\pi \rho^2 \). The shells can
have different densities, as they do in the Earth, and still Newton is correct. This also
applies to the force integral (3)—each separate shell acts as if its mass were concentra-
ted at the center. Then the final \( \rho \) integral yields this property for the solid sphere.

**Note 2** Physicists also know that force is minus the derivative of potential. The
derivative of (2) with respect to \( D \) produces the force integral (3). Problem 27 explains
this shortcut to equation (3).

**EXAMPLE 8** **Everywhere inside a hollow sphere the force of gravity is zero.**

When \( D \) is smaller than \( \rho \), the lower limit \( \sqrt{u} \) in the integral (4) changes from \( D - \rho \) to \( \rho - D \). That way the square root stays positive. This changes the answer in (5) to
\( 4\pi \rho^2 / \rho \), so the potential no longer depends on \( D \). The potential is constant inside the
hollow shell. Since the force comes from its derivative, the force is zero.

A more intuitive proof is in the second figure. The infinitesimal areas on the surface
are proportional to \( q^2 \) and \( Q^2 \). But the distances to those areas are \( q \) and \( Q \), so the
forces involve $1/q^2$ and $1/Q^2$ (the inverse square law). Therefore the two areas exert equal and opposite forces on the inside point, and they cancel each other. The total force from the shell is zero.

I believe this zero integral is the reason that the inside of a car is safe from lightning. Of course a car is not a sphere. But electric charge distributes itself to keep the surface at constant potential. The potential stays constant inside—therefore no force. The tires help to prevent conduction of current (and electrocution of driver).

P.S. Don't just step out of the car. Let a metal chain conduct the charge to the ground. Otherwise you could be the conductor.

**CHANGE OF COORDINATES—STRETCHING FACTOR $J$**

Once more we look to calculus for a formula. We need the volume of a small curved box in any $uvw$ coordinate system. The $r \theta z$ box and the $\rho \phi \theta$ box have right angles, and their volumes were read off from the geometry (stretching factors $J = r$ and $J = \rho^2 \sin \phi$ in Figures 14.15 and 14.17). Now we change from $xyz$ to other coordinates $uvw$—which are chosen to fit the problem.

Going from $xy$ to $uv$, the area $dA = J \, du \, dv$ was a 2 by 2 determinant. In three dimensions the determinant is 3 by 3. The matrix is always the "Jacobian matrix," containing first derivatives. There were four derivatives from $xy$ to $uv$, now there are nine from $xyz$ to $uvw$.

**EXAMPLE 9** Find the factor $J$ for $x = \rho \sin \phi \cos \theta$, $y = \rho \sin \phi \sin \theta$, $z = \rho \cos \phi$.

$$J = \left| \begin{array}{ccc} \frac{\partial x}{\partial \rho} & \frac{\partial x}{\partial \phi} & \frac{\partial x}{\partial \theta} \\ \frac{\partial y}{\partial \rho} & \frac{\partial y}{\partial \phi} & \frac{\partial y}{\partial \theta} \\ \frac{\partial z}{\partial \rho} & \frac{\partial z}{\partial \phi} & \frac{\partial z}{\partial \theta} \end{array} \right| = \frac{\partial(x, y, z)}{\partial(u, v, w)}.$$

Remember that a 3 by 3 determinant is the sum of six terms (Section 11.5). One term in $J$ is $(\partial x/\partial \rho)(\partial y/\partial \phi)(\partial z/\partial \theta)$, along the main diagonal. This comes from pure stretching, and the other five terms allow for rotation. The best way to exhibit the formula is for spherical coordinates—where the nine derivatives are easy but the determinant is not:

The determinant has six terms, but two are zero—because of the zero in the corner. The other four terms are $\rho^2 \sin \phi \cos^2 \phi \sin^2 \theta$ and $\rho^2 \sin \phi \cos^2 \phi \cos^2 \theta$ and $\rho^2 \sin^3 \phi \sin^2 \theta$ and $\rho^2 \sin^3 \phi \cos^2 \theta$. Add the first two (note $\sin^2 \theta + \cos^2 \theta$) and separately add the second two. Then add the sums to reach $J = \rho^2 \sin \phi$.

Geometry already gave this answer. For most $uvw$ variables, use the determinant.
14.4 Cylindrical and Spherical Coordinates

Read-through questions

The three ___ coordinates are rθz. The point at x = y = z = 1 has r = ____, θ = ____, z = ____. The volume integral is ____. The solid region 1 ≤ r ≤ 2, 0 ≤ θ ≤ 2π, 0 ≤ z ≤ 4 is a ____. Its volume is ____. From the r and θ integrals the area of a ____ equals ____. From the z and θ integrals the area of a ____ equals ____. In rθz coordinates the shapes of ____ are convenient, while ____ are not.

The three ___ coordinates are ρφθ. The point at x = y = z = 1 has ρ = ____, φ = ____, θ = ____. The angle φ is measured from ____, θ is measured from ____, ρ is the distance to ____. If ρφθ are known then x = ____, y = ____, z = ____. The stretching factor J is a 3 by 3 ____ and volume is ____. The solid region 1 ≤ ρ ≤ 2, 0 ≤ φ ≤ π, 0 ≤ θ ≤ 2π is a ____. Its volume is ____. From the φ and θ integrals the area of a ____ at radius ρ equals ____. Newton discovered that the outside gravitational attraction of a ____ is the same as for an equal mass located at ____.

Convert the xyz coordinates in 1–4 to rθz and ρφθ.

1 (D, 0, 0) 2 (0, −D, 0)
3 (0, 0, D) (watch θ) 4 (3, 4, 5)

Convert the spherical coordinates in 5–7 to xyz and rθz.

5 ρ = 4, φ = π/4, θ = −π/4
6 ρ = 2, φ = π/3, θ = π/6
7 ρ = 1, φ = π, θ = anything.
8 Where does x = r and y = θ?
9 Find the polar angle φ for the point with cylindrical coordinates rθz.
10 What are x(t), y(t), z(t) on the great circle from ρ = 1, φ = π/2, θ = 0 with speed 1 to ρ = 1, φ = π/4, θ = π/2?

From the limits of integration describe each region in 11–20 and find its volume. The inner integral has the inner limits.

11 ∫0^2π ∫0^1 ∫r^1−r r dz dr dθ
12 ∫0^π ∫0^1 ∫0^r+1 r dz dr dθ
13 ∫0^2π ∫0^1 ∫0^1−z r dz dθ dθ
14 ∫0^π ∫0^1 ∫0^r r dθ dr dz
15 ∫0^π ∫0^2 ∫0^1 ρ^2 sin φ dρ dφ dθ
16 ∫0^2π ∫0^1 ∫0^ρ sin φ dρ dφ dθ
17 ∫0^π ∫0^1 ∫0^ρ^2 sin φ dφ dθ dθ
18 ∫0^2π ∫0^1 ∫0^ρ^2 sin φ dφ dθ dθ
19 ∫0^π ∫0^1 ∫0^ρ^2 sin φ dφ dθ dθ
20 ∫0^π ∫0^1 ∫0^ρ^2 sin φ dφ dθ dθ

21 Example 5 gave the volume integral for a sphere in rθz coordinates. What is the area of the circular slice at height z? The outside gravitational attraction of a ____ is the same as for an equal mass located at ____.

22 Describe the solid with 0 ≤ ρ ≤ 1 − cos φ and find its volume.

23 A cylindrical tree has radius a. A saw cuts horizontally, ending halfway in at the x axis. Then it cuts on a sloping plane (angle θ with the horizontal), also ending at the x axis. What is the volume of the wedge that falls out?

24 Find the mass of a planet of radius R, if its density at each radius ρ is ρ = (ρ + 1)/ρ. Notice the infinite density at the center, but finite mass M = ||δ dV. Here ρ is radius, not density.

25 For the cone out to r = 1 − z, the average distance from the z axis is r̄ = ™. For the triangle out to r = 1 − z the average is ™ = ™. How can they be different when rotating the triangle produces the cone?

Problems 26–32, on the attraction of a sphere, use Figure 14.18 and the law of cosines ρ^2 = D^2 − 2ρD cos φ + ρ^2 = u.

26 Newton's achievement Show that ||(cos α) dV/ρ^2 equals volume/D^2. One hint only: Find cos α from a second law of cosines ρ^2 = D^2 − 2ρD cos α + ρ^2. The φ integral should involve 1/ρ and 1/ρ^3. Equation (2) integrates 1/ρ, leaving ||∫(cos α) dV/ρ^3 still to do.

27 Compute δq/δD in the first cosine law and show from Figure 14.18 that it equals cos α. Then the derivative of equation (2) with respect to D is a shortcut to Newton's equation (3).

28 The lines of length D and q meet at the angle α. Move the meeting point up by ΔD. Explain why the other line stretches by Δq ≈ ΔD cos α. So δq/δD = cos α as before.
29 Show that the average distance is \( q = 4R/3 \), from the North Pole \((D = R)\) to points on the Earth's surface \((p = R)\).
To compute: \( q = \iiint [q R^2 \sin \phi \ d\phi \ d\theta] / (\text{area } 4\pi R^2) \). Use the same substitution \( u \).

30 Show as in Problem 29 that the average distance is \( q = D + \frac{1}{3} p^2 / D \), from the outside point \((0, 0, D)\) to points on the shell of radius \( p \). Then integrate \( \iiint q \, dV \) and divide by \( 4\pi R^3 / 3 \) to find \( q \) for the solid sphere.

31 In Figure 14.18b, it is not true that the areas on the surface are exactly proportional to \( q^2 \) and \( Q^2 \). Why not? What happens to the second proof in Example 8?

32 For two solid spheres attracting each other (sun and planet), can we concentrate both spheres into point masses at their centers?

33 Compute \( \iiint \cos \alpha \, dV / q^3 \) to find the force of gravity at \((0, 0, D)\) from a cylinder \( x^2 + y^2 \leq a^2 \), \( 0 \leq z \leq h \). Show from a figure why \( q^2 = r^2 + (D - z)^2 \) and \( \cos \alpha = (D - z)/q \).

34 A linear change of variables has \( x = au + bv + cw \), \( y = du + ev + fw \), and \( z = gu + hv + iw \). Write down the six terms in the determinant \( J \). Three terms have minus signs.

35 A pure stretching has \( x = au \), \( y = bv \), and \( z = cw \). Find the 3 by 3 matrix and its determinant \( J \). What is special about the \( xyz \) box in this case?

36 (a) The matrix in Example 9 has three columns. Find the lengths of those three vectors (sum of squares, then square root). Compare with the edges of the box in Figure 14.17.
(b) Take the dot product of every column in \( J \) with every other column. Zero dot products mean right angles in the box. So \( J \) is the product of the column lengths.
37 Find the stretching factor \( J \) for cylindrical coordinates from the matrix of first derivatives.
38 Follow Problem 36 for cylindrical coordinates—find the length of each column in \( J \) and compare with the box in Figure 14.15.
39 Find the moment of inertia around the \( z \) axis of a spherical shell (radius \( p \), density \( 1 \)). The distance from the axis to a point on the shell is \( r = \ldots \). Substitute for \( r \) to find
\[
I(p) = \int_0^R \int_0^\pi r^2 \rho^2 \sin \phi \ d\phi \ d\theta.
\]
Divide by \( mr^2 \) (which is \( 4\pi \rho^4 \)) to compute the number \( J \) for a hollow ball in the rolling experiment of Section 8.5.

40 The moment of inertia of a solid sphere (radius \( R \), density \( 1 \)) adds up the hollow spheres of Problem 39:
\[
I = \int_0^R \int_0^\pi r^2 (p^2) \, d\phi \, dp = \ldots.
\]
Divide by \( mR^2 \) (which is \( \frac{4}{3} \pi R^3 \)) to find \( J \) in the rolling experiment. A solid ball rolls faster than a hollow ball because _______.

41 Inside the Earth, the force of gravity is proportional to the distance \( \rho \) from the center. Reason: The inner ball of radius \( \rho \) has mass proportional to ________ (assume constant density). The force is proportional to that mass divided by ________. The rest of the Earth (sphere with hole) exerts no force because ________.

42 Dig a tunnel through the center to Australia. Drop a ball in the tunnel at \( y = R \); Australia is \( y = -R \). The force of gravity is \(-cy\) by Problem 41. Newton's law is \( my'' = -cy \).
What does the ball do when it reaches Australia?
Chapter 14 introduced double and triple integrals. We went from \( \int dx \) to \( \iint dx\,dy \) and \( \iiint dx\,dy\,dz \). All those integrals add up small pieces, and the limit gives area or volume or mass. What could be more natural than that? I regret to say, after the success of those multiple integrals, that something is missing. It is even more regrettable that we didn't notice it. The missing piece is nothing less than the Fundamental Theorem of Calculus.

The double integral \( \iint dx\,dy \) equals the area. To compute it, we did not use an antiderivative of 1. At least not consciously. The method was almost trial and error, and the hard part was to find the limits of integration. This chapter goes deeper, to show how the step from a double integral to a single integral is really a new form of the Fundamental Theorem—when it is done right.

Two new ideas are needed early, one pleasant and one not. You will like vector fields. You may not think so highly of line integrals. Those are ordinary single integrals like \( \int v(x)\,dx \), but they go along curves instead of straight lines. The nice step \( dx \) becomes the confusing step \( ds \). Where \( \int dx \) equals the length of the interval, \( \int ds \) is the length of the curve. The point is that regions are enclosed by curves, and we have to integrate along them. The Fundamental Theorem in its two-dimensional form (Green's Theorem) connects a double integral over the region to a single integral along its boundary curve.

The great applications are in science and engineering, where vector fields are so natural. But there are changes in the language. Instead of an antiderivative, we speak about a potential function. Instead of the derivative, we take the "divergence" and "curl." Instead of area, we compute flux and circulation and work. Examples come first.

15.1 Vector Fields

For an ordinary scalar function, the input is a number \( x \) and the output is a number \( f(x) \). For a vector field (or vector function), the input is a point \( (x, y) \) and the output is a two-dimensional vector \( F(x, y) \). There is a "field" of vectors, one at every point.
In three dimensions the input point is \((x, y, z)\) and the output vector \(\mathbf{F}\) has three components.

**DEFINITION**  Let \(R\) be a region in the \(xy\) plane. A **vector field** \(\mathbf{F}\) assigns to every point \((x, y)\) in \(R\) a vector \(\mathbf{F}(x, y)\) with two components:

\[
\mathbf{F}(x, y) = M(x, y)i + N(x, y)j. \tag{1}
\]

This plane vector field involves two functions of two variables. They are the components \(M\) and \(N\), which vary from point to point. A vector has fixed components, a vector field has varying components.

A three-dimensional vector field has components \(M(x, y, z)\) and \(N(x, y, z)\) and \(P(x, y, z)\). Then the vectors are \(\mathbf{F} = M\mathbf{i} + N\mathbf{j} + P\mathbf{k}\).

**EXAMPLE 1** The **position vector** at \((x, y)\) is \(\mathbf{R} = xi + yj\). Its components are \(M = x\) and \(N = y\). The vectors grow larger as we leave the origin (Figure 15.1a). Their direction is outward and their length is \(|\mathbf{R}| = \sqrt{x^2 + y^2} = r\). The vector \(\mathbf{R}\) is boldface, the number \(r\) is lightface.

**EXAMPLE 2** The vector field \(\mathbf{R}/r\) consists of **unit vectors** \(\mathbf{u}_r\), pointing outward. We divide \(\mathbf{R} = xi + yj\) by its length, at every point except the origin. The components of \(\mathbf{R}/r\) are \(M = x/r\) and \(N = y/r\). Figure 15.1 shows a third field \(\mathbf{R}/r^2\), whose length is \(1/r\).

![Fig. 15.1](image)

**EXAMPLE 3** The **spin field** or rotation field or turning field goes around the origin instead of away from it. The field is \(\mathbf{S}\). Its components are \(M = -y\) and \(N = x\):

\[
\mathbf{S} = -yi + xj \text{ also has length } |\mathbf{S}| = \sqrt{(-y)^2 + x^2} = r. \tag{2}
\]

\(\mathbf{S}\) is perpendicular to \(\mathbf{R}\)—their dot product is zero: \(\mathbf{S} \cdot \mathbf{R} = (-y)(x) + (x)(y) = 0\). The spin fields \(\mathbf{S}/r\) and \(\mathbf{S}/r^2\) have lengths \(1\) and \(1/r\):

\[
\mathbf{S} = -yi + xj \text{ has } \left|\frac{\mathbf{S}}{r}\right| = 1 \quad \frac{\mathbf{S}}{r} = -\frac{y}{r^2}i + \frac{x}{r^2}j \text{ has } \left|\frac{\mathbf{S}}{r^2}\right| = \frac{1}{r}.
\]

The unit vector \(\mathbf{S}/r\) is \(\mathbf{u}_s\). Notice the blank at \((0, 0)\), where this field is not defined.

![Fig. 15.2](image)
EXAMPLE 4  A gradient field starts with an ordinary function \( f(x, y) \). The components \( M \) and \( N \) are the partial derivatives \( \frac{\partial f}{\partial x} \) and \( \frac{\partial f}{\partial y} \). Then the field \( F \) is the gradient of \( f \):

\[
F = \text{grad } f = \nabla f = \frac{\partial f}{\partial x} \mathbf{i} + \frac{\partial f}{\partial y} \mathbf{j}.
\]

(3)

This vector field \( \text{grad } f \) is everywhere perpendicular to the level curves \( f(x, y) = c \). The length \( |\text{grad } f| \) tells how fast \( f \) is changing (in the direction it changes fastest). Invent a function like \( f = x^2y \), and you immediately have its gradient field \( F = 2xy \mathbf{i} + x^2 \mathbf{j} \). To repeat, \( M = \frac{\partial f}{\partial x} \) and \( N = \frac{\partial f}{\partial y} \).

For every vector field you should ask two questions: Is it a gradient field? If so, what is \( f \)? Here are answers for the radial fields and spin fields:

<table>
<thead>
<tr>
<th>15A</th>
<th>The radial fields ( \mathbf{R} ) and ( \mathbf{R}/r ) and ( \mathbf{R}/r^2 ) are all gradient fields.</th>
</tr>
</thead>
<tbody>
<tr>
<td></td>
<td>The spin fields ( \mathbf{S} ) and ( \mathbf{S}/r ) are not gradients of any ( f(x, y) ).</td>
</tr>
<tr>
<td></td>
<td>The spin field ( \mathbf{S}/r^2 ) is the gradient of the polar angle ( \theta = \tan^{-1}(y/x) ).</td>
</tr>
</tbody>
</table>

The derivatives of \( f = \frac{1}{2}(x^2 + y^2) \) are \( x \) and \( y \). Thus \( \mathbf{R} \) is a gradient field. The gradient of \( f = r \) is the unit vector \( \mathbf{R}/r \) pointing outwards. Both fields are perpendicular to circles around the origin. Those are the level curves of \( f = \frac{1}{2} r^2 \) and \( f = r \).

**Question** Is every \( \mathbf{R}/r^2 \) a gradient field?

**Answer** Yes. But among the spin fields, the only gradient is \( \mathbf{S}/r^2 \).

A major goal of this chapter is to recognize gradient fields by a simple test. The rejection of \( \mathbf{S} \) and \( \mathbf{S}/r \) will be interesting. For some reason \( -y \mathbf{i} + x \mathbf{j} \) is rejected and \( y \mathbf{i} + x \mathbf{j} \) is accepted. (It is the gradient of \( ________ \).) The acceptance of \( \mathbf{S}/r^2 \) as the gradient of \( f = \theta \) contains a surprise at the origin (Section 15.3).

Gradient fields are called *conservative*. The function \( f \) is the *potential function*. These words, and the next examples, come from physics and engineering.

EXAMPLE 5  The velocity field is \( \mathbf{V} \) and the flow field is \( \rho \mathbf{V} \).

Suppose fluid moves steadily down a pipe. Or a river flows smoothly (no waterfall). Or the air circulates in a fixed pattern. The velocity can be different at different points, but there is no change with time. The velocity vector \( \mathbf{V} \) gives the *direction of flow* and *speed of flow* at every point.

In reality the velocity field is \( \mathbf{V}(x, y, z) \), with three components \( M, N, P \). Those are the velocities \( v_1, v_2, v_3 \) in the \( x, y, z \) directions. The speed \( |\mathbf{V}| \) is the length: \( |\mathbf{V}|^2 = v_1^2 + v_2^2 + v_3^2 \). In a “plane flow” the \( k \) component is zero, and the velocity field is \( v_1 \mathbf{i} + v_2 \mathbf{j} = M \mathbf{i} + N \mathbf{j} \).

[Fig. 15.3] A steady velocity field \( \mathbf{V} \) and two force fields \( \mathbf{F} \).
15 Vector Calculus

For a compact disc or a turning wheel, $V$ is a spin field ($V = \omega S$, $\omega =$ angular velocity). A tornado might be closer to $V = S/r^2$ (except for a dead spot at the center). An explosion could have $V = R/r^2$. A quieter example is flow in and out of a lake with steady rain as a source term.

The flow field $\rho V$ is the density $\rho$ times the velocity field. While $V$ gives the rate of movement, $\rho V$ gives the rate of movement of mass. A greater density means a greater rate $|\rho V|$ of "mass transport." It is like the number of passengers on a bus times the speed of the bus.

**EXAMPLE 6** Force fields from gravity: $F$ is downward in the classroom, $F$ is radial in space.

When gravity pulls downward, it has only one nonzero component: $F = -mgk$. This assumes that vectors to the center of the Earth are parallel—almost true in a classroom. Then $F$ is the gradient of $-mgz$ (note $\partial f/\partial z = -mg$).

In physics the usual potential is not $-mgz$ but $+mgz$. The force field is minus grad $f$ also in electrical engineering. Electrons flow from high potential to low potential. The mathematics is the same, but the sign is reversed.

In space, the force is radial inwards: $F = -mMGR/r^3$. Its magnitude is proportional to $1/r^2$ (Newton's inverse square law). The masses are $m$ and $M$, and the gravitational constant is $G = 6.672 \times 10^{-11}$—with distance in meters, mass in kilograms, and time in seconds. The dimensions of $G$ are (force)(distance)$^2$(mass). This is different from the acceleration $g = 9.8m/sec^2$, which already accounts for the mass and radius of the Earth.

Like all radial fields, gravity is a gradient field. It comes from a potential $f$:

$$f = \frac{mMG}{r} \quad \text{and} \quad \frac{\partial f}{\partial x} = -\frac{mMGx}{r^3} \quad \text{and} \quad \frac{\partial f}{\partial y} = -\frac{mMGy}{r^3} \quad \text{and} \quad \frac{\partial f}{\partial z} = -\frac{mMGz}{r^3}. \quad (4)$$

**EXAMPLE 7** (a short example) Current in a wire produces a magnetic field $B$. It is the spin field $S/r^2$ around the wire, times the strength of the current.

**STREAMLINES AND LINES OF FORCE**

Drawing a vector field is not always easy. Even the spin field looks messy when the vectors are too long (they go in circles and fall across each other). The circles give a clearer picture than the vectors. In any field, the vectors are tangent to "field lines"—which in the spin case are circles.

**DEFINITION** $C$ is a field line or integral curve if the vectors $\mathbf{F}(x, y)$ are tangent to $C$. The slope $dy/dx$ of the curve $C$ equals the slope $N/M$ of the vector $\mathbf{F} = Mi + NJ$:

$$\frac{dy}{dx} = \frac{N(x, y)}{M(x, y)} \quad \left( = \frac{x}{y} \text{ for the spin field} \right). \quad (6)$$

We are still drawing the field of vectors, but now they are infinitesimally short. They are connected into curves! What is lost is their length, because $S$ and $S/r$ and $S/r^2$ all have the same field lines (circles). For the position field $\mathbf{R}$ and gravity field $\mathbf{R}/r^3$, the field lines are rays from the origin. In this case the "curves" are actually straight.

**EXAMPLE 8** Show that the field lines for the velocity field $V = yi + xj$ are hyperbolas.

$$\frac{dy}{dx} = \frac{N}{M} = \frac{x}{y} \quad \Rightarrow \quad y \frac{dy}{dx} = x \quad \Rightarrow \quad \frac{1}{2} y^2 - \frac{1}{2} x^2 = \text{constant.}$$
At every point these hyperbolas line up with the velocity \( \mathbf{V} \). Each particle of fluid travels on a field line. In fluid flow those hyperbolas are called streamlines. Drop a leaf into a river, and it follows a streamline. Figure 15.4 shows the streamlines for a river going around a bend.

Don't forget the essential question about each vector field. Is it a gradient field? For \( \mathbf{V} = y \mathbf{i} + x \mathbf{j} \) the answer is yes, and the potential is \( f = xy \):

\[
\frac{\partial f}{\partial x} = y, \quad \frac{\partial f}{\partial y} = x.
\]

When there is a potential, it has level curves. They connect points of equal potential, so the curves \( f(x, y) = c \) are called equipotentials. Here they are the curves \( xy = c \)—also hyperbolas. Since gradients are perpendicular to level curves, the streamlines are perpendicular to the equipotentials. Figure 15.4 is sliced one way by streamlines and the other way by equipotentials.

A gradient field \( \mathbf{F} = \frac{\partial f}{\partial x} \mathbf{i} + \frac{\partial f}{\partial y} \mathbf{j} \) is tangent to the field lines (streamlines) and perpendicular to the equipotentials (level curves of \( f \)).

In the gradient direction \( f \) changes fastest. In the level direction \( f \) doesn't change at all. The chain rule along \( f(x, y) = c \) proves these directions to be perpendicular:

\[
\frac{\partial f}{\partial x} \frac{dx}{dt} + \frac{\partial f}{\partial y} \frac{dy}{dt} = 0 \quad \text{or} \quad (\text{grad } f)^\times (\text{tangent to level curve}) = 0.
\]

**EXAMPLE 9** The streamlines of \( \mathbf{S}/r^2 \) are circles around \((0, 0)\). The equipotentials are rays \( \theta = c \). Add rays to Figure 15.2 for the gradient field \( \mathbf{S}/r^2 \).

For the gravity field those are reversed. A body is pulled in along the field lines (rays). The equipotentials are the circles where \( f = 1/r \) is constant. The plane is crisscrossed by "orthogonal trajectories"—curves that meet everywhere at right angles.

If you bring a magnet near a pile of iron filings, a little shake will display the field lines. In a force field, they are "lines of force." Here are the other new words.

<table>
<thead>
<tr>
<th>Vector field ( \mathbf{F}(x, y, z) = M \mathbf{i} + N \mathbf{j} + P \mathbf{k} )</th>
<th>Plane field ( \mathbf{F} = M(x, y) \mathbf{i} + N(x, y) \mathbf{j} )</th>
</tr>
</thead>
<tbody>
<tr>
<td>Radial field: multiple of ( \mathbf{R} = x \mathbf{i} + y \mathbf{j} + z \mathbf{k} )</td>
<td>Spin field: multiple of ( \mathbf{S} = -y \mathbf{i} + x \mathbf{j} )</td>
</tr>
<tr>
<td>Gradient field = conservative field: ( M = \frac{\partial f}{\partial x}, N = \frac{\partial f}{\partial y}, P = \frac{\partial f}{\partial z} )</td>
<td></td>
</tr>
<tr>
<td>Potential ( f(x, y) ) (not a vector)</td>
<td>Equipotential curves ( f(x, y) = c )</td>
</tr>
<tr>
<td>Streamline = field line = integral curve: a curve that has ( \mathbf{F}(x, y) ) as its tangent vectors.</td>
<td></td>
</tr>
</tbody>
</table>
15 Vector Calculus

15.1 Exercises

Read-through questions

A vector field assigns a \( \mathbf{a} \) to each point \((x, y)\) or \((x, y, z)\). In two dimensions \( \mathbf{F}(x, y) = -y\mathbf{i} + x\mathbf{j} \). An example is the position field \( \mathbf{R} = \mathbf{a} \). Its magnitude is \( |\mathbf{R}| = \mathbf{a} \) and its direction is \( \mathbf{a} \). It is the gradient field for \( f = y \). The level curves are \( y = \) constant. They are \( y \) to the vectors \( \mathbf{R} \).

Reversing this picture, the spin field \( \mathbf{S} = \mathbf{i} \). Its magnitude is \( |\mathbf{S}| = \mathbf{k} \) and its direction is \( \mathbf{k} \). It is not a gradient field, because no function has \( \partial f/\partial x = \mathbf{m} \) and \( \partial f/\partial y = \mathbf{n} \). \( \mathbf{S} \) is the velocity field for flow going \( \mathbf{n} \). The streamlines or \( \mathbf{m} \) lines or integral \( \mathbf{a} \) are \( \mathbf{a} \). The flow field \( \rho \mathbf{V} \) gives the rate at which \( \mathbf{a} \) is moved by the flow.

A gravity field from the origin is proportional to \( \mathbf{F} = \mathbf{i} \) which has \( |\mathbf{F}| = \mathbf{u} \). This is Newton’s \( \mathbf{y} \) square law. It is a gradient field, with potential \( f = \mathbf{w} \). The equipotential curves \( f(x, y) = c \) are \( \mathbf{u} \). They are \( \mathbf{u} \) to the field lines which are \( \mathbf{y} \). This illustrates that the \( \mathbf{A} \) of a function \( f(x, y) \) is \( \mathbf{b} \) to its level curves.

The velocity field \( y\mathbf{i} + x\mathbf{j} \) is the gradient of \( f = \mathbf{c} \). Its streamlines are \( \mathbf{y} \). The slope \( dy/dx \) of a streamline equals the ratio \( \mathbf{e} \) of velocity components. The field is \( \mathbf{f} \) to the streamlines. Drop a leaf onto the flow, and it goes along \( \mathbf{c} \).

Find a potential \( f(x, y) \) for the gradient fields 1–8. Draw the streamlines perpendicular to the equipotentials \( f(x, y) = c \).

1. \( \mathbf{F} = i + 2j \) (constant field) 2. \( \mathbf{F} = xi + j \)
3. \( \mathbf{F} = \cos(x + y)\mathbf{i} + \cos(x + y)\mathbf{j} \)
4. \( \mathbf{F} = (1/y)\mathbf{i} - (x/y^2)\mathbf{j} \)
5. \( \mathbf{F} = (2x + 2y)/(x^2 + y^2) \)
6. \( \mathbf{F} = x^2\mathbf{i} + y^2\mathbf{j} \)
7. \( \mathbf{F} = xy\mathbf{i} + \mathbf{j} \)
8. \( \mathbf{F} = \sqrt{y}\mathbf{i} + \mathbf{j} \)

9. Draw the shear field \( \mathbf{F} = x\mathbf{j} \). Check that it is not a gradient field: If \( \partial f/\partial x = 0 \) then \( \partial f/\partial y = x \) is impossible. What are the streamlines (field lines) in the direction of \( \mathbf{F} \)?

10. Find all functions that satisfy \( \partial f/\partial x = -y \) and show that none of them satisfy \( \partial f/\partial y = x \). Then the spin field \( \mathbf{S} = -yi + xj \) is not a gradient field.

Compute \( \partial^2 f/\partial x^2 \) and \( \partial^2 f/\partial y^2 \) in 11–18. Draw the gradient field \( \mathbf{F} = \text{grad} f \) and the equipotentials \( f(x, y) = c \).

11. \( f = 3x + y \) 12. \( f = x - 3y \)
13. \( f = x + y^2 \) 14. \( f = (x - 1)^2 + y^2 \)
15. \( f = x^2 - y^2 \) 16. \( f = e^x \cos y \)
17. \( f = e^{-y} \) 18. \( f = y/x \)

Find equations for the streamlines in 19–24 by solving \( dy/dx = N/M \) (including a constant \( C \)). Draw the streamlines.

19. \( \mathbf{F} = i - j \) 20. \( \mathbf{F} = i + xj \)
21. \( \mathbf{F} = \mathbf{S} \) (spin field) 22. \( \mathbf{F} = \mathbf{S}/r \) (spin field)
23. \( \mathbf{F} = \text{grad} (x/y) \) 24. \( \mathbf{F} = \text{grad} (2x + y) \).

25. The Earth’s gravity field is radial, but in a room the field lines seem to go straight down into the floor. This is because nearby field lines always look _______.

26. A line of charges produces the electrostatic force field \( \mathbf{F} = \mathbf{R}/r^2 = (xi + yj)/(x^2 + y^2) \). Find the potential \( f(x, y) \). (\( \mathbf{F} \) is also the gravity field for a line of masses.)

In 27–32 write down the vector fields \( M\mathbf{i} + N\mathbf{j} \).

27. \( \mathbf{F} \) points radially away from the origin with magnitude 5.

28. The velocity is perpendicular to the curves \( x^3 + y^3 = c \) and the speed is 1.

29. The gravitational force \( \mathbf{F} \) comes from two unit masses at \((0, 0)\) and \((1, 0)\).

30. The streamlines are in the 45° direction and the speed is 4.

31. The streamlines are circles clockwise around the origin and the speed is 1.

32. The equipotentials are the parabolas \( y = x^2 + c \) and \( \mathbf{F} \) is a gradient field.

33. Show directly that the hyperbolas \( xy = 2 \) and \( x^2 - y^2 = 3 \) are perpendicular at the point \((2, 1)\), by computing both slopes \( dy/dx \) and multiplying to get \(-1\).

34. The derivative of \( f(x, y) = c \) is \( f_x + f_y(dy/dx) = 0 \). Show that the slope of this level curve is \( dy/dx = -M/N \). It is perpendicular to streamlines because \( (-M/N)(N/M) = 1 \).

35. The \( x \) and \( y \) derivatives of \( f(r) \) are \( \partial f/\partial x = \mathbf{a} \) and \( \partial f/\partial y = \mathbf{b} \) by the chain rule. (Test \( f = r^2 \).) The equipotentials are _______.

36. \( \mathbf{F} = (ax + by)\mathbf{i} + (bx + cy)\mathbf{j} \) is a gradient field. Find the potential \( f \) and describe the equipotentials.

37. True or false:

1. The constant field \( i + 2k \) is a gradient field.
2. For non-gradient fields, equipotentials meet streamlines at non-right angles.
3. In three dimensions the equipotentials are surfaces instead of curves.
4. \( \mathbf{F} = x^2\mathbf{i} + y^2\mathbf{j} + z^2\mathbf{k} \) points outward from \((0, 0, 0)\)—a radial field.

38. Create and draw \( f \) and \( \mathbf{F} \) and your own equipotentials and streamlines.
How can different vector fields have the same streamlines? Can they have the same equipotentials? Can they have the same $f$?

**15.2 Line Integrals**

A line integral is an integral along a curve. It can equal an area, but that is a special case and not typical. Instead of area, here are two important line integrals in physics and engineering:

$$\text{Work along a curve } = \int_C F \cdot T \, ds \quad \text{Flow across a curve } = \int_C F \cdot n \, ds.$$  

In the first integral, $F$ is a force field. In the second integral, $F$ is a flow field. Work is done in the direction of movement, so we integrate $F \cdot T$. Flow is measured through the curve $C$, so we integrate $F \cdot n$. Here $T$ is the unit tangent vector, and $F \cdot T$ is the force component along the curve. Similarly $n$ is the unit normal vector, at right angles with $T$. Then $F \cdot n$ is the component of flow perpendicular to the curve.

We will write those integrals in several forms. They may never be as comfortable as $\int y(x) \, dx$, but eventually we get them under control. I mention these applications early, so you can see where we are going. This section concentrates on work, and flow comes later. (It is also called flux—the Latin word for flow.) You recognize $ds$ as the step along the curve, corresponding to $dx$ on the $x$ axis. Where $\int dx$ gives the length of an interval (it equals $b-a$), $\int ds$ is the length of the curve.

**EXAMPLE 1** Flight from Atlanta to Los Angeles on a straight line and a semicircle.

According to Delta Airlines, the distance straight west is 2000 miles. Atlanta is at $(1000, 0)$ and Los Angeles is at $(-1000, 0)$, with the origin halfway between. The semicircle route $C$ has radius 1000. *This is not a great circle route*. It is more of a "flat circle," which goes north past Chicago. No plane could fly it (it probably goes into space).

The equation for the semicircle is $x^2 + y^2 = 1000^2$. Parametrically this path is $x = 1000 \cos t$, $y = 1000 \sin t$. For a line integral the parameter is better. The plane leaves Atlanta at $t = 0$ and reaches L.A. at $t = \pi$, more than three hours later. On the straight 2000-mile path, Delta could almost do it. Around the semicircle $C$, the distance is 1000$a$ miles and the speed has to be 1000 miles per hour. Remember that speed is distance $ds$ divided by time $dt$:

$$\frac{ds}{dt} = \sqrt{(dx/dt)^2 + (dy/dt)^2} = 1000 \sqrt{(-\sin t)^2 + (\cos t)^2} = 1000. \quad (1)$$

The tangent vector to $C$ is proportional to $(dx/dt, dy/dt) = (-1000 \sin t, 1000 \cos t)$. But $T$ is a unit vector, so we divide by 1000—which is the speed.

Suppose the wind blows due east with force $F = Mi$. The components are $M$ and zero. For $M =$ constant, compute the dot product $F \cdot T$ and the work $-2000 M$:

$$\int_C F \cdot T \, ds = \int_{t=0}^{t=\pi} (-M \sin t) \left( \frac{ds}{dt} \right) dt = \int_{t=0}^{t=\pi} -1000 M \sin t \, dt = -2000 M.$$
Work is force times distance moved. It is negative, because the wind acts against the movement. You may point out that the work could have been found more simply—go 2000 miles and multiply by \(-M\). I would object that this straight route is a different path. But you claim that the path doesn't matter—the work of the wind is \(-2000M\) on every path. I concede that this time you are right (but not always).

Most line integrals depend on the path. Those that don’t are crucially important. For a gradient field, we only need to know the starting point \(P\) and the finish \(Q\).

When \(F\) is the gradient of a potential function \(f(x, y)\), the work \(\int_C F \cdot T \, ds\) depends only on the endpoints \(P\) and \(Q\). The work is the change in \(f\):

\[
\text{if } F = \frac{\partial f}{\partial x} i + \frac{\partial f}{\partial y} j \text{ then } \int_C F \cdot T \, ds = f(Q) - f(P). \tag{2}
\]

When \(F = Mi\), its components \(M\) and zero are the partial derivatives of \(f = Mx\). To compute the line integral, just evaluate \(f\) at the endpoints. Atlanta has \(x = 1000\), Los Angeles has \(x = -1000\), and the potential function \(f = Mx\) is like an antiderivative:

\[
\text{work} = f(Q) - f(P) = M(-1000) - M(1000) = -2000M. \tag{3}
\]

May I give a rough explanation of the work integral \(\int F \cdot T \, ds\)? It becomes clearer when the small movement \(T \, ds\) is written as \(dx \, i + dy \, j\). The work is the dot product with \(F\):

\[
F \cdot T \, ds = \left(\frac{\partial f}{\partial x} \, i + \frac{\partial f}{\partial y} \, j\right) \cdot (dx \, i + dy \, j) = \frac{\partial f}{\partial x} \, dx + \frac{\partial f}{\partial y} \, dy = df. \tag{4}
\]

The infinitesimal work is \(df\). The total work is \(\int df = f(Q) - f(P)\). This is the Fundamental Theorem for a line integral. Only one warning: When \(F\) is not the gradient of any \(f\) (Example 2), the Theorem does not apply.

**EXAMPLE 2** Fly these paths against the non-constant force field \(F = Myi\). Compute the work.

There is no force on the straight path where \(y = 0\). Along the \(x\) axis the wind does no work. But the semicircle goes up where \(y = 1000 \sin t\) and the wind is strong:

\[
F \cdot T = (Myi) \cdot (-\sin t \, i + \cos t \, j) = -My \, \sin t = -1000M \sin^2 t
\]

\[
\int_C F \cdot T \, ds = \int_0^\pi (-1000M \sin^2 t) \frac{ds}{dt} \, dt = \int_0^\pi -10^5M \sin^2 t \, dt = -\frac{\pi}{2} \cdot 10^5M.
\]

This work is enormous (and unrealistic). But the calculations make an important point—everything is converted to the parameter \(t\). The second point is that \(F = Myi\) is not a gradient field. First reason: The work was zero on the straight path and...
nonzero on the semicircle. **Second reason:** No function has $\frac{\partial f}{\partial x} = My$ and $\frac{\partial f}{\partial y} = 0$. (The first makes $f$ depend on $y$ and the second forbids it. This $F$ is called a shear force.) Without a potential we cannot substitute $P$ and $Q$—and the work depends on the path.

**THE DEFINITION OF LINE INTEGRALS**

We go back to the start, to define $\int F \cdot T \, ds$. We can think of $F \cdot T$ as a function $g(x, y)$ along the path, and define its integral as a limit of sums:

$$\int_c g(x, y) \, ds = \lim_{\Delta s \to 0} \sum_{i=1}^N g(x_i, y_i) \Delta s_i$$

The points $(x_i, y_i)$ lie on the curve $C$. The last point $Q$ is $(x_N, y_N)$; the first point $P$ is $(x_0, y_0)$. The step $\Delta s_i$ is the distance to $(x_i, y_i)$ from the previous point. As the steps get small ($\Delta s \to 0$) the straight pieces follow the curve. Exactly as in Section 8.2, the special case $g = 1$ gives the arc length. As long as $g(x, y)$ is piecewise continuous (jumps allowed) and the path is piecewise smooth (corners allowed), the limit exists and defines the line integral.

When $g$ is the density of a wire, the line integral is the total mass. When $g$ is $F \cdot T$, the integral is the work. But nobody does the calculation by formula (5). We now introduce a parameter $t$—which could be the time, or the arc length $s$, or the distance $x$ along the base.

**The differential $ds$ becomes** $(ds/\,dt)\,dt$. Everything changes over to $t$:

$$\int g(x, y) \, ds = \int_{t=a}^{t=b} g(x(t), y(t)) \sqrt{(dx/\,dt)^2 + (dy/\,dt)^2} \, dt.$$  

The curve starts when $t = a$, runs through the points $(x(t), y(t))$, and ends when $t = b$. The square root in the integral is the speed $ds/\,dt$. In three dimensions the points on $C$ are $(x(t), y(t), z(t))$ and $(dz/\,dt)^2$ is in the square root.

**EXAMPLE 3**  The points on a coil spring are $(x, y, z) = (\cos t, \sin t, t)$. Find the mass of two complete turns (from $t = 0$ to $t = 4\pi$) if the density is $\rho = 4$.

Solution  The key is $(dx/\,dt)^2 + (dy/\,dt)^2 + (dz/\,dt)^2 = \sin^2 t + \cos^2 t + 1 = 2$. Thus $ds/\,dt = \sqrt{2}$. To find the mass, integrate the mass per unit length which is $g = \rho = 4$:

$$\text{mass} = \int_0^{2\pi} \frac{ds}{\,dt} \, dt = \int_0^{2\pi} 4\sqrt{2} \, dt = 16\sqrt{2}\pi.$$  

That is a line integral in three-dimensional space. It shows how to introduce $t$. But it misses the main point of this section, because it contains no vector field $F$. This section is about work, not just mass.

**DIFFERENT FORMS OF THE WORK INTEGRAL**

The work integral $\int F \cdot T \, ds$ can be written in a better way. The force is $F = Mi + Nj$. A small step along the curve is $ds \, i + dy \, j$. Work is force times distance, but it is only the force component along the path that counts. The dot product $F \cdot T \, ds$ finds that component automatically.
15 Vector Calculus

15C The vector to a point on $C$ is $\mathbf{R} = xi + yj$. Then $d\mathbf{R} = \mathbf{T} ds = dx \mathbf{i} + dy \mathbf{j}$:

$$\text{work} = \int_C \mathbf{F} \cdot d\mathbf{R} = \int_C M \, dx + N \, dy.$$  \hfill (7)

Along a space curve the work is $\int \mathbf{F} \cdot \mathbf{T} \, ds = \int \mathbf{F} \cdot d\mathbf{R} = \int M \, dx + N \, dy + P \, dz$.

The product $M \, dx$ is (force in $x$ direction)(movement in $x$ direction). This is zero if either factor is zero. When the only force is gravity, pushing a piano takes no work. It is friction that hurts. Carrying the piano up the stairs brings in $P dz$, and the total work is the piano weight $P$ times the change in $z$.

To connect the new $\int \mathbf{F} \cdot d\mathbf{R}$ with the old $\int \mathbf{F} \cdot \mathbf{T} \, ds$, remember the tangent vector $\mathbf{T}$. It is $d\mathbf{R}/ds$. Therefore $T \, ds$ is $d\mathbf{R}$. The best for computations is $d\mathbf{R}$, because the unit vector $\mathbf{T}$ has a division by $ds/dt = \sqrt{(dx/dt)^2 + (dy/dt)^2}$. Later we multiply by this square root, in converting $ds$ to $(ds/dt) \, dt$. It makes no sense to compute the square root, divide by it, and then multiply by it. That is avoided in the improved form $\int M \, dx + N \, dy$.

**EXAMPLE 4** Vector field $\mathbf{F} = -yi + xj$, path from $(1, 0)$ to $(0, 1)$: Find the work.

**Note 1** This $\mathbf{F}$ is the spin field $\mathbf{S}$. It goes around the origin, while $\mathbf{R} = xi + yj$ goes outward. Their dot product is $\mathbf{F} \cdot \mathbf{R} = -yx + xy = 0$. This does not mean that $\mathbf{F} \cdot d\mathbf{R} = 0$. The force is perpendicular to $\mathbf{R}$, but not to the change in $\mathbf{R}$. The work to move from $(1, 0)$ to $(0, 1)$, $x$ axis to $y$ axis, is not zero.

**Note 2** We have not described the path $C$. That must be done. The spin field is not a gradient field, and the work along a straight line does not equal the work on a quarter-circle:

$$\text{straight line } x = 1 - t, \ y = t \quad \text{quarter-circle } x = \cos t, \ y = \sin t.$$

**Calculation of work** Change $\mathbf{F} \cdot d\mathbf{R} = M \, dx + N \, dy$ to the parameter $t$:

- **Straight line**: $\int -y \, dx + x \, dy = \int_0^1 -t(-dt) + (1-t)dt = 1$

- **Quarter-circle**: $\int -y \, dx + x \, dy = \int_0^{\pi/2} -\sin t(-\sin t \, dt) + \cos t(\cos t \, dt) = \frac{\pi}{2}$.

**General method** The path is given by $x(t)$ and $y(t)$. Substitute those into $M(x, y)$ and $N(x, y)$—then $\mathbf{F}$ is a function of $t$. Also find $dx/dt$ and $dy/dt$. Integrate $M \, dx/dt + N \, dy/dt$ from the starting time $t$ to the finish.

![Fig. 15.6](image-url) Three paths for $\int \mathbf{F} \cdot d\mathbf{R} = \int -y \, dx + x \, dy = 1, \pi/2, 0$. 


15.2 Line Integrals

For practice, take the path down the x axis to the origin \((x = 1 - t, y = 0)\). Then go up the y axis \((x = 0, y = t - 1)\). The starting time at \((1, 0)\) is \(t = 0\). The turning time at the origin is \(t = 1\). The finishing time at \((0, 1)\) is \(t = 2\). The integral has two parts because this new path has two parts:

Bent path: \(\int (-y)dx + xdy = 0 + 0 \quad (y = 0 \text{ on one part, then } x = 0)\).

**Note 3** The answer depended on the path, for this spin field \(F = S\). The answer did not depend on the choice of parameter. If we follow the same path at a different speed, the work is the same. We can choose another parameter \(t\), since \((ds/dt)dt\) and \((ds/dt)dt\) both equal \(ds\). Traveling twice as fast on the straight path \((x = 1 - 2t, y = 2t)\) we finish at \(t = \frac{1}{2}\) instead of \(t = 1\). The work is still 1:

\[
\int (-y)dx + xdy = \int_0^{1/2} 2\tau (2\tau) + 1 (2\tau) 2\tau = \int_0^{1/2} 4\tau = 1.
\]

**CONSERVATION OF TOTAL ENERGY (KINETIC + POTENTIAL)**

When a force field does work on a mass \(m\), it normally gives that mass a new velocity. Newton's Law is \(F = ma = m\frac{dv}{dt}\). (It is a vector law. Why write out three components?) The work \(\int F \cdot dR\) is

\[
\int (m\frac{dv}{dt}) \cdot (v \cdot dt) = \frac{1}{2}mv^2 - \frac{1}{2}mv_0^2 = \frac{1}{2}m|v(Q)|^2 - \frac{1}{2}m|v(P)|^2.
\]

The work equals the change in the kinetic energy \(\frac{1}{2}mv^2\). But for a gradient field the work is also the change in potential—with a minus sign from physics:

\[
\text{work} = \int F \cdot dR = -\int df = f(P) - f(Q).
\]

Comparing (8) with (9), the combination \(\frac{1}{2}m|v|^2 + f\) is the same at \(P\) and \(Q\). The total energy, kinetic plus potential, is conserved.

**INDEPENDENCE OF PATH: GRADIENT FIELDS**

The work of the spin field \(S\) depends on the path. Example 4 took three paths—straight line, quarter-circle, bent line. The work was 1, \(\pi/2\), and 0, different on each path. This happens for more than 99.99% of all vector fields. It does not happen for the most important fields. Mathematics and physics concentrate on very special fields—forwhich the work depends only on the endpoints. We now explain what happens, when the integral is independent of the path.

Suppose you integrate from \(P\) to \(Q\) on one path, and back to \(P\) on another path. Combined, that is a closed path from \(P\) to \(P\) (Figure 15.7). But a backward integral is the negative of a forward integral, since \(dR\) switches sign. If the integrals from \(P\) to \(Q\) are equal, the integral around the closed path is zero:

\[
\oint_P F \cdot dR = \int_P^Q F \cdot dR + \int_Q^P F \cdot dR = \int_P^Q F \cdot dR - \int_P^Q F \cdot dR = 0.
\]

The circle on the first integral indicates a closed path. Later we will drop the \(P\)'s.

Not all closed path integrals are zero! For most fields \(F\), different paths yield different work. For "conservative" fields, all paths yield the same work. Then zero
work around a closed path conserves energy. The big question is: **How to decide which fields are conservative, without trying all paths?** Here is the crucial information about conservative fields, in a plane region \( R \) with no holes:

<table>
<thead>
<tr>
<th>15D</th>
<th><strong>( \mathbf{F} = M(x, y) \mathbf{i} + N(x, y) \mathbf{j} )</strong> is a conservative field if it has these properties:</th>
</tr>
</thead>
<tbody>
<tr>
<td>A.</td>
<td>The work ( \int_C \mathbf{F} \cdot d\mathbf{R} ) around every closed path is zero.</td>
</tr>
<tr>
<td>B.</td>
<td>The work ( \oint_C \mathbf{F} \cdot d\mathbf{R} ) depends only on ( P ) and ( Q ), not on the path.</td>
</tr>
<tr>
<td>C.</td>
<td>( \mathbf{F} ) is a <strong>gradient field</strong>: ( M = \partial f/\partial x ) and ( N = \partial f/\partial y ) for some potential ( f(x, y) ).</td>
</tr>
<tr>
<td>D.</td>
<td>The components satisfy ( \partial M/\partial y = \partial N/\partial x ).</td>
</tr>
</tbody>
</table>

A field with one of these properties has them all. D is the quick test.

These statements A–D bring everything together for conservative fields (alias gradient fields). A closed path goes one way to \( Q \) and back the other way to \( P \). The work cancels, and statements A and B are equivalent. We now connect them to C. **Note**: Test D says that the "curl" of \( \mathbf{F} \) is zero. That can wait for Green's Theorem in the next section—the full discussion of the curl comes in 15.6.

First, a gradient field \( \mathbf{F} = \nabla f \) is conservative. The work is \( f(Q) - f(P) \), by the fundamental theorem for line integrals. It depends only on the endpoints and not the path. Therefore statement C leads back to B.

Our job is in the other direction, to show that conservative fields \( M \mathbf{i} + N \mathbf{j} \) are gradients. Assume that the work integral depends only on the endpoints. We must construct a potential \( f \), so that \( \mathbf{F} \) is its gradient. In other words, \( \partial f/\partial x \) must be \( M \) and \( \partial f/\partial y \) must be \( N \).

**Fix the point \( P \). Define \( f(Q) \) as the work to reach \( Q \). Then \( \mathbf{F} \) equals \( \nabla f \).**

Check the reasoning. At the starting point \( P, f \) is zero. At every other point \( Q, f \) is the work \( \int M \, dx + N \, dy \) to reach that point. **All paths from \( P \) to \( Q \) give the same** \( f(Q) \), because the field is assumed conservative. After two examples we prove that \( \nabla f \) agrees with \( \mathbf{F} \)—the construction succeeds.

**EXAMPLE 5** Find \( f(x, y) \) when \( \mathbf{F} = M \mathbf{i} + N \mathbf{j} = 2xy \mathbf{i} + x^2 \mathbf{j} \). We want \( \partial f/\partial x = 2xy \) and \( \partial f/\partial y = x^2 \).

**Solution 1** Choose \( P = (0, 0) \). Integrate \( M \, dx + N \, dy \) along to \((x, 0)\) and up to \((x, y)\):

\[
\int_{(0,0)}^{(x,0)} 2xy \, dx = 0 \text{ (since } y = 0) \quad \int_{(x,0)}^{(x,y)} x^2 \, dy = x^2 y \text{ (which is } f). 
\]

Certainly \( f = x^2 y \) meets the requirements: \( f_x = 2xy \) and \( f_y = x^2 \). Thus \( \mathbf{F} = \nabla f \). Note that \( dy = 0 \) in the first integral (on the \( x \) axis). Then \( dx = 0 \) in the second integral (\( x \) is fixed). The integrals add to \( f = x^2 y \).
15.2 Line Integrals

Solution 2 Integrate 2 x y d x + x2dy on the straight line (xt, yt) from t = 0 to t = 1:

Iol

+

2(xt)(yt)(xdt) ( ~ t ) ~dt)( =
y

So1

3 x2yt2dt= x2yt3]: = x2y.

Most authors use Solution 1. I use Solution 2. Most students use Solution 3:
Solution 3 Directly solve df/dx = M = 2xy and then fix up dfldy = N = x2:

af/dx = 2xy gives f = x2y (plus any function of y).

In this example x2y already has the correct derivative dfldy = x2. No additional
function of y is necessary. When we integrate with respect to x, the constant of
integration (usually C ) becomes a function C(y).
You will get practice in finding f. This is only possible for conservative fields! I
tested M = 2xy and N = x2 in advance (using D) to be sure that dM/dy = dN/dx.

+ Nj is the spin field -yi + xj.
-y d x + x dy from (0,O) to ( x , 0 ) to ( x , y):

EXAMPLE 6 Look for f ( x , y) when Mi

Attempted solution 1 Integrate

I

(x, 0)

I

(x. Y)

-ydx=O

(0,O)

and

x dy = x y (which seems like f ) .

(x. 0)

Attempted solution 2 Integrate - y d x + x dy on the line (xt, yt) from t = 0 to 1 :

So1

-( y t ) ( xdt) + ( x t ) ( ydt) = 0 (a different f, also wrong).

Aitempted solution 3 Directly solve dfldx = - y and try to fix up af/dy = x :

af/dx = - y gives f = - x y (plus any function C(y)).

The y derivative of this f is - x + dC/dy. That does not agree with the required
dfldy = x. Conclusion: The spin field -yi + xj is not conservative. There is no f.
Test D gives dM/dy = - 1 and dN/dx = + 1.
To finish this section, we move from examples to a proof. The potential f(Q) is
defined as the work to reach Q. We must show that its partial derivatives are M and
N. This seems reasonable from the formula f (Q)= M d x N dy, but we have to
think it through.
Remember statement A, that all paths give the same f(Q). Take a path that goes
from P to the left of Q. It comes in to Q on a line y = constant (so dy = 0). As the
path reaches Q, we are only integrating M dx. The derivative of this integral, at Q, is
df/ax = M. That is the Fundamental Theorem of Calculus.
To show that af/ay = N, take a different path. G o from P to a point below Q. The
path comes up to Q on a vertical line (so d x = 0). Near Q we are only integrating
N dy, so i?f/dy = N.
The requirement that the region must have no holes will be critical for test D.

I

EXAMPLE 7 Find f ( x ,y) =

+

x d x + y dy. Test D is passed: a N / a x = 0 = dM/dy.

j:",: x d x = +x2 is added to j:;:; y dy = fy2.
Solution 2 1; ( x t ) ( xdt) + ( y t ) ( ydt) = 1; ( x 2 + y2)t dt = f ( x 2+ y2).
Solution 1

+

Solution 3 afjax = x gives f = +x2 C(y). Then af/dy = y needs C ( y )= :y2.


15.2 Exercises

Read-through questions

Work is the _c_ of F · dR. Here F is the _b_ and R is the _c_. The _d_ product finds the component of _e_ in the direction of movement dR = dx i + dy j. The straight path (x, y) = (t, 2t) goes from _f_ to _g_ at t = 0 to _g_ at t = 1 with dR = dt i + _h_. The work of F = 3i + j is \( \int F \cdot dR = \int_1^1 dt = 1 \).

Another form of dR is T ds, where T is the _k_ vector to the path and ds = _k_. For the path (t, 2t), the unit vector T is _m_ and ds = _n_ dt. For F = 3i + j, F · T ds is still _p_. This F is the gradient of f = _q_. The change in f = 3x + y from (0, 0) to (1, 2) is _q_.

When F = grad f, the dot product F · dR is (\( \partial f/\partial x \))dx + (\( \partial f/\partial y \))dy. The work integral from P to Q is \( \int_{P}^{Q} f \). In this case the work depends on the path. The path (x, y) = (3 cos t, 3 sin t) is a circle with S · dR = _q_. Formally \( \int f(x, y) ds \) is the limit of the sum _c_.

The four equivalent properties of a conservative field F = M i + N j are A: _D_, B: _E_, C: _F_, and D: _G_. Test D (passed)(not passed) by F = (y + 1)j + xj. The work \( \int F \cdot dR \) around the circle (cos t, sin t) is _h_. The work on the upper semicircle equals the work on _i_. This field is the gradient of f = _j_, so the work to (1, 0) is _k_.

Compute the line integrals in 1–6.

1. \( \int_C ds \) and \( \int_C dy \): x = t, y = 2t, 0 ≤ t ≤ 1.
2. \( \int_C dx \) and \( \int_C y dx \): x = cos t, y = sin t, 0 ≤ t ≤ \( \pi/2 \).
3. \( \int_C xy dx - x dy \): any square path, sides of length 3.
4. \( \int_C ds \) and \( \int_C dx \): any closed circle of radius 3.
5. \( \int_C ds \) and \( \int_C y dx \): any closed circle of radius 3.
6. \( \int_C ds \) and \( \int_C y dx \): any path of length 7.
7. Does \( \int_C xy dx \) = \( 1/2 \)\( \int_C xy^2 \)?
8. Does \( \int_C x dx \) = \( 1/2 \)\( \int_C x^2 \)?
9. Does \( (\int_C dx)^2 = (\int_C dx)^2 + (\int_C dy)^2 ?
10. Does \( (\int_C dx)^2 \) make sense?

In 11–16 find the work in moving from (1, 0) to (0, 1). When F is conservative, construct f. Choose your own path when F is not conservative.

11. F = i + y j
12. F = y i + j
13. F = xy i + y x j
14. F = e^x i + xe^x j
15. F = (x/r)i + (y/r) j
16. F = - y^2 i + x^2 j

17. For which powers n is S/r^n a gradient by test D?
18. For which powers n is R/r^n a gradient by test D?

19. A wire hoop around a vertical circle \( x^2 + y^2 = a^2 \) has density \( \rho = a + z \). Find its mass M = \( \int ds \).
20. A wire of constant density \( \rho \) lies on the semicircle \( x^2 + y^2 = a^2 \), \( y ≥ 0 \). Find its mass M and also its moment \( M_x = \int \rho y ds \). Where is its center of mass \( \bar{x} = M_x/M, \bar{y} = M_y/M \)?
21. If the density around the circle \( x^2 + y^2 = a^2 \) is \( \rho = x^2 \), what is the mass and where is the center of mass?
22. Find \( \int F \cdot dR \) along the space curve x = t, y = t^2, z = t^3, 0 ≤ t ≤ 1.

(a) F = grad (xy + xz)  
(b) F = yi - xj + zk

23. (a) Find the unit tangent vector T and the speed ds/dt along the path R = 2ti + 12j.
(b) For F = 3xi + 4j, find F · T ds using (a) and F · dR directly.
(c) What is the work from (2, 1) to (4, 4)?

24. If M(x, y, z)i + N(x, y, z)j is the gradient of f(x, y, z), show that none of these functions can depend on z.
25. Find all gradient fields of the form M(y)j + N(x)j.
26. Compute the work W(x, y) = \( \int M dx + N dy \) on the straight line path (xt, yt) from t = 0 to t = 1. Test to see if \( \partial W/\partial x = M \) and \( \partial W/\partial y = N \).

