Discrete Mathematics: An Open Introduction
Course Notes for Math 228 at the University of Northern Colorado

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Chapter 1

Logic and Set Theory

Logic is the study of consequence. Given a few mathematical statements or facts, we would like to be able to draw some conclusions. For example, if I told you that a particular real valued function was continuous on the interval \([0, 1]\), and \(f(0) = -1\) and \(f(1) = 5\), can we conclude that there is some point between \([0, 1]\) where the graph of the function crosses the \(x\)-axis? Yes, we can, thanks to the Intermediate Value Theorem from Calculus. Can we conclude that there is exactly one point? No. Whenever we find an “answer” in math, we really have a (perhaps hidden) argument - given the situation we are in, we can conclude the answer is the case. Of course real mathematics is about proving general statements (like the Intermediate Value Theorem), and this too is done via an argument, usually called a proof. We start with some given conditions - these are the premises of our argument. From these we find a consequence of interest - our conclusion.

The problem is, as I’m sure you are aware from arguing with friends, not all arguments are good arguments. A “bad” argument is one in which the conclusion does not follow from the premises - the conclusion is not a consequence of the premises. Logic is the study of what makes an argument good or bad. In other words, logic aims to determine in which cases a conclusion is, or is not, a consequence of a set of premises.

By the way, “argument” is actually a technical term in math (and philosophy, another discipline which studies logic):

**Definition 1.** An *argument* is a set of statements, on of which is called the *conclusion* and the rest of which are called *premises*. An argument is said to be *valid* if the conclusion must be true whenever the premises are all true. An argument is *invalid* if it is not valid - if it is possible for all the premises to be true and the conclusion to be false.

For example, consider the following two arguments:

If Edith eats her vegetables, then she can have a cookie.
Edith eats her vegetables.
\[
\therefore \text{ Edith gets a cookie.}
\]

Florence must eat her vegetables in order to get a cookie.
Florence eats her vegetables.
\[
\therefore \text{ Florence gets a cookie.}
\]
(The symbol “\(\therefore\)” means “therefore”)

Are these arguments valid? Hopefully you agree that the first one is but the second one is not. Logic tells us why. How? By analyzing the structure of the statements in the argument. Notice that the two arguments above look almost identical. Edith and Florence both eat their vegetables. In both cases there is a connection between eating of vegetables and cookies. But
we claim that it is valid to conclude that Edith gets a cookie, but not that Florence does. The difference must be in the connection between eating vegetables and getting cookies. We need to be good at reading and comprehending these English sentences. Do the two sentences mean the same thing? Unfortunately, when talking in everyday language, we are often sloppy, and you might be tempted to say they are equivalent. But notice that just because Florence must eat her vegetables, we have not said that doing so would be enough - she might also need to clean her room, for example. In everyday (non-mathematical) practice, you might say that the “other direction” was implied. We don’t ever get to say that.

Our goal in studying logic is to gain intuition on which arguments are valid and which are invalid. This will require us to become better at reading and writing mathematics - a worthy goal in its own right. So let’s get started.

1.1 Propositional Logic

A proposition is simply a statement. Propositional logic studies the ways statements can interact with each other. It is important to remember that propositional logic does not really care about the content of the statements. For example, from a propositional logic statement point, the claims, “if the moon is made of cheese then basketballs are round,” and, “if spiders have eight legs then Sam walks with a limp” are exactly the same. They are both statements of the form, “if <something>, then <something else>.”

Here’s a question: is it true that when playing Monopoly, if you get more doubles than any other player you will lose, or that if you lose you must have bought the most properties? We will answer this question, and won’t need to know anything about Monopoly. Instead we will look at the logical form of the statement. First though, let’s back up and make sure we are very clear on some basics.

Definition 2. A statement is any declarative sentence which is either true or false.

Example: These are statements:
  • Telephone numbers in the USA have 10 digits.
  • The moon is made of cheese.
  • 42 is a perfect square.
  • Every even number greater than 2 can be expressed as the sum of two primes.

And these are not:
  • Would you like some cake?
  • The sum of two squares.
  • 1 + 3 + 5 + 7 + · · · + 2n + 1.
  • Go to your room!
  • This sentence is false.\(^1\)
  • That’s what she said.

\(^1\)This is a tricky one. Remember, a sentence is only a statement if it is either true or false. Here, the sentence is not false, for if it were, it would be true. It is not true, for that would make it false.
The reason the last sentence is not a statement is because it contains variables ("that" and "she"). Unless those are specified, the sentence cannot be true or false, and as such not a statement. Other examples of this: \( x + 3 = 7 \). Depending on \( x \), this is either true or false, but as it stands it it neither.

You can build more complicated statements out of simpler ones using logical connectives. For example, this is a statement: Telephone numbers in the USA have 10 digits and 42 is a perfect square. Note that we can break this down into two smaller statements. The two shorter statements are connected by an “and”. We will consider 5 connectives: “and” (Sam is a man and Chris is a woman), “or” (Sam is a man or Chris is a woman), “if...then...” (if Sam is a man, then Chris is a woman), “if and only if” (Sam is a man if and only if Chris is a woman), and “not” (Sam is not a man).

Since we rarely care about the content of the individual statements, we can replace them with variables. We use capital letters in the middle of the alphabet for these propositional (or sentential) variables: \( P, Q, R, S \ldots \). We also have symbols for the logical connectives: \( \land, \lor, \rightarrow, \leftrightarrow, \neg \).

**Logical Connectives**

- \( P \land Q \) means \( P \) and \( Q \), called a **conjunction**.
- \( P \lor Q \) means \( P \) or \( Q \), called a **disjunction**.
- \( P \rightarrow Q \) means if \( P \) then \( Q \), called an **implication** or **conditional**.
- \( P \leftrightarrow Q \) means \( P \) if and only if \( Q \), called a **biconditional**.
- \( \neg P \) means not \( P \), called a **negation**.

The logical connectives allow us to construct longer statements out of simpler statements. But the result is still a statement - it is either true or false. The truth value can be determined by the truth or falsity of the parts, depending on the connectives.

**Truth Conditions for Connectives**

- \( P \land Q \) is true when both \( P \) and \( Q \) are true.
- \( P \lor Q \) is true when \( P \) or \( Q \) or both are true.
- \( P \rightarrow Q \) is true when \( P \) is false or \( Q \) is true or both.
- \( P \leftrightarrow Q \) is true when \( P \) and \( Q \) are both true, or both false.
- \( \neg P \) is true when \( P \) is false.
I think all of these are obvious, except for $P \rightarrow Q$. Consider the statement, “If Bob gets a 90 on the final, then Bob will pass the class.” This is definitely an implication: $P$ is the statement “Bob gets a 90 on the final” and $Q$ is the statement “Bob will pass the class.” Suppose I said that to you - in what circumstances would it be fair to call me a liar? What if Bob really did get a 90 on the final, and he did pass the class? Then I have not lied to you, so my statement is true. But if Bob did get a 90 on the final and did not pass the class, then I lied, making the statement false. The tricky case is this: what if Bob did not get a 90 on the final? Maybe he passes the class, maybe he doesn’t - did I lie in either case? I think not. In these last two cases, $P$ was false, and so was the statement $P \rightarrow Q$. In the first case, $Q$ was true, and so was $P \rightarrow Q$. So $P \rightarrow Q$ is true when either $P$ is false or $Q$ is true. Perhaps an easier way to look at it is this: $P \rightarrow Q$ is *false* in only one case: if $P$ is true and $Q$ is false. Otherwise, $P \rightarrow Q$ is true. Admittedly, there are times in English when this is not how “if. . . , then. . . ” works. However, in mathematics, we *define* the implication to work this way.

### 1.1.1 Truth Tables

In order to answer our question about monopoly, we need to decide when the statement $(P \rightarrow Q) \lor (Q \rightarrow R)$ is true. Using the rules above, we see that either $P \rightarrow Q$ is true or $Q \rightarrow R$ is true (or both). Those are true if either $P$ is false or $Q$ is true (in the first case) and $Q$ is false or $R$ is true (in the second case). So... yeah, it gets kind of messy. Luckily, we can make a chart to keep track of all the possibilities. Enter truth tables. The idea is this: on each row, we list a possible combination of T’s and F’s (standing of course, for true and false) for each of the sentential variables, and then mark down whether the statement in question is true or false in that case. We do this for every possible combination of T’s and F’s. Then we can clearly see in which cases the statement is true or false. For complicated statements, we will first fill in values for each part of the statement, as a way of breaking up our task into smaller, more manageable pieces.

All you really need to know is the truth tables for each of the logical connectives. Here they are:

<table>
<thead>
<tr>
<th>$P$</th>
<th>$Q$</th>
<th>$P \land Q$</th>
<th>$P$</th>
<th>$Q$</th>
<th>$P \lor Q$</th>
<th>$P$</th>
<th>$Q$</th>
<th>$P \rightarrow Q$</th>
<th>$P$</th>
<th>$Q$</th>
<th>$P \leftrightarrow Q$</th>
</tr>
</thead>
<tbody>
<tr>
<td>T</td>
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</table>

The truth table for negation looks like this:

<table>
<thead>
<tr>
<th>$P$</th>
<th>$\neg P$</th>
</tr>
</thead>
<tbody>
<tr>
<td>T</td>
<td>F</td>
</tr>
<tr>
<td>F</td>
<td>T</td>
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</tbody>
</table>

None of these truth tables should come as a surprise - they are all just restating the definitions of the connectives. Let’s try another one:
**Example:** Make a truth table for the statement $\neg P \lor Q$.

**Solution:** Note that this statement is not $\neg(P \lor Q)$, the negation belongs to $P$ alone. Here is the truth table:

<table>
<thead>
<tr>
<th>$P$</th>
<th>$Q$</th>
<th>$\neg P$</th>
<th>$\neg P \lor Q$</th>
</tr>
</thead>
<tbody>
<tr>
<td>T</td>
<td>T</td>
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</tbody>
</table>

We added a column for $\neg P$ to make filling out the last column easier. The entries in the $\neg P$ column were determined by the entries in the $P$ column. Then to fill in the final column, I look only at the column for $Q$ and the column for $\neg P$ and use the rule of $\lor$.

You might notice that the final column is identical to the final column in the truth table for $P \rightarrow Q$. Since we listed the possible values for $P$ and $Q$ in the same (in fact, standard) order, this says that $\neg P \lor Q$ and $P \rightarrow Q$ are logically equivalent.

Now let’s answer our question about monopoly:

**Example:** Analyze the statement, “if you get more doubles than any other player you will lose, or that if you lose you must have bought the most properties,” using truth tables.

**Solution:** Represent the statement in symbols as $(P \rightarrow Q) \lor (Q \rightarrow R)$, where $P$ is the statement “you get more doubles than any other player,” $Q$ is the statement “you will lose,” and $R$ is the statement “you must have bought the most properties.” Now make a truth table.

The truth table needs to contain 8 rows in order to account for every possible combination of truth and falsity among the three statements. Here is the full truth table:

<table>
<thead>
<tr>
<th>$P$</th>
<th>$Q$</th>
<th>$R$</th>
<th>$P \rightarrow Q$</th>
<th>$Q \rightarrow R$</th>
<th>$(P \rightarrow Q) \lor (Q \rightarrow R)$</th>
</tr>
</thead>
<tbody>
<tr>
<td>T</td>
<td>T</td>
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</tbody>
</table>

The first three columns are simply a systematic listing of all possible combinations of $T$ and $F$ for the three statements (do you see how you would list the 16 possible combinations for four statements?). The next two columns are determined by the values of $P$, $Q$, and $R$ and the definition of implication. Then, the last column is determined by the values in the previous two columns and the definition of $\lor$. It is this final column we care about.
Notice that in each of the eight possible cases, the statement in question is true. So our statement about monopoly is true (regardless of how many properties you own, how many doubles you roll, or whether you win or lose).

The statement about monopoly is an example of a tautology - a statement which is true on the basis of its logical form alone. Tautologies are always true but they don’t tell us much about the world. No knowledge about monopoly was required to determine that the statement was true. In fact, it is equally true that “If the moon is made of cheese, then Elvis is still alive, or if Elvis is still alive, then unicorns have 5 legs.”

1.1.2 Deductions

In the introduction, we claimed that the following was a valid argument:

If Edith eats her vegetables, then she can have a cookie. Edith ate her vegetables. Therefore Edith gets a cookie.

How do we know this is valid? Let’s look at the form of the statements. Let P denote “Edith eats her vegetables” and Q denote “Edith can have a cookie.” The logical form of the argument is then:

\[ P \rightarrow Q \]
\[ P \quad \therefore \quad Q \]

This is an example of a deduction rule - a logical form of an argument which is always valid. This one is a particularly famous rule called modus ponens. Are you convinced that it is a valid deduction rule? If not, consider the following truth table:

<table>
<thead>
<tr>
<th>P</th>
<th>Q</th>
<th>P → Q</th>
</tr>
</thead>
<tbody>
<tr>
<td>T</td>
<td>T</td>
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</table>

In this case, the truth table is identical to the truth table of \( P \rightarrow Q \), but what matters here is that all the lines in the deduction rule have their own column in the truth table. Now remember that an argument is valid provided the conclusion must be true given that the premises are true. The premises in this case are \( P \rightarrow Q \) and \( P \). Which rows of the truth table correspond to both of these being true? \( P \) is true in the first two rows, and of those, only the first row has \( P \rightarrow Q \) true as well. And low-and-behold, in this one case, \( Q \) is true as well. So if \( P \rightarrow Q \) and \( P \) are both true, we see that \( Q \) must be true as well.

Here are a few more examples.

**Example:** Show that

\[ P \rightarrow Q \]
\[ \neg P \rightarrow Q \]
\[ \therefore \quad Q \]
is a valid deduction rule.

**Solution**: We make a truth table which contains all the lines of the argument form:

<p>| | | | | | |</p>
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</thead>
<tbody>
<tr>
<td>P</td>
<td>Q</td>
<td>P → Q</td>
<td>¬P</td>
<td>¬P → Q</td>
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</table>

(we include a column for ¬P just as a step to help getting the column for ¬P → Q).

Now look at all the rows for which both P → Q and ¬P → Q are true. This happens only in rows 1 and 3. Hey! In those rows Q is true as well, so the argument form is valid (it is a valid deduction rule).

**Example**: Decide whether

\[(P → R) ∨ (Q → R)\]

\[\therefore (P ∨ Q) → R\]

is a valid deduction rule.

**Solution**: Let’s make a truth table containing both statements.

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<tr>
<td>P</td>
<td>Q</td>
<td>R</td>
<td>P ∨ Q</td>
<td>P → R</td>
<td>Q → R</td>
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</table>

Look at the fourth row. In this case, (P → R) ∨ (Q → R) is true, but (P ∨ Q) → R is false. Therefore the argument form is not valid (it is not a valid deduction rule). The truth table tells us more: our premise is true when P is true and Q and R are both false, but the conclusion is false in this case. The same problem occurs when Q is true and P and R are false (row 6).

Notice that if we switch the premise and conclusion, then we do have a valid argument - whenever (P ∨ Q) → R is true, so is (P → R) ∨ (Q → R).

Another way to say this is to state that the statement

\[[(P ∨ Q) → R] → [(P → R) ∨ (Q → R)]\]

is a tautology.
1.2 Rephrasing - Logical Equivalence

When reading or writing a proof, or even just trying to understand a mathematical statement, it can be very helpful to rephrase the statement. But how do you know you are doing so correctly? How do you know the two statements are equivalent?

One way is to make a truth table for each and ensure that the final columns of both are identical. We saw earlier that $P \rightarrow Q$ is logically equivalent to $\neg P \lor Q$ because their truth tables agreed. Now though we can just remember this fact - if we see the statement, “if Sam is a man then Chris is a woman,” we can instead think of it as “Sam is a woman or Chris is a woman.” You might also be tempted to rephrase further: “Sam or Chris is a woman.” This is okay of course, but that this second rephrasing is allowed is due to the meaning of “is,” not any of our logical connectives.

Here are some common logical equivalences which can help rephrase mathematical statements:

**Double Negation**

$\neg\neg P$ is logically equivalent to $P$

Example: “It is not the case that $c$ is not odd” means “$c$ is odd.”

No surprise there. Now let’s see how negation plays with conjunctions and disjunctions.

**De Morgan’s Laws**

$\neg(P \land Q)$ is logically equivalent to $\neg P \lor \neg Q$

$\neg(P \lor Q)$ is logically equivalent to $\neg P \land \neg Q$

Example: “$c$ is not even or $c$ is not prime” means “$c$ is not both odd and prime”

Do you believe De Morgan’s laws? If not, make a truth table for each of them. I think most of us get these right most of the time without thinking about them too hard. If I told you that I had popcorn and goobers at the movies, but then you found out it was opposite day (so my statement was false) then you would agree, I hope, that I either did not have popcorn or did not have goobers (or didn’t have either). You would not insist that I could not have had either.

I should warn you that often in English, we are sloppy about our and’s, or’s and not’s. When you write about mathematics, you should be careful and write what you mean. If you are not sure what to write, rephrasing carefully using De Morgan’s can help you make sure that statement matches your intended meaning.

Now some rules for implications:
Negation of Implication

\[ \neg (P \rightarrow Q) \] is logically equivalent to \( P \land \neg Q \)

In words: the only way for an implication to be false is for the “if” part to be true and the “then” part to be false.

This is very important, and not obvious - implications are tricky. But look at the truth table for \( P \rightarrow Q \):

<table>
<thead>
<tr>
<th>( P )</th>
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</tbody>
</table>

There is only one way for the implication to be false - \( P \) is true and \( Q \) is false. Another way to see that this is true is by using De Morgan’s Laws. We saw earlier that \( P \rightarrow Q \) can be rephrased as \( \neg P \lor Q \) so we have

\[ \neg (P \rightarrow Q) \text{ is logically equivalent to } \neg (\neg P \lor Q) \]

But by De Morgan’s laws, \( \neg (\neg P \lor Q) \) is equivalent to \( \neg \neg P \land \neg Q \). By double negation \( \neg \neg P \) is the same as \( P \).

While we are thinking about implications, we should talk about the converse and contrapositive:

Converse and Contrapositive

- The **converse** of an implication \( P \rightarrow Q \) is the implication \( Q \rightarrow P \). The converse is NOT logically equivalent to the original implication.
- The **contrapositive** of an implication \( P \rightarrow Q \) is the statement \( \neg Q \rightarrow \neg P \). An implication and its contrapositive are logically equivalent.

A related, but lesser used, term is the **inverse** of an implication. The inverse of \( P \rightarrow Q \) is \( \neg P \rightarrow \neg Q \). Notice that the inverse of an implication is the contrapositive of the converse. Read that one more time. Good. So since implications and their contrapositives are logically equivalent, the inverse and converse of an implication are logically equivalent to each other, but not to the original implication.