(a) M = y^3, N = 3x y^2  
(b) M = x^3, N = 3xy^2  
(c) M = x/y, N = y/x  
(d) M = e^{xy}, N = e^{x+y}  

27. Find a field F whose work around the unit square (y = 0 then x = 1 then y = 1 then x = 0) equals 4.
28. Find a nonconservative F whose work around the unit circle \( x^2 + y^2 = 1 \) is zero.

In 29–34 compute \( \int F \cdot dR \) along the straight line R = ti + j and the parabola R = ti + t^2j, from (0, 0) to (1, 1). When F is a gradient field, use its potential f(x, y).

29. F = i - j  
30. F = x^2 j
31. F = 2xy^2i + 2yx^2 j  
32. F = x^2yi + xy^2 j
33. F = yi - xj  
34. F = (xi + yj)/((x^2 + y^2) + 1)

35. For which numbers a and b is F = axyi + (x^2 + by)j a gradient field?
36. Compute \( -y dx + x dy \) from (1, 0) to (0, 1) on the line x = 1 - t^2, y = t^2 and the quarter-circle x = cos 2t, y = sin 2t. Example 4 found 1 and \( \pi/2 \) with different parameters.
Apply the test $N_x = M_y$ to 37–42. Find $f$ when test $D$ is passed.

37 $F = y^2e^{-x}i - 2ye^{-x}j$

38 $F = y^2e^{xz} - 2ye^{xz}j$

39 $F = \frac{xi + yj}{|xi + yj|}$

40 $F = \frac{\text{grad } xy}{|\text{grad } xy|}$

41 $F = R + S$

42 $F = (ax + by)i + (cx + dy)j$

43 Around the unit circle find $\oint ds$ and $\oint dx$ and $\oint xds$.

44 True or false, with reason:

(a) When $F = yi$ the line integral $\int F \cdot d\mathbf{R}$ along a curve from $P$ to $Q$ equals the usual area under the curve.

(b) That line integral depends only on $P$ and $Q$, not on the curve.

(c) That line integral around the unit circle equals $\pi$.

15.3 Green's Theorem

This section contains the Fundamental Theorem of Calculus, extended to two dimensions. That sounds important and it is. The formula was discovered 150 years after Newton and Leibniz, by an ordinary mortal named George Green. His theorem connects a double integral over a region $R$ to a line integral along its boundary $C$.

The integral of $df/dx$ equals $f(b) - f(a)$. This connects a one-dimensional integral to a zero-dimensional integral. The boundary only contains two points $a$ and $b$. The answer $f(b) - f(a)$ is some kind of a "point integral." It is this absolutely crucial idea—to integrate a derivative from information at the boundary—that Green's Theorem extends into two dimensions.

There are two important integrals around $C$. The work is $\int \mathbf{F} \cdot T \, ds = \int M \, dx + N \, dy$. The flux is $\int \mathbf{F} \cdot n \, ds = \int M \, dy - N \, dx$ (notice the switch). The first is for a force field, the second is for a flow field. The tangent vector $T$ turns 90° clockwise to become the normal vector $n$. Green's Theorem handles both, in two dimensions. In three dimensions they split into the Divergence Theorem (15.5) and Stokes' Theorem (15.6).

Green's Theorem applies to "smooth" functions $M(x, y)$ and $N(x, y)$, with continuous first derivatives in a region slightly bigger than $R$. Then all integrals are well defined. $M$ and $N$ will have a definite and specific meaning in each application—to electricity or magnetism or fluid flow or mechanics. The purpose of a theorem is to capture the central ideas once and for all. We do that now, and the applications follow.

**15E Green's Theorem** Suppose the region $R$ is bounded by the simple closed piecewise smooth curve $C$. Then an integral over $R$ equals a line integral around $C$:

$$\oint_C M \, dx + N \, dy = \iint_R \left( \frac{\partial N}{\partial x} - \frac{\partial M}{\partial y} \right) \, dx \, dy. \tag{1}$$

A curve is "simple" if it doesn't cross itself (figure 8's are excluded). It is "closed" if its endpoint $Q$ is the same as its starting point $P$. This is indicated by the closed circle on the integral sign. The curve is "smooth" if its tangent $T$ changes continuously—the word "piecewise" allows a finite number of corners. Fractals are not allowed, but all reasonable curves are acceptable (later we discuss figure 8's and rings). First comes an understanding of the formula, by testing it on special cases.
**Special case 1:** $M = 0$ and $N = x$. Green's Theorem with $\partial N/\partial x = 1$ becomes

$$\oint_C x \, dy = \iint_R 1 \, dx \, dy$$

(which is the area of $R$). (2)

The integrals look equal, because the inner integral of $dx$ is $x$. Then both integrals have $x \, dy$—but we need to go carefully. The area of a layer of $R$ is $dy$ times the difference in $x$ (the length of the strip). The line integral in Figure 15.8 agrees. It has an upward $dy$ times $x$ (at the right) plus a downward $-dy$ times $x$ (at the left). The integrals add up the strips, to give the total area.

**Special case 2:** $M = y$ and $N = 0$ and $\oint_C y \, dx = \iint_R (-1) \, dx \, dy = -(\text{area of } R)$.

Now Green's formula has a minus sign, because the line integral is *counter-clockwise*. The top of each slice has $dx < 0$ (going left) and the bottom has $dx > 0$ (going right). Then $y \, dx$ at the top and bottom combine to give *minus* the area of the slice in Figure 15.8b.

**Special case 3:** $\oint_C 1 \, dx = 0$. The $dx$'s to the right cancel the $dx$'s to the left (the curve is closed). With $M = 1$ and $N = 0$, Green's Theorem is $0 = 0$.

**Most important case:** $Mi + Nj$ is a *gradient field*. It has a potential function $f(x, y)$. Green's Theorem is $0 = 0$, because $\partial M/\partial y = \partial N/\partial x$. This is test D:

$$\frac{\partial M}{\partial y} = \frac{\partial f}{\partial y} \left( \frac{\partial f}{\partial x} \right)$$

is the same as

$$\frac{\partial N}{\partial x} = \frac{\partial f}{\partial x} \left( \frac{\partial f}{\partial y} \right)$$

(3)

The cross derivatives always satisfy $f_{xy} = f_{yx}$. That is why gradient fields pass test D.

When the double integral is zero, the line integral is also zero: $\oint_C M \, dx + N \, dy = 0$. The work is zero. *The field is conservative!* This last step in $A \Rightarrow B \Rightarrow C \Rightarrow D \Rightarrow A$ will be complete when Green's Theorem is proved.

Conservative examples are $\oint_C x \, dx = 0$ and $\oint_C y \, dy = 0$. Area is not involved.

**Remark** The special cases $\oint_C x \, dy$ and $-\oint_C y \, dx$ led to the area of $R$. As long as $1 = \partial N/\partial x - \partial M/\partial y$, the double integral becomes $\iint 1 \, dx \, dy$. This gives a way to compute area by a line integral.

**The area of $R$ is**

$$\oint_C x \, dy = -\oint_C y \, dx = \frac{1}{2} \oint_C (x \, dy - y \, dx).$$

(4)

**EXAMPLE 1** The area of the triangle in Figure 15.9 is 2. Check Green's Theorem.

The last area formula in (4) uses $\frac{1}{2}S$, half the spin field. $N = \frac{1}{2}x$ and $M = -\frac{1}{2}y$ yield $N_x - M_y = \frac{1}{2} + \frac{1}{2} = 1$. On one side of Green's Theorem is $\iint 1 \, dx \, dy = \text{area of triangle}$. On the other side, the line integral has three pieces.
15.3  Green's Theorem

Two pieces are zero: \( x \, dy - y \, dx = 0 \) on the sides where \( x = 0 \) and \( y = 0 \). The sloping side \( x = 2 - y \) has \( dx = -dy \). The line integral agrees with the area, confirming Green's Theorem:

\[
\frac{1}{2} \int_{C} x \, dy - y \, dx = \frac{1}{2} \int_{y=0}^{2} (2 - y) \, dy + y \, dy = \frac{1}{2} \int_{0}^{2} 2 \, dy = 2.
\]

**EXAMPLE 2**  The area of an ellipse is \( \pi ab \) when the semiaxes have lengths \( a \) and \( b \).

This is a classical example, which all authors like. The points on the ellipse are \( x = a \cos t \), \( y = b \sin t \), as \( t \) goes from 0 to \( 2\pi \). (The ellipse has \( (x/a)^2 + (y/b)^2 = 1 \).) By computing the boundary integral, we discover the area inside. Note that the differential \( x \, dy - y \, dx \) is just \( ab \, dt \):

\[
(a \cos t)(b \cos t \, dt) - (b \sin t)(-a \sin t \, dt) = ab(\cos^2 t + \sin^2 t) \, dt = ab \, dt.
\]

The line integral is \( \frac{1}{2} \int_{0}^{2\pi} ab \, dt = \pi ab \). This area \( \pi ab \) is \( \pi r^2 \), for a circle with \( a = b = r \).

**Proof of Green's Theorem:** In our special cases, the two sides of the formula were equal. We now show that they are always equal. The proof uses the Fundamental Theorem to integrate \( \partial N / \partial x \) \( dx \) and \( \partial M / \partial y \) \( dy \). Frankly speaking, this one-dimensional theorem is all we have to work with—since we don't know \( M \) and \( N \). The proof is a step up in mathematics, to work with symbols \( M \) and \( N \) instead of specific functions. The integral in (6) below has no numbers. The idea is to deal with \( M \) and \( N \) in two separate parts, which added together give Green's Theorem:

\[
\oint_{C} M \, dx = \int_{R} - \frac{\partial M}{\partial y} \, dx \, dy \quad \text{and separately} \quad \oint_{C} N \, dy = \int_{R} \frac{\partial N}{\partial x} \, dx \, dy. \quad (5)
\]

Start with a "very simple" region (Figure 15.10a). Its top is given by \( y = f(x) \) and its bottom by \( y = g(x) \). In the double integral, integrate \( - \partial M / \partial y \) first with respect to \( y \). The inner integral is

\[
\int_{g(x)}^{f(x)} - \frac{\partial M}{\partial y} \, dy = -M(x, y) \bigg|_{g(x)}^{f(x)} = -M(x, f(x)) + M(x, g(x)). \quad (6)
\]

The Fundamental Theorem (in the \( y \) variable) gives this answer that depends on \( x \). If we knew \( M \) and \( f \) and \( g \), we could do the outer integral—from \( x = a \) to \( x = b \). But we have to leave it and go to the other side of Green's Theorem—the line integral:

\[
\oint M \, dx = \int_{\text{top}} M(x, y) \, dx + \int_{\text{bottom}} M(x, y) \, dx = \int_{a}^{b} M(x, f(x)) \, dx + \int_{a}^{b} M(x, g(x)) \, dx. \quad (7)
\]
Compare (7) with (6). The integral of $M(x, g(x))$ is the same for both. The integral of $M(x, f(x))$ has a minus sign from (6). In (7) it has a plus sign but the integral is from $b$ to $a$. So life is good.

The part for $N$ uses the same idea. Now the $x$ integral comes first, because $(\partial N/\partial x)dx$ is practically asking to be integrated—from $x = G(y)$ at the left to $x = F(y)$ at the right. We reach $N(F(y), y) - N(G(y), y)$. Then the $y$ integral matches $\frac{1}{2} N dy$ and completes (5). Adding the two parts of (5) proves Green's Theorem.

Finally we discuss the shape of $R$. The broken ring in Figure 15.10 is not "very simple," because horizontal lines go in and out and in and out. Vertical lines do the same. The $x$ and $y$ strips break into pieces. Our reasoning assumed no break between $y = f(x)$ at the top and $y = g(x)$ at the bottom.

There is a nice idea that saves Green's Theorem. Separate the broken ring into three very simple regions $R_1$, $R_2$, $R_3$. The three double integrals equal the three line integrals around the $R$'s. Now add these separate results, to produce the double integral over all of $R$. When we add the line integrals, the crosscuts $CC$ are covered twice and they cancel. The cut between $R_1$ and $R_2$ is covered upward (around $R_1$) and downward (around $R_2$). That leaves the integral around the boundary equal to the double integral inside—which is Green's Theorem.

When $R$ is a complete ring, including the piece $R_4$, the theorem is still true. The integral around the outside is still counterclockwise. But the integral is clockwise around the inner circle. Keep the region $R$ to your left as you go around $C$. The complete ring is "doubly" connected, not "simply" connected. Green's Theorem allows any finite number of regions $R_i$ and crosscuts $CC$ and holes.

**EXAMPLE 3** The area under a curve is $\int_a^b y \, dx$, as we always believed.

In computing area we never noticed the whole boundary. The true area is a line integral $-\frac{1}{2} y \, dx$ around the closed curve in Figure 15.11a. But $y = 0$ on the $x$ axis. Also $dx = 0$ on the vertical lines (up and down at $b$ and $a$). Those parts contribute zero to the integral of $y \, dx$. The only nonzero part is back along the curve—which is the area $-\int_a^b y \, dx$ or $\int_a^b y \, dx$ that we know well.

What about signs, when the curve dips below the $x$ axis? That area has been counted as negative since Chapter 1. I saved the proof for Chapter 15. The reason lies in the arrows on $C$.

The line integral around that part goes the other way. The arrows are clockwise, the region is on the right, and the area counts as negative. With the correct rules, a figure 8 is allowed after all.
CONSERVATIVE FIELDS

We never leave gradients alone! They give conservative fields—the work around a closed path is \( f(P) - f(P) = 0 \). But a potential function \( f(x, y) \) is only available when test \( D \) is passed: If \( \frac{\partial f}{\partial x} = M \) and \( \frac{\partial f}{\partial y} = N \) then \( \partial M/\partial y = \partial N/\partial x \). The reason is that \( f_{xy} = f_{yx} \).

Some applications prefer the language of “differentials.” Instead of looking for \( f(x, y) \), we look for \( df \):

**DEFINITION** The expression \( M(x, y)dx + N(x, y)dy \) is a *differential form*. When it agrees with the differential \( df = (\partial f/\partial x)dx + (\partial f/\partial y)dy \) of some function, the form is called *exact*. The test for an exact differential is \( D: \partial N/\partial x = \partial M/\partial y \).

Nothing is new but the language. Is \( y \, dx \) an exact differential? No, because \( M_y = 1 \) and \( N_x = 0 \). Is \( y \, dx + x \, dy \) an exact differential? Yes, it is the differential of \( f = xy \). That is the product rule! Now comes an important example, to show why \( R \) should be *simply connected* (a region with no holes).

**EXAMPLE 4** The spin field \( S/r^2 = (-yi + xj)/(x^2 + y^2) \) almost passes test \( D \).

\[
N_x = \frac{\partial}{\partial x} \left( \frac{x}{x^2 + y^2} \right) = \frac{x^2 + y^2 - x(2x)}{(x^2 + y^2)^2} = M_y = \frac{\partial}{\partial y} \left( \frac{-y}{x^2 + y^2} \right) = \frac{-(x^2 + y^2) + y(2y)}{(x^2 + y^2)^2}.
\]

Both numerators are \( y^2 - x^2 \). Test \( D \) looks good. To find \( f \), integrate \( M = \partial f/\partial x \):

\[
f(x, y) = \int -y \, dx/(x^2 + y^2) = \tan^{-1}(y/x) + C(y).
\]

The extra part \( C(y) \) can be zero—the \( y \) derivative of \( \tan^{-1}(y/x) \) gives \( N \) with no help from \( C(y) \). The *potential \( f \) is the angle \( \theta \) in the usual \( x, y, r \) right triangle*.

Test \( D \) is passed and \( F \) is grad \( \theta \). What am I worried about? It is only this, that *Green’s Theorem on a circle seems to give \( 2\pi = 0 \).* The double integral is \( \oint (N_x - M_y)dx \, dy \). According to (8) this is the integral of zero. But the line integral is \( 2\pi \):

\[
\oint F \cdot dR = \oint (-y \, dx + x \, dy)/(x^2 + y^2) = 2 \text{(area of circle)}/a^2 = 2\pi a^2/2a^2 = 2\pi.
\]

With \( x = a \cos t \) and \( y = a \sin t \) we would find the same answer. The *work is \( 2\pi \) (not zero!) when the path goes around the origin.*

We have a paradox. If Green’s Theorem is wrong, calculus is in deep trouble. Some requirement must be violated to reach \( 2\pi = 0 \). Looking at \( S/r^2 \), the problem is at the origin. The field is not defined when \( r = 0 \) (it blows up). The derivatives in (8) are not continuous. Test \( D \) does not apply at the origin, and was not passed. *We could remove \( (0, 0) \), but then the region where test \( D \) is passed would have a hole.*
It is amazing how one point can change everything. When the path circles the origin, the line integral is not zero. The potential function \( f = \theta \) increases by \( 2\pi \). That agrees with \( \int F \cdot dR = 2\pi \) from (9). It disagrees with \( \iint 0 \, dx \, dy \). The \( 2\pi \) is right, the zero is wrong. \( N_x - M_y \) must be a "delta function of strength \( 2\pi \)."

The double integral is \( 2\pi \) from an infinite spike over the origin—even though \( N_x = M_y \) everywhere else. In fluid flow the delta function is a "vortex."

**FLOW ACROSS A CURVE: GREEN'S THEOREM TURNED BY 90°**

A flow field is easier to visualize than a force field, because something is really there and it moves. Instead of gravity in empty space, water has velocity \( M(x, y)i + N(x, y)j \). At the boundary \( C \) it can flow in or out. The new form of Green's Theorem is a fundamental "balance equation" of applied mathematics:

*Flow through \( C \) (out minus in) = replacement in \( R \) (source minus sink).*

The flow is *steady*. Whatever goes out through \( C \) must be replaced in \( R \). When there are no sources or sinks (negative sources), the total flow through \( C \) must be zero. This balance law is Green's Theorem in its "normal form" (for \( n \)) instead of its "tangential form" (for \( T \)):

\[
\oint C F \cdot n \, ds = \iint_R \left( \frac{\partial M}{\partial x} + \frac{\partial N}{\partial y} \right) \, dx \, dy.
\]  (10)

Figure 15.12 shows the 90° turn. \( T \) becomes \( n \) and "circulation" along \( C \) becomes flux through \( C \). In the original form of Green's Theorem, change \( N \) and \( M \) to \( M \) and \( -N \) to obtain the flux form:

\[
\int M \, dx + N \, dy \rightarrow \int -N \, dx + M \, dy \quad \iint (N_x - M_y) \, dx \, dy \rightarrow \iint (M_x + N_y) \, dx \, dy.
\]  (11)

Playing with letters has proved a new theorem! The two left sides in (11) are equal, so the right sides are equal—which is Green's Theorem (10) for the flux. The components \( M \) and \( N \) can be chosen freely and named freely.

The change takes \( Mi + Nj \) into its perpendicular field \( -Ni + Mj \). The field is turned at every point (we are not just turning the plane by 90°). The spin field \( S = -yi + xj \) changes to the position field \( R = xi + yj \). The position field \( R \) changes to \( -S \). Streamlines of one field are equipotentials of the other field. The new form (10) of Green's

![Diagram](image.png)

**Fig. 15.12** The perpendicular component \( F \cdot n \) flows through \( C \). Note \( n \, ds = dy \, i - dx \, j \).
15.3 Green's Theorem

Theorem is just as important as the old one—in fact I like it better. It is easier to visualize flow across a curve than circulation along it.

The change of letters was just for the proof. From now on $F = Mi + Nj$.

**EXAMPLE 5** Compute both sides of the new form (10) for $F = 2xi + 3yj$. The region $R$ is a rectangle with sides $a$ and $b$.

**Solution** This field has $\frac{\partial M}{\partial x} + \frac{\partial N}{\partial y} = 2 + 3$. The integral over $R$ is $\iint_R 5 \, dx \, dy = 5ab$. The line integral has four parts, because $R$ has four sides. Between the left and right sides, $M = 2x$ increases by $2a$. Down the left and up the right, $\int M \, dy = 2ab$ (those sides have length $b$). Similarly $N = 3y$ changes by $3b$ between the bottom and top. Those sides have length $a$, so they contribute $3ab$ to a total line integral of $5ab$.

**Important:** The "divergence" of a flow field is $\frac{\partial M}{\partial x} + \frac{\partial N}{\partial y}$. The example has divergence $= 5$. To maintain this flow we must replace 5 units continually—not just at the origin but everywhere. (A one-point source is in example 7.) The divergence is the source strength, because it equals the outflow. To understand Green's Theorem for any vector field $Mi + Nj$, look at a tiny rectangle (sides $dx$ and $dy$):

Flow out the right side minus flow in the left side = (change in $M$) times $dy$
Flow out the top minus flow in the bottom = (change in $N$) times $dx$
Total flow out of rectangle: $dM \, dy + dN \, dx = (\frac{\partial M}{\partial x} + \frac{\partial N}{\partial y})dx \, dy$.

The divergence times the area $dx \, dy$ equals the total flow out. Section 15.5 gives more detail with more care in three dimensions. The divergence is $M_x + N_y + P_z$.

**Fig. 15.13** $M_x + N_y = 2 + 3 = 5$ yields flux $= 5(area) = 5ab$. The flux is $dM \, dy + dN \, dx = (M_x + N_y) \, dx \, dy$. The spin field has no flux.

**EXAMPLE 6** Find the flux through a closed curve $C$ of the spin field $S = -yi + xj$.

**Solution** The field has $M = -y$ and $N = x$ and $M_x + N_y = 0$. The double integral is zero. Therefore the total flow (out minus in) is also zero—through any closed curve. Figure 15.13 shows flow entering and leaving a square. No fluid is added or removed. There is no rain and no evaporation. When the divergence $M_x + N_y$ is zero, there is no source or sink.

**FLOW FIELDS WITHOUT SOURCES**

This is really quite important. Remember that conservative fields do no work around $C$, they have a potential $f$, and they have "zero curl." Now turn those statements through $90^\circ$, to find their twins. Source-free fields have no flux through $C$, they have stream functions $g$, and they have "zero divergence." The new statements E-F-G-H describe fields without sources.
The field \( F = M(x, y)i + N(x, y)j \) is source-free if it has these properties:

E The total flux \( \oint F \cdot n \, ds \) through every closed curve is zero.

F Across all curves from \( P \) to \( Q \), the flux \( \int_F F \cdot n \, ds \) is the same.

G There is a stream function \( g(x, y) \), for which \( M = \frac{\partial g}{\partial y} \) and \( N = -\frac{\partial g}{\partial x} \).

H The components satisfy \( \frac{\partial M}{\partial x} + \frac{\partial N}{\partial y} = 0 \) (the divergence is zero).

A field with one of these properties has them all. H is the quick test.

The spin field \( -yi + xj \) passed this test (Example 6 was source-free). The field \( 2xi + 3yj \) does not pass (Example 5 had \( M_x + N_y = 5 \)). Example 7 almost passes.

**EXAMPLE 7** The radial field \( R/r^2 = (xi + yj)/(x^2 + y^2) \) has a point source at \((0, 0)\).

The new test \( H \) is divergence \( = \frac{\partial M}{\partial x} + \frac{\partial N}{\partial y} = 0 \). Those two derivatives are

\[
\frac{\partial}{\partial x} \left( \frac{x}{x^2 + y^2} \right) = \frac{x^2 + y^2 - x(2x)}{(x^2 + y^2)^2} \quad \text{and} \quad \frac{\partial}{\partial y} \left( \frac{y}{x^2 + y^2} \right) = \frac{x^2 + y^2 - y(2y)}{(x^2 + y^2)^2}.
\]

They add to zero. There seems to be no source (if the calculation is correct). The flow through a circle \( x^2 + y^2 = a^2 \) should be zero. But it's not:

\[
\oint M \, dy - N \, dx = \oint (x \, dy - y \, dx)/(x^2 + y^2) = 2(\text{area of circle})/a^2 = 2\pi.
\]

A source is hidden somewhere. Looking at \( R/r^2 \), the problem is at \((0, 0)\). The field is not defined when \( r = 0 \) (it blows up). The derivatives in (12) are not continuous. Test \( H \) does not apply, and was not passed. The divergence \( M_x + N_y \) must be a "delta function" of strength \( 2\pi \). There is a point source sending flow out through all circles.

I hope you see the analogy with Example 4. The field \( S/r^2 \) is curl-free except at \( r = 0 \). The field \( R/r^2 \) is divergence-free except at \( r = 0 \). The mathematics is parallel and the fields are perpendicular. A potential \( f \) and a stream function \( g \) require a region without holes.

**THE BEST FIELDS: CONSERVATIVE AND SOURCE-FREE**

What if \( F \) is conservative and also source-free? Those are outstandingly important fields. The curl is zero and the divergence is zero. Because the field is conservative, it comes from a potential. Because it is source-free, there is a stream function:

\[
M = \frac{\partial f}{\partial x} = \frac{\partial g}{\partial y} \quad \text{and} \quad N = \frac{\partial f}{\partial y} = -\frac{\partial g}{\partial x}.
\]

Those are the Cauchy-Riemann equations, named after a great mathematician of his time and one of the greatest of all time. I can't end without an example.

**EXAMPLE 8** Show that \( y i + x j \) is both conservative and source-free. Find \( f \) and \( g \).

Solution With \( M = y \) and \( N = x \), check first that \( \frac{\partial M}{\partial y} = 1 = \frac{\partial N}{\partial x} \). There must be a potential function. It is \( f = xy \), which achieves \( \frac{\partial f}{\partial x} = y \) and \( \frac{\partial f}{\partial y} = x \). Note that \( f_{xx} + f_{yy} = 0 \).

Check next that \( \frac{\partial M}{\partial x} + \frac{\partial N}{\partial y} = 0 + 0 \). There must be a stream function. It is \( g = \frac{1}{4}(y^2 - x^2) \), which achieves \( \frac{\partial g}{\partial y} = y \) and \( \frac{\partial g}{\partial x} = -x \). Note that \( g_{xx} + g_{yy} = 0 \).
The curves $f =$ constant are the equipotentials. The curves $g =$ constant are the streamlines (Figure 15.4). These are the twin properties—a conservative field with a potential and a source-free field with a stream function. They come together into the fundamental partial differential equation of equilibrium—Laplace’s equation $f_{xx} + f_{yy} = 0$.

15.3 Green’s Theorem

The work integral $\oint_C M \, dx + N \, dy$ equals the double integral $\iint_R (N_x - M_y) \, dx \, dy$ by Green’s Theorem. For $F = 3i + 4j$ the work equals the area of $R$. When $M = \delta f/\delta x$ and $N = \delta f/\delta y$, the double integral is zero because $f$ is a gradient field. The line integral is zero because $g$ is a potential. An example is $f = -xy$. The direction on $C$ is tangent around the outside and along the boundary of a hole. If $R$ is broken into very simple pieces with crosscuts between them, the integrals of $M$ cancel along the crosscuts.

Test $D$ for gradient fields is tangent. A field that passes this test has $\iint_C F \cdot dR = \alpha$. There is a solution to $f = \alpha$ and $g = \beta$. Then $df = M \, dx + N \, dy$ is an $\alpha$ differential. The spin field $S/r^2$ passes test $D$ except at $r = 0$. Its potential $f = 1/r$ increases by $\alpha$ going around the origin. The integral $\iint_R (N_x - M_y) \, dx \, dy$ is not zero but $1/r$.

The flow form of Green’s Theorem is $\int_C M \, dx - N \, dy$. The divergence of $M + Nj$ is $x$. For $F = x$ the double integral is $1/2$. There is a source. For $F = y$ the divergence is $1$. The divergence of $R/r^2$ is zero except at $r = 0$. This field has a $C$ source.

A field with no source has properties $E = -D$, $F = -E$, $G = -F$, $H = 0$. The stream function $g$ satisfies the equations $G$. Then $\partial M/\partial x + \partial N/\partial y = 0$ because $g_x = 1$ and $g_y = 0$. The example $F = y$ has $g = -x$. There is a potential function. The example $F = x_i - y_j$ has $g = -j$ and also $f = -k$. This $f$ satisfies Laplace’s equation $\Delta f = 0$, because the field $F$ is both $M$ and $N$. The functions $f$ and $g$ are connected by the equations $\partial f/\partial x = \partial g/\partial y$ and $\partial f/\partial y = -\partial g/\partial x$.

Compute the line integrals 1–6 and (separately) the double integrals in Green’s Theorem (1). The circle has $x = a \cos t$, $y = a \sin t$. The triangle has sides $x = 0$, $y = 0$, $x + y = 1$.

1 $\int_C x \, dy$ along the circle
2 $\int_C x^2 \, y \, dy$ along the circle

15.3 EXERCISES

Read-through questions

The work integral $\oint_C M \, dx + N \, dy$ equals the double integral $\alpha$ by $\beta$’s Theorem. For $F = 3i + 4j$ the work equals the area of $R$. When $M = \delta f/\delta x$ and $N = \delta f/\delta y$, the double integral is zero because $f$ is a gradient field. The line integral is zero because $g$ is a potential. An example is $f = -xy$. The direction on $C$ is $\alpha$ around the outside and $\beta$ around the boundary of a hole. If $R$ is broken into very simple pieces with crosscuts between them, the integrals of $M$ cancel along the crosscuts.

Test $D$ for gradient fields is $\gamma$. A field that passes this test has $\iint_C F \cdot dR = \alpha$. There is a solution to $f = \alpha$ and $g = \beta$. Then $df = M \, dx + N \, dy$ is an $\alpha$ differential. The spin field $S/r^2$ passes test $D$ except at $r = 0$. Its potential $f = 1/r$ increases by $\alpha$ going around the origin. The integral $\iint_R (N_x - M_y) \, dx \, dy$ is not zero but $\alpha$.

The flow form of Green’s Theorem is $\gamma = \alpha$. The normal vector in $F \cdot n \, ds$ points $\alpha$ and $|n| = \beta$ and $n \, ds$ equals $\gamma$. The divergence of $M + Nj$ is $\gamma$. For $F = x$ the double integral is $\gamma$. There is a source. For $F = y$ the divergence is $\gamma$. The divergence of $R/r^2$ is zero except at $r = 0$. This field has a $\gamma$ source.

A field with no source has properties $E = -D$, $F = -E$, $G = -F$, $H = 0$. The stream function $g$ satisfies the equations $G$. Then $\partial M/\partial x + \partial N/\partial y = 0$ because $g_x = 1$ and $g_y = 0$. The example $F = y$ has $g = -x$. There is a potential function. The example $F = x_i - y_j$ has $g = -j$ and also $f = -k$. This $f$ satisfies Laplace’s equation $\Delta f = 0$, because the field $F$ is both $M$ and $N$. The functions $f$ and $g$ are connected by the equations $\partial f/\partial x = \partial g/\partial y$ and $\partial f/\partial y = -\partial g/\partial x$.

Compute the line integrals 1–6 and (separately) the double integrals in Green’s Theorem (1). The circle has $x = a \cos t$, $y = a \sin t$. The triangle has sides $x = 0$, $y = 0$, $x + y = 1$.

1 $\int_C x \, dy$ along the circle
2 $\int_C x^2 \, y \, dy$ along the circle
Compute both sides of Green's Theorem in the form (10):
(a) \( F = xi + yj, \) \( R \) = upper half of the disk \( x^2 + y^2 \leq 1. \)
(b) \( F = x^2i + y^3j, C \) = square with sides \( y = 0, x = 1, y = 1, x = 0. \)

8 Show that \( \oint_C (x^2y + 2x)dy + x^2y^2dx \) depends only on the area of \( R. \) Does it equal the area?

9 Find the area inside the hypocycloid \( x = \cos^3 t, y = \sin^3 t \) from \( \oint_C xdy - ydx. \)

10 For constants \( b \) and \( c, \) how is \( \oint_C bydx + cx\,dy \) related to the area inside \( C? \) If \( b = 7, \) which \( c \) makes the integral zero?

11 For \( F = \text{grad} \sqrt{x^2 + y^2}, \) show in three ways that \( \oint_C F \cdot d\mathbf{R} = 0 \) around \( x = \cos t, y = \sin t. \)
   (a) \( F \) is a gradient field so ________.
   (b) Compute \( \oint_C F \cdot d\mathbf{R} \) and compare with the area inside your path.
   (c) Compute \( \oint_C (y\,dx - x\,dy) / (x^2 + y^2) \) and compare with \( 2\pi \) in Example 7.

12 Devise a way to find the one-dimensional theorem \( \int_C (df / dx)\,dx = f(b) - f(a) \) as a special case of Green's Theorem when \( R \) is a square.

13 (a) Choose \( x(t) \) and \( y(t) \) so that the path goes from \( (1, 0) \) to \( (1, 0) \) after circling the origin twice.
   (b) Compute \( \oint_C y\,dx \) and compare with the area inside your path.
   (c) Compute \( \oint_C (y\,dx - x\,dy) / (x^2 + y^2) \) and compare with \( 2\pi \) in Example 7.

14 In Example 4 of the previous section, the work \( \int_C \mathbf{S} \cdot d\mathbf{R} \) between \( (1, 0) \) and \( (0, 1) \) was 1 for the straight path and \( \pi/2 \) for the quarter-circle path. Show that the work is always twice the area between the path and the axes.

Compute both sides of \( \oint_C F \cdot \mathbf{n}\,ds = \int_C (M_x + N_y)\,dx\,dy \) in 15-20.

15 \( F = yi + xj \) in the unit circle
16 \( F = xyi \) in the unit square \( 0 \leq x, y \leq 1 \)
17 \( F = \mathbf{R}/r \) in the unit circle
18 \( F = \mathbf{S}/r \) in the unit square
19 \( F = x^2yi \) in the unit triangle (sides \( x = 0, y = 0, x + y = 1 \))
20 \( F = \text{grad} r \) in the top half of the unit circle.

21 Suppose \( \text{div} \, F = 0 \) except at the origin. Then the flux \( \oint_C F \cdot \mathbf{n}\,ds \) is the same through any two circles around the origin, because _________. (What is \( \int_C (M_x + N_y)\,dx\,dy \) between the circles?)

22 Example 7 has \( \text{div} \, F = 0 \) except at the origin. The flux through every circle \( x^2 + y^2 = a^2 \) is \( 2\pi. \) The flux through a square around the origin is also \( 2\pi \) because _________. (Compare Problem 21.)

23 Evaluate \( \oint_C (x, y)dx + b(x, y)\,dy \) by both forms of Green's Theorem. The choice \( M = a, N = b \) in the work form gives the double integral _________. The choice \( M = b, N = -a \) in the flux form gives the double integral _________. There was only one Green.

24 Evaluate \( \oint_C \cos^3 y\,dy - \sin^3 x\,dx \) by Green's Theorem.

25 The field \( R/r^3 \) in Example 7 has zero divergence except at \( r = 0. \) Solve \( \partial g / \partial y = y(x^2 + y^2) \) to find an attempted stream function \( g. \) Does \( g \) have trouble at the origin?

26 Show that \( S/r^2 \) has zero divergence (except at \( r = 0). \) Find a stream function by solving \( \partial g / \partial y = y/(x^2 + y^2). \) Does \( g \) have trouble at the origin?

27 Which differentials are exact: \( y\,dx - x\,dy, \) \( x^2\,dx + y^2\,dy, \) \( y^2\,dx + x^2\,dy? \)

28 If \( M_x + N_y = 0 \) then the equations \( \partial g / \partial y = _____ \) and \( \partial g / \partial x = _____ \) yield a stream function. If also \( N_x = M_y, \) show that \( g \) satisfies Laplace's equation.

Compute the divergence of each field in 29-36 and solve \( g_y = M \) and \( g_x = -N \) for a stream function (if possible).

29 \( 2xyi - y^3j \)
30 \( 3xy^3i - y^3j \)
31 \( x^2i + y^2j \)
32 \( y^3i + x^2j \)
33 \( e^x \cos y \,i - e^x \sin y \,j \)
34 \( e^{x+y}(i-j) \)
35 \( 2yji + y^2j/x^2 \)
36 \( xyi - xyj \)
37 Compute \( N_x - M_y \) for each field in 29-36 and find a potential function \( f \) when possible.

38 The potential \( f(Q) \) is the work \( \int_P^Q F \cdot \mathbf{T}\,ds \) to reach \( Q \) from a fixed point \( P \) (Section 15.2). In the same way, the stream function \( g(Q) \) can be constructed from the integral _______. Then \( g(Q) - g(P) \) represents the flux across the path from \( P \) to \( Q. \) Why do all paths give the same answer?

39 The real part of \( (x + iy)^3 = x^3 + 3ix^2y - 3xy^2 - iy^3 \) is \( f = x^3 - 3xy^2. \) Its gradient field is \( F = \text{grad} f = _____ \). The divergence of \( F \) is _____ . Therefore \( f \) satisfies Laplace's equation. The Cauchy-Riemann equations.

40 Since \( \text{div} \, F = 0 \) in Problem 39, we can solve \( \partial g / \partial y = _____ \) and \( \partial g / \partial x = _____ \). The stream function is \( g = _____ . \) It is the imaginary part of the same \( (x + iy)^3. \) Check that \( f \) and \( g \) satisfy the Cauchy-Riemann equations.

41 The real part \( f \) and imaginary part \( g \) of \( (x + iy)^3 \) satisfy the Laplace and Cauchy-Riemann equations for \( n = 1, 2, ..., (\text{They give all the polynomial solutions.) Compute} f \) and \( g \) for \( n = 4. \)

42 When is \( M\,dy - N\,dx \) an exact differential \( dg? \)

43 The potential \( f = e^x \cos y \) satisfies Laplace's equation. There must be a \( g. \) Find the field \( F = \text{grad} f \) and the stream function \( g(x, y). \)
44 Show that the spin field $S$ does work around every simple closed curve.

45 For $F = f(x)\mathbf{j}$ and $R = \text{unit square } 0 \leq x \leq 1, 0 \leq y \leq 1$, integrate both sides of Green's Theorem (1). What formula is required from one-variable calculus?

46 A region $R$ is "simply connected" when every closed curve inside $R$ can be squeezed to a point without leaving $R$. Test these regions:

1. $xy$ plane without $(0, 0)$
2. $xyz$ space without $(0, 0, 0)$
3. sphere $x^2 + y^2 + z^2 = 1$
4. a torus (or doughnut)
5. a sweater
6. a human body
7. the region between two spheres
8. $xyz$ space with circle removed.

15.4 Surface Integrals

The double integral in Green's Theorem is over a flat surface $R$. Now the region moves out of the plane. It becomes a curved surface $S$, part of a sphere or cylinder or cone. When the surface has only one $z$ for each $(x, y)$, it is the graph of a function $z(x, y)$. In other cases $S$ can twist and close up—a sphere has an upper $z$ and a lower $z$. In all cases we want to compute area and flux. This is a necessary step (it is our last step) before moving Green's Theorem to three dimensions.

First a quick review. The basic integrals are $\int dx$ and $\int \int dx
dy$ and $\int \int \int dx
dy
dz$. The one that didn't fit was $\int ds$—the length of a curve. When we go from curves to surfaces, $ds$ becomes $dS$. **Area is $\int \int dS$ and flux is $\int \int F \cdot n
dS$, with double integrals because the surfaces are two-dimensional. The main difficulty is in $dS$.**

**All formulas are summarized in a table at the end of the section.**

There are two ways to deal with $ds$ (along curves). The same methods apply to $dS$ (on surfaces). The first is in $xyz$ coordinates; the second uses parameters. Before this subject gets complicated, I will explain those two methods.

**Method 1 is for the graph of a function: curve $y(x)$ or surface $z(x, y)$.

A small piece of the curve is almost straight. It goes across by $dx$ and up by $dy$:

$$
\text{length } ds = \sqrt{(dx)^2 + (dy)^2} = \sqrt{1 + \left(\frac{dy}{dx}\right)^2} dx.
$$

A small piece of the surface is practically flat. Think of a tiny sloping rectangle. One side goes across by $dx$ and up by $(\partial z/\partial x)dx$. The neighboring side goes along by $dy$ and up by $(\partial z/\partial y)dy$. Computing the area is a linear problem (from Chapter 11), because the flat piece is in a plane.

Two vectors $A$ and $B$ form a parallelogram. **The length of their cross product is the area.** In the present case, the vectors are $A = i + (\partial z/\partial x)k$ and $B = j + (\partial z/\partial y)k$. Then $A_{\text{dx}}$ and $B_{\text{dy}}$ are the sides of the small piece, and we compute $A \times B$:

$$
A \times B = \begin{vmatrix}
i & j & k \\
1 & 0 & \partial z/\partial x \\
0 & 1 & \partial z/\partial y
\end{vmatrix} = -\partial z/\partial x i - \partial z/\partial y j + k.
$$

This is exactly the normal vector $N$ to the tangent plane and the surface, from Chapter 13. Please note: The small flat piece is actually a parallelogram (not always...
a rectangle). Its area $dS$ is much like $ds$, but the length of $\mathbf{N} = \mathbf{A} \times \mathbf{B}$ involves two derivatives:

$$\text{area } dS = |\mathbf{A} dx \times \mathbf{B} dy| = |\mathbf{N}| dx dy = \sqrt{1 + (\partial z/\partial x)^2 + (\partial z/\partial y)^2} dx dy. \quad (3)$$

**EXAMPLE 1** Find the area on the plane $z = x + 2y$ above a base area $A$.

*This is the example to visualize.* The area down in the $xy$ plane is $A$. The area up on the sloping plane is greater than $A$. A roof has more area than the room underneath it. If the roof goes up at a $45^\circ$ angle, the ratio is $\sqrt{2}$. Formula (3) yields the correct ratio for any surface—including our plane $z = x + 2y$.

![Fig. 15.14 Roof area = base area times $|\mathbf{N}|$. Cone and cylinder with parameters $u$ and $v$.](image)

The derivatives are $\partial z/\partial x = 1$ and $\partial z/\partial y = 2$. They are constant (planes are easy). The square root in (3) contains $1 + 1^2 + 2^2 = 6$. Therefore $dS = \sqrt{6} dx dy$. An area in the $xy$ plane is multiplied by $\sqrt{6}$ up in the surface (Figure 15.14a). The vectors $\mathbf{A}$ and $\mathbf{B}$ are no longer needed—their work was done when we reached formula (3)—but here they are:

$$\mathbf{A} = \mathbf{i} + (\partial z/\partial x)\mathbf{k} = \mathbf{i} + \mathbf{k} \quad \mathbf{B} = \mathbf{j} + (\partial z/\partial y)\mathbf{k} = \mathbf{j} + 2\mathbf{k} \quad \mathbf{N} = -\mathbf{i} - 2\mathbf{j} + \mathbf{k}.$$  

The length of $\mathbf{N} = \mathbf{A} \times \mathbf{B}$ is $\sqrt{6}$. The angle between $\mathbf{k}$ and $\mathbf{N}$ has $\cos \theta = 1/\sqrt{6}$. That is the angle between base plane and sloping plane. Therefore the sloping area is $\sqrt{6}$ times the base area. For curved surfaces the idea is the same, except that the square root in $|\mathbf{N}| = 1/\cos \theta$ changes as we move around the surface.

**Method 2 is for curves $x(t)$, $y(t)$ and surfaces $x(u, v)$, $y(u, v)$, $z(u, v)$ with parameters.**

A curve has one parameter $t$. A surface has two parameters $u$ and $v$ (it is two-dimensional). One advantage of parameters is that $x, y, z$ get equal treatment, instead of picking out $z$ as $f(x, y)$. Here are the first two examples:

- **cone** $x = u \cos v$, $y = u \sin v$, $z = u$  
  **cylinder** $x = \cos v$, $y = \sin v$, $z = u$. \quad (4)

Each choice of $u$ and $v$ gives a point on the surface. By making all choices, we get the complete surface. Notice that a parameter can equal a coordinate, as in $z = u$. Sometimes both parameters are coordinates, as in $x = u$ and $y = v$ and $z = f(u, v)$. That is just $z = f(x, y)$ in disguise—the surface without parameters. In other cases we find the $xyz$ equation by eliminating $u$ and $v$:

- **cone** $(u \cos v)^2 + (u \sin v)^2 = u^2$ or $x^2 + y^2 = z^2$ or $z = \sqrt{x^2 + y^2}$

- **cylinder** $(\cos v)^2 + (\sin v)^2 = 1$ or $x^2 + y^2 = 1$.
The cone is the graph of \( f = \sqrt{x^2 + y^2} \). The cylinder is not the graph of any function. There is a line of \( z \)'s through each point on the circle \( x^2 + y^2 = 1 \). That is what \( z = u \) tells us: Give \( u \) all values, and you get the whole line. Give \( u \) and \( v \) all values, and you get the whole cylinder. Parameters allow a surface to close up and even go through itself—which the graph of \( f(x, y) \) can never do.

Actually \( z = \sqrt{x^2 + y^2} \) gives only the top half of the cone. (A function produces only one \( z \).) The parametric form gives the bottom half also. Similarly \( y = \sqrt{1 - x^2} \) gives only the top of a circle, while \( x = \cos t, y = \sin t \) goes all the way around.

Now we find \( dS \), using parameters. Small movements give a piece of the surface, practically flat. One side comes from the change \( du \), the neighboring side comes from \( dv \). The two sides are given by small vectors \( \mathbf{A}du \) and \( \mathbf{B}dv \):

\[
\mathbf{A} = \left( \frac{\partial x}{\partial u} i + \frac{\partial y}{\partial u} j + \frac{\partial z}{\partial u} k \right) \text{ and } \mathbf{B} = \left( \frac{\partial x}{\partial v} i + \frac{\partial y}{\partial v} j + \frac{\partial z}{\partial v} k \right). \tag{5}
\]

To find the area \( dS \) of the parallelogram, start with the cross product \( \mathbf{N} = \mathbf{A} \times \mathbf{B} \):

\[
\left| \begin{array}{ccc}
i & j & k \\
x_u & y_u & z_u \\
x_v & y_v & z_v
\end{array} \right| = \left( \frac{\partial y \partial z}{\partial u \partial v} - \frac{\partial z \partial y}{\partial u \partial v} \right)i + \left( \frac{\partial z \partial x}{\partial u \partial v} - \frac{\partial x \partial z}{\partial u \partial v} \right)j + \left( \frac{\partial x \partial y}{\partial u \partial v} - \frac{\partial y \partial x}{\partial u \partial v} \right)k. \tag{6}
\]

Admittedly this looks complicated—actual examples are often fairly simple. The area \( dS \) of the small piece of surface is \( |\mathbf{N}| du dv \). The length \( |\mathbf{N}| \) is a square root:

\[
dS = \sqrt{\left( \frac{\partial y \partial z}{\partial u \partial v} - \frac{\partial z \partial y}{\partial u \partial v} \right)^2 + \left( \frac{\partial z \partial x}{\partial u \partial v} - \frac{\partial x \partial z}{\partial u \partial v} \right)^2 + \left( \frac{\partial x \partial y}{\partial u \partial v} - \frac{\partial y \partial x}{\partial u \partial v} \right)^2} du dv. \tag{7}
\]

**EXAMPLE 2** Find \( \mathbf{A} \) and \( \mathbf{B} \) and \( \mathbf{N} = \mathbf{A} \times \mathbf{B} \) and \( dS \) for the cone and cylinder.

The cone has \( x = u \cos v, y = u \sin v, z = u \). The \( u \) derivatives produce \( \mathbf{A} = \partial \mathbf{R} / \partial u = \cos v \mathbf{i} + \sin v \mathbf{j} + \mathbf{k} \). The \( v \) derivatives produce the other tangent vector \( \mathbf{B} = \partial \mathbf{R} / \partial v = -u \sin v \mathbf{i} + u \cos v \mathbf{j}. \) The normal vector is \( \mathbf{A} \times \mathbf{B} = -u \cos v \mathbf{i} - u \sin v \mathbf{j} + u \mathbf{k} \). Its length gives \( dS \):

\[
dS = |\mathbf{A} \times \mathbf{B}| du dv = \sqrt{(u \cos v)^2 + (u \sin v)^2 + u^2} du dv = \sqrt{2} u du dv.
\]

The cylinder is even simpler: \( dS = du dv \). In these and many other examples, \( \mathbf{A} \) is perpendicular to \( \mathbf{B} \). The small piece is a rectangle. Its sides have length \( |\mathbf{A}| du \) and \( |\mathbf{B}| dv \). (The cone has \( |\mathbf{A}| = u \) and \( |\mathbf{B}| = \sqrt{2} \), the cylinder has \( |\mathbf{A}| = |\mathbf{B}| = 1 \). The cross product is hardly needed for area, when we can just multiply \( |\mathbf{A}| du \) times \( |\mathbf{B}| dv \).

**Remark on the two methods** Method 1 also used parameters, but a very special choice—\( u \) is \( x \) and \( v \) is \( y \). The parametric equations are \( x = x, y = y, z = f(x, y) \). If you go through the long square root in (7), changing \( u \) to \( x \) and \( v \) to \( y \), it simplifies to the square root in (3). (The terms \( \partial y / \partial x \) and \( \partial x / \partial y \) are zero; \( \partial x / \partial x \) and \( \partial y / \partial y \) are 1.) Still it pays to remember the shorter formula from Method 1.

Don't forget that after computing \( dS \), you have to integrate it. Many times the good way is with polar coordinates. Surfaces are often symmetric around an axis or a point. Those are the surfaces of revolution—which we saw in Chapter 8 and will come back to.

Strictly speaking, the integral starts with \( \Delta S \) (not \( dS \)). A flat piece has area \( |\mathbf{A} \times \mathbf{B}| \Delta \mathbf{A} \Delta \mathbf{y} \) or \( |\mathbf{A} \times \mathbf{B}| \Delta u \Delta v \). The area of a curved surface is properly defined as a limit. The key step of calculus, from sums of \( \Delta S \) to the integral of \( dS \), is safe for
smooth surfaces. In examples, the hard part is computing the double integral and substituting the limits on x, y or u, v.

**EXAMPLE 3** Find the surface area of the cone $z = \sqrt{x^2 + y^2}$ up to the height $z = a$.

We use Method 1 (no parameters). The derivatives of $z$ are computed, squared, and added:

$$
\frac{\partial z}{\partial x} = \frac{x}{\sqrt{x^2 + y^2}} \quad \frac{\partial z}{\partial y} = \frac{y}{\sqrt{x^2 + y^2}} \quad |N|^2 = 1 + \frac{x^2}{x^2 + y^2} + \frac{y^2}{x^2 + y^2} = 2.
$$

Conclusion: $|N| = \sqrt{2}$ and $dS = \sqrt{2} \, dx \, dy$. The cone is on a $45^\circ$ slope, so the area $dx \, dy$ in the base is multiplied by $\sqrt{2}$ in the surface above it (Figure 15.15). The square root in $dS$ accounts for the extra area due to slope. A horizontal surface has $dS = \sqrt{1} \, dx \, dy$, as we have known all year.

Now for a key point. **The integration is down in the base plane.** The limits on $x$ and $y$ are given by the “shadow” of the cone. To locate that shadow set $z = \sqrt{x^2 + y^2}$ equal to $z = a$. The plane cuts the cone at the circle $x^2 + y^2 = a^2$. We integrate over the inside of that circle (where the shadow is):

$$
\text{surface area of cone} = \int_{\text{shadow}} \sqrt{2} \, dx \, dy = \sqrt{2} \pi a^2.
$$

**EXAMPLE 4** Find the same area using $dS = \sqrt{2} \, u \, du \, dv$ from Example 2.

With parameters, $dS$ looks different and the shadow in the base looks different. The circle $x^2 + y^2 = a^2$ becomes $u^2 \cos^2 v + u^2 \sin^2 v = a^2$. In other words $u = a$. (The cone has $z = u$, the plane has $z = a$, they meet when $u = a$.) The angle parameter $v$ goes from 0 to $2\pi$. The effect of these parameters is to switch us “automatically” to polar coordinates, where area is $r \, dr \, d\theta$:

$$
\text{surface area of cone} = \int_{0}^{2\pi} \int_{0}^{a} \sqrt{2} \, u \, du \, dv = \sqrt{2} \pi a^2.
$$

Fig. 15.15 Cone cut by plane leaves shadow in the base. Integrate over the shadow.

**EXAMPLE 5** Find the area of the same cone up to the sloping plane $z = 1 - \frac{1}{2} x$.

Solution The cone still has $dS = \sqrt{2} \, dx \, dy$, but the limits of integration are changed. The plane cuts the cone in an ellipse. Its shadow down in the $xy$ plane is another ellipse (Figure 15.15c). To find the edge of the shadow, set $z = \sqrt{x^2 + y^2}$ equal to $z = 1 - \frac{1}{2} x$. We square both sides:

$$
x^2 + y^2 = 1 - x + \frac{1}{4} x^2 \quad \text{or} \quad \frac{3}{4} (x + \frac{3}{2})^2 + y^2 = \frac{3}{4}.
$$
This is the ellipse in the base—where height makes no difference and \( z \) is gone. The area of an ellipse is \( \pi ab \), when the equation is in the form \( \frac{x^2}{a^2} + \frac{y^2}{b^2} = 1 \). After multiplying by \( 3\pi/4 \) we find \( a = 4 \) and \( b = 4\sqrt{2}/3 \). Then \( \int \int \sqrt{2} \, dx \, dy = \sqrt{2} \pi ab \) is the surface area of the cone.

The hard part was finding the shadow ellipse (I went quickly). Its area \( \pi ab \) came from Example 15.3.2. The new part is \( \sqrt{2} \) from the slope.

**Example 6** Find the surface area of a sphere of radius \( a \) (known to be \( 4\pi a^2 \)).

This is a good example, because both methods almost work. The equation of the sphere is \( x^2 + y^2 + z^2 = a^2 \). Method 1 writes \( z = \sqrt{a^2 - x^2 - y^2} \). The \( x \) and \( y \) derivatives are \( -\frac{x}{z} \) and \( -\frac{y}{z} \):

\[
1 + \left( \frac{\partial z}{\partial x} \right)^2 + \left( \frac{\partial z}{\partial y} \right)^2 = \frac{z^2}{z^2} + \frac{x^2}{z^2} + \frac{y^2}{z^2} = \frac{a^2}{a^2 - x^2 - y^2}.
\]

The square root gives \( dS = a \, dx \, dy / \sqrt{a^2 - x^2 - y^2} \). Notice that \( z \) is gone (as it should be). Now integrate \( dS \) over the shadow of the sphere, which is a circle. Instead of \( dx \, dy \), switch to polar coordinates and \( r \, dr \, d\theta \):

\[
\int_{\text{shadow}} dS = \int_0^{2\pi} \int_0^a \frac{r \, dr \, d\theta}{\sqrt{a^2 - r^2}} = -2\pi a \sqrt{a^2 - r^2} \bigg|_0^a = 2\pi a^2.
\]

This calculation is successful but wrong. \( 2\pi a^2 \) is the area of the half-sphere above the \( xy \) plane. The lower half takes the negative square root of \( z^2 = a^2 - x^2 - y^2 \). This shows the danger of Method 1, when the surface is not the graph of a function.

**Example 7** (same sphere by Method 2: use parameters) The natural choice is spherical coordinates. Every point has an angle \( \theta = \phi \) down from the North Pole and an angle \( \phi = \theta \) around the equator. The \( xyz \) coordinates from Section 14.4 are \( x = a \sin \phi \cos \theta, y = a \sin \phi \sin \theta, z = a \cos \phi \). The radius \( \rho = a \) is fixed (not a parameter). Compute the first term in equation (6), noting \( \partial z / \partial \theta = 0 \):

\[
(\partial y / \partial \theta)(\partial z / \partial \phi) - (\partial z / \partial \phi)(\partial y / \partial \theta) = -(a \sin \phi)(a \sin \phi \cos \theta) = a^2 \sin^2 \phi \cos \theta.
\]

The other terms in (6) are \( a^2 \sin^2 \phi \sin \theta \) and \( a^2 \sin \phi \cos \phi \). Then \( dS \) in equation (7) squares these three components and adds. We factor out \( a^4 \) and simplify:

\[
a^4(\sin^4 \phi \cos^2 \theta + \sin^4 \phi \sin^2 \theta + \sin^2 \phi \cos^2 \phi) = a^4(\sin^4 \phi + \sin^2 \phi \cos^2 \phi) = a^4 \sin^2 \phi.
\]

Conclusion: \( dS = a^2 \sin \phi \, d\phi \, d\theta \). A spherical person will recognize this immediately. It is the volume element \( dV = \rho^2 \sin \phi \, d\rho \, d\phi \, d\theta \), except \( d\rho \) is missing. The small box has area \( dS \) and thickness \( d\rho \) and volume \( dV \). Here we only want \( dS \):

\[
\text{area of sphere} = \int_0^{2\pi} \int_0^a dS = \int_0^{2\pi} \int_0^a a^3 \sin \phi \, d\phi \, d\theta = 4\pi a^2.
\]

Figure 15.16a shows a small surface with sides \( a \, d\phi \) and \( a \sin \phi \, d\theta \). Their product is \( dS \). Figure 15.16b goes back to Method 1, where equation (8) gave \( dS = (a^2 \frac{z}{a}) \, dx \, dy \).

I doubt that you will like Figure 15.16c—and you don't need it. With parameters \( \phi \) and \( \theta \), the shadow of the sphere is a rectangle. The equator is the line down the middle, where \( \phi = \pi/2 \). The height is \( z = a \cos \phi \). The area \( d\phi \, d\theta \) in the base is the shadow of \( dS = a^2 \sin \phi \, d\phi \, d\theta \) up in the sphere. Maybe this figure shows what we don't have to know about parameters.
**EXAMPLE 8**  Rotate \( y = x^2 \) around the \( x \) axis. Find the surface area using parameters.

The first parameter is \( x \) (from \( a \) to \( b \)). The second parameter is the rotation angle \( \theta \) (from 0 to \( 2\pi \)). The points on the surface in Figure 15.17 are \( x = x, y = x^2 \cos \theta, z = x^2 \sin \theta \). Equation (7) leads after much calculation to \( dS = x^2 / \sqrt{1 + 4x^2} \, dx \, d\theta \).

*Main point:* \( dS \) agrees with Section 8.3, where the area was \( \int 2\pi y \sqrt{1 + (dy/dx)^2} \, dx \). The \( 2\pi \) comes from the \( \theta \) integral and \( y \) is \( x^2 \). Parameters give this formula automatically.

**VECTOR FIELDS AND THE INTEGRAL OF \( F \cdot n \)**

Formulas for surface area are dominated by square roots. There is a square root in \( dS \), as there was in \( ds \). Areas are like arc lengths, one dimension up. The good point about line integrals \( \int F \cdot ds \) is that the square root disappears. It is in the denominator of \( n \), where \( ds \) cancels it: \( F \cdot nds = M \, dy - N \, dx \). The same good thing will now happen for surface integrals \( \iint F \cdot ndS \).

Through the surface \( z = f(x, y) \), the vector field \( F(x, y, z) = Mi + Nj +Pk \) has

\[
\text{flux} = \iint_{\text{surface}} F \cdot ndS = \iint_{\text{shadow}} \left( -M \frac{\partial f}{\partial x} - N \frac{\partial f}{\partial y} + P \right) \, dx \, dy. \tag{10}
\]

This formula tells what to integrate, given the surface and the vector field \( (f \) and \( F) \). The \( xy \) limits come from the shadow. Formula (10) takes the normal vector from Method 1:

\[
N = -\frac{\partial f}{\partial x} \mathbf{i} - \frac{\partial f}{\partial y} \mathbf{j} + \mathbf{k} \quad \text{and} \quad |N| = \sqrt{1 + (\partial f/\partial x)^2 + (\partial f/\partial y)^2}.
\]

For the unit normal vector \( n \), divide \( N \) by its length: \( n = N/|N| \). The square root is in the denominator, and the same square root is in \( dS \). See equation (3):

\[
F \cdot ndS = \frac{F \cdot N}{\sqrt{|N|}} \, dx \, dy = \left( -M \frac{\partial f}{\partial x} - N \frac{\partial f}{\partial y} + P \right) dx \, dy. \tag{11}
\]

That is formula (10), with cancellation of square roots. The expression \( F \cdot ndS \) is often written as \( F \cdot dS \), again relying on boldface to make \( dS \) a vector. Then \( dS \) equals \( ndS \), with direction \( n \) and magnitude \( dS \).
EXAMPLE 9  Find \( ndS \) for the plane \( z = x + 2y \). Then find \( F \cdot ndS \) for \( F = k \).