**Example:** Suppose I tell Sue that if she gets a 93% on her final, she will get an A in the class. Assuming that what I said is true, what can you conclude in the following cases:

(a) Sue gets a 93% on her final.
(b) Sue gets an A in the class.
(c) Sue does not get a 93% on her final.
(d) Sue does not get an A in the class.

Solution: Note first that whenever $P \rightarrow Q$ and $P$ are both true statements, $Q$ must be true as well. For this problem, take $P$ to mean “Sue gets a 93% on her final” and $Q$ to mean “Sue will get an A in the class.”

(a) We have $P \rightarrow Q$ and $P$ so $Q$ follows - Sue gets an A.
(b) You cannot conclude anything - Sue could have gotten the A because she did extra credit for example. Notice that we do not know that if Sue gets an A, then she gets a 93% on her final - that is the converse of the original implication, so it might or might not be true.
(c) The inverse of $P \rightarrow Q$ is $\neg P \rightarrow \neg Q$, which states that if Sue does not get a 93% on the final then she will not get an A in the class. But this does not follow from the original implication. Again, we can conclude nothing - Sue could have done extra credit.
(d) What would happen if Sue does not get an A but did get a 93% on the final. Then $P$ would be true and $Q$ would be false. But this makes the implication $P \rightarrow Q$ false! So it must be that Sue did not get a 93% on the final. Notice now we have the implication $\neg Q \rightarrow \neg P$ - the contrapositive of $P \rightarrow Q$. Since $P \rightarrow Q$ is assumed to be true, we know $\neg Q \rightarrow \neg P$ is true as well - they are equivalent.

As we said above, an implication is not logically equivalent to its converse. Given particular statements $P$ and $Q$, the statements $P \rightarrow Q$ and $Q \rightarrow P$ could both be true, both be false, or one could be true and the other false (in either order). Now if both $P \rightarrow Q$ and $Q \rightarrow P$ are true, then we say that $P$ and $Q$ are equivalent. In fact, we have:

**If and only if**

$$P \leftrightarrow Q \text{ is logically equivalent to } (P \rightarrow Q) \land (Q \rightarrow P).$$

Example: given an integer $n$, it is true that $n$ is even if and only if $n^2$ is even. That is, if $n$ is even, then $n^2$ is even, as well as the converse: if $n^2$ is even then $n$ is even.

You can think of “if and only if” statements as having two parts: an implication as its converse. We might say one is the “if” part, and the other is the “only if” part. We also sometimes say that “if and only if” statements have two directions: a forward direction ($P \rightarrow Q$) and a backwards direction ($P \leftarrow Q$, which is really just sloppy notation for $Q \rightarrow P$).

Let’s think a little about which part is which. Is $P \rightarrow Q$ the “if” part or the “only if” part? Perhaps we should look at an example.
Example: Suppose it is true that I sing if and only if I’m in the shower. We know this means that both if I sing, then I’m in the shower, and also the converse - that if I’m in the shower, then I sing. Let $P$ be the statement, “I sing,” and $Q$ be, “I’m in the shower.” So $P \rightarrow Q$ is the statement “if I sing, then I’m in the shower.” Which part of the if and only if statement is this?

What we are really asking is what is the meaning of “I sing if I’m in the shower” and “I sing only if I’m in the shower.” When is the first one (the “if” part) false? Well if I am in the shower but not singing. That is the same condition on being false as the statement “if I’m in the shower, then I sing.” So the “if” part is $Q \rightarrow P$. On the other hand, to say, “I sing only if I’m in the shower” is equivalent to saying “if I sing, then I’m in the shower,” so the only if part is $P \rightarrow Q$.

It is not terribly important to know which part is the if or only if part, but this does get at something very very important: THERE ARE MANY WAYS TO STATE AN IMPLICATION! The problem is, since these are all different ways of saying the same implication, we cannot use truth tables to analyze the situation. Instead, we just need good English skills.

Example: Rephrase the implication, “if I dream, then I am asleep” in as many different ways as possible. Then do the same for the converse.

Solution: The following are all equivalent to the original implication:

1. I am asleep if I dream.
2. I dream only if I am asleep.
3. In order to dream, I must be asleep.
4. To dream, it is necessary that I am asleep.
5. To be asleep, it is sufficient to dream.
6. I am not dreaming unless I am asleep.

The following are equivalent to the converse - if I am asleep, then I dream:

1. I dream if I am asleep.
2. I am asleep only if I dream.
3. It is necessary that I dream in order to be asleep.
4. It is sufficient that I be asleep in order to dream.
5. If I don’t dream, then I’m not asleep.

Hopefully you agree with the above example. We include the “necessary and sufficient” versions because those are common when discussing mathematics. In fact, let’s agree once and for all what they mean:
Necessary and Sufficient

- “P is necessary for Q” means \( Q \rightarrow P \).
- “P is sufficient for Q” means \( P \rightarrow Q \).
- If P is necessary and sufficient for Q, then \( P \leftrightarrow Q \).

To be honest, I have trouble with these if I’m not very careful. I find it helps to have an example in mind:

**Example:** Recall from calculus, if a function is differentiable at a point \( c \), then it is continuous at \( c \), but that the converse of this statement is not true (for example, \( f(x) = |x| \) at the point 0). Restate this fact using necessary and sufficient language.

**Solution:** It is true that in order for a function to be differentiable at a point \( c \), it is necessary for the function to be continuous at \( c \). However, it is not necessary that a function be differentiable at \( c \) for it to be continuous at \( c \).

It is true that to be continuous at a point \( c \), is is sufficient that the function be differentiable at \( c \). However, it is not the case that being continuous at \( c \) is sufficient for a function to be differentiable at \( c \).

**Quantifiers and Sets**

So far we have seen how statements can be combined with logical symbols. This is helpful when trying to understand a complicated mathematical statement - you can determine under which conditions the complicated statement is true. Additionally, we have been able to analyze the logical form of arguments to decide which arguments are valid and which are not. However, the types of statements we have been able to make so far has be sorely limited.

For example, consider a classic argument:

All men are mortal.
Socrates is a man.
Therefore, Socrates is mortal.

This is clearly a valid argument - it is an example of a *syllogism*. Historically, the study of logic began with Aristotle who worked out all possible forms of syllogisms and decided which were valid and which were not. We will not do that here. However, this is an important example because it highlights a limitation of the propositional logic we have studied so far. Can we use propositional logic to analyze the argument?

The trouble is that we don’t have a way to translate “all men are mortal.” It looks like an implication - being a man implies you are mortal. So maybe it is \( P \rightarrow Q \). But what is
P? We could rephrase: “if Socrates is a man, then Socrates is mortal.” Now we have a valid argument form we have seen before. But it is not quite the same. (What if the argument was: All men are mortal, all mortals have hair, therefore all men have hair - also valid and Socrates has nothing to do with it.) Or perhaps we could go with, “for every thing there is, if the thing is a man, then the thing is mortal.” Looks promising, but we still can’t let P be “the thing is a man” - that is not a statement because “thing” is a variable.

One way to sort this mess out is to introduce a new sort of logic called predicate logic. This is the logic of properties. Doing so will allow us to discuss how properties of various things are related. Above, if a thing has the property of being a man, then it has the property of being mortal. We can then quantify over what things we talk about. The the example above, all things.

Another way to accomplish much of the same goal is to use set theory. This will also be useful in solving other types of problems later on. The idea here is that we want talk about collections of things - for example, the collection of all men, and the collection of all mortal things. We can then express “all men are mortal” by saying that the set of men is a subset of the set of mortals. Then we claim that Socrates is a member of the set of men, so therefore is also a member of the set of mortals.

One last example to highlight these two different approaches before delving into the details of each. We all agree that all squares are rectangles. The set theory approach would be to consider the set of squares and the set of rectangles, and point out that one is a subset of the other (the squares are a subset of the rectangles). The predicate logic approach would be to consider the properties of “being a square” and of “being a rectangle” and assign these to predicates - say S and R. We would then say \( \forall x (S(x) \rightarrow R(x)) \) - for all things, if the thing is a square, then it is a rectangle. So having the property of being a square implies having the property of being a rectangle.

Now some details.

### 1.3 Quantifiers and Predicate Logic

Consider the statement “for all integers \( a \) and \( b \), if \( ab \) is even, then \( a \) is even or \( b \) is even.” If we use propositional logic to analyze this statement, what should \( P \) be?

You might want to say \( P \) is “\( ab \) is even” so \( Q \) can be “\( a \) is even” and \( R \) can be “\( b \) is even” and say that the whole statement is therefore of the form \( P \rightarrow (Q \lor R) \). This does not work! One reason is that “\( ab \) is even” is not a statement - it contains free variables, so it is not true or false (until the variables have values). Another reason is that we have just lost the “for all integers \( a \) and \( b \)” from our statement. We can remedy both these problems using Predicate Logic.

#### 1.3.1 Predicates

We can think of predicates as properties of objects. For example, consider the predicate \( E \) which we will use to mean “is even.” Being even is a property of some numbers, so \( E \) needs to be applied to something. We will adopt the notation \( E(x) \) to mean \( x \) is even. (Some books would write \( Ex \) instead.) Notice that if we put a number in for \( x \), then this becomes
a statement - and as such can be true or false. So \( E(2) \) is true, and \( E(3) \) is false. On the other hand \( E(x) \) is not true or false, since we don’t know what \( x \) is. If we have a variable floating around like that, we say the expression is merely a formula, and not a statement.

Since \( E(2) \) is a statement (a proposition), we can apply propositional logic to it. Consider

\[
E(2) \land \lnot E(3)
\]

which is a true statement, because it is both the case that 2 is even and that 3 is not even.

What we have done here is capture the logical form (using connectives) of the statement “2 is even and 3 is not” as well as the mathematical content (using predicates).

Notice that we can only assert even-ness of a single number at a time. That is to say, \( E \) is a one-place predicate. There are also predicates which assert a property of two or more numbers (or other objects) at the same time. Consider the two-place predicate “is less than.” Perhaps we will use the variable \( L \). Now we can say \( L(2, 3) \), which is true because 2 is less than 3. Of course we are already have a symbol for this: \( 2 < 3 \). However, what about “divides evenly into” as a predicate? We can say \( D(2, 10) \) is true because 2 divides evenly into 10, while \( D(3, 10) \) is false since there is a remainder when you divide 10 by 3. Incidentally, there is a standard mathematical symbol for this: \( 2 | 10 \) is read “2 divides 10.”

Predicates can be as complicated and have as many places as we want or need. For example, we could \( R(x, y, z, u, v, w) \) be the predicate asserting that \( x, y, z \) are distinct natural numbers whose only common factor is \( u \), the difference between \( x \) and \( y \) is \( v \) and the difference between \( y \) and \( z \) is \( w \). This is a silly and most likely useless example, but it is an example of a predicate. It is true of some ordered lists of six numbers (6-tuples), and false of others. Additionally, predicates need not have anything to do with numbers: we could let \( F(a, b, c, d) \) be the predicate that asserts that \( a \) and \( b \) are the only two children of mother \( c \) and father \( d \).

1.3.2 Quantifiers

Perhaps the most important reason to use predicate logic is that doing so allows for quantification. We can now express statements like “ever natural number is either even or odd,” and “there is a natural number such that no number is less than it.” Think back to Calculus and the Mean Value Theorem. It states that for every function \( f \) and every interval \((a, b)\), if \( f \) is continuous on the interval \([a, b]\) and differentiable on the interval \((a, b)\), then there exists a number \( c \) such that \( a \leq c \leq b \) and \( f'(c)(b - a) = f(b) - f(a) \). Using the correct predicates and quantifiers, we could express this statement entirely in symbols.

There are two quantifiers we will be interested in: existential and universal.

### Quantifiers

- The existential quantifier is \( \exists \) and is read “there exists” or “there is.” For example,

\[
\exists x(x < 0)
\]

asserts that there is a number less than 0.
The universal quantifier is \( \forall \) and is read “for all” or “every.” For example,

\[
\forall x(x \geq 0)
\]

asserts that every number is greater than or equal to 0.

Are these statements true? Well, first notice that they cannot both be true. In fact, they assert exactly the opposite of each other. (Note that \( x < y \leftrightarrow \neg(x \geq y) \) – although you might wonder what \( x \) and \( y \) are here, so it might be better to say \( \forall x \forall y (x < y \leftrightarrow \neg(x \geq y)) \).

Which one is it though? The answer depends entirely on our domain of discourse - the universe over which we quantify. Usually, this universe is clear from the context. If we are only discussing the natural numbers, then \( \forall x \ldots \) means “for every natural number \( x \ldots \).” On the other hand, in calculus we care about the real numbers, so it would mean “for every real number \( x \ldots \).” If the context is not clear, we might right \( \forall x \in \mathbb{N} \ldots \) to mean “for every natural number \( x \ldots \).” Of course, for the two statements above, the second is true of the natural numbers, the first is true for any larger universe.\(^2\)

Some more examples: to say “every natural number is either even or odd,” we would write, using \( E \) and \( O \) as the predicates for even and odd respectively:

\[
\forall x(E(x) \vee O(x))
\]

To say “there is a number such that no number is less than it” we would write:

\[
\exists x \forall y(y \geq x)
\]

Actually, I did a little translation before I wrote that down. The above statement would be literally read “there is a number such that every number is greater than or equal to it.” This of course amounts to the same thing. However, if I wanted to be exact, I could have also written:

\[
\exists x \neg \exists y(y < x).
\]

Notice also that say that there is a number for which no number is smaller is equivalent to saying that it is not the case that for every number there is a number smaller than it:

\[
\neg \forall x \exists y(y < x).
\]

That these three statements are equivalent is no coincidence. To understand what is going on, we will need to better understand how quantification interacts with the logical connectives, specifically negation.

### 1.3.3 Quantifiers and Connectives

What does it mean to say that it is false that there is something that has a certain property? Well, it means that everything does not have that property. What does it mean for it to be

\(^2\)In this class we take the natural numbers to be 0, 1, 2, 3, . . .
false that everything has a certain property? It means that there is something that doesn’t have the property. So in symbols, we have the following

<table>
<thead>
<tr>
<th>Quantifiers and Negation</th>
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<tbody>
<tr>
<td>$\neg \forall x P(x)$ is equivalent to $\exists x \neg P(x)$</td>
</tr>
<tr>
<td>and</td>
</tr>
<tr>
<td>$\neg \exists x P(x)$ is equivalent to $\forall x \neg P(x)$.</td>
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</tbody>
</table>

In other words, to move a negation symbol past a quantifier, you must switch the quantifier. This can be done multiple times:

$$
\neg \exists x \forall y \exists z P(x, y, z) \text{ is equivalent to } \forall x \exists y \forall z \neg P(x, y, z).
$$

Now we also know how to move negation symbols through other connectives (using De Morgan’s Laws) so it is always possible to rewrite a statement so that the only negation symbols that appear are right in front of a predicate. This hints at the possibility of having a standard form for all predicate statements. However, to get this we must also understand how to move quantifiers through connectives.

Before we get too excited, note that we only need to worry about two connectives: $\land$ and $\lor$. This is because we can rewrite $p \rightarrow q$ as $\neg p \lor q$ (they are logically equivalent) and $p \leftrightarrow q$ as $(p \land q) \lor (\neg p \land \neg q)$ (also logically equivalent).

Let us consider an example to see what can happen.

**Example:** Let $E$ be the predicate for being even, and $O$ for being odd.

Consider:

$$
\exists x E(x) \land \exists x O(x),
$$

which says that there is a number which is even and a number which is odd. This is of course true. However there is no number which is both even and odd, so

$$
\exists x (E(x) \land O(x))
$$

is false. Note also that

$$
\exists x (E(x) \lor O(x))
$$

while true, is not really the same thing – if $O$ is instead the predicate for “is less than 0” then the original statement is false, but this new one is true (of the natural numbers). Changing the quantifier also doesn’t help:

$$
\forall x (E(x) \land O(x))
$$

is false. So what can we do?

The problem is that in the original sentence, the variable $x$ is doing double duty. We want to express the fact that there is an even number and an odd
number. But that even number is in no way related to that odd number. So we might as well have said
\[ \exists x E(x) \land \exists y O(y). \]

Now we can move the quantifiers out:
\[ \exists x \exists y (E(x) \land O(y)). \]

The same thing works with \( \lor \) and for \( \forall \) with either connective. As long as there is no repeat in quantified variables we can move the quantifiers outside of conjunctions and disjunctions.

A warning though: you cannot do this for \( \rightarrow \), at least not directly.\(^3\) Let’s see what happens.

**Example:** Consider
\[ \forall x P(x) \rightarrow \exists y Q(y) \]
for some predicates \( P \) and \( Q \). This sentence is **not** the same as
\[ \forall x \exists y (P(x) \rightarrow Q(y)). \]

Remember that \( p \rightarrow q \) is the same as \( \neg p \lor q \). So the original sentence is really
\[ \neg \forall x P(x) \lor \exists y Q(y). \]

Before we move the quantifiers out, we must move the \( \forall x \) past the negation sign, which switches it to a \( \exists x \):
\[ \exists x \neg P(x) \lor \exists y Q(y). \]

Then we can finish by writing,
\[ \exists x \exists y (\neg P(x) \lor Q(y)) \]
or equivalently
\[ \exists x \exists y (P(x) \rightarrow Q(y)). \]

### 1.4 Sets and Set Notation

For us, a set will simply be an unordered collection of objects. For example, we could consider the set of all students enrolled at UNC this semester. Or the set of natural numbers between 1 and 10 inclusive. In the first case, each student here is a element (or member) of the set, while Barack Obama, among many others, is not an element of the set. Also, the

\(^3\)We must be similarly careful with \( \leftrightarrow \)
two example are of different sets. Two sets are equal exactly if they contain the exact same elements.

Because we will want to consider many examples, we should have some notation to make talking about sets easier. Consider,

\[ A = \{1, 2, 3\} \]

This is read, “\(A\) is the set containing the elements 1, 2 and 3.” We use curly braces “\{, \}” to enclose elements of a set. Some more notation:

\[ a \in \{a, b, c\} \]

The symbol “\(\in\)” is read “is in” or “is an element of.” Thus the above means that \(a\) is an element of the set containing the letters \(a\), \(b\), and \(c\). Note that this is a true statement. It would also be true to say that \(d\) is not in that set:

\[ d \notin \{a, b, c\} \]

Be warned: we say “\(x \in A\)” when we wish to express that “one of the elements of the set \(A\) is \(x\).” For example, consider the set,

\[ A = \{1, b, \{x, y, z\}, \emptyset\} \]

This is a strange set, to be sure. It contains four elements: the number 1, the letter \(b\), the set \(\{x, y, z\}\) and the empty set (\(\emptyset = \{\}\), the set containing no elements). Is \(x\) in \(A\)? The answer is no. None of the four elements in \(A\) are the letter \(x\), so we must conclude that \(x \notin A\). Similarly, if we considered the set \(B = \{1, b\}\), then again \(B \notin A\) - even though the elements of \(B\) are also elements of \(A\), we cannot say that the thing \(B\) is one of the things in the collection \(A\).