This plane produced \( \sqrt{6} \) in Example 1 (for area). For flux the \( \sqrt{6} \) disappears:

\[ ndS = \frac{N}{|N|} dS = \frac{-i - 2j + k}{\sqrt{6}} d\sqrt{6} \, dx \, dy = (-i - 2j + k) \, dx \, dy. \]

For the flow field \( F = k \), the dot product \( k \cdot ndS \) reduces to \( 1 dx \, dy \). The slope of the plane makes no difference! The flow through the base also flows through the plane. The areas are different, but flux is like rain. Whether it hits a tent or the ground below, it is the same rain (Figure 15.18). In this case \( \iiint F \cdot ndS = \iint dx \, dy = \) shadow area in the base.

EXAMPLE 10  Find the flux of \( F = xi + yj + zk \) through the cone \( z = \sqrt{x^2 + y^2} \).

Solution  \( F \cdot ndS = \left[ -x \left( \frac{x}{\sqrt{x^2 + y^2}} \right) - y \left( \frac{y}{\sqrt{x^2 + y^2}} \right) + \sqrt{x^2 + y^2} \right] dx \, dy = 0. \)

The zero comes as a surprise, but it shouldn't. The cone goes straight out from the origin, and so does \( F \). The vector \( n \) that is perpendicular to the cone is also perpendicular to \( F \). There is no flow through the cone, because \( F \cdot n = 0 \). The flow travels out along rays.

\[ \iiint F \cdot ndS \] FOR A SURFACE WITH PARAMETERS

In Example 10 the cone was \( z = f(x, y) = \sqrt{x^2 + y^2} \). We found \( dS \) by Method 1. Parameters were not needed (more exactly, they were \( x \) and \( y \)). For surfaces that fold and twist, the formulas with \( u \) and \( v \) look complicated but the actual calculations can be simpler. This was certainly the case for \( dS = du \, dv \) on the cylinder.

A small piece of surface has area \( dS = |A \times B| du \, dv \). The vectors along the sides are \( A = x_i + y_j + z_k \) and \( B = x_i + y_j + z_k \). They are tangent to the surface. Now we put their cross product \( N = A \times B \) to another use, because \( F \cdot ndS \) involves not only area but direction. We need the unit vector \( n \) to see how much flow goes through.

The direction vector is \( n = N/|N| \). Equation (7) is \( dS = |N| du \, dv \), so the square root \( |N| \) cancels in \( ndS \). This leaves a nice formula for the "normal component" of flow:

\[ \iiint F \cdot ndS = \iiint F \cdot N \, du \, dv = \iiint F \cdot (A \times B) \, du \, dv. \]
EXAMPLE 11  Find the flux of \( \mathbf{F} = x\mathbf{i} + y\mathbf{j} + z\mathbf{k} \) through the cylinder \( x^2 + y^2 = 1, \ 0 \leq z \leq b \).

Solution. The surface of the cylinder is \( x = \cos u, y = \sin u, z = v \). The tangent vectors from (5) are \( \mathbf{A} = (-\sin u)\mathbf{i} + (\cos u)\mathbf{j} \) and \( \mathbf{B} = \mathbf{k} \). The normal vector in Figure 15.19 goes straight out through the cylinder:

\[
\mathbf{N} = \mathbf{A} \times \mathbf{B} = \cos u \mathbf{i} + \sin u \mathbf{j} \quad \text{(check \( \mathbf{A} \cdot \mathbf{N} = 0 \) and \( \mathbf{B} \cdot \mathbf{N} = 0 \)).}
\]

To find \( \mathbf{F} \cdot \mathbf{N} \), switch \( \mathbf{F} = x\mathbf{i} + y\mathbf{j} + z\mathbf{k} \) to the parameters \( u \) and \( v \). Then \( \mathbf{F} \cdot \mathbf{N} = 1 \):

\[
\mathbf{F} \cdot \mathbf{N} = (\cos u \mathbf{i} + \sin u \mathbf{j} + v\mathbf{k}) \cdot (\cos u \mathbf{i} + \sin u \mathbf{j}) = \cos^2 u + \sin^2 u.
\]

For the flux, integrate \( \mathbf{F} \cdot \mathbf{N} = 1 \) and apply the limits on \( u = 0 \) and \( v = z \):

\[\text{flux} = \int_0^b \int_0^{2\pi} 1 \ du \ dv = 2\pi b = \text{surface area of the cylinder.}\]

Note that the top and bottom were not included! We can find those fluxes too. The outward direction is \( \mathbf{n} = \mathbf{k} \) at the top and \( \mathbf{n} = -\mathbf{k} \) down through the bottom. Then \( \mathbf{F} \cdot \mathbf{n} = +z = b \) at the top and \( -z = 0 \) at the bottom. The bottom flux is zero, the top flux is \( b \) times the area (or \( \pi b \)). The total flux is \( 2\pi b + \pi b = 3\pi b \). Hold that answer for the next section.

Apology: I made \( u \) the angle and \( v \) the height. Then \( \mathbf{N} \) goes outward not inward.

EXAMPLE 12  Find the flux of \( \mathbf{F} = \mathbf{k} \) out the top half of the sphere \( x^2 + y^2 + z^2 = a^2 \).

Solution. Use spherical coordinates. Example 7 had \( u = \phi \) and \( v = \theta \). We found

\[
\mathbf{N} = \mathbf{A} \times \mathbf{B} = a^2 \sin^2 \phi \cos \theta \mathbf{i} + a^2 \sin^2 \phi \sin \theta \mathbf{j} + a^2 \sin \phi \cos \phi \mathbf{k}.
\]

The dot product with \( \mathbf{F} = \mathbf{k} \) is \( \mathbf{F} \cdot \mathbf{N} = a^2 \sin \phi \cos \phi \). The integral goes from the pole to the equator, \( \phi = 0 \) to \( \phi = \pi/2 \), and around from \( \theta = 0 \) to \( \theta = 2\pi \):

\[\text{flux} = \int_0^{2\pi} \int_0^{\pi/2} a^2 \sin \phi \cos \phi \ d\phi \ d\theta = 2\pi a^2 \frac{\sin^2 \phi}{2} \bigg|_0^{\pi/2} = \pi a^2.\]

The next section will show that the flux remains at \( \pi a^2 \) through any surface (!) that is bounded by the equator. A special case is a flat surface—the disk of radius \( a \) at the equator. Figure 15.18 shows \( \mathbf{n} = \mathbf{k} \) pointing directly up, so \( \mathbf{F} \cdot \mathbf{n} = \mathbf{k} \cdot \mathbf{k} = 1 \). The flux is \( \iint 1 \ dS = \text{area of disk} = \pi a^2 \). All fluid goes past the equator and out through the sphere.
I have to mention one more problem. It might not occur to a reasonable person, but sometimes a surface has only one side. The famous example is the Möbius strip, for which you take a strip of paper, twist it once, and tape the ends together. Its special property appears when you run a pen along the “inside.” The pen in Figure 15.20 suddenly goes “outside.” After another round trip it goes back “inside.” Those words are in quotation marks, because on a Möbius strip they have no meaning.

Suppose the pen represents the normal vector. On a sphere \( n \) points outward. Alternatively \( n \) could point inward; we are free to choose. But the Möbius strip makes the choice impossible. After moving the pen continuously, it comes back in the opposite direction. **This surface is not orientable.** We cannot integrate \( \mathbf{F} \cdot \mathbf{n} \) to compute the flux, because we cannot decide the direction of \( \mathbf{n} \).

A surface is **oriented** when we can and do choose \( \mathbf{n} \). This uses the final property of cross products, that they have length and direction and also a right-hand rule. We can tell \( \mathbf{A} \times \mathbf{B} \) from \( \mathbf{B} \times \mathbf{A} \). Those give the two orientations of \( \mathbf{n} \). For an open surface (like a wastebasket) you can select either one. For a closed surface (like a sphere) it is conventional for \( \mathbf{n} \) to be outward. By making that decision once and for all, the sign of the flux is established: outward flux is positive.

### Formulas for Surface Integrals

**Method 1:** Parameters \( x, y \)

<table>
<thead>
<tr>
<th>Coordinates ( x, y, z(x, y) )</th>
<th>( A = \mathbf{i} + \partial z/\partial x \mathbf{k} )</th>
</tr>
</thead>
<tbody>
<tr>
<td>( B = \mathbf{j} + \partial z/\partial y \mathbf{k} )</td>
<td>( N = A \times B )</td>
</tr>
<tr>
<td>( dS =</td>
<td>N</td>
</tr>
</tbody>
</table>

**Method 2:** Parameters \( u, v \)

<table>
<thead>
<tr>
<th>( x(u, v), y(u, v), z(u, v) ) on surface</th>
<th>( A = \partial x/\partial u \mathbf{i} + \partial y/\partial u \mathbf{j} + \partial z/\partial u \mathbf{k} )</th>
</tr>
</thead>
<tbody>
<tr>
<td>( B = \partial x/\partial v \mathbf{i} + \partial y/\partial v \mathbf{j} + \partial z/\partial v \mathbf{k} )</td>
<td>( N = A \times B )</td>
</tr>
<tr>
<td>( dS =</td>
<td>N</td>
</tr>
</tbody>
</table>

### Exerises

**Read-through questions**

A small piece of the surface \( z = f(x, y) \) is nearly \( \mathbf{0} \). When we go across by \( dx \), we go up by \( b \). That movement is \( A \, dx \), where the vector \( A = \mathbf{i} + c \). The other side of the piece is \( B \, dy \), where \( B = \mathbf{j} + d \). The cross product \( A \times B \) is \( \mathbf{N} = e \). The area of the piece is \( dS = |\mathbf{N}| \, dx \, dy \).

For the surface \( z = xy \), the vectors are \( \mathbf{A} = f \) and \( \mathbf{B} = g \) and \( \mathbf{N} = h \). The area integral is \( \iint dS = f \cdot \mathbf{n} \, dx \, dy \).

With parameters \( u \) and \( v \), a typical point on a 45\(^\circ\) cone is \( x = u \cos v, y = u \sin v, z = u \). The tangent vectors are \( \mathbf{A} = w \) and \( \mathbf{B} = x \). This cone has \( ndS = A \times B \, du \, dv \). For \( F = xi + yj + zk \), the flux element through the cone is \( F \cdot ndS = z \).

The reason for this answer is \( A \). The reason we don't compute flux through a Möbius strip is \( B \).

In 1–14 find \( \mathbf{n} \) and \( dS = |\mathbf{N}| \, dx \, dy \) and the surface area \( \iint dS \).

Integrate over the \( xy \) shadow which ends where the \( z \)'s are equal \( x^2 + y^2 = 4 \) in Problem 1.

1. Paraboloid \( z = x^2 + y^2 \) below the plane \( z = 4 \).
2. Paraboloid \( z = x^2 + y^2 \) between \( z = 4 \) and \( z = 8 \).
3. Plane \( z = x - y \) inside the cylinder \( x^2 + y^2 = 1 \).
4. Plane \( z = 3x + 4y \) above the square \( 0 \leq x \leq 1, 0 \leq y \leq 1 \).
In 5-18 compute the surface integrals $\iint g(x, y, z) dS$.

15. $g = xy$ over the triangle $x + y + z = 1$, $x, y, z \geq 0$.

16. $g = x^2 + y^2$ over the top half of $x^2 + y^2 + z^2 = 1$ (use $\phi, \theta$).

17. $g = xyz$ on $x^2 + y^2 + z^2 = 1$ above $z^2 = x^2 + y^2$ (use $\phi, \theta$).

18. $g = x$ on the cylinder $x^2 + y^2 = 4$ between $z = 0$ and $z = 3$.

In 19-22 calculate $A$, $B$, $N$, and $dS$.

19. $x = u$, $y = v + u$, $z = v + 2u + 1$.

20. $x = u$, $y = u + v$, $z = u - v$.

21. $x = (3 + \cos u) \cos v$, $y = (3 + \cos u) \sin v$, $z = \sin u$.

22. $x = u \cos v$, $y = u \sin v$, $z = v$ (not $z = u$).

23-26 In Problems 1-4 respectively find the flux $\iint F \cdot n dS$ for $F = xi + yj + zk$.

27-28 In Problems 19-20 respectively compute $\iint F \cdot n dS$ for $F = yi - xj$ through the region $u^2 + v^2 \leq 1$.

29. A unit circle is rotated around the $z$ axis to give a torus (see figure). The center of the circle stays a distance 3 from the $z$ axis. Show that Problem 21 gives a typical point $(x, y, z)$ on the torus and find the surface area $\iint dS = \iint |N| du dv$.

30. The surface $x = r \cos \theta$, $y = r \sin \theta$, $z = a^2 - r^2$ is bounded by the equator ($r = a$). Find $N$ and the flux $\iint k \cdot n dS$, and compare with Example 12.

31. Make a “double Möbius strip” from a strip of paper by twisting it twice and taping the ends. Does a normal vector (use a pen) have the same direction after a round trip?

32. Make a “triple Möbius strip” with three twists. Is it orientable—does the normal vector come back in the same or opposite direction?

33. If a very wavy surface stays close to a smooth surface, are their areas close?

34. Give the equation of a plane with roof area $dS = 3$ times base area $dx dy$.

35. The points $(x, f(x) \cos \theta, f(x) \sin \theta)$ are on the surface of revolution: $y = f(x)$ revolved around the $x$ axis, parameters $u = x$ and $v = \theta$. Find $N$ and compare $dS = |N| dx d\theta$ with Example 8 and Section 8.3.

15.5 The Divergence Theorem

This section returns to the fundamental law (flow out) - (flow in) = (source). In two dimensions, the flow was in and out through a closed curve $C$. The plane region inside was $R$. In three dimensions, the flow enters and leaves through a closed surface $S$. The solid region inside is $V$. Green's Theorem in its normal form (for the flux of a smooth vector field) now becomes the great three-dimensional balance equation—the Divergence Theorem:
15.5 The Divergence Theorem

The flux of \( F = M \mathbf{i} + N \mathbf{j} + P \mathbf{k} \) through the boundary surface \( S \) equals the integral of the divergence of \( F \) inside \( V \). The Divergence Theorem is

\[
\iiint_V \nabla \cdot F \, dV = 
\iiint_S \mathbf{F} \cdot n \, dS.
\]

In Green's Theorem the divergence was \( \frac{\partial M}{\partial x} + \frac{\partial N}{\partial y} \). The new term \( \frac{\partial P}{\partial z} \) accounts for upward flow. Notice that a constant upward component \( P \) adds nothing to the divergence (its derivative is zero). It also adds nothing to the flux (flow up through the top equals flow up through the bottom). When the whole field \( F \) is constant, the theorem becomes \( 0 = 0 \).

There are other vector fields with \( \nabla \cdot F = 0 \). They are of the greatest importance. The Divergence Theorem for those fields is again \( 0 = 0 \), and there is conservation of fluid. When \( \nabla \cdot F = 0 \), flow in equals flow out. We begin with examples of these "divergence-free" fields.

**EXAMPLE 1** The spin fields \(-yi + xj + 0k\) and \(6i - zj + yk\) have zero divergence.

The first is an old friend, spinning around the \( z \) axis. The second is new, spinning around the \( x \) axis. Three-dimensional flow has a great variety of spin fields. The separate terms \( \frac{\partial M}{\partial x}, \frac{\partial N}{\partial y}, \frac{\partial P}{\partial z} \) are all zero, so \( \nabla \cdot F = 0 \). The flow goes around in circles, and whatever goes out through \( S \) comes back in. (We might have put a circle on \( \Sigma \), as we did on \( \Sigma_c \), to emphasize that \( S \) is closed.)

**EXAMPLE 2** The position field \( \mathbf{R} = xi + yj + zk \) has \( \nabla \cdot \mathbf{R} = 1 + 1 + 1 = 3 \).

This is radial flow, straight out from the origin. Mass has to be added at every point to keep the flow going. On the right side of the divergence theorem is \( \iiint 3 \, dV \). Therefore the flux is three times the volume.

Example 11 in Section 15.4 found the flux of \( \mathbf{R} \) through a cylinder. The answer was \( 3\pi b \). Now we also get \( 3\pi b \) from the Divergence Theorem, since the volume is \( \pi b \). This is one of many cases in which the triple integral is easier than the double integral.

**EXAMPLE 3** An electrostatic field \( \mathbf{R}/\rho^3 \) or gravity field \( -\mathbf{R}/\rho^3 \) almost has \( \nabla \cdot F = 0 \).

The vector \( \mathbf{R} = xi + yj + zk \) has length \( \sqrt{x^2 + y^2 + z^2} = \rho \). Then \( F \) has length \( \rho/\rho^3 \) (inverse square law). Gravity from a point mass pulls inward (minus sign). The electric field from a point charge repels outward. The three steps almost show that \( \nabla \cdot F = 0 \):

Step 1. \( \frac{\partial \rho}{\partial x} = x/\rho, \frac{\partial \rho}{\partial y} = y/\rho, \frac{\partial \rho}{\partial z} = z/\rho \) — but do not add those three. \( F \) is not \( \rho \) or \( 1/\rho^2 \) (these are scalars). The vector field is \( \mathbf{R}/\rho^3 \). We need \( \frac{\partial M}{\partial x}, \frac{\partial N}{\partial y}, \frac{\partial P}{\partial z} \).

Step 2. \( \frac{\partial M}{\partial x} = \frac{\partial}{\partial x}(x\rho/\rho^3) \) is equal to \( 1/\rho^3 - (3x \rho/\rho^3)/3x^2/\rho^5 \). For \( \frac{\partial N}{\partial y} \) and \( \frac{\partial P}{\partial z} \), replace \( 3x^2 \) by \( 3y^2 \) and \( 3z^2 \). Now add those three.

Step 3. \( \nabla \cdot F = 3/\rho^3 - 3(x^2 + y^2 + z^2)/\rho^5 = 3/\rho^3 - 3/\rho^3 = 0 \).

The calculation \( \nabla \cdot F = 0 \) leaves a puzzle. One side of the Divergence Theorem seems to give \( \iiint 0 \, dV = 0 \). Then the other side should be \( \iiint \mathbf{F} \cdot n \, dS = 0 \). But the flux is not zero when all flow is outward:

The unit normal vector to the sphere \( \rho = \) constant is \( n = \mathbf{R}/\rho \).

The outward flow \( \mathbf{F} \cdot n = (\mathbf{R}/\rho^3) \cdot (\mathbf{R}/\rho) = \rho^2/\rho^4 \) is always positive.

Then \( \iiint \mathbf{F} \cdot n \, dS = \iiint dS/\rho^2 = 4\pi \rho^2/\rho^2 = 4\pi \). We have reached \( 4\pi = 0 \).
This paradox in three dimensions is the same as for $R/r^2$ in two dimensions. Section 15.3 reached $2\pi = 0$, and the explanation was a point source at the origin. Same explanation here: $M, N, P$ are infinite when $\rho = 0$. The divergence is a “delta function” times $4\pi$, from the point source. The Divergence Theorem does not apply (unless we allow delta functions). That single point makes all the difference.

**Every surface enclosing the origin has** flux $= 4\pi$. Our calculation was for a sphere. The surface integral is much harder when $S$ is twisted (Figure 15.21a). But the Divergence Theorem takes care of everything, because $\text{div } F = 0$ in the volume $V$ between these surfaces. Therefore $\iint F \cdot ndS = 0$ for the two surfaces together. The flux $\iiint F \cdot ndS = -4\pi$ into the sphere must be balanced by $\iiint F \cdot ndS = 4\pi$ out of the twisted surface.

Fig. 15.21 Point source: flux $4\pi$ through all enclosing surfaces. Net flux upward $= \iint (\partial P/\partial z)dV$.

Instead of a paradox $4\pi = 0$, this example leads to Gauss's Law. A mass $M$ at the origin produces a gravity field $F = -GMR/\rho^3$. A charge $q$ at the origin produces an electric field $E = (q/4\pi\varepsilon_0)R/\rho^3$. The physical constants are $G$ and $\varepsilon_0$, the mathematical constant is the relation between divergence and flux. Equation (1) yields equation (2), in which the mass densities $M(x, y, z)$ and charge densities $q(x, y, z)$ need not be concentrated at the origin:

| 15L | Gauss's law in differential form: $\text{div } F = -4\pi GM$ and $\text{div } E = q/\varepsilon_0$. |
| 15L | Gauss's law in integral form: Flux is proportional to total mass or charge: $\iiint F \cdot ndS = -\iiint 4\pi GMdV$ and $\iiint E \cdot ndS = \iiint qdV/\varepsilon_0$. |

(2)

**THE REASONING BEHIND THE DIVERGENCE THEOREM**

The general principle is clear: Flow out minus flow in equals source. Our goal is to see why the divergence of $F$ measures the source. In a small box around each point, we show that $\text{div } F \ dV$ balances $F \cdot ndS$ through the six sides. So consider a small box. Its center is at $(x, y, z)$. Its edges have length $\Delta x, \Delta y, \Delta z$. Out of the top and bottom, the normal vectors are $k$ and $-k$. The dot product with $F = M_i + Nj + Pk$ is $+P$ or $-P$. The area $\Delta S$ is $\Delta x\Delta y$. So the two fluxes are close to $P(x, y, z + \frac{1}{2}\Delta z)\Delta x\Delta y$ and $-P(x, y, z - \frac{1}{2}\Delta z)\Delta x\Delta y$. When the top is combined with the bottom, the difference of those $P$'s is $\Delta P$:

$\text{net flux upward } \approx \Delta P \Delta x\Delta y = (\Delta P/\Delta z)\Delta x\Delta y\Delta z \approx (\partial P/\partial z)\Delta V$.  

(3)
Similarly, the combined flux on two side faces is approximately \((\partial M/\partial x)\Delta V\). On the front and back it is \((\partial N/\partial y)\Delta V\). Adding the six faces, we reach the key point:

\[
\text{flux out of the box} \approx (\partial M/\partial x + \partial N/\partial y + \partial P/\partial z)\Delta V. \tag{4}
\]

This is \((\text{div } F)\Delta V\). For a constant field both sides are zero—the flow goes straight through. For \(F = xi + yj + zk\), a little more goes out than comes in. The divergence is 3, so \(3\Delta V\) is created inside the box. By the balance equation the flux is also \(3\Delta V\).

The approximation symbol \(\approx\) means that the leading term is correct (probably not the next term). The ratio \(\Delta P/\Delta z\) is not exactly \(\partial P/\partial z\). The difference is of order \(\Delta z\), so the error in (3) is of higher order \(\Delta V\Delta z\). Added over many boxes (about \(1/\Delta V\) boxes), this error disappears as \(\Delta z \to 0\).

The sum of \((\text{div } F)\Delta V\) over all the boxes approaches \(\iiint (\text{div } F) dV\). On the other side of the equation is a sum of fluxes. There is \(F \cdot n\Delta S\) out of the top of one box, plus \(F \cdot n\Delta S\) out of the bottom of the box above. The first has \(n = k\) and the second has \(n = -k\). They cancel each other—the flow goes from box to box. This happens every time two boxes meet. The only fluxes that survive (because nothing cancels them) are at the outer surface \(S\). The final step, as \(\Delta x, \Delta y, \Delta z \to 0\), is that those outside terms approach \(\iint F \cdot n dS\). Then the local divergence theorem (4) becomes the global Divergence Theorem (1).

**Remark on the proof** That "final step" is not easy, because the box surfaces don't line up with the outer surface \(S\). A formal proof of the Divergence Theorem would imitate the proof of Green's Theorem. On a very simple region \(\iiint (\partial P/\partial z) dx\, dy\, dz\) equals \(\iint P\, dx\, dy\) over the top minus \(\iint P\, dx\, dy\) over the bottom. After checking the orientation this is \(\iint P k \cdot ndS\). Similarly the volume integrals of \(\partial M/\partial z\) and \(\partial N/\partial y\) are the surface integrals \(\iint M_i \cdot ndS\) and \(\iint N_j \cdot ndS\). Adding the three integrals gives the Divergence Theorem. Since Green's Theorem was already proved in this way, the reasoning behind (4) is more helpful than repeating a detailed proof.

The discoverer of the Divergence Theorem was probably Gauss. His notebooks only contain the outline of a proof—but after all, this is Gauss. Green and Ostrogradsky both published proofs in 1828, one in England and the other in St. Petersburg (now Leningrad). As the theorem was studied, the requirements came to light (smoothness of \(F\) and \(S\), avoidance of one-sided Möbius strips).

New applications are discovered all the time—when a scientist writes down a balance equation in a small box. The source is known. The equation is \(\text{div } F = \text{source}\). After Example 5 we explain \(F\).

**EXAMPLE 4** If the temperature inside the sun is \(T = \ln 1/\rho\), find the heat flow \(F = -\text{grad } T\) and the source \(\text{div } F\) and the flux \(\iint F \cdot ndS\). The sun is a ball of radius \(\rho_0\).

**Solution** \(F = -\text{grad } \ln 1/\rho = +\text{grad } \ln \rho\). Derivatives of \(\ln \rho\) bring division by \(\rho\):

\[
F = (\partial \rho/\partial x i + \partial \rho/\partial y j + \partial \rho/\partial z k)/\rho = (xi + yj + zk)/\rho^2.
\]

This flow is radially outward, of magnitude \(1/\rho\). The normal vector \(n\) is also radially outward, of magnitude 1. The dot product on the sun's surface is \(1/\rho_0\):

\[
\iint F \cdot ndS = \iiint dS/\rho_0 = \text{(surface area)}/\rho_0 = 4\pi\rho_0^2/\rho_0 = 4\pi\rho_0. \tag{5}
\]

Check that answer by the Divergence Theorem. Example 5 will find \(\text{div } F = 1/\rho^2\). Integrate over the sun. In spherical coordinates we integrate \(dp\), \(\sin \phi d\phi\), and \(d\theta\):

\[
\iiint_{\text{sun}} F \, dV = \int_0^{2\pi} \int_0^\pi \int_0^{\rho_0} \rho^2 \sin \phi \, d\rho \, d\phi \, d\theta/\rho^2 = (\rho_0)(2\pi) \text{ as in (5).}
\]
This example illustrates the basic framework of equilibrium. The pattern appears everywhere in applied mathematics—electromagnetism, heat flow, elasticity, even relativity. There is usually a constant \( c \) that depends on the material (the example has \( c = 1 \)). The names change, but we always take the divergence of the gradient:

- \( \text{potential } f \to \text{force field } -c \text{ grad } f \). Then \( \text{div}(-c \text{ grad } f) = \text{electric charge} \)
- \( \text{temperature } T \to \text{flow field } -c \text{ grad } T \). Then \( \text{div}(-c \text{ grad } T) = \text{heat source} \)
- \( \text{displacement } u \to \text{stress field } -c \text{ grad } u \). Then \( \text{div}(-c \text{ grad } u) = \text{outside force} \).

You are studying calculus, not physics or thermodynamics or elasticity. But please notice the main point. The equation to solve is \( \text{div}(-c \text{ grad } f) = \text{known source} \). The divergence and gradient are exactly what the applications need. Calculus teaches the right things.

This framework is developed in many books, including my own text Introduction to Applied Mathematics (Wellesley-Cambridge Press). It governs equilibrium, in matrix equations and differential equations.

**PRODUCT RULE FOR VECTORS: INTEGRATION BY PARTS**

May I go back to basic facts about the divergence? First the definition:

\[
F(x, y, z) = M\mathbf{i} + N\mathbf{j} + P\mathbf{k} \text{ has } \text{div } F = \nabla \cdot F = \frac{\partial M}{\partial x} + \frac{\partial N}{\partial y} + \frac{\partial P}{\partial z}.
\]

The divergence is a scalar (not a vector). At each point \( \text{div } F \) is a number. In fluid flow, it is the rate at which mass leaves—the "flux per unit volume" or "flux density."

The symbol \( \nabla \) stands for a vector whose components are operations not numbers:

\[
\nabla = \text{"del" } = i \frac{\partial}{\partial x} + j \frac{\partial}{\partial y} + k \frac{\partial}{\partial z}.
\]

This vector is illegal but very useful. First, apply it to an ordinary function \( f(x, y, z) \):

\[
\nabla f = \text{"del } f" = i \frac{\partial f}{\partial x} + j \frac{\partial f}{\partial y} + k \frac{\partial f}{\partial z} = \text{gradient of } f.
\]

Second, take the dot product \( \nabla \cdot F \) with a vector function \( F(x, y, z) = M\mathbf{i} + N\mathbf{j} + P\mathbf{k} \):

\[
\nabla \cdot F = \text{"del dot } F" = \frac{\partial M}{\partial x} + \frac{\partial N}{\partial y} + \frac{\partial P}{\partial z} = \text{divergence of } F.
\]

Third, take the cross product \( \nabla \times F \). This produces the vector curl \( F \) (next section):

\[
\nabla \times F = \text{"del cross } F" = \ldots (\text{to be defined}) \ldots = \text{curl of } F.
\]

The gradient and divergence and curl are \( \nabla \) and \( \nabla \cdot \) and \( \nabla \times \). The three great operations of vector calculus use a single notation! You are free to write \( \nabla \) or not—to make equations shorter or to help the memory. Notice that Laplace's equation shrinks to

\[
\nabla \cdot \nabla f = \frac{\partial^2 f}{\partial x^2} + \frac{\partial^2 f}{\partial y^2} + \frac{\partial^2 f}{\partial z^2} = 0.
\]

Equation (10) gives the potential when the source is zero (very common). \( F = \text{grad } f \) combines with \( \text{div } F = 0 \) into Laplace's equation \( \text{div } \text{grad } f = 0 \). This equation is so important that it shrinks further to \( \nabla^2 f = 0 \) and even to \( \Delta f = 0 \). Of course \( \Delta f = f_{xx} + f_{yy} + f_{zz} \) has nothing to do with \( \Delta f = f(x + \Delta x) - f(x) \). Above all, remember that \( f \) is a scalar and \( F \) is a vector: gradient of scalar is vector and divergence of vector is scalar.
15.5 The Divergence Theorem

Underlying this chapter is the idea of extending calculus to vectors. So far we have emphasized the Fundamental Theorem. The integral of \( df/dx \) is now the integral of \( \text{div} \, F \). Instead of endpoints \( a \) and \( b \), we have a curve \( C \) or surface \( S \). But it is the rules for derivatives and integrals that make calculus work, and we need them now for vectors. Remember the derivative of \( u \) times \( v \) and the integral (by parts) of \( u \, dv/dx \):

Scalar functions \( u(x, y, z) \) and vector fields \( V(x, y, z) \) obey the product rule:

\[
\text{div}(uV) = u \, \text{div} \, V + V \cdot (\text{grad} \, u).
\] (11)

The reverse of the product rule is integration by parts (Gauss's Formula):

\[
\int \int \int u \, \text{div} \, V \, dx \, dy \, dz = - \int \int \int V \cdot (\text{grad} \, u) \, dx \, dy \, dz + \int \int u \, V \cdot n \, dS.
\] (12)

For a plane field this is Green's Formula (and \( u = 1 \) gives Green's Theorem):

\[
\int \int \left( \frac{\partial M}{\partial x} + \frac{\partial N}{\partial y} \right) \, dx \, dy = - \int \int \left( M \frac{\partial u}{\partial x} + N \frac{\partial u}{\partial y} \right) \, dx \, dy + \int u(Mi + Nj) \cdot n \, ds.
\] (13)

Those look like heavy formulas. They are too much to memorize, unless you use them often. The important point is to connect vector calculus with "scalar calculus," which is not heavy. Every product rule yields two terms:

\[
(uM)_x = u \, \frac{\partial M}{\partial x} + M \, \frac{\partial u}{\partial x} \quad (uN)_y = u \, \frac{\partial N}{\partial y} + N \, \frac{\partial u}{\partial y} \quad (uP)_z = u \, \frac{\partial P}{\partial z} + P \, \frac{\partial u}{\partial z}.
\]

Add those ordinary rules and you have the vector rule (11) for the divergence of \( uV \).

Integrating the two parts of \( \text{div}(uV) \) gives \( \int uV \cdot n \, dS \) by the Divergence Theorem. Then one part moves to the other side, producing the minus signs in (12) and (13). Integration by parts leaves a boundary term, in three and two dimensions as it did in one dimension: \( \int uv' \, dx = - \int u'v \, dx + [uv]_a^b \).

EXAMPLE 5 Find the divergence of \( F = R/p^2 \), starting from \( \text{grad} \, \rho = R/\rho \).

Solution Take \( V = R \) and \( u = 1/p^2 \) in the product rule (11). Then \( \text{div} \, F = (\text{div} \, R)/p^2 + R \cdot (\text{grad} \, 1/p^2) \). The divergence of \( R = xi + yj + zk \) is 3. For \( 1/p^2 \) apply the chain rule:

\[
R \cdot (\text{grad} \, 1/p^2) = -2R \cdot (\text{grad} \, \rho)/\rho^3 = -2R \cdot R/\rho^4 = -2/\rho^2.
\]

The two parts of \( \text{div} \, F \) combine into \( 3/\rho^2 - 2/\rho^2 = 1/\rho^2 \)—as claimed in Example 4.

EXAMPLE 6 Find the balance equation for flow with velocity \( V \) and fluid density \( \rho \).

\( V \) is the rate of movement of fluid, while \( \rho V \) is the rate of movement of mass. Comparing the ocean to the atmosphere shows the difference. Air has a greater velocity than water, but a much lower density. So normally \( F = \rho V \) is larger for the ocean. (Don't confuse the density \( \rho \) with the radial distance \( \rho \). The Greeks only used 24 letters.)

There is another difference between water and air. Water is virtually incompressible (meaning \( \rho = \text{constant} \)). Air is certainly compressible (its density varies). The balance equation is a fundamental law—the conservation of mass or the "continuity equation" for fluids. This is a mathematical statement about a physical flow without sources or sinks:

\[
\text{Continuity Equation:} \quad \text{div}(\rho V) + \frac{\partial \rho}{\partial t} = 0.
\] (14)
**15 Vector Calculus**

Explanation: The mass in a region is \( \iiint \rho \, dV \). Its rate of decrease is \( -\iiint \partial \rho / \partial t \, dV \). The decrease comes from flow out through the surface (normal vector \( n \)). The dot product \( F \cdot n = \rho V \cdot n \) is the rate of mass flow through the surface. So the integral \( \iint F \cdot ndS \) is the total rate that mass goes out. By the Divergence Theorem this is \( \iiint \text{div} \, F \, dV \).

To balance \( -\iiint \partial \rho / \partial t \, dV \) in every region, \( \text{div} \, F \) must equal \( -\partial \rho / \partial t \) at every point.

The figure shows this continuity equation (14) for flow in the \( x \) direction.

\[
\text{mass in} \quad \rho V \, ds \, dt \quad \to \quad \boxed{\text{mass } \rho ds \, dx} \quad \text{extra mass out} \quad \text{mass loss} \quad d(\rho V) \, ds \, dt = -d \rho \, ds \, dx
\]

Fig. 15.22 Conservation of mass during time \( dt \): \( d(\rho V) \, ds \, dx + d\rho / dt \, dt = 0 \).

### 15.5 Exercises

#### Read-through questions

In words, the basic balance law is **0**. The flux of \( F \) through a closed surface \( S \) is the double integral **0**. The divergence of \( M \) + \( N \) + \( P \) is **0**, and it measures **0**. The total source is the triple integral **0**. That equals the flux by the **0** Theorem.

For \( F = 5 \hat{k} \) the divergence is **0**. If \( V \) is a cube of side \( a \) then the triple integral equals **0**. The top surface where \( z=a \) has \( n = \hat{k} \) and \( F \cdot n = -a \). The bottom and sides have \( F \cdot n = 0 \). The integral \( \iint F \cdot ndS \) equals **0**.

The field \( F = R, \rho^3 \) has \( \text{div} \, F = 0 \) except **0**. \( \iint F \cdot ndS \) equals **0** over any surface around the origin. This illustrates Gauss's Law **0**. The field \( F = x \hat{i} + y \hat{j} - 2 \hat{k} \) has \( \text{div} \, F = -y \) and \( \iint F \cdot ndS = 0 \). For this \( F \), the flux out through a pyramid and in through its base are **0**.

The symbol \( V \) stands for **0**. In this notation \( \text{div} \, F \) is **0**. The gradient of \( f \) is **0**. The divergence of \( g \) is **0**. The equation \( \text{div} \, \text{grad} \, f = 0 \) is **0**'s equation.

The divergence of a product is \( \text{div}(uv) = u \text{div} \, v + v \text{div} \, u \). Integration by parts is \( \iint u \text{div} \, V \, dx \, dy \, dz = \gamma + \zeta \). In two dimensions this becomes **0**. In one dimension it becomes **0**. For steady fluid flow the continuity equation is \( \text{div} \, \rho V = 0 \).

In 1-10 compute the flux \( \iint F \cdot ndS \) by the Divergence Theorem.

1. \( F = x \hat{i} + y \hat{j} + z \hat{k} \); \( S \): unit sphere \( x^2 + y^2 + z^2 = 1 \).

2. \( F = -y \hat{i} + x \hat{j} \); \( V \): unit cube \( 0 \leq x \leq 1, 0 \leq y \leq 1, 0 \leq z \leq 1 \).

3. \( F = x \hat{i} + y \hat{j} + z \hat{k} \); \( S \): unit sphere.

4. \( F = x \hat{i} + 2y \hat{j} + z \hat{k} \); \( V \): unit cube.

5. \( F = x \hat{i} - 2y \hat{j} \); \( S \): sides \( x = 0, 0 \leq y \leq 1, 0 \leq z \leq 1 \).

6. \( F = u \), \( (x \hat{i} + y \hat{j} + z \hat{k}) \rho \); \( S \): sphere \( \rho = a \).

7. \( F = \rho (x \hat{i} + y \hat{j} + z \hat{k}) \); \( S \): sphere \( \rho = a \).

8. \( F = x \hat{i} + y \hat{j} + z \hat{k} \); \( S \): sphere \( \rho = a \).

9. \( F = z \hat{k} \); \( V \): upper half of ball \( \rho < a \).

10. \( F = \text{grad} \, (x^2 \sin \theta) \); \( S \): sphere \( \rho < a \).

11. Find \( \iint \text{div}(x^2 \hat{i} + y \hat{j} + 2 \hat{k}) \, dV \) in the cube \( 0 \leq x, y, z \leq a \).

Also compute \( n \) and \( \iint F \cdot ndS \) for all six faces and add.

12. When \( a \) is small in problem 11, the answer is close to \( ca^3 \). Find the number \( c \). At what point does \( \text{div} \, F = c? \)

(a) Integrate the divergence of \( F = \rho i \) in the ball \( \rho < a \).

(b) Compute \( \iint \text{F} \cdot ndS \) over the spherical surface \( \rho = a \).

13. Evaluate \( \iint F \cdot ndS \) when \( F = x \hat{i} + y \hat{j} + z \hat{k} \) and:

(a) \( S \) is the cone \( z^2 = x^2 + y^2 \) bounded above by the plane \( z = 1 \).

(b) \( S \) is the pyramid with corners \( (0, 0, 0), (1, 0, 0), (0, 1, 0), (0, 0, 1) \).

14. Compute all integrals in the Divergence Theorem when \( F = x \hat{i} + y \hat{j} - k \hat{k} \).

15. Following Example 5, compute the divergence of \( (x \hat{i} + y \hat{j} - z \hat{k}) \rho^2 \).

16. Give an example where \( \iint F \cdot ndS \) is easier than \( \iint \text{div} \, F \, dV \).
15.6 Stokes' Theorem and the Curl of F

21 Suppose \( \mathbf{F} = M(x, y)i + N(x, y)j \), \( R \) is a region in the \( xy \) plane, and \( (x, y, z) \) is in \( V \) if \( (x, y) \) is in \( R \) and \( |z| \leq 1 \).
(a) Describe \( V \) and reduce \( \iiint \text{div} \mathbf{F} \, dV \) to a double integral.
(b) Reduce \( \iint \mathbf{F} \cdot n \, dS \) to a line integral (check top, bottom, side).
(c) Whose theorem says that the double integral equals the line integral?

22 Is it possible to have \( \mathbf{F} \cdot \mathbf{n} = 0 \) at all points of \( S \) and also \( \text{div} \mathbf{F} = 0 \) at all points in \( V \)? \( \mathbf{F} = 0 \) is not allowed.

23 Inside a solid ball (radius \( a \), density \( 1 \), mass \( M = 4\pi a^3 / 3 \)) the gravity field is \( \mathbf{F} = -GMR/a^3 \).
(a) Check \( \text{div} \mathbf{F} = -4\pi G \) in Gauss's Law.
(b) The force at the surface is the same as if the whole mass \( M \) were ________.
(c) Find a gradient field with \( \text{div} \mathbf{F} = 6 \) in the ball \( \rho \leq a \) and \( \text{div} \mathbf{F} = 0 \) outside.

24 The outward field \( \mathbf{F} = \mathbf{R}/\rho^3 \) has magnitude \( |\mathbf{F}| = 1/\rho^2 \). Through an area \( A \) on a sphere of radius \( \rho \), the flux is _________. A spherical box has faces at \( \rho_1 \) and \( \rho_2 \) with \( A = \rho_1^2 \sin \phi \, d\phi \, d\theta \) and \( A = \rho_2^2 \sin \phi \, d\phi \, d\theta \). Deduce that the flux out of the box is zero, which confirms \( \text{div} \mathbf{F} = 0 \).

25 In Gauss's Law, what charge distribution \( q(x, y, z) \) gives the unit field \( \mathbf{E} = \mathbf{u} \)? What is the flux through the unit sphere?

26 If a fluid with velocity \( \mathbf{V} \) is incompressible (constant density \( \rho \)), then its continuity equation reduces to _______. If it is irrotational then \( \mathbf{F} = \text{grad} f \). If it is both then \( f \) satisfies _______ equation.

27 True or false, with a good reason.
(a) If \( \iint \mathbf{F} \cdot n \, dS = 0 \) for every closed surface, \( \mathbf{F} \) is constant.
(b) If \( \mathbf{F} = \text{grad} f \) then \( \text{div} \mathbf{F} = 0 \).
(c) If \( |\mathbf{F}| \leq 1 \) at all points then \( \iiint |\text{div} \mathbf{F}| \, dV \leq \text{area of the surface} \).
(d) If \( |\mathbf{F}| \leq 1 \) at all points then \( |\text{div} \mathbf{F}| \leq 1 \) at all points.

28 Write down statements \( E-F-G-H \) for source-free fields \( \mathbf{F}(x, y, z) \) in three dimensions. In statement \( F \), paths sharing the same endpoint become surfaces sharing the same boundary curve. In \( G \), the stream function becomes a vector field such that \( \mathbf{F} = \text{curl} \mathbf{g} \).

29 Describe two different surfaces bounded by the circle \( x^2 + y^2 = 1, z = 0 \). The field \( \mathbf{F} \) automatically has the same flux through both if ________.

30 The boundary of a bounded region \( R \) has no boundary. Draw a plane region and explain what that means. What does it mean for a solid ball?

---

**15.6 Stokes' Theorem and the Curl of F**

For the Divergence Theorem, the surface was closed. \( S \) was the boundary of \( V \). Now the surface is not closed and \( S \) has its own boundary—a curve called \( C \). We are back near the original setting for Green's Theorem (region bounded by curve, double integral equal to work integral). But Stokes' Theorem, also called Stokes's Theorem, is in three-dimensional space. There is a curved surface \( S \) bounded by a space curve \( C \). This is our first integral around a space curve.

The move to three dimensions brings a change in the vector field. The plane field \( \mathbf{F}(x, y) = Mi + Nj \) becomes a space field \( \mathbf{F}(x, y, z) = Mi + Nj + Pk \). The work \( Mdx + Ndy \) now includes \( Pdz \). The critical quantity in the double integral (it was \( \partial N/\partial x - \partial M/\partial y \) must change too. We called this scalar quantity "curl \( \mathbf{F} \)," but in reality it is only the third component of a vector. Stokes' Theorem needs all three components of that vector—which is curl \( \mathbf{F} \).

**DEFINITION** The curl of a vector field \( \mathbf{F}(x, y, z) = Mi + Nj + Pk \) is the vector field

\[
\text{curl} \ \mathbf{F} = \left( \frac{\partial P}{\partial y} - \frac{\partial N}{\partial z} \right) \mathbf{i} + \left( \frac{\partial M}{\partial z} - \frac{\partial P}{\partial x} \right) \mathbf{j} + \left( \frac{\partial N}{\partial x} - \frac{\partial M}{\partial y} \right) \mathbf{k}.
\]
The symbol \( \nabla \times \mathbf{F} \) stands for a determinant that yields those six derivatives:

\[
\text{curl } \mathbf{F} = \nabla \times \mathbf{F} = \begin{vmatrix}
\mathbf{i} & \mathbf{j} & \mathbf{k} \\
\frac{\partial}{\partial x} & \frac{\partial}{\partial y} & \frac{\partial}{\partial z} \\
M & N & P
\end{vmatrix}.
\]

The three products \( i \frac{\partial}{\partial y} P \) and \( j \frac{\partial}{\partial z} M \) and \( k \frac{\partial}{\partial x} N \) have plus signs. The three products like \( k \frac{\partial}{\partial y} M \), down to the left, have minus signs. There is a cyclic symmetry. This determinant helps the memory, even if it looks and is illegal. A determinant is not supposed to have a row of vectors, a row of operators, and a row of functions.

**EXAMPLE 1** The plane field \( M(x, y)i + N(x, y)j \) has \( P = 0 \) and \( \partial M/\partial z = 0 \) and \( \partial N/\partial z = 0 \). Only two terms survive: curl \( \mathbf{F} = (\partial N/\partial x - \partial M/\partial y)k \). Back to Green.

**EXAMPLE 2** The cross product \( \mathbf{a} \times \mathbf{R} \) is a spin field \( \mathbf{S} \). Its axis is the fixed vector \( \mathbf{a} = a_1 \mathbf{i} + a_2 \mathbf{j} + a_3 \mathbf{k} \). The flow in Figure 15.23 turns around \( \mathbf{a} \), and its components are

\[
\mathbf{S} = \mathbf{a} \times \mathbf{R} = \begin{vmatrix}
\mathbf{i} & \mathbf{j} & \mathbf{k} \\
a_1 & a_2 & a_3 \\
x & y & z
\end{vmatrix} = (a_2 z - a_3 y) \mathbf{i} + (a_3 x - a_1 z) \mathbf{j} + (a_1 y - a_2 x) \mathbf{k}.
\]

Our favorite spin field \(-yi + xj\) has \((a_1, a_2, a_3) = (0, 0, 1)\) and its axis is \( \mathbf{a} = \mathbf{k} \).

The divergence of a spin field is \( M_x + N_y + P_z = 0 + 0 + 0 \). Note how the divergence uses \( M_x \) while the curl uses \( N_y \) and \( P_z \). **The curl of \( \mathbf{S} \) is the vector \( 2\mathbf{a} \):**

\[
\left( \frac{\partial P}{\partial y} - \frac{\partial N}{\partial z} \right) \mathbf{i} + \left( \frac{\partial M}{\partial z} - \frac{\partial P}{\partial x} \right) \mathbf{j} + \left( \frac{\partial N}{\partial x} - \frac{\partial M}{\partial y} \right) \mathbf{k} = 2a_1 \mathbf{i} + 2a_2 \mathbf{j} + 2a_3 \mathbf{k} = 2\mathbf{a}.
\]

This example begins to reveal the meaning of the curl. It measures the spin! The direction of curl \( \mathbf{F} \) is the \textit{axis of rotation}—in this case along \( \mathbf{a} \). The magnitude of curl \( \mathbf{F} \) is twice the speed of rotation. In this case \( |\text{curl } \mathbf{F}| = 2|\mathbf{a}| \) and the angular velocity is \( |\mathbf{a}| \).

**EXAMPLE 3** (II) \textit{Every gradient field} \( \mathbf{F} = \partial f/\partial x \mathbf{i} + \partial f/\partial y \mathbf{j} + \partial f/\partial z \mathbf{k} \) \textit{has curl} \( \mathbf{F} = 0 \):

\[
\text{curl } \mathbf{F} = \left( \frac{\partial}{\partial y} \frac{\partial f}{\partial z} - \frac{\partial}{\partial z} \frac{\partial f}{\partial y} \right) \mathbf{i} + \left( \frac{\partial}{\partial z} \frac{\partial f}{\partial x} - \frac{\partial}{\partial x} \frac{\partial f}{\partial z} \right) \mathbf{j} + \left( \frac{\partial}{\partial x} \frac{\partial f}{\partial y} - \frac{\partial}{\partial y} \frac{\partial f}{\partial x} \right) \mathbf{k} = 0.
\]

Always \( f_{yz} = f_{zy} \). They cancel. Also \( f_{xz} = f_{zx} \) and \( f_{yx} = f_{xy} \). So curl grad \( f = 0 \).
15.6 Stokes' Theorem and the Curl of $F$

**EXAMPLE 4** (twin of Example 3) The divergence of curl $F$ is also automatically zero:

$$\text{div} \, \text{curl} \, F = \frac{\partial}{\partial x} \left( \frac{\partial P}{\partial y} - \frac{\partial N}{\partial z} \right) + \frac{\partial}{\partial y} \left( \frac{\partial M}{\partial z} - \frac{\partial P}{\partial x} \right) + \frac{\partial}{\partial z} \left( \frac{\partial N}{\partial x} - \frac{\partial M}{\partial y} \right) = 0. \quad (5)$$

Again the mixed derivatives give $P_{xy} = P_{yx}$ and $N_{xz} = N_{zx}$ and $M_{zy} = M_{yz}$. The terms cancel in pairs. In "curl grad" and "div curl", everything is arranged to give zero.

---

**THE MEANING OF CURL $F$**

Example 5 put a paddlewheel into the flow. This is possible for any vector field $F$, and it gives insight into curl $F$. The turning of the wheel (if it turns) depends on its location $(x, y, z)$. The turning also depends on the orientation of the wheel. We could put it into a spin field, and if the wheel axis $n$ is perpendicular to the spin axis $a$, the wheel won’t turn! The general rule for turning speed is this: **the angular velocity of the wheel is $\frac{1}{2} \text{curl} \, F \cdot n$.** This is the "directional spin," just as $(\text{grad} \, f) \cdot u$ was the "directional derivative"—and $n$ is a unit vector like $u$.

There is no spin anywhere in a gradient field. It is **irrotational**: $\text{curl} \, \text{grad} \, f = 0$.

The pure spin field $a \times R$ has curl $F = 2a$. The angular velocity is $a \cdot n$ (note that $\frac{1}{2}$ cancels 2). This turning is everywhere, **not just at the origin**. If you put a penny on a compact disk, it turns once when the disk rotates once. That spin is "around itself," and it is the same whether the penny is at the center or not.

The turning speed is greatest when the wheel axis $n$ lines up with the spin axis $a$. Then $a \cdot n$ is the full length $|a|$. The gradient gives the direction of fastest growth, and the curl gives the direction of fastest turning:

- maximum growth rate of $f$ is $|\text{grad} \, f|$ in the direction of $\text{grad} \, f$
- maximum rotation rate of $F$ is $\frac{1}{2} |\text{curl} \, F|$ in the direction of $\text{curl} \, F$.

---

**THE MEANING OF CURL $F$**

The spin field $S$ has no divergence. The position field $R$ has no curl. $R$ is the gradient of $f = \frac{1}{2}(x^2 + y^2 + z^2)$. $S$ is the curl of a suitable $F$. Then $\text{div} \, S = \text{div} \, \text{curl} \, F$ and $\text{curl} \, R = \text{curl} \, \text{grad} \, f$ are automatically zero.

You correctly believe that curl $F$ measures the "spin" of the field. You may expect that curl $(F + G)$ is curl $F + \text{curl} \, G$. Also correct. Finally you may think that a field of parallel vectors has no spin. That is wrong. Example 5 has parallel vectors, but their different lengths produce spin.

---

**EXAMPLE 5** The field $V = zi$ in the $x$ direction has curl $V = j$ in the $y$ direction.

If you put a wheel in the $xz$ plane, **this field will turn it**. The velocity $zi$ at the top of the wheel is greater than $zi$ at the bottom (Figure 15.23c). So the top goes faster and the wheel rotates. The axis of rotation is curl $V = j$. The turning speed is $\frac{1}{2}$, because this curl has magnitude 1.

Another velocity field $v = -xk$ produces the same spin: curl $v = j$. The flow is in the $z$ direction, it varies in the $x$ direction, and the spin is in the $y$ direction. Also interesting is $V + v$. The two "shear fields" add to a perfect spin field $S = zi - xk$, whose curl is $2j$. 

---
Finally we come to the big theorem. It will be like Green’s Theorem—a line integral equals a surface integral. The line integral is still the work $\int \mathbf{F} \cdot d\mathbf{R}$ around a curve. The surface integral in Green’s Theorem is $\iint (N_x - M_y) \, dx \, dy$. The surface is flat (in the $xy$ plane). Its normal direction is $\mathbf{k}$, and we now recognize $N_x - M_y$ as the $k$ component of the curl. Green’s Theorem uses only this component because the normal direction is always $\mathbf{k}$. For Stokes’ Theorem on a curved surface, we need all three components of curl $\mathbf{F}$.

Figure 15.24 shows a hat-shaped surface $S$ and its boundary $C$ (a closed curve). Walking in the positive direction around $C$, with your head pointing in the direction of $\mathbf{n}$, the surface is on your left. You may be standing straight up ($\mathbf{n} = k$ in Green’s Theorem). You may even be upside down ($\mathbf{n} = -k$ is allowed). In that case you must go the other way around $C$, to keep the two sides of equation (6) equal. The surface is still on the left. A Möbius strip is not allowed, because its normal direction cannot be established. The unit vector $\mathbf{n}$ determines the “counterclockwise” direction along $C$.

\[ \oint_C \mathbf{F} \cdot d\mathbf{R} = \iint_S (\text{curl} \, \mathbf{F}) \cdot \mathbf{n} \, dS. \]  

(Stokes’ Theorem) (6)

The right side adds up small spins in the surface. The left side is the total circulation (or work) around $C$. That is not easy to visualize—this may be the hardest theorem in the book—but notice one simple conclusion. If $\text{curl} \, \mathbf{F} = 0$ then $\oint_C \mathbf{F} \cdot d\mathbf{R} = 0$. This applies above all to gradient fields—as we know.

A gradient field has no curl, by (4). A gradient field does no work, by (6). In three dimensions as in two dimensions, gradient fields are conservative fields. They will be the focus of this section, after we outline a proof (or two proofs) of Stokes’ Theorem.

The first proof shows why the theorem is true. The second proof shows that it really is true (and how to compute). You may prefer the first.

First proof  Figure 15.24 has a triangle $ABC$ attached to a triangle $ACD$. Later there can be more triangles. $S$ will be piecewise flat, close to a curved surface. Two triangles are enough to make the point. In the plane of each triangle (they have different $\mathbf{n}$’s) Green’s Theorem is known:

\[ \oint_{AB+BC+CA} \mathbf{F} \cdot d\mathbf{R} = \iint_{ABC} \text{curl} \, \mathbf{F} \cdot d\mathbf{S} \quad \oint_{AC+CD+DA} \mathbf{F} \cdot d\mathbf{R} = \iint_{ACD} \text{curl} \, \mathbf{F} \cdot d\mathbf{S}. \]

Now add. The right sides give $\iint \text{curl} \, \mathbf{F} \cdot d\mathbf{S}$ over the two triangles. On the left, the integral over $CA$ cancels the integral over $AC$. The “crosscut” disappears. That leaves $AB + BC + CD + DA$. This line integral goes around the outer boundary $C$—which is the left side of Stokes’ Theorem.

Fig. 15.24 Surfaces $S$ and boundary curves $C$. Change in $\mathbf{B} \rightarrow \text{curl} \, \mathbf{E} \rightarrow$ current in $C$. 

\[ \]
15.6 Stokes' Theorem and the Curl of $F$

**Second proof** Now the surface can be curved. A new proof may seem excessive, but it brings formulas you could compute with. From $z = f(x, y)$ we have

$$dz = \frac{\partial f}{\partial x} \, dx + \frac{\partial f}{\partial y} \, dy$$

and

$$ndS = (-\frac{\partial f}{\partial x} \mathbf{i} - \frac{\partial f}{\partial y} \mathbf{j} + k) \, dx \, dy.$$  

For $ndS$, see equation 15.4.11. With this $dz$, the line integral in Stokes' Theorem is

$$\oint_{C} F \cdot d\mathbf{R} = \int_{\text{shadow of } C} M \, dx + N \, dy + P(\frac{\partial f}{\partial x} \, dx + \frac{\partial f}{\partial y} \, dy).$$

(7)

The dot product of $\text{curl } F$ and $ndS$ gives the surface integral $\iint_{S} \text{curl } F \cdot ndS$:

$$\iint_{S} \left[ (P_x - N_y)(-\frac{\partial f}{\partial x}) + (M_z - P_z)(-\frac{\partial f}{\partial y}) + (N_x - M_x) \right] \, dx \, dy.$$  

(8)

To prove $(7) = (8)$, change $M$ in Green's Theorem to $M + P\frac{\partial f}{\partial x}$. Also change $N$ to $N + P\frac{\partial f}{\partial y}$. Then $(7) = (8)$ is Green's Theorem down on the shadow (Problem 47). This proves Stokes' Theorem up on $S$. Notice how Green's Theorem (flat surface) was the key to both proofs of Stokes' Theorem (curved surface).

**EXAMPLE 6** Stokes' Theorem in electricity and magnetism yields Faraday's Law.

Stokes' Theorem is not heavily used for calculations—equation (8) shows why. But the spin or curl or vorticity of a flow is absolutely basic in fluid mechanics. The other important application, coming now, is to electric fields. Faraday's Law is to Gauss's Law as Stokes' Theorem is to the Divergence Theorem.

Suppose the curve $C$ is an actual wire. We can produce current along $C$ by varying the magnetic field $B(t)$. The flux $\varphi = \iint B \cdot ndS$, passing within $C$ and changing in time, creates an electric field $E$ that does work:

$$\text{Faraday's Law (integral form): work} = \oint_{C} E \cdot d\mathbf{R} = -d\varphi/dt.$$

That is physics. It may be true, it may be an approximation. Now comes mathematics (surely true), which turns this integral form into a differential equation. Information at points is more convenient than information around curves. Stokes converts the line integral of $E$ into the surface integral of $\text{curl } E$:

$$\oint_{C} E \cdot d\mathbf{R} = \iint_{S} \text{curl } E \cdot ndS$$

and also

$$-d\varphi/dt = \iint_{S} - (\partial B/\partial t) \cdot ndS.$$  

These are equal for any curve $C$, however small. So the right sides are equal for any surface $S$. We squeeze to a point. The right hand sides give one of Maxwell's equations:

$$\text{Faraday's Law (differential form): curl } E = -\partial B/\partial t.$$

**CONSERVATIVE FIELDS AND POTENTIAL FUNCTIONS**

The chapter ends with our constant and important question: Which fields do no work around closed curves? Remember test $D$ for plane curves and plane vector fields:

if $\partial M/\partial y = \partial N/\partial x$ then $F$ is conservative and $F = \text{grad } f$ and $\oint F \cdot d\mathbf{R} = 0$.

Now allow a three-dimensional field like $F = 2xy \mathbf{i} + (x^2 + z) \mathbf{j} + y \mathbf{k}$. Does it do work around a space curve? Or is it a gradient field? That will require $\partial f/\partial x = 2xy$ and $\partial f/\partial y = x^2 + z$ and $\partial f/\partial z = y$. We have three equations for one function $f(x, y, z)$. Normally they can't be solved. When test $D$ is passed (now it is the three-dimensional test: curl $F = 0$) they can be solved. This example passes test $D$, and $f$ is $x^2y + yz$. 
A detailed proof of A $\Rightarrow$ B $\Rightarrow$ C $\Rightarrow$ D $\Rightarrow$ A is not needed. Only notice how C $\Rightarrow$ D: curl grad F is always zero. The newest part is D $\Rightarrow$ A. If curl F = 0 then $\oint F \cdot dR = 0$. But that is not news. It is Stokes' Theorem.

The interesting problem is to solve the three equations for $f$, when test D is passed. The example above had
\[ \frac{\partial f}{\partial x} = 2xy \Rightarrow f = \int 2xy \, dx = x^2y + \text{any function } C(y, z) \]
\[ \frac{\partial f}{\partial y} = x^2 + z = x^2 + \partial C/\partial y \Rightarrow C = yz + \text{any function } C(z) \]
\[ \frac{\partial f}{\partial z} = y = y + dc/dz = C(z) \text{ can be zero.} \]

The first step leaves an arbitrary $C(y, z)$ to fix the second step. The second step leaves an arbitrary $C(z)$ to fix the third step (not needed here). Assembling the three steps, $f = x^2y + C = x^2y + yz + c = x^2y + yz$. Please recognize that the "fix-up" is only possible when curl F = 0. Test D must be passed.

EXAMPLE 7 Is $F = (z - y)i + (x - z)j + (y - x)k$ the gradient of any $f$?

Test D says no. This F is a spin field $\mathbf{a} \times \mathbf{R}$. Its curl is $2\mathbf{a} = (2, 2, 2)$, which is not zero. A search for $f$ is bound to fail, but we can try. To match $\partial f/\partial x = z - y$, we must have $f = xz - yx + C(y, z)$. The $y$ derivative is $-x + \partial C/\partial y$. That never matches $N = x - z$, so $f$ can't exist.

EXAMPLE 8 What choice of $P$ makes $F = yz^2i + xz^2j + Pk$ conservative? Find $f$.

Solution We need curl F = 0, by test D. First check $\partial M/\partial y = z^2 = \partial N/\partial x$. Also
\[ \partial P/\partial x = \partial M/\partial z = 2yz \quad \text{and} \quad \partial P/\partial y = \partial N/\partial z = 2xz. \]

$P = 2xyz$ passes all tests. To find $f$ we can solve the three equations, or notice that $f = xyz^2$ is successful. Its gradient is $F$.

A third method defines $f(x, y, z)$ as the work to reach $(x, y, z)$ from $(0, 0, 0)$. The path doesn't matter. For practice we integrate $F \cdot d\mathbf{R} = M \, dx + N \, dy + P \, dz$ along the straight line $(xt, yt, zt)$:
\[ f(x, y, z) = \int_0^1 \left( (yt)(zt^2)(x \, dt) + (xt)(zt^2)(y \, dt) + 2(xt)(yt)(zt)(z \, dt) \right) = xyz^2. \]

EXAMPLE 9 Why is div curl grad $f$ automatically zero (in two ways)?

Solution First, curl grad $f$ is zero (always). Second, div curl F is zero (always). Those are the key identities of vector calculus. We end with a review.

\textit{Green's Theorem:} \[ \oint F \cdot d\mathbf{R} = \iint (\partial N/\partial x - \partial M/\partial y) \, dx \, dy \]
\[ \oint F \cdot n \, ds = \iint (\partial M/\partial x + \partial N/\partial y) \, dx \, dy \]
15.6 Stokes' Theorem and the Curl of \( \mathbf{F} \)

**Divergence Theorem:** \[ \iiint \mathbf{F} \cdot d\mathbf{S} = \iiint (\partial M/\partial x + \partial N/\partial y + \partial P/\partial z) dx \, dy \, dz \]

**Stokes' Theorem:** \[ \int C \mathbf{F} \cdot d\mathbf{R} = \int \text{curl } \mathbf{F} \cdot d\mathbf{S}. \]

The first form of Green's Theorem leads to Stokes' Theorem. The second form becomes the Divergence Theorem. You may ask, why not go to three dimensions in the first place? The last two theorems contain the first two (take \( P = 0 \) and a flat surface). We could have reduced this chapter to two theorems, not four. I admit that, but a fundamental principle is involved: "It is easier to generalize than to specialize."

For the same reason, \( \frac{df}{dx} \) came before partial derivatives and the gradient.

15.6 **EXERCISES**

**Read-through questions**

The curl of \( \mathbf{M} + \mathbf{N} + \mathbf{P} \) is the vector \( \mathbf{\omega} \). It equals the 3 by 3 determinant \( \mathbf{b} \). The curl of \( x^2 \mathbf{i} + x^2 \mathbf{k} \) is \( \mathbf{c} \).

For \( \mathbf{S} = y \mathbf{i} - (x + 2y) \mathbf{j} + y \mathbf{k} \) the curl is \( \mathbf{d} \). This \( \mathbf{S} \) is a \( \mathbf{e} \) field \( \mathbf{a} \times \mathbf{R} = \frac{1}{2} (\text{curl } \mathbf{F}) \times \mathbf{R} \), with axis vector \( \mathbf{a} = \mathbf{f} \). For any gradient field \( f_x \mathbf{i} + f_y \mathbf{j} + f_z \mathbf{k} \) the curl is \( \mathbf{g} \).

That is the important identity \( \text{curl } \text{grad } f = \mathbf{h} \). It is based on \( f_{xx} = f_{yy} \) and \( i \) and \( j \).

The curl measures the \( \mathbf{m} \) of a vector field. A paddlewheel in the field with its axis along \( \mathbf{n} \) has turning speed \( \mathbf{m} \).

**Problems**

1-6 find curl \( \mathbf{F} \).

1. \( \mathbf{F} = z \mathbf{i} + x \mathbf{j} + y \mathbf{k} \)
2. \( \mathbf{F} = \text{grad}(xe^z \sin z) \)
3. \( \mathbf{F} = (x + y + z)(\mathbf{i} + \mathbf{j} + \mathbf{k}) \)
4. \( \mathbf{F} = (x + y) \mathbf{i} - (x + y) \mathbf{k} \)
5. \( \mathbf{F} = \rho^2 (\mathbf{x} \mathbf{i} + \mathbf{y} \mathbf{j} + \mathbf{z} \mathbf{k}) \)
6. \( \mathbf{F} = (\mathbf{i} + \mathbf{j}) \times \mathbf{R} \)
7. Find a potential \( f \) for the field in Problem 3.
8. Find a potential \( f \) for the field in Problem 5.
9. When do the fields \( x^2 \mathbf{i} \) and \( x^2 \mathbf{j} \) have zero curl?
10. When does \( a_1 \mathbf{x} + a_2 \mathbf{y} + a_3 \mathbf{z} \mathbf{k} \) have zero curl?

In 11-14, compute \( \text{curl } \mathbf{F} \) and find \( \oint_C \mathbf{F} \cdot d\mathbf{R} \) by Stokes' Theorem.

11. \( \mathbf{F} = x^2 \mathbf{i} + y^2 \mathbf{k} \), \( C \) is circle \( x^2 + z^2 = 1 \), \( y = 0 \).
12. \( \mathbf{F} = \mathbf{i} \times \mathbf{R} \), \( C \) is circle \( x^2 + z^2 = 1 \), \( y = 0 \).
13. \( \mathbf{F} = (\mathbf{i} + \mathbf{j}) \times \mathbf{R} \), \( C \) is circle \( y^2 + z^2 = 1 \), \( x = 0 \).
14. \( \mathbf{F} = (y + x) \mathbf{j} \), \( C \) is circle \( x^2 + y^2 = 1 \), \( z = 0 \).

15 (important) Suppose two surfaces \( S \) and \( T \) have the same boundary \( C \), and the direction around \( C \) is the same.

(a) Prove \( \iiint_S \text{curl } \mathbf{F} \cdot d\mathbf{S} = \iiint_T \text{curl } \mathbf{F} \cdot d\mathbf{S} \).
(b) Second proof: The difference between those integrals is \( \iiint \text{div} (\text{curl } \mathbf{F}) dV \). By what Theorem? What region is \( V \)? Why is this integral zero?

16 In 15, suppose \( S \) is the top half of the earth (\( \mathbf{n} \) goes out) and \( T \) is the bottom half (\( \mathbf{n} \) comes in). What are \( C \) and \( \partial V \)?

17 Explain why \( \iiint \text{curl } \mathbf{F} \cdot d\mathbf{S} = 0 \) over the closed boundary of any solid \( V \).

18 Suppose \( \text{curl } \mathbf{F} = \mathbf{0} \) and \( \text{div } \mathbf{F} = \mathbf{0} \). (a) Why is \( \mathbf{F} \) the gradient of a potential? (b) Why does the potential satisfy Laplace's equation \( f_{xx} + f_{yy} + f_{zz} = 0 \)?

In 19-22, find a potential \( f \) if it exists.

19. \( \mathbf{F} = \mathbf{i} + \mathbf{j} + \mathbf{x} \mathbf{k} \)
20. \( \mathbf{F} = 2xy \mathbf{i} + x^2 \mathbf{j} + x^4 \mathbf{y} \mathbf{k} \)
21. \( \mathbf{F} = e^x \mathbf{i} - e^{-x} \mathbf{k} \)
22. \( \mathbf{F} = y \mathbf{i} + x \mathbf{j} + (xy + z^2) \mathbf{k} \)
23. Find a field with \( \text{curl } \mathbf{F} = (1, 0, 0) \).
24. Find all fields with \( \text{curl } \mathbf{F} = (1, 0, 0) \).
25. \( \mathbf{S} = \mathbf{a} \times \mathbf{R} \) is a spin field. Compute \( \mathbf{F} = \mathbf{b} \times \mathbf{S} \) (constant vector \( \mathbf{b} \)) and find its curl.
26 How fast is a paddlewheel turned by the field \( \mathbf{F} = yi - xk \) (a) if its axis direction is \( \mathbf{n} = j \)? (b) if its axis is lined up with \( \text{curl } \mathbf{F} \)? (c) if its axis is perpendicular to \( \text{curl } \mathbf{F} \)?

27 How is \( \text{curl } \mathbf{F} \) related to the angular velocity \( \omega \) in the spin field \( \mathbf{F} = \rho(-yi + xj) \)? How fast does a wheel spin, if it is in the plane \( x + y + z = 1 \)?

28 Find a vector field \( \mathbf{F} \) whose curl is \( \mathbf{S} = yi - xj \).

29 Find a vector field \( \mathbf{F} \) whose curl is \( \mathbf{S} = a \times \mathbf{R} \).

30 True or false: when two vector fields have the same curl at all points: (a) their difference is a constant field (b) their difference is a gradient field (c) they have the same divergence.

In 31–34, compute \( \iiint \text{curl } \mathbf{F} \cdot n \, dS \) over the top half of the sphere \( x^2 + y^2 + z^2 = 1 \) and (separately) \( \int \int \int \mathbf{F} \cdot d\mathbf{R} \) around the equator.

31 \( \mathbf{F} = yi - xj \)

32 \( \mathbf{F} = R/r^2 \)

33 \( \mathbf{F} = a \times \mathbf{R} \)

34 \( \mathbf{F} = (a \times \mathbf{R}) \times \mathbf{R} \)

35 The circle \( C \) in the plane \( x + y + z = 6 \) has radius \( r \) and center at \((1, 2, 3)\). The field \( \mathbf{F} \) is \( 3zj + 2y k \). Compute \( \int \int \int \mathbf{F} \cdot d\mathbf{R} \) around \( C \).

36 \( \mathbf{S} \) is the top half of the unit sphere and \( \mathbf{F} = zi + xj + x y z k \). Find \( \iiint \text{curl } \mathbf{F} \cdot n \, dS \).

37 Find \( \mathbf{g}(x, y) \) so that \( \text{curl } \mathbf{g}k = yi + x^2 j \). What is the name of \( \mathbf{g} \) in Section 15.3? It exists because \( yi + x^2 j \) has zero divergence.

38 Construct \( \mathbf{F} \) so that \( \text{curl } \mathbf{F} = 2xi + 3xy - 5zk \) (which has zero divergence).

39 Split the field \( \mathbf{F} = xyi \) into \( \mathbf{V} + \mathbf{W} \) with \( \text{curl } \mathbf{V} = 0 \) and \( \text{div } \mathbf{W} = 0 \).

40 Ampère's law for a steady magnetic field \( \mathbf{B} \) is \( \text{curl } \mathbf{B} = \mu \mathbf{J} \) (\( \mathbf{J} = \text{current density}, \mu = \text{constant} \)). Find the work done by \( \mathbf{B} \) around a space curve \( C \) from the current passing through it.

Maxwell allows varying currents which brings in the electric field.

41 For \( \mathbf{F} = (x^2 + y^2)i \), compute \( \text{curl } (\text{curl } \mathbf{F}) \) and \( \text{grad } (\text{div } \mathbf{F}) \) and \( \mathbf{F}_{xx} + \mathbf{F}_{yy} + \mathbf{F}_{zz} \).

42 For \( \mathbf{F} = \rho(x, y, z)i \), prove these useful identities:
   (a) \( \text{curl } (\text{curl } \mathbf{F}) = \text{grad } (\text{div } \mathbf{F}) - (\mathbf{F}_{xx} + \mathbf{F}_{yy} + \mathbf{F}_{zz}) \).
   (b) \( \text{curl } (f \mathbf{F}) = f \text{curl } \mathbf{F} + (\text{grad } f) \times \mathbf{F} \).

43 If \( \mathbf{B} = a \cos \tau \) (constant direction \( a \)), find \( \text{curl } \mathbf{E} \) from Faraday's Law. Then find the alternating spin field \( \mathbf{E} \).

44 With \( \mathbf{G}(x, y, z) = mi + nj + pk \), write out \( \mathbf{F} \times \mathbf{G} \) and take its divergence. Match the answer with \( \mathbf{G} \), \( \text{curl } \mathbf{F} \) and \( \text{curl } \mathbf{G} \).

45 Write down Green's Theorem in the \( xz \) plane from Stokes' Theorem.

46 True or false: \( \nabla \times \mathbf{F} \) is perpendicular to \( \mathbf{F} \).

47 (a) The second proof of Stokes' Theorem took \( M^* = M(x, y, f(x, y)) + P(x, y, f(x, y)) \partial f/ \partial x \) as the \( M \) in Green's Theorem. Compute \( \partial M^*/ \partial y \) from the chain rule and product rule (there are five terms).
   (b) Similarly \( N^* = N(x, y, f) + P(x, y, f) \partial f/ \partial y \) has the \( x \) derivative \( N_x + N_{xx} + P_x f_x + P_{xx} f_x + P_{xx} f_{xx} \). Check that \( N^* \) matches the right side of equation (8), as needed in the proof.

48 "The shadow of the boundary is the boundary of the shadow." This fact was used in the second proof of Stokes' Theorem, going to Green's Theorem on the shadow. Give two examples of \( S \) and \( C \) and their shadows.

49 Which integrals are equal when \( C = \text{boundary of } S \) or \( S = \text{boundary of } V \)?

\[
\begin{align*}
\int \int \int (\text{curl } \mathbf{F}) \cdot d\mathbf{R} & = \int \int (\text{curl } \mathbf{F}) \cdot \mathbf{n} \, dS \\
\int \int \int \mathbf{F} \cdot n \, dS & = \int \int \int \mathbf{F} \cdot dS \\
\int \int \int (\text{div } \mathbf{F}) \cdot n \, dS & = \int \int (\text{grad } \text{div } \mathbf{F}) \cdot n \, dS \\
\int \int \int \mathbf{F} \cdot dV & = \int \int \int \mathbf{F} \cdot dS \\
\end{align*}
\]

50 Draw the field \( \mathbf{V} = -zk \) spinning a wheel in the \( xz \) plane. What wheels would not spin?
### CHAPTER 14  
**Multiple Integrals**

<table>
<thead>
<tr>
<th>Section</th>
<th>Title</th>
<th>Page</th>
</tr>
</thead>
<tbody>
<tr>
<td>14.1</td>
<td>Double Integrals</td>
<td>521</td>
</tr>
<tr>
<td>14.2</td>
<td>Changing to Better Coordinates</td>
<td>527</td>
</tr>
<tr>
<td>14.3</td>
<td>Triple Integrals</td>
<td>536</td>
</tr>
<tr>
<td>14.4</td>
<td>Cylindrical and Spherical Coordinates</td>
<td>541</td>
</tr>
</tbody>
</table>

### CHAPTER 15  
**Vector Calculus**

<table>
<thead>
<tr>
<th>Section</th>
<th>Title</th>
<th>Page</th>
</tr>
</thead>
<tbody>
<tr>
<td>15.1</td>
<td>Vector Fields</td>
<td>549</td>
</tr>
<tr>
<td>15.2</td>
<td>Line Integrals</td>
<td>555</td>
</tr>
<tr>
<td>15.3</td>
<td>Green's Theorem</td>
<td>563</td>
</tr>
<tr>
<td>15.4</td>
<td>Surface Integrals</td>
<td>573</td>
</tr>
<tr>
<td>15.5</td>
<td>The Divergence Theorem</td>
<td>582</td>
</tr>
<tr>
<td>15.6</td>
<td>Stokes' Theorem and the Curl of ( \mathbf{F} )</td>
<td>589</td>
</tr>
</tbody>
</table>

### CHAPTER 16  
**Mathematics after Calculus**

<table>
<thead>
<tr>
<th>Section</th>
<th>Title</th>
<th>Page</th>
</tr>
</thead>
<tbody>
<tr>
<td>16.1</td>
<td>Linear Algebra</td>
<td>599</td>
</tr>
<tr>
<td>16.2</td>
<td>Differential Equations</td>
<td>603</td>
</tr>
<tr>
<td>16.3</td>
<td>Discrete Mathematics</td>
<td>611</td>
</tr>
</tbody>
</table>

- Study Guide For Chapter 1  
- Answers to Odd-Numbered Problems  
- Index  
- Table of Integrals
I would like this book to do more than help you pass calculus. (I hope it does that too.) After calculus you will have choices—Which mathematics course to take next?—and these pages aim to serve as a guide. Part of the answer depends on where you are going—toward engineering or management or teaching or science or another career where mathematics plays a part. The rest of the answer depends on where the courses are going. This chapter can be a useful reference, to give a clearer idea than course titles can do:

Linear Algebra   Differential Equations   Discrete Mathematics

Advanced Calculus (with Fourier Series)   Numerical Methods   Statistics

Pure mathematics is often divided into analysis and algebra and geometry. Those parts come together in the “mathematical way of thinking”—a mixture of logic and ideas. It is a deep and creative subject—here we make a start.

Two main courses after calculus are linear algebra and differential equations. I hope you can take both. To help you later, Sections 16.1 and 16.2 organize them by examples. First a few words to compare and contrast those two subjects.

Linear algebra is about systems of equations. There are \( n \) variables to solve for. A change in one affects the others. They can be prices or velocities or currents or concentrations—outputs from any model with interconnected parts.

Linear algebra makes only one assumption—the model must be linear. A change in one variable produces proportional changes in all variables. Practically every subject begins that way. (When it becomes nonlinear, we solve by a sequence of linear equations. Linear programming is nonlinear because we require \( x \geq 0 \).) Elsewhere I wrote that “Linear algebra has become as basic and as applicable as calculus, and fortunately it is easier.” I recommend taking it.

A differential equation is continuous (from calculus), where a matrix equation is discrete (from algebra). The rate \( dy/dt \) is determined by the present state \( y \)—which changes by following that rule. Section 16.2 solves \( y' = cy + s(t) \) for economics and life sciences, and \( y'' + by' + cy = f(t) \) for physics and engineering. Please keep it and refer to it.
A third key direction is discrete mathematics. Matrices are a part, networks and algorithms are a bigger part. Derivatives are not a part—this is closer to algebra. It is needed in computer science. Some people have a knack for counting the ways a computer can send ten messages in parallel—and for finding the fastest way.

Typical question: Can 25 states be matched with 25 neighbors, so one state in each pair has an even number of letters? New York can pair with New Jersey, Texas with Oklahoma, California with Arizona. We need rules for Hawaii and Alaska. This matching question doesn't sound mathematical, but it is.

Section 16.3 selects four topics from discrete mathematics, so you can decide if you want more.

Go back for a moment to calculus and differential equations. A completely realistic problem is seldom easy, but we can solve models. (Developing a good model is a skill in itself.) One method of solution involves complex numbers:

any function \( u(x + iy) \) solves \( u_{xx} + u_{yy} = 0 \) (Laplace equation)

any function \( e^{i(kx + \alpha)} \) solves \( u_{tt} - c^2 u_{xx} = 0 \) (wave equation).

From those building blocks we assemble solutions. For the wave equation, a signal starts at \( t = 0 \). It is a combination of pure oscillations \( e^{ikx} \). The coefficients in that combination make up the Fourier transform—to tell how much of each frequency is in the signal. A lot of engineers and scientists would rather know those Fourier coefficients than \( f(x) \).

A Fourier series breaks the signal into \( \Sigma a_k \cos kx \) or \( \Sigma b_k \sin kx \) or \( \Sigma c_k e^{ikx} \). These sums can be infinite (like power series). Instead of values of \( f(x) \), or derivatives at the basepoint, the function is described by \( a_k, b_k, c_k \). Everything is computed by the “Fast Fourier Transform.” This is the greatest algorithm since Newton’s method.

A radio signal is near one frequency. A step function has many frequencies. A delta function has every frequency in the same amount: \( \delta(x) = \Sigma \cos kx \). Channel 4 can’t broadcast a perfect step function. You wouldn’t want to hear a delta function.

We mentioned computing. For nonlinear equations this means Newton’s method. For \( Ax = b \) it means elimination—algorithms take the place of formulas. Exact solutions are gone—speed and accuracy and stability become essential. It seems right to make scientific computing a part of applied mathematics, and teach the algorithms with the theory. My text Introduction to Applied Mathematics is one step in this direction, trying to present advanced calculus as it is actually used.

We cannot discuss applications and forget statistics. Our society produces oceans of data—somebody has to draw conclusions. To decide if a new drug works, and if oil spills are common or rare, and how often to have a checkup, we can’t just guess. I am astounded that the connection between smoking and health was hidden for centuries. It was in the data! Eventually the statisticians uncovered it. Professionals can find patterns, and the rest of us can understand (with a little mathematics) what has been found.

One purpose in studying mathematics is to know more about your own life. Calculus lights up a key idea: Functions. Shapes and populations and heart signals and profits and growth rates, all are given by functions. They change in time. They have integrals and derivatives. To understand and use them is a challenge—mathematics takes effort. A lot of people have contributed, in whatever way they could—as you and I are doing. We may not be Newton or Leibniz or Gauss or Einstein, but we can share some part of what they created.
You have met the idea of a matrix. An $m$ by $n$ matrix $A$ has $m$ rows and $n$ columns (it is square if $m = n$). It multiplies a vector $x$ that has $n$ components. The result is a vector $Ax$ with $m$ components. The central problem of linear algebra is to go backward: From $Ax = b$ find $x$. That is possible when $A$ is square and invertible. Otherwise there is no solution $x$—or there are infinitely many.

The crucial property of matrix multiplication is linearity. If $Ax = b$ and $AY = B$ then $A$ times $x + X$ is $b + B$. Also $A$ times $2x$ is $2b$. In general $A$ times $cx$ is $cb$. In particular $A$ times $0$ is $0$ (one vector has $n$ zeros, the other vector has $m$ zeros). The whole subject develops from linearity. Derivatives and integrals obey linearity too.

**Question 1** What are the solutions to $Ax = 0$? One solution is $x = 0$. There may be other solutions and they fill up the "nullspace":

$$
\begin{bmatrix}
1 & 2 \\
0 & 3
\end{bmatrix}
\begin{bmatrix}
x \\
y
\end{bmatrix} =
\begin{bmatrix}
0 \\
0
\end{bmatrix}
$$

requires $x = 0, y = 0$ also allows $y = -1, z = 3$

When there are more unknowns than equations—when $A$ has more columns than rows—the system $Ax = 0$ has many solutions. They are not scattered randomly around! Another solution is $X = 4, Y = -2, Z = 6$. This lies on the same line as $(2, -1, 3)$ and $(0, 0, 0)$. Always the solutions to $Ax = 0$ form a "space" of vectors—which brings us to a central idea of linear algebra.

**Note** These pages are not concentrating on the mechanics of multiplying or inverting matrices. Those are explained in all courses. My own teaching emphasizes that $Ax$ is a combination of the columns of $A$. The solution $x = A^{-1}b$ is computed by elimination. Here we explain the deeper idea of a vector space—and especially the particular spaces that control $Ax = b$. I cannot go into the same detail as in my book on *Linear Algebra and Its Applications*, where examples and exercises develop the new ideas. Still these pages can be a useful support.

All vectors with $n$ components lie in $n$-dimensional space. You can add them and subtract them and multiply them by any $c$. (Don’t multiply two vectors and never write $1/x$ or $1/A$). The results $x + X$ and $x - X$ and $cx$ are still vectors in the space. Here is the important point:

**The line of solutions to $Ax = 0$ is a "subspace"**—a vector space in its own right. The sum $x + X$ has components $6, -3, 9$—which is another solution. The difference $x - X$ is a solution, and so is $4x$. These operations leave us in the subspace.

The nullspace consists of all solutions to $Ax = 0$. It may contain only the zero vector (as in the first example). It may contain a line of vectors (as in the second example). It may contain a whole plane of vectors (Problem 5). In every case $x + X$ and $x - X$ and $cx$ are also in the nullspace. We are assigning a new word to an old idea—the equation $x - 2y = 0$ has always been represented by a line (its nullspace). Now we have 6-dimensional subspaces of an 8-dimensional vector space.

Notice that $x^2 - y = 0$ does not produce a subspace (a parabola instead). Even the $x$ and $y$ axes together, from $xy = 0$, do not form a subspace. We go off the axes when we add $(1, 0)$ to $(0, 1)$. You might expect the straight line $x - 2y = 1$ to be a subspace, but again it is not so. When $x$ and $y$ are doubled, we have $X - 2Y = 2$. Then $(X, Y)$ is on a different line. Only $Ax = 0$ is guaranteed to produce a subspace.
Figure 16.1 shows the nullspace and "row space." Check dot products (both zero).

**Question 2** When \( A \) multiplies a vector \( x \), what subspace does \( Ax \) lie in? The product \( Ax \) is a combination of the columns of \( A \)—hence the name "column space":

\[
\begin{bmatrix}
2 & 1 \\
0 & 1 \\
0 & 0
\end{bmatrix}
\begin{bmatrix}
3 \\
2
\end{bmatrix} = \begin{bmatrix}
8 \\
2 \\
0
\end{bmatrix}
\]

This is in the column space; \( \begin{bmatrix} 0 \\ 0 \\ 1 \end{bmatrix} \) is not.

No choice of \( x \) can produce \( Ax = (0, 0, 1) \). For this \( A \), all combinations of the columns end in a \( (0, 0, 1) \). The column space is like the xy plane within xyz space. It is a subspace of \( m \)-dimensional space, containing every vector \( b \) that is a combination of the columns:

The system \( Ax = b \) has a solution exactly when \( b \) is in the column space.

When \( A \) has an inverse, the column space is the whole \( n \)-dimensional space. The nullspace contains only \( x = 0 \). There is exactly one solution \( x = A^{-1}b \). This is the good case—and we outline four more key topics in linear algebra.

1. **Basis and dimension of a subspace.** A one-dimensional subspace is a line. A plane has dimension two. The nullspace above contained all multiples of \( (2, -1, 3) \)—by knowing that "basis vector" we know the whole line. The column space was a plane containing column 1 and column 2. Again those vectors are a "basis"—by knowing the columns we know the whole column space.

   Our 2 by 3 matrix has three columns: \((1, 0)\) and \((2, 3)\) and \((0, 1)\). Those are not a basis for the column space! This space is only a plane, and three vectors are too many. The dimension is two. By combining \((1, 0)\) and \((0, 1)\) we can produce the other vector \((2, 3)\). There are only two independent columns, and they form a basis for this column space.

   In general: When a subspace contains \( r \) independent vectors, and no more, those vectors are a basis and the dimension is \( r \). "Independent" means that no vector is a combination of the others. In the example, \((1, 0)\) and \((2, 3)\) are also a basis. A subspace has many bases, just as a plane has many axes.

2. **Least squares.** If \( Ax = b \) has no solution, we look for the \( x \) that comes closest. Section 11.4 found the straight line nearest to a set of points. We make the length of \( Ax - b \) as small as possible, when zero length is not possible. No vector solves
Ax = b, when b is not in the column space. So b is projected onto that space. This leads to the “normal equations” that produce the best x:

$$A^T Ax = A^T b.$$  

(1)

When a rectangular matrix appears in applications, its transpose generally comes too. The columns of $A$ are the rows of $A^T$. The rows of $A$ are the columns of $A^T$. Then $A^T A$ is square and symmetric—equal to its transpose and vital for applied mathematics.

3. **Eigenvalues (for square matrices only).** Normally $Ax$ points in a direction different from $x$. For certain special *eigenvectors*, $Ax$ is parallel to $x$. Here is a 2 by 2 matrix with two eigenvectors—in one case $Ax = 5x$ and in the other $Ax = 2x$:

$$Ax = \lambda x: \begin{bmatrix} 3 & 2 \\ 1 & 4 \end{bmatrix} \begin{bmatrix} 1 \\ 1 \end{bmatrix} = \begin{bmatrix} 5 \\ 5 \end{bmatrix} = 5 \begin{bmatrix} 1 \\ 1 \end{bmatrix} \quad \text{and} \quad \begin{bmatrix} 3 & 2 \\ 1 & 4 \end{bmatrix} \begin{bmatrix} 2 \\ -1 \end{bmatrix} = \begin{bmatrix} 4 \\ -2 \end{bmatrix} = 2 \begin{bmatrix} 2 \\ -1 \end{bmatrix}.$$

The multipliers 5 and 2 are the *eigenvalues* of $A$. An 8 by 8 matrix has eight eigenvalues, which tell what the matrix is doing (to the eigenvectors). The eigenvectors are uncoupled, and they go their own way. A system of equations $dy/dt = Ay$ acts like one equation—when $y$ is an eigenvector:

$$dy_1/dt = 3y_1 + 2y_2 \quad \text{has the solution} \quad y_1 = e^{5t} \quad \text{which is} \quad e^{5t} \begin{bmatrix} 1 \\ 1 \end{bmatrix}.$$  

The eigenvector is (1, 1). The eigenvalue $\lambda = 5$ is in the exponent. When you substitute $y_1$ and $y_2$ the differential equations become $5e^{5t} = 5e^{5t}$. The fundamental principle for $dy/dt = cy$ still works for the system $dy/dt = Ay$: *Look for pure exponential solutions.* The eigenvalue “lambda” is the growth rate in the exponent.

I have to add: Find the eigenvectors also. The second eigenvector $(2, -1)$ has eigenvalue $\lambda = 2$. A second solution is $y_1 = 2e^{2t}$, $y_2 = -e^{2t}$. Substitute those into the equation—they are even better at displaying the general rule:

*If* $Ax = \lambda x$ *then* $d/dt(e^{\lambda t}x) = A(e^{\lambda t}x)$. *The pure exponentials are* $y = e^{\lambda t}x$.

The four entries of $A$ pull together for the eigenvector. So do the 64 entries of an 8 by 8 matrix—again $e^{\lambda t}x$ solves the equation. Growth or decay is decided by $\lambda > 0$ or $\lambda < 0$. When $\lambda = k + io$ is a complex number, growth and oscillation combine in $e^{\lambda t} = e^{kt} e^{i\omega t} = e^{kt} (\cos \omega t + i \sin \omega t)$.

Subspaces govern static problems $Ax = b$. Eigenvalues and eigenvectors govern dynamic problems $dy/dt = Ay$. Look for exponentials $y = e^{\lambda t}x$.

4. **Determinants and inverse matrices.** A 2 by 2 matrix has determinant $D = ad - bc$. This matrix has no inverse if $D = 0$. Reason: $A^{-1}$ divides by $D$:

$$A^{-1} = \frac{1}{D} \begin{bmatrix} d & -b \\ -c & a \end{bmatrix} \quad \text{times} \quad A = \begin{bmatrix} a & b \\ c & d \end{bmatrix} \quad \text{equals} \quad I = \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix}.$$  

This pattern extends to $n$ by $n$ matrices, but $D$ and $A^{-1}$ become more complicated.

For 3 by 3 matrices $D$ has six terms. Section 11.5 identified $D$ as a triple product $a \cdot (b \times c)$ of the columns. Three events come together in the singular case: $D$ is zero and $A$ has no inverse and the columns lie in a plane. The opposite events produce the “nonsingular” case: $D$ is nonzero and $A^{-1}$ exists. Then $Ax = b$ is solved by $x = A^{-1}b$. 
16 Mathematics after Calculus

D is also the product of the pivots and the product of the eigenvalues. The pivots arise in elimination—the practical way to solve $Ax = b$ without $A^{-1}$. To find eigenvalues we turn $Ax = \lambda x$ into $(A - \lambda I)x = 0$. By a nice twist of fate, this matrix $A - \lambda I$ has $D = 0$. Go back to the example:

$$
\begin{bmatrix}
3 & 2 \\
1 & 4
\end{bmatrix} - \lambda
\begin{bmatrix}
1 & 0 \\
0 & 1
\end{bmatrix} =
\begin{bmatrix}
3 - \lambda & 2 \\
1 & 4 - \lambda
\end{bmatrix}
$$

has $D = (3 - \lambda)(4 - \lambda) - 2 = \lambda^2 - 7\lambda + 10$.

The equation $\lambda^2 - 7\lambda + 10 = 0$ gives $\lambda = 5$ and $\lambda = 2$. The eigenvalues come first, to make $D = 0$. Then $(A - 5I)x = 0$ and $(A - 2I)x = 0$ yield the eigenvectors. These $x$'s go into $y = e^{\mu x}$ to solve differential equations—which come next.

16.1 Exercises

Read-through questions

If $Ax = b$ and $AX = B$, then $A$ times $2x + 3X$ equals ____. If $Ax = 0$ and $AX = 0$ then $A$ times $2x + 3X$ equals ___. In this case $x$ and $X$ are in the __ of $A$, and so is the combination ___. The nullspace contains all solutions to ___. It is a subspace, which means ___. If $x = (1, 1, 1)$ is in the nullspace then the columns add to ___, so they are (independent)(dependent).

Another subspace is the ___ space of $A$, containing all combinations of the columns. The system $Ax = b$ can be solved when $b$ is ___. Otherwise the best solution comes from $A^TAx = x$. Here $A^T$ is the ___ matrix, whose rows are ___. The nullspace of $A^T$ contains all solutions to ___. The ___ space of $A^T$ (row space of $A$) is the fourth fundamental subspace. Each subspace has a basis containing as many ___ vectors as possible. The number of vectors in the basis is the ___ of the subspace.

When $Ax = \lambda x$, the number $\lambda$ is an ___ and $x$ is an ___. The equation $dy/dt = Ay$ has the exponential solution $y = e^{\lambda t}$. A $7$ by $7$ matrix has ___ eigenvalues, whose product is the ___ $D$. If $D$ is nonzero the matrix $A$ has an ___. Then $Ax = b$ is solved by $x = w$. The formula for $D$ contains $7! = 5040$ terms, so $x$ is better computed by ___. On the other hand $Ax = \lambda x$ means that $A - \lambda I$ has determinant ___. The eigenvalue is computed before the ___.

Find the nullspace in 1–6. Along with $x$ go all $cx$.

1 $A = \begin{bmatrix} 1 & 2 \\ 2 & 4 \end{bmatrix}$ (solve $Ax = 0$) 2 $B = \begin{bmatrix} 12 & -6 \\ -6 & 3 \end{bmatrix}$

3 $C = \begin{bmatrix} 1 & 0 \\ 0 & 1 \\ 1 & 2 \end{bmatrix}$ (solve $Cx = 0$) 4 $C^T = \begin{bmatrix} 1 & 0 & 1 \\ 0 & 1 & 2 \end{bmatrix}$

5 $E = \begin{bmatrix} 1 & 1 & 1 \\ 1 & 1 & 1 \\ 1 & 1 & 1 \end{bmatrix}$ 6 $F = \begin{bmatrix} 2 & 1 \\ 1 & 2 \end{bmatrix}$

7 Change Problem 1 to $Ax = \begin{bmatrix} 3 \\ 6 \end{bmatrix}$ (a) Find any particular solution $x_p$. (b) Add any $x_0$ from the nullspace and show that $x_p + x_0$ is also a solution.

8 Change Problem 1 to $Ax = \begin{bmatrix} 1 \\ 0 \end{bmatrix}$ and find all solutions.

Graph the lines $x_1 + 2x_2 = 1$ and $2x_1 + 4x_2 = 0$ in a plane.

9 Suppose $Ax_p = b$ and $Ax_0 = 0$. Then by linearity $A(x_p + x_0) =$ ___ Conclusion: The sum of a particular solution $x_p$ and any nullvector $x_0$ is ___.

10 Suppose $Ax = b$ and $Ax_p = b$. Then by linearity $A(x - x_p) =$ ___. The difference between solutions is a vector in ___. Conclusion: Every solution has the form $x = x_p + x_0$, one particular solution plus a vector in the nullspace.

11 Find three vectors $b$ in the column space of $E$. Find all vectors $b$ for which $Ex = b$ can be solved.

12 If $Ax = 0$ then the rows of $A$ are perpendicular to $x$. Draw the row space and nullspace (lines in a plane) for $A$ above.

13 Compute $CC^T$ and $C^TC$. Why not $C^T$?

14 Show that $Cx = b$ has no solution, if $b = (-1, 1, 1)$. Find the best solution from $C^TCx = C^Tb$.

15 $C^T$ has three columns. How many are independent? Which ones?

16 Find two independent vectors that are in the column space of $C$ but are not columns of $C$.

17 For which of the matrices $ABCCEF$ are the columns a basis for the column space?
18. Explain the reasoning: If the columns of a matrix $A$ are independent, the only solution to $Ax = 0$ is $x = 0$.

19. Which of the matrices $ABCEF$ have nonzero determinants?

20. Find a basis for the full three-dimensional space using only vectors with positive components.

21. Find the matrix $F^{-1}$ for which $FF^{-1} = I = \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix}$.

22. Verify that $(\text{determinant of } F)^2 = (\text{determinant of } F^2)$.

23. (Important) Write down $F - \lambda I$ and compute its determinant. Find the two numbers $\lambda$ that make this determinant zero. For those two numbers find eigenvectors $x$ such that $Fx = \lambda x$.

24. Compute $G = F^2$. Find the determinant of $G - \lambda I$ and the two $\lambda$'s that make it zero. For those two $\lambda$'s find eigenvectors $x$ such that $Gx = \lambda x$. Conclusion: if $Fx = \lambda x$ then $F^2x = \lambda^2 x$.

25. From Problem 23 find two exponential solutions to the equation $dy/dt = Fy$. Then find a combination of those solutions that starts from $y_0 = (1, 0)$ at $t = 0$.

26. From Problem 24 find two solutions to $dy/dt = Gx$. Then find the solution that starts from $y_0 = (2, 1)$.

27. Compute the determinant of $E - \lambda I$. Find all $\lambda$'s that make this determinant zero. Which eigenvalue is repeated?

28. Which previous problem found eigenvectors for $Ex = 0x$? Find an eigenvector for $Ex = 3x$.

29. Find the eigenvalues and eigenvectors of $A$.

30. Explain the reasoning: A matrix has a zero eigenvalue if and only if its determinant is zero.

31. Find the matrix $H$ whose eigenvalues are 0 and 4 with eigenvectors $(1, 1)$ and $(1, -1)$.

32. If $Fx = \lambda x$ then multiplying both sides by __________ gives $F^{-1}x = \lambda^{-1}x$. If $F$ has eigenvalues 1 and 3 then $F^{-1}$ has eigenvalues __________. The determinants of $F$ and $F^{-1}$ are __________.

33. True or false, with a reason or an example.

(a) The solutions to $Ax = b$ form a subspace.

(b) $\begin{bmatrix} 0 & 2 \\ 0 & 0 \end{bmatrix}$ has $\begin{bmatrix} 1 \\ 0 \end{bmatrix}$ in its nullspace and column space.

(c) $A^T A$ has the same entry in its upper right and lower left corners.

(d) If $Ax = \lambda x$ then $y = e^{\lambda t}$ solves $dy/dt = Ay$.

(e) If the columns of $A$ are not independent, their combinations still form a subspace.

---

16.2 Differential Equations

We just solved differential equations by linear algebra. Those were special systems $dy/dt = Ay$, linear with constant coefficients. The solutions were exponentials, involving $e^{\lambda t}$. The eigenvalues of $A$ were the "growth factors" $\lambda$. This section solves other equations—by no means all. We concentrate on a few that have important applications.

Return for a moment to the beginning—when direct integration was king:

1. $dy/dt = s(t)$  
2. $dy/dt = cy$  
3. $dy/dt = u(y)c(t)$

In 1, $y(t)$ is the integral of $s(t)$. In 2, $y(t)$ is the integral of $cy(t)$. That sounds circular—it only made sense after the discovery of $y = e^{\lambda t}$. This exponential has the correct derivative $cy$. To find it by integration instead of inventing it, separate $y$ from $t$:

Separation and integration also solve 3: $\int dy/ud(y) = \int c(t)dt$. The model logistic equation has $u = y - y^2 = \text{quadratic}$. Equation 2 has $u = y = \text{linear}$. Equation 1 is also a special case with $u = 1 = \text{constant}$. But 2 and 1 are very different, for the following reason.

The compound interest equation $y' = cy$ is growing from inside. The equation $y' = s(t)$ is growing from outside. Where $c$ is a "growth rate," $s$ is a "source." They don't have the same meaning, and they don't have the same units. The combination $y' = cy + s$ was solved in Chapter 6, provided $c$ and $s$ are constant—but applications force us to go further.
In three examples we introduce non-constant source terms.

**EXAMPLE 1** Solve \( \frac{dy}{dt} = cy + s \) with the new source term \( s = e^{kt} \).

**Method** Substitute \( y = Be^{kt} \), with an “undetermined coefficient” \( B \) to make it right:
\[
kBe^{kt} = cBe^{kt} + e^{kt} \quad \text{yields} \quad B = \frac{1}{k-c}.
\]

The source \( e^{kt} \) is the driving term. The solution \( Be^{kt} \) is the response. The exponent is the same! The key idea is to expect \( e^{kt} \) in the response.

**Initial condition** To match \( y_0 \) at \( t = 0 \), the solution needs another exponential. It is the free response \( Ae^{kt} \), which satisfies \( \frac{dy}{dt} = cy \) with no source. To make \( y = Ae^{kt} + Be^{kt} \) agree with \( y_0 \), choose \( A = y_0 - B \):

**Final solution**
\[
y = (y_0 - B)e^{kt} + Be^{kt} = y_0e^{kt} + (e^{kt} - e^{kt})(k-c).
\]

**Exceptional case** \( B = \frac{1}{k-c} \) grows larger as \( k \) approaches \( c \). When \( k = c \) the method breaks down—the response \( Be^{kt} \) is no longer correct. The solution (1) approaches 0, and in the limit we get a derivative. It has an extra factor \( t \):

\[
\frac{e^{kt} - e^{kt}}{k-c} \quad \text{change in} \quad e^{kt} \quad \text{change in} \quad c \quad \rightarrow \quad \frac{d}{dc}(e^{kt}) = te^{kt}.
\]

**The correct response is** \( te^{kt} \) **when** \( k = c \). This is the form to substitute, when the driving rate \( k \) equals the natural rate \( c \) (called resonance).

Add the free response \( y_0e^{kt} \) to match the initial condition.

**EXAMPLE 2** Solve \( \frac{dy}{dt} = cy + s \) with the new source term \( s = \cos kt \).

Substitute \( y = B \sin kt + D \cos kt \). This has two undetermined coefficients \( B \) and \( D \):
\[
kB \cos kt - kD \sin kt = c(B \sin kt + D \cos kt) + \cos kt.
\]

Matching cosines gives \( kB = cD + 1 \). The sines give \(-kD = cB \). Algebra gives \( B, D, y \):

\[
B = \frac{c}{k^2 + c^2} \quad D = \frac{k}{k^2 + c^2} \quad y = \frac{c \sin kt + k \cos kt}{k^2 + c^2}.
\]

**Question** Why do we need both \( B \sin kt \) and \( D \cos kt \) in the response to \( \cos kt \)?

First Answer Equation (3) is impossible if we leave out \( B \) or \( D \).

Second Answer \( \cos kt \) is \( \frac{1}{2}e^{ikt} + \frac{1}{2}e^{-ikt} \). So \( e^{ikt} \) and \( e^{-ikt} \) are both in the response.

**EXAMPLE 3** Solve \( \frac{dy}{dt} = cy + s \) with the new source term \( s = te^{kt} \).

**Method** Look for \( y = Be^{kt} + Dte^{kt} \). Problem 13 determines \( B \) and \( D \). Add \( Ae^{kt} \) as needed, to match the initial value \( y_0 \).

**SECOND-ORDER EQUATIONS**

The equation \( \frac{dy}{dt} = cy \) is first-order. The equation \( \frac{d^2y}{dt^2} = -cy \) is second-order. The first is typical of problems in life sciences and economics—the rate \( \frac{dy}{dt} \) depends on the present situation \( y \). The second is typical of engineering and physical sciences—the acceleration \( \frac{d^2y}{dt^2} \) enters the equation.

If you put money in a bank, it starts growing immediately. If you turn the wheels of a car, it changes direction gradually. The path is a curve, not a sharp corner. Newton’s law is \( F = ma \), not \( F = mv \).
A mathematician compares a straight line to a parabola. The straight line crosses the x axis no more than once. The parabola can cross twice. The equation \( ax^2 + bx + c = 0 \) has two solutions, provided we allow them to be complex or equal. These are exactly the possibilities we face below: two real solutions, two complex solutions, or one solution that counts twice. The quadratic could be \( x^2 - 1 \) or \( x^2 + 1 \) or \( x^2 \). The roots are 1 and \(-1\), \(i\) and \(-i\), 0 and 0.

In solving differential equations the roots appear in the exponent, and are called \( \lambda \).

**EXAMPLE 4** \( y'' = y' + y \): solutions \( y = e^t \) and \( y = e^{-t} \) \( \lambda = 1, -1 \)

**EXAMPLE 5** \( y'' = 0 \): solutions \( y = 1 \) and \( y = t \) \( \lambda = 0, 0 \)

**EXAMPLE 6** \( y'' = -y \): solutions \( y = \cos t \) and \( y = \sin t \) \( \lambda = i, -i \)

Where are the complex solutions? They are hidden in Example 6, which could be written \( y = e^{it} \) and \( y = e^{-it} \). These satisfy \( y'' = -y \) since \( i^2 = -1 \). The use of sines and cosines avoids the imaginary number \( i \), but it breaks the pattern of \( e^\lambda \).

Example 5 also seems to break the pattern—again \( e^\lambda \) is hidden. The solution \( y = 1 \) is \( e^0 \). The other solution \( y = t \) is \( te^0 \). The zero exponent is repeated—another exceptional case that needs an extra factor \( t \).

Exponentials solve every equation with constant coefficients and zero right hand side:

To solve \( ay'' + by' + cy = 0 \) substitute \( y = e^{\lambda t} \) and find \( \lambda \).

This method has three steps, leading to the right exponents \( \lambda = r \) and \( \lambda = s \):

1. With \( y = e^{\lambda t} \) the equation is \( a\lambda^2 e^{\lambda t} + b\lambda e^{\lambda t} + c e^{\lambda t} = 0 \). Cancel \( e^{\lambda t} \).
2. Solve \( a\lambda^2 + b\lambda + c = 0 \). Factor or use the formula \( \lambda = \frac{-b \pm \sqrt{b^2 - 4ac}}{2a} \).
3. Call those roots \( \lambda = r \) and \( \lambda = s \). The complete solution is \( y = Ae^{\lambda t} + Be^{st} \).

The pure exponentials are \( y = e^{\lambda t} \) and \( y = e^{st} \). Depending on \( r \) and \( s \), they grow or decay or oscillate. They are combined with constants \( A \) and \( B \) to match the two conditions at \( t = 0 \). The initial state \( y_0 \) equals \( A + B \). The initial velocity \( y_b \) equals \( rA + sB \) (the derivative at \( t = 0 \)).

**EXAMPLE 7** Solve \( y'' - 3y' + 2y = 0 \) with \( y_0 = 5 \) and \( y_b = 4 \).

Step 1 substitutes \( y = e^{\lambda t} \). The equation becomes \( \lambda^2 e^{\lambda t} - 3\lambda e^{\lambda t} + 2e^{\lambda t} = 0 \). Cancel \( e^{\lambda t} \).
Step 2 solves \( \lambda^2 - 3\lambda + 2 = 0 \). Factor into \( (\lambda - 1)(\lambda - 2) = 0 \). The exponents \( r, s \) are 1, 2.
Step 3 produces \( y = Ae^t + Be^{2t} \). The initial conditions give \( A + B = 5 \) and \( 1A + 2B = 4 \).
The constants are \( A = 6 \) and \( B = -1 \). The solution is \( y = 6e^t - e^{2t} \).

This solution grows because there is a positive \( \lambda \). The equation is "unstable." It becomes stable when the middle term \(-3y'\) is changed to \(+3y'\). When the damping is positive the solution decays. The \( \lambda \)'s are negative:

**EXAMPLE 8** \((\lambda^2 + 3\lambda + 2)\) factors into \((\lambda + 1)(\lambda + 2)\). The exponents are \(-1\) and \(-2\).
The solution is \( y = Ae^{-t} + Be^{-2t} \). It decays to zero for any initial condition.

**EXAMPLES 9-10** Solve \( y'' + 2y' + 2y = 0 \) and \( y'' + 2y' + y = 0 \). How do they differ?

Key difference \( \lambda^2 + 2\lambda + 2 \) has **complex roots**, \( \lambda^2 + 2\lambda + 1 \) has a **repeated root**:

\[
\lambda^2 + 2\lambda + 2 = 0 \quad \text{gives} \quad \lambda = -1 \pm i \quad (\lambda + 1)^2 = 0 \quad \text{gives} \quad \lambda = -1, -1.
\]

The \(-1\) in all these \( \lambda \)'s means decay. The \( i \) means oscillation. The first exponential is \( e^{-1 + it} \), which splits into \( e^{-t} \) (decay) times \( e^{it} \) (oscillation). Even better, change \( e^\lambda \)
and $e^{-it}$ into cosines and sines:

$$y = Ae^{-1+it} + Be^{-1-it} = e^{-t}(a \cos t + b \sin t).$$

(5)

At $t = 0$ this produces $y_0 = a$. Then matching $y'_0$ leads to $b$.

Example 10 has $r = s = -1$ (repeated root). One solution is $e^{-t}$ as usual. The second solution cannot be another $e^{-t}$. Problem 21 shows that it is $te^{-t}$—again the exceptional case multiplies by $t$. The general solution is $y = Ae^{-t} + Bte^{-t}$.

Without the damping term $2y'$, these examples are $y'' + 2y = 0$ or $y'' + y = 0$—pure oscillation. A small amount of damping mixes oscillation and decay. Large damping gives pure decay. The borderline is when $\lambda$ is repeated ($r = s$). That occurs when $b^2 - 4ac$ in the square root is zero. The **borderline between two real roots and two complex roots is two repeated roots**.

The method of solution comes down to one idea: **Substitute $y = e^{\lambda t}$**. The equations apply to mechanical vibrations and electrical circuits (also other things, but those two are of prime importance). While describing these applications I will collect the information that comes from $\lambda$.

**SPRINGS AND CIRCUITS: MECHANICAL AND ELECTRICAL ENGINEERING**

A mass is hanging from a spring. We pull it down an extra distance $y_0$ and give it a starting velocity $y'_0$. The mass moves up or down, obeying Newton's law: *mass times acceleration equals spring force plus damping force*:

$$my'' = -ky - dy' \quad \text{or} \quad my'' + dy' + ky = 0.$$  

(6)

This is **free oscillation**. The spring force $-ky$ is proportional to the stretching $y$ (Hooke's law). The damping acts like a shock absorber or air resistance—it takes out energy. Whether the system goes directly toward zero or swings back and forth is decided by the three numbers $m, d, k$. They were previously called $a, b, c$.

16A The solutions $e^{\lambda t}$ to $my'' + dy' + ky = 0$ are controlled by the roots of $m\lambda^2 + d\lambda + k = 0$. With $d > 0$ there is damping and decay. From $\sqrt{d^2 - 4mk}$ there may be oscillation:

- **overdamping**: $d^2 > 4mk$ gives real roots and pure decay (Example 8)
- **underdamping**: $d^2 < 4mk$ gives complex roots and oscillation (Example 9)
- **critical damping**: $d^2 = 4mk$ gives a real repeated root $-d/2m$ (Example 10)

We are using letters when the examples had numbers, but the results are the same:

$$m\lambda^2 + d\lambda + k = 0 \quad \text{has roots} \quad r, s = -\frac{d}{2m} \pm \frac{1}{2m} \sqrt{d^2 - 4mk}.$$

Overdamping has no imaginary parts or oscillations: $y = Ae^{\lambda t} + Be^{\mu t}$. Critical damping has $r = s$ and an exceptional solution with an extra $t$: $y = Ae^{\lambda t} + Bte^{\mu t}$. (This is only a solution when $r = s$.) Underdamping has decay from $-d/2m$ and oscillation from the imaginary part. An undamped spring ($d = 0$) has pure oscillation at the natural frequency $\omega_0 = \sqrt{k/m}$.

*All these possibilities are in Figure 16.2*, created by Alar Toomre. At the top is pure oscillation ($d = 0$ and $y = \cos 2t$). The equation is $y'' + dy' + 4y = 0$ and $d$ starts to grow. When $d$ reaches 4, the quadratic is $\lambda^2 + 4\lambda + 4$ or $(\lambda + 2)^2$. The repeated root
yields $e^{-2t}$ and $te^{-2t}$. After that the oscillation is gone. There is a smooth transition from one case to the next—as complex roots join in the repeated root and split into real roots.

At the bottom right, the final value $y(2\pi)$ increases with large damping. This was a surprise. At $d = 5$ the roots are $-1$ and $-4$. At $d = 8.5$ the roots are $-1.5$ and $-8$. The small root gives slow decay (like molasses). As $d \to \infty$ the solution approaches $y = 1$.

If we are serious about using mathematics, we should take advantage of anything that helps. For second-order equations, the formulas look clumsy but the examples are quite neat. The idea of $e^{zt}$ is absolutely basic. The good thing is that electrical circuits satisfy the same equation. There is a beautiful analogy between springs and circuits:

\[
\begin{align*}
\text{mass } m & \leftrightarrow \text{inductance } L \\
\text{damping constant } d & \leftrightarrow \text{resistance } R \\
\text{elastic constant } k & \leftrightarrow 1/(\text{capacitance } C)
\end{align*}
\]

The resistor takes out energy as the shock absorber did—converting into heat by friction. Without resistance we have pure oscillation. Electric charge is stored in the capacitor (like potential energy). It is released as current (like kinetic energy). It is stored up again (like a stretched spring). This continues at a frequency $\omega_0 = 1/\sqrt{LC}$ (like the spring’s natural frequency $\sqrt{k/m}$). These analogies turn mechanical engineers into electrical engineers and vice versa.

The equation for the current $y(t)$ now includes a driving term on the right:

\[
Ly'' + Ry' + \frac{1}{C} \int y \, dt = \text{applied voltage } = V \sin \omega t.
\]

(7)

To match networks with springs, differentiate both sides of (7):

\[
Ly'' + Ry' + y/C = V\omega \cos \omega t.
\]

(8)

The oscillations are free when $V = 0$ and forced when $V \neq 0$. The free oscillations $e^{zt}$ are controlled by $L\lambda^2 + R\lambda + 1/C = 0$. Notice the undamped case $R = 0$ when
\[ \lambda = \pm i\sqrt{LC}. \] This shows the natural frequency \( \omega_0 = 1/\sqrt{LC}. \) Damped free oscillations are in the exercises—what is new and important is the forcing from the right hand side. Our last step is to solve equation (8).

**PARTICULAR SOLUTIONS—THE METHOD OF UNDETERMINED COEFFICIENTS**

The forcing term is a multiple of \( \cos \omega t. \) The "particular solution" is a multiple of \( \cos \omega t \) plus a multiple of \( \sin \omega t. \) To discover the undetermined coefficients in \( y = a \cos \omega t + b \sin \omega t, \) substitute into the differential equation (8):

\[
-\lambda \omega^2 (a \cos \omega t + b \sin \omega t) + \lambda \omega (-a \sin \omega t + b \cos \omega t) + (a \cos \omega t + b \sin \omega t)/C = V \omega \cos \omega t.
\]

The terms in \( \cos \omega t \) and the terms in \( \sin \omega t \) give two equations for \( a \) and \( b: \)

\[
-alomega^2 + bR\omega + a/C = V\omega \quad \text{and} \quad -b\omega^2 - aR\omega + b/C = 0. \tag{9}
\]

**EXAMPLE 11** Solve \( y'' + y = \cos \omega t. \) The oscillations are forced at frequency \( \omega. \) The oscillations are free \( (y'' + y = 0) \) at frequency 1. The solution contains both.

**Particular solution** Set \( y = a \cos \omega t + b \sin \omega t \) at the driving frequency \( \omega, \) and (9) becomes

\[
-alomega^2 + bR\omega + a/C = V\omega \quad \text{and} \quad -b\omega^2 - aR\omega + b/C = 0.
\]

The second equation gives \( b = 0. \) No sines are needed because the problem has no \( dy/dt. \) The first equation gives \( a = 1/(1 - \omega^2), \) which multiplies the cosine:

\[
y = (\cos \omega t)(1 - \omega^2) \quad \text{solves} \quad y'' + y = \cos \omega t. \tag{10}
\]

**General solution** Add to this particular solution any solution to \( y'' + y = 0: \)

\[
y = y_{\text{particular}} + y_{\text{homogeneous}} = (\cos \omega t)(1 - \omega^2) + C \cos t + D \sin t. \tag{11}
\]

**Problem of resonance** When the driving frequency is \( \omega = 1, \) the solution (11) becomes meaningless—its denominator is zero. **Reason:** The natural frequency in \( \cos t \) and \( \sin t \) is also 1. A new particular solution comes from \( t \cos t \) and \( t \sin t. \)

The key to success is to know the form for \( y. \) The table displays four right hand sides and the correct \( y \)'s for any constant-coefficient equation:

<table>
<thead>
<tr>
<th>Right hand side</th>
<th>Particular solution</th>
</tr>
</thead>
<tbody>
<tr>
<td>( e^{kt} )</td>
<td>( y = Be^{kt} ) (same exponent)</td>
</tr>
<tr>
<td>( \cos \omega t ) or ( \sin \omega t )</td>
<td>( y = a \cos \omega t + b \sin \omega t ) (include both)</td>
</tr>
<tr>
<td>polynomial in ( t )</td>
<td>( y = \text{polynomial of the same degree} )</td>
</tr>
<tr>
<td>( e^{kt} \cos \omega t ) or ( e^{kt} \sin \omega t )</td>
<td>( y = ae^{kt} \cos \omega t + be^{kt} \sin \omega t )</td>
</tr>
</tbody>
</table>

**Exception** If one of the roots \( \lambda \) for free oscillation equals \( k \) or \( i\omega \) or 0 or \( k + i\omega, \) the corresponding \( y \) in the table is wrong. The proposed solution would give zero on the right hand side. The correct form for \( y \) includes an extra \( t. \) All particular solutions are computed by substituting into the differential equation.

**Apology** Constant-coefficient equations hardly use calculus (only \( e^{kt}. \) They reduce directly to algebra (substitute \( y, \) solve for \( \lambda \) and \( a \) and \( b). \) I find the S-curve from the logistic equation much more remarkable. The nonlinearity of epidemics or heartbeats or earthquakes demands all the calculus we know. The solution is not so predictable. The extreme of unpredictability came when Lorenz studied weather prediction and discovered chaos.
NUMERICAL METHODS

Those four pages explained how to solve linear equations with constant coefficients: Substitute $y = e^{kt}$. The list of special solutions becomes longer in a course on differential equations. But for most nonlinear problems we enter another world—where solutions are numerical and approximate, not exact.

In actual practice, numerical methods for $\frac{dy}{dt} = f(t, y)$ divide in two groups:

1. Single-step methods like Euler and Runge-Kutta
2. Multistep methods like Adams-Bashforth

The unknown $y$ and the right side $f$ can be vectors with $n$ components. The notation stays the same: the step is $\Delta t = h$, the time $t_n$ is $nh$, and $y_n$ is the approximation to the true $y$ at that time. We test the first step, to find $y_1$ from $y_0 = 1$. The equation is $\frac{dy}{dt} = y$, so the right side is $f = y$ and the true solution is $y = e^t$.

Notice how the first value of $f$ (in this case 1) is used inside the second $f$:

**Improved Euler**

$$y_{n+1} = y_n + \frac{1}{2}h[f(t_n, y_n) + f(t_{n+1}, y_n + hf(y_n, t_n))]$$

**TEST**

$$y_1 = 1 + \frac{1}{2}h[1 + (1 + h)] = 1 + h + \frac{1}{2}h^2$$

At time $h$ the true solution equals $e^h$. Its infinite series is correct through $h^2$ for Improved Euler (a second-order method). The ordinary Euler method $y_{n+1} = y_n + hf(t_n, y_n)$ is first-order. **TEST**: $y_1 = 1 + h$. Now try Runge-Kutta (a fourth-order method):

**Runge-Kutta**

$$y_{n+1} = y_n + \frac{1}{6}[k_1 + 2k_2 + 2k_3 + k_4]$$

$$k_1 = hf(t_n, y_n)$$

$$k_2 = hf(t_n + \frac{1}{2}h, y_n + \frac{1}{2}k_1)$$

$$k_3 = hf(t_n + \frac{1}{2}h, y_n + \frac{1}{2}k_2)$$

$$k_4 = hf(t_n + h, y_n + k_3)$$

Now the first value of $f$ is used in the second (for $k_2$), the second is used in the third, and then $k_3$ is used in $k_4$. The programming is easy. Check the accuracy with another test on $\frac{dy}{dt} = y$:

**TEST**

$$y_1 = 1 + \frac{1}{6}\left[ h + 2h\left(1 + \frac{h}{2}\right) + 2h\left(1 + \frac{h}{2}\right) + h\left(1 + \frac{h}{2}\right)\right]$$

$$= 1 + h + \frac{h^2}{2} + \frac{h^3}{6} + \frac{h^4}{24}.$$ This answer agrees with $e^h$ through $h^4$.