If a set is finite, then we can describe it by simply listing the elements. Infinite sets exists though, so we need to be able to describe them as well. For instance, if we want \(A\) to be the set of all even natural numbers, would could write,

\[ A = \{2, 4, 6, \ldots\} \]

but this is a little imprecise. Better would be

\[ A = \{x \in \mathbb{N} : \exists n \in \mathbb{N}(x = 2n)\} \]

Breaking that down: \(x \in \mathbb{N}\) means \(x\) is in the set \(\mathbb{N}\) (the set of natural numbers), \(:\) is read “such that” and \(\exists n(x = 2n)\) is read “there exists an \(n\) in the natural numbers for which \(x\) is two times \(n\)” (in other words, \(x\) is even). Slightly easier might be,

\[ A = \{x : x \text{ is even}\} \]

Note: sometime people use \(|\) for the “such that” symbol instead of \(:\).
## Set Theory Notation

<table>
<thead>
<tr>
<th>Symbol:</th>
<th>Read:</th>
<th>Example:</th>
</tr>
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<tbody>
<tr>
<td>{, }</td>
<td>braces</td>
<td>{1,2,3}. The braces enclose the elements of a set. This is the set which contains the numbers 1, 2 and 3.</td>
</tr>
<tr>
<td>:</td>
<td>such that</td>
<td>{x : x &gt; 2} is the set of all x such that x is greater than 2.</td>
</tr>
<tr>
<td>∈</td>
<td>is an element of</td>
<td>2 ∈ {1,2,3} asserts that 2 is one of the elements in the set {1,2,3}. However, 4 ∉ {1,2,3}.</td>
</tr>
<tr>
<td>⊆</td>
<td>is a subset of</td>
<td>A ⊆ B asserts that every element of A is also an element of B.</td>
</tr>
<tr>
<td>⊂</td>
<td>is a proper subset of</td>
<td>A ⊂ B asserts that every element of A is also an element of B, but A ≠ B.</td>
</tr>
<tr>
<td>∩</td>
<td>intersection</td>
<td>A ∩ B is the set of all elements which are elements of both A and B.</td>
</tr>
<tr>
<td>∪</td>
<td>union</td>
<td>A ∪ B is the set of all elements which are elements of A or B or both.</td>
</tr>
<tr>
<td>\</td>
<td>set difference</td>
<td>A \ B is the set of all elements of A which are not elements of B.</td>
</tr>
<tr>
<td>A</td>
<td>compliment (of A)</td>
<td>A is the set of everything which is not an element of A. The A can be any set here.</td>
</tr>
<tr>
<td></td>
<td>cardinality (of A)</td>
<td>(</td>
</tr>
</tbody>
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### Special sets

- **∅** The *empty set* is the set which contains no elements.
- **U** The *universe set* is the set of all elements.
- **N** The set of natural numbers. That is, \(N = \{0,1,2,3\ldots\}\)
- **Z** The set of integers. \(Z = \{\ldots,-2,-1,0,1,2,3\ldots\}\)
- **Q** The set of rational numbers.
- **R** The set of real numbers.
- **P(A)** The *power set* of any set A is the set of all subsets of A.
1.4.1 Relationships between sets

We have already said what it means for two sets to be equal: they have exactly the same elements. Thus, for example,

\{1, 2, 3\} = \{2, 1, 3\}.

(Remember, the order the elements are written down in does not matter.) Also,

\{1, 2, 3\} = \{I, II, III\}.

Now what about the sets \(A = \{1, 2, 3\}\) and \(B = \{1, 2, 3, 4\}\)? Clearly \(A \neq B\). However, we can notice that every element of \(A\) is also an element of \(B\). Because of this, we say that \(A\) is a subset of \(B\), or in symbols \(A \subset B\) or \(A \subseteq B\). (Both symbols are read “is a subset of.” The difference is that sometimes we want to say that \(A\) is either equal to or a subset of \(B\), in which case we use \(\subseteq\). Compare the difference between \(<\) and \(\leq\).)

**Example:** Let \(A = \{1, 2, 3, 4, 5, 6\}\), \(B = \{2, 4, 6\}\), \(C = \{1, 2, 3\}\) and \(D = \{7, 8, 9\}\). Determine which of the following are true, false, or meaningless.

1. \(A \subset B\)
2. \(B \subset A\)
3. \(B \in C\)
4. \(\emptyset \in A\)
5. \(\emptyset \subset A\)
6. \(A < D\)
7. \(3 \in C\)
8. \(3 \subset C\)
9. \(\{3\} \subset C\)

**Solution:**

1. False.
2. True: every element in \(B\) is an element in \(A\).
3. False: the elements in \(C\) are 1, 2, and 3. The set \(B\) is not equal to 1, 2, or 3.
4. False: \(A\) has exactly 6 elements, and none of them are the empty set.
5. True: Everything in the empty set (nothing) is also an element of \(A\). Notice that the empty set is a subset of every set.
6. Meaningless. A set cannot be less that another set.
7. True.
8. Meaningless. 3 is not a set, so it cannot be a subset of another set.
9. True. 3 is the only element of the set \(\{3\}\), and is an element of \(C\), so every element in \(\{3\}\) is an element of \(C\).

In the example above, \(B\) is a subset of \(A\). You might wonder what other sets are subsets of \(A\). If you collect all these subsets of \(A\), they themselves for a set - a set of sets. We call the set of all subsets of \(A\) the **power set** of \(A\), and write it \(\mathcal{P}(A)\).
Example: Let \( A = \{1, 2, 3\} \). Find \( \mathcal{P}(A) \).

Solution: \( \mathcal{P}(A) \) is a set of sets - all of which are subsets of \( A \). So

\[
\mathcal{P}(A) = \{\emptyset, \{1\}, \{2\}, \{3\}, \{1, 2\}, \{1, 3\}, \{2, 3\}, \{1, 2, 3\}\}
\]

Notice that while \( 2 \in A \), it is wrong to write \( 2 \in \mathcal{P}(A) \) - none of the elements in \( \mathcal{P}(A) \) are numbers! On the other hand we do have \( \{2\} \in \mathcal{P}(A) \) because \( \{2\} \subseteq A \).

What does a subset of \( \mathcal{P}(A) \) look like? Notice that \( \{2\} \not\subseteq \mathcal{P}(A) \) because not everything in \( \{2\} \) is in \( \mathcal{P}(A) \). But we do have \( \{\{2\}\} \subseteq \mathcal{P}(A) \). The only element of \( \{\{2\}\} \) is the set \( \{2\} \) which is also an element of \( \mathcal{P}(A) \). We could take the collection of all subsets of \( \mathcal{P}(A) \) and call that \( \mathcal{P}(\mathcal{P}(A)) \). Or even the power set of that set of sets of sets.

Another way to compare sets is by their size. Notice that in the example above, \( A \) has 6 elements, \( B, C, \) and \( D \) all have 3 elements. The size of a set is called the set’s cardinality. We would write \( |A| = 6, \ |B| = 3 \) and so on. For sets that have a finite number of elements, the cardinality of the set is simply the number of elements in the set. Note that the cardinality of \( \{1, 2, 3, 2, 1\} \) is 3 – we do not count repeats (in fact, \( \{1, 2, 3, 2, 1\} \) is exactly the same set as \( \{1, 2, 3\} \)). There are sets with infinite cardinality, such as \( \mathbb{N} \), the set of rational numbers (written \( \mathbb{Q} \)), the set of even natural numbers, the set of real numbers (\( \mathbb{R} \)). It is possible to distinguish between different infinite cardinalities, but that is beyond the scope of these notes. For us, a set will either be infinite, or finite, and if it is finite, we can determine it’s cardinality by counting elements.

Example: Find the cardinality of \( \{23, 24, \ldots, 37, 38\} \).

Solution: Since \( 38 - 23 = 15 \), we can conclude that the cardinality of the set is 16 (you need to add one since 23 is included).

1.4.2 Operations on sets

Is it possible to add two sets? Not really, however there is something similar. If we want to combine two sets – to get the collection of objects that are in either set, then we can take the union of the two sets. Symbolically,

\[
C = A \cup B
\]

means \( C \) is the union of \( A \) and \( B \). Every element of \( C \) is either an element of \( A \) or an element of \( B \) (or an element of both). For example, if \( A = \{1, 2, 3\} \) and \( B = \{2, 3, 4\} \), then \( A \cup B = \{1, 2, 3, 4\} \).

The other common operation on sets is intersection. We write,

\[
C = A \cap B
\]

to mean that \( C \) is the intersection of \( A \) and \( B \); everything in \( C \) is in both \( A \) and in \( B \). So if \( A = \{1, 2, 3\} \) and \( B = \{2, 3, 4\} \), then \( A \cap B = \{2, 3\} \).
Often when dealing with sets, we will have some understanding as to what “everything” is. Perhaps we are only concerned with natural numbers. We would say that our universe is $\mathbb{N}$. Sometimes we call denote this universe by $U$. Given this context, we might wish to speak of all the elements which are not in a particular set. We call this the compliment of the set, and write, 

$$B = A^c$$

when $B$ contains every element not contained in $A$. So if our universe is $\{1, 2, \ldots, 9, 10\}$, and $A = \{2, 3, 5, 7\}$, then $A^c = \{1, 4, 6, 8, 9, 10\}$.

Of course we can perform more than one operation at a time. For example, consider $A \cap B$

This is the set of all elements which are both elements of $A$ and not elements of $B$. What have we done? We’ve started with $A$ and removed all of the elements which were in $B$. Another way to write this is the set difference:

$$A \cap B = A \setminus B$$

It is important to remember that these operations (union, intersection, compliment and difference) on sets produce other sets. Don’t confuse these with the symbols from the previous section (element of and subset of). $A \cap B$ is a set, while $A \subseteq B$ is true or false. This is the same difference as between $3 + 2$ (which is a number) and $3 \leq 2$ (which is in this case false).

**Example:** Let $A = \{1, 2, 3, 4, 5, 6\}$, $B = \{2, 4, 6\}$, $C = \{1, 2, 3\}$ and $D = \{7, 8, 9\}$. The universe is $U = \{1, 2, \ldots, 10\}$. Find:

1. $A \cup B$
2. $A \cap B$
3. $B \cap C$
4. $A \cap D$
5. $B \cup C$
6. $A \cap B$
7. $(D \cap C) \cup A \cap B$
8. $\emptyset \cup C$
9. $\emptyset \cap C$

**Solution:**

1. $A \cup B = \{1, 2, 3, 4, 5, 6\} = A$ since everything in $B$ is already in $A$.
2. $A \cap B = \{2, 4, 6\} = B$ since everything in $B$ is in $A$.
3. $B \cap C = \{2\}$ - the only element of both $B$ and $C$ is 2.
4. $A \cap D = \emptyset$ since $A$ and $D$ have no common elements
5. $B \cup C = \{5, 7, 8, 9, 10\}$. First we find that $B \cup C = \{1, 2, 3, 4, 6\}$, then we take everything not in that set.
6. $A \cap B = \{1, 3, 5\}$. Everything that is in $A$ which is not in $B$. This is the same as $A \setminus B$.
7. $(D \cap C) \cup A \cap B = \{1, 3, 5, 7, 8, 9\}$. The set contains all elements that are either in $D$ but not in $C$ or not in both $A$ and $B$.
8. $\emptyset \cup C = C$ - nothing is added by the emptyset.
9. $\emptyset \cap C = \emptyset$ - nothing can be both in a set and in the emptyset.

You might notice that the symbols for union and intersection slightly resemble the logic symbols for “or” and “and.” This is no accident. What does it mean for $x$ to be an element of $A \cup B$? It means that $x$ is an element of $A$ or $x$ is an element of $B$ (or both). That is,

$$x \in A \cup B \iff x \in A \lor x \in B.$$  

Similarly,

$$x \in A \cap B \iff x \in A \land x \in B.$$  

Also,

$$x \in \overline{A} \iff \neg(x \in A)$$

which says $x$ is an element of the compliment of $A$ if $x$ is not an element of $A$.

Given all this, you should not be surprised to find out that there is a version of De Morgan’s laws for sets:

**De Morgan’s Laws (for sets)**

For any sets $A$ and $B$ in some universe $U$:

$$A \cap B = \overline{A} \cup \overline{B}$$

$$A \cup B = \overline{A} \cap \overline{B}$$

Do you believe these equations? To check De Morgan’s laws in logic we could just make a truth table for each statement. We don’t have truth tables for sets, but we do have....

### 1.4.3 Venn Diagrams

Union, intersection, compliment, set difference - operations can get complicated (see part 7 in the above example). Luckily, there is a very nice visual tool we can use to clarify things. Venn diagrams represent sets as intersecting circles. We can shade the region we are talking about when we carry out an operation. We can also represent cardinality of a particular set by putting the number in the corresponding region.

![Venn Diagrams](image)

Each circle represents a set. The rectangle containing the circles represents the universe. To represent combinations of these sets, we shade the corresponding region. For example, we could draw $A \cap B$ as:
Here is a representation of $A \cap \overline{B}$, or equivalently $A \setminus B$:

A more complicated example is $(B \cap C) \cup (C \cap \overline{A})$, as seen below.

Notice that the shaded regions above could also be arrived at in another way. We could have started with all of $C$, then excluded the region where $C$ and $A$ overlap (without $B$). That region is $(A \cap B) \cap \overline{B}$. So the above Venn diagram also represents $C \cap \left( (A \cap B) \cap \overline{B} \right)$. So using just the picture, we have determined that

$$(B \cap C) \cup (C \cap \overline{A}) = C \cap \left( (A \cap B) \cap \overline{B} \right).$$

1.5 Logic or Set Theory? Yes.

Understanding predicate logic (including quantifiers) and set theory is useful throughout mathematics. Certainly we will return to sets when we study combinatorics (counting) and elsewhere - sets are one of the basic building blocks of mathematics. That said, our primary focus at the moment is to use logic and set theory for the purpose of communicating mathematics: how to read and write. It is often useful to have multiple ways of expressing ideas. Predicate logic and set theory give us two such ways. Sometimes switching viewpoints is all it takes to get the insight needed to solve a problem.

The truth is, everything you can express using sets, you can also express using predicates, and \textit{visa versa}. This is because “being an element of the set $A$” is a predicate. On the other hand, if we start with a predicate $P(x)$, we can consider the set of all $x$ such that $P(x)$ holds.
Logic Exercises

We saw above that there is a parallel between conjunctions and disjunctions in logic and intersections and unions in set theory. Is there something similar we can do for \( \rightarrow \)? In fact there is. Recall \( A \subseteq B \) means that every element of \( A \) is also an element of \( B \). In other words, if \( x \in A \), then \( x \in B \). In symbols, we would write this as \( x \in A \rightarrow x \in B \). Of course if we had \( x \in A \leftrightarrow x \in B \), then we could conclude \( A = B \).

**Example:** Explain the relationship between even integers and multiples of 4 using both predicate logic and set theory.

**Solution:** In English, you might say that while every multiple of 4 is even, not every even number is a multiple of 4. Using predicate logic, we could use \( E(x) \) to say \( x \) is even, and \( F(x) \) to say \( x \) is a multiple of 4. We then translate the English sentences as follows:

\[
\forall x (F(x) \rightarrow E(x))
\]

and

\[
\exists x (E(x) \land \neg F(x)).
\]

(The second sentence could also have been \( \neg \forall x (E(x) \rightarrow F(x)) \). Do you see why this is equivalent?)

Alternatively, let \( E \) be the set of all even numbers, and \( F \) be the set of all multiples of 4. Since the set of all multiples of 4 is a proper subset of the set of all even numbers, we have \( F \subset E \). We use \( \subset \) and not \( \subseteq \) to represent the fact that \( E \neq F \). We can also express the fact that there are even numbers which are not multiples of 4 by saying that the set difference of \( E \) and \( F \) is non-empty. So \( E \setminus F \neq \emptyset \), or alternatively \( E \cap \overline{F} \neq \emptyset \). In fact, we have \( E \cap F = F \).

**Exercises**

1. Make a truth table for the statement \((P \lor Q) \rightarrow (P \land Q)\).
2. Make a truth table for the statement \(\neg P \land (Q \rightarrow P)\). What can you conclude about \( P \) and \( Q \) if you know the statement is true?
3. Make a truth table for the statement \(\neg P \rightarrow (Q \land R)\).
4. Determine whether the following two statements are logically equivalent: \( \neg(P \rightarrow Q) \) and \( P \land \neg Q \). Explain how you know you are correct.
5. Are the statements \( P \rightarrow (Q \lor R) \) and \( (P \rightarrow Q) \lor (P \rightarrow R) \) logically equivalent?
6. Determine if the following argument form is valid:

\[
\begin{align*}
P \lor Q & \\
\neg P & \\
\therefore & Q
\end{align*}
\]

7. Determine if the following argument form is valid:

\[
\begin{align*}
P \rightarrow (Q \lor R) & \\
\neg (P \rightarrow Q) & \\
\therefore & R
\end{align*}
\]
8. Determine if the following argument form is valid:

\[
\begin{align*}
(P \land Q) & \rightarrow R \\
\neg P \lor \neg Q & \\
\therefore \neg R
\end{align*}
\]

9. Consider the statement about a party, “If it’s your birthday or there will be cake, then there will be cake.”

(a) Translate the above statement into symbols. Clearly state which statement is \( P \) and which is \( Q \).

(b) Make a truth table for the statement.

(c) Assuming the statement is true, what (if anything) can you conclude if there will be cake?

(d) Assuming the statement is true, what (if anything) can you conclude if there will not be cake?

(e) Suppose you found out that the statement was a lie. What can you conclude?

10. Suppose \( P \) and \( Q \) are the statements: \( P \): Jack passed math. \( Q \): Jill passed math.

(a) Translate “Jack and Jill both passed math” into symbols.

(b) Translate “If Jack passed math, then Jill did not” into symbols.

(c) Translate “\( P \lor Q \)” into English.

(d) Translate “\( \neg(P \land Q) \rightarrow Q \)” into English.

(e) Suppose you know that if Jack passed math, then so did Jill. What can you conclude if you know that:

   (a) Jill passed math?