These formulas are included in the book so that you can apply them directly—for example to see the S-shape from the logistic equation with $f = cy - hy^2$.

Multistep formulas are simpler and quicker, but they need a single-step method to get started. Here is $y_n$ in a fourth-order formula that needs $y_0, y_1, y_2, y_3$. Just shift all indices for $y_5, y_6,$ and $y_{n+1}$:

**Multistep**

$$y_4 = y_3 + \frac{h}{24}[55y_3 - 59y_2 + 37y_1 - 9y_0].$$

The advantage is that each step needs only one new evaluation of $y_0 = f(t_n, y_n)$. Runge-Kutta needs four evaluations for the same accuracy.

**Stability** is the key requirement for any method. Now the good test is $y' = -y$. The solution should decay and not blow up. Section 6.6 showed how a large time step makes Euler's method unstable—the same will happen for more accurate formulas. The price of total stability is an "implicit method" like $y_1 = y_0 + \frac{1}{2}h(y' + y_1)$, where the unknown $y_1$ appears also in $y'_1$. There is an equation to be solved at every step. Calculus is ending as it started—with the methods of Isaac Newton.
The solution to \( y' - 5y = 10 \) is \( y = Ae^{5t} + B \). The homogeneous part \( Ae^{5t} \) satisfies \( y' - 5y = 0 \). The particular solution \( B \) equals \( 1 \). The initial condition \( y_0 \) is matched by \( A = \frac{y_0 - B}{e^{5t}} \). For \( y' - 5y = e^{5t} \) the right form is \( y = Ae^{5t} + c \). For \( y' - 5y = \cos t \) the form is \( y = Ae^{5t} + \cos t \).

The equation \( y'' + 4y' + 5y = 0 \) is second-order because \( y'' \). The pure exponential solutions come from the roots of \( \lambda^2 + 4\lambda + 5 = 0 \), which are \( \lambda = -2 \pm i \). The general solution is \( y = e^{-2t}(c \cos t + d \sin t) \). Changing \( 4y' \) to \( 1 \) yields pure oscillation. Changing to \( 2y' \) yields \( \lambda = -1 \pm 2i \), when the solutions become \( y = m e^{-t} \). This oscillation is over(under)(critically) damped. A spring with \( m = 1 \), \( k = 5 \) goes (back and forth)(directly to zero). An electrical network with \( L = 1 \), \( R = 2 \), \( C = 1 \) also \( s = \ldots \).

One particular solution of \( y'' + 4y' + 5y = e^{5t} \) is \( e^{5t} \) times \( c \). If the right side is \( \cos t \), the form of \( y_p = \ldots \). If the right side is \( 1 \) then \( y_p = \ldots \). If the right side is \( \ldots \) we have resonance and \( y_p \) contains an extra factor \( s \).

### Problems 1–14 are about first-order linear equations.

1. Substitute \( y = Be^{5t} \) into \( y' - 5y = e^{5t} \) to find a particular solution.
2. Substitute \( y = a \cos 2t + b \sin 2t \) into \( y' + y = 4 \sin 2t \) to find a particular solution.
3. Substitute \( y = a \cos t + b \sin t \) into \( y' + y = 1 + t^2 \) to find a particular solution.
4. Substitute \( y = Ae^t \) into \( y'' + 2y' + 3y = 0 \) to find a particular solution.
5. In Problem 1 we can add \( Ae^{-t} \) because this solves the equation \( \ldots \). Choose \( A \) so that \( y(0) = 7 \).
6. In Problem 2 we can add \( Ae^{-t} \), which solves \( \ldots \). Choose \( A \) to match \( y(0) = 0 \).
7. In Problem 3 we add \( \ldots \) to match \( y(0) = 2 \).
8. In Problem 4 we can add \( y = A \). Why?
9. Starting from \( y_0 \) = \( 0 \) solve \( y' = e^{kt} \) and also solve \( y' = \ldots \). Show that the first solution approaches the second as \( k \to 0 \).
10. Solve \( y' - y = e^{kt} \) starting from \( y_0 = 0 \). What happens to your formula as \( k \to 1 \)? By l'Hôpital's rule show that \( y \) approaches \( e^t \) as \( k \to 1 \).
11. Solve \( y' - y = e^t + \cos t \). What form do you assume for \( y \) with two terms on the right side?
12. Solve \( y' + y = e^t + t \). What form do you assume for \( y \)?
13. Solve \( y' = cy + te^t \) following Example 3 (\( c \neq 1 \)).
14. Solve \( y' = y + t \) following Example 3 (\( c = 1 \) and \( k = 0 \)).

### Problems 15–28 are about second-order linear equations.

15. Substitute \( y = e^{kt} \) into \( y'' + 6y' + 5y = 0 \). (a) Find all \( \lambda \)'s. (b) The solution decays because \( \ldots \). (c) The general solution with constants \( A \) and \( B \) is \( \ldots \).
16. Substitute \( y = e^{kt} \) into \( y'' + 2y' + 3y = 0 \). (a) Find all \( \lambda \)'s. (b) The solution oscillates because \( \ldots \). (c) The general solution with constants \( a \) and \( b \) is \( \ldots \).
17. Substitute \( y = e^{kt} \) into \( y'' + 2y' + 3y = 0 \). (a) Find all \( \lambda \)'s. (b) The solution oscillates as it decays because \( \ldots \). The general solution with \( A \) and \( B \) and \( e^{kt} \) is \( \ldots \). The general solution with \( e^{-t} \) times cosine is \( \ldots \).
18. Substitute \( y = e^{kt} \) into \( y'' + 6y' + 9y = 0 \). (a) Find all \( \lambda \)'s. (b) The general solution with \( e^{kt} \) and \( e^{-kt} \) is \( \ldots \).
19. For \( y'' + dy' + y = 0 \) find the type of damping at \( d = 0, 1, 2, 3 \).
20. For \( y'' + 2y' + ky = 0 \) find the type of damping at \( k = 0, 1, 2 \).
21. If \( \lambda^2 + 2 \lambda + 1 = 0 \) has a repeated root prove it is \( \lambda = -1 \). In this case compute \( y'' + by' + cy \) when \( y = te^t \).
22. \( \lambda^2 + 3 \lambda + 2 = 0 \) has roots \( -1 \) and \( -2 \) (not repeated). Show that \( te^{-t} \) does not solve \( y'' + 3y' + 2y = 0 \).
23. Find \( y = a \cos t + b \sin t \) to solve \( y'' + y' + y = \cos t \).
24. Find \( y = a \cos wt \) to solve \( y'' + y' + y = \sin wt \).
25. Solve \( y'' + 9y = \cos 3t \), \( y(0) = 0 \) and \( y'(0) = 0 \). The solution contains \( \cos 3t \) and \( \cos 2t \).
26. The difference \( \cos 3t - \cos 2t \) equals \( 2 \sin 4t \). Graph it to see fast oscillations inside slow oscillations (beats).
27. The solution to \( y'' + \omega^2 y = \cos \omega t \) with \( y(0) = 0 \) and \( y'(0) = 0 \) is what multiple of \( \cos \omega t - \cos \omega t \)? The formula breaks down when \( \omega = \ldots \).
28. Substitute \( y = Ae^{int} \) into the circuit equation \( Ly'' + Ry' + \int y dt/C = Ve^{-int} \). Cancel \( e^{-int} \) to find \( A \). Its denominator is the impedance.

### Problems 29–32 have the four right sides in the table (end of section). Find \( y_{\text{particular}} \) by using the correct form.

29. \( y'' + 3y = e^{kt} \) \hspace{1cm} 30. \( y'' + 3y = \sin t \)
31. \( y'' + 2y = 1 + t \) \hspace{1cm} 32. \( y'' + 2y = e^t \cos t \)
33. Find the coefficients of \( y \) in Problems 29–31 for which the forms in the table are wrong. Why are they wrong? What new forms are correct?
34 The magic factor $i$ entered equation (2). The series for $e^{kt} - e^{ct}$ starts with $1 + kt + \frac{1}{2}k^2t^2$ minus $1 + ct + \frac{1}{2}c^2t^2$. Divide by $k - c$ and set $k = c$ to start the series for $te^t$.

35 Find four exponentials $y = e^{at}$ for $d^4y/dt^4 - y = 0$.

36 Find a particular solution to $d^4y/dt^4 + y = 0$.

37 The solution is $y = Ae^{-2t} + Be^{-3t}$ when $d = 4$ in Figure 16.2. Choose $A$ and $B$ to match $y_0 = 1$ and $y'_0 = 0$. How large is $y(2\pi)$?

38 When $d$ reaches 4 the quadratic for Figure 16.2 is $\lambda^2 + 4\lambda + 4 = (\lambda + 1)(\lambda + 4)$. Match $y = Ae^{-\lambda} + Be^{-4\lambda}$ to $y_0 = 1$ and $y'_0 = 0$. How large is $y(2\pi)$?

39 When the quadratic for Figure 16.2 has roots $-r$ and $-4/r$, the solution is $y = Ae^{-r} + Be^{-4r}$.

- (a) Match the initial conditions $y_0 = 1$ and $y'_0 = 0$.
- (b) Show that $y$ approaches 1 as $r \to 0$.

40 In one sentence tell why $y'' = 6y$ has exponential solutions but $y'' = 6y^2$ does not. What power $y = x^n$ solves this equation?

41 The solution to $dy/dt = f(t)$, with no $y$ on the right side, is $y = \int f(t) \, dt$. Show that the Runge-Kutta method becomes Simpson's Rule.

42 Test all methods on the logistic equation $y' = y - y^2$ to see which gives $y_0 = 1$ most accurately. Start at the inflection point $y_0 = \frac{1}{2}$ with $h = 1/\pi$. Begin the multistep method with exact values of $y = (1 + e^{-t})^{-1}$.

43 Extend the tests of Improved Euler and Runge-Kutta to $y' = -y$ with $y_0 = 1$. They are stable if $|y_1| \leq 1$. How large can $h$ be?

44 Apply Runge-Kutta to $y' = -100y + 100\sin t$ with $y_0 = 0$ and $h = .02$. Increase $h$ to .03 to see that instability is no joke.

---

16.3 Discrete Mathematics: Algorithms

Discrete mathematics is not like calculus. Everything is finite. I can start with the 50 states of the U.S. I ask if Maine is connected to California, by a path through neighboring states. You say yes. I ask for the shortest path (fewest states on the way). You get a map and try all possibilities (not really all—but your answer is right). Then I close all boundaries between states like Illinois and Indiana, because one has an even number of letters and the other has an odd number. Is New York still connected to Washington? You ask what kind of game this is—but I hope you will read on.

Far from being dumb, or easy, or useless, discrete mathematics asks good questions. It is important to know the fastest way across the country. It is more important to know the fastest way through a phone network. When you call long distance, a quick connection has to be found. Some lines are tied up, like Illinois to Indiana, and there is no way to try every route.

The example connects New York to New Jersey (7 letters and 9). Washington is connected to Oregon (10 letters and 6). As you read those words, your mind jumps to this fact—there is no path from New York with 7 letters to Washington with 10. Somewhere you must get stuck. There might be a path between all states with an odd number of letters—I doubt it. Graph theory gives a way to find out.

**GRAPHS**

A model for a large part of finite mathematics is a graph. It is not the graph of $y = f(x)$. The word "graph" is used in a totally different way, for a collection of nodes and edges. The nodes are like the 50 states. The edges go between two nodes—the neighboring states. A network of computers fits this model. So do the airline connections between cities. A pair of cities may or may not have an edge between them—depending on flight schedules. The model is determined by $V$ and $E$. 
A **graph** is a set \( V \) of nodes (or vertices) and a set \( E \) of edges.

**Example 1** How many edges are possible with \( n \) nodes, in a complete graph?

The first node has edges to the \( n - 1 \) other nodes. (An edge to itself is not allowed.) The second node has \( n - 2 \) new edges. The third node has another \( n - 3 \). The total count of edges, when none are missing, is the sum from Section 5.3:

\[
1 + 2 + \cdots + (n - 1) = \frac{1}{2} n(n - 1) \quad \text{edges in a complete graph.}
\]

Fifty states have \( 25 \cdot 49 = 1225 \) possible edges. The “neighboring states graph” has less than 200. A line of 6 nodes has 5 edges, out of \( \frac{1}{2} \cdot 6 \cdot 5 = 15 \) possible.

**Example 2** Which states with an odd number of letters are reachable from New York? Boundaries to states like Pennsylvania (12 letters) are closed.

Method of solution Start from New York (7). There is an edge to Connecticut (11). That touches Massachusetts (13), which is a neighbor of Vermont (7). But we missed Rhode Island, and how do we get back? *The order depends on our search method—* and two methods are specially important.

**Depth first search** (DFS) “From the current state, go to one new state if possible.” But what do we do from Vermont, when New Hampshire (12) is not allowed? The answer is: backtrack to Massachusetts. That becomes the next current state.

We label every state as we reach it, to show which state we came from. Then VT has the label MA, and we easily cross back. From MA we go to RI. Then backtrack to MA and CT and NY. At every step I searched for a new state with no success. From NY we see NJ (9). Finally we are in a corner.

The depth first search is ended, by a barrier of even states. Unless we allow Ontario and keep going to Minnesota.

**Breadth first search** (BFS) “From the current state, add all possible new states to the bottom of the list. But take the next current state from the top of the list.” There is no need to backtrack.

From NY we reach VT and MA and CT and NJ. What comes next?

Where DFS moves from the last possible state, breadth first search moves from the first possible state. No move from VT is possible—so we “scan” from Massachusetts. We see Rhode Island (barely). That ends BFS.

The same six states are reached both ways. Only the order is different. DFS is **last in-first out**. BFS is **first in-first out**. You have the same choice in drawing a family tree—follow a path as far as it goes and backtrack, or list all brothers and sisters before their children. The BFS graph in Figure 16.3 is a tree. So is the DFS graph, using forward edges only.

---

**Fig. 16.3** Search trees from New York. The minimum spanning tree.
DEFINITION A tree is a connected graph with no loops. Its \( N \) nodes are connected by \( N - 1 \) edges. If \( N = n \), so every node is in the tree, it is a spanning tree.

The path from VA to KY to TN to NC to VA is a loop (or cycle). If one of those four edges is removed, we have a tree. If two edges are removed, we have two trees (a small forest).

EXAMPLE 3 Allow an edge between neighboring states only when one state is even and the other is odd. Are the lower 48 states connected?

Start anywhere—say California. Apply either type of search—maybe DFS. Go to Arizona (7) then Utah (4) then WY (9) then CO (8) then NM then OK then TX. (I am writing this on an airplane, looking at the map.) We will never get to Florida! It is blocked by Alabama and Georgia.

The search creates a tree, but not a spanning tree. This graph is not connected.

An odd-to-even graph is special and important. It is called "bipartite," meaning two parts. The odd states are in one part, the even states are in the other. All edges go between parts. No edges are within a part.†

EXAMPLE 4 Is there a "complete matching" between 25 even and 25 odd states?

This requires neighboring states to be paired off (with no repetition).

Method 1 Start pairing them off: CA–AZ, UT–WY, NV–ID, NE–SD, WA–MT. What about Oregon? Maybe it should have been paired with Idaho. Then Nevada could pair with Arizona. Trial and error goes nowhere fast.

Method 2 Think first. The four states CA–OR–WA–NV are even. This whole group is only connected to three odd states (AZ, ID, MT). The matching is impossible.

\[\text{16B A complete matching is impossible if a group of nodes in one part is connected to fewer nodes in the other part. If every group has enough connections, the matching can be achieved.}\]

This is Hall's Theorem. In a course on graphs, it would be proved. Our purpose here is to see the ideas and questions in discrete mathematics, more than the proofs.

THE GREEDY ALGORITHM

Put back all edges between neighboring states. The nodes could be provinces of Canada or states of Australia. If they are countries of Europe–Asia–Africa (or the Americas), we need a new map. The essential thing is the new problem.

In a network each edge has a "length." A positive number \( c_{ij} \) is assigned to the edge from node \( i \) to node \( j \). In an economics problem, \( c_{ij} \) is the cost. In a flow problem it is the capacity, in an electrical circuit it is the conductance. We look for paths that minimize these "lengths."

PROBLEM Find the minimum spanning tree. Connect all nodes by a tree with the smallest possible total length.

The six cheapest highways connecting seven cities form a minimum spanning tree. It is cheapest to build, not cheapest to drive—you have to follow the tree. Where there

†Exactly half the states have an even number of letters (a real trivia question). This is the little-known reason for admitting Alaska and Hawaii.
is no edge we set $c_{ij} = \infty$ (or an extremely large value, in an actual code). Then the algorithm works with a complete network—all $n(n-1)/2$ edges are allowed. How does it find the minimum spanning tree in Figure 16.3c?

**Method 1**  Always add the shortest edge that goes out from the current tree.

Starting from node $s$, this rule chooses edges of length 1, 2, 7, 4, 3. Now it skips 5, which would close a loop. It chooses 6, for total length 23.

**Method 2**  Add edges in order, from shortest to longest. Reject an edge that closes a loop. Several trees grow together (a forest). At the end we have a minimum spanning tree.

This variation chooses edge lengths in the order $1, 2, 3, 4, 6$ (rejecting 5), 7. In our network both methods produce the same tree. When many edges have equal length, there can be many shortest trees.

These methods are examples of the **Greedy Algorithm**: **Do the best thing at every step**. Don’t look ahead. Stick to a decision once it is made. In most network problems the Greedy Algorithm is not optimal—in this spanning tree problem it is.

Method 2 looks faster than Method 1. Sort the edges by length, and go down the list. But sorting takes time! It is a fascinating problem in itself—bubble sort or insertion sort or heapsort. We go on to a final example of discrete mathematics and its algorithms.

**PROBLEM**  Find the shortest path from the source node $s$ to each other node.

The shortest path may not go along the minimum spanning tree. In the figure, the best path going east has length $1 + 8$. There is a new shortest path tree, in which the source plays a special role as the “root.”

How do we find shortest paths? Listing all possibilities is more or less insane. A good algorithm builds out from the source, selecting one new edge at every step. After $k$ steps we know the distances $d_1, \ldots, d_k$ to the $k$ nearest nodes.

*Algorithm.*  Minimize $d_i + c_{ij}$ over all settled nodes $i$ and all remaining nodes $j$.

The best new node $j$ is a distance $c_{ij}$ from a settled node, which is a distance $d_i$ from the source. In the example network, the first edges are 1, 2, 7. Next is 8. The northeast node is closest to the source at this step. The final tree does not use edges 3, 5, 6—even though they are short.

These pages were written to show you the algorithmic part of discrete mathematics. The other part is **algebra**—permutations, partitions, groups, counting problems, generating functions. There is no calculus, but that’s fair. The rest of the book was written to show what calculus can do—**I hope very much that you enjoyed it.**

Thank you for reading, and thinking, and working.

---

**16.3 EXERCISES**

**Read-through questions**

A graph is a set $V$ of _a_ and a set $E$ of _b_. With _c_ nodes, a complete graph has _d_ edges. A spanning tree has only _e_. A tree is defined as _f_, and it is spanning if _g_. It has _h_ path between each pair of nodes.

To find a path from node _i_ to node _j_, two search methods are _h_. As nodes are reached, DFS looks out from the _i_. Next is _j_. BFS looks out from _j_. DFS must be prepared to _k_ to earlier nodes. In case of fire, BFS locates all doors from the room you are in before _l_.

---
In a bipartite graph, all edges go from one part to \( m \). A matching is impossible if \( k \) nodes in one part are connected to \( n \) nodes in the other part. The edges in a network have \( c_{ij} \). A minimum spanning tree is \( p \). It can be found by the \( a \) algorithm, which accepts the shortest edge to a new node without worrying about \( t \).

1. Start from one node of a hexagon (six nodes, six edges). Number the other nodes by (a) breadth first search (b) depth first search.

2. Draw two squares with one node in common (7-node graph). From that node number all others by DFS and BFS. Indicate backtracks.

3. How many spanning trees in the hexagon graph?

4. Draw a spanning tree in the two-square graph. How many spanning trees does it have?

5. Define a connected graph. If a graph has 7 edges and 9 nodes, prove that it is not connected.

6. Define a loop. If a connected graph has 8 edges and 9 nodes, prove that it has no loops.

7. Find the shortest path (minimum number of edges) from Maine to California.

8. Which state is farthest (how many edges are needed) from the state you are in? Why would it come last in BFS?

9. List the steps of BFS from your state to Georgia or Colorado or New Jersey. (There are edges Hawaii-California and Alaska-Washington.)

10. With edges between odd neighboring states and between even neighbors, what is the largest connected set of states? Map required.

11. With edges only from odd to even neighbors, how many states can be matched? (Answer unknown to author—please advise.)

12. A matching is a forest of two-node trees. Give another description.

13. Find the minimum spanning tree for network A.

14. Find the shortest path tree from the center of network A.

15. Is there a complete matching between left and right nodes in graph B? If not, which group of nodes has too few connections?

16. Find the loop in network B. Then find a minimum spanning tree by Method 1 and Method 2.

17. How many spanning trees in graph B? It has one loop.

18. Show that a graph cannot have 0, 1, 2, 3, and 4 edges going into its five nodes.

19. If the only edges into a node have lengths 6 and 8, can they both be in a minimum spanning tree?

20. In Problem 19, prove that a minimum spanning tree contains edge (6) if it contains edge (8).

21. True or false, with reason or example.
   (a) In a complete network, the minimum spanning tree contains the \( n - 1 \) shortest edges.
   (b) If a graph has 9 nodes and 9 edges, it has a loop.
   (c) A graph with a complete matching must be connected.

22. Draw a tree that is perfect for (a) DFS; (b) BFS.

23. The adjacency matrix has \( a_{ij} = 1 \) if there is an edge from node \( i \) to node \( j \). Write down this matrix for graphs A and B.

24. In a complete network start with \( d_{ij} = c_{ij} \). Show that the \( d_{ij} \) at the end of this program are shortest distances:
   
   \[
   d_{ij} = \max(d_{ij}, d_{ik} + d_{kj})
   \]

25. How many spanning trees in graph A?

26. A maximum spanning tree has greatest possible length. Give an algorithm to find it.

27. Write a code that will find a spanning tree (or stop), given a list of edges like \( (1, 2), (1, 3), (4, 7) \), ....
CHAPTER 1 INTRODUCTION TO CALCULUS

Section 1.1 Velocity and Distance (page 6)

1 \( v = 30, 0, -30; v = -10, 20 \)

3 \( v(t) = \begin{cases} \begin{array}{ll} 2 & \text{for } \ 0 < t < 10 \\ 1 & \text{for } \ 10 < t < 20 \\ -3 & \text{for } \ 20 < t < 30 \end{array} \end{cases} \)

\( v(t) = \frac{1}{t} \) for \( T < t < 2T \)

\( v(t) = 0 \) for \( 2T < t < 3T \)

5 \( t = 22; t + 10 \)

7 \( v(t) = \begin{cases} \begin{array}{ll} 20 & \text{for } \ t < . \ 2 \\ 0 & \text{for } \ t > . \ 2 \end{array} \end{cases} \)

\( f(t) = \begin{cases} \begin{array}{ll} 20t & \text{for } \ t \leq . \ 2 \\ 4 & \text{for } \ t \geq . \ 2 \end{array} \end{cases} \)

10 \% 12.1%}

13 \( f(t) = 0, 30(t - 1), 30; f(t) = -30t, -60, 30(t - 6) \)

15 Average \( 8, 20 \)

17 \( 40t - 80 \) for \( 1 \leq t \leq 2.5 \)

21 \( 0 \leq t \leq 3, -40 \leq f \leq 20; 0 \leq t \leq 3T, 0 \leq f \leq 60T \)

23 \( 3 - 7t \)

25 \( 6t - 2 \)

27 \( 3t + 7 \)

29 \( -2; 1 \leq f \leq 9 \)

31 \( v(t) = \begin{cases} \begin{array}{ll} 8 & \text{for } \ 0 < t < T \\ -2 & \text{for } \ T < t < 5T \end{array} \end{cases} \)

\( f(t) = \begin{cases} \begin{array}{ll} 8t & \text{for } \ 0 \leq t \leq T \\ 10T - 2t & \text{for } \ T \leq t \leq 5T \end{array} \end{cases} \)

33 \( \frac{9}{5}C + 32; \) slope \( \frac{9}{5} \)

35 \( \frac{f(w)}{1000}; \) slope = conversion factor

37 \( 1 \leq t \leq 5, 0 \leq f \leq 2 \)

39 \( 0 \leq t \leq 5, 0 \leq f \leq 4 \)

41 \( 0 \leq t \leq 5, 1 \leq t \leq 32 \)

43 \( \frac{1}{2}t + 4; \frac{1}{2}t + \frac{7}{2}; 2t + 12; 2t + 3 \)

45 \( -1 \leq t \leq 1: \) ranges \( 0 \leq 2t + 2 \leq 4, \ -3 \leq t - 2 \leq -1, \ -2 \leq f(t) \leq 0, \ 0 \leq f(-t) \leq 2 \)

47 \( \frac{3}{2}V; \frac{3}{2}V \)

49 \( \text{input input} \rightarrow A \)

51 \( 3t + 5, 6t - 1, -3t + 1, 9t - 4; \) slopes \( 3, 3, 6, 6, -3, 9 \)

53 \( \text{The graph goes up and down twice. } f(f(t)) = \begin{cases} \begin{array}{ll} 2(2t) & \text{for } \ 0 \leq t \leq 1.5 \\ 12 - 2(12 - 2t) & \text{for } \ 1.5 \leq t \leq 3 \end{array} \end{cases} \)

Section 1.2 Calculus Without Limits (page 14)

1 \( 2 + 5 + 3 = 10; \) \( f = 1, 3, 8, 11, 10 \)

3 \( f = 3, 4, 6, 7, 7, 6; \) max \( f \) at \( v = 0 \) or at break from \( v = 1 \) to \( -1 \)

5 \( 1.1, -2.5; f(6) = 6.6, -11, 4; f(7) = 7.7, -13, 9 \)

7 \( f(t) = 2t \) for \( t \leq 5, 10 + 3(t - 5) \) for \( t \geq 5; f(10) = 25 \)

9 \( 2t, 2t + 4; \) multiply slopes

11 \( f(8) = 8.8, -15, 14; \) \( \frac{dA}{dt} = 1.1, -2, 5 \)

13 \( f(x) = 3052.50 + .28(x - 20, 350); \) then \( 11, 158.50 \) is \( f(49, 300) \)

17 \( \text{Credit subtracts } 1,000, \) deduction only subtracts \( 15\% \) of \( 1000 \)

19 \( \text{All } v_j = 2; \) \( v_j = (-1)^{-1}; \) \( v_j = (\frac{1}{2})^j \)

21 \( \text{L's have area } 1, 3, 5, 7 \)

23 \( f_j = \text{sum } j^2 + j; \text{ sum } \frac{j^2}{2} + \frac{j}{2} \)

25 \( (101^2 - 99^2)/2 = \frac{400}{2} \)

27 \( v_j = 2j \)

29 \( f_{31} = 5 \)

31 \( a_j = -f_j \)

33 \( 0, 1, 1 \)

35 \( v = 2, 6, 18, 54; 2 \cdot 3^j - 1 \)

37 \( \frac{dA}{dt} = 1.7177, .6956, .6934 \rightarrow \ln 2 = .6931 \) in Chapter 6

39 \( v_j = -(\frac{1}{2})^j \)

41 \( v_j = 2((-1)^j, \text{ sum is } f_j - 1 \)

43 \( v = 1000, t = 10/V \)

47 \( M, N \)

51 \( \sqrt{9} < 2 \cdot 9 < 2^9; \) \( (\frac{1}{2})^2 < 2(\frac{1}{2}) < \sqrt{1/9} < 2^{1/9} \)

Section 1.3 The Velocity at an Instant (page 21)

1 \( 6.6, \frac{13}{2}a, -12, 0, 13 \)

3 \( 4, 3.1, 3 + h, 2.9 \)

5 \( \text{Velocity at } t = 1 \) is \( 3 \)

7 \( \text{Area } f = t + t^2, \) slope of \( f \) is \( 1 + 2t \)

9 \( F; F; F; T \)

11 \( 2; 2t \)

13 \( 12 + 10t^2; 2 + 10t^2 \)

15 \( \text{Time } 2, \) height \( 1, \) stays above \( \frac{3}{4} \) from \( t = \frac{1}{2} \) to \( \frac{3}{2} \)

17 \( f(6) = 18 \)

21 \( v(t) = -2t \) then \( 2t \)

23 \( \text{Average to } t = 5 \) is \( 2; v(5) = 7 \)

25 \( 4v(4t) \)

27 \( v_{\text{ave}} = t, v(t) = 2t \)

Section 1.4 Circular Motion (page 28)

1 \( 10\pi, (0, -1), (-1, 0) \)

3 \( 4 \cos t, 4 \sin t; 4 \) and \( 4t; 4 \cos t \) and \( -4 \sin t \)

5 \( 3t; (\cos 3t, \sin 3t); -3 \sin 3t \) and \( 3 \cos 3t \)

7 \( x = \cos t; \sqrt{2}/2; -\sqrt{2}/2 \)

9 \( 2\pi/3; 1; 2\pi \)

11 \( \text{Clockwise starting at } (1, 0) \)

13 \( \text{Speed } \frac{2}{\pi} \)

15 \( \text{Area } 2 \)

17 \( \text{Area } 0 \)
19 4 from speed, 4 from angle  
21 \frac{1}{4} from radius times 4 from angle gives 1 in velocity  
23 Slope \frac{1}{2}; average \frac{(1 - \sqrt{3})}{\left(\frac{\sqrt{3}}{2}\right)} = \frac{3(2 - \sqrt{3})}{\pi} = 256  
25 Clockwise with radius 1 from (1,0), speed 3  
27 Clockwise with radius 5 from (0,5), speed 10  
29 Counterclockwise with radius 1 from (\cos 1, \sin 1), speed 1  
31 Left and right from (1,0) to (-1,0), v = -\sin t  
33 Up and down between 2 and -2; start 2 \sin \theta, v = 2 \cos(\theta + \theta)  
35 Up and down from (0, -2) to (0,2); v = \frac{1}{2}t  
37 x = \cos \frac{2\pi}{360}, y = \sin \frac{2\pi}{360}, speed \frac{2\pi}{360}, u_{up} = \cos \frac{2\pi}{360}

Section 1.5 A Review of Trigonometry (page 33)

1 Connect corner to midpoint of opposite side, producing 30° angle  
3 \pi  
7 \frac{\theta}{2\pi} \rightarrow \text{area } \frac{1}{2}r^2\theta  
9 \cos(2t + t) = \cos 2t \cos t - \sin 2t \sin t = 4 \cos^3 t - 3 \cos t  
13 \cos x = \sec x = \pm 1 at \theta = n\pi  
17 \cos \theta = \sec \theta = \pm 1 at \theta = n\pi  
19 Use \cos(\frac{\pi}{2} - s - t) = \cos(\frac{\pi}{2} - s) \cos t + \sin(\frac{\pi}{2} - s) \sin t  
23 \theta = \frac{3\pi}{2} + \text{multiple of } 2\pi  
25 \theta = \frac{\pi}{4} + \text{multiple of } \pi  
27 No \theta  
29 \phi = \frac{\pi}{4}  
31 |OP| = a, |OQ| = b

CHAPTER 2 DERIVATIVES

Section 2.1 The Derivative of a Function (page 49)

1 (b) and (c)  
3 12 + 3h; 13 + 3h; 3; 3  
5 f(x) + 1  
7 -6  
9 2x + \Delta x + 1; 2x + 1  
11 \frac{4}{i + \Delta t} - \frac{4}{i} = \frac{-\frac{4}{i(\Delta t + i)}}{i(\Delta t + i)} \rightarrow \frac{-4}{i^2}  
13 7; 9; corner  
15 A = 1, B = -1  
17 F; F; T; F  
19 b = B; m and M; m or undefined  
21 Average \frac{x_2 + x_1}{2} \rightarrow 2x_1  
25 \frac{1}{2}; \text{no limit (one-sided limits } 1, -1); 1; 1 \text{ if } t \neq 0, -1 \text{ if } t = 0  
27 f'(3); f(4) - f(3)  
29 2x^4(4x^3) = 8x^7  
31 \frac{d}{dx} = \frac{1}{2u} = \frac{1}{2\sqrt{x}}  
33 \frac{d}{dx} = -\frac{1}{2}; f'(2) \text{ doesn't exist}  
35 2f \frac{d}{dx} = 4u^3 \frac{d}{dx}

Section 2.2 Powers and Polynomials (page 56)

1 6x^5; 30x^4; f^{iii} = 720 = 6!  
3 2x + 7  
5 1 + 2x + 3x^2 + 4x^3  
7 nx^{n-1} - nx^{-n-1}  
9 1 + x + \frac{1}{2}x^2 + \frac{1}{6}x^3  
11 -\frac{1}{x}, (-\frac{1}{x}) + 5  
13 \frac{x^{-2/3}}{x^{4/3}}; \frac{x^{-2/3}}{x^{-4/3}}; \frac{1}{x^2}  
15 3x^2 - 1 = 0 \text{ at } x = \frac{1}{\sqrt{3}} \text{ and } \frac{1}{\sqrt{3}}  
17 8 \text{ ft/sec; } -8 \text{ ft/sec; } 0  
19 Decreases for \(-1 < x < \frac{1}{3}\)  
21 \frac{(x + h) - x}{h/\sqrt{x + h + h/\sqrt{3}}} \rightarrow \frac{1}{2\sqrt{x}}  
23 5 10 10 5 1 \text{ adds to } (1 + 1)^5  
25 3x^2; 2h \text{ is difference of } x'\text{s}  
27 \frac{d}{dx} = 2x + \Delta x + 3x^2 + 3x\Delta x + (\Delta x)^2} \rightarrow 2x + 3x^2 \text{ is sum of separate derivatives}  
29 7x^6; (x+1)^6  
31 \frac{1}{24}x^4 \text{ plus any cubic}  
33 x + \frac{1}{2}x^2 + \frac{1}{3}x^3 + \frac{1}{4}x^4 + C  
35 \frac{1}{24}x^4, \frac{1}{120}x^5  
37 F; F; F; T; T  
39 \frac{7}{2} = 1.2 \text{ so } \frac{d}{dx} = \frac{7}{2}.12 \text{; six cents}  
41 \frac{d}{dx} = \frac{1}{\Delta x}(\frac{e}{x + \Delta x} - \frac{e}{x})  
43 E = \frac{2x}{2x + 3} \rightarrow \frac{1}{2\sqrt{2}}  
47 \frac{1}{10}x^{10}; \frac{1}{n+1}x^{n+1}; \text{ divide by } n + 1 = 0  
49 .7913, -3.7913, 1.618, -1.618, 0, 1.266, -2.766

Section 2.3 The Slope and the Tangent Line (page 63)

1 \frac{\sqrt{3}}{2}; y - 6 = -3(x - 2); y - 6 = \frac{3}{2}(x - 2); y - 6 = -\frac{3}{2}(x - 2)  
3 y + 1 = 3(x - 1); y = 3x - 4  
5 y = x; (3,3)  
7 y - a^2 = (c + a)(x - a); y - a^2 = 2a(x - a)  
9 y = \frac{1}{6}x^2 + 2; y - 7 = -\frac{1}{2}(x - 5)  
11 y = 1; x = \frac{\pi}{2}  
13 y = \frac{1}{a} - \frac{1}{a^2}(a - x); y = \frac{1}{a}, x = 2a; 2  
15 c = 0, \text{ tangent at } x = 2
17 \((-3, 19)\) and \(\left(\frac{1}{3}, \frac{13}{2}\right)\) \(e = 4, y = 3 - x\) tangent at \(x = 1\)

21 \((1 + h)^3; 3h + 3h^2 + h^3; 3 + 3h + h^2; 3\) \(23\) Tangents parallel, same normal

25 \(y = 2ax - a^2, Q = (0, -a^2)\); distance \(a^2 + \frac{1}{4}\); angle of incidence = angle of reflection

27 \(x = 2p; \) focus has \(y = \frac{x^2}{4p} = p\) \(29\) \(y - \frac{1}{\sqrt{2}} = x + \frac{1}{\sqrt{2}}; x = -\frac{2}{\sqrt{2}} = -\sqrt{2}\)

31 \(y = a^2 - \frac{1}{2a}(x - a); y = a^2 + \frac{1}{2} a = \frac{\sqrt{2}}{2} \) \(33\) \((\frac{1}{2x})(1000) = 10\) at \(x = 10\) hours

37 \(1.01004512; 1 + 10(0.001) = 1.01 \) \(39\) \((2 + \Delta x)^3 - (8 + 6\Delta x) = 6(\Delta x)^2 + (\Delta x)^3\)

41 \(x_1 = \frac{5}{4}, x_2 = \frac{41}{40}\)

43 \(T = 8\) sec; \(f(T) = 96\) meters

Section 2.4 The Derivative of the Sine and Cosine (page 70)

1 \(a\) and \(b\) \(3\) 0; 1; 5; \(\frac{1}{6}\) 5 \(\sin(x + 2\pi); (\sin h)/h \to 2\pi\) \(7\) \(\cos^2 \theta \approx 1 - \theta^2 + \frac{1}{4} \theta^4; \frac{1}{4} \theta^4\) is small

9 \(\sin \frac{1}{2} \theta \approx \frac{1}{2} \theta\) 11 \(\frac{3\pi}{4}\) 13 \(PS = \sin h; \) area \(OPR = \frac{1}{2} \sin h < \) curved area \(\frac{1}{2} h\)

15 \(\cos x = 1 - \frac{x^2}{2!} + \frac{x^4}{4!} - \cdots\) \(17\) \(\frac{1}{2h}(\cos(x + h) - \cos(x - h)) = \frac{1}{h}(\sin x \sin h) \to -\sin x\)

19 \(y' = \cos x - \sin x = 0\) at \(x = \frac{\pi}{4}\) and \(n\pi\) \(21\) \((\tan h)/h = \sin h/h \cos h < \frac{1}{\cos h} \to 1\)

23 \(\tan \frac{1}{2} x = \frac{1}{2}, 0, -\frac{1}{2}, \frac{1}{2}\); no \(25\) \(y = 2\cos x + \sin x; y'' = -y\) \(27\) \(y = -\frac{1}{3} \cos 3x; y = \frac{1}{3} \sin 3x\)

29 In degrees \((\sin h)/h \to 2x_{1360} = .01745\)

Section 2.5 The Product and Quotient and Power Rules (page 77)

1 \(2x^3 (\frac{-1}{(1+x)^2} - \frac{\cos x}{(1+\sin x)^2})\) \(5\) \((x-2)(x-3) + (x-1)(x-3) + (x-1)(x-2)\)

7 \(-2x^2 \sin x + 4x \cos x + 2 \sin x\) 9 \(2x - 1 - \frac{1}{\sin^2 x}\) \(11\) \(2 \sqrt{x} \sin x \cos x + \frac{1}{2} x^{-1/2} \sin^2 x + \frac{1}{2} (\sin x)^{-1/2} \cos x\)

13 \(4x^3 \cos x - x^4 \sin x + 4x^4 \sin x - 4x^3 \cos x\) sin \(x\) \(15\) \(\frac{1}{2} x^2 \cos x + 2x \sin x\)

17 \(0.19 - \frac{8}{3} (x-5)^{-5/3} + \frac{8}{3} (5-x)^{-5/3} (= 0?)\)

21 \(3(\sin x \cos x)^2(2 \cos^2 x - \sin^2 x) + 2 \cos 2x\) \(23\) \(u'vwx + u'uvwz + w'uvwz + x'uw\)

25 \(-\csc^2 x - \sec^2 x\)

27 \(V = \frac{t \cos t}{1+t}, V' = \frac{t \cos t - s \sin t - t^2 \sin t}{(1+t)^2}\)

\(A = 2(\frac{1}{1+t} + t \cos t + \frac{\cos t}{1+t})\)

\(A' = 2(t \cos t - t \sin t + \frac{1-\cos t}{(1+t)^2} - \sin t)\)

10 \(t\) for \(t < 10\) \(\frac{50}{\sqrt{t-10}}\) \(\) for \(t > 10\) \(31\) \(2 \sqrt{r + 3g^2 + 6t^2 + 6t} \)

33 \(u''v + 2u'v' + uv''; u'''v + 3u''v' + 3u'v'' + v''''\)

35 \(\frac{1}{2} \sin^2 t; \frac{1}{2} \tan^2 t; \frac{1}{2} [1 + t]^{3/2} - 1\)

39 \(T; F; F; F\) \(41\) degree \(2n - 1\) degree \(2n\)

43 \(v(t) = \cos t - t \sin t (t \leq \frac{\pi}{2}); v(t) = -\frac{\pi}{2} (t \geq \frac{\pi}{2})\)

45 \(y = \frac{2h^3}{L^3} + \frac{3h^2}{L} \) has \(\frac{dy}{dx} = 0\) at \(x = 0\) (no crash) and at \(x = -L\) (no dive). Then \(\frac{dy}{dx} = \frac{6v}{h} (\frac{h}{L} + \frac{1}{2})\) and \(\frac{d^2y}{dx^2} = \frac{6v^2}{L} (\frac{2h}{L} + 1)\)

Section 2.6 Limits (page 84)

1 \(\frac{1}{4}, L = 0,\) after \(N = 10; \frac{25}{24}, \) infinite \(\) no \(N; \frac{1}{4}, 0,\) after \(5; 1.1111, \frac{10}{9},\) all \(n; \sqrt{2}, 1,\) after \(38; \sqrt{20} - 4, \frac{1}{2},\) all \(n; \frac{625}{256}, e = 2.718 \ldots, \) after \(N = 12.\)

3 \((c)\) and \((d)\)

5 Outside any interval around zero there are only a finite number of \(a\)'s

7 \(\frac{5}{2}\) \(9\) \(\frac{L(h) - f(0)}{h}\) \(11\) \(1\)

13 \(15\) \(17\) \(No\) limit \(19\) \(\frac{1}{2}\) \(21\) \(Zero\) if \(f(x)\) is continuous at \(a\)

23 \(2\)

25 \(.001, .0001, .005, .1\) \(27\) \(|f(x) - L|; \frac{4x}{1+x^2}\)

29 \(0; X = 100\)

33 \(4; \) infinite \(7; 7\)

35 \(3; \) no limit; \(0; 1\)

37 \(\frac{1}{1-r} \) if \(|r| < 1;\) no limit if \(|r| \geq 1\)

39 \(.0001;\) after \(N = 7\) (or 8?) \(41\) \(\frac{1}{2}\)

43 \(9; \frac{3}{2}; a_n - 8 = \frac{1}{2}(a_{n-1} - 8) \to 0\)

45 \(a_n - L \leq b_n - L \leq c_n - L\) so \(|b_n - L| < \epsilon\) if \(|a_n - L| < \epsilon\) and \(|c_n - L| < \epsilon\)
Section 2.7 Continuous Functions  (page 89)

1 \( c = \sin 1; \) no \( c \)
3 Any \( c; \) \( c = 0 \)
5 \( c = 0 \) or 1; no \( c \)
7 \( c = 1; \) no \( c \)
9 no \( c; \) no \( c \)
11 \( c = \frac{1}{e^4}; c = \frac{1}{e^4} \)
13 \( c = -1; c = -1 \)
15 \( c = 1; c = 1 \)
17 \( c = -1; c = -1 \)
19 \( c = 2, 1, 0, -1, \cdots; \) same \( c \)
21 \( f(x) = 0 \) except at \( x = 1 \)
23 \( \sqrt{x - 1} \)
25 \( -\frac{2x}{|x|} \)
27 \( \frac{5}{x - 1} \)
29 One; two; two
31 No; yes; no
33 \( x f(x), (f(x))^2, x, f(x), 2(f(x) - x), f(x) + 2x \)
35 F; F; F; T
37 Step; \( f(x) = \sin \)
39 Yes; no; no; yes \( (f(4) = 1) \)
41 \( g(\frac{1}{2}) = f(1) - f(\frac{1}{2}) = f(0) - f(\frac{1}{2}) = -g(0); \) zero is an intermediate value between \( g(0) \) and \( g(\frac{1}{2}) \)
43 \( f(x) - x \) is 20 at \( x = 0 \) and 50 at \( x = 1 \)

CHAPTER 3 APPLICATIONS OF THE DERIVATIVE

Section 3.1 Linear Approximation  (page 95)

1 \( Y = x \)
3 \( Y = 1 + 2(x - \frac{\pi}{4}) \)
5 \( Y = 2\pi(x - 2\pi) \)
7 \( 2^6 + 6 \cdot 2^5 \cdot .001 \)
9 \( 1 \)
11 \( 1(-.02) = 1.02 \)
13 Error .000301 vs. \( \frac{1}{2} (.0001)^6 \)
15 .0001 - \( \frac{1}{3} 10^{-8} \) vs. \( \frac{1}{2} (.0001)(2) \)
17 Error .59 vs. \( \frac{1}{2} (.01)(90) \)
19 \( \frac{d}{du} \sqrt{2^2 + u} = \frac{1}{2}{\sqrt{2^2 + u}} = \frac{1}{2c} \) at \( u = 0 \), \( c + \frac{u}{2c} = c + \frac{x^2}{2c} \)
23 \( dV = 3(10)^2(1) \)
25 \( A = 4\pi r^2, dA = 8\pi r \) \( dr \)
27 \( V = \pi r^2 h, dV = 2\pi rh \) \( dr \) (plus \( \pi r^2 dh \))
29 \( 1 + \frac{1}{2}x \)
31 32nd root

Section 3.2 Maximum and Minimum Problems  (page 103)

1 \( x = -2: \) abs min
3 \( x = -1: \) rel max, \( x = 0: \) abs min, \( x = 4: \) abs max
5 \( x = -1: \) abs max, \( x = 0, 1: \) abs min, \( x = \frac{1}{2}: \) rel max
7 \( x = -3: \) abs min, \( x = 1: \) rel max, \( x = 0: \) stationary (not min or max)
9 \( x = 1, 9: \) abs min, \( x = 5: \) abs max
11 \( x = \frac{1}{2}: \) rel max, \( x = 1: \) rel min, \( x = 0: \) stationary (not min or max)
13 \( x = 0, 1, 2, \cdots: \) abs min, \( x = \frac{1}{2}, \frac{3}{2}, \frac{5}{2}, \cdots: \) abs max
15 \( |x| \leq 1: \) all min, \( x = -3 \) abs max, \( x = 2 \) rel max
17 \( x = 0: \) rel min, \( x = \frac{1}{3}: \) abs max, \( x = 4: \) abs min
19 \( x = 0: \) abs min, \( x = \pi: \) stationary (not min or max), \( x = 2\pi: \) abs max
21 \( \theta = 0: \) rel min, \( \tan \theta = -\frac{4}{3} \) \( (\sin \theta = \frac{4}{5} \) and \( \cos \theta = -\frac{3}{5} \) abs max, \( \sin \theta = -\frac{4}{5} \) and \( \cos \theta = \frac{3}{5} \) abs min), \( \theta = 2\pi: \) rel max
23 \( h = \frac{1}{2}(62'' \) or 158 cm); cube
25 \( \frac{25}{3} 2\sqrt{ab} \) gallons/mile, \( \frac{3}{2} \sqrt{ab} \) miles/gallon at \( v = \sqrt{\frac{k}{a}} \)
27 \( \theta = \frac{3\pi}{6} = 67.5^0 \)
29 \( x = \frac{a}{\sqrt{5}} \); compare Example 7; \( \frac{a}{\sqrt{3}} = \sqrt{5} \)
31 \( R(x) - C(x); \frac{(R(x) - C(x))}{x}; \frac{dR}{dx} \)
33 \( x = \frac{d-a}{2(b-c)}; \) zero
35 \( x = 2 \)
37 \( V = x(6 - \frac{3x}{2})(12 - 2x); x \approx 1.6 \)
39 \( A = \pi r^2 + x^2, x = \frac{1}{4}(4 - 2\pi r); r_{\text{min}} = \frac{2}{2+\pi} \)
41 max area 2500 vs 10000 = 3185
43 \( x = 2, y = 3 \)
45 \( P(x) = 12 - x; \) thin rectangle up y axis
47 \( h = \frac{3}{4}, r = \frac{3}{4}, V = \frac{4}{3}\pi R^2 \frac{3}{4} = \frac{1}{9} \) of cone volume
49 \( r = \frac{HR}{2(H-R)}; \) best cylinder has no height, area \( 2\pi R^2 \) from top and bottom (?)
51 \( r = 2, h = 4 \)
53 25 and 0
55 8 and \(-\infty\)
57 \( \sqrt{r^2 + x^2} + \sqrt{q^2 + (s-x)^2}; \frac{dq}{dx} = \frac{2}{\sqrt{r^2 + x^2}} - \frac{s-x}{\sqrt{q^2 + (s-x)^2}} = 0 \) when \( \sin a = \sin c \)
59 \( y = x^2 = \frac{3}{2} \)
61 \( (1, -1), (\frac{13}{5}, \frac{1}{3}) \)
63 \( m = 1 \) gives nearest line
65 \( m = \frac{1}{3} \)
67 equal; \( x = \frac{1}{2} \)
69 \( x \) equal
71 True (use sign change of \( f'' \))
73 Radius \( R, \) swim \( 2R \cos \theta, \) run \( 2R \theta, \) time \( \frac{2R \cos \theta}{v} + \frac{2R \theta}{10v}; \) max when \( \sin \theta = \frac{1}{10}, \) min all run
Section 3.3 Second Derivatives: Bending and Acceleration  (page 110)

3 \[ y = -1 - x^2; \text{ no} \] 5 False 7 True 9 True \( f' \) has 8 zeros, \( f'' \) has 7
11 \[ x = 3 \text{ is min: } f''(3) = 2 \] 15 \[ x = \frac{3\pi}{4} \text{ is max: } f''(\frac{3\pi}{4}) = -\sqrt{2}; x = \frac{7\pi}{4} \text{ is min: } f''(\frac{7\pi}{4}) = \sqrt{2} \]
17 \( x = 3 \text{ is max: } f''(3) = -4; x = 2, 4 \text{ are min but } f'' = 0 \] 21 \( f(\Delta x) = f(-\Delta x) \] 23 \( 1 + x - \frac{x^2}{2} \)
25 \( 1 - \frac{x^2}{6} \) 27 \( 1 - \frac{1}{2} x - \frac{1}{8} x^2 \) 29 Error \( \frac{1}{2} f''(x) \Delta x \) 31 Error \( 0 \Delta x + \frac{1}{6} f'''(x)(\Delta x)^2 \)
37 \( \frac{1}{39} = 1.010101; \frac{1}{1.1} = .90099 \) 39 Inflection 41 18 vs. 17 43 Concave up; below

Section 3.4 Graphs  (page 119)

1 120; 150; 60 3 Odd; \( x = 0, y = x \) 5 Even; \( x = 1, x = -1, y = 0 \) 7 Even; \( y = 1 \) 9 Even
11 Even; \( x = 1, x = -1, y = 0 \) 13 \( x = 0, x = -1, y = 0 \) 15 \( x = 1, y = 1 \) 17 Odd 19 \( \frac{2x}{x-1} \)
21 \( x + \frac{1}{4} \) 23 \( \sqrt{x^2 + 1} \) 25 Of the same degree 27 Have degree \( P < \) degree \( Q \); none
29 \( x = 1 \) and \( y = 3x + C \) if \( f \) is a polynomial; but \( f(x) = (x - 1)^{1/3} + 3x \) has no asymptote \( x = 1 \)
31 \( (x - 3)^2 \) 39 \( x = \sqrt{2}, x = -\sqrt{2}, y = x \) 41 \( Y = 100 \sin \frac{3x}{30} \) 45 \( c = 3, d = 10; c = 4, d = 20 \)
47 \( x^* = \sqrt{5} = 2.236 \) 49 \( y = x - 2; Y = X; y = 2x \) 51 \( x_{\text{max}} = .281, x_{\text{min}} = 6.339; x_{\text{inf}} = 4.724 \)
53 \( x_{\text{min}} = .393, x_{\text{max}} = 1.53, x_{\text{min}} = 3.33; x_{\text{inf}} = .896, 2.604 \)
55 \( x_{\text{min}} = -1.793, x_{\text{max}} = .8135; x_{\text{inf}} = .04738; x_{\text{blowup}} = \pm 2.38 \) 57 8 digits

Section 3.5 Parabolas, Ellipses, and Hyperbolas  (page 128)

1 \( \frac{dy}{dx} = 0 \) at \( \frac{b}{2a} \) 3 \( V = (1, -4), F = (1, -3.75) \) 5 \( V = (0,0), F = (0, -1) \) 7 \( F = (1,1) \)
9 \( V = (0, \pm 3); F = (0, \pm \sqrt{8}) \) 11 \( V = (0, \pm 1); F = (0, \pm \sqrt{\frac{5}{4}}) \) 13 Two lines, \( a = b = c = 0; V = F = (0,0) \)
15 \( y = 5x^2 - 4x \) 17 \( y + p = \sqrt{x^2 + (y - p)^2 - 4py} = x^2; F = (0, \frac{1}{12}), y = -\frac{12}{12}; (\pm \frac{11}{6}, \frac{11}{12}) \)
19 \( x = ay^2 \) with \( a > 0 \); \( y = \frac{(a+y)^2}{4p}; y = -ax^2 + ax \) with \( a > 0 \)
21 \( \frac{a^2}{4} + y^2 = 1; (x-1)^2 + (y-1)^2 = 1 \) 23 \( \frac{a^2}{25} + \frac{y^2}{9} = 1; \frac{(x-3)^2}{36} + \frac{(y-1)^2}{32} = 1; x^2 + y^2 = 25 \)
25 Circle, hyperbola, ellipse, parabola 27 \( \frac{dy}{dx} = -\frac{4}{5}; y = -\frac{4}{5}x + 5 \) 29 \( 5\frac{4}{4}; \frac{40}{41} = \frac{1}{2}(\frac{5}{4} - \frac{4}{5}) \)
31 Circle; \( (3, 1); 2; X = \frac{2-x}{2}, Y = \frac{y-1}{2} \) 33 \( 3x^2 + y^2 = 2 \) 35 \( y^2 - \frac{1}{3}x^2 = 1; \frac{y^2}{9} - \frac{4x^2}{9} = 1; y^2 - x^2 = 5 \)
37 \( \frac{x^2}{25} - \frac{y^2}{30} = 1 \) 39 \( y^2 - 4y + 4, 2x^2 + 12x + 18; -14, (-3,2), \text{ right-left} \)
41 \( F = (\pm \sqrt{\frac{2}{2}}, 0); y = \pm \sqrt{2} \) 43 \( x + y + 1)^2 = 0 \)
45 \( (a^2 - 1)x^2 + 2abby + (b^2 - 1)y^2 + 2acx + 2bcy + c^2 = 0; 4(a^2 + b^2 - 1); \text{ if } a^2 + b^2 < 1 \) then \( B^2 - 4AC < 0 \)

Section 3.6 Iterations \( x_{n+1} = F(x_n) \)  (page 136)

1 \( -1.366; \infty \) 3 \( 1; 1 \) 5 \( \frac{2}{3}; \infty \) 7 \( -2; -2 \)
9 \( \frac{1+\sqrt{3}}{2} \) attracts, \( \frac{1+\sqrt{3}}{2} \) repels; \( \frac{1}{2} \) attracts, 0 repels; \( 1 \) attracts, 0 repels; \( 1 \) attracts, 0 repels; \( \pm \sqrt{2} \) repel
11 Negative 13 .900 15 .679 17 \( |a| < 1 \) 19 Unstable \( |F'| > 1 \) 21 \( x^* = \frac{x}{1-a}; |a| < 1 \)
Answers to Odd-Numbered Problems


Section 3.7 Newton's Method and Chaos (page 145)

1 $2000; $2000 5 $2000; $2000


31 $2000; $2000 33 $2000; $2000

43 $2000; $2000

Section 3.8 The Mean Value Theorem and l'Hôpital's Rule (page 152)

1 $2000; $2000 3 $2000; $2000


31 $2000; $2000 33 $2000; $2000

41 $2000; $2000 43 $2000; $2000

Chapter 4 DERIVATIVES BY THE CHAIN RULE

Section 4.1 The Chain Rule (page 158)

1 $2000; $2000 3 $2000; $2000

5 $2000; $2000 7 $2000; $2000


41 $2000; $2000 43 $2000; $2000


Section 4.2 Implicit Differentiation and Related Rates (page 163)

1. \(-x^{-1}/y^{n-1}\)  
2. \(3 \frac{dy}{dx} = 1\)  
3. \(5 \frac{dy}{dx} = \frac{1}{F(y)}\)  
4. \((y^2 - 2xy)/(x^2 - 2xy)\) or \(9 \frac{1}{\sec^2 y}\) or \(\frac{1}{1+x^2}\)  
11. First \(\frac{dy}{dx} = -\frac{y}{x}\), second \(\frac{dy}{dx} = \frac{y}{x}\)  
13. Faster, faster  
15. \(2\frac{z^3}{} = 2y' \Rightarrow x' = \frac{y'y'}{y} = y'\sin \theta\)  
17. \(\sec^2 \theta = \frac{c}{200\pi}\)  
19. \(500\frac{dV}{dt} = 500\sqrt{1 + \left(\frac{dV}{dx}\right)^2}\)  
21. \(\frac{dy}{dt} = -\frac{8}{3} \frac{dy}{dt} = -2\sqrt{3}; \infty \text{ then } 0\)  
23. \(V = \pi r^2h; \frac{dh}{dt} = \frac{1}{4\pi} \frac{dV}{dt} = -\frac{1}{4\pi} \text{ in/sec}\)  
25. \(A = \frac{1}{2} \int \sin \theta \frac{dA}{dt} = 7\)  
27. \(1.6 \text{ m/sec}; 9 \text{ m/sec}; 12.8 \text{ m/sec}\)  
29. \(-\frac{7}{5}\)  
31. \(\frac{dx}{dt} = \sqrt{2} \frac{dx}{dt}; \frac{dy}{dt} = \frac{1}{10} \cos^2 \theta \frac{dθ}{dt}; \theta'' = \frac{\cos \theta y'}{10} - \frac{1}{50} \cos^3 \theta \sin \theta (y')^2\)

Section 4.3 Inverse Functions and Their Derivatives (page 170)

1. \(x = \frac{y+6}{3}\)  
3. \(x = \sqrt{y+1} \text{ (x unrestricted \rightarrow no inverse)}\)  
5. \(x = \frac{1}{y-1}\)  
7. \(x = (1+y)^{1/3}\)  
9. \((x \text{ unrestricted} \rightarrow \text{no inverse)}\)  
11. \(y = \frac{1}{x-a}\)  
13. \(2 < f^{-1}(x) < 3\)  
15. \(f \text{ goes up and down}\)  
17. \(f(x)g(x)\) and \(\frac{1}{f(x)}\)  
19. \(m \neq 0; m \geq 0; |m| \geq 1\)  
21. \(\frac{dy}{dx} = 5x^4, \frac{dz}{dy} = \frac{1}{5} y^{-4/5}\)  
23. \(\frac{dy}{dx} = 3x^2, \frac{dz}{dy} = \frac{1}{3}(1+y)^{-2/3}\)  
25. \(\frac{dz}{dy} = \frac{1}{(x-1)^3}, \frac{dx}{dy} = \frac{1}{(y-1)^3}\)  
27. \(y; \frac{1}{2}y^2 + C\)  
29. \(f(g(x)) = -1/3x^3; g^{-1}(y) = y^3; g(-1)(x) = x\)  
39. \(2/\sqrt{3}\)  
41. \(1/6\cos 9\)  
43. Decreasing: \(\frac{dy}{dx} < 0\)  
45. \(F; T; F\)  
47. \(g(x) = x^m, f(y) = y^n, x = (z^{1/n})^{1/m}\)  
49. \(g(x) = x^3, f(y) = y+6, x = (x-6)^{1/3}\)  
51. \(g(x) = 10^x, f(y) = \log y, x = \log(10^y) = y\)  
53. \(y = x^3, y' = 6x, d^2x/dy^2 = -\frac{2}{3}y^{-5/3}; m/\text{sec}^2, \text{sec}^2/m^2\)  
55. \(p = \frac{1}{\sqrt{y} - 1}; 0 < y \leq 1\)  
57. \(\max = G = \frac{3}{8}y^{4/3}, G' = \frac{1}{2}y^{1/3}\)  
59. \(y^2/100\)  