   (b) Jill did not pass math?

11. Geoff Poshingten is out at a fancy pizza joint, and decides to order a calzone. When the waiter asks what he would like in it, he replies, “I want either pepperoni or sausage, and if I have sausage, I must also include quail. Oh, and if I have pepperoni or quail then I must also have ricotta cheese.”

(a) Translate Geoff’s order into logical symbols.

(b) The waiter knows that Geoff is either a liar or a truth-teller (so either everything he says is false, or everything is true). Which is it?

(c) What, if anything, can the waiter conclude about the ingredients in Geoff’s desired calzone?

12. Consider the statement “If Oscar eats Chinese food, then he drinks milk.”

(a) Write the converse of the statement.

(b) Write the contrapositive of the statement.

(c) Is it possible for the contrapositive to be false? If it was, what would that tell you?
(d) Suppose the original statement is true, and that Oscar drinks milk. Can you conclude anything (about his eating Chinese food)? Explain.
(e) Suppose the original statement is true, and that Oscar does not drink milk. Can you conclude anything (about his eating Chinese food)? Explain.

13. Simplify the following statements (so that negation only appears right before variables).
   (a) \( \neg(P \to \neg Q) \)
   (b) \( (\neg P \lor \neg Q) \to \neg(\neg Q \land R) \)
   (c) \( \neg((P \to \neg Q) \lor \neg(R \land \neg R)) \)
   (d) It is false that if Sam is not a man then Chris is a woman, and that Chris is not a woman.

14. Which of the following statements are equivalent to the implication, “if you win the lottery, then you will be rich,” and which are equivalent to the converse of the implication?
   (a) Either you win the lottery or else you are not rich.
   (b) Either you don’t win the lottery or else you are rich.
   (c) You will win the lottery and be rich.
   (d) You will be rich if you win the lottery.
   (e) You will win the lottery if you are rich.
   (f) It is necessary for you to win the lottery to be rich.
   (g) It is sufficient to with the lottery to be rich.
   (h) You will be rich only if you win the lottery.
   (i) Unless you win the lottery, you won’t be rich.
   (j) If you are rich, you must have one the lottery.
   (k) If you are not rich, then you did not win the lottery.
   (l) You will win the lottery if and only if you are rich.

15. Consider the implication, “if you clean your room, then you can watch TV.” Rephrase the implication in as many ways as possible. Then do the same for the converse.

16. Translate into symbols. Use \( E(x) \) for “\( x \) is even” and \( O(x) \) for “\( x \) is odd.”
   (a) No number is both even and odd.
   (b) One more than any even number is an odd number.
   (c) There is prime number that is even.
   (d) Between any two numbers there is a third number.
   (e) There is no number between a number and one more than that number.
17. Translate into English:
(a) \( \forall x (E(x) \rightarrow E(x + 2)) \)
(b) \( \forall x \exists y (\sin(x) = y) \)
(c) \( \forall y \exists x (\sin(x) = y) \)
(d) \( \forall x \forall y (x^3 = y^3 \rightarrow x = y) \)

18. Simplify the statements (so negation appears only directly next to predicates).
(a) \( \neg \exists x \forall y (\neg O(x) \vee E(y)) \)
(b) \( \neg \forall x \neg \forall y (x < y \land \exists z (x < z \vee y < z)) \)
(c) There is a number \( n \) for which no other number is either less \( n \) than or equal to \( n \).
(d) It is false that for every number \( n \) there are two other numbers which \( n \) is between.

19. Let \( A = \{1, 2, 3, 4, 5\} \), \( B = \{3, 4, 5, 6, 7\} \) and \( C = \{2, 3, 5\} \).
(a) Find \( A \cap B \).
(b) Find \( A \cup B \).
(c) Find \( A \setminus B \).
(d) Is \( C \subseteq A \)?
(e) Is \( C \subseteq B \)?

20. Let \( A = \{x \in \mathbb{N} : 3 \leq x \leq 13\} \), \( B = \{x \in \mathbb{N} : x \text{ is even}\} \), and \( C = \{x \in \mathbb{N} : x \text{ is odd}\} \).
(a) Find \( A \cap B \).
(b) Find \( A \cup B \).
(c) Find \( B \cap C \).
(d) Find \( B \cup C \).

21. Find an example of sets \( A \) and \( B \) such that \( A \cap B = \{3, 5\} \) and \( A \cup B = \{2, 3, 5, 7, 8\} \).

22. Find an example of sets \( A \) and \( B \) such that \( A \subseteq B \) and \( A \in B \).

23. Recall \( \mathbb{Z} = \{\ldots, -2, -1, 0, 1, 2, \ldots\} \) (the integers). Let \( \mathbb{Z}^+ \) be the positive integers. Let \( 2\mathbb{Z} \) be the even integers, \( 3\mathbb{Z} \) be the multiples of 3, and so on.
(a) Is \( \mathbb{Z}^+ \subseteq 2\mathbb{Z} \)?
(b) Is \( 2\mathbb{Z} \subseteq \mathbb{Z}^+ \)?
(c) Find \( 2\mathbb{Z} \cap 3\mathbb{Z} \). Describe the set in words, and also in symbols (using a \( : \) symbol).
(d) Express \( \{x \in \mathbb{Z} : \exists y \in \mathbb{Z} (x = 2y \vee x = 3y)\} \) as a union or intersection of two sets above.
24. Let $A_2$ be the set of all multiples of 2 except for 2. Let $A_3$ be the set of all multiples of 3 except for 3. And so on, so that $A_n$ is the set of all multiples of $n$ except for $n$, for any $n \geq 2$. Describe (in words) the set $A_2 \cup A_3 \cup A_4 \cup \cdots$.

25. Draw a Venn diagram to represent each of the following:
   
   (a) $A \cup \overline{B}$
   (b) $(A \cup B)$
   (c) $A \cap (B \cup C)$
   (d) $(A \cap B) \cup C$
   (e) $\overline{A} \cap B \cap \overline{C}$
   (f) $(A \cup B) \setminus C$

26. Describe a set in terms of $A$ and $B$ which has the following Venn diagram:

27. Find the cardinalities:
   
   (a) Find $|A|$ when $A = \{4, 5, 6, \ldots, 37\}$
   (b) Find $|A|$ when $A = \{x \in \mathbb{Z} : -2 \leq x \leq 100\}$
   (c) Find $|A \cap B|$ when $A = \{x \in \mathbb{N} : x \leq 20\}$ and $B = \{x \in \mathbb{N} : x \text{ is prime}\}$

28. Let $A = \{a, b, c\}$. Find $\mathcal{P}(A)$.

29. Let $A = \{1, 2, \ldots, 10\}$. How many subsets of $A$ contain exactly one element (i.e., how many singleton subsets are there). How many doubleton (containing exactly two elements) are there?

30. Let $A = \{1, 2, 3, 4, 5, 6\}$. Find all sets $B \in \mathcal{P}(A)$ which have the property $\{2, 3, 5\} \subseteq B$.

31. Find an example of sets $A$ and $B$ such that $|A| = 4$, $|B| = 5$ and $|A \cup B| = 9$.

32. Find an example of sets $A$ and $B$ such that $|A| = 3$, $|B| = 4$ and $|A \cup B| = 5$.

33. In a regular deck of playing cards there are 26 red cards and 12 face cards. Explain in terms of sets why there are only 32 cards which are either red or a face card.
Chapter 2

Introduction to Graph Theory

Graph Theory is a relatively new area of mathematics, first studied by the super famous mathematician Leonhard Euler in 1735. Since then it has blossomed in to a powerful tool used in nearly every branch of science and is currently an active area of mathematics research. We will begin our study with the problem that started it all: The Seven Bridges of Königsberg.

In the time of Euler, in the town of Königsberg in Prussia, there was a river containing two islands. The islands were connected to the banks of the river by seven bridges (as seen below). The bridges were very beautiful, and on their days off, townspeople would spend time walking over the bridges. As time passed, a question arose: was it possible to plan a walk so that you cross each bridge once and only once? Euler was able to answer this question. Are you?

Try finding a path which uses each bridge exactly once.

Here is another problem: below is a drawing of four dots connected by some lines. Is it possible to trace over each line once and only once (without lifting up your pencil)? You must start and end on one of the dots.
There is an obvious connection between these two problems - any path through the dot and line drawing corresponds exactly to the path over the bridges of Königsberg.

Pictures like the dot and line drawing are called graphs. Graphs are made up of a collection of dots (called vertices) and lines connecting those dots (called edges). The nice thing about looking at graphs instead of pictures of rivers, islands and bridges is that we now have a mathematical object to study. A set of vertices and a set of edges (in fact, we can take the set of edges to be a set of ordered pairs from the set of vertices). We have distilled the “important” parts of the bridge picture for the purposes of the problem - it does not matter how big the islands are, what the bridges are made out of, if the river contains alligators, etc. All that matters is which land masses are connected to which other land masses, and how many times. This was the great insight that Euler had.

We will return to the question of finding paths through graphs later. But first, here are a few other situations you can represent with graphs.

**Example:** Al, Bob, Cam, Dan, and Euclid are all members of the social networking website Facebook. The site allows members to be “friends” with each other. It turns out that Al and Cam are friends, as are Bob and Dan. Euclid is friends with everyone. Represent this situation with a graph.

**Solution:** Each person will be represented by a vertex and each friendship will be represented by an edge - that is, there will be an edge between two vertices if and only if the people represented by those vertices are friends. We get the following graph:

![Graph Example](image)

**Example:** Each of three houses must be connected to each of three utilities. Is it possible to do this without any of the utility lines crossing?

**Solution:** We will answer this question later. For now, notice how we would ask this question in the context of graph theory. We are really asking whether it is possible to redraw the graph below without any edges crossing (except at vertices).

![Graph Example](image)

### 2.1 Basics

While we almost always think of graphs as pictures, these are really just visual representations of mathematical objects. In fact, all a graph is is a set of vertices some pairs of which
are “related” by an edge. For example, we can describe one graph like this: the vertices are the letters \{a, b, c, d\} and the edges are the pairs \{(a, b), (a, c), (b, c), (b, d), (c, d)\}. In other words, we have four vertices, \(a, b, c, \text{ and } d\), and \(a\) is connected (by an edge) to \(b\) and \(c\), \(b\) is connected to both \(c\) and \(d\), and \(c\) is also connected to \(d\). One way to draw this graph is this:

![Graph](image)

However we could also have drawn the graph differently. For example either of these:

![Graph](image)

Viewed as pictures, the graphs above are the same whether or not the vertices are labeled as they are, or at all. In fact, two graphs are equal precisely when there is a way to label the vertices so that all the pairs of vertices which have edges in one graph also have edges in the other graph, and \textit{vice versa}. Two vertices which are connected by an edge are called \textit{adjacent}.

The graph above is also particularly nice in that no pair of vertices is connected more than once, and no vertex is connected to itself. The technical name for graphs like this (no double edges or single edge loops) is called \textit{simple}. The graph is also \textit{connected} - you can get from any vertex to any other vertex by following some path of edges (\(a\) and \(d\) are connected by a path two edges long). A graph that is not connected can be thought of as two separate graphs drawn close together. Unless otherwise stated, we will assume all graphs are connected. However, our graphs might or might not be simple (the graph for the Königsberg bridge problem is \textit{not} simple).

The graph above does not have an edge between \(a\) and \(d\). Thus it is possible to add an edge to the graph (and keep it simple). If we add all possible edges (keeping the graph simple), then the resulting graph is said to be \textit{complete}. That is, a graph is complete if every pair of vertices is connected by an edge. Since a graph is determined completely by which vertices are adjacent to which other vertices, there is only one complete graph with a given number of vertices. We give these a special name: \(K_n\) is the complete graph on \(n\) vertices.

Each vertex in \(K_n\) is adjacent to \(n - 1\) other vertices. We call the number of edges emanating from a given vertex the \textit{degree} of that vertex. So every vertex in \(K_n\) has degree \(n - 1\). How many edges does \(K_n\) have? One might think the answer should be \(n(n - 1)\), since we count \(n - 1\) edges \(n\) times (once for each vertex). However, each edge is adjacent to 2 vertices, so we count every edge exactly twice in this way. Thus there are \(n(n - 1)/2\) edges in \(K_n\).

In general, if we know the degrees of all the vertices in a graph we can find the number of edges. The sum of the degrees of all vertices will always be twice the number of edges,
since each edge adds to the degree of two vertices. Notice this means that the sum of the degrees of all vertices in any graph must be even!

**Example:** At a recent math seminar, 9 mathematicians greeted each other by shaking hands. Is it possible that each mathematician shook hands with exactly 7 people at the seminar?

**Solution:** It seems like this should be possible - each mathematician chooses one person to not shake hands with. But it is impossible. We are asking whether a graph with 9 vertices can have each vertex have degree 7. If such a graph existed, the sum of the degrees of the vertices would be $9 \cdot 7 = 63$. This would be twice the number of edges (handshakes) so this says that the graph would have 31.5 edges. That is impossible - edges only come in whole numbers. Thus at least one (in fact an odd number) of the mathematicians must have shaken hands with an even number of people at the seminar.

One final definition: we say a graph is *bipartite* if the vertices can be divided into two sets, $A$ and $B$, with no two vertices in $A$ adjacent and no two vertices in $B$ adjacent. Of course the vertices in $A$ can be adjacent to some or all of the vertices in $B$. If all vertices in $A$ and adjacent to all the vertices in $B$, then the graph is a complete bipartite graph, and gets a special name: $K_{m,n}$, where $|A| = m$ and $|B| = n$. The graph in the houses and utilities puzzle is $K_{3,3}$.

### Named Graphs
Some graphs are used more than others, and get special names.

- $K_n$: the complete graph on $n$ vertices.
- $K_{m,n}$: the complete bipartite graph with sets of $m$ and $n$ vertices.
- $C_n$: the cycle graph on $n$ vertices - just one big loop.
- $P_n$: the path graph on $n$ vertices - just one long path.

Here are some typical examples:
Graph Theory Definitions

- **Graph**: A collection of vertices, some of which are connected by edges.
- **Adjacent**: Two vertices are adjacent if they are connected by an edge.
- **Bipartite graph**: A graph for which it is possible to divide the vertices into two disjoint sets such that there are no edges between any two vertices in the same set.
- **Complete bipartite graph**: A bipartite graph for which every vertex in the first set is adjacent to every vertex in the second set.
- **Complete graph**: A simple graph with edges connecting every pair of vertices.
- **Connected**: A graph is connected if there is a path from any vertex to any other vertex.
- **Degree of a vertex**: The number of edges connected to a vertex is called the degree of the vertex.
- **Euler path**: A path which uses each edge exactly once
- **Euler circuit**: An Euler path which starts and stops at the same vertex
- **Planar**: A graph is planar if it is possible to draw it without any edges crossing.
- **Simple**: A graph is simple if it does not contain any double edges or single edge loops (that is an edge from a vertex to itself).
- **Tree**: A graph with no cycles.
- **Vertex coloring**: An assignment of colors to each of the vertices of a graph.
- **Proper vertex coloring**: A vertex coloring is proper if adjacent vertices are always colored differently.
- **Chromatic number**: The minimum number of colors required in a proper vertex coloring of the graph.
2.2 Euler Paths and Circuits

If we start at a vertex and trace along edges to get to other vertices, we create a path on the graph. If the path travels along every edge exactly once, then the path is called an Euler path (or Eulerian path). If in addition, the starting and ending vertices are the same (so you trace along every edge exactly once and end up where you started) then the path is called an Euler circuit. Of course if a graph is not connected, there is no hope of finding such a path or circuit. For the rest of this section, assume all graphs are connected.

The bridges of Königsberg problem is really a question about the existence of Euler paths. There will be a route that crosses every bridge exactly once if and only if the graph below has an Euler path:

This graph is small enough that we could actually check every possible path and in doing so convince ourselves that there is no Euler path (let alone an Euler circuit). On small graphs which do have an Euler path, it is usually not difficult to find one. Our goal is to find a quick way to check whether a graph has an Euler path or circuit, even if the graph is quite large.

If you have not already, be sure to try the worksheet on Euler paths. It is quite instructive to build graphs which do and do not have Euler paths. One way to guarantee that a graph does not have an Euler circuit is to include a “spike” - a vertex of degree 1.

The vertex $a$ has degree 1, and if you try to make an Euler circuit, you see that you will get stuck at the vertex. It is a dead end. That is, unless you start there. But then there is no way to return, so there is no hope of finding an Euler circuit. There is however an Euler path - you can start at the vertex $a$, then loop around the triangle. You will end at the vertex of degree 3.

You run into a similar problem whenever you have a vertex of odd degree. If you start at such a vertex, you will not be able to end there (after visiting every edge exactly once). After using one edge to leave the starting vertex, you will be left with an even number of edges emanating from the vertex. Half of these could be used for returning to the vertex, the other half for leaving. So you return, then leave. Return, then leave. The only way to use up all the edges is to use the last one by leaving the vertex. On the other hand, if you have a vertex with odd degree that you do not start a path at, then you will eventually get stuck at that vertex. The path will use pairs of edges connected the the vertex to arrive and
leave again. Eventually all but one of these edges will be used up, leaving only an edge to arrive by, and none to leave again.

What all this says is that if a graph has an Euler path and two vertices with odd degree, then the Euler path must start at one of the odd degree vertices and end at the other. In such a situation, every other vertex must have an even degree - since we need an equal number of edges to get to those vertices as to leave them. How could we have an Euler circuit? The graph could not have an odd degree vertex - an Euler path would have to start there or end there, but not both. Thus for a graph to have an Euler circuit, all vertices must have even degree.

The converse is also true: if all the vertices of a graph have even degree, then the graph has an Euler circuit, and if there are exactly two vertices with odd degree, the graph has an Euler path. To prove this is a little tricky, but the basic idea is that you will never get stuck because there is an “outbound” edge for every “inbound” edge at every vertex. If you try to make an Euler path and miss some edges, you will always be able to “splice in” a circuit using the edges you previously missed.

### Euler Paths and Circuits

- A graph has an Euler circuit if and only if the degree of every vertex is even.
- A graph has an Euler path if and only if there are at most two vertices with odd degree.

Since the bridges of Königsberg graph has all four vertices with odd degree, there is no Euler path through the graph.

### Hamilton paths

Suppose you wanted to tour Königsberg in such a way where you visit each land mass (the two islands and both banks) exactly once. This can be done. In graph theory terms we are asking whether there is a path which visits every vertex exactly once. Such a path is called a Hamilton path (or Hamiltonian path). It appears that finding Hamilton paths would be easier - graphs often have more edges than vertices, so there are fewer requirements to be met. However, nobody knows whether this is true. There is no known simple test for whether a graph has a Hamilton path. For small graphs this is not a problem, but as the size of the graph grows, it gets harder and harder to check wither there is a Hamilton path. In fact, this is an example of a question which as far as we know is too difficult for computers to solve - it is an example of a problem which is NP-complete.

### 2.3 Planar Graphs

We leave the question of finding paths for a completely different type of question: when is it possible to draw a graph so that none of the edges cross, except for at vertices? If this is
possible, we say the graph is planar (since you can draw it on the plane).

Notice that the definition of planar includes the phrase “it is possible to.” This means that even if a graph does not look like it is planar, it still might be. Perhaps you can redraw it in a way in which no edges cross. For example, this is a planar graph:

That is because we can redraw it like this:

The graphs are the same graph, so if one is planar, the other must be too. The original drawing of the graph was not a planar representation of the graph.

Now when a planar graph is drawn without edges crossing, the graph divides the plane into regions. We will call each region a face. The graph above has 3 faces (yes, we do include the “outside” region as a face). The number of faces does not change no matter how you draw the graph (as long as you do so without the edges crossing), so it makes sense to ascribe the number of faces as a property of the planar graph.

A warning: you can only count faces when the graph is drawn in a planar way. For example, consider these two representations of the same graph:

If you try to count faces using the graph on the left, you might say there are 5 faces (including the outside). But drawing the graph with a planar representation shows that in fact there are only 4 faces.

There is a connection between the number of vertices \(V\), the number of edges \(E\) and the number of faces \(F\) in any connected planar graph. This relationship is called Euler’s Formula.

\[
V - E + F = 2
\]
Why is Euler’s formula true? One way to convince yourself of its validity is to draw a planar graph step by step. Start with the graph $P_2$:

![Graph P2](image)

Any connected graph (besides just a single isolated vertex) must contain this subgraph. Now build up to your graph by adding edges and vertices. Each step will consist of either adding a new vertex connected by a new edge to part of your graph (so creating a new “spike”) or by connecting two vertices already in the graph with a new edge (completing a circuit).

![Graph Construction](image)

What do these “moves” do? When adding the spike, the number of edges increases by 1, the number of vertices increases by one, and the number of faces remains the same. But this means that $V - E + F$ does not change. Completing a circuit adds one edge, adds one face, and keeps the number of vertices the same. So again, $V - E + F$ does not change.

Since we can build any graph using a combination of these two moves, and doing so never changes the quantity $V - E + F$, that quantity will be the same for all graphs. But notice that our starting graph $P_2$ has $V = 2$, $E = 1$ and $F = 1$, so $V - E + F = 2$.

**Non-planar graphs**

Not all graphs are planar. If there are too many edges and too few vertices, then some of the edges will need to intersect. The first time this happens is in $K_5$.

![Graph K5](image)

If you try to redraw this without edges crossing, you quickly get into trouble. There seems to be one edge too many. In fact, we can prove that no matter how you draw it, $K_5$ will always have edges crossing.

**Theorem.** $K_5$ is not planar.
Proof. The proof is by contradiction. So assume that $K_5$ were planar. Then the graph would satisfy Euler’s formula for planar graphs. $K_5$ has 5 vertices and 10 edges, so we get

$$5 - 10 + F = 2$$

which says that if the graph were drawn without any edges crossing, there would be $F = 7$ faces.

Now consider how many edges surround each face. Each face must be surrounded by at least 3 edges (since $K_5$ is simple - it contains no double edges or loops). Let $B$ be the total number of boundaries around all the faces in the graph. Thus we have that $B \geq 3F$. But also $B = 2E$, since each edge is used as a boundary exactly twice. Putting this together we get

$$3F \leq 2E$$

But this is impossible, since we have already determined that $F = 7$ and $E = 10$, and $21 \not\leq 20$. This is a contradiction so in fact $K_5$ is not planar. QED

The other simplest graph which is not planar is $K_{3,3}$.

Proving that $K_{3,3}$ is not planar answers the houses and utilities puzzle - it is not possible to connect each of three houses to each of three utilities without the lines crossing.

**Theorem.** $K_{3,3}$ is not planar.

**Proof.** Again, we proceed by contradiction. Suppose $K_{3,3}$ were planar. Then by Euler’s formula there will be 5 faces, since $V = 6$, $E = 9$, and $6 - 9 + F = 2$.

How many boundaries surround these 5 faces? Let $B$ be this number. Since each edge is used as a boundary twice, we have $B = 2E$. Also, $B \geq 4F$ since each face is surrounded by 4 or more boundaries. We know this is true because the graph is simple (so there are no faces surrounded by 1 or 2 boundaries) and because $K_{3,3}$ is bipartite, so does not contain any 3-edge cycles. Thus

$$4F \leq 2E$$

But this would say that $20 \leq 18$, which is clearly false. Thus $K_{3,3}$ is not planar. QED

Note the similarities and differences in these proofs. Both are proofs by contradiction, and both start with using Euler’s formula to derive the (supposed) number of faces in the graph. Then we find a relationship between the number of faces and the number of edges based on how many edges surround each face. This is the part that changes - in the proof for $K_5$, we got $3F \leq 2E$ and for $K_{3,3}$ we go $4F \leq 2E$. The coefficient of $F$ is the key. It is the smallest number of edges which could surround any face. If some number of edges surround a face, then these edges form a circuit. So that number is the size of the smallest circuit in the graph.

In general, if we let $g$ be the size of the smallest cycle in a graph ($g$ stands for girth, which is the technical term for this) then for any planar graph we have $gF \leq 2E$. When this disagrees with Euler’s formula, we know for sure that the graph cannot be planar.
2.4 Coloring

We conclude with perhaps the most famous graph theory problem - how to color maps.

**Question:** Given any map of countries, states, counties, etc. how many colors are needed to color each region on the map so that neighboring regions are colored differently?

Actual map makers usually use around seven colors - for one thing, they require watery regions to be a specific color, and with a lot of colors it is easier to find a permissible coloring. But we want to know whether there is a smaller palette that will work for any map.

How is this related to graph theory? Well, if we place a vertex in the center of each region (say in the capital of each state) and then connect two vertices if their states share a border, we get a graph. The coloring regions on the map corresponds to coloring the vertices of the graph. Since neighboring regions cannot be colored the same, our graph cannot have vertices colored the same when those vertices are adjacent (connected by an edge).

In general, given any graph, a coloring of the vertices is called (not surprisingly) a *vertex coloring*. If the vertex coloring has the property that adjacent vertices are colored differently, then the coloring is called *proper*. Every graph has a proper vertex coloring - for example, you could color every vertex with a different color. But often you can do better. The smallest number of colors needed to get a proper vertex coloring is called the *chromatic number* of the graph.

**Example:** Find the chromatic number of the graphs below.

![Graphs](image)

**Solution:** The graph on the left is $K_6$. The only way to properly color the graph is to give every vertex a different color (since every vertex is adjacent to every other vertex). Thus the chromatic number is 6.

The middle graph can be properly colored with just 3 colors (Red, Blue, and Green). For example:

![Graph Coloring](image)

There is no way to color it with just two colors, since there are three vertices mutually adjacent (i.e., a triangle). Thus the chromatic number is 3.

The graph on the right is just $K_{2,3}$. As with all bipartite graphs, this graph has chromatic number 2 - color the vertices in the top row red and the vertices on the bottom row blue.
It appears that graphs can have any chromatic number. It should not come as a surprise that $K_n$ has chromatic number $n$. So how could there possibly be an answer to the original map coloring question? If the chromatic number of graph can be arbitrarily large, then it seems like there would be no upper bound to the number of colors needed for any map. But there is.

The key observation is that while it is true that for any number $n$, there is a graph with chromatic number $n$, only some graphs arrive as representations of maps. If you convert a map to a graph, the edges between vertices correspond to borders between the countries. So you should be able to connect vertices in such a way where the edges do not cross. In other words, the graphs representing maps are all planar!

So the question is, what is the largest chromatic number of any planar graph? The answer is one of the best known theorems of mathematics:

**Theorem.** The Four Color Theorem If $G$ is a planar graph, then the chromatic number of $G$ is less than or equal to 4. Thus any map can be colored with 4 or fewer colors.

We will not prove this theorem. Really. Even though the theorem is easy to state and understand, the proof is not. In fact, there is currently no “easy” known proof of the theorem. The current best proof still requires powerful computers to check an unavoidable set of 633 reducible configurations. The idea is that every graph must contain one of these reducible configurations (this fact also needs to be checked by a computer) and that reducible configurations can in fact be colored in 4 or fewer colors.

The chromatic number of a graph tells us about coloring vertices - but we could also ask about coloring edges. What if we colored every edge of a graph either red or blue. Can we do so without, say, creating a triangle of same like colored edges (i.e., an all red or all blue triangle - we say the triangle is mono-chromatic)? Certainly for some graphs the answer is yes - try doing so for $K_4$. What about $K_5$? $K_6$? How far can we go?

The answer the above problem is known - I encourage you to try to solve it. We could extend the question though - what if we had three colors? What if we were trying to avoid other graphs. The surprising fact is that very little is known about these questions. For example, we know that you need to go up to $K_{17}$ in order to force a mono-chromatic triangle using three colors, but nobody knows how big you need to go with more colors. Similarly, we know that using two colors $K_{18}$ is the smallest graph that forces a mono-chromatic copy of $K_4$, but the best we have to force a mono-chromatic $K_5$ is a range - somewhere from $K_{43}$ to $K_{49}$.

**Exercises**

1. If 10 people each shake hands with each other, how many handshakes took place? What does this question have to do with graph theory?

2. Among a group of 5 people, is it possible for everyone to be friends with exactly 2 of the people in the group? What about 3 of the people in the group?

3. You and your friends want to tour the southwest by car. You will visit the nine states below, with the following rather odd rule: you must cross each border between
neighboring states exactly once (so, for example, you must cross the Colorado-Utah border exactly once). Can you do it? If so, does it matter where you start your road trip? What fact about graph theory solves this problem?

4. Which (if any) of the graphs below are the same? Which are different? Explain.

5. Which of the graphs in the previous question contain Euler paths or circuits? Which of the graphs are planar?

6. Is it possible for two different graphs to have the same number of vertices and the same number of edges? What if the degrees of the vertices in the two graphs are the same (so both graphs have vertices with degrees 1, 2, 2, 3, and 4, for example)? Draw two such graphs or explain why not.

7. Which of the following graphs contain an Euler path? Which contain an Euler circuit?
   (a) $K_4$
   (b) $K_5$.
   (c) $K_{5,7}$
   (d) $K_{2,7}$
   (e) $C_7$
   (f) $P_7$

8. For which $n$ does the graph $K_n$ contain an Euler circuit? Explain.

9. For which $m$ and $n$ does the graph $K_{m,n}$ contain an Euler path? An Euler circuit? Explain.

10. Which of the graphs below are bipartite?
11. For which $n$ is the graph $C_n$ bipartite?

12. Draw a graph which has an Euler circuit but is not planar.

13. Draw a graph which does not have an Euler path and is also not planar.

14. Is it possible for a planar graph to have 6 vertices, 10 edges and 5 faces? Explain.

15. If a graph has 10 vertices and 10 edges and contains an Euler circuit, must it be planar? How many faces would it have?

16. The graph $G$ has 6 vertices with degrees 2, 2, 3, 4, 4, 5. How many edges does $G$ have? Could $G$ be planar? If so, how many faces would it have.

17. What is the smallest number of colors you need to properly color the vertices of $K_{4,5}$? That is, find the chromatic number of the graph.

18. Draw a graph with chromatic number 6 (i.e., which requires 6 colors to properly color the vertices). Could your graph be planar? Explain.

19. Find the chromatic number of each of the following graphs.

![Graphs](image.png)

20. Suppose $G$ is a simple graph with $n$ vertices, each having degree 5.
   
   (a) For which values of $n$ does this make sense?
   
   (b) For which values of $n$ does the graph have an Euler path?
   
   (c) What is the smallest value of $n$ for which the graph might be planar? (tricky)
Chapter 3

How to Count

One of the first things you learn in mathematics is how to count. Now we are going to
learn that all over again, and you will find that counting is a lot harder than you remember.
The problem is that we want to count large collections of things quickly and precisely. For example:

- How many different lotto tickets are possible?
- In a group of 10 people, if everyone shakes hands with everyone else exactly once, how
  many handshakes took place?\footnote{We know the answer to this thanks to graph theory - we are asking how many edges there are in $K_{10}$. What if your group had a secrete 3-person handshake? If every group of three people participated in such a handshake, how many handshakes take place? Equivalently, how many triangles are there in $K_{10}$?}
- How many ways can you distribute 10 girl scout cookies to 7 boy scouts?
- How many 10 digit numbers contain exactly 4 prime digits?
- How many functions $f : \{1, 2, 3, 4\} \to \{1, 2, 3\}$ are onto?
- How many subsets of $\{1, 2, 3, \ldots, 10\}$ have cardinality 7?

Before tackling these difficult questions, let’s look at the basics of counting.

3.1 Additive and Multiplicative Principles

Consider this rather simple counting problem: at Red Dogs and Donuts, there are 14 varieties
of donuts, and 16 types of hot-dogs. If you want either a donut or a dog, how many options
do you have? This is an easy question - you just add 14 and 16. Will that always work?
What is important here?

Additive Principle
The additive principle states that if an event $A$ can occur in $m$ ways, and event $B$ can
occur in $n$ disjoint ways, then the event “$A$ or $B$” can occur in $m + n$ ways.

It is important that the events be disjoint. For example, a standard deck of 52 cards
contains 26 red cards and 12 face cards. However, the number of ways to select a card which
is either red or a face card is not $26 + 12 = 38$. This is because there are 6 cards which are
both red and face cards.
The additive principle works with more than two events. Say you would also consider eating one of 15 waffles? How many choices do you have now? You would have $14 + 16 + 15 = 45$ options.

**Example:** How many two letter “words” start with either A or B? How many start with one of the 5 vowels? (A word is just a strings of letters - they don’t have to be English words, or even pronounceable).

**Solution:** First, how many two letter words start with A? We just need to select the second letter, which can be accomplished in 26 ways. So there are 26 words starting with A. There are also 26 words that start with B. So to select a word which starts with either A or B, we can pick the word from the first 26 or the second 26, for a total of 52 words. The additive principle is at work here.

Now what about all the two letter words starting with a vowel? Well there are 26 starting with A, another 26 starting with E, and so on. We will have 5 groups of 26. So we add 26 to itself 5 times. Of course it would be easier to just multiply $5 \cdot 26$ - we are really using the additive principle again, just using multiplication as a shortcut.

**Example:** Suppose you are going for some FroYo - you can pick one of 6 yogurt choices, and one of 4 toppings. How many choices do you have?

**Solution:** Break your choices up into disjoint events: $A$ are the choices with the first topping, $B$ the choices featuring the second topping, and so on. So we have events. Each can occur in 6 ways (one for each yogurt flavor). The events are disjoint, so the total number of choices is $6 + 6 + 6 + 6$.

Note that in both of the previous examples, when using the additive principle on a bunch of sets all the same size, it is quicker to multiply. This really is the same - not just because $6 + 6 + 6 + 6 = 4 \cdot 6$. We can first select the topping in 4 ways (that is we first select which of the disjoint events we will take). For each of those first 4 choices, we now have 6 choices of yogurt. We have:

### Multiplicative Principle

The *multiplicative principle* states that if event $A$ can occur in $m$ ways, and each possibility for $A$ allows for exactly $n$ ways for event $B$, then the event “$A$ and $B$” can occur in $m \cdot n$ ways.

The product rule generalizes to more than two events as well.

**Example:** How many license plates can you make out of three letters followed by three numerical digits?

**Solution:** Here we have six events: the first letter, the second letter, the third letter, the first digit, the second digit and the third digit. The first three events
can each happen in 26 ways, the last three can each happen in 10 ways. So the total number of license plates will be $26 \cdot 26 \cdot 26 \cdot 10 \cdot 10 \cdot 10$, using the multiplicative principle.

Does this make sense? Think about how we would pick a license plate - how many choices we would have. First, we need to pick the first letter. There are 26 choices. Now for each of those, there are 26 choices for the second letter. So 26 second letters with first letter A, 26 second letters with first letter B, and so on. So we add 26 to itself 26 times. Or quicker: there are $26 \cdot 26$ choices for the first two letters.

Now for each choice of the first two letters, we have 26 choices for the third letter. That is, 26 third letters for the first two letters AA, 26 choices for the third letter after starting AB, and so on. There are $26 \cdot 26$ of these 26 third letter choices, for a total of $(26 \cdot 26) \cdot 26$ choices for the first three letters. And for each of these $26 \cdot 26 \cdot 26$ choices of letters, we have a bunch of choices for the remaining digits.

In fact, there are going to be exactly 1000 choices for the numbers. We can see this because there are 1000 three-digit numbers (000 through 999). This is 10 choices for the first digit, 10 for the second, and 10 for the third. The multiplicative principle says we multiply: $10 \cdot 10 \cdot 10 = 1000$.

So there were $26^3$ choices for the three letters, and $10^3$ choices for the numbers, so to we have a total of $26^3 \cdot 10^3$ choices of license plates.

Careful: “and” doesn’t mean “times.” For example, how many playing cards are both red and a face card? Not $26 \cdot 12!$ The answer is 6, and we needed to know something about cards to answer that question.

Another caution: how many ways can you select two cards, so that the first one is a red card and the second one is a face card? This looks more like the multiplicative principle (you are counting two separate events) but the answer is not $26 \cdot 12$ here either. The problem is that while there are 26 ways for the first card to be selected, it is not the case that for each of those there are 12 ways to select the second card. If the first card was both red and a face card then there would be only 11 choices for the second card. The moral of this story? The multiplicative principle only works if the events are independent.²

### 3.2 Counting with sets

Do you believe the additive and multiplicative principles? How would you convince someone they are correct? This is surprisingly difficult. They seem so simple, so obvious. But why do they work?

To make things clearer, and more mathematically rigorous, we will use sets. Do not skip this section! It might seem like we are just trying to give a proof of these principles, but we

²To solve this problem, you could break into two cases - first count how many ways there are to select the two cards when the first card is a red non-face card and second count how many ways when the first card is a red face card. Doing so makes the events in each separate case independent, so the multiplicative principle can be applied.
are doing a lot more. If we understand the additive and multiplicative principles rigorously, we will be better at applying them, and knowing when and when not to apply them at all.