Section 4.4 Inverses of Trigonometric Functions (page 175)

1. \(0, \frac{\pi}{2}, 0\)  
3. \(\frac{\pi}{2}, 0, \frac{\pi}{4}\)  
5. \(\pi \text{ is outside } [-\frac{\pi}{2}, \frac{\pi}{2}]\)  
7. \(y = -\sqrt{3}/2 \text{ and } \sqrt{3}/2\)  
9. \(\sin x = \sqrt{1 - y^2}; \sqrt{1 - y^2} \text{ and } 1\)  
11. \(\frac{d(sin^{-1}y)}{dy} \cos x = 1 \rightarrow \frac{d(sin^{-1}y)}{dy} = \frac{1}{\cos x} = \frac{1}{\sqrt{1-y^2}}\)  
13. \(y = 0; 1, -1, 1; y = 1: 0, 0, \frac{1}{2}\)  
15. \(F; F; T; T; F; F\)  
17. \(\frac{dy}{dx} = \frac{1}{1-x^2}\)  
19. \(\frac{dx}{dz} = 3\)  
21. \(\frac{dz}{dx} = \frac{2\sin^{-1}x}{1-x^2}\)  
23. \(1-\frac{1}{\sin^{-1}y}\)  
25. \(\frac{dz}{dx} = \frac{1}{|y+1|\sqrt{y^2+2y}}\)  
27. \(u = 1 \text{ so } \frac{dy}{dx} = 0\)  
31. \(\sec x = \sqrt{y^2 + 1}\)  
33. \(\frac{1}{10}, 1, \frac{1}{2}\)  
35. \(-y/\sqrt{1-y^2}\)  
37. \(\frac{1}{2} \sec \frac{\pi}{2}\)  
39. \(\frac{mz^{-n-1}}{m^n\sqrt{z^{2n-1}}}\)  
41. \(\frac{dy}{dz} = \frac{1}{1+z^2}\)  
43. \(\frac{dy}{dx} = \pm \frac{1}{1+x^2}\)  
47. \(u = 4\sin^{-1}y\)  
49. \(\pi\)  
51. \(-\pi/4\)  

CHAPTER 5 INTEGRALS

Section 5.1 The Idea of the Integral (page 181)

1. \(1,3,7,15,127\)  
3. \(-\frac{1}{2} - \frac{1}{4} - \frac{1}{8} = \frac{1}{8} - 1\)  
5. \(f_j - f_0 = \frac{r_{j-1}}{r-1}\)  
7. \(3z\text{ for } x \leq 7, 7z - 4 \text{ for } x \geq 1\)  
9. \(\frac{1}{6}, \frac{1}{52}, \frac{1}{72}, \frac{1}{52}\)  
11. \(\text{Lower by 2}\)  
13. \(\text{Up, down; rectangle}\)  
15. \(\sqrt{x + \Delta x - \sqrt{z}; \Delta z; \frac{df}{dz}; \sqrt{z}}\)  
17. \(6; 18; \text{triangle}\)  
19. \(18 \text{ rectangles}\)  
21. \(6x - \frac{1}{2}x^2 - 10; 6 - x\)  
23. \(\frac{14}{27}\)  
25. \(x^2 - x^2; \frac{1}{3}x^3\)
Section 5.2 Antiderivatives (page 186)

1 \(x^5 + \frac{2}{3}x^6; \frac{5}{3}\)  3 \(2\sqrt{x}; 2\)  5 \(\frac{2}{3}x^{4/3}(1 + 2^{1/3}); \frac{3}{2}(1 + 2^{1/3})\)  7 \(-2\cos x - \frac{1}{3}\cos 2x; \frac{5}{2} - 2\cos 1 - \frac{1}{2}\cos 2\)
9 \(x\sin x + \cos x; \sin 1 + \cos 1\)  11 \(\frac{1}{2}\sin^2 x; \frac{1}{2}\sin^2 1\)  13 \(f = C; 0\)  15 \(f(b) - f(a); f_7 - f_2\)
17 \(8 + \frac{5}{N}\)  19 \(\frac{\pi}{3}(1 + \sqrt{3}); \frac{\pi}{2}(3 + \sqrt{3}); 2\)  21 \(\frac{205}{36}i\infty\)  23 \(f(x) = 2\sqrt{x}\)  25 \(\frac{1}{2}, \text{below} -1; \frac{1}{4}, \frac{5}{4}\)
27 Increase - decrease; increase - decrease - increase
29 Area under \(B\) - area under \(D\); time when \(B = D\); time when \(B - D\) is largest

Section 5.3 Summation Versus Integration (page 194)

1 \(\frac{25}{12}; 16\)  3 \(127; 2^{n+1} - 1\)  5 \(\sum_{j=1}^{50} 2j = 2550; \sum_{j=1}^{100} (2j - 1) = 10,000; \sum_{k=1}^{4} (-1)^{k+1}/k = \frac{7}{12}\)
7 \(\sum_{k=0}^{n} a_k x^k; \sum_{j=1}^{n} \sin \frac{2\pi j}{n}\)  9 \(5.18738; 7.48547\)  11 \(2(a_1^2 + b_1^2)\)  13 \(2^n - 1; \frac{1}{11} - \frac{1}{1}\)  15 \(F; T\)
17 \(\frac{df}{dx} + C; f_0 - f_b - f_1 + f_0\)  19 \(f_1 = 1; n^2 + (2n + 1) = (n + 1)^2\)
21 \(a + b + c = 1, 2a + 4b + 8c = 5, 3a + 6b + 12c = 7; \text{sum of squares}\) 23 \(S_{400} = 80200; E_{400} = .0025 = \frac{1}{n}\)
25 \(S_{100,1/3} \approx 350, E_{100,1/3} \approx .00587; S_{100,3} = 25502500, E_{100,3} = .0201\)  27 \(v_1\) and \(v_2\) have the same sign
29 \(v_1 = 9, v_2 = 12, \Sigma v = 21\)  31 \(\text{At } N = 1, 2^{N-2}\) is not 1  33 \(0; \frac{1}{n}(v_1 + \cdots + v_n)\)
35 \(\Delta x \sum_{j=1}^{n} v(j\Delta x)\)  37 \(f(1) - f(0) = \int_{0}^{1} \frac{df}{dx} dx\)

Section 5.4 Indefinite Integrals and Substitutions (page 200)

1 \(\frac{2}{3}(2 + x)^{3/2} + C\)  3 \((x + 1)^{n+1}/(n+1) + C(x \neq -1)\)  5 \(\frac{12}{12}(x^2 + 1)^6 + C\)  7 \(-\frac{1}{4}\cos^4 x + C\)
9 \(-\frac{1}{8}\cos^4 2x + C\)  11 \(\sin^{-1} t + C\)  13 \(\frac{1}{3}(1 + t^2)^{3/2} - (1 + t^2)^{1/2} + C\)  15 \(2\sqrt{x} + x + C\)
17 \(\sec x + C\)  19 \(-\cos x + C\)  21 \(\frac{1}{3}x^3 + \frac{2}{3}x^{3/2}\)  23 \(-\frac{1}{3}(1 - 2x)^{3/2}\)  25 \(y = \sqrt{2x}\)
27 \(\frac{1}{3}x^2\)  29 \(a\sin x + b\cos x\)  31 \(\frac{4}{15}x^{5/2}\)  33 \(F; F; F; F; F; F\)
35 \(f(x - 1); 2f\left(\frac{x}{2}\right)\)
37 \(x - \tan^{-1} x\)  39 \(\int_{1/4}^{1} du\)  41 \(4.9t^2 + C_1 t + C_2\)  43 \(f(t + 3); f(t) + 3t; 3f(t); \frac{1}{3}f(3t)\)

Section 5.5 The Definite Integral (page 205)

1 \(C = -f(2)\)  3 \(C = f(3)\)  5 \(f(t)\) is wrong  7 \(C = 0\)  9 \(C = f(-a) - f(-b)\)
11 \(u = x^2 + 1; \int_{1/2}^{x} u^{10} du = \frac{u^{11}}{11}|_{\frac{1}{2}} = \frac{2^{11} - 1}{22}\)  13 \(u = \tan x; \int_{0}^{1} u du = \frac{1}{2}\)
15 \(u = \sec x; \int_{\frac{1}{2}}^{\sqrt{2}} u du = \frac{1}{2}\) (same as 13)  17 \(u = \frac{1}{x}; x = \frac{1}{u}; dx = -\frac{du}{u}; \int_{1/2}^{1} \frac{du}{u}\)
19 \(S = \frac{1}{2}\left(\frac{1}{4} + 1\right)^4 + \frac{1}{2}\left(1 + 1\right)^4; s = \frac{1}{2}(0) + \frac{1}{2}(\frac{1}{4} + 1)^4\)
21 \(S = \frac{1}{2}\left[\left(\frac{1}{2}\right)^3 + 1^3 + \left(\frac{3}{2}\right)^3 + 2^3\right]; s = \frac{1}{2}[0^3 + (\frac{1}{2})^3 + 1^3 + (\frac{3}{2})^3]\)
23 \(S = \frac{1}{4}\left[\left(\frac{17}{16}\right)^4 + (\frac{5}{4})^4 + \left(\frac{30}{16}\right)^4 + 2^4\right]\)  25 \(\text{Last rectangle minus first rectangle}\)
27 \(S = .07\) since 7 intervals have points where \(W(x)\) exists and equals zero.
29 \(M\) is increasing so Problem 25 gives \(S - s = \Delta x(1 - 0); \text{area from graph up to } y = 1 = \frac{1}{2} \cdot 1 + \frac{1}{2} \cdot \frac{1}{2} + \cdots = \frac{1}{2}(1 + \frac{1}{4} + \frac{1}{16} + \cdots); \text{area under graph is } \frac{1}{2}\).
31 \(f(x) = 3 + \int_{0}^{x} v(x) dx; f(x) = \int_{0}^{x} v(x) dx\)
33 \(T; T; F; T; F; T\)
Section 5.6 Properties of the Integral and Average Value (page 212)

1 \( v = \frac{1}{2} \int_{-1}^{1} x^4 \, dx = \frac{1}{8} \) equals \( c^4 \) at \( c = \pm \left( \frac{1}{2} \right)^{1/4} \)
3 \( v = \frac{1}{\pi} \int_{0}^{\pi} \cos^2 x \, dx = \frac{1}{2} \) equals \( \cos^2 c \) at \( c = \frac{\pi}{4} \) and \( \frac{3\pi}{4} \)
5 \( v = \int_{0}^{1} \frac{dx}{x} = \frac{1}{2} \) equals \( \frac{1}{4} \) at \( c = \sqrt{2} \)
7 \( \int_{3}^{5} v(x) \, dx \)
9 False, take \( v(x) < 0 \)
11 True; \( \frac{1}{3} \int_{0}^{1} v(x) \, dx + \frac{2}{3} \cdot \frac{1}{2} \int_{0}^{1} v(x) \, dx = \frac{1}{3} \int_{0}^{1} v(x) \, dx \)
13 False; when \( v(x) = x^2 \) the function \( x^2 - \frac{1}{3} \) is even
15 False; take \( v(x) = 1 \); factoring is missing
17 \( \int_{1}^{2} v(x) \, dx \)
19 \( \int_{1}^{2} v(x) \, dx \)
21 \( v(x) = Cx^2; \) \( v(x) = C \). This is "constant elasticity" in economics (Section 2.2)
23 \( V \to 0, V \to 1 \)
25 \( \frac{1}{2} \int_{0}^{a} (a - x) \, dx = a + 1 \) if \( a > 2; \)
27 \( \frac{1}{2} \int_{0}^{a} |a - x| \, dx = \frac{a^2}{2} - a + 1 \) if \( a < 2; \) distance = absolute value
29 Area under \( \cos \theta \) is 1. Rectangle \( 0 < \theta < \frac{\pi}{2}, 0 \leq y \leq 1 \) has area \( \frac{\pi}{2} \). Chance of falling across a crack is \( \frac{1}{\pi/2} = \frac{2}{\pi} \).
31 \( \frac{1}{6} \times 3, \frac{3}{6}, \ldots, \frac{1}{6}; \) \( \frac{10.5}{
33 \int_{0}^{\pi} \sin^2 \theta \, d\theta \)
35 Any \( v(x) = v_{\text{even}}(x) + v_{\text{odd}}(x); \) \( (a + 1)^2 = (3x^2 + 1) + (x^3 + 3x); \)
37 \( 16 \) per class; \( \frac{6}{64}; E(x) = \frac{1800}{64} = \frac{225}{8} \)
39 \( f(5) - f(0) = \frac{9}{2} + \frac{4}{2} = \frac{13}{2} \)

Section 5.7 The Fundamental Theorem and Its Applications (page 219)

1 \( \int_{0}^{\pi} \cos^2 x \, du \)
5 \( (x^2)^3(2x) = 2x^7 \)
7 \( v(x + 1) - v(x) \)
9 \( \sin^2 x - \frac{1}{4} \int_{0}^{\pi} \sin^2 t \, dt \)
11 \( \int_{0}^{\pi} u(x) \, du \)
13 \( 0 \)
15 \( 2 \sin x^2 \)
17 \( u(x) v(x) \)
19 \( \sin^{-1}(\sin x) \cos x = x \cos x \)
21 False; \( \int_{0}^{\pi} \cos^2 t \, dt = \frac{1}{4} \cdot 220 \cdot \frac{60}{2\pi} \sin \frac{2\pi}{60} = V_{\text{ave}} \)
23 Taking derivatives \( v(x) = (x \cos x)' = \cos x - x \sin x \)
25 Taking derivatives \( -v(-x)(-1) = v(x) \) so \( v \) is even
27 \( f; T; T; F \)
29 \( \int_{0}^{\pi} v(t) \, dt = \int_{0}^{\pi} v(t) \, dt - \int_{0}^{\pi} v(t) \, dt = \frac{x}{x+2} - \frac{1}{1+2} \) (in revised printing)
31 \( V = s^3; A = 3s^2; \) half of hollow cube; \( \Delta V \approx 3s^2 dS; 3s^2 \) (which is \( A \))
33 \( dH/dr = 2\pi^2 r^3 \)
35 Wedge has length \( r \approx \) height of triangle; \( \int_{0}^{\pi/2} \frac{1}{2} r^2 \, d\theta = \frac{\pi r^2}{4} \)
37 \( r = \frac{1}{\cos \theta}; \frac{d\theta}{2\cos^2 \theta}; \int_{0}^{\pi/4} \frac{d\theta}{2\cos^2 \theta} = \frac{1}{2} \tan \frac{\pi}{4} \)
39 \( x = y^2; \int_{0}^{a} y^2 \, dy = \frac{a^3}{3} = \frac{3}{3}; \) vertical strips have length \( 2 - \sqrt{x} \)
41 Length \( \sqrt{2}a; \) width \( \frac{d\theta}{\sqrt{2}}; \int_{0}^{\pi} a \, da = \frac{1}{2} \)
43 The differences of the sums \( f_j = v_1 + v_2 + \cdots + v_j \) are \( f_j - f_{j-1} = v_j \)
45 No, \( \int_{0}^{\pi} a(t) \, dt = \frac{df}{dx}(x) - \frac{df}{dx}(0) \) and \( \int_{0}^{1} (f(x) a(t) \, dt) \, dx = f(1) - f(0) - \frac{df}{dx}(0) \)

Section 5.8 Numerical Integration (page 226)

1 \( \frac{1}{2} \Delta x(u_0 - u_n) \)
3 \( 1 \times 0.5625, 0.3025; 0, 0.0625, 0.2025 \)
5 \( L_8 \approx 0.1427, T_8 \approx 0.2052, S_8 \approx 0.0000 \)
7 \( p = 2; \) for \( y = x^2, \frac{1}{4} \cdot \frac{1}{2} - \frac{1}{4} \cdot \frac{1}{2} + \frac{1}{2} \cdot \frac{1}{2} + \frac{1}{2} \cdot \frac{1}{2} - \frac{1}{2} \cdot \frac{1}{2} \neq \frac{1}{4} \)
9 For \( y = x^2, \) error \( \frac{1}{6} (\Delta x)^2 \) from \( \frac{1}{2} - \frac{1}{3}, y' = 2\Delta x \)
13 8 intervals give \( \frac{(0.2)^2}{12} \cdot \frac{[-1 + \frac{1}{3}]}{12} = \frac{1}{1024} < 0.001 \)
15 \( f''(c) = y'(c) \)
17 \( \infty; \) \( 0.683, 0.749, 0.772 \to \frac{\pi}{4} \)
19 \( A + B + C = 1, \frac{1}{2} B + C = \frac{1}{2}, \frac{1}{3} B + C = \frac{1}{2} \); Simpson
21 \( y = 1 \) and \( x \) on \( [0,1]; \) \( L_n = 1 \) and \( \frac{1}{2} - \frac{1}{2n} \), \( R_n = 1 \) and \( \frac{1}{2} + \frac{1}{2n} \), so only \( \frac{1}{2} L_n + \frac{1}{2} R_n \) gives 1 and \( \frac{1}{2} \)
23 \( T_{10} \approx 500,000,000; T_{100} \approx 50,000,000; 25,000\pi \)
25 \( a = 4, b = 2, c = 1; \) \( \int_{0}^{1} (4x^2 + 2x + 1) \, dx = \frac{10}{3} \); Simpson fits parabola
27 \( \frac{1}{4220} \)
CHAPTER 6  EXPONENTIALS AND LOGARITHMS

Section 6.1 An Overview  (page 234)

1  5; -5; -1; $\frac{1}{3}$; 2
  5 1; -10; 80; 1; 4; -1
  7 $n \log_b x$  9 $\frac{3}{10}$; $\frac{3}{1}$  13 105
  15 0; $I_{EF} = 10^7 I_0$; 8.3 + log10 4
  17 $A = 7$, $b = 2.5$  19 $A = 4, k = 1.5$
  21 $\frac{1}{x^2}$; $\frac{2}{x^2}$; log 2
  25 $y - 1 = cx; y - 10 = c(x - 1)$
  27 $y'' = c^2 b x^2, x'' = -1/cy^2$

Section 6.2 The Exponential $e^x$  (page 241)

1 $49e^{7x}$  3 $8e^{3x}$  5 $3e^x$  7 $(\frac{2}{3})x \ln 2 + 3 \frac{e^x}{1+e^x}$
  9 $\frac{e^x}{1+e^x} 11 2$  13 $xe^x$
  15 $(e^x + e^{-x})^2$
  17 $e^{ln x} \cos x + e^x \cos e^x$
  19 $1.246, .0135$, .0014 are close to $\frac{x}{2n}$  21 $\frac{1}{x}; \frac{1}{e}$
  23 $Y(h) = 1 + \frac{1}{10}, Y(1) = (1 + \frac{1}{10})^{10}$
  25 $(1 + \frac{1}{2})^x < e < e^x < e^{3x/2} < e^{2x} < 10^x < x^x$
  27 $\frac{x^2}{3} + \frac{x^3}{7}$
  29 $x + \frac{x^2}{2n} + \frac{x^3}{3n}$
  31 $(\frac{e^x}{1+e^x}) + 2x^2$
  33 $\frac{x^2}{2} - \frac{x^3}{2}$
  35 $2e^{x^2} + \frac{x^3}{2}$
  37 $e^{-x} < 0$ faster at $x = 0$ (slope -1); meet at $x = 1$; $e^{-x}/e^{-x} < e^{-9}/e^{-3} < \frac{1}{100}$ for $x > 3$
  39 $y - e^a = e^a(x - a)$; need $-e^a = -ae^a$ or $a = 1$
  41 $y' = x^2(x + 1) = 0$ at $x_{min} = \frac{1}{e}; y'' = x^2[ln(x + 1)^2 + \frac{1}{x}] > 0$
  43 $\frac{d}{dx}(e^{-x} y) = e^{-x} \frac{dy}{dx} e^{-x} y = 0$ so $e^{-x} y = Constant$ or $y = Ce^x$
  45 $\frac{e^x}{x^2} = 0$
  47 $\frac{e^x}{x^2} = 0$
  49 $e^{-x}|_{x = 0} = 1$
  51 $e^{1+x}|_{x = 0} = 1$
  53 $\frac{e^{ln x}}{ln x} = 0$
  55 $\int e^{u} \frac{du}{dx} = -e^{-u} + C$
  57 $yy' = 1$ gives $\frac{1}{2}y^2 = x + C$ or $y = \sqrt{2x + 2C}$
  59 $\frac{dF}{dx} = (n - x)x^{n-1}/e^x < 0$ for $x > n$; $F(2x) < \frac{constant}{e^x} \rightarrow 0$

Section 6.3 Growth and Decay in Science and Economics  (page 250)

1 $t^2 + y_0$  3 $y_0e^{2t}$
  5 $10e^{4t}; t = \frac{ln 10}{4}$
  7 $\frac{1}{4} e^{4t} + 9.75; t = \frac{ln 361}{4}$
  11 $c = \frac{ln 2}{2}; t = \frac{ln 10}{c}$
  13 $\frac{5568}{-7} \ln \left(\frac{1}{2}\right)$
  15 $c = \frac{ln 2}{20}; t = \frac{1}{c} \ln \left(\frac{5}{3}\right)$
  17 $t = \frac{ln(1/240)}{ln(1/98)}$
  19 $e^c = 3$ so $y_0 = e^{-3} 1000 = 1000$
  21 $p = 1013 e^{ch}; 50 = 1013 e^{20c}; c = \frac{1}{20} \ln (\frac{50}{1013})$
  23 $c = \frac{ln 2}{3}; \frac{1}{2} = \frac{1}{8}$
  25 $y = y_0 - at, y(t) = \frac{y_0 - y_1}{a}$
  27 $F; F; T; T$
  29 $A = \frac{1}{3}, B = -\frac{1}{3}$
  31 $e^t - 1$
  33 $1 - e^{-t}$
  35 $6; 6 + Ae^{-2t}; 6 - 6e^{-2t} + 6e^{-2t}; 6$
  37 $4; -\frac{1}{e}; 4$
  39 $ye^{-t}, y(t) = te^t$
  41 $A = 1, B = -1, C = -1$
  43 $c^{0.725} > 0.75$
  45 $s(e - 1); \frac{s(e - 1)}{e}$
  47 $(1.02)(1.03) \rightarrow 5.06%; 5%$ by Problem 27
  49 $20,000 e^{20-T}(1.05) = 34,400$ (it grows for 20 - T years)
  51 $s = -cy_0 e^{ct}/(e^{ct} - 1) = -0.01(1000)e^{60}/(e^{60} - 1)$
  53 $y_0 = 0.005 (1 - e^{-0.005})$
  55 $e^{ct} = 1.20$ so $c = \frac{ln 1.20}{4}$
  57 $24e^{36.5}$; $? = 59$ To $-\infty$; constant; to $+\infty$
  61 $\frac{dy}{dx} = 60Y; \frac{dy}{dx} = 60(-Y + 5)$; still $Y_{\infty} = 5$
  63 $y = 60e^{ct} + 20, 60 = 60e^{12c} + 20, c = \frac{1}{12} \ln (\frac{40}{60})$
  100 = 60e^{ct} + 20 at $t = 1/4 \ln (\frac{50}{60})$  65 0
### Section 6.4 Logarithms (page 258)

<table>
<thead>
<tr>
<th>Problem</th>
<th>Solution</th>
</tr>
</thead>
<tbody>
<tr>
<td>1. ( \frac{1}{x} )</td>
<td>( \frac{3 - \ln(x^2)}{2} )</td>
</tr>
<tr>
<td>3. ( \ln x )</td>
<td>( 7 \cos(x) )</td>
</tr>
<tr>
<td>5. ( \frac{\ln x}{x} )</td>
<td>( 3 \tan x )</td>
</tr>
<tr>
<td>7. ( \ln \cos x )</td>
<td>( 9 \sec x )</td>
</tr>
<tr>
<td>9. ( \frac{1}{x} )</td>
<td>( \frac{1}{x} )</td>
</tr>
<tr>
<td>11. Tangent at ( x = \ln \frac{1}{3} )</td>
<td>( \ln t + C )</td>
</tr>
<tr>
<td>13. ( \ln \frac{4}{3} )</td>
<td>( \frac{1}{2} )</td>
</tr>
<tr>
<td>15. ( \frac{1}{2} \ln 5 )</td>
<td>( 17 - \ln(\ln 2) )</td>
</tr>
<tr>
<td>17. ( \ln(\sin x) + C )</td>
<td>( 21 - \frac{1}{3} \ln(\cos 3x) + C )</td>
</tr>
<tr>
<td>19. ( \frac{1}{2} \ln \frac{4}{3} )</td>
<td>( \frac{1}{2} \ln(z^3) + C )</td>
</tr>
<tr>
<td>21. ( \ln \cos 3x + C )</td>
<td>( 23 - \frac{1}{3} \ln(x^2 + 1) )</td>
</tr>
<tr>
<td>23. ( \frac{1}{2} \ln(x^2 + 1) )</td>
<td>( \frac{1}{2} )</td>
</tr>
<tr>
<td>25. ( \ln(y) )</td>
<td>( e^x )</td>
</tr>
<tr>
<td>27. ( \frac{dy}{dx} = e^{\sin x} \cos x )</td>
<td>( 29 )</td>
</tr>
<tr>
<td>31. ( \frac{dy}{dx} = \frac{e^x e^{-x}}{\sqrt{x^2 + 1}} )</td>
<td>( 33 )</td>
</tr>
<tr>
<td>35. ( \ln y = -1 ) so ( y = \frac{1}{e} )</td>
<td>( \frac{dy}{dx} = 0 )</td>
</tr>
<tr>
<td>37. ( \ln(\ln x) + C )</td>
<td>( 0 )</td>
</tr>
</tbody>
</table>

### Section 6.5 Separable Equations Including the Logistic Equation (page 266)

<table>
<thead>
<tr>
<th>Problem</th>
<th>Solution</th>
</tr>
</thead>
<tbody>
<tr>
<td>1. ( 3 + (-3)^x )</td>
<td>( -5 + \frac{3}{2} )</td>
</tr>
<tr>
<td>3. ( 2 )</td>
<td>( x^3 )</td>
</tr>
<tr>
<td>5. ( x )</td>
<td>( 7 )</td>
</tr>
<tr>
<td>7. ( e^{1 - \cos t} )</td>
<td>( 9 )</td>
</tr>
<tr>
<td>9. ( \frac{\ln x}{x} )</td>
<td>( 11 )</td>
</tr>
<tr>
<td>11. Tangent at ( x = \ln \frac{1}{3} )</td>
<td>( 13 )</td>
</tr>
<tr>
<td>13. Tangent at ( x = \ln \frac{1}{3} )</td>
<td>( 15 )</td>
</tr>
<tr>
<td>17. ( x = \ln \frac{1}{3} )</td>
<td>( 19 )</td>
</tr>
<tr>
<td>21. ( x = \ln \frac{1}{3} )</td>
<td>( 23 )</td>
</tr>
<tr>
<td>25. ( x = \ln \frac{1}{3} )</td>
<td>( 27 )</td>
</tr>
<tr>
<td>29. ( x = \ln \frac{1}{3} )</td>
<td>( 31 )</td>
</tr>
<tr>
<td>33. ( x = \ln \frac{1}{3} )</td>
<td>( 35 )</td>
</tr>
<tr>
<td>37. ( x = \ln \frac{1}{3} )</td>
<td>( 39 )</td>
</tr>
</tbody>
</table>

### Section 6.6 Powers Instead of Exponentials (page 276)

<table>
<thead>
<tr>
<th>Problem</th>
<th>Solution</th>
</tr>
</thead>
<tbody>
<tr>
<td>1. ( 1 - x + \frac{x^2}{2} - \frac{x^3}{6} + \cdots )</td>
<td>( 3 )</td>
</tr>
<tr>
<td>3. ( 1 + x + \frac{x^2}{2} + \frac{x^3}{3} + \cdots )</td>
<td>( 5 )</td>
</tr>
<tr>
<td>7. ( 1 + n(-\frac{1}{n})^2 )</td>
<td>( 9 )</td>
</tr>
<tr>
<td>9. ( 1 + n(-\frac{1}{n})^2 )</td>
<td>( 11 )</td>
</tr>
<tr>
<td>13. ( y(3) = 8 )</td>
<td>( 15 )</td>
</tr>
<tr>
<td>17. ( y(t) = t^3 )</td>
<td>( 19 )</td>
</tr>
<tr>
<td>21. ( y(t) = a^t ) if ( a \neq 0 )</td>
<td>( 23 )</td>
</tr>
<tr>
<td>25. ( y(t) = t^3 )</td>
<td>( 27 )</td>
</tr>
<tr>
<td>29. ( P = \frac{b}{c+d} )</td>
<td>( 31 )</td>
</tr>
<tr>
<td>33. ( P = \frac{b}{c+d} )</td>
<td>( 35 )</td>
</tr>
<tr>
<td>37. ( P = \frac{b}{c+d} )</td>
<td>( 39 )</td>
</tr>
</tbody>
</table>

### Section 6.7 Hyperbolic Functions (page 280)

<table>
<thead>
<tr>
<th>Problem</th>
<th>Solution</th>
</tr>
</thead>
<tbody>
<tr>
<td>1. ( e^x, e^{-x}, e^{x/2} - e^{-x/2} )</td>
<td>( \frac{1}{2} \sinh 2x )</td>
</tr>
<tr>
<td>3. ( \sinh x )</td>
<td>( \sinh nx )</td>
</tr>
<tr>
<td>5. ( \cosh x )</td>
<td>( \sinh(3x + 1) )</td>
</tr>
<tr>
<td>7. ( \cosh x )</td>
<td>( -\sinh x \cosh^2 x )</td>
</tr>
<tr>
<td>9. ( \tanh x )</td>
<td>( \sech x )</td>
</tr>
<tr>
<td>11. ( \frac{1}{2} )</td>
<td>( \frac{1}{2} )</td>
</tr>
<tr>
<td>13. ( \frac{1}{2} )</td>
<td>( \frac{1}{2} )</td>
</tr>
<tr>
<td>15. ( \frac{1}{2} )</td>
<td>( \frac{1}{2} )</td>
</tr>
<tr>
<td>17. ( \frac{1}{2} )</td>
<td>( \frac{1}{2} )</td>
</tr>
<tr>
<td>19. ( \frac{1}{2} )</td>
<td>( \frac{1}{2} )</td>
</tr>
<tr>
<td>21. ( \frac{1}{2} )</td>
<td>( \frac{1}{2} )</td>
</tr>
<tr>
<td>23. ( \frac{1}{2} )</td>
<td>( \frac{1}{2} )</td>
</tr>
<tr>
<td>25. ( \frac{1}{2} )</td>
<td>( \frac{1}{2} )</td>
</tr>
</tbody>
</table>

---

A-10 Answers to Odd-Numbered Problems
Answers to Odd-Numbered Problems

33 ∫ y dx = ∫ sinh t sinh t dt; A = \frac{1}{2} sinh t cosh t - ∫ y dx; A' = \frac{1}{2}; A = 0 \text{ at } t = 0 \text{ so } A = \frac{1}{2}t.

41 e^y = x + √(x^2 + 1), y = ln[x + √(x^2 + 1)]

47 \frac{1}{4} \ln\left[\frac{4x + 4}{2 - x}\right]

49 sinh^{-1} x \text{ (see 41)}

51 -sech^{-1} x

53 \frac{1}{2} \ln 3; \infty

55 y(x) = \frac{1}{c} \cosh cx; \frac{1}{c} \cosh cL - \frac{1}{c}

57 y'' = y - 3y^2; \frac{1}{2}(y')^2 = \frac{1}{2}y^2 - y^3 \text{ is satisfied by } y = \frac{1}{2}\text{sech}^2 \frac{x}{2}

CHAPTER 7 TECHNIQUES OF INTEGRATION

Section 7.1 Integration by Parts (page 287)

1 \int x \cos x + x + C \quad 3 \int xe^{-x} - e^{-x} + C \quad 5 \int x^2 \sin x + 2x \cos x - 2 \sin x + C

7 \frac{1}{2}(2x + 1) \ln(2x + 1) - x + C

9 \frac{1}{2}e^{x}(\sin x - \cos x) + C

11 \int \frac{e^{x}}{2^x + 3^x}(\sin bx - b \cos bx) + C

13 \frac{3}{2}(\sin(\ln x) - \cos(\ln x)) + C

15 \int x (\ln x)^2 - 2x \ln x + 2x + C

17 x \sin^{-1} x + \sqrt{1 - x^2} + C

19 \frac{1}{2}(x^2 + 1) \tan^{-1} x - \frac{x}{2} + C

21 x^3 \sin x + 3x^2 \cos x - 6x \sin x - 6 \cos x + C

23 e^x(x^3 - 3x^2 + 6x - 6) + C

25 \int x \tan x + \ln(\cos x) + C

27 -1

29 -\frac{3}{4} e^{-2} + \frac{1}{4} C

31 -2

33 \int 2 \pi 10 - 6 + 2 \tan^{-1} x

35 u = x^n, v = e^x

37 u = x^n, v = \sin x

39 u = (\ln x)^n, v = x

41 u = x \sin x, v = e^x \rightarrow \int e^x \sin x \, dx \text{ in 9 and } - \int \cos x \, e^x \, dx \text{. Then } u = -x \cos x, v = e^x \rightarrow \int e^x \cos x \, dx

\text{ in 10 and } - \int x \sin x \, e^x \, dx \text{ (move to left side): } e^x (\sin x - \cos x + \cos x). \text{ Also try } u = xe^x, v = -\cos x.

43 \int \frac{1}{2} u \sin u \, du = \frac{1}{2}(u - u \cos u) = \frac{1}{2}(x^2 - x^2 \cos x^2); \text{ odd}

45 \text{ 3: step function; 3e^x: step function

49 0; x \delta(x) - \int \delta(x) \, dx = -1; v(x) \delta(x) - \int v(x) \delta(x) \, dx

51 v(x) = \int_0^1 f(x) \, dx

53 u(x) = \frac{k}{\frac{k}{2} \int_0^x v(z) \, dz}; \frac{1}{\frac{k}{2} \int_0^x v(z) \, dz}

55 u = x^2, v = -\cos x \rightarrow -x^2 \cos x + (2x) \sin x - \int 2 \sin x \, dx

57 \text{ Compare 23}

59 uu' |^a_0 - \int_0^a u'u' - u'u' |^a_0 + \int_0^a u'u' = [u'u' - u'u']|^a_0

61 \text{ No mistake: } e^x \cos x - e^x \sin h x = 1 \text{ is part of the constant } C

Section 7.2 Trigonometric Integrals (page 293)

1 \int (1 - \cos^2 x) \sin x \, dx = -\cos x + \frac{1}{3} \cos^3 x + C

3 \frac{1}{2} \sin^2 x + C

5 \int (1 - u^2) \, du = -\frac{1}{3} \cos^3 x + \frac{2}{5} \cos^5 x - \frac{1}{7} \cos^7 x + C

7 \int \frac{1}{2} (\sin x)^{3/2} + C

9 \int \frac{\sin^3 2x \, dx}{2} = \frac{1}{16}(- \cos 2x + \frac{1}{2} \cos^3 2x) + C

11 \frac{\pi}{2}

13 \frac{1}{3}(\frac{3x}{2} + \sin 6x) + C

15 x + C

17 \frac{1}{6} \cos^5 x \sin x + \frac{5}{6} \int \cos^4 x \, dx \text{; use equation (5)}

19 \int_0^{\pi/2} \cos^n x \, dx = \frac{n-1}{n} \int_0^{\pi/2} \cos^{n-2} x \, dx = \cdots = \frac{n-1}{n} \frac{n-3}{n-2} \cdots \frac{1}{2} \int_0^{\pi/2} \, dx

21 I = -\sin^{-1} x \cos x + (n - 1) \int \sin^{-1} x \cos^2 x \, dx = -\sin^{-1} x \cos x + (n - 1) \int \sin^{n-2} x \, dx - (n - 1) I.

\text{ So } nI = -\sin^{-1} x \cos x + (n - 1) \int \sin^{n-2} x \, dx.

23 0, +0, 0, 0, -

25 \frac{2}{3} \cos^3 x, 0

27 \frac{1}{2} (\cos^2 x + \frac{x}{200} \cos 200x), 0

29 \frac{1}{2} (\sin 200x + \frac{1}{200} \sin 200x), 0

31 -\frac{1}{2} \cos x, 0

33 \int_0^a \sin x \, dx = \int_0^a A \sin^2 x \, dx \rightarrow A = 2

35 \text{ Sum = zero = } \frac{1}{2} \text{ (left + right)}

37 p \text{ is even}

39 \frac{1}{2} - q \text{ is even}

41 \sec x + C

43 \frac{1}{3} \tan^3 x + C

45 \frac{1}{3} \sec^3 x + C

47 \frac{1}{3} \tan^3 x - \tan x + x + C

49 \int |\sin x| \, dx

51 \frac{1}{2} \cos^2 x + C

53 A = \sqrt{2}, -\sqrt{2} \sin(x + \frac{\pi}{4})

55 \frac{1}{2}\sqrt{2} \quad 57 \frac{1000}{\sqrt{3}} \quad 59 \frac{1 - \cos x + \sin x}{1 + \cos x + \sin x} + C

61 p \text{ and } q \text{ are 10 and 1}
Section 7.3 Trigonometric Substitutions  (page 299)

1  \[ x = 2 \sin \theta; \int d\theta = \sin^{-1} \frac{x}{2} + C \]
2  \[ x = 2 \sin \theta; \int 4 \cos^2 \theta \ d\theta = 2 \sin^{-1} \frac{x}{2} + x \sqrt{1 - \frac{x^2}{4}} + C \]
3  \[ x = 2 \sin \theta; \int d\theta = \sin^{-1} \frac{x}{2} + C \]
4  \[ x = 2 \sin \theta; \int 4 \cos^2 \theta \ d\theta = 2 \sin^{-1} \frac{x}{2} + x \sqrt{1 - \frac{x^2}{4}} + C \]
5  \[ x = 2 \sin \theta; \int \sin^2 \theta \ d\theta = \frac{1}{2} \sin \theta - \frac{1}{2} \cos \theta + C \]
6  \[ x = 2 \sin \theta; \int \sin^2 \theta \ d\theta = \frac{1}{2} \sin \theta - \frac{1}{2} \cos \theta + C \]
7  \[ x = 2 \sin \theta; \int \cos^2 \theta \ d\theta = \frac{1}{2} \sin \theta + \frac{1}{2} \cos \theta + C \]
8  \[ x = 2 \sin \theta; \int \cos^2 \theta \ d\theta = \frac{1}{2} \sin \theta + \frac{1}{2} \cos \theta + C \]
9  \[ x = 2 \sin \theta; \int \sec^2 \theta \ d\theta = \frac{1}{2} \tan \theta + \frac{1}{2} \cot \theta + C \]
10 \[ x = 2 \sin \theta; \int \sec^2 \theta \ d\theta = \frac{1}{2} \tan \theta + \frac{1}{2} \cot \theta + C \]

Section 7.4 Partial Fractions  (page 304)

1  \[ A = -1, B = 1, -\ln x + \ln(x - 1) + C \]
2  \[ A = -1, B = 1, -\ln x + \ln(x - 1) + C \]
3  \[ A = -1, B = 1, -\ln x + \ln(x - 1) + C \]
4  \[ A = 1, B = -1, -\ln \left( x - \frac{1}{2} \right) + \ln \left( x + \frac{1}{2} \right) + C \]
5  \[ A = 1, B = -1, -\ln \left( x - \frac{1}{2} \right) + \ln \left( x + \frac{1}{2} \right) + C \]
6  \[ A = 1, B = -1, -\ln \left( x - \frac{1}{2} \right) + \ln \left( x + \frac{1}{2} \right) + C \]
7  \[ A = 1, B = -1, -\ln \left( x - \frac{1}{2} \right) + \ln \left( x + \frac{1}{2} \right) + C \]
8  \[ A = 1, B = -1, -\ln \left( x - \frac{1}{2} \right) + \ln \left( x + \frac{1}{2} \right) + C \]

All correct with complex numbers.
27 \( x + 1 = u^2, \, dx = 2u \, du; \) \( \int \frac{2u \, du}{1 + u} = \int \frac{2}{1 + u} \, du = 2 \ln(1 + u) + C = 2\sqrt{x + 1} - 2 \ln(1 + \sqrt{x + 1}) + C \).

29 Note \( Q(a) = 0. \) Then \( \frac{\delta a}{Q(x)} = \frac{\frac{\delta a}{Q(a)}}{Q(a)} \rightarrow \frac{1}{Q'(a)} \) by definition of derivative. At a double root \( Q'(a) = 0. \)

Section 7.5 Improper Integrals (page 309)

1 \( \int_{-\infty}^{\infty} \frac{e^{-x}}{1 - e} \, dx = \frac{1}{e - 1} \)

3 \( -2(1 - x^{1/2}) \bigg|_0^1 = 2 \)

5 \( \tan^{-1}x \bigg|_0^{\pi/2} = \frac{\pi}{2} \)

7 \( \frac{1}{2} (\ln x)^2 \bigg|_0^\infty = -\infty \)

9 \( x \ln x - x \bigg|_0^0 = 0 \)

11 \( \ln(\ln x)) \bigg|_0^\infty = \infty \)

13 \( \frac{1}{2} (x + \sin x \cos x) \bigg|_0^\infty = \infty \)

15 \( \int_{-\infty}^{\infty} \frac{1}{x^p} \, dx \) diverges for every \( p! \)

17 Less than \( \int_1^\infty \frac{dx}{x^p} = \frac{1}{p} \)

21 Less than \( \int_1^\infty \frac{e^x}{x} \, dx = e^2 + e \int_1^\infty e^{-u^2} \, du = e^2 + \frac{\pi}{2} \sqrt{\pi} \)

23 Less than \( \int_0^1 x^2 \, dx + e \int_0^\infty e^{-x^2 - 1} \, dx = e^2 - e \int_1^\infty e^{-u^2} \, du = e^2 - \frac{\pi}{2} \sqrt{\pi} \)

25 \( \int_0^1 \sin \frac{x^2}{2} \, dx + \int_0^\infty \sin \frac{x^2}{2} \, dx \) less than \( 1 + \int_1^\infty \frac{dx}{x^2} = 2 \)

27 \( p! = p \) times \( (p - 1)!; \) \( 1 = 1 \) times \( 0! \)

29 \( u = x, \, dv = x e^{-x^2} \, dx \)

31 \( \int_0^\infty 10000 e^{-14} \, dt = -10,000 e^{-14} \bigg|_0^\infty = \$10,000 \)

33 \( W = -\frac{GMm}{x} \bigg|_R = \frac{GMm}{R} = \frac{1}{2} mu^2 \) if \( v_0 = \sqrt{\frac{2GM}{R}} \)

35 \( \int_0^\infty \frac{dx}{2x} = \int_0^\infty e^{-x} \ln 2 \, dx = \frac{e^{-x} \ln 2}{-\ln 2} \bigg|_0^\infty = \frac{1}{\ln 2} \)

37 \( \int_0^{\pi/2} (\sec x - \tan x) \, dx = [\ln(\sec x + \tan x) + \ln(\cos x)]^{\pi/2}_0 = [\ln(1 + \sin x)]^{\pi/2}_0 = \ln 2. \)

The areas under \( \sec x \) and \( \tan x \) separately are infinite

Chapter 8 Applications of the Integral

Section 8.1 Areas and Volumes by Slices (page 318)

1 \( x^2 - 3 = 1 \) gives \( x = \pm 2; \) \( \int_{-2}^2 [(1 - (x^2 - 3)] \, dx = \frac{32}{3} \)

3 \( y^2 = x = 9 \) gives \( y = \pm 3; \) \( \int_{-3}^3 [9 - y^2] \, dy = 36 \)

5 \( x^4 - 2x^2 = 2x^2 \) gives \( x = \pm 2 \) (or \( x = 0); \) \( \int_{-2}^2 [2x^2 - (x^4 - 2x^2)] \, dx = \frac{128}{15} \)

7 \( y = x^2 = -x^2 + 18x \) gives \( x = 0, 9; \) \( \int_0^9 [(-x^2 + 18x) - x^2] \, dx = 243 \)

9 \( y = \cos x = \cos^2 x \) when \( \cos x = 1 \) or \( 0, x = 0 \) or \( \frac{\pi}{2} \) or \( \cdots \)

11 \( e^x = e^{2x - 1} \) gives \( x = 1; \) \( \int_1^x [e^x - e^{2x - 1}] \, dx = (e - 1) - \left(\frac{e - e^{-1}}{2}\right) \)

13 Intersections \( (0, 0), (1, 3), (2, 2); \) \( \int_0^1 [3x - x] \, dx + \int_1^2 [4 - x - x] \, dx = 2 \)

15 Inside, since \( 1 - x < \sqrt{1 - x^2}; \) \( \int_1^\pi \sqrt{1 - x^2} - (1 - x^2) \, dx = \frac{\pi}{2} - \frac{4}{3} \)

17 \( V = \int_{-a}^a \pi y^2 \, dx = \int_{-a}^a \pi b^2 [1 - \frac{x^2}{a^2}] \, dx = \frac{4\pi ab^2}{3}; \) around \( y \) axis \( V = \frac{4\pi ab^2}{3}; \) rotating \( x = 2, y = 0 \) around \( y \) axis gives a circle not in the first football

19 \( V = \int_0^\pi 2\pi x \sin x \, dx = 2\pi^2 \)

21 \( \int_0^\pi (8 - x^2) \, dx = \frac{512\pi}{3}; \) \( \int_0^\pi 2\pi x (8 - x) \, dx = \frac{512\pi}{3} \) (same cone tipped over)

23 \( \int_0^\pi 1 \cdot 2x - \int_0^\pi \pi (x^2)^2 \, dx = \frac{8\pi}{9}; \) \( \int_0^\pi 2\pi (1 - x^4) \, dx = \frac{2\pi}{3} \)

25 \( \int_{1/3}^2 \pi (8^2) \, dx - \int_{1/3}^2 \pi (1/2)^2 \, dx = \frac{25\pi}{2}; \) \( \int_{1/3}^2 \pi (3 - x) \, dx = \frac{25\pi}{3} \)

27 \( \int_0^\pi \pi [(x^2/3)^2 - (x^2/3)^2] \, dx = \frac{5\pi}{28}; \) \( \int_0^\pi 2\pi x (x^2/3 - x^2/3) \, dx = \frac{5\pi}{28} \) (notice \( xy \) symmetry)

29 \( x^2 = x^2 + y^2, \) \( V = \int_{R-h}^R \pi (R^2 - y^2) \, dy = \pi (R^2 - h^2) \)

31 \( \int_a^b [2(\sqrt{a^2 - x^2})^2 \, dx = \frac{16}{3} a^3 \)

33 \( \int_0^\pi (2 \sqrt{1 - y^2})^2 \, dy = 2 \)

37 \( \int A(x) \, dx \) or in this case \( \int a(y) \, dy \)

39 Ellipse; \( \sqrt{1 - x^2} \) \( \frac{1}{2} \) \( (1 - x^2) \) \( \tan \theta; \) \( \frac{1}{2} \) \( \tan \theta; \) \( \frac{3}{2} \) \( \tan \theta \)

41 Half of \( \pi r^2 h; \) rectangles \( \int_0^3 \pi (5^2 - 2^2) \, dx = 42 \pi \)

45 \( \int_0^3 \pi (4^2 - 1^2) \, dx = 30 \pi \)
47  \int_0^b \pi((b - y)^2 - a^2) dy = \frac{\pi}{6}(b^3 - 3a^2b + 2a^3)  
49  \int_0^2 \pi(3 - x)^2 dx; \int_0^1 2\pi y(2) dy + \int_1^3 2\pi y(3 - y) dy

51  \int_a^b \pi \left(\frac{b}{m}\right)^2 dy = \frac{\pi(b-a)^2}{3m}  
53  960 \pi \text{ cm}   
55  \frac{\pi}{2}  
57  2\pi^3 

59  2\pi  
61  \int_0^4 2\pi y(2 - \sqrt{2}) dy = \frac{32\pi}{3}  
63  3\pi e  
65  \text{Height 1; } \int_0^a 2\pi x e dx = \pi a^2; \text{ cylinder}

67  \text{Length of hole is } 2\sqrt{b^2 - a^2}, \text{ so } b^2 - a^2 = 1 \text{ and volume is } \frac{4\pi}{3} 
69  \text{F; } T(?) \text{; } F \text{; } T

Section 8.2 Length of a Plane Curve (page 324)

1  \int_0^1 \sqrt{1 + (\frac{3}{2}z^{1/2})^2} dz = \frac{8}{27}(1^{3/2} - 1) = \frac{13\sqrt{13} - 8}{27}  
3  \int_0^1 \sqrt{1 + x^2(x^2 + 2)} dx = \int_0^1 (1 + x^2) dx = \frac{4}{3}

5  \int_0^3 \sqrt{1 + (z^2 - \frac{1}{4}z^2)^2} dz = \int_0^3 (z^2 + \frac{1}{4}z^4) dz = \frac{53}{6}

7  \int_0^4 \sqrt{1 + (\frac{1}{4}z^{-1/2})^2} dz = \int_0^4 (\frac{1}{4}z^{1/2} + \frac{1}{4}z^{-1/2}) dz = \frac{31}{6}

9  \int_0^{\pi/2} \sqrt{9\cos^2 t \sin^2 t + 9\sin^2 t \cos^2 t} \ dt = \int_0^{\pi/2} 3\cos t \sin t \ dt = \frac{3}{2}

11  \int_0^{\pi/2} \sqrt{\sin^2 t + (1 - \cos t)^2} dt = \int_0^{\pi/2} \sqrt{2 - 2\cos t} \ dt = \int_0^{\pi/2} 2\sin t \ dt = 4 - 2\sqrt{2}

13  \int_0^1 \sqrt{2 + t + 1} dt = \int_0^1 (t + 1) dt = \frac{3}{2}  
15  \int_0^\pi \sqrt{1 + \cos^2 z} \ dz = 3.820  
17  \int_0^\pi \sqrt{1 + \frac{1}{2}z^2} \ dz = 2.003

19  \text{Graphs are flat toward } (1,0) \text{ then steep up to } (1,1); \text{ limiting length is 2}

21  \frac{ds}{dt} = \sqrt{36\sin^2 3t + 36\cos^2 3t} = 6  
23  \int_0^1 \sqrt{26} \ dy = \sqrt{26}

25  \int_0^1 \sqrt{\frac{1}{4}(e^y - e^{-y})^2 + 1} \ dy = \int_0^1 \frac{1}{2}(e^y + e^{-y}) \ dy = \frac{1}{2}(e^y - e^{-y}) \bigg|_0^1 = e - \frac{1}{e}

\text{Using } x = \cosh y \text{ this is } \int \sqrt{1 + \sinh^2 y} \ dy = \int \cosh y \ dy = \sinh y \bigg|_1^1 = 2 \sinh 1

27  \text{Ellipse; two } y's \text{ for the same } x  
29  \text{Carpet length } 2 \neq \text{ straight distance } \sqrt{2}

31  (ds)^2 = (dx)^2 + (dy)^2 + (dz)^2; ds = \sqrt{\left(\frac{dx}{dt}\right)^2 + \left(\frac{dy}{dt}\right)^2 + \left(\frac{dz}{dt}\right)^2} dt

\text{surface area } \int_0^{2\pi} 2\pi(x + a \sin t) dt = 4\pi^2 ab

17  \int_0^2 2\pi x \sqrt{1 + \left(\frac{-s}{2x - s^2}\right)^2} dx = \int_0^1 2\pi x \sqrt{2s - x^2} = \pi^2 + 2\pi \text{ (write } 2x - s^2 = 1 - (x - 1)^2 \text{ and set } x - 1 = \sin \theta)

19  \int_1^2 2\pi x \sqrt{1 + \frac{1}{x^2}} dx \text{ (can be done)}

21  \text{Surface area } = \int_0^\infty 2\pi \frac{1}{2} \sqrt{1 + \frac{1}{z^2}} \ dz > \int_1^\infty \frac{2z dz}{\sqrt{z}} = 2\pi \ln z \bigg|_1^\infty = \infty \text{ but volume } = \int_1^\infty \pi \left(\frac{1}{a}\right)^2 dx = \pi

23  \int_0^\pi 2\pi \sin \sqrt{2\sin^2 t + \cos^2 t} \ dt = \int_0^\pi 2\pi \sin t \sqrt{2 - \cos^2 t} \ dt = \int_0^1 2\pi \sqrt{2 - u^2} du = \pi u \sqrt{2 - u^2} + 2\pi \sin^{-1} \frac{u}{\sqrt{2}} \bigg|_1^1 = 2\pi + \pi^2

Section 8.3 Area of a Surface of Revolution (page 327)

1  \int_0^6 2\pi \sqrt{2 + \left(\frac{1}{2x}\right)^2} dx = \int_0^6 2\pi \sqrt{2 + \frac{1}{4}x^2} \ dx = \frac{49\pi}{3}  
3  \int_0^1 2\pi(7x) \sqrt{50} \ dx = 14\pi \sqrt{50}

5  \int_0^1 2\pi \sqrt{4 - z^2} \left(\frac{1}{4}z^{-2}\right) \ dz = \int_0^1 4\pi dx = 8\pi  
7  \int_2^3 2\pi x \sqrt{1 + (2x)^2} \ dx = \frac{\pi}{6} (1 + 4x^2)^{3/2} \bigg|_0^3 = \frac{\pi}{6}[17^{3/2} - 1]

9  \int_0^3 2\pi x \sqrt{2} = 9\pi\sqrt{2}  
11  \text{Figure shows radius } s \text{ times angle } \theta = \arccos 2\pi R

13  2\pi \Delta s = \pi(R + R')(s - s') = \pi Rs - \pi R's' \text{ because } R's - R's' = 0

15  \text{Radius } a, \text{ center at } (0, b); \left(\frac{dx}{dt}\right)^2 + \left(\frac{dy}{dt}\right)^2 = a^2, \text{ surface area } \int_0^{2\pi} 2\pi(b + a \sin t) a \ dt = 4\pi^2 ab

17  \int_0^1 2\pi x \sqrt{1 + \left(\frac{1}{2x - s^2}\right)^2} dx = \int_0^1 2\pi x \sqrt{2s - x^2} = \pi^2 + 2\pi \text{ (write } 2x - s^2 = 1 - (x - 1)^2 \text{ and set } x - 1 = \sin \theta)

19  \int_1^2 2\pi x \sqrt{1 + \frac{1}{x^2}} dx \text{ (can be done)}

21  \text{Surface area } = \int_0^\infty 2\pi \frac{1}{2} \sqrt{1 + \frac{1}{z^2}} \ dz > \int_1^\infty \frac{2z dz}{\sqrt{z}} = 2\pi \ln z \bigg|_1^\infty = \infty \text{ but volume } = \int_1^\infty \pi \left(\frac{1}{a}\right)^2 dx = \pi

23  \int_0^\pi 2\pi \sin t \sqrt{2\sin^2 t + \cos^2 t} \ dt = \int_0^\pi 2\pi \sin t \sqrt{2 - \cos^2 t} \ dt = \int_1^1 2\pi \sqrt{2 - u^2} du = \pi u \sqrt{2 - u^2} + 2\pi \sin^{-1} \frac{u}{\sqrt{2}} \bigg|_1^1 = 2\pi + \pi^2

Section 8.4 Probability and Calculus (page 334)

1  P(X < 4) = \frac{7}{8}, P(X = 4) = \frac{1}{16}, P(X > 4) = \frac{1}{16}  
3  \int_0^\infty p(x) dx \text{ is not 1; } p(x) \text{ is negative for large } x

5  \int_2^\infty e^{-z} dx = \frac{1}{e^6}; \int_1^{0.1} e^{-z} dz \approx (0.1)^\frac{1}{e}  
7  p(x) = \frac{1}{x}; F(x) = \frac{x}{e} \text{ for } 0 \leq x \leq \pi \text{ (} F = 1 \text{ for } x > \pi \text{)}
$\int_0^\infty \frac{2e^{-ax}}{\pi(1+x^2)} dx = \frac{1}{\pi} \ln(1+x^2)|_0^\infty = +\infty$

13 $\int_0^\infty axe^{-ax} dx = \frac{2}{\pi} \tan^{-1} x$; $\int_0^\infty e^{-ax} dx = 1 - e^{-x}$; $\int_0^\infty axe^{-ax} dx = 1 - e^{-ax}$

17 $\int_0^\infty \frac{1}{10} e^{-x/10} dx = -e^{-x/10}|_0^\infty = \frac{1}{e}$

19 Exponential better than Poisson: 60 years $\int_0^{60} .01 e^{-0.01x} dx = 1 - e^{-6} = .45$

21 $y = \frac{x}{\sigma}$; three areas $\approx \frac{1}{3}$ each because $\mu - \sigma$ to $\mu + \sigma$ and areas add to 1

23 $-2\mu \int xp(x) dx + \mu^2 \int p(x) dx = -2\mu \cdot \mu + \mu^2 \cdot 1 = -\mu^2$

25 $\mu = 0 + \frac{1}{5} + \frac{1}{3} + \frac{2}{3} = \frac{11}{3}$; $\sigma^2 = (0 - 1)^2 \cdot \frac{1}{3} + (1 - 1)^2 \cdot \frac{1}{3} = \frac{8}{3}$.

Also $\sum n^2 p_n - \mu^2 = 0 - \frac{1}{3} + 1 - \frac{1}{3} + 4 - \frac{1}{3} = \frac{2}{3}$.

27 $\mu = \int_0^\infty \frac{x e^{-x^2/2}}{2} dx = 2; 1 - \int_0^4 e^{-x^2/2} dx = 1 + [e^{-2}]_{10} = e^{-2}$

29 Standard deviation (yes – no poll) $\leq \frac{1}{2\sqrt{N}} = \frac{1}{2\sqrt{900}} = \frac{1}{30}$ Poll showed $870/900 = \frac{29}{30}$ peaceful.

95% confidence interval is from $\frac{29}{30} - \frac{4}{60}$ to $\frac{29}{30} + \frac{2}{60}$, or 93% to 100% peaceful.

31 95% confidence of unfair if more than $\frac{2\sigma}{\sqrt{N}} = \frac{1}{\sqrt{2500}} = .02$ away from 50% heads.

2% of 2500 = 50. So unfair if more than 1300 or less than 1200.

33 55 is $1.5\sigma$ below the mean, and the area up to $\mu - 1.5\sigma$ is about 8% so 24 students fail.

A grade of 57 is $1.3\sigma$ below the mean and the area up to $\mu - 1.3\sigma$ is about 10%.

35 .99; .99\textsuperscript{1000} = (1 - \frac{1}{100})\textsuperscript{1000} $\approx \frac{1}{e}$ because $(1 - \frac{1}{n})^n \to \frac{1}{e}$.