We will look at the previous problems in a slightly different way. Instead of thinking about event \( A \) and event \( B \), we want to think of a set \( A \) and a set \( B \). The sets will contain all the different ways the event can happen. (It will be helpful to be able to switch back and forth between these two models when checking that we have counted correctly.) Here’s what we mean:

**Example:** Suppose you own 9 shirts and 5 pairs of pants.
1. How many outfits can you make?
2. If today is half-naked-day, and you will wear only a shirt or only a pair of pants, how many choices do you have?

**Solution:** By now you should agree that the answer to the first question is \( 9 \cdot 5 = 45 \) and the answer to the second question is \( 9 + 5 = 14 \). These are the multiplicative and additive principles at work. There are two events: picking a shirt and picking a pair of pants. The first event can happen in 9 ways and the second event can happen in 5 ways. To get both a shirt and a pair of pants, you multiply. To get just one article of clothing, you add.

Now let’s look at this using sets. There are two sets, call them \( S \) and \( P \). The set \( S \) contains all 9 shirts so \( |S| = 9 \) while \( |P| = 5 \), since there are 5 elements in the set \( P \) (namely your 5 pairs of pants). What are we asking in terms of these sets? Well in question 2, we really want \( |S \cup P| \); how many elements are there in the union of shirts and pants. Clearly this is just \( |S| + |P| \) (since there is no overlap — in other words, since \( |S \cap P| = 0 \)). Question 1 is a little more complicated. Your first guess might be to find \( |S \cap P| \) - but this is not right (there is nothing in the intersection). We are not asking for how many clothing items are both a shirt and a pair of pants. Instead, we want one of each. We could think of this as asking how many pairs \((x, y)\) there are, where \( x \) is a shirt and \( y \) is a pair of pants. As we will see, to find this number, we would take \( |S| \cdot |P| \).

From this example we can see right away how to rephrase our additive principle in terms of sets:

**Additive Principle (with sets)**
Given two sets \( A \) and \( B \), if \( A \cap B = \emptyset \) (that is, if there is no element in common to both \( A \) and \( B \)), then
\[
|A \cup B| = |A| + |B|
\]

This hardly needs a proof - to find \( A \cup B \) you take everything in \( A \) and throw in everything in \( B \). Since there is no element in both sets already, you will have \( |A| \) things and add \( |B| \).
new things to it. This is what adding does! Of course, we can easily extend this to any number of (disjoint) sets.

From the example above, we see that in order to investigate the multiplicative principle carefully, we need to consider ordered pairs. We should define this carefully.

**Definition 3.** Given sets $A$ and $B$, we can form the set $A \times B = \{(x, y) : x \in A, y \in B\}$ to be the set of all ordered pairs $(x, y)$ where $x$ is an element of $A$ and $y$ is an element of $B$. We call $A \times B$ the Cartesian product of $A$ and $B$.

The question is, what is $|A \times B|$ - the cardinality of the Cartesian product? To figure this out, let’s write out $A \times B$.

Let $A = \{a_1, a_2, a_3, \ldots, a_m\}$ and $B = \{b_1, b_2, b_3, \ldots, b_n\}$ (so $|A| = m$ and $|B| = n$). The set $A \times B$ contains all pairs with the first half of the pair being $a_i$ for some $i$ and the second being $b_j$ for some $j$. In other words:

$$A \times B = \{(a_1, b_1), (a_1, b_2), (a_1, b_3), \ldots, (a_1, b_n), (a_2, b_1), (a_2, b_2), (a_2, b_3), \ldots, (a_2, b_n), (a_3, b_1), (a_3, b_2), (a_3, b_3), \ldots, (a_3, b_n), \ldots, (a_m, b_1), (a_m, b_2), (a_m, b_3), \ldots, (a_m, b_n)\}$$

Notice what we have done here: we made $m$ rows of $n$ pairs - so that is a total of $m \cdot n$ pairs.

Each row above is really $\{a_i\} \times B$ for some $a_i \in A$. That is, we fixed the $A$-element. Clearly we have

$$A \times B = \{a_1\} \times B \cup \{a_2\} \times B \cup \{a_3\} \times B \cup \cdots \cup \{a_m\} \times B$$

So $A \times B$ is really the union of $m$ disjoint sets. Each of those sets has $n$ elements in them. So the total (using the additive principle) is $n + n + n + \cdots + n = m \cdot n$.

To summarize:

---

**Multiplicative Principle (with sets)**
Given two sets $A$ and $B$, we have $|A \times B| = |A| \cdot |B|$.

Again, we can easily extend this to any number of sets.

### 3.2.1 Principle of Inclusion/Exclusion

While we are thinking about sets, let’s consider what happens to the additive principle when the sets are NOT disjoint. Suppose we know that $|A| = 10$ and $|B| = 8$ (there are 10 elements in $A$ and 8 elements in $B$). This is not enough information though. We do not know how many of the 8 elements in $B$ are also element of $A$. But if we do know that $|A \cap B| = 6$, then we can say exactly how many elements are in $A$, and of those how many are in $B$ and how many are not (6 of the 10 elements are in $B$, so 4 are in $A$ but not in $B$). We would fill in a Venn diagram as follows:
This says there are 6 elements in $A \cap B$, 4 elements in $A \setminus B$ and 2 elements in $B \setminus A$. Now these three sets are disjoint, so we can use the additive principle to find the number of elements in $A \cup B$. It is $6 + 4 + 2 = 12$.

This will always work, but drawing a Venn diagram is more than we need to do. In fact, it would be nice to relate this problem to the case where $A$ and $B$ are disjoint. Is there one rule we can make that works in either case?

Here is another way to get the answer to the problem above. We start by just adding $|A| + |B|$. This is $10 + 8 = 18$, which would be the answer if $|A \cap B| = 0$. And we see that we are off by exactly 6, which just so happens to be $|A \cap B|$. So perhaps we guess

$$|A \cap B| = |A| + |B| - |A \cap B|.$$ 

This works for this one example. Will it always work? Think about what we are doing here. We want to know how many things are either in $A$ or $B$ (or both). We can throw in everything in $A$, and everything in $B$. This would give $|A| + |B|$ many elements. But of course when you actually take the union, you do not repeat elements that are in both. So far we have counted every element in $A \cap B$ exactly twice - once when we put in the elements from $A$ and once when we included the elements from $B$. We correct by subtracting out the number of elements we have counted twice. So we added them in twice, subtracted once, leaving them counted only one time.

In other words, we have:

<table>
<thead>
<tr>
<th>Cardinality of a union (2 sets)</th>
</tr>
</thead>
<tbody>
<tr>
<td>For any finite sets $A$ and $B$,</td>
</tr>
<tr>
<td>$</td>
</tr>
</tbody>
</table>

We can do something similar with three sets. First here is an example of how you could use Venn diagrams.

**Example:** An examination in three subjects, Algebra, Biology, and Chemistry, was taken by 41 students. The following table shows how many students failed in each single subject and in their various combinations.

<table>
<thead>
<tr>
<th>Subject:</th>
<th>A</th>
<th>B</th>
<th>C</th>
<th>AB</th>
<th>AC</th>
<th>BC</th>
<th>ABC</th>
</tr>
</thead>
<tbody>
<tr>
<td>Failed:</td>
<td>12</td>
<td>5</td>
<td>8</td>
<td>2</td>
<td>6</td>
<td>3</td>
<td>1</td>
</tr>
</tbody>
</table>

How many students failed at least one subject?
Solution:
The answer is not 37, even though the sum of the numbers above is 37. The reason is that while 12 students failed algebra, 2 of those students also failed biology and 1 failed chemistry as well. In fact, that 1 student who failed all three subjects is counted a total of 7 times in the total 37. To clarify things, let us think of the students who failed algebra as the elements of the set $A$, and similarly for sets $B$ and $C$. The one student who failed all three subjects is the lone element of the set $A \cap B \cap C$. Thus in Venn diagrams:

Now let’s fill in the other intersections. We know $A \cap B$ contains 2 elements, but one element has already been counted. So we should put a one in the region where $A$ and $B$ intersect (but $C$ does not). Similarly, we calculate the cardinality of $(A \cap C) \cap \overline{B}$, and $(B \cap C) \cap \overline{A}$:

Next we determine the numbers which should go in the remaining regions, including outside of all three circles. This last number is the number of students who did not fail any subject – the number we were asked to find:

We found that 5 so go in the $A$ only region because the entire circle for $A$ needed to have a total of 12, and 7 were already accounted for. Thus the number of students who passed all three classes is 26. The number who failed at least one class is 15.
Note that we can also answer other questions. For example, now many students failed just chemistry? None. How many passed biology but failed both algebra and chemistry? 5.

Can we solve the problem in an algebraic way? Note that while the additive principle generalizes to any number of sets, when we add a third set here, we must be careful. With two sets, we needed to know the cardinalities of $A$, $B$, and $A \cap B$ in order to find the cardinality of $A \cup B$. With three sets we need more information - there are more ways the sets can combine. Not surprisingly then, the formula for cardinality of the union of three non-disjoint sets is more complicated:

**Cardinality of a union (3 sets)**
For any finite sets $A$, $B$, and $C$,

$$|A \cup B \cup C| = |A| + |B| + |C| - |A \cap B| - |A \cap C| - |B \cap C| + |A \cap B \cap C|$$

To determine how many elements are in at least one of $A$, $B$, or $C$ we add up all the elements in each of those sets. However, when we do that, any element in both $A$ and $B$ is counted twice. Also each element in both $A$ and $C$ is counted twice, as are elements in $B$ and $C$. So we take each of those out of our sum once. But now what about the elements which are in $A \cap B \cap C$ (in all three sets). We added them in three times, but also removed them three times. So they have not yet been counted. Thus we add those elements back in at the end.

Returning to our example above, we have $|A| = 12$, $|B| = 5$, $|C| = 8$. We also have $|A \cap B| = 2$, $|A \cap C| = 6$, $|B \cap C| = 3$ and $|A \cap B \cap C| = 1$. So:

$$|A \cup B \cup C| = 12 + 5 + 8 - 2 - 6 - 3 + 1 = 15$$

This is what we got when we solved the problem using Venn diagrams.

This process of adding in, then taking out, then adding back in, and so on is called the **Principle of Inclusion/Exclusion**, or simply PIE. We will return to this counting technique later to solve for more complicated problems (involving many more sets).

### 3.3 Four things to count

Here are some apparently different discrete objects we can count: subsets, bit strings, lattice paths, and binary coefficients. We will give an example of each type of counting problem (and say what these things even are). As we will see, these counting problems are surprisingly similar.

**Subsets**

Subsets should be familiar - if not, read over the notes on set theory again. Our example problem is this: Let $A = \{1, 2, 3, 4, 5\}$. How many subsets of $A$ contain exactly 3 elements?
First, a simpler question. How many subsets of $A$ are there total? In other words, what is $|\mathcal{P}(A)|$ (the cardinality of the power set of $A$)? Think about how we would build a subset - we need to decide, for each of the elements of $A$, whether or not to include the element in our subset. So we need to decide “yes” or “no” for the element 1. And for each choice we make, we need to decide “yes” or “no” for the element 2. And so on. For each of the 5 elements, we have 2 choices. Therefore the number of subsets is simply $2^5$ (by the multiplicative principle).

Of those 32 subsets, how many have 3 elements? This is a tricky question. Note that we cannot just use the multiplicative principle. Maybe we want to say we have 2 choices for the first element, 2 choices for the second, 2 choices for the third, and then only 1 choice for the other two. But what if we said “no” to one of the first three elements? Then we would have two choices for the 4th element - this is not working.

Another (bad) idea: we need to pick three elements to be in our subset. There are 5 elements to choose from. So there are 5 choices for the first element, and for each of those 4 choices for the second, and then 3 for the third (last) element. The multiplicative principle would say then that there are a total of $5 \cdot 4 \cdot 3 = 60$ ways to select the 3 element subset. But this is wrong! One of the outcomes we would get from these choices would be the set $\{3, 2, 5\}$ - choosing the element 3 first, then the element 2, then the element 5. Another outcome would be $\{5, 2, 3\}$ - choose the element 5 first, then the element 2 then the element 3. But these are the same set! We can correct this by dividing the supposed 60 outcomes by the number of different outcomes which count as the same for each three elements - there happen to be 6 of these. So we expect there to be 10 3-element subsets of $A$.

Is this right? Well, we could list out all 10 of them, being very systematic in doing so to make sure we don’t miss any or list any twice. Or we could try to count how many subsets of $A$ don’t have 3 elements in them. How many have no elements? Just 1 (the empty set). How many have 5? Again, just one. These are the cases in which we say “no” to all elements, or “yes” to all elements. Okay, what about the subsets which contain a single element? There are 5 of these - we need to say “yes” to exactly one element, and there are 5 to choose from. Now this is also the number of subsets containing 4 elements - those are the ones for which we must say “no” to exactly one element.

So far we have counted 12 of the 32 subsets. We have not yet counted the subsets with cardinality 2 and with cardinality 3. There are a total of 20 subsets left to split up between these two groups. But the numbers must be the same! If we say “yes” to exactly two elements, that can be accomplished in exactly the same number of ways as the number of ways we can say “no” to exactly two elements. So the number of 2-element subsets is equal to the number of 3-element subsets. Together there are 20 of these subsets, so 10 each.

<table>
<thead>
<tr>
<th>Number of elements:</th>
<th>0</th>
<th>1</th>
<th>2</th>
<th>3</th>
<th>4</th>
<th>5</th>
</tr>
</thead>
<tbody>
<tr>
<td>Number of subsets:</td>
<td>1</td>
<td>5</td>
<td>10</td>
<td>10</td>
<td>5</td>
<td>1</td>
</tr>
</tbody>
</table>

**Bit Strings**

“Bit” is short for “binary digit,” so a bit string is a string of binary digits. The binary digits are simply the numbers 0 and 1. So all of the following are bit strings:

1001 0 1111 1010101010
The number of bits (0’s or 1’s) in the string is the \textit{length} of the string. So the strings above have lengths 4, 1, 4, and 10 respectively. We also can ask how many of the bits are 1’s. The number of 1’s in a bit string is the \textit{weight} of the string. So the weights of the above strings are 2, 0, 4, and 5 respectively.

\begin{itemize}
  \item A \(n\)-bit string is a bit string of length \(n\). That is, it is a string containing \(n\) symbols, each of which is a bit - a 0 or a 1.
  \item The \textit{weight} of a bit string is the number of 1’s in it.
  \item \(B^n\) is the set of all \(n\)-bit strings.
  \item \(B^n_k\) is the set of all \(n\)-bit strings of weight \(k\).
\end{itemize}

For example the elements of the set \(B^3_2\) are the bit strings 011, 101, and 110. Those are the only strings containing three bits exactly two of which are 1’s.

The counting questions: How many bit strings have length 5? How many of those have weight 3? In other words, we are asking for the cardinalities \(|B^5|\) and \(|B^3_3|\).

To find the number of 5-bit strings is easy. We have 5 bits, and each can either be a 0 or a 1. So there are 2 choices for the first bit, 2 choices for the second, and so on. By the multiplicative principle, there are \(2 \cdot 2 \cdot 2 \cdot 2 \cdot 2 = 2^5 = 32\) such strings.

Finding the number of 5-bit strings of weight 3 is harder. Think about how such a string could start. The first bit must be either a 0 or a 1. In the first case (the string starts with a 0) we must then decide on four more bits. To have a total of three 1’s, of those four remaining bits there must be three 1’s. In other words, we must include all 4-bit strings of weight 3. In the second case (the string starts with a 1) we still have four bits to choose, but now only two of them can be 1’s. So we must use all the 4-bit strings of weight 2. In other words:

\[
|B^5_3| = |B^4_2| + |B^4_3|
\]

This is because the strings in \(B^3_3\) all have the form 1\(B^4_2\) (that is, a 1 followed by a string from \(B^4_2\)) or 0\(B^4_3\). We have a recurrence relation! If only we knew the cardinalities of \(B^4_2\) and \(B^4_3\).

Well of course

\[
|B^4_2| = |B^3_1| + |B^3_2| \quad \text{and} \quad |B^4_3| = |B^3_2| + |B^3_3|
\]

We can keep going down, but this should be good enough - \(B^3_1\) and \(B^3_2\) both contain 3 bit strings - we need to pick one of the three bits to be a 1 (three ways to do that) or one of the three bits to be a 0 (three ways to do that). \(B^3_3\) contains just one string: 111. Thus \(|B^4_2| = 6\) and \(|B^4_3 + 3| = 4\), which puts \(B^5_3\) at a total of 10 strings.

But wait - 32 and 10 were the answers to the counting questions about subsets. Coincidence? Not at all. Each bit string can be thought of as a \textit{code} for a subset. For the set \(A = \{1, 2, 3, 4, 5\}\) we would use 5-bit strings - one bit for each element of \(A\). Each bit in the string is a 0 if its corresponding element of \(A\) is not in the subset, and a 1 if the element of
A is in the subset. Remember, deciding the subset amounted to a sequence of five yes/no votes for the elements of \( A \). Instead of yes, we put a 1; instead of no we put a 0.

For example, the bit string 11001 represents the subset \( \{1, 2, 5\} \) since the first, second and fifth bits are 1’s. The subset \( \{3, 5\} \) would be coded by the string 00101.

Now for a subset to contain exactly three elements, the corresponding bit string must contain exactly three 1’s. In other words, the weight must be 3. So counting the number of 3-element subsets of \( A \) is the same as counting the number 5-bit strings of weight 3.

**Lattice Paths**

The (integer) lattice is the set of all points in the Cartesian plain for which both the \( x \) and \( y \) coordinates are integers. If you like to draw graphs on graph paper, the lattice is all of the intersections of the grid lines.

A lattice path is one of the shortest possible paths connecting two points on the lattice, moving only horizontally and vertically. For example, here are three possible lattice paths from the points \((0,0)\) to \((3,2)\):

Notice that to ensure the path is the shortest possible, each move must be either to the right or up. Additionally, in this case, each path has length 5 - there are a total of 5 steps, each either right or up, which must be taken.

The counting question: how many lattice paths are there between \((0,0)\) and \((3,2)\)? We could try to draw all of these, or instead of drawing them, maybe just list which direction we travel on each of the 5 steps. One path might be RRUUR, or maybe UURRR, or perhaps RURRU (those correspond to the three paths drawn above). So how many such strings of R’s and U’s are there?

Notice that each of these strings must contain 5 symbols. Exactly 3 of them must be R’s (since our destination is 3 units to the right). This seems awfully familiar. In fact, what if we used 1’s instead of R’s and 0’s instead of U’s? Then we would just have 5-bit strings of weight 3. There are 10 of those. So there are 10 lattice paths from \((0,0)\) to \((3,2)\).

The correspondence between bit strings and lattice paths does not stop there. Here is another way to count lattice paths. Consider the lattice shown below:
Now any lattice path from (0,0) to (3,2) must pass through exactly one of A and B. The point A is 4 steps away from (0,0) and two of them are right. So the number of lattice paths to A is the same as the number of 4-bit strings of weight 2 - that is, 6. The point B is 4 steps away from (0,0) but now 3 of them are right. So the number of paths to point B is the same as the number of 4-bit strings of weight 3, namely 4. So the total number of paths to (3,2) is just 6 + 4. This is the same way we calculated the number of 5-bit strings of weight 3. The point: exact same recurrence relation exists for bit strings and for lattice paths.

### Binomial Coefficients

Binomial coefficients are the coefficients in the expanded version of a binomial such as \((x+y)^5\). What happens when we multiply such a binomial out? We will expand \((x+y)^n\) for various values of \(n\). Each of these are done by multiplying everything out (e.g., FOIL-ing) and then collecting like terms.

\[
(x+y)^1 = x + y \\
(x+y)^2 = x^2 + 2xy + y^2 \\
(x+y)^3 = x^3 + 3x^2y + 3xy^2 + y^3 \\
(x+y)^4 = x^4 + 4x^3y + 6x^2y^2 + 4xy^3 + y^4 \\
\]

In fact, there is a quicker way to expand the above binomials. For example, consider the next one, \((x+y)^5\). What we are really doing is multiplying out, \((x+y)(x+y)(x+y)(x+y)(x+y)\).

Now in the expansion, there will be only one \(x^5\) term and one \(y^5\) term. That is because to get an \(x^5\), we need to use the \(x\) term in each of the copies of the binomial \((x+y)\), and similarly for \(y^5\). What about \(x^4y\)? To get terms like this, we need to use four \(x\)'s and one \(y\). So we need exactly one of the five binomials to contribute a \(y\). There are 5 choices for this, so there are 5 ways to get \(x^4y\), so the coefficient of \(x^4y\) is 5. This is also the coefficient for \(xy^4\) for the same (but opposite) reason - there are 5 ways to pick which of the 5 binomials contribute the single \(x\). So far we have

\[
(x+y)^5 = x^5 + 5x^4y + 10x^3y^2 + 10x^2y^3 + 5xy^4 + y^5 \\
\]

We still need the coefficients of \(x^3y^2\) and \(x^2y^3\). In both cases, we need to pick exactly 3 of the 5 binomials to contribute one variable, the other two to contribute the other. Wait. This sounds familiar. We have 5 things, each can be one of two things, and we need a total of 3 of one of them. That’s just like taking 5 bits and making sure exactly 3 of them are 1’s. So the coefficient of \(x^3y^2\) (and also \(x^2y^3\)) will be exactly the same as the number of bit strings of length 5 and weight 3, which we found earlier to be 10. So we have:

\[
(x+y)^5 = x^5 + 5x^4y + 10x^3y^2 + 10x^2y^3 + 5xy^4 + y^5 \\
\]

These numbers we keep seeing over and over again - the number of subsets of a particular size, the number of bit strings of a particular weight, the number of lattice paths, and the
coefficients of these binomial products - these numbers will all be called \textit{binomial coefficients}. We even have a special symbol for them: \( \binom{n}{k} \).

\begin{center}
\textbf{Binomial Coefficients}
\end{center}

For each integer \( n \geq 0 \) and integer \( k \) with \( 0 \leq k \leq n \) we have a number

\[
\binom{n}{k}
\]

read “\( n \) choose \( k \).” We have:

- \( \binom{n}{k} = |B^n_k| \), the number of \( n \)-bit strings of weight \( k \).
- \( \binom{n}{k} \) is the number of subsets of a set of size \( n \) each with cardinality \( k \).
- \( \binom{n}{k} \) is the number of lattice paths of length \( n \) containing \( k \) steps to the right.
- \( \binom{n}{k} \) is the coefficient of \( x^k y^{n-k} \) in the expansion of \((x + y)^n\).
- \( \binom{n}{k} \) is the number of ways to select \( k \) objects from a total of \( n \) objects.

The last bullet point is usually taken as the definition of \( \binom{n}{k} \) - we have \( n \) objects, and must choose \( k \) of them, so there are \( n \) choose \( k \) ways of doing this. Each of our counting problems above can be viewed in this way:

- How many subsets of \( \{1, 2, 3, 4, 5\} \) contain exactly 3 elements? We must choose 3 of the 5 elements to be in our subset. There are \( \binom{5}{3} \) ways to do this, so there are \( \binom{5}{3} \) such subsets.

- How many bit strings have length 5 and weight 3? We must choose 3 of the 5 bits to be 1’s. There are \( \binom{5}{3} \) ways to do this, so there are \( \binom{5}{3} \) such bit strings.

- How many lattice paths are there from \((0,0)\) to \((3,2)\)? We must choose 3 of the 5 steps to be towards the right. There are \( \binom{5}{3} \) ways to do this, so there are \( \binom{5}{3} \) such paths.

- What is the coefficient of \( x^3 y^2 \) in the expansion of \((x + y)^5\)? We must choose 3 of the 5 copies of the binomial to contribute an \( x \). There are \( \binom{5}{3} \) ways to do this, so the coefficient is \( \binom{5}{3} \).

It should be clear that in each case above, we have the right answer. All we had to do is phrase the question correctly and it became obvious that \( \binom{5}{3} \) is correct. However, this does not tell us that the answer is in fact 10 in each case. We will eventually find a formula for \( \binom{n}{k} \) but for now, look back at how we arrived at the answer 10 in our counting problems above. It all came down to bit strings, and we have a recurrence relation for bit strings:

\[
|B^n_k| = |B^{n-1}_{k-1}| + |B^{n-1}_k|
\]
Remember, this is because we can start the bit string with either a 1 or a 0. In both cases, we have now have \( n - 1 \) more bits to pick. The strings starting with 1 must contain \( k - 1 \) more 1’s, while the strings starting with 0 still need \( k \) more 1’s.

Now since \( |B_k^n| = \binom{n}{k} \), we can use the same recurrence relation for binomial coefficients:

\[
\binom{n}{k} = \binom{n-1}{k-1} + \binom{n-1}{k}
\]

### 3.4 Pascal’s Triangle

Let’s arrange the binomial coefficients \( \binom{n}{k} \) into a triangle like follows:

\[
\begin{array}{ccccccc}
(0) & & & & & & \\
& 0 & & & & & \\
& & 1 & & & & \\
& & & 2 & & & \\
& (3) & & (3) & & & \\
& 0 & & 1 & & 2 & \\
& & & & & & \\
& (4) & (4) & (4) & (4) & & \\
& 0 & 1 & 2 & 3 & 4 & \\
\end{array}
\]

Of course we can continue this as far down as we like. The recurrence relation for \( \binom{n}{k} \) tells us that each entry in the triangle is the sum of the two entries above it. The entries on the sides of the triangle are always 1. This is because \( \binom{n}{0} = 1 \) for all \( n \) - there is only one way to pick 0 of \( n \) objects - and \( \binom{n}{n} = 1 \) since there is one way to select all \( n \) out of \( n \) objects. Using the recurrence relation, and the fact that the sides of the triangle are 1’s, we can easily replace all the entries above with the correct values of \( \binom{n}{k} \). Doing so gives us Pascal’s triangle.

We can use Pascal’s triangle to calculate binomial coefficients. For example, using the triangle on the next page, we can find \( \binom{12}{6} = 924 \).
### 3.5 Combinations and Permutations

A *permutation* is a (possible) rearrangement of objects. For example, there are 6 permutations of the letters *abc*:

\[ abc, acb, bac, bca, cab, cba \]

We know that we have them all listed above - there are 3 choices for which letter we put first, then 2 choices for which letter comes next, which leaves only 1 choice for the last letter. The multiplicative principle says we multiply these numbers.

**Example:** How many permutations are there of the letters *abcdef*?

**Solution:** We do NOT want to try to list all of these out. However, if we did, we would need to pick a letter to write down first. There are 6 choices for that letter. For each choice of first letter, there are 5 choices for the second letter (we cannot repeat the first letter), and for each of those, there are 4 choices for the third, 3 choices for the fourth, 2 choices for the fifth and finally only 1 choice for the last letter. So there are \(6 \cdot 5 \cdot 4 \cdot 3 \cdot 2 \cdot 1 = 720\) permutations of the 6 letters. We could also write this as \(6!\) (that is, 6 factorial).

Sometimes we do not want to permute all of the letters.

**Example:** How many 4 letter “words” can you make from the letters a-f, with no repeated letters?

**Solution:** This is just like the problem of permuting 4 letters. Only now we have more choices for each letter. For the first letter, there are 6 choices. For each of those, there are 5 choices for the second letter. Then there are 4 choices for the third letter, and 3 choices for the last letter. The total number of words is \(6 \cdot 5 \cdot 4 \cdot 3 = 360\). This is not \(6!\), but it is almost. The difference is that in \(6!\) we multiply by 2 and 1, where as here we did not. So to cancel the 2 and 1, we could write \(\frac{6!}{2!}\).

In general, we can ask how many permutations exist of \(k\) objects choosing those objects from a larger \(n\)-many objects. (In the example above, \(k = 4\), and \(n = 6\)). We write this number \(P(n, k)\). From the example above we see that to compute \(P(n, k)\) we must apply the multiplicative principle to \(k\) numbers, starting with \(n\) and counting backwards. So for example

\[ P(10, 4) = 10 \cdot 9 \cdot 8 \cdot 7 \]

Notice again that \(P(10, 4)\) starts out looking like \(10!\), but we stop after 7. So

\[ P(10, 4) = \frac{10 \cdot 9 \cdot 8 \cdot 7 \cdot 6 \cdot 5 \cdot 4 \cdot 3 \cdot 2 \cdot 1}{6 \cdot 5 \cdot 4 \cdot 3 \cdot 2 \cdot 1} = \frac{10!}{6!} \]

Careful - the factorial in the denominator is not \(4!\) but rather \((10 - 4)!\).
Permutations

\( P(n, k) \) is the number of permutations of \( k \) out of \( n \) objects.

\[
P(n, k) = \frac{n!}{(n-k)!}
\]

Note that when \( n = k \), we have \( P(n, n) = \frac{n!}{(n-n)!} = n! \) (since 0! = 1). This makes sense - we already had that \( n! \) gives the number of permutations of all \( n \) objects.

Here is another way to find the number of permutations of \( k \) out of \( n \) objects: first select which \( k \) objects will be in the permutation, then count how many ways there are to arrange them. Once you have selected the \( k \) objects, we know there are \( k! \) ways to arrange (permute) them. But how do you select \( k \) objects from the \( n \)? You have \( n \) objects, and you need to choose \( k \) of them. You can do that in \( \binom{n}{k} \) ways. This says:

\[
P(n, k) = \binom{n}{k} \cdot k!
\]

Now since we have a closed formula for \( P(n, k) \) already, we can substitute that in:

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

If we divide both sides by \( k! \) we get a closed formula for \( \binom{n}{k} \).

Closed formula for \( \binom{n}{k} \)

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

We say \( P(n, k) \) counts permutations, and \( \binom{n}{k} \) counts combinations. The formulas for each are very similar - there is just an extra \( k! \) in the denominator of \( \binom{n}{k} \). That extra \( k! \) accounts for the fact that \( \binom{n}{k} \) does not distinguish between the different orders that the \( k \) objects can appear in - we are just selecting (or choosing) the \( k \) objects, not arranging them. Perhaps “combination” is a misleading label - we don’t mean it like a combination lock (where the order would definitely matter). Perhaps better: a combination of flavors - you just need to decide which flavors to combine, not the order in which to combine them.

To further illustrate the connection between combinations and permutations, we close with an example.

**Example:** You decide to have a dinner party. Even though you are incredibly popular and have 14 different friends, you only have enough chairs to invite 6 of them.

1. How many choices do you have for which 6 friends to invite?
2. What if you need to decide not only which friends to invite but also in which order to invite them in? How many choices do you have then?

Solution:

1. You must simply choose 6 friends from a group of 14. This can be done in \( \binom{14}{6} \) ways. We can find this number either by using Pascal’s triangle or the closed formula \( \frac{14!}{8! \cdot 6!} = 3003 \).

2. Here you must count all the ways you can permute 6 friends chosen from a group of 14. So the answer is \( P(14, 6) \), which can be calculated as \( \frac{14!}{8!} = 2192190 \).

How are these numbers related? Notice that \( P(14, 6) \) is much larger than \( \binom{14}{6} \). This makes sense - \( \binom{14}{6} \) picks 6 friends, but \( P(14, 6) \) arranges them 6 friends as well as picks them. In fact, we can say exactly how much larger \( P(14, 6) \) is. In both counting problems we choose 6 out of 14 friends. For the first one, we stop there - at 3003 ways. But for the second counting problem, each of those 3003 choices of 6 friends can be arranged in exactly 6! ways. So now we have 3003 \cdot 6! choices - and that is exactly 2192190.

Alternatively, look at the first problem another way. We want to select 6 out of 14 friends, but we do not care about the order they are selected in. To select 6 out of 14 friends, we might try this:

\[
14 \cdot 13 \cdot 12 \cdot 11 \cdot 10 \cdot 9
\]

This is a reasonable guess, since we have 14 choices for the first guest, then 13 for the second, and so on. But the guess is wrong (in fact, that product is exactly 2192190 = \( P(14, 6) \)). It distinguishes between the different orders in which we could invite the guests. To correct for this, we could divide by the number of different arrangements of the 6 guests (so that all of these would count as just one outcome). There are precisely 6! ways to arrange 6 guests, so the correct answer to the first question is

\[
\frac{14 \cdot 13 \cdot 12 \cdot 11 \cdot 10 \cdot 9}{6!}
\]

Note that another way to write this is

\[
\frac{14!}{8! \cdot 6!}
\]

which is what we had originally.

### 3.6 Patterns in Pascal’s Triangle

Have a look again at Pascal’s triangle. Forget for a moment where it comes from - just look at it as a mathematical object. What do you notice?
There are lots of patterns hidden away in the triangle - enough to fill a reasonably sized book. Here are just a few of the most obvious ones:

1. The entries on the border of the triangle are all 1.
2. Any entry not on the border is the sum of the two entries above it.
3. The triangle is symmetric - on any row, entries on the left side are mirrored on the right side.
4. The sum of all entries on a given row is a power of 2. (You should check this!)

We would like to state these observations in a more precise way, and then prove that they are correct. Now each entry in Pascal’s triangle is in fact a binomial coefficient. The 1 on the very top of the triangle is \( \binom{0}{0} \). The next row (which we will call row 1, even though it is not the top-most row) consists of \( \binom{1}{0} \) and \( \binom{1}{1} \). Row 4 (the row 1, 4, 6, 4, 1) consists of the binomial coefficients

\[
\binom{4}{0}, \binom{4}{1}, \binom{4}{2}, \binom{4}{3}, \binom{4}{4}
\]

Given this description of the elements in Pascal’s triangle, we can rewrite the above observations as follows:

1. \( \binom{n}{0} = 1 \) and \( \binom{n}{n} = 1 \)
2. \( \binom{n}{k} = \binom{n-1}{k-1} + \binom{n-1}{k} \)
3. \( \binom{n}{k} = \binom{n}{n-k} \)
4. \( \binom{n}{0} + \binom{n}{1} + \binom{n}{2} + \cdots + \binom{n}{n} = 2^n \)

Each of these are an example of a binomial identity - a identity (i.e., equation) involving binomial coefficients.
Our goal is to establish these identities - prove that they hold for all values of \( n \) and \( k \). These proofs can be done in many ways. One option would be to give algebraic proofs, using the formula for \( \binom{n}{k} \):

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

Here’s how you might do that for the second identity above.

**Example:** Give an algebraic proof for the binomial identity

\[
\binom{n}{k} = \binom{n-1}{k-1} + \binom{n-1}{k}
\]

**Proof.** By the definition of \( \binom{n}{k} \), we have

\[
\binom{n-1}{k-1} = \frac{(n-1)!}{(n-1-(k-1))!(k-1)!} = \frac{(n-1)!}{(n-k)!(k-1)!} \quad \text{and} \quad \binom{n-1}{k} = \frac{(n-1)!}{(n-1-k)!k!}
\]

Thus, starting with the right hand side of the equation we are trying to establish:

\[
\binom{n-1}{k-1} + \binom{n-1}{k} = \frac{(n-1)!}{(n-k)!(k-1)!} + \frac{(n-1)!}{(n-1-k)!k!} = \frac{(n-1)!k}{(n-k)!k!} + \frac{(n-1)!(n-k)}{(n-k)!k!} = \frac{(n-1)!(k+n-k)}{(n-k)!k!} = \frac{n!}{(n-k)!k!} = \binom{n}{k}
\]

The second line above (where the common denominator is found) works because \( k(k-1)! = k! \) and \( (n-k)(n-k-1)! = (n-k)! \). QED

This is certainly a valid proof, but also is entirely useless. Even if you understand the proof perfectly, it does not tell you *why* the identity is true. In other words, the proof does not have any explanatory power. A better approach would be to explain what \( \binom{n}{k} \) means and then say why that is also what \( \binom{n-1}{k-1} + \binom{n-1}{k} \) means. Let’s see how this works for the four identities we have, um, identified.

**Example:** Explain why \( \binom{n}{0} = 1 \) and \( \binom{n}{n} = 1 \).

**Solution:** What do these binomial coefficients tell us? Well, \( \binom{n}{0} \) gives the number of ways to select 0 objects from a collection of \( n \) objects. There is only one way to do this - to not select any of the objects. Thus \( \binom{n}{0} = 1 \). Similarly,
\( \binom{n}{k} \) gives the number of ways to select \( n \) objects from a collection of \( n \) objects. There is only one way to do this - select all \( n \) objects. Thus \( \binom{n}{n} = 1 \).

Alternatively, we know that \( \binom{n}{0} \) is the number of \( n \)-bit strings with weight \( 0 \) - that is, the number of strings made up of \( n \) bits, \( 0 \) of which are \( 1 \)'s. There is only one such string - the string of all \( 0 \)'s. So \( \binom{n}{0} = 1 \). Similarly \( \binom{n}{n} \) is the number of \( n \)-bit strings with weight \( n \). There is only one string - the string of all \( 1 \)'s - with this property.

Another way: \( \binom{n}{0} \) gives the number of subsets of a set of size \( n \) containing \( 0 \) elements. There is only one such subset - the empty set. \( \binom{n}{n} \) gives the number of subsets containing \( n \) elements. The only such subset is the original set (of all elements).