Section 8.5 Masses and Moments (page 340)

1 $x = \frac{10}{6}$ 3 $x = \frac{4}{2}$ 5 $x = \frac{3}{5}$ 7 $x = \frac{2}{3} = \frac{\pi}{3}$ 9 $x = \frac{7/2}{2} = \frac{\pi}{2}$ 11 $x = \frac{1/3}{\pi/4} = \frac{\pi}{9}$ 13 $x = \frac{1/4}{\pi/2}, \rho = \frac{1/3}{\pi/2}$

15 $x = \frac{0}{3\pi} = \frac{\rho}{\pi}$ 21 $I = \int x^2 \rho dx - 2t \int \rho dx + t^2 \int \rho dx; \frac{dx}{dt} = -2 \int \rho dx + 2t \int \rho dx = 0$ for $t = x$

23 South Dakota 25 $2\pi a^2 b$ 27 $M_x = 0, M_y = \frac{\pi}{2}$ 29 $\frac{2}{\pi}$ 31 Moment

33 $I = \sum m_n r_n^2 \frac{1}{2} \frac{1}{2} \sum m_n r_n^2 \omega_n^2 \cdot 0$ 35 $14\pi \ell^2/2; 14\pi \ell^2/4 \cdot 1/2$

37 $\frac{1}{3}$; solid ball, solid cylinder, hollow ball, hollow cylinder 39 No

41 $T \approx \sqrt{1 + J}$ by Problem 40 so $T \approx \sqrt{1.4}, \sqrt{1.5}, \sqrt{5/3}, \sqrt{2}$

Section 8.6 Force, Work, and Energy (page 346)

1 2.4 ft lb; 2.424 \ldots ft lb 3 24000 lb/ft; 83\frac{1}{3} ft lb 5 10x ft lb; 10x ft lb 7 25000 ft lb; 20000 ft lb

9 864,000 Nkm 11 5.6 \cdot 10^7 Nkm 13 $k = 10$ lb/ft; $W = 25$ ft lb 15 $\int 60wh \ dh = 48000w, 12000w$

17 $\frac{1}{2} w AH^2, \frac{3}{8} w AH^2$ 19 9600w 21 $1 - \frac{x^2}{8} - 3/2$ 23 $800$ (9800) kg 25 $\pm$ force

CHAPTER 9 POLAR COORDINATES AND COMPLEX NUMBERS

Section 9.1 Polar Coordinates (page 350)

1 $1, \frac{\pi}{2}$ 3 $2, \frac{\pi}{3}$ 5 $\left(\sqrt{2}, \frac{5\pi}{4}\right)$ 7 $(0, 2)$ 9 $(\sqrt{10}, \sqrt{10})$ 11 $(\sqrt{3}, -1)$ 13 $2\sqrt{2}$

15 $\sqrt{r^2 + R^2 - 2rrR \cos(\theta - \phi)}$

17 $0 < r < \infty, -\frac{\pi}{2} < \theta < \frac{\pi}{2}; 0 \leq r < \infty, \pi \leq \theta \leq 2\pi; \sqrt{r} < r < \sqrt{\frac{r}{2}}, 0 \leq \theta < 2\pi; 0 \leq r < \infty, -\frac{\pi}{4} \leq \theta \leq \frac{\pi}{4}$

19 $y = x \tan \theta, r = x \sec \theta$ 21 $\theta = \frac{\pi}{2}, r = \frac{\sin \theta + \cos \theta}{\sin \theta}$

23 $x^2 + y^2 = y$ 25 $x = r \sin \theta \cos \theta, y = r \sin^2 \theta, x^2 + y^2 = y$

27 $x^2 + y^2 = x + y, (x - \frac{1}{2})^2 + (y - \frac{1}{2})^2 = (\frac{\sqrt{2}}{2})^2$ 29 $x = \frac{\cos \theta}{\cos \theta + \sin \theta}, y = \frac{\sin \theta}{\cos \theta + \sin \theta}$ 31 $(x^2 + y^2)^3 = x^4$
Section 9.2 Polar Equations and Graphs  
(page 355)

1 Line \( y = 1 \)  
3 Circle \( x^2 + y^2 = 2x \)  
5 Ellipse \( 3x^2 + 4y^2 = 1 - 2x \)  
7 \( x, y, r \) symmetries
9 \( x \) symmetry only  
11 No symmetry  
13 \( x, y, r \) symmetries!

15 \( x^2 + y^2 = 6y + 8x \to (y - 4)^2 + (x - 3)^2 = 5^2, \) center \((4,3)\)  
17 \((2,0), (0,0)\)

19 \( r = 1 - \sqrt{\frac{a}{2}}, \theta = \frac{3\pi}{4}; r = 1 + \sqrt{\frac{a}{2}}, \theta = \frac{7\pi}{4}; (0,0) \)  
21 \( r = 2, \theta = \pm \frac{\pi}{12}, \pm \frac{5\pi}{12}, \pm \frac{7\pi}{12}, \pm \frac{11\pi}{12} \)

23 \((x, y) = (1, 1)\)  
25 \( r = \cos 5\theta \) has 5 petals  
27 \((x^2 + y^2 - x)^2 = x^2 + y^2\)

31 \( \cos \theta = -\frac{1}{2}, \sin \theta = \frac{\sqrt{3}}{2} \to y = \frac{2\sqrt{3}}{3}, x = -\frac{2}{3} \)  
33 \( x = \frac{4}{3}, r = -\frac{5}{3} \)  
35 .967

Section 9.3 Slope, Length, and Area for Polar Curves  
(page 359)

1 Area \( \frac{3\pi}{2} \)  
3 Area \( \frac{9\pi}{2} \)  
5 Area \( \frac{\pi}{8} \)  
7 Area \( \frac{\pi}{6} - \frac{1}{4} \)  
9 \( \int_{-\pi/3}^{\pi/3} \left( \frac{\sqrt{3}}{2} \cos \theta - \frac{(1 + \cos \theta)^2}{2} \right) d\theta = \pi \)

11 Area \( 8\pi \)  
13 Only allow \( r > 0 \), then \( 4 \int_{\pi/4}^{\pi/2} \frac{1}{2} \cos 2\theta \, d\theta = 1 \)  
15 \( 2 + \frac{\pi}{4} \)

17 \( \theta = 0 \); left points \( r = \frac{1}{2}, \theta = \pm \frac{3\pi}{2}, x = -\frac{1}{4}, y = \pm \frac{\sqrt{3}}{4} \)

19 \( \frac{\pi^2}{25} \) is 40,000; \( \frac{\pi^2}{25} \) is 40,000.001

21 \( \tan \psi = \tan \theta \)  
23 \( x = 0, y = 1 \) is on limacon but not circle  
25 \( \frac{1}{2} \ln(2\pi + \sqrt{1 + 4\pi^2}) + n\sqrt{1 + 4\pi^2} \)

27 \( \frac{\pi}{2} \)  
29 \( \frac{1}{2} \) (base)(height) \( \approx \frac{1}{2}(r\Delta \theta)r \)  
31 \( \frac{4x}{5} \sqrt{2} \)  
33 \( 2\pi(2 - \sqrt{2}) \)  
35 \( \frac{8\pi}{3} \)  
39 \( \sec \theta \)

Section 9.4 Complex Numbers  
(page 364)

1 Sum = 4, product = 5  
5 Angles \( \frac{3\pi}{4}, \frac{3\pi}{2}, \frac{\pi}{4} \)  
7 Real axis; imaginary axis; \( \frac{1}{2} \) axis \( x \geq 0 \); unit circle

9 \( \frac{cd}{d} = 5 + 10i, \frac{5}{d} = \frac{11 - 2i}{25} \)

11 \( 2 \cos \theta, 1, -1, 1 \)  
13 Sum = 0, product = -1  
15 \( r^4 e^{4i\theta}, \frac{1}{r} e^{-i\theta}, \frac{1}{r} e^{-4i\theta} \)

17 Evenly spaced on circle around origin  
19 \( e^{it}, e^{-it} \)  
21 \( e^t, e^{-t}, e^0 \)  
23 \( \cos 7t, \sin 7t \)

29 \( t = -\frac{2\pi}{\sqrt{3}} \)  
31 \( F; T \) at most 2; \( \Re c < 0 \)  
33 \( \frac{1}{r} e^{-i\theta}, x = \frac{1}{r} \cos \theta, y = -\frac{1}{r} \sin \theta; \pm \frac{1}{\sqrt{r}} e^{-i\theta/2} \)

CHAPTER 10 INFINITE SERIES

Section 10.1 The Geometric Series  
(page 373)

1 Subtraction leaves \( G - xG = 1 \) or \( G = \frac{1}{1-x} \)  
3 \( \frac{1}{2}, \frac{4}{5}, \frac{100}{11}, \frac{3\times9}{99} \)  
5 \( 2 \cdot 1 + 3 \cdot 2x + 4 \cdot 3x^2 + \cdots = \frac{2}{1-x^3} \)

7 .142857 repeats because the next step divides 7 into 1 again

9 If \( q \) (prime, not 2 or 5) divides \( 10^N - 10^M \) then it divides \( 10^{N-M} - 1 \)  
11 This decimal does not repeat

13 \( \frac{7}{99}, \frac{123}{999} \)

15 \( \frac{x}{1 - x^3} \)

17 \( \frac{x^2}{1 - x^3} \)

19 \( \frac{\ln x}{1 - \ln x} \)  
21 \( \frac{1}{x-1} \)  
23 \( \tan^{-1}(\tan x) = x \)

25 \( (1 + x + x^2 + x^3 + \cdots)(1 - x + x^2 - x^3 + \cdots) = 1 + x^2 + x^4 + \cdots \)

27 \( 2(1.234...) \) is 2 \( \cdot \frac{10}{1} \cdot \frac{1}{(1 - 10)^2} = \frac{20}{81} ; 1 - .0123 \ldots \) is 1 - \( \frac{1}{100}(1 - \frac{1}{10})^2 = \frac{80}{81} \)

29 \( \frac{2}{3}, \frac{1}{3} = 1 \)

31 \( -\ln(1 - .1) = -\ln .9 \)  
33 \( \frac{1}{2} \ln \frac{1}{1.9} \)  
35 \( n(n+1)! \)  
37 \( y = \frac{b}{b - bx} \)

39 All products like \( a_1b_2 \) are missed; \( (1 + 1)(1 + 1) \neq 1 + 1 \)  
41 Take \( x = \frac{1}{2} \) in (13); \( \ln 3 = 1.0986 \)

43 In 3 seconds the ball goes 78 feet  
45 \( \tan z = \frac{2}{3} \) (18) is slower with \( x = \frac{2}{3} \)
Section 10.2 Convergence Tests: Positive Series  

1 \( \frac{1}{2} + \frac{1}{4} + \cdots \) is smaller than \( 1 + \frac{1}{2} + \cdots \)  
5 \( a_n = s_n - s_{n-1} = \frac{1}{n^2} - \frac{1}{(n-1)^2}, s = 1; a_n = 2^{n+1} - 2^n \)  
9 No decision on \( \sum b_n \)  
11 Converges: \( \sum \frac{1}{n^2} \) is larger  
13 Diverges: \( \sum \frac{1}{2n} \) is smaller  
17 Converges: \( \sum \frac{2}{n^2} \) is larger  
19 Converges: \( \sum \frac{3}{n^2} \) is larger  
21 \( \lim \frac{\Delta f}{\Delta x} = 1 \)  
29 \( \lim \frac{f}{\Delta x} = 1 \)  
33 \( \sum e^{-n} \leq \int_0^\infty e^{-x} \, dx = 1; \frac{1}{e} + \frac{1}{e^2} + \cdots = \frac{1}{1-e} \)  
35 Converges faster than \( \sum \)  
37 Diverges because \( \int_0^\infty e^{-x} \, dx = 1 \)  
39 Diverges because \( \int_0^\infty x^{n-1} \, dx = \frac{1}{n-1} \)  
41 \( \sum e^{-n} \leq \int_0^\infty e^{-x} \, dx = 1; \frac{1}{e} + \frac{1}{e^2} + \cdots = \frac{1}{1-e} \)  
43 (b) \( \int_0^\infty e^{-x} \, dx = 1 \)  
45 After adding we have \( 1 + \frac{1}{2} + \cdots + \frac{1}{n} \) (close to \( \ln 2n \)); thus originally close to \( \ln 2n \)  
47 \( \lim_{x \to 0} \frac{x}{\sin x} = \lim_{x \to 0} \frac{\sin x}{x} = 1 \)  
51 After adding we have \( 1 + \frac{1}{2} + \cdots + \frac{1}{n} \) (close to \( \ln 2n \)); thus originally close to \( \ln 2n \)  
53 \( \lim_{x \to 0} \frac{x}{\sin x} = \lim_{x \to 0} \frac{\sin x}{x} = 1 \)  
57 Text proves: If \( \sum |a_n| \) converges so does \( \sum a_n \)  
59 \( \sum a_n = \frac{2}{n+1} \)  
61 Divergence: \( N \) terms add to \( \ln \frac{N+1}{2} \)  
63 Diverge (compare 1)  
65 \( \sum a_n \) converges  
67 Beyond some point \( \frac{a_n}{b_n} < 1 \) or \( a_n < b_n \)  

Section 10.3 Convergence Tests: All Series  

1 Terms don’t approach zero  
3 Absolutely  
5 Conditionally not absolutely  
7 No convergence  
9 Absolutely  
11 No convergence  
13 By comparison with \( \sum \)  
15 Even sums \( \sum a_n \) are not decreasing  
17 (b) If \( a_n > 0 \) then \( s_n \) is too large so \( s_n > s \)  
19 \( s = 1 - \frac{1}{e} \); below by less than \( \frac{1}{5!} \)  
21 Subtract \( 2\left(\frac{1}{2^2} + \frac{1}{3^2} + \cdots\right) = \frac{2}{3}(\frac{1}{2^2} + \frac{1}{3^2} + \cdots) = \frac{\pi^2}{12} \) from positive series to get alternating series  
23 Text proves: If \( \sum |a_n| \) converges so does \( \sum a_n \)  
25 New series \( (\frac{1}{2}) - \frac{1}{4} + (\frac{1}{6}) - \frac{1}{8} \cdots = \frac{1}{2}(1 - \frac{1}{2} + \frac{1}{3} - \frac{1}{4} \cdots) \)  
27 \( \frac{3}{2} \ln 2 \) : add \( \ln 2 \) series to \( \ln \frac{3}{2} \)  
29 Terms alternate and decrease to zero; partial sums are \( 1 + \frac{1}{2} + \cdots + \frac{1}{n} \)  
31 \( \frac{1}{e} \)  
33 Hint + comparison test  
35 Partial sums \( a_n - a_0 \); sum \(-a_0\) if \( a_n \to 0 \)  
37 \( \frac{1}{1-x} = 3 \) but product is not \( 1 + \frac{3}{2} \cdots \)  
39 Write \( x \) to base 2, as in 1.0010 which keeps \( 1 + \frac{1}{8} \) and deletes \( \frac{1}{2}, \frac{1}{4}, \cdots \)  
41 \( \frac{1}{6} + \frac{1}{27} + \cdots \) adds to \( \frac{1}{1-\frac{1}{3}} = \frac{1}{2} \) and can’t cancel \( \frac{1}{3} \)  
43 \( \frac{1}{\cos \frac{1}{2}} = \cot \frac{1}{2} \) (trig identity)  
45 Add imaginary part  

Section 10.4 The Taylor Series for \( e^x, \sin x, \) and \( \cos x \)  

1 \( 1 + 2x + \frac{(2x)^2}{2!} + \cdots \); derivatives \( 2^n; 1 + 2 + \frac{2^2}{2!} + \cdots \)  
3 Derivatives \( i^n; 1 + ix + \cdots \)  
5 Derivatives \( 2^n n!; 1 + 2x + 4x^2 + \cdots \)  
7 Derivatives \( -(n-1)!; -x - \frac{x^3}{3!} - \frac{x^5}{5!} - \cdots \)  
9 \( y = 2 - e^x = 1 - x - \frac{x^2}{2!} - \cdots \)  
11 \( y = x - \frac{x^3}{3!} - \cdots = \sin x \)  
13 \( y = x e^x = x + x^2 + \frac{x^3}{2!} + \cdots \)  
15 \( 1 + 2x + x^2 + 4(x-1)^2 + (x-1)^4 \)  
17 \( \frac{1}{1-x} = 1 - (x-1)^2 - (x-1)^4 + \cdots \)  
19 \( 1 - (x-1)^2 - (x-1)^4 + \cdots \)  
21 \( (x-1)^2 + \frac{(x-1)^3}{3!} - \cdots = \ln(1 + (x-1)) \)  
23 \( e^{-1} e^{1-x} = e^{-1} (1 - (x-1) + \frac{(x-1)^2}{2!} - \cdots) \)  
25 \( x + 2x^2 + 2x^3 \)  
27 \( \frac{1}{2} + \frac{x^2}{12} + \frac{x^4}{720} \)  
29 \( x + \frac{x^3}{3!} + \frac{x^5}{600} \)  
31 \( 1 + x^2 + \frac{x^4}{2!} \)  
33 \( 1 + x - \frac{x^3}{2} \)
Section 10.5 Power Series (page 395)

1. \(1 + 4x + (4x)^2 + \cdots; r = \frac{1}{4}; x = \frac{1}{4}\)
2. \(3 e(1-x+x^2+\cdots); r = \infty\)
3. \(5 e^{1-x-x^2} + (1-x)\)
4. \(7 |x-a| < 1; -\ln(1-|x-a|)\)
5. \(11 \frac{\pi}{2} + \frac{\pi^2}{3} + \cdots; add to 1 at x = 0\)
6. \(13 a_1, a_3, \ldots \) are all zero
7. \(15 \frac{1}{x^2} - \frac{1}{x^3} - \frac{1}{x^4} - \cdots; r = 0\)
8. \(17 f^{(8)}(c) = \cos c < 1; \) alternating terms might not decrease (as required)
9. \(19 f = \frac{1}{1-x}, |R_n| \leq \frac{\pi^{n+1}}{(1-x)^{n+1}} + \frac{(1-c)^4}{2}; (1-c)^2 = 1 - \frac{1}{2}\)
10. \(21 f^{(n+1)}(x) = \frac{n!}{(1-x)^{n+1}}, |R_n| \leq \frac{\pi^{n+1}}{(1-x)^{n+1}} + \frac{(1-c)^{n+1}}{2} = 0\) when \(x = \frac{1}{2} \) and \(1-c > \frac{1}{2}\)
11. \(23 R_2 = f(x) - f(a) - f'(a)(x-a) - \frac{1}{2} f''(a)(x-a)^2\) so \(R_2 = R_2' = R_2'' = 0\) at \(x = a, R_2'' = f''\)
12. Generalized Mean Value Theorem in 3.8 gives \(a < c < c_2 < c_1 < x\)
13. \(25 1 + \frac{1}{2} x^2 + \frac{3}{8} (x^2)^2\)
14. \(27 (1-x^n); (1-n)(n+1)\)
15. \(29 (a) \) one friend \(k\) times, the other \(n-k\) times, \(0 \leq k \leq n; 21\)
16. \(33 (16 - 1)^{1/4} \approx 1.968\)
17. \(35 (1 + 1.1^1 = 1(1.1)(1) + (1.1)(1.1) = 1.1105\)
18. \(37 1 + \frac{\pi^2}{2} + \frac{5\pi^4}{24}; r = \frac{\pi}{2}\)
19. \(41 x + x^2 + \frac{5}{6} x^3 + \frac{5}{6} x^4\)
20. \(43 x^2 - \frac{1}{2} x^4 + \frac{2}{45} x^6\)
21. \(45 1 + \frac{\pi}{2} + \frac{3\pi^2}{8} + \frac{5\pi^4}{16}\)
22. \(47 .2727\)
23. \(49 - \frac{1}{3} = -\frac{1}{2}\)
24. \(51 r = 1, r = \frac{\pi}{2} - 1\)

CHAPTER 11 VECTORS AND MATRICES

Section 11.1 Vectors and Dot Products (page 405)

1. \((0,0,0); (5,5,5); 3; -3; \cos \theta = -1\)
2. \(2(1 - x + \frac{x^2}{2} - \cdots); r = \infty\)
3. \(5 (v_2, -v_1); (v_2, -v_1, 0), (v_3, 0, -v_1)\)
4. \(7 (0,0); (0,0,0)\)
5. \(9 \cos \theta; \) projection of \(w\) on \(v\)
6. \(11 \text{, } F; T; F\)
7. \(13 \text{, } Z; \) sum = 10 o'clock vector; sum = 8 o'clock vector times \(\frac{1+\sqrt{2}}{2}\)
8. \(15 \text{, } 45^{\circ}\)
9. \(17 \text{, } \text{Circle } x^2 + y^2 = 4; (x-1)^2 + y^2 = 4; \) vertical line \(x = 2; \) half-line \(x \geq 0\)
10. \(19 \text{, } v = -3i + 2j, w = 2i - j - k; -i - 7j + 8k; 6; -3; \cos \theta = -\frac{1}{2}\)
11. \(23 \cos \theta = \frac{1}{\sqrt{2}}; \cos \theta = \frac{2}{\sqrt{6}}; \cos \theta = \frac{1}{3}\)
12. \(25 A \cdot (A + B) = 1 + A \cdot B = 1 + B \cdot A = B \cdot (A + B); \) equilateral, \(60^{\circ}\)
13. \(27 a = A \cdot i, b = A \cdot J\)
14. \(29 \text{, } \cos(t, \sin(t)) \text{ and } (-\sin(t, \cos(t)); (\cos(2t, \sin(2t))) \text{ and } (-\sin(2t, 2\cos(2t))\)
15. \(31 C = A + B, D = A - B; C \cdot D = A \cdot A + B \cdot A - A \cdot B - B \cdot B = r^2 - r^2 = 0\)
16. \(33 U + V - W = (2, 5, 8), U - V + W = (0, -1, -2), -U + V + W = (4, 3, 6)\)
17. \(35 c \text{ and } \sqrt{a^2 + b^2}; c/a \text{ and } \sqrt{a^2 + b^2 + c^2}\)
18. \(37 M_1 = \frac{1}{2} A + C, M_2 = A + \frac{1}{2} B, M_3 = B + \frac{1}{2} C, M_3 = \frac{3}{2} (A + B + C) = 0\)
19. \(39 8 \leq 3 \cdot 3; 2\sqrt{xy} \leq x + y\)
20. \(41 \text{, } \text{Cancel } a^2c^2 \text{ and } b^2d^2; \text{ then } b^2c^2 + a^2d^2 \geq 2abcd \text{ because } (bc - ad)^2 \geq 0\)
21. \(43 F; T; F\)
22. \(45 \text{, } \text{all } 2\sqrt{2}; \cos \theta = -\frac{1}{3}\)

Section 11.2 Planes and Projections (page 414)

1. \((0,0,0) \text{ and } (2, -1, 0)\)
2. \(N = (1, 2, 3)\)
3. \((0,5,6) \text{ and } (0,6,7)\)
4. \(N = (1,0,0)\)
5. \((1,1,1) \text{ and } (1,2,2)\)
6. \(N = (1,1,-1)\)
7. \(x + y = 3\)
8. \(9x + 2y + z = 2\)
11 Parallel if \( N \cdot V = 0 \); perpendicular if \( V = \) multiple of \( N \)

13 \( i + j + k \) (vector between points) is not perpendicular to \( N \); \( V \cdot N \) is not zero; plane through first three is \( x + y + z = 1; x + y - z = 3 \) succeeds; right side must be zero

15 \( ax + by + cz = 0; a(x - x_0) + b(y - y_0) + c(z - z_0) = 0 \)

17 \( \cos \theta = \frac{\sqrt{6}}{3}, \frac{\sqrt{6}}{3}, \frac{1}{3} \)

19 \( \frac{2}{\sqrt{56}} A \) has length \( \frac{1}{3} \) \( 21 P = \frac{1}{3} A \) has length \( \frac{1}{2} |A| \)

23 \( -A \) has length \( |A| \)

27 Projection on \( A = (1,2,2) \) has length \( \frac{5}{3} \); force down is 4; mass moves in the direction of \( F \)

29 \( |P|_{\text{min}} = \frac{5}{6} \) distance from plane to origin

31 Distances \( |\frac{1}{\sqrt{3}}| \) and \( \frac{2}{\sqrt{3}} \) both reached at \( (\frac{1}{3}, \frac{1}{3}, -\frac{1}{3}) \)

33 \( i + j + k; t = -\frac{4}{3}; (\frac{2}{3}, -\frac{1}{3}, \frac{1}{3}) \)

35 Same \( N = (2,-2,1) \); for example \( Q = (0,0,1) \); then \( Q + \frac{2}{5} N = (\frac{2}{5}, -\frac{4}{5}, \frac{1}{5}) \) is on second plane; \( \frac{2}{5} N \) is zero

37 \( 3i + 4j; (3t, 4t) \) is on the line if \( 3(3t) + 4(4t) = 10 \) or \( t \) = \( \frac{5}{18} \); \( P = (\frac{5}{18}, \frac{5}{18}, -\frac{5}{18}) \)

39 \( 3i + 2j - 3k \) \( 5 - 2i + 3j - 5k \) \( 727i + 12j - 17k \)

39 \( \cos B = \frac{6}{25} \), \( \frac{6}{25} \), \( \frac{10}{25} \)

43 \( t = \frac{A \cdot B}{|A|^2} \); \( B \) onto \( A \)

45 \( aV = \frac{1}{2} L_I - \frac{1}{2} L_{III}; aV = \frac{1}{2} L_{II} + \frac{1}{2} L_{III} \)

47 \( V \cdot L_I = 2; V \cdot L_{II} = -3; V \cdot L_{III} = -3; V \cdot L_{III} = -3; V \cdot (i - \sqrt{3}j) = -4 \), and \( V = \frac{1}{2} i + \frac{3}{2} \sqrt{3} j \)

Section 11.3 Cross Products and Determinants

1 \( 0 \) \( 8 \) \( 3 \) \( 1 \)

9 \( A \) perpendicular to \( B; A, B, C \) mutually perpendicular

11 \( |A \times B| = \sqrt{2}, A \times B = j - k \)

13 \( A \times B = 0 \)

15 \( |A \times B|^2 = (a_1^2 + a_2^2)(b_1^2 + b_2^2) - (a_1 b_1 + a_2 b_2)^2 = (a_1 b_2 - a_2 b_1)^2 \)

17 \( F; T; F; T \)

19 \( N = (2,1,0) \) or \( 21 + j \)

21 \( x - y + z = 2 \) so \( N = i - j + k \)

23 \( (1,2,1) - (2,1,1) = N = i + j + k; x + y + z = 4 \)

25 \( (1,1,1) \) \( (a, b, c) = N = (c - b) + (a - c)j + (b - a)k \); points on a line if \( a = b = c \) (many planes)

27 \( N = i + j, x + y = \) constant

31 \( i - 3j; -i + 3j; -3i - j \)

33 \( i - 3j; -i + 3j; -3i - j \)

35 \( 1,4,9 \)

39 \( |a_2 b_3 - a_3 b_2| \)

41 \( \text{area}^2 = \left( \frac{1}{2} ab \right)^2 + \left( \frac{1}{2} ac \right)^2 + \left( \frac{1}{2} bc \right)^2 = \left( \frac{1}{2} |A \times B| \right)^2 \) when \( A = ai - bj; B = ai - ck \)

43 \( A = \frac{1}{2}(2 \cdot 1 - (-1)1) = \frac{3}{2} \); fourth corner can be \( (3,3) \)

45 \( a_1 i + a_2 j \) and \( b_1 i + b_2 j; |a_1 b_2 - a_2 b_1|; A \times B = \cdots = (a_1 b_2 - a_2 b_1)k \)

47 \( A \times B \); from Eq. (6), \( (A \times B) \times i = -a_3 b_1 - a_1 b_3) + (a_1 b_2 - a_2 b_1)j; (A \cdot i)B - (B \cdot i)A = a_1 (b_1 i + b_2 j + b_3 k) - b_1 (a_1 i + a_2 j + a_3 k) \)

49 \( N = (Q - P) \times (R - P) = i + j + k; \) area \( \frac{1}{3} \sqrt{3}; x + y + z = 2 \)

Section 11.4 Matrices and Linear Equations

1 \( z = 5, y = 2, D = -2 \)

3 \( x = 3, y = 1, 8 \)

5 \( x = 2y, y = \) anything, \( D = 0, 2y \)

7 \( \text{no solution}, D = 0 \)

9 \( x = \frac{1}{D} 8 \) \( -1 \) \( 0 \) \( -3 \) \( = \frac{24}{8} = 3, y = \frac{1}{D} 3 \) \( 8 \) \( 0 \) \( 11 \) \( 0 \)
15 \( ad - bc = -2 \) so \( A^{-1} = \begin{bmatrix} .5 & .5 \\ .5 & -5 \end{bmatrix} \)

17 Are parallel; multiple; the same; infinite

19 Multiples of each other; in the same direction as the columns; infinite

21 \( d_1 = .34, d_2 = 4.91 \)

23 \( .96x + .02y = .58, .04x + .98y = 4.92; D = .94, x = .5, y = 5 \)

25 \( a = 1 \) gives any \( x = -y; a = -1 \) gives any \( x = y \)

27 \( D = -2, A^{-1} = \begin{bmatrix} -1 & -4 \\ -3 & 2 \end{bmatrix}; D = -8, (2A)^{-1} = \frac{1}{2} A^{-1}; D = \frac{1}{-2}, (A^{-1})^{-1} = \text{original } A; \)

\( D = -2 \) (not +2), \( (-A)^{-1} = -A^{-1}; D = 1, I^{-1} = I \)

29 \( AB = \begin{bmatrix} 7 & 5 \\ 5 & 1 \end{bmatrix}, BA = \begin{bmatrix} 5 & 11 \\ 3 & 5 \end{bmatrix}, BC = \begin{bmatrix} 3 & 5 \\ 1 & 3 \end{bmatrix}, CB = \begin{bmatrix} 4 & 2 \\ 2 & 2 \end{bmatrix} \)

31 \( AB = \begin{bmatrix} ae + bg & af + bh \\ ce + dg & cf + dh \end{bmatrix}, aecf + aedh + bgcf + bgdh \)

\( -afce - afdg - bhec - bhfdg = (ad - bc)(eh - fg) \)

33 \( A^{-1} = \begin{bmatrix} 1 & -2 \\ 0 & \frac{1}{2} \end{bmatrix}, B^{-1} = \begin{bmatrix} \frac{1}{2} & -1 \\ 0 & 1 \end{bmatrix} \)

35 Identity; \( B^{-1} A^{-1} \)

37 Perpendicular; \( u = v \times w \)

39 Line 4 + \( t \), errors \(-1, 2, -1 \)

41 \( d_1 - 2d_2 + d_3 = 0 \)

43 \( A^{-1} \) can't multiply \( O \) and produce \( u \)

Section 11.5 Linear Algebra (page 443)

1 \[
\begin{bmatrix}
1 & 2 & 2 \\
2 & 3 & 5 \\
0 & 2 & 2
\end{bmatrix}
\begin{bmatrix}
x \\
y \\
z
\end{bmatrix}
= \begin{bmatrix}
0 \\
0 \\
8
\end{bmatrix}
\]

3 \[
\begin{bmatrix}
1 & -1 & 0 \\
0 & 1 & -1 \\
1 & 0 & -1
\end{bmatrix}
\begin{bmatrix}
x \\
y \\
z
\end{bmatrix}
= \begin{bmatrix}
0 \\
0 \\
1
\end{bmatrix}
\]

5 \( \det A = 0, \text{ add 3 equations } \rightarrow 0 = 1 \)

7 \( 5a + 1b + 0c = d, A^{-1} = \begin{bmatrix}
1 & 0 & 0 \\
0 & 1/2 & 0 \\
0 & 0 & 1/3
\end{bmatrix} \)

9 \( b \times c; a \cdot b \times c = 0; \text{ determinant is zero} \)

11 \( 6, 2, 0; \text{ product of diagonal entries} \)

13 \( A^{-1} = \begin{bmatrix}
1 & -2 & 0 \\
0 & 1/2 & -1 \\
0 & 0 & 1/3
\end{bmatrix}, B^{-1} = \begin{bmatrix}
0 & 2 & -1/2 \\
0 & -3 & 1 \\
1 & 0 & 0
\end{bmatrix} \)

15 \( \text{Zero; same plane; } D \text{ is zero} \)

17 \( d = (1, -1, 0); u = (1, 0, 0) \text{ or } (7, 3, 1) \)

19 \( AB = \begin{bmatrix}
8 & 4 & 1 \\
40 & 26 & 0 \\
18 & 12 & 0
\end{bmatrix}, \text{ det } AB = 12 = (\text{det } A) \times (\text{det } B) \)

21 \( A + C = \begin{bmatrix}
2 & 3 & -3 \\
-1 & 4 & 6 \\
0 & -1 & 6
\end{bmatrix}, \text{ det}(A + C) \text{ is not det } A \text{ + det } C \)

23 \( p = \frac{(2)(3) - (0)(6)}{6} = 1, q = \frac{-(4)(3) + (0)(0)}{6} = -2 \)

25 \( (A^{-1})^{-1} \text{ is always } A \)

27 \(-1, -1, 1, 1; (y, x, z), (x, y, x), (y, z, x), (x, x, y) \)

29 \( 2! = 2, 4! = 24 \)

31 \( z = \frac{1}{2}, y = -\frac{3}{2}, x = 3; z = \frac{7}{2}, y = \frac{3}{2}, x = -\frac{1}{2} \)

33 \( \text{New second equation } 3z = 0 \text{ doesn't contain } y; \text{ exchange with third equation; there is a solution} \)

35 \( \text{Pivots }, 1, 2, 4, D = 8; \text{ pivots } 1, -1, 2, D = -2 \)

37 \( a_{12} = 1, a_{21} = 0, \sum a_{ij} b_{jk} = \text{ row } i, \text{ column } k \text{ in AB} \)

39 \( a_{11} a_{22} - a_{12} a_{21} \neq 0; D = 0 \)

CHAPTER 12 MOTION ALONG A CURVE

Section 12.1 The Position Vector (page 452)

1 \( v(1) = i + 3j; \text{ speed } \sqrt{10}; \)

3 \( \frac{dy}{dx} = \frac{dy}{dt} = \frac{\cos t}{\sin t}; \text{ tangent to circle is perpendicular to } \frac{x}{y} = \frac{\cos t}{\sin t} \)

5 \( v = e^t \text{ i } - e^{-t} \text{ j } = i - j; y - 1 = -(x - 1); xy = 1 \)
Answers to Odd-Numbered Problems

A-21

7 \( R = (1, 2, 4) + (4, 3, 0)t; R = (1, 2, 4) + (8, 6, 0)t; R = (5, 5, 4) + (8, 6, 0)t \)
8 \( R = (2 + t, 3, 4 - t); R = (2 + \frac{t^2}{2}, 3, 4 - \frac{t^2}{2}); \) the same line
10 Line; \( y = 2 + 2t, z = 2 + 3t; y = 2 + 4t, z = 2 + 6t \)
11 Line; \( \sqrt{36} + 9 + 4 = 7; (6, 3, 2); \) line segment
12 \( \sqrt{2}; 1; \sqrt{2} \)
13 \( x = t, y = mt + b \)
14 \( \mathbf{v} = 1 - \frac{1}{2}j, |\mathbf{v}| = \sqrt{1 + t^2}, \mathbf{T} = \mathbf{v}/|\mathbf{v}|; \mathbf{v} = (\cos t - t \sin t)\mathbf{i} + (\sin t + t \cos t)\mathbf{j}; |\mathbf{v}| = \sqrt{1 + t^2}; \)
15 \( \mathbf{T} = \mathbf{v}/|\mathbf{v}|; \mathbf{v} = i + 2j + 2k, |\mathbf{v}| = 3, T = \frac{1}{3} \mathbf{v} \)
16 \( R = - \sin t \mathbf{i} + \cos t \mathbf{j} + \mathbf{k}; \) same \( R \) plus any \( w \)
17 \( \mathbf{v} = (1 - \sin t)\mathbf{i} + (1 - \cos t)\mathbf{j}; |\mathbf{v}| = \sqrt{3 - 2 \sin t - 2 \cos t}, |\mathbf{v}|_{\text{min}} = \sqrt{3 - 2 \sqrt{2}}, |\mathbf{v}|_{\text{max}} = \sqrt{3 + 2 \sqrt{2}}; \)
18 \( a = - \cos t \mathbf{i} + \sin t \mathbf{j}, |a| = 1; \) center is on \( x = t, y = t \)
19 Leaves at \( \left( \sqrt{2}, \sqrt{2} \right); \mathbf{v} = (-\sqrt{2}, \sqrt{2}); \mathbf{R} = (\sqrt{2}, \sqrt{2}) + v(t - \pi) \)
20 \( R = \cos \frac{\sqrt{2}}{2} \mathbf{i} + \sin \frac{\sqrt{2}}{2} \mathbf{j} + \frac{1}{\sqrt{2}} \mathbf{k} \)
21 \( v = \sec^2 t \mathbf{i} + \sec t \tan t \mathbf{j}; |\mathbf{v}| = \sec^2 t \sqrt{1 + \sin^2 2}; a = 2 \sec^2 t \tan t \mathbf{i} + (\sec^3 t + \sec t \tan^2 t) \mathbf{j}; \)
curve is \( y^2 - x^2 = 1; \) hyperbola has asymptote \( y = x \)
22 If \( T = \mathbf{v} \) then \( |\mathbf{v}| = 1; \) line \( R = ti \) or helix in Problem 27
23 \( (x(t), y(t)) = (2t, 0) \quad 0 \leq t \leq \frac{1}{2} \quad (3 - 2t, 1) \quad 1 \leq t \leq 3 \quad (0, 4 - 2t) \quad \frac{3}{2} \leq t \leq 2 \)
24 \( x(t) = 4 \cos \frac{t}{2}, y(t) = 4 \sin \frac{t}{2} \quad 37 \mathbf{F}; \mathbf{F}; \mathbf{T}; \mathbf{T}; \mathbf{F} \quad 39 \frac{v}{x} = \tan \theta \) but \( \frac{v}{x} \neq \tan t \)
25 \( v \) and \( w; \) \( v \) and \( w \) and \( u; \) not zero
26 \( u = (8, 3, 2); \) projection perpendicular to \( \mathbf{v} = (1, 2, 2) \) is \( (6, -1, -2) \) which has length \( \sqrt{41} \)
27 \( x = G(t), y = F(t); y = x^2/3; t = 1 \) and \( t = -1 \) give the same \( x \) so they would give the same \( y; y = G(F^{-1}(x)) \)

Section 12.2 Plane Motion: Projectiles and Cycloids (page 457)

1 \( a \quad T = 16/g \sec, R = 128\sqrt{3}/g \text{ ft, } Y = 32/g \text{ ft } \quad b \quad 16\sqrt{2}/g, 128\sqrt{3}/g, 96/g \quad c \quad 32/g, 0, 128/g \quad 3 \quad x = 1.2 \) or 33.5
2 \( y = x - \frac{1}{2}x^2 = 0 \) at \( x = 2; y = x \tan x - \frac{1}{2} \left( \frac{x}{v_0 \cos \alpha} \right)^2 = 0 \) at \( x = R \quad 7 \quad x = v_0 \sqrt{\frac{2h}{g}} \)
3 \( v_0 \approx 11.2, \tan \alpha \approx 4.32 \quad 11 \quad v_0 = \sqrt{gR} = \sqrt{980} \) m/\sec; \) larger \( 13 \quad v_0^2/2g = 40 \) meters
4 \( \mathbf{R} \) and \( H \) by \( 4; dH = 2v^2 \cos 2\alpha \alpha /g, dH = v^2 \sin \alpha \alpha \cos \alpha \alpha \alpha /g \)
5 \( t = \frac{12\sqrt{2}}{10} \) sec; \( y = 12 - \frac{144g}{100} \approx -2.1 \) m; \( +.21 \) m \quad 19 \quad T = \left( \frac{1 - \cos \theta}{1 + \sin \theta} \right) \)
6 \( \frac{1}{2}\sqrt{1 - \cos \theta} \quad 17 \quad \kappa = \frac{3}{16}, N = i \quad 19 \quad (0, 0); (-3, 0) \) with \( \frac{1}{\kappa} = 4; (-1, 2) \) with \( \frac{1}{\kappa} = 2\sqrt{2} \)
20 Top of circle \( 25 ca(1 - \cos \theta), ca \sin \theta; \theta = \pi, \frac{\pi}{2} \quad 27 \) After \( \theta = \pi; x = x \alpha + v_0 t \) and \( y = 2a - \frac{1}{2}gt^2 \quad 29 \) 2; 3
21 \( 64\pi g^2; 5\pi a^2 \frac{3}{3} \quad 33 \quad x = \cos \theta + \sin \theta, y = \sin \theta - \cos \theta \quad 35 \quad (a = 4) \quad 6\pi \)
22 \( \mathbf{y} = 2 \sin \theta - \sin 2\theta = 2 \sin \theta(1 - \cos \theta); x^2 + y^2 = 4(1 - \cos \theta)^2; r = 2(1 - \cos \theta) \)

Section 12.3 Curvature and Normal Vector (page 463)

1 \( e^{\frac{t}{(1 + t^2)^{3/2}}} \quad 3 \quad \frac{1}{2} \quad 5 \quad 0 \quad \text{(line) } \quad 7 \quad \frac{2 + t^2}{(1 + t^2)^{3/2}} \quad 9 \quad (-\sin t^2, \cos t^2); (-\cos t^2, -\sin t^2) \)
11 \( \cos t, \sin t; (-\sin t, -\cos t) \quad 13 \quad (-\frac{3}{5} \sin t, \frac{3}{5} \cos t, \frac{2}{5}); |\mathbf{v}| = 5, \kappa = \frac{3}{25} \frac{5}{3} \text{ longer; tan \( \theta = \frac{3}{4} \frac{4}{3} \)} \)
15 \( 2\sqrt{2a(1 - \cos \theta)} \quad 17 \quad \kappa = \frac{3}{16}, N = i \quad 19 \quad (0, 0); (-3, 0) \) with \( \frac{1}{\kappa} = 4; (-1, 2) \) with \( \frac{1}{\kappa} = 2\sqrt{2} \)
21 Radius \( \frac{1}{2}, \) center \( (1, \pm \sqrt{\frac{1}{x^2} - 1}) \) for \( \kappa \leq 1 \quad 23 \quad \mathbf{U} \cdot \mathbf{V} \quad 25 \quad \frac{1}{\sqrt{2}}(\sin t i - \cos t j + k) \quad 27 \quad \frac{1}{2} \)
27 \( 29 \quad \mathbf{N} \) in the plane, \( \mathbf{B} = \mathbf{k}, r = 0 \quad 31 \quad \frac{d^2y}{dx^2} \quad 33 \quad a = 0 \quad T + 5\omega^2 \mathbf{N} \quad 35 \quad a = \frac{4}{\sqrt{1 + t^2}} \mathbf{T} + \frac{2 + t^2}{\sqrt{1 + t^2}} \mathbf{N} \)
37 \( a = \frac{4t}{\sqrt{1 + 4t^2}} \mathbf{T} + \frac{2}{\sqrt{1 + 4t^2}} \mathbf{N} \quad 39 \quad |\mathbf{F}^2 + (2F'')[x] - \frac{F''}{(F^2 + F'')^{3/2}} \)
Section 12.4 Polar Coordinates and Planetary Motion  

1 $ij; -i + j = u_r - u_\theta$  
2 $(2, -1); (1, 2)$  
3 $5 \mathbf{v} = 3e^3(u_r + u_\theta) = 3e^3(cos 3 - sin 3)i + 3e^3(sin 3 + cos 3)j$  
4 $\mathbf{v} = -20 \sin 5t + 20 \cos 5t; \mathbf{j} = 20 \mathbf{T} = 20 \mathbf{u}_r; \mathbf{a} = -100 \cos 5t i - 100 \sin 5t j = 100 \mathbf{N} = -100 \mathbf{u}_r$  
5 $r \frac{d^2 \theta}{dt^2} + 2 \frac{dr}{dt} \frac{d \theta}{dt} = 0$  
6 $11 \frac{d \theta}{dt} = 0 = \frac{1}{r} \frac{d}{dt}(r \frac{d \theta}{dt})$  
7 $\frac{d \theta}{dt} = 0.0004 \text{ radians/sec}; h = \frac{r^2 \frac{d \theta}{dt}}{40,000}$  
8 $\mathbf{r} = (\frac{2\pi}{10})^2 \mathbf{u}_r - 4(\frac{2\pi}{10}) \mathbf{u}_\theta; \int_0^\infty |v| \, dt$  

CHAPTER 13 PARTIAL DERIVATIVES  

Section 13.1 Surfaces and Level Curves  

3 $x$ derivatives $\infty, -1, -2, -4e^{-4}$ (flattest)  
4 Straight lines  
5 Logarithm curves  
6 Parabolas  
7 No: $f(x, y) = (x + y)^4$ or $(ax + by)^4$ or any function of $ax + by$  
8 Ellipses $4x^2 + y^2 = c^2$  
9 Saddle  
10 Ellipses $5x^2 + y^2 = c^2 + 4cx + x^2$  
11 Ellipses not reaching $(1, 2)$  
12 Center $(1, 1); f = x^2 + y^2 - 1$  
13 Four, three, planes, spheres  
14 Less than 1, equal to 1, greater than 1  
15 Parallel lines, hyperbolas, parabolas  

Section 13.2 Partial Derivatives  

1 $f_x = \frac{1}{y} f_{xy} + \frac{1}{x} f_{y} - \frac{1}{y} f_{x} + \frac{1}{x} f_{y}$  
2 $f_{xy} = \frac{2}{y^2} f_{x} + \frac{2}{x^2} f_{y}$  
3 $f_{xy} = \frac{2}{y^2} f_{x} + \frac{2}{x^2} f_{y}$  
4 $f_{xy} = \frac{2}{y^2} f_{x} + \frac{2}{x^2} f_{y}$  
5 $f_{xx} = 6y^3, f_{yy} = 6x^3, f_{xxy} = f_{xyx} = f_{xzz} = 18x^2 y, f_{yyx} = f_{xyy} = f_{yzz} = 18x^2 y$  
6 $g(y) = Ae^{cv/2} + Be^{-cv/2}$  
7 $f_t = -2f, f_{xx} = f_{yy} = -e^{-2t} \sin x \sin y; e^{-13t} \sin 2x \sin 3y$  
8 $\sin(x + ct)$ moves left  
9 $\sin(x + ct), \cos(x + ct), e^x - c^t$  

43 $(B - A)h_y(C^*) = (B - A)[f_y(b, C^*) - f_y(a, C^*)] = (B - A)(b - a)f_{yz}(e^*, C^*); \text{continuous } f_{xy} \text{ and } f_{yz}$  
44 $\mathbf{y}$ converges to $b; \text{inside and stay inside}; d_n = \sqrt{(x_n - a)^2 + (y_n - b)^2} \to \text{zero}; d_n < \epsilon \text{ for } n > N$  
45 $\epsilon$, less than $\delta$  

47 $f(a, b) = \frac{1}{x - 1} \frac{1}{y - 2}$  

51 $(0, 0) = 1; f(0, 0) = 1; \text{not defined for } x < 0$
Answers to Odd-Numbered Problems

Section 13.3 Tangent Planes and Linear Approximations (page 488)

1 \[ z - 1 = y - 1; N = j - k \]
2 \[ 3 z - 2 = \frac{1}{3} (x - 6) - \frac{2}{3} (y - 3); N = \frac{1}{3} i - \frac{2}{3} j - k \]
5 \[ 2(x - 1) + 4(y - 2) + 2(z - 1) = 0; N = 2i + 4j + 2k \]
7 \[ z - 1 = x - 1; N = i - k \]
9 Tangent plane \[ 2z_0(x - x_0) - 2x_0(y - y_0) = 0; (0, 0, 0) \] satisfies this equation because \[ x_0^2 - x_0^2 - y_0^2 = 0 \] on the surface; \[ \cos \theta = \frac{N \cdot \langle x_0, y_0, z_0 \rangle}{|N||\langle x_0, y_0, z_0 \rangle|} = \frac{-1}{\sqrt{2}} \] (surface is the 45° cone)

11 \[ dx = 3dx - 2dy \] for both; \[ dx = 0 \] for both; \[ \Delta x = 0.0029 \] for \([x^3]/y^2\); tangent plane

13 \[ z = x_0 + F_t \]; plane \[ 6(x - 4) + 12(y - 2) + 8(z - 3) = 0; \] normal line \[ x = 4 + 6t, y = 2 + 12t, z = 3 + 8t \]

15 Tangent plane \[ 4(x - 2) + 2(y - 1) + 4(z - 2) = 0; \] normal line \[ x = 2 + 4t, y = 1 + 2t, z = 2 + 4t; (0, 0, 0) \] at \( t = -\frac{1}{2} \)

17 \[ dw = y_0 dx + x_0 dy \]; product rule; \[ \Delta w - dw = (x - x_0)(y - y_0) \]

19 \[ dI = 4000dR + 0.08dP; dP = $100; I = (7.8)(4100) = $319.80 \]

21 Increase \[ = \frac{25}{100} - \frac{25}{100} = \frac{3}{404} \]; decrease \[ = \frac{25}{100} - \frac{25}{100} = \frac{1}{404} \]; \[ dA = \frac{1}{y} dx - \frac{x}{y} dy; 3 \]

23 \[ \Delta \theta \approx \frac{v_{z^+z^+} + v_{z^-z^-}}{\sqrt{y^2 + y^2}} \]

25 \[ Q \] increases; \[ Q_s = -\frac{250}{3}, Q_t = \frac{5}{3}, P_s = -2Q_s, P_t = -2Q_t = \frac{1}{3}; Q = 50 - \frac{250}{3} (s - \frac{4}{5} - \frac{3}{5} (t - 10) \]

27 \[ s = 1, t = 10 \] gives \[ Q = 40; P_s = -Q_s = sQ_s = Q_s + 40; P_t = -Q_t = sQ_t + 1 = Q_t + 1 \] \[ Q_s = -20, Q_t = -\frac{1}{2}, P_s = 20, P_t = \frac{1}{2} \]

29 \[ x = 1, y = \frac{1}{2}, z = 0 \]

31 \[ \Delta x = -\frac{1}{2}, \Delta y = \frac{1}{2}, x_1 = \frac{1}{2}, y_1 = -\frac{1}{2}, \text{line } x + y = 0 \]

33 \[ 3a^2 \Delta x - \Delta y = -a - a^3 \] gives \[ \Delta y = -\Delta x = \frac{a + a^3}{1 + 3a^2}; \] lemon starts at \( (1/\sqrt{3}, -1/\sqrt{3}) \)

35 \[ \Delta x + 3a^2 \Delta y = a + a^3 \]

37 \[ \Delta x = -x_0 + 1, \Delta y = -y_0 + 2, (x_1, y_1) = (1, 2) \]

39 \[ G = H = \frac{x^2}{2x-1} \]

41 \[ J = \begin{bmatrix} e^x & 0 \\ 0 & e^y \end{bmatrix}, \Delta x = -1 + e^{-x}, \Delta y = -1 - (x_1 - 1 + e^{-x})e^{-y} \]

43 \[ (x_1, y_1) = (0, \frac{5}{4}), (\frac{3}{4}, \frac{5}{4}), (0, 0), (0, \frac{5}{4}) \]

Section 13.4 Directional Derivatives and Gradients (page 495)

1 \[ \text{grad } f = 2zi - 2yj, D_uf = (5\sqrt{x} - y, D_uf(P) = \sqrt{3} \]

3 \[ \text{grad } f = e^x \cos y i - e^x \sin y j, D_uf = -e^x \sin y, D_uf(P) = -1 \]

5 \[ f = \sqrt{x^2 + (y - 3)^2}, \text{grad } f = \frac{x}{f} i + \frac{y - 3}{f} j, D_uf = \frac{\vec{e}}{f}, D_uf(P) = \frac{1}{\sqrt{6}} \]

7 \[ \text{grad } f = \frac{2x}{x^2 + y^2} i + \frac{2y}{x^2 + y^2} j \]

9 \[ \text{grad } f = 6zi + 4yj, \text{level curve is point } P; \text{level direction } -8i + 6j \]

11 \[ T; F \] (grad is a vector); \[ F, T \]

13 \[ u = \left( \frac{a}{\sqrt{a^2 + b^2}}, \frac{b}{\sqrt{a^2 + b^2}} \right), D_uf = \sqrt{a^2 + b^2} \]

15 \[ \text{grad } f = (e^{-x} - e^{-y}, -e^{-x} - e^{-y}) = (e^{-1} - e^{-1}) \] at \( P; u = \left( \frac{1}{\sqrt{2}}, \frac{-1}{\sqrt{2}} \right), D_uf = \sqrt{2e^{-1}} \]

17 \[ \text{grad } f = \text{O at maximum; level curve is point } \]

19 \[ \text{gradient is one point } \]

21 \[ \text{Direction } -U = (-2, 0, -4) \]

23 \[ -U = \left( \frac{-x}{\sqrt{1 - x^2 - y^2}}, \frac{y}{\sqrt{1 - x^2 - y^2}}, \frac{z^2}{1 - x^2 - y^2} \right) \]

25 \[ f = (x + 2y) \text{ and } (x + 2y)^2; i + 2j \] straight lines \( x + 2y = \text{constant} \) (perpendicular to \( i + 2j \))

27 \[ \text{grad } f = \pm \left( \frac{\sqrt{2}}{\sqrt{3}}, \frac{-\sqrt{2}}{\sqrt{3}} \right) \]; \[ g = \pm (2\sqrt{5}, \sqrt{5}), f = \pm \left( \frac{x}{\sqrt{5}} - \frac{2y}{\sqrt{5}} \right) + C, \] \[ g = \pm \left( 2\sqrt{5}x + \sqrt{5}y \right) + C \]

29 \[ \theta = \text{constant along ray in direction } u = \frac{3i + 4j}{5}; \] \[ \text{grad } \theta = \frac{-y + z}{x^2 + y^2} = \frac{3i + 4j}{25} \]; \[ \text{grad } \theta = 0 \]

31 \[ U = (f_x, f_y, f_x^2 + f_y^2) = (-1, -2, 5), \text{tangent at the point } (2, 1, 6) \]

33 \[ \text{grad } f \text{ toward } 2i + j \text{ at } P; j \text{ at } Q, \text{2i} + j \text{ at } R; \left( 2, \frac{1}{2} \right) \text{ and } \left( 2, \frac{1}{2} \right); \text{largest upper left, smallest lower right}; z_{\text{max}} > 9; \text{z} \text{ goes from 2 to 8 and back to 6} \]
Section 13.5 The Chain Rule  (page 503)

1 \( f_y = cf_x = c \cos(x + cy) \)  
3 \( f_y = 7f_x = 7e^{x+y} \)  
5 \( 3g^2 \frac{\partial u}{\partial x} \frac{dx}{dt} + 3g^2 \frac{\partial u}{\partial y} \frac{dy}{dt} \)  
7 Moves left at speed 2

9 \( \frac{dx}{dt} = 1 \) (wave moves at speed 1)

11 \( \frac{\partial^2 f}{\partial x \partial y} f(x + iy) = f''(x + iy), \frac{\partial^2 f}{\partial y \partial x} f(x + iy) = i^2 f''(x + iy) \)

so \( f_{xx} + f_{yy} = 0; (x + iy)^2 = (x^2 - y^2) + i(2xy) \)

13 \( \frac{df}{dt} = 2x(1) + 2y(2t) = 2t + 4t^3 \)  
15 \( \frac{df}{dt} = y \frac{dx}{dt} + x \frac{dy}{dt} = -1 \)  
17 \( \frac{df}{dt} = \frac{1}{x+y} \frac{dx}{dt} + \frac{1}{x+y} \frac{dy}{dt} = 1 \)

19 \( V = \frac{1}{3} \pi r^2 h, \frac{dV}{dt} = \frac{2\pi rh}{3} \frac{dx}{dt} + \frac{\pi r^2}{3} \frac{dh}{dt} = 36 \pi \)

21 \( \frac{dP}{dt} = \frac{90}{\sqrt{90^2 + 90^2}}(60) + \frac{90}{\sqrt{90^2 + 90^2}}(45) = \frac{105}{\sqrt{2}} \) mph; \( \frac{dP}{dt} = \frac{60}{\sqrt{45^2 + 60^2}}(60) + \frac{45}{\sqrt{45^2 + 60^2}}(45) \approx 74 \) mph

23 \( \frac{df}{dt} = u_1 \frac{\partial f}{\partial x} + u_2 \frac{\partial f}{\partial y} + u_3 \frac{\partial f}{\partial z} \)  
25 \( \frac{df}{dt} = 1 \) with \( x \) and \( y \) fixed; \( \frac{df}{dt} = 6 \)

27 \( f_t = f_x t + f_y (2t) \); \( f_t = f_{xt} t + f_x + 2f_y t + 2f_y = (f_{xx} t + f_y (2t)) t + f_x + 2(f_{xy} t + f_y (2t)) t + 2f_y \)

29 \( \frac{df}{dr} = \frac{df}{dx} \frac{dx}{dr} + \frac{df}{dy} \frac{dy}{dr} = \frac{\partial f}{\partial x} \cos \theta + \frac{\partial f}{\partial y} \sin \theta, \theta \) is fixed

31 \( r_{xx} = \frac{1}{(x^2 + y^2)^{3/2}} \)  
33 \( (\frac{\partial f}{\partial x})_y = \frac{1}{(x^2 + y^2)^{3/2}} \) \( \cos (x) \) \( \frac{\partial f}{\partial x} = 1 \); first answer is also \( \frac{1}{\sqrt{1 - \sin^2 x}} = \frac{1}{\cos x} \)

35 \( f_r = f_x \cos \theta + f_y \sin \theta, f_r = -f_x \sin \theta + f_y \cos \theta + f_{xx} (-r \sin \theta \cos \theta) + f_{xy} (-r \sin^2 \theta + r \cos^2 \theta) + f_{yy} (r \cos \theta \sin \theta) \)

37 Yes (with \( y \) constant): \( \frac{\partial x}{\partial y} = ye^{x+y}, \frac{\partial x}{\partial z} = \frac{1}{y} \)  
39 \( f_t = f_x x_t + f_y y_t; f_t = f_{xx} x_t^2 + 2f_{xy} x_t y_t + f_{yy} y_t^2 \)

41 \( \frac{df}{dz} = \frac{df}{dx} \frac{dx}{dz} + \frac{df}{dy} \frac{dy}{dz} = a - \frac{2}{5} b; (\frac{df}{dz})_y = a; (\frac{df}{dz})_x = \frac{1}{5} b \)

43 \( y = y^2 \) so \( f_x = 0, f_y = 2y = 2r \sin \theta; r = 0 \) so \( f_r = 2r = 2 \sqrt{x^2 + y^2}, f_\theta = 0 \)

47 \( g_{uu} = f_{xx} u_x + f_{yy} u_y = f_x + f_y; g_{uv} = f_{xx} u_x + f_{yy} u_y = f_x + f_y; g_{vv} = f_{xx} u_x + f_{xx} u_y + f_{yy} x_u + f_{yy} y_v \)

Section 13.6 Maxima, Minima, and Saddle Points  (page 512)

1 \((0,0)\) is a minimum  
3 \((3,0)\) is a saddle point  
5 No stationary points  
7 \((0,0)\) is a maximum

9 \((0,0,2)\) is a minimum  
11 All points on the line \( z = y \) are minima  
13 \((0,0)\) is a saddle point

15 \((0,0)\) is a saddle point; \((2,0)\) is a minimum; \((0,-2)\) is a maximum; \((2,-2)\) is a saddle point

17 Maximum of area \((x - 3y)\) is 12

19 \( 2(x + y) + 2(x + 2y - 5) + 2(x + 2y - 5) = 0 \)  
21 Minimum at \((0, \frac{1}{2}); (0,1) \)
Answers to Odd-Numbered Problems

23 \frac{df}{dt} = 0 \text{ when } t \tan t = \sqrt{3}; f_{\text{max}} = 2 \text{ at } \left(\frac{1}{2}, \frac{\sqrt{3}}{2}\right), f_{\text{min}} = -2 \text{ at } \left(-\frac{1}{2}, -\frac{\sqrt{3}}{2}\right)

25 (ax + by)_{\text{max}} = \sqrt{a^2 + b^2}; (x^2 + y^2)_{\text{min}} = \frac{1}{a^2 + b^2}

27 0 < c < \frac{1}{4}

29 The vectors head-to-tail form a 60-60-60 triangle. The outer angle is 120°

31 \sqrt{2} + \sqrt{3}; 1 + \sqrt{3}; 1 + \frac{\sqrt{3}}{10}

35 Steiner point where the arcs meet

39 Best point for \( p = \infty \) is equidistant from corners

41 \text{grad} f = (\sqrt{2} \frac{\partial}{\partial x} + \frac{\partial}{\partial y}) \ln(1 - xy)

43 Third derivatives all 0; \( f = \frac{6}{5} x^5 + \frac{6}{5} x^2 y^3 + \frac{6}{5} y^3 \)

45 \left(\frac{\partial}{\partial x}\right)^n \left(\frac{\partial}{\partial y}\right)^m \ln(1 - xy)_{|0,0} = n!(n - 1)! \text{ for } m = n > 0, \text{ other derivatives zero}

47 All derivatives are \( e^2 \) at (1,1); \( f \approx e^2 (1 + (x - 1) + (y - 1) + \frac{1}{2} (x - 1)^2 + (x - 1)(y - 1) + \frac{1}{2} (y - 1)^2) \)

49 \( x = 1, y = -1 : f_x = 2, f_y = -2, f_{xx} = 2, f_{xy} = 0, f_{yy} = 2; \) series must recover \( x^2 + y^2 \)

51 \text{Line } x - 2y = \text{constant}; x + y = \text{constant}

53 \frac{x^2}{2} f_{xx} + xy f_{xy} + \frac{y^2}{2} f_{yy} \text{ at } (0,0) ; f_{xx} > 0 \text{ and } f_{xx} f_{yy} > f_{xy}^2 \text{ at } (0,0); f_x = f_y = 0

55 \Delta x = -1, \Delta y = -1

57 \( f = x^2 (12 - 4z) \) has \( f_{\text{max}} = 16 \text{ at } (2,4); \) line has slope \(-4, y = \frac{16}{5} \) has slope \(-\frac{32}{8} = -4 \)

59 If the fence were not perpendicular, a point to the left or right would be closer

Section 13.7 Constraints and Lagrange Multipliers

1 \( f = x^2 + (k - 2x)^2; \frac{df}{dx} = 2x - 4(k - 2x) = 0; \left(\frac{2k}{5}, \frac{k^2}{5}\right) \)

3 \( \lambda = -4, x_{\text{min}} = 2, y_{\text{min}} = 2 \)

5 \( \lambda = \frac{3}{4} (x, y) : (x, y) = (\pm 2^{1/6}, 0) \text{ or } (0, \pm 2^{1/6}), f_{\text{min}} = 2^{1/3}; \lambda = \frac{1}{3} : (x, y) = (\pm 1, \pm 1), f_{\text{max}} = 2 \)

7 \( \lambda = \frac{1}{2}, (x, y) = (2, -3); \) tangent line is \( 2x - 3y = 13 \)

9 \( (1 - c)^2 + (-a - c)^2 + (2 - a - b - c)^2 + (2 - b - c)^2 \) is minimized at \( a = -\frac{1}{2}, b = \frac{3}{2}, c = \frac{3}{4} \)

11 \( (1, -1) \text{ and } (-1, 1); \lambda = -\frac{1}{2} \)

13 \( f \) is not a minimum when \( C \) crosses to lower level curve; stationary point when \( C \) is tangent to level curve

15 Substituting \( \frac{\partial L}{\partial x} = \frac{\partial L}{\partial y} = \frac{\partial L}{\partial \lambda} = 0 \) and \( L = f_{\text{min}} \) leaves \( \frac{df_{\text{max}}}{dk} = \lambda \)

17 \( x^2 \) is never negative; \( (0,0), 1 = \lambda(-3y^2) \) but \( y = 0, g = 0 \) has a cusp at \( (0,0) \)

19 \( 2x = \lambda_1 + \lambda_2, 4y = \lambda_1, 2z = \lambda_1 - \lambda_2, x + y + z = 0, x - z = 1 \) gives \( \lambda_1 = 0, \lambda_2 = 1, f_{\text{min}} = \frac{1}{2} \text{ at } (\frac{1}{2}, 0, -\frac{1}{2}) \)

21 \( (1, 0, 0); (0, 1, 0); (\lambda_1, \lambda_2, 0); x = y = 0 \)

23 \( \frac{df}{dx} \text{ and } \frac{df}{dy}; \lambda = 0 \)

25 \( (1, 0, 0), (0, 1, 0), (0, 0, 1); \) at these points \( f = 4 \text{ and } -2 \) (min) and \( 5 \) (max)

27 By increasing \( k \), more points are available so \( f_{\text{max}} \) goes up. Then \( \lambda = \frac{df_{\text{max}}}{dk} \geq 0 \)

29 \( (0, 0); \lambda = 0; f_{\text{min}} \) stays at 0

31 \( 5 = \lambda_1 + \lambda_2, 6 = \lambda_1 + \lambda_3, \lambda_2 \geq 0, \lambda_3 \leq 0; \) subtraction \( 5 - 6 = \lambda_2 - \lambda_3 \) or \(-1 \geq 0 \) (impossible);

\[ x = 2004, y = -2000 \text{ gives } 5x + 6y = -1980 \]

33 \( 2x = 4\lambda_1 + \lambda_2, 2y = 4\lambda_1 + \lambda_3, \lambda_2 \geq 0, \lambda_3 \geq 0, 4x + 4y = 40; \) max area 100 at \( (10,0)(0,10); \) min 25 at \( (5,5) \)

CHAPTER 14 MULTIPLE INTEGRALS

Section 14.1 Double Integrals

1 \( \frac{8}{3} \); \( \frac{2}{3} \); \( \frac{3}{2} \); \( \ln \frac{3}{2} \); \( \frac{5}{2} \); \( \frac{7}{2} \); \( \frac{9}{4} \); \( \frac{11}{1} \)

11 \( \int_{y=1}^{2} \int_{x=1}^{2} f_{x=1} f_{y=2} dx dy + \int_{y=2}^{4} \int_{y=2}^{2} dx dy \)

13 \( \int_{y=0}^{1} \int_{x=-1}^{1} f_{x=0} \ln y \ dx \ dy \)

15 \( \int_{x=0}^{1} \int_{y=-\sqrt{x}}^{\sqrt{x}} f_{x=0} \ dy \ dx \)

17 \( \int_{x=0}^{1} \int_{y=x}^{(1/2)} f_{x=0} \ dy \ dx = \int_{0}^{x} \int_{0}^{1} \ dy \ dx = \frac{1}{4} \)

19 \( \int_{x=0}^{3} f_{y} \ dx \ dy = \int_{0}^{1} \int_{y=x}^{3} f_{x} \ dy \ dx + \int_{0}^{1} \int_{y=x}^{3} f_{x} \ dy \ dx \)

21 \( \int_{y=0}^{1} \int_{x=0}^{1} f_{y} \ dy \ dx + \int_{y=0}^{1} \int_{x=0}^{1} f_{y} \ dy \ dx = \int_{0}^{f} \int_{0}^{1} f_{x} \ dy \ dx \)

23 \( \int_{0}^{1} \int_{0}^{b} f_{x} \ dy \ dx + \int_{0}^{1} \int_{0}^{b} f_{x} \ dy \ dx = \int_{0}^{b} \int_{0}^{b} f_{x} \ dy \ dx = b \)

25 \( f(a, b) = f(a, 0) - f(0, b) + f(0, 0) \)
Section 14.2 Change to Better Coordinates (page 534)

1 \[ \int_{\pi/4}^{\pi} \int_0^1 r \, dr \, d\theta = \frac{\pi}{8} \quad 3 \quad \text{S is quarter-circle with } u \geq 0 \text{ and } v \geq 0; \int_0^1 \int_0^{\sqrt{1-v^2}} u \, du \, dv \]

5 \( R \) is symmetric across the y axis; \( \int_0^1 \int_0^{\sqrt{1-x^2}} u \, du \, dv = \frac{1}{2} \) divided by area gives \((u, v) = (4/3\pi, 4/3\pi)\)

7 \( 2 \int_0^1 \int_{1+x}^{1+\sqrt{1-x^2}} dy \, dx; \) \( y \text{ region } R^* \text{ becomes } R \text{ in the } x^*y^* \text{ plane; } dy \, dx = dx^*dy^* \) when region moves

9 \( J = \begin{vmatrix} \frac{\partial x}{\partial x^*} & \frac{\partial x}{\partial y^*} \\ \frac{\partial y}{\partial x^*} & \frac{\partial y}{\partial y^*} \end{vmatrix} = \left| \begin{array}{cc} \cos \theta^* & -r \sin \theta^* \\ r \sin \theta^* & r^* \cos \theta^* \end{array} \right| = r^*; \int_{\pi/4}^{\pi/2} \int_0^{r^*} r^* \, dr^* \, d\theta^* \)

11 \( I_y = \int_R x^2 \, dx \, dy = \int_{\pi/4}^{\pi/2} \int_0^1 r^2 \cos^2 \theta \, r \, dr \, d\theta = \frac{\pi}{16} - \frac{1}{8}; I_x = \frac{\pi}{16} + \frac{1}{8}; I_0 = \frac{\pi}{8} \)

13 \((0, 0), (1, 2), (3, 0), (0, 1); \) area of parallelogram is 1

15 \( x = u, y = u + 3v + uv; \) then \((u, v) = (1, 0), (1, 1), (0, 1)\) give corners \((x, y) = (1, 0), (1, 5), (0, 3)\)

17 Corners \((0, 0), (2, 1), (3, 3), (1, 2); \) sides \( y = \frac{1}{2}x, y = 2x - 3, y = \frac{1}{2}x + \frac{3}{4}, y = 2x \)

19 Corners \((1, 1), (e^2, e), (e^3, e^3), (e, e^2); \) sides \( x = y, x = e^2/e^3, y = x^2/e^3, y = e^2 \)

21 Corners \((0, 0), (1, 0), (1, 2), (0, 1); \) sides \( y = 0, x = 1, y = 1 + x^2, y = x = 0 \)

23 \( J = \begin{vmatrix} 1 & 1 \\ 1 & 2 \end{vmatrix} = 3, \text{ area } \int_0^1 \int_0^1 3 \, du \, dv = 3; J = \begin{vmatrix} 2e^{u+v} & e^{2u+v} \\ e^{u+2v} & 2e^{u+2v} \end{vmatrix} = 3e^{3u+3v}, \int_0^1 \int_0^1 3e^{3u+3v} \, du \, dv = \int_0^1 (e^{3u+3v} - 3e^v) \, dv = \frac{3}{4}(e^6 - 2e^3 + 1) \)

25 Corners \((x, y) = (0, 0), (1, 0), (1, 1), (1, f(1)), (0, f(0)); (1/2, 1) \) gives \( x = \frac{1}{2}, y = f(\frac{1}{2}); \)

27 \( B^2 = 2 \int_0^{\pi/4} \int_0^{1/\sin \theta} e^{-r^2/2} \, r \, dr \, d\theta = \frac{\pi}{8} \)

29 \( f = \int \int r^2 \, dr \, d\theta / \int \int r \, dr \, d\theta = \frac{r^3}{3} = \frac{28}{27} \int_0^{\pi/4} \int_0^{1/\sin \theta} \cos^3 \theta \, d\theta / \pi a^2 = \frac{32a}{9\pi} \quad 31 \quad \int_0^{2\pi} \int_0^1 r^2 \, r \, dr \, d\theta = \frac{\pi}{2} \)

33 Along the right side; along the bottom; at the bottom right corner

35 \( \int \int xy \, dx \, dy = \int_0^1 \int_0^1 (u \cos \alpha - v \sin \alpha)(u \sin \alpha + v \cos \alpha) \, du \, dv = \frac{1}{4}(\cos^2 \alpha - \sin^2 \alpha) \)

37 \( \int_0^{2\pi} \int_0^1 r^2 \, r^2 \, dr \, d\theta = \frac{2\pi}{6}(5^6 - 4^6) \quad 39 \quad x = \cos \alpha - \sin \alpha, y = \sin \alpha + \cos \alpha \) goes to \( u = 1, v = 1 \)

Section 14.3 Triple Integrals (page 540)

1 \( \int_0^1 \int_0^1 \int_0^y dx \, dy \, dz = \frac{1}{6} \)

3 \( 0 \leq y \leq x \leq z \leq 1 \) and all other orders \( xzy, yzx, zyx, zyx; \) all six contain \((0, 0, 0); \) to contain \((1,0,1)\)

5 \( \int_0^1 \int_0^1 \int_0^1 dx \, dy \, dz = 8 \quad 7 \quad \int_0^1 \int_0^1 \int_0^1 dx \, dy \, dz = 4 \quad 9 \quad \int_0^1 \int_0^1 \int_0^x dx \, dy \, dz = \frac{4}{3} \)

11 \( \int_0^1 \int_0^{2-z} \int_0^{2-y-2z} dx \, dy \, dz = \frac{2}{3} \)

13 \( \int_0^{1/2} \int_0^{2-2x} \int_0^{2-y-2z} dx \, dy \, dz = \frac{7}{12} \)

15 \( \int_0^1 \int_0^{1-z} \int_0^{(1-z)^2-y^2} dx \, dy \, dz = \frac{\pi}{3} \quad 17 \quad \int_0^6 \int_0^1 \int_0^{\sqrt{1-y^2}} dx \, dy \, dz = 6\pi \quad 19 \quad \int_0^1 \int_0^1 \int_0^{\sqrt{2-4y^2}} dx \, dy \, dz = \pi \)

21 Corner of cube at \((\frac{1}{\sqrt{3}}, \frac{1}{\sqrt{3}}, \frac{1}{\sqrt{3}}); \) sides \( \frac{2}{\sqrt{3}} \); area \( \frac{8}{3\sqrt{3}} \)

23 Horizontal slices are circles of area \( \pi r^2 = \pi(4 - x); \) volume = \( \int_0^4 \pi(4-x) \, dx = 8\pi; \) centroid has \( z = 0, y = 0, z = \int_0^4 \pi(4-x) \, dx / 8\pi = \frac{4}{3} \)
Answers to Odd-Numbered Problems

25 \( I = \frac{z^2}{2} \) gives zeros; \( \frac{\partial I}{\partial x} = \int_0^x f \, dy \, dx \), \( \frac{\partial I}{\partial y} = \int_0^x f \, dx \, dz \), \( \frac{\partial^2 I}{\partial y \partial x} = \int_0^x f \, dx \)

27 \( \int_{-1}^1 \int_{-1}^1 (y^2 + z^2) \, dx \, dy \, dz = \frac{16}{3} \); \( \iiint 2^2 \, dV = \frac{5}{3} \); \( \iiint (x - \frac{z + y^2}{2}) \, dV = \frac{16}{3} \)

29 \( \int_0^3 \int_0^y \int_0^z \, dx \, dy \, dz = 6 \)

31 Trapezoidal rule is second-order; correct for \( 1, x, y, z, xy, xz, yz, xyz \)

Section 14.4 Cylindrical and Spherical Coordinates (page 547)

1 \((r, \theta, z) = (D, 0, 0); (\rho, \phi, \theta) = (D, \frac{\pi}{2}, 0)\)

3 \((r, \theta, z) = (0, \text{any angle}, D); (\rho, \phi, \theta) = (D, 0, \text{any angle})\)

5 \((x, y, z) = (2, -2, 2\sqrt{2}); (\theta, z) = (2\sqrt{2}, -\frac{\pi}{4}, 2\sqrt{2})\)

7 \((x, y, z) = (0, 0, -1); (\theta, z) = (0, \text{any angle}, -1)\)

9 \(\phi = \tan^{-1}\left(\frac{z}{r}\right)\)

11 \(r = 45^\circ\) cone in unit sphere: \(\frac{2\pi}{3}(1 - \frac{1}{\sqrt{2}})\)

13 cone without top: \(\frac{2\pi}{3}\)

17 Hemispheres of radius \(\frac{2}{3}x^2\)

21 \(\pi(R^2 - z^2); 4\pi r \sqrt{R^2 - r^2}\)

23 \(a^3\) tan \(\alpha\) (see 8.1.39)

25 \(\cos \theta - r \sin \theta \quad 0\)

27 \(\sin \theta \quad r \cos \theta \quad 0\)

29 \(\cos \theta \quad r \cos \theta \quad 0\)

31 Wedges are not exactly similar; the error is higher order \(\Rightarrow\) proof is correct

33 Proportional to \(1 + \frac{1}{6}(\sqrt{a^2 + (D - h)^2} - \sqrt{a^2 + D^2})\)

35 \(\rho \pm \sqrt{3}; \frac{2}{3}\)

41 \(\rho^3, \rho^2\); force = 0 inside hollow sphere

CHAPTER 15 VECTOR CALCULUS

Section 15.1 Vector Fields (page 554)

1 \(f(x, y) = x + 2y\)

3 \(f(x, y) = \sin(x + y)\)

5 \(f(x, y) = \ln(x^2 + y^2) = 2 \ln r\)

7 \(F = xy \pm \frac{3}{2} j, f(x, y) = \frac{3}{2} y\)

9 \(\frac{\partial f}{\partial z} = 0 \) so \( f \) cannot depend on \( z \); streamlines are vertical \((y = \text{constant})\)

11 \(F = 3i + j\)

13 \(F = i + 2y j\)

15 \(F = 2x i - 2y j\)

17 \(F = e^{x-y} i - e^{x+y} j\)

19 \(\frac{dy}{dz} = -1; y = -x + C\)

21 \(\frac{dy}{dz} = \mp \frac{z}{y}; x^2 + y^2 = C\)

23 \(\frac{dy}{dz} = \mp \frac{z^2}{y^2} = \mp \frac{z}{y}; x^2 + y^2 = C\)

25 \(\parallel\)

27 \(F = \frac{5z}{r} + \frac{5y}{j}\)

29 \(F = -m M G (x + y j) - (m M G) j \alpha ((x - 1) i + y j)\)

31 \(F = 2y i - 2x j\)

33 \(\frac{dy}{dz} = \mp \frac{z}{y} = -1; \frac{dy}{dz} = \mp \frac{z}{\sqrt{z^2 - 3}} = 2\)

35 \(\frac{dy}{dz} = \frac{\partial f}{\partial r}, \frac{\partial f}{\partial r}; \frac{\partial f}{\partial y}, \frac{\partial f}{\partial y}; f(r) = C\) gives circles

37 \(T; F \) (no equipotentials); \(T; F \) (not multiple of \(xi + yj + zk\))

39 \(F\) and \(F + i\) and \(2F\) have the same streamlines (different velocities) and equipotentials (different potentials).

But if \( f \) is given, \( F \) must be grad \( f \).

Section 15.2 Line Integrals (page 562)

1 \(\int_0^1 \sqrt{1 + 2^2} \, dt = \sqrt{5} ; \int_0^1 2 \, dt = 2\)

3 \(\int_0^1 t^2 \sqrt{2^2} \, dt + \int_1^2 1 \cdot (2 - t) \, dt = \sqrt{2^2 + \frac{1}{2}}\)

5 \(\int_0^{2\pi} (-3 \sin t) \, dt = 0\) (gradient field); \(\int_0^{2\pi} -9 \sin^2 t \, dt = -9 \pi = -\text{area}\)

7 \(f(y j)\) is not a gradient field; take line \(x = t, y = t \) from \((0, 0)\) to \((1, 1)\) and \(\int t^2 \, dt \neq \frac{1}{2}\)

9 \(f(x, y) = x + \frac{1}{2}y^2; f(0, 1) - f(1, 0) = -\frac{1}{2}\)

11 \(f = r = \sqrt{x^2 + y^2} ; f(0, 1) - f(1, 0) = 0\)

15 \(f = r = \sqrt{x^2 + y^2} ; f(0, 1) - f(1, 0) = 0\)

17 \(\parallel\)

19 \(x = \cos t, z = \sin t, ds = a \, dt, M = \int_0^{2\pi} (a + a \sin t) \, a \, dt = 2\pi a^2\)
21 \( x = a \cos t, y = a \sin t, ds = dt, M = \int_0^{2\pi} a^3 \cos^2 t \, dt = \pi a^3, (x, y) = (0, 0) \) by symmetry

23 \( T = \frac{21 + 2i}{\sqrt{4 + 4t}} \), \( F = 3x \mathbf{i} + 4j \), \( ds = 2\sqrt{1 + t^2} \, dt \), \( F \cdot dR = (6t \mathbf{i} + 4j) \cdot (2 dt \mathbf{i} + 2 dt j) = 20 dt \); work = \( \int_1^2 20 \, dt = 30 \)

25 If \( \frac{\partial M}{\partial y} = \frac{\partial N}{\partial x} \) then \( M = ay + b, N = ax + c \), constants \( a, b, c \)

27 \( F = 4xj \) (work = 4 from (1,0) up to (1,1))

33 Not conservative; \( \int_0^{2\pi} (t - t^2 \mathbf{j}) \cdot dt + dt j = 0 \)

35 \( \frac{\partial M}{\partial y} = ax, \frac{\partial N}{\partial x} = bx \), \( f\mathbf{i} + g\mathbf{j} \)

39 \( \frac{\partial M}{\partial y} = -\frac{xy}{r^2}, f = r = \sqrt{x^2 + y^2} = |x + iy| \)

41 \( F = (x - y)\mathbf{i} + (x + y)\mathbf{j} \) has \( \frac{\partial M}{\partial y} = -1, \frac{\partial N}{\partial x} = 1 \), no \( f \)

\[ f = |x - 2y|^{(1,1)}_{(0,0)} = -1 \quad f = |xy|^{(1,1)}_{(0,0)} = 1 \]

\section{15.3 Green's Theorem (page 571)}

1 \( \int_0^{2\pi} (a \cos t) a \cos t \, dt = \pi a^2; N_x - M_y = 1 \)

3 \( \int_0^1 x \, dx + \int_0^1 x \, dx = 0, N_x - M_y = 0 \)

5 \( \int x^2y \, dx = \int_0^{2\pi} (a \cos t)^2 (a \sin t) (-a \sin t \, dt) = -\frac{a^4}{4} \int_0^{2\pi} (\sin 2t)^2 \, dt = -\frac{\pi a^4}{4} \)

7 \( \int x \, dy - y \, dx = \int_0^{2\pi} (\cos^2 t + \sin^2 t) \, dt = \pi \)

9 \( \frac{1}{2} \int_0^{2\pi} (3 \cos^4 t \sin^2 t + 3 \sin^4 t \cos^2 t) \, dt = \frac{1}{2} \int_0^{2\pi} 3 \cos^2 t \sin^2 t \, dt = \frac{3\pi}{2} \) (see Answer 5)

11 \( \int F \cdot dR = 0 \) around any loop; \( F = \frac{\pi}{3} \mathbf{i} + \frac{\pi}{3} \mathbf{j} \) and \( \int F \cdot dR = \int_0^{2\pi} [-\sin t \cos t + \sin t \cos t] \, dt = 0 \)

13 \( x = \cos 2t, y = \sin 2t, t \) from 0 to \( 2\pi \); \( \int_0^{2\pi} -2 \sin^2 2t \, dt = -2\pi = -2 \) (area);

15 \( \int M \, dy - N \, dx = \int_0^{2\pi} 2 \sin t \cos t \, dt = 0 \)

17 \( M = \frac{\pi}{2}, N = \frac{\pi}{2}, \int M \, dy - N \, dx = \int_0^{2\pi} (\cos^2 t + \sin^2 t) \, dt = \pi \)

19 \( \int M \, dy - N \, dx = \int -x^2 y \, dx = \int_0^1 -x^2 (1-x) \, dx = \frac{1}{12}; \int_0^1 1 - y^2 x \, dx = \frac{1}{12} \)

21 \( \int (M_x + N_y) \, dx \, dy = \int \int \text{div} F \, dx \, dy = 0 \) between the circles

23 Work: \( \int a \, dx + b \, dy = \int \left( \frac{\partial b}{\partial x} - \frac{\partial a}{\partial y} \right) \, dx \, dy \); Flux: same integral

25 \( g = \tan^{-1} \left( \frac{y}{x} \right) = \theta \) is undefined at \( (0,0) \)

27 Test \( M_y = x^2 \, dx + y^2 \, dy \) is exact = \( d\left( \frac{1}{3} x^3 + \frac{1}{3} y^3 \right) \)

29 \( \text{div} F = 2y - 2y = 0; g = xy^2 \)

31 \( \text{div} F = 2x + 2y; \) no \( g \)

33 \( \text{div} F = 0; g = e^x \sin y \)

35 \( \text{div} F = 0; g = \frac{y^2}{x} \)

37 \( N_x - M_y = -2x, -6xy, 0, 2x - 2y, 0, -2xe^y; \) in 31 and 33 \( f = \frac{1}{3} (x^3 + y^3) \) and \( f = e^x \cos y \)

39 \( F = (3x^2 - 3y^2)i - 6xy j \); \( \text{div} F = 0 \)

41 \( f = x^4 - 6x^2 y^2 + 3y^4; g = 3x^3 y - 4x y^3 \)

43 \( F = e^x \cos y i - e^x \sin y j; g = e^x \sin y \)

45 \( N = f(x), \int M \, dx + N \, dy = \int_0^1 f(1)\, dy + \int_0^1 f(0)\, dy = f(1) - f(0); \int (N_x - M_y) \, dx \, dy = \int \int \frac{\partial f}{\partial x} \, dx \, dy = \int_0^1 \int_0^1 \text{div} F \, dx \, dy \) (Fundamental Theorem of Calculus)
Section 15.4 Surface Integrals  (page 581)

1. $\mathbf{N} = -2\mathbf{i} - 2\mathbf{j} + \mathbf{k}; \, dS = \sqrt{1 + 4x^2 + 4y^2} \, dx \, dy; \, \int_0^{2\pi} \int_0^r \sqrt{1 + 4r^2} \, r \, dr \, d\theta = \frac{\pi}{8} (173/2 - 1)$

3. $\mathbf{N} = -\mathbf{i} + \mathbf{j} + \mathbf{k}; \, dS = \sqrt{3} \, dx \, dy; \, \text{area } \sqrt{3}\pi$

5. $\mathbf{N} = -\frac{\mathbf{y}}{\sqrt{1-x^2-y^2}} + \mathbf{k}; \, dS = \frac{dx \, dy}{\sqrt{1-x^2-y^2}}; \, \int_0^{2\pi} \int_0^r \frac{r \, dr \, d\theta}{\sqrt{1-r^2}} = \pi (2 - \sqrt{2})$

7. $\mathbf{N} = -7\mathbf{j} + \mathbf{k}; \, dS = 5\sqrt{2} \, dx \, dy; \, \text{area } 5\sqrt{2}A$

9. $\mathbf{N} = (y^2 - x^2)\mathbf{i} - 2xy\mathbf{j} + \mathbf{k}; \, dS = \sqrt{1 + (y^2 - x^2)^2 + 4x^2y^2} \, dx \, dy = \sqrt{1 + (y^2 + x^2)^2} \, dx \, dy; \, \int_0^{2\pi} \int_0^r \frac{r \, dr \, d\theta}{\sqrt{1+r^2}} = \pi (2 - \sqrt{2})$

11. $\mathbf{N} = 2\mathbf{i} + 2\mathbf{j} + \mathbf{k}; \, dS = 3dx \, dy; \, \text{3(area of triangle with } 2x + 2y \leq 1) = \frac{3}{8}$

13. $\pi a\sqrt{a^2 + h^2} \quad 15. \int_0^1 \int_0^{1-y} \frac{x}{y} \sqrt{\frac{1}{2}} \, dx \, dy = \frac{\sqrt{3}}{24}$

17. $\int_0^{2\pi} \int_0^{\pi/4} \sin^2 \phi \cos \phi \sin \theta \cos \theta (\sin \phi \, d\phi \, d\theta) = 0 \quad 19. \mathbf{A} = \mathbf{i} + 2\mathbf{j}; \, \mathbf{B} = \mathbf{j} + \mathbf{k}; \, \mathbf{N} = -\mathbf{i} - \mathbf{j} - \mathbf{k}; \, dS = \sqrt{3} \, du \, dv$

21. $\mathbf{A} = -\sin u (\cos v \mathbf{i} + \sin v \mathbf{j}) + \cos u \mathbf{k}; \, \mathbf{B} = -3 + \cos u \sin v \mathbf{i} + (3 + \cos u) \cos v \mathbf{j}; \, \mathbf{N} = -(3 + \cos u) \cos v \mathbf{u} + \cos u \sin v \mathbf{j} + \sin u \mathbf{k}; \, dS = (3 + \cos u) \, du \, dv$

23. $\int \int (-M \frac{\partial F}{\partial y} - N \frac{\partial F}{\partial y} + P) \, dx \, dy = \int \int (-2x^2 - 2y^2 + x) \, dx \, dy = \int \int -r^2 (r \, dr \, d\theta) = -8\pi$

25. $\mathbf{F} \cdot \mathbf{N} = -x + y + z = 0 \, \text{on plane}$

27. $\mathbf{N} = -\mathbf{i} - \mathbf{j} + \mathbf{k}, \, \mathbf{F} = (v + u) \mathbf{i} - u \mathbf{j}, \, \mathbf{J} = -\mathbf{i} - \mathbf{j} - \mathbf{k}; \, \mathbf{dS} = \frac{\mathbf{F} \cdot \mathbf{N} \, dS}{\mathbf{F} \cdot \mathbf{N}} = \int \int -u \, du \, dv = 0$

29. $\int \int dS = \int_0^{2\pi} \int_0^r (3 + \cos u) \, du \, dv = 12\pi^2 \quad 31. \text{Yes} \quad 33. \text{No}$

35. $\mathbf{A} = 1 + f' \cos \theta \mathbf{j} + f' \sin \theta \mathbf{k}; \, \mathbf{B} = -f \sin \theta \mathbf{j} + f \cos \theta \mathbf{k}; \, \mathbf{N} = f f' \mathbf{i} - f \cos \theta \mathbf{j} - f \sin \theta \mathbf{k}; \, dS = |\mathbf{N}| \, dx \, dy = f(x) \sqrt{1 + f'^2} \, dx \, d\theta$

Section 15.5 The Divergence Theorem  (page 588)

1. $\text{div } \mathbf{F} = 1, \int \int \int dV = \frac{4\pi}{3} \quad 3. \text{div } \mathbf{F} = 2x + 2y + 2z, \int \int \int \text{div } \mathbf{F} \, dV = 0 \quad 5. \text{div } \mathbf{F} = 3, \int \int \int 3dV = \frac{3}{6} \pi \cdot \frac{3}{2} \quad 7. \mathbf{F} \cdot \mathbf{N} = \rho^2, \int_0^a \rho^2 \, dS = 4\pi a^4 \quad 9. \text{div } \mathbf{F} = 2x, \int_0^{2\pi} \int_0^a \int_0^a 2\rho \cos \phi (\rho^2 \sin \phi \, d\phi \, d\theta) = \frac{1}{4} \pi a^4$

11. $\int_0^a \int_0^a (2x + 1) \, dx \, dy = a^4 + a^3; -a^2 + 2a^2 + 0 + a^4 + 0 + a^3$

13. $\text{div } \left( \frac{\mathbf{R}}{\rho^2} \right) = \frac{\text{div } \mathbf{R}}{\rho^2} + \mathbf{R} \cdot \text{grad } \frac{1}{\rho} = \frac{7}{3} - \frac{7}{3}\mathbf{R} \cdot \text{grad } \rho$

19. Two spheres, $\mathbf{n}$ radial out, $\mathbf{n}$ radial in, $\mathbf{n} = \mathbf{k}$ on top, $\mathbf{n} = -\mathbf{k}$ on bottom, $\mathbf{n} = -\frac{z \mathbf{i} + y \mathbf{j}}{\sqrt{x^2 + y^2}}$ on side;

$\mathbf{n} = -\mathbf{i} - \mathbf{j} - \mathbf{k}, 1 + 2 \mathbf{j} + 3 \mathbf{k}$ on 4 faces; $\mathbf{n} = \mathbf{k}$ on top, $\mathbf{n} = \frac{1}{\sqrt{2}} (\mathbf{z} \mathbf{i} + 2 \mathbf{j} - \mathbf{k})$ on cone

21. $V$ = cylinder, $\int \int \int \text{div } \mathbf{F} \, dV = \int \int \int \left( \frac{\partial F}{\partial z} + \frac{\partial N}{\partial y} \right) \, dx \, dy$ (integral = 1); $\int \int \mathbf{F} \cdot \mathbf{n} \, dS = \int M \, dy - N \, dx, \text{z integral} = 1 \, \text{on side}, \mathbf{F} \cdot \mathbf{n} = 0 \, \text{top and bottom}; \text{Green's flux theorem}$

23. $\text{div } \mathbf{F} = -\frac{2GM}{a^3} = -4\pi G; \text{at the center}; \mathbf{F} = 2\mathbf{R}$ inside, $\mathbf{F} = 2(\mathbf{R} \mathbf{a})^3 \mathbf{R}$ outside

25. $\text{div } \mathbf{u} = \frac{2a}{\rho}, \int \int \mathbf{E} \cdot \mathbf{n} \, dS = \int \int 1 \, dS = 4\pi \quad 27. \mathbf{F} (\text{div } \mathbf{F} = 0); \mathbf{F}, \mathbf{T}(\mathbf{F} \cdot \mathbf{n} \leq 1); \mathbf{F}$

29. Plane circle; top half of sphere; $\text{div } \mathbf{F} = 0$

Section 15.6 Stokes' Theorem and the Curl of $\mathbf{F}$  (page 595)

1. $\text{curl } \mathbf{F} = \mathbf{i} + \mathbf{j} + \mathbf{k} \quad 3. \text{curl } \mathbf{F} = 0 \quad 5. \text{curl } \mathbf{F} = 0 \quad 7. f = \frac{1}{2}(x + y + z)^2$

9. $\text{curl } x^m \mathbf{i} = 0; \text{cylindrical} \quad 11. \text{curl } \mathbf{F} = 2\mathbf{y}; \mathbf{n} = \mathbf{j} \text{ on circle so } \int \int \mathbf{F} \cdot \mathbf{n} \, dS = 0$

13. $\text{curl } \mathbf{F} = 2\mathbf{i} + 2\mathbf{j}, \mathbf{n} = \mathbf{i}, \int \int \text{curl } \mathbf{F} \cdot \mathbf{n} \, dS = \int \int 2 \, dS = 2\pi$
A-30 Answers to Odd-Numbered Problems

15 Both integrals equal \( \int F \cdot dR \); Divergence Theorem, \( V \) = region between \( S \) and \( T \), always \( \text{div} \ \text{curl} \ F = 0 \\
17 \ \text{Always} \ \text{div} \ \text{curl} \ F = 0 \quad 19 \ f = xz + y \quad 21 \ f = e^{x-y} \quad 23 \ F = yk

25 \ \text{curl} \ F = (a_3 b_2 - a_2 b_3) i + (a_1 b_3 - a_3 b_1) j + (a_2 b_1 - a_1 b_2) k \quad 27 \ \text{curl} \ F = 2\omega k; \ \text{curl} \ F \cdot \frac{i+j+k}{\sqrt{3}} = 2\omega/\sqrt{3}

29 \ F = x(a_3 + a_2 y)i + y(a_1 x + a_3 z) j + z(a_1 x + a_2 y) k

31 \ \text{curl} \ F = -2k, \int \int (-2k \cdot dS) = \int_0^{2\pi} \int_0^{\pi/2} -2 \cos \phi \sin \phi d\phi d\theta = -2\pi; \int y \ dx - x \ dy = \int_0^{2\pi} (-\sin^2 t - \cos^2 t) dt = -2\pi

33 \ \text{curl} \ F = 2a, 2 \int (a_1 x + a_2 y + a_3 z) dS = 0 + 0 + 2a_3 \int_0^{2\pi} \int_0^{\pi/2} \cos \phi \sin \phi d\phi d\theta = 2\pi a_3

35 \ \text{curl} \ F = -i, n = \frac{i+j+k}{\sqrt{3}}, \int \int F \cdot ndS = -\frac{1}{\sqrt{3}} \pi r^2

37 \ g = \frac{y^2}{2} - \frac{z^2}{3} \ \text{is stream function; zero divergence}

39 \ \text{div} \ F = \text{div} \ (V + W) = \text{div} \ V \ \text{so} \ y = \text{div} \ V \ \text{so} \ V = \frac{y^2}{2} j \ \text{has zero curl}. \ \text{Then} \ W = F - V = x y i - \frac{y^2}{2} j

41 \ \text{curl} \ (\text{curl} \ F) = \text{curl} \ (-2yk) = -2i; \ \text{grad} \ (\text{div} \ F) = \text{grad} \ 2x = 2i; \ F_{xx} + F_{yy} + F_{zz} = 4i

43 \ \text{curl} \ E = \frac{\partial B}{\partial t} = a \sin t \ \text{so} \ E = \frac{1}{2} (a \times R) \sin t

45 \ n = j \ \text{so} \ \int M \ dx + P \ dz = \int (\frac{2M}{dx} - \frac{P}{dz}) \ dx \ dz \quad 47 \ M_y^* = M_y + M_z f_y + P_y f_z + P_z f_y f_z + P_{fyz}

49 \ \int \int F \cdot dR = \int \int \text{curl} \ F \cdot n \ dS; \ \int \int F \cdot ndS = \int \int \text{div} \ F \ dV

CHAPTER 16  MATHEMATICS AFTER CALCULUS

Section 16.1 Linear Algebra (page 602)

1 All vectors \( c \begin{bmatrix} 2 \\ -1 \end{bmatrix} \)  
3 Only \( x = 0 \)  
5 Plane of vectors with \( x_1 + x_2 + x_3 = 0 \)

7 \( x_p = \begin{bmatrix} 3 \\ 0 \end{bmatrix}, A(x_p + x_0) = \begin{bmatrix} 3 \\ 6 \end{bmatrix} + \begin{bmatrix} 0 \\ 0 \end{bmatrix} \)  
9 \( A(x_p + x_0) = b + 0 = b; \) another solution

11 \( \begin{bmatrix} 1 \\ 1 \\ 1 \end{bmatrix}, \begin{bmatrix} 0 \\ 0 \\ 2 \end{bmatrix}, \begin{bmatrix} 2 \\ 2 \\ 2 \end{bmatrix}; b = \begin{bmatrix} c \\ c \\ c \end{bmatrix} \)

13 \( CC^T = \begin{bmatrix} 1 & 0 & 1 \\ 0 & 1 & 2 \\ 1 & 2 & 5 \end{bmatrix}; C^T C = \begin{bmatrix} 2 & 2 \\ 2 & 5 \end{bmatrix}; \) (2 by 3) (2 by 3) is impossible

15 Any two are independent  
17 \( C \) and \( F \) have independent columns

19 \( \det F = 3 \)  
21 \( F^{-1} = \frac{1}{3} \begin{bmatrix} 2 & -1 \\ -1 & 2 \end{bmatrix} \)

23 \( \det(F - \lambda I) = \det \begin{bmatrix} 2 - \lambda & 1 \\ 1 & 2 - \lambda \end{bmatrix} = (2 - \lambda)^2 - 1 = 3 - 4\lambda + \lambda^2 = 0 \) if \( \lambda = 1 \) or \( \lambda = 3; \)

\( F \begin{bmatrix} 1 \\ -1 \end{bmatrix} = \begin{bmatrix} 1 \\ -1 \end{bmatrix}, F \begin{bmatrix} 1 \\ 1 \end{bmatrix} = \begin{bmatrix} 1 \\ 1 \end{bmatrix} \)

25 \( y = e^t \begin{bmatrix} 1 \\ -1 \end{bmatrix}, y = e^{3t} \begin{bmatrix} 1 \\ 1 \end{bmatrix}, y = e^{t/2} \begin{bmatrix} 1 \\ -1 \end{bmatrix} + e^{3t/2} \begin{bmatrix} 1 \\ 1 \end{bmatrix} \)

27 \( \det \begin{bmatrix} 1 - \lambda & 1 & 1 \\ 1 & 1 - \lambda & 1 \\ 1 & 1 & 1 - \lambda \end{bmatrix} = (1 - \lambda)^3 - 3(1 - \lambda) + 2 = \lambda^3 - 3\lambda^2 = 0 \) if \( \lambda = 3 \) or \( \lambda = 0 \) (repeated)
Answers to Odd-Numbered Problems

**Section 16.2 Differential Equations (page 610)**

1. \(3Be^{3t} - Be^{3t} = 8e^{3t}\) gives \(B = 4\); \(y = 4e^{3t}\)
2. \(y = 3 - 2t + t^2\)
3. \(5Ae^t + 4e^{3t} = 7\) at \(t = 0\) if \(A = 3\)
4. Add \(y = Ae^{-t}\) because \(y' + y = 0\); choose \(A = -1\) so \(-e^{-t} + 3 - 2t + t^2 = 2\) at \(t = 0\)
5. \(e^{\lambda t}\) is not a vector; \(T\)
6. \(\lambda^2 - 5\lambda = 0\) if \(\lambda = 0\) or \(\lambda = 5\); \(A\left[\begin{array}{c} 2 \\ -1 \end{array}\right] = 0\), \(A\left[\begin{array}{c} 1 \\ 2 \end{array}\right] = 5\left[\begin{array}{c} 2 \\ 12 \end{array}\right] \)

**Section 16.3 Discrete Mathematics (page 615)**

1. Two then two then last one; go around hexagon
2. Six (each deletes one edge)
3. Connected: there is a path between any two nodes; connecting each new node requires an edge
4. Edge lengths 1, 2, 4
5. No; 1,3,4 on left connect only to 2,3 on right; 1,3 on right connect only to 2 on left
6. Yes
7. \(21 F\) (may loop); \(T\)
8. \(25 16\)

---

The text above is a continuation of the problems and solutions found in the section on odd-numbered problems of a textbook. It includes solutions to differential equations, discrete mathematics problems, and other mathematical exercises. Each problem is numbered and some are marked with yes or no based on the context provided. The text is structured to follow the sequence of the problems, ensuring that each is addressed in order.
Table of Integrals

1 \int u^n \, dx = \frac{u^{n+1}}{a(n+1)} \quad \text{except for } \int \frac{du}{u} = \ln|u| \quad \text{All the integrals 1 - 17 involve } u = ax + b

2 \int xu^n \, dx = \frac{u^{n+2}}{a^2(n+2)} - \frac{bu^{n+1}}{a^2(n+1)} \quad \text{except for } \int \frac{x \, du}{u} = \frac{x}{a} - \frac{b \ln|u|}{a^2} \quad \text{and } \int \frac{u \, dx}{u^2} = \frac{b}{a^2u} + \ln|u| \quad \text{for } a^2(b^2 - 4ac) > 0

3 \int \frac{e^2 \, dx}{u} = \frac{1}{a^2}(u^2 - 2bu + b^2) \ln|u|

4 \int \frac{e^2 \, dx}{u^2} = \frac{1}{ab}(u - 2b \ln|u| - \frac{b^2}{u})

5 \int \frac{e^2 \, dx}{u^2} = \frac{1}{a^2}(\ln|u| + \frac{2b}{u} - \frac{b^2}{2u^2})

6 \int \frac{e^2 \, dx}{x} = \frac{1}{b} \ln|\frac{x}{u}|

7 \int \frac{e^2 \, dx}{x^2} = -\frac{1}{b} + \frac{a}{b^2} \ln|\frac{x}{u}|

8 \int \frac{e^2 \, dx}{x^2} = \frac{1}{bu} - \frac{1}{b^2} \ln|\frac{x}{u}|

9 \int \frac{e^2 \, dx}{x^2} = \frac{b^2 - 2ax}{b^2u^2} + \frac{2a}{b^2} \ln|\frac{x}{u}|

10 \int \sqrt{u} \, dx = \frac{2}{3a}u^{3/2}

11 \int x \sqrt{u} \, dx = \frac{2}{15a^2} \sqrt{u^3/2}

12 \int x^2 \sqrt{u} \, dx = \frac{2}{105a^3} \sqrt{u^4/2}

13 \int \frac{\sqrt{u}}{x} \, dx = 2 \sqrt{u} + b \int \frac{dx}{\sqrt{u}}

14 \int \frac{\sqrt{u}}{x} \, dx = \frac{2}{3a^2} \sqrt{u^3/2}

15 \int \frac{\sqrt{u}}{x} \, dx = \frac{2}{15a^3} \sqrt{u^3/2}

16 \int \frac{\sqrt{u}}{x} \, dx = \frac{1}{b} \ln|\sqrt{u} + \sqrt{v}| \quad (b > 0) \quad \text{or } \frac{2}{b} \tan^{-1} \left( \frac{\sqrt{u}}{\sqrt{v}} \right) \quad (b < 0)

17 \int \frac{\sqrt{u}}{x} \, dx = -\frac{\sqrt{u}}{x} + \frac{a}{b} \int \frac{dx}{x \sqrt{u}}

18 \int \frac{dx}{(ax + b)(cx + d)} = \frac{1}{bc-ad} \ln|\frac{cx+d}{ax+b}|

19 \int \frac{dx}{(ax + b)(cx + d)} = \frac{1}{bc-ad} \ln|ax + b| - \frac{d}{c} \ln|cx + d|

20 \int \sqrt{x^2 + a^2} \, dx = \frac{x}{2} \sqrt{x^2 + a^2} + \frac{a^2}{2} \ln|x + \sqrt{x^2 + a^2}|

21 \int \frac{dx}{\sqrt{x^2 + a^2}} = \ln|x + \sqrt{x^2 + a^2}|

22 \int \frac{dx}{\sqrt{x^2 + a^2}} = \sqrt{x^2 + a^2} - a \ln \left( \frac{x + \sqrt{x^2 + a^2}}{a} \right)

23 \int \frac{dx}{\sqrt{x^2 - a^2}} = \sqrt{x^2 - a^2} - a \sec^{-1} \frac{x}{a}

24 \int \frac{dx}{\sqrt{x^2 + a^2}} = \frac{x}{2} \sqrt{x^2 + a^2} + \frac{a^2}{2} \ln|x + \sqrt{x^2 + a^2}|

25 \int \frac{dx}{x \sqrt{x^2 + a^2}} = -\frac{x}{2} \sqrt{x^2 + a^2} + \ln|x + \sqrt{x^2 + a^2}|

26 \int \frac{dx}{x \sqrt{x^2 + a^2}} = \frac{x}{2} \sqrt{x^2 + a^2} - \frac{a^2}{2} \ln|x + \sqrt{x^2 + a^2}|

27 \int \frac{dx}{x \sqrt{x^2 + a^2}} = \tan^{-1} \left( \frac{x}{a} \right)

28 \int \frac{dx}{x \sqrt{x^2 + a^2}} = \frac{x}{2} \sqrt{x^2 + a^2} + \frac{a^2}{2} \ln|x + \sqrt{x^2 + a^2}|

29 \int \frac{dx}{(x^2 + a^2)^{3/2}} = \frac{x}{2a} \sqrt{x^2 + a^2}

30 \int \frac{dx}{(x^2 + a^2)^{3/2}} = \frac{x}{2a} \sqrt{x^2 + a^2} - a \ln|\frac{a^2 + x^2}{x}|

31 \int \frac{dx}{(x^2 + a^2)^{3/2}} = \sin^{-1} \frac{x}{a}

32 \int \frac{dx}{(x^2 + a^2)^{3/2}} = \frac{x}{a \sqrt{x^2 + a^2}}

33 \int \frac{dx}{x \sqrt{x^2 + a^2}} = \sqrt{x^2 + a^2} - a \ln|\frac{a^2 + x^2}{x}|

34 \int \frac{dx}{x \sqrt{x^2 + a^2}} = \frac{x}{2} \sqrt{x^2 + a^2} - \frac{a^2}{2} \sin^{-1} \frac{x}{a}

35 \int \frac{dx}{x \sqrt{x^2 - a^2}} = -\frac{a^2}{x} \sqrt{x^2 - a^2}

36 \int \frac{dx}{x \sqrt{x^2 - a^2}} = -\frac{a^2}{x} \sqrt{x^2 - a^2} - \sin^{-1} \frac{x}{a}

37 \int \frac{dx}{x \sqrt{x^2 - a^2}} = -\frac{1}{a} \ln|x + \sqrt{x^2 - a^2}|

38 \int \frac{dx}{x \sqrt{x^2 - a^2}} = \frac{x}{2} \sqrt{x^2 - a^2} + \frac{a^2}{2} \sin^{-1} \frac{x}{a}

39 \int (a^2 - x^2)^{3/2} \, dx = \frac{x}{2} \sqrt{a^2 - x^2} - \frac{a^4}{8} \sin^{-1} \frac{x}{a}

40 \int \frac{dx}{b \sin ax} = \frac{1}{b \sin^2 ax} \tan^{-1} \left( \frac{b \sin ax}{a \sin ax} \right), \quad b > 0

41 \int \frac{dx}{1 + \sin ax} = -\frac{1}{a} \tan \left( \frac{x}{4} + \frac{a x}{2} \right)

42 \int \frac{dx}{b \cos ax} = \frac{1}{a \cos^2 ax} \tan^{-1} \left( \frac{b \cos ax}{a \cos ax} \right), \quad b < 0

43 \int \frac{dx}{1 - \cos ax} = \frac{1}{a} \tan \left( \frac{x}{2} + \frac{a x}{2} \right)

44 \int \frac{dx}{b \cos ax} = \frac{1}{a \cos^2 ax} \tan^{-1} \left( \frac{b \cos ax}{a \cos ax} \right), \quad b > 0

45 \int \frac{dx}{1 - \cos ax} = \frac{1}{a} \tan \left( \frac{x}{2} + \frac{a x}{2} \right)

46 \int \frac{dx}{b \cos ax} = -\frac{1}{a \cos^2 ax} \ln|\frac{b \cos ax + a^2 \sin ax}{b \cos ax}|

47 \int \frac{dx}{b \sin ax} = \frac{1}{a \cos^2 ax} \ln|\frac{b \sin ax + a^2 \cos ax}{b \cos ax}|

48 \int \sin^{-1} ax \, dx = x \sin^{-1} ax + \frac{x}{a} \sqrt{1 - a^2 x^2}

49 \int \frac{x^2 \sin^{-1} ax \, dx}{ax} = \frac{z^{n+1}}{n+1} \sin^{-1} ax - \frac{a}{n+1} \int \frac{z^{n+1} \, dz}{\sqrt{1 - a^2 z^2}}

50 \int \tan^{-1} ax \, dx = \frac{1}{a} \ln(1 + a^2 x^2)

51 \int x^n \tan^{-1} ax \, dx = \frac{x^{n+1}}{n+1} \tan^{-1} ax - \frac{a}{n+1} \int \frac{x^{n+1} \, dz}{1 + a^2 z^2}

52 \int e^{ax} \, dx = \frac{e^{ax}}{a}

53 \int xe^{ax} \, dx = \frac{e^{ax}}{a^2} (ax - 1)

54 \int x^2 e^{ax} \, dx = \frac{e^{ax}}{a^2} (a^2 x^2 - 2ax + 2)

55 \int \frac{dx}{z \ln ax} = \ln|\ln ax| \quad \text{Not elementary: } \int e^{z^2} \, dz, \int e^{z} \ln \, dx, \int \frac{dx}{z} \, dz, \int \frac{\sin z}{z} \, dx, \int \frac{\sin^{-1} z}{z} \, dx
Exponentials and Logarithms

\[ y = b^x \iff x = \log_b y \quad \text{and} \quad e^x \iff x = \ln y \]
\[ e = \lim_{n \to \infty} \left(1 + \frac{1}{n}\right)^n = \sum_{n=0}^{\infty} \frac{1}{n!} = 2.71828 \ldots \]
\[ e^x = \lim_{n \to \infty} \left(1 + \frac{x}{n}\right)^n = \sum_{n=0}^{\infty} \frac{x^n}{n!} \]
\[ \ln y = \int_1^y \frac{dx}{x} \quad \ln 1 = 0 \quad \ln e = 1 \]
\[ \ln xy = \ln x + \ln y \quad \ln x^n = n \ln x \]
\[ \log_a y = \left(\log_a b\right) \left(\log_b y\right) \quad \log_a b = 1/\log_b a \]
\[ e^{x+y} = e^x e^y \quad e^x e^y = e^{x+y} \quad e^{\ln y} = y \]

Vectors and Determinants

\[ \mathbf{A} = a_1 \mathbf{i} + a_2 \mathbf{j} + a_3 \mathbf{k} \]
\[ |\mathbf{A}|^2 = \mathbf{A} \cdot \mathbf{A} = a_1^2 + a_2^2 + a_3^2 \quad \text{(length squared)} \]
\[ \mathbf{A} \cdot \mathbf{B} = a_1 b_1 + a_2 b_2 + a_3 b_3 = |\mathbf{A}| |\mathbf{B}| \cos \theta \]
\[ |\mathbf{A} \cdot \mathbf{B}| \leq |\mathbf{A}| |\mathbf{B}| \quad \text{(Schwarz inequality: } \cos \theta \leq 1) \]
\[ |\mathbf{A} + \mathbf{B}| \leq |\mathbf{A}| + |\mathbf{B}| \quad \text{(triangle inequality)} \]
\[ |\mathbf{A} \times \mathbf{B}| = |\mathbf{A}| |\mathbf{B}| \sin \theta \quad \text{(cross product)} \]
\[ \mathbf{A} \times \mathbf{B} = \begin{vmatrix} i & j & k \\ a_1 & a_2 & a_3 \\ b_1 & b_2 & b_3 \end{vmatrix} = i(a_2 b_3 - a_3 b_2) - j(a_1 b_3 - a_3 b_1) + k(a_1 b_2 - a_2 b_1) \]

Right hand rule \( \mathbf{i} \times \mathbf{j} = \mathbf{k}, \mathbf{j} \times \mathbf{k} = \mathbf{i}, \mathbf{k} \times \mathbf{i} = \mathbf{j} \)

Parallelogram area \( |a_1 b_2 - a_2 b_1| = |\text{Det}| \)

Triangle area \( \frac{1}{2} |a_1 b_2 - a_2 b_1| = \frac{1}{2} |\text{Det}| \)

Box volume \( |\mathbf{A} \cdot (\mathbf{B} \times \mathbf{C})| = |\text{Determinant}| \)

<table>
<thead>
<tr>
<th>SI Units</th>
<th>Symbols</th>
<th>From</th>
<th>To</th>
<th>Multiply by</th>
</tr>
</thead>
<tbody>
<tr>
<td>length</td>
<td>meter</td>
<td>m</td>
<td>degrees</td>
<td>radians</td>
</tr>
<tr>
<td>mass</td>
<td>kilogram</td>
<td>kg</td>
<td>calories</td>
<td>joules</td>
</tr>
<tr>
<td>time</td>
<td>second</td>
<td>s</td>
<td>BTU</td>
<td>joules</td>
</tr>
<tr>
<td>current</td>
<td>ampere</td>
<td>A</td>
<td>foot-pounds</td>
<td>joules</td>
</tr>
<tr>
<td>frequency</td>
<td>hertz</td>
<td>Hz \sim 1/s</td>
<td>feet</td>
<td>meters</td>
</tr>
<tr>
<td>force</td>
<td>newton</td>
<td>N \sim \text{kg-m/s}^2</td>
<td>miles</td>
<td>km</td>
</tr>
<tr>
<td>pressure</td>
<td>pascal</td>
<td>Pa \sim \text{N/m}^2</td>
<td>feet/sec</td>
<td>km/hr</td>
</tr>
<tr>
<td>energy, work</td>
<td>joule</td>
<td>J \sim \text{N-m}</td>
<td>pounds</td>
<td>kg</td>
</tr>
<tr>
<td>power</td>
<td>watt</td>
<td>W \sim \text{J/s}</td>
<td>ounces</td>
<td>kg</td>
</tr>
<tr>
<td>charge</td>
<td>coulomb</td>
<td>C \sim \text{A-s}</td>
<td>gallons</td>
<td>liters</td>
</tr>
<tr>
<td>temperature</td>
<td>kelvin</td>
<td>K</td>
<td>horsepower</td>
<td>watts</td>
</tr>
</tbody>
</table>

Speed of light \( c = 2.9979 \times 10^8 \text{m/s} \)

Gravity \( G = 6.6720 \times 10^{-11} \text{Nm}^2/\text{kg}^2 \)

Equations and Their Solutions

\[ y' = cy \quad y_0 e^{ct} \]
\[ y' = cy + s \quad y_0 e^{ct} + \frac{s}{c} (e^{ct} - 1) \]
\[ y' = cy - b y^2 \quad \frac{c}{b + de^{-ct}} \quad d = \frac{c-by_0}{y_0} \]
\[ y'' = -\lambda^2 y \quad \cos \lambda t \text{ and } \sin \lambda t \]
\[ my'' + dy' + ky = 0 \quad e^{\lambda_1 t} \text{ and } e^{\lambda_2 t} \text{ or } te^{\lambda_1 t} \]
\[ y_{n+1} = ay_n \quad a^n y_0 \]
\[ y_{n+1} = ay_n + s \quad a^n y_0 + s \frac{a^n - 1}{a - 1} \]

Matrices and Inverses

\[ Ax = \text{combination of columns } = b \]

Solution \( x = A^{-1} b \) if \( A^{-1} A = I \)

Least squares \( A^T A \bar{x} = A^T b \)

\[ Ax = \lambda x \quad (\lambda \text{ is an eigenvalue}) \]

\[ \begin{bmatrix} a & b \\ c & d \end{bmatrix}^{-1} = \frac{1}{ad-bc} \begin{bmatrix} d & -b \\ -c & a \end{bmatrix} \]

\[ (AB)^{-1} = B^{-1} A^{-1}, (AB)^T = B^T A^T \]

\[ \begin{bmatrix} a & b & c \\ a & b & c \end{bmatrix}^{-1} = \frac{1}{d} \begin{bmatrix} b \times c \\ a \times b \end{bmatrix} \]

\[ \begin{vmatrix} a_1 & a_2 & a_3 \\ b_1 & b_2 & b_3 \end{vmatrix} = \begin{vmatrix} +a_1 b_2 c_3 + a_2 b_3 c_1 + a_3 b_1 c_2 \\ -a_1 b_3 c_2 - a_2 b_1 c_3 - a_3 b_2 c_1 \end{vmatrix} \]

Radius at Equator \( R = 6378 \text{ km} = 3964 \text{ miles} \)

Acceleration \( g = 9.8067 \text{ m/s}^2 = 32.174 \text{ ft/s}^2 \)
Sums and Infinite Series

\[ 1 + x + \ldots + x^{n-1} = \frac{1 - x^n}{1 - x} \]
\[ 1 + n z + \frac{n(n-1)}{2} z^2 + \ldots + z^n = (1 + z)^n \]
\[ 1 + 2 + \ldots + n = \frac{1}{2} n(n+1) \approx \frac{n^2}{2} \]
\[ 1^2 + 2^2 + \ldots + n^2 = \frac{n(n+1)(2n+1)}{6} \approx \frac{n^3}{3} \]
\[ 1 + \frac{1}{2} + \ldots + \frac{1}{n} \approx \ln n \rightarrow \infty \text{ (harmonic)} \]
\[ 1 - \frac{1}{2} + \frac{1}{3} - \ldots = \ln 2 \text{ (alternating)} \]
\[ 1 - \frac{1}{3} + \frac{1}{5} - \ldots = \frac{\pi}{4} \]
\[ \sum \frac{1}{n^2} = \frac{\pi^2}{6} \quad \sum \frac{1}{n^4} = \frac{\pi^4}{90} \]

\[ \frac{1}{1-x} = 1 + x + x^2 + \ldots \text{ (geometric: } |x| < 1) \]
\[ \frac{1}{1+x} = 1 - x + x^2 - \ldots \text{ (geometric for } -1) \]
\[ \ln(1+z) = z - \frac{z^2}{2} + \frac{z^3}{3} - \ldots = \int \frac{dz}{1+z} \]
\[ \sin x = x - x^3/6 + x^5/120 - \ldots \text{ (all } x) \]
\[ \cos x = 1 - x^2/2 + x^4/24 - \ldots \text{ (all } x) \]
\[ e^z = 1 + z + \frac{z^2}{2!} + \ldots \text{ (all } x) \]
\[ e^{iz} = \cos z + i \sin z \text{ (Euler's formula)} \]
\[ \cosh z = \frac{1}{2}(e^z + e^{-z}) = 1 + \frac{z^2}{2!} + \ldots \]
\[ \sinh z = \frac{1}{2}(e^z - e^{-z}) = z + \frac{z^3}{3!} + \ldots \]
\[ \cos \theta + i \sin \theta)^n = \cos n\theta + i \sin n\theta \]
\[ f(x) = f(0) + f'(0)x + f''(0)\frac{x^2}{2!} + \ldots \text{ (Taylor)} \]
\[ f(x,y) = f + x f_x + y f_y + \frac{x^2}{2!} f_{xx} + xy f_{xy} + \ldots \]

Partial Derivatives of \( z = f(x,y) \)

Tangent plane \( z - z_0 = \left( \frac{\partial f}{\partial x} \right)(x - x_0) + \left( \frac{\partial f}{\partial y} \right)(y - y_0) \)

Cone or pyramid \[ \text{Volume } \frac{1}{3} \text{ (base area) (height)} \]

Length of curve \[ \int ds = \int \sqrt{1 + (dy/dx)^2} \ dx \]

Area between curves \[ \int (u(x) - v(x)) \ dx \]

Surface area of revolution \[ \int 2\pi r \ ds = \int r \sqrt{1 + (dy/dx)^2} \ dx \]

Volume of revolution: Slices \[ \int \pi y^2 \ dx \]

Area of surface \[ \int \int \sqrt{1 + x^2 + z^2} \ dx dy \]

Mass \[ M = \int \int \rho \ dA \quad \text{Moment } M_y = \int \int x \rho \ dA \]

\[ \vec{a} = M_y/M, \vec{b} = M_z/M \quad \text{Moment of Inertia } I_y = \int \int x^2 \ dA \]

Work \[ W = \int_a^b F(x) \ dx = V(b) - V(a) \quad \text{Force } F = dV/dx \]

Vector field \( \mathbf{F}(x,y,z) = M \mathbf{i} + N \mathbf{j} + P \mathbf{k} \)

Divergence of \( \mathbf{F} = \nabla \cdot \mathbf{F} = \frac{\partial M}{\partial x} + \frac{\partial N}{\partial y} + \frac{\partial P}{\partial z} \)

Chain rule: \[ \frac{dz}{dt} = \frac{\partial z}{\partial x} \frac{dx}{dt} + \frac{\partial z}{\partial y} \frac{dy}{dt} + \frac{\partial z}{\partial t} \frac{dt}{dt} \]

Stokes' Theorem \[ \int \int \mathbf{F} \cdot \mathbf{n} \ ds = \int \int \text{div} \mathbf{F} \ dV \]

An additional table of integrals is included just after the index.