**Example:** Explain why \( \binom{n}{k} = \binom{n-1}{k-1} + \binom{n-1}{k} \)

*Solution:* The easiest way to see this is to consider bit strings. \( \binom{n}{k} \) is the number of bit strings of length \( n \) containing \( k \) \( 1 \)'s. Of all of these strings, some start with a \( 1 \) and the rest start with a \( 0 \). First consider all the bit strings which start with a \( 1 \). After the \( 1 \), there must be \( n - 1 \) more bits (to get the total length up to \( n \)) and exactly \( k - 1 \) of them must be \( 1 \)'s (as we already have one, and we need \( k \) total). How many strings are there like that? There are exactly \( \binom{n-1}{k-1} \) such bit strings, so of all the length \( n \) bit strings containing \( k \) \( 1 \)'s, \( \binom{n-1}{k-1} \) of them start with a \( 1 \). Similarly, there are \( \binom{n-1}{k} \) which start with a \( 0 \) - we still need \( n - 1 \) bits and now \( k \) of them must be \( 1 \)'s. Since there are \( \binom{n-1}{k} \) bit strings containing \( n - 1 \) bits with \( k \) \( 1 \)'s, that is the number of length \( n \) bit strings with \( k \) \( 1 \)'s which start with a \( 0 \).

Another way: consider the question - how many ways can you select \( k \) pizza toppings from a menu containing \( n \) choices? One way to do this is just \( \binom{n}{k} \). Another way to answer the same question is to first decide whether or not you want anchovies. If you do want anchovies, you still need to pick \( k - 1 \) toppings, now from just \( n - 1 \) choices. That can be done in \( \binom{n-1}{k-1} \) ways. If you do not want anchovies, then you still need to select \( k \) toppings from \( n - 1 \) choices (the anchovies are out). You can do that in \( \binom{n-1}{k} \) ways. Since the choices with anchovies are disjoint from the choices without anchovies, the total choices are \( \binom{n-1}{k-1} + \binom{n-1}{k} \). But wait. We answered the same question in two different ways - so the two answers must be the same. Thus \( \binom{n}{k} = \binom{n-1}{k-1} + \binom{n-1}{k} \).

You can also explain (prove) this identity by counting subsets, or even lattice paths.

**Example:** Prove the binomial identity \( \binom{n}{k} = \binom{n}{n-k} \).

*Solution:* Why is this true? \( \binom{n}{k} \) counts the number of ways to select \( k \) things from \( n \) choices. On the other hand, \( \binom{n}{n-k} \) counts the number of ways to select \( n - k \) things from \( n \) choices. Are these really the same? Well, what if instead of selecting the \( n - k \) things you choose to exclude them. How many ways are
there to choose \( n - k \) things to exclude from \( n \) choices. Clearly this is \( \binom{n-k}{n-k} \)
as well (it doesn’t matter whether you include or exclude the things once you have chosen them). And if you exclude \( n - k \) things, then you are including the other \( n \) things. So the set of outcomes should be the same.

Let’s do the pizza counting example like we did above. How many ways are there to pick \( k \) toppings from a list of \( n \) choices? On the one hand, then answer is simply \( \binom{n}{k} \). Alternatively, you could make a list of all the toppings you don’t want. To end up with a pizza containing exactly \( k \) toppings, you need to pick \( n - k \) toppings to not put on the pizza. You have \( \binom{n}{n-k} \) choices for the toppings you don’t want. Both of these ways give you a pizza with \( k \) toppings - in fact all the ways to get a pizza with \( k \) toppings. Thus these two answers must be the same: \( \binom{n}{k} = \binom{n}{n-k} \).

You can also prove (explain) this identity using bit strings, subsets, or lattice paths. The bit string argument is nice: \( \binom{n}{k} \) counts the number of bit strings of length \( n \) with \( k \) 1’s. This is also the number of bit string of length \( n \) with \( n-k \) 0’s though - just replace each 1 with a 0 and each 0 with a 1. But if a string of length \( n \) has \( k \) 0’s, it must have \( n - k \) 1’s. And there are exactly \( \binom{n}{n-k} \) strings of length \( n \) with \( n - k \) 1’s.

**Example:** Prove the binomial identity \( \binom{n}{0} + \binom{n}{1} + \binom{n}{2} + \cdots + \binom{n}{n} = 2^n \)

**Solution:** Let’s do a “pizza proof” again. We need to find a question about pizza toppings which has \( 2^n \) as the answer. How about this: If a pizza joint offers \( n \) toppings, how many pizzas can you build using any number of toppings from no toppings to all toppings, using each topping at most once?

On one hand, the answer is \( 2^n \) - for each topping you can say “yes” or “no,” so you have two choices for each topping.

On the other hand, divide the possible pizzas into disjoint groups - the pizzas with no toppings, the pizzas with one topping, the pizzas with two toppings, etc. If we want no toppings, there is only one pizza like that (the empty pizza, if you will) but it would be better to think of that number as \( \binom{n}{0} \) - we choose 0 of the \( n \) toppings. How many pizzas have 1 topping? We need to choose 1 of the \( n \) toppings, so \( \binom{n}{1} \). We have:

- Pizzas with 0 toppings: \( \binom{n}{0} \)
- Pizzas with 1 topping: \( \binom{n}{1} \)
- Pizzas with 2 toppings: \( \binom{n}{2} \)
  :
- Pizzas with \( n \) toppings: \( \binom{n}{n} \).

The total number of possible pizzas will be the sum of these, which is exactly the left hand side of the identity we are trying to prove.

Again, we could have proved the identity using subsets, bit strings, or lattice paths (although the lattice path argument is a little tricky).
Hopefully this gives some idea of how explanatory proofs of binomial identities can go. It is worth pointing out that sometimes more traditional proofs are also nice.\(^3\) For example, consider the following rather slick proof of the last identity.

Recall the binomial theorem:

\[
(x + y)^n = \binom{n}{0}x^n + \binom{n}{1}x^{n-1}y + \binom{n}{2}x^{n-2}y^2 + \cdots + \binom{n}{n-1}x \cdot y^n + \binom{n}{n}y^n
\]

Let \(x = 1\) and \(y = 1\). We get:

\[
(1 + 1)^n = \binom{n}{0}1^n + \binom{n}{1}1^{n-1}1 + \binom{n}{2}1^{n-2}1^2 + \cdots + \binom{n}{n-1}1 \cdot 1^n + \binom{n}{n}1^n
\]

Of course this simplifies to:

\[
(2)^n = \binom{n}{0} + \binom{n}{1} + \binom{n}{2} + \cdots + \binom{n}{n-1} + \binom{n}{n}
\]

Something fun to try: Let \(x = 1\) and \(y = 2\). Neat huh?

### 3.7 Combinatorial Proof

The explanatory proofs given in the above examples are typically called *combinatorial proofs*. In general, to give a combinatorial proof for a binomial identity, say \(A = B\) you do the following:

1. Find a counting problem you will be able to answer in two ways.
2. Explain why one answer to the counting problem is \(A\).
3. Explain why the other answer to the counting problem is \(B\).

Since both \(A\) and \(B\) are the answers to the same question, we must have \(A = B\).

The tricky thing is coming up with the question. This is not always obvious, but it gets easier the more counting problems you solve. You will start to recognize types of answers as the answers to types of questions. More often what will happen is you will be solving a counting problem and happen to think up two different ways of finding the answer. Now you have a binomial identity and the proof is right there - the proof is the problem you just solved together with your two solutions.

For example, consider this counting question:

How many 10-letter words contain exactly four A’s, three B’s, two C’s and one D?

\(^3\)Most every binomial identity can be proved using mathematical induction, using the recursive definition for \(\binom{n}{k}\). We will discuss induction later in the course.
Let’s try to solve this problem. We have 10 spots for letters to go. Four of those need to be A’s. We can pick the four A-spots in \( \binom{10}{4} \) ways. Now where can we put the B’s? Well there are only 6 spots left, we need to pick 3 of them. This can be done in \( \binom{6}{3} \) ways. The two C’s need to go in two of the 3 remaining spots, so we have \( \binom{3}{2} \) ways of doing that. That leaves just one spot of the D, but we could write that 1 choice as \( \binom{1}{1} \). Thus the answer is:

\[
\binom{10}{4} \binom{6}{3} \binom{3}{2} \binom{1}{1}
\]

But why stop there? We can find the answer another way too. First let’s decide where to put the one D: we have 10 spots, we need to choose 1 of them, so this can be done in \( \binom{10}{1} \) ways. Next, choose one of the \( \binom{9}{2} \) ways to place the two C’s. We now have 7 spots left, and three of them need to be filled with B’s. There are \( \binom{7}{3} \) ways to do this. Finally the A’s can be placed in \( \binom{4}{4} \) (that is, only one) ways. So another answer to the question is

\[
\binom{10}{1} \binom{9}{2} \binom{7}{3} \binom{4}{4}
\]

Interesting. This gives us the binomial identity:

\[
\binom{10}{4} \binom{6}{3} \binom{3}{2} \binom{1}{1} = \binom{10}{1} \binom{9}{2} \binom{7}{3} \binom{4}{4}
\]

Here are a couple of other binomial identities with combinatorial proofs.

**Example:** Prove the identity

\[
1n + 2(n − 1) + 3(n − 2) + \cdots + (n − 1)2 + n1 = \binom{n + 2}{3}
\]

**Solution:** To give a combinatorial proof we need to think up a question we can answer in two ways: one way needs to give the left-hand-side of the identity, the other way needs to be the right-hand-side of the identity. Our clue to what question to ask comes from the right hand side: \( \binom{n + 2}{3} \) counts the number of ways to select 3 things from a group of \( n + 2 \) things. Let’s name those things \( 1, 2, 3, \ldots, n + 2 \). In other words, we want to find 3-element subsets of those numbers (since order should not matter, subsets are exactly the right thing to think about). We will have to be a bit clever to explain why the left-hand-side also gives the number of these subsets. Here’s the proof.

**Proof.** Consider the question “How many 3-element subsets are there of the set \( \{1, 2, 3, \ldots, n + 2\}\)?” We answer this in two ways:

**Answer 1:** We must select 3 elements from the collection of \( n + 2 \) elements. This can be done in \( \binom{n + 2}{3} \) ways.

**Answer 2:** Break this problem up into cases, by what the middle number in the subset is. That is, each subset is \( \{a, b, c\} \), say written in increasing order.
We count how many such subsets there are for each distinct value of $b$. The smallest possible value of $b$ is 2, and the largest is $n + 1$.

When $b = 2$, there are $1 \cdot n$ subsets - 1 choice for $a$ and $n$ choices (3 through $n + 2$) for $c$.

When $b = 3$, there are $2 \cdot (n - 1)$ subsets - 2 choices for $a$ and $n - 1$ choices for $c$.

When $b = 4$, there are $3 \cdot (n - 2)$ subsets - 3 choices for $a$ and $n - 2$ choices for $c$.

And so on. When $b = n + 1$, there are $n$ choices for $a$ and only 1 choice for $c$, so $n \cdot 1$ subsets.

Therefore the total number of subsets is

$$1n + 2(n - 1) + 3(n - 2) + \cdots + (n - 1)2 + n1.$$

Since Answer 1 and Answer 2 are answers to the same question, they must be equal. Therefore

$$1n + 2(n - 1) + 3(n - 2) + \cdots + (n - 1)2 + n1 = \binom{n + 2}{3}$$

QED

**Example:** Prove the binomial identity

$$\binom{n}{0}^2 + \binom{n}{1}^2 + \binom{n}{2}^2 + \cdots + \binom{n}{n}^2 = \binom{2n}{n}$$

**Solution:** We will give two different proofs of this fact. The first will be very similar to the previous example (counting subsets). The second proof is a little slicker - it will use lattice paths.

**Proof.** Consider the question: “How many pizzas can you make using $n$ toppings when there are $2n$ toppings to choose from?”

**Answer 1:** There are $2n$ toppings, from which you must choose $n$. This can be done in $\binom{2n}{n}$ ways.

**Answer 2:** Divide the toppings into two groups of $n$ toppings (perhaps $n$ meats and $n$ veggies). Any choice of $n$ toppings must include some number from the first group and some number from the second group. Consider each possible number of meat toppings separately:

- 0 meats: $\binom{n}{0} \binom{n}{n}$, since you need to choose 0 of the $n$ meats and $n$ of the $n$ veggies.

- 1 meat: $\binom{n}{1} \binom{n}{n-1}$, since you need 1 of $n$ meats so $n - 1$ of $n$ veggies.

- 2 meats: $\binom{n}{2} \binom{n}{n-2}$. Choose 2 meats and the remaining $n - 2$ toppings from the $n$ veggies.

And so on. The last case is $n$ meats, which can be done in $\binom{n}{n} \binom{n}{0}$ ways.
Thus the total number of pizzas possible is
\[
\binom{n}{0} \binom{n}{n} + \binom{n}{1} \binom{n}{n-1} + \binom{n}{2} \binom{n}{n-2} + \cdots + \binom{n}{n} \binom{n}{0}
\]
This is not quite the left hand side . . . yet. Notice that \(\binom{n}{n} = \binom{n}{0}\) and \(\binom{n}{n-1} = \binom{n}{1}\) and so on, by the symmetry formula. Thus we do indeed get
\[
\binom{n}{0}^2 + \binom{n}{1}^2 + \binom{n}{2}^2 + \cdots + \binom{n}{n}^2
\]
QED

Here’s the second proof:

*Proof.* Consider the question: “How many lattice paths are there from \((0,0)\) to \((n,n)\)?”

**Answer 1:** We must travel \(2n\) units, and \(n\) of them must be in the up direction. Thus there are \(\binom{2n}{n}\) paths.

**Answer 2:** Note that any path from \((0,0)\) to \((n,n)\) must cross the line \(x + y = n\). That is, must pass through exactly one of the points: \((0,n)\), \((1,n-1)\), \((2,n-2)\), . . . , \((n,0)\). For example, this is what happens in the case \(n = 4\):

How many paths pass through \((0,n)\)? To get to that point, you must travel \(n\) units, and 0 of them are to the right, so there are \(\binom{n}{0}\) ways to get to \((0,n)\). From \((0,n)\) to \((n,n)\) takes \(n\) steps, and 0 of them are up. So there are \(\binom{n}{0}\) ways to get from \((0,n)\) to \((n,n)\). Therefore there are \(\binom{n}{0} \binom{n}{0}\) paths from \((0,0)\) to \((n,n)\) through the point \((0,n)\).

What about through \((1,n-1)\). There are \(\binom{n}{1}\) paths to get there \((n\) steps, 1 to the right) and \(\binom{n}{1}\) paths to complete the journey to \((n,n)\) \((n\) steps, 1 up). So there are \(\binom{n}{1} \binom{n}{1}\) paths from \((0,0)\) to \((n,n)\) through \((1,n-1)\).

In general, to get to \((n,n)\) through the point \((k,n-k)\) we have \(\binom{n}{k}\) paths to the mid point and then \(\binom{n}{k}\) paths from the mid point to \((n,n)\). So there are \(\binom{n}{k} \binom{n}{k}\) paths from \((0,0)\) to \((n,n)\) through \((k,n-k)\).
All together then the total paths from \((0, 0)\) to \((n, n)\) passing through exactly one of these mid points is

\[
\binom{n}{0}^2 + \binom{n}{1}^2 + \binom{n}{2}^2 + \cdots + \binom{n}{n}^2
\]

QED

3.8 Stars and Bars

Consider the following counting problem:

You have 7 cookies to give to 4 kids. How many ways can you do this?

Take a moment to think about how you might solve this problem. You may assume that it is acceptable to give a kid no cookies. Also, the cookies are all identical. And the order in which you give out the cookies does not matter.

Before solving the problem, here is a wrong answer. You might guess that the answer should be \(4^7\) because for each of the 7 cookies, there are 4 choices of kids to give the cookie too. This is reasonable, but wrong. To see why, consider a few possible outcomes: we could assign the first six cookies to kid A, and the seventh cookie to kid B. Another outcome would assign the first cookie to kid B and the six remaining cookies to kid A. Both outcomes are included in the \(4^7\) answer. But for our counting problem, both outcomes are really the same - kid A gets six cookies and kid B gets one cookie.

What do outcomes actually look like? How can we represent them? One approach would be to write a outcome as a string of four numbers like this:

\[3112\]

which represent the outcome in which the first kid gets 3 cookies, the second and third kid each get 1 cookie, and the fourth kid gets 2 cookies. Represented this way, the order in which the numbers occur matters - 1312 is a different outcome, because the first kid gets a one cookie instead of 3. Each number in the string can be any integer between 0 and 7. But the answer is not \(7^4\) - we need the sum of the numbers to be 7.

Another way we might represent outcomes is to write a string of seven letters:

\[ABAADCD\]

which represents that the first cookie goes to kid A, the second cookie goes to kid B, the third and fourth cookies go to kid A, and so on. In fact, this outcome is identical to the previous one - A gets 3 cookies, B and C get 1 each and D gets 2. Each of the seven letters in the string can be any of the 4 possible letters (one for each kid), but the number of such strings is not \(4^7\), because here order does not matter. In fact, another way to write the same outcome is

\[AAABCDD\]
This will be the preferred representation of the outcome. Since we can write the letters in any order, we might as well write them in non-decreasing order for the purposes of counting. So we will write all the A’s first, then all the B’s, and so on.

Now think about how you could specify such an outcome. All we really need to do is say when to switch from one letter to the next. In terms of cookies, we need to say after how many cookies do we stop giving cookies to the first kid and start giving cookies to the second kid. And then after how many do we switch to the third kid? And after how many do we switch to the fourth? So yet another way to represent an outcome is like this:

```
*** | * | * | *
```

Three cookies go to the first kid, then we switch and give one cookie to the second kid, then switch, one to the third kid, switch, two to the fourth kid. Notice that we need 7 stars and 3 bars - one star for each cookie, and one bar for each switch between kids, so one fewer bars than there are kids (we don’t need to switch after the last kid - we are done).

Why have we done all of this? Simple: to count the number of ways to distribute 7 cookies to 4 kids, all we need to do is count how many stars and bars charts there are. But a stars and bars chart is just a string of symbols, some stars and some bars. If instead of stars and bars we would use 0’s and 1’s, it would just be a bit string. We know how to count those.

Before we get too excited, we should make sure that really any string of (in our case) 7 stars and 3 bars corresponds to a different way to distribute cookies to kids. In particular consider a string like this:

```
| *** | | ***
```

Does that correspond to a cookie distribution? Yes. It represents the distribution in which kid A gets 0 cookies (because we switch to kid B before any stars), kid B gets three cookies (three stars before the next bar), kid C gets 0 cookies (no stars before the next bar) and kid D gets the remaining 4 cookies. No matter how the stars and bars are arranged, we can distribute cookies in that way. Also, given any way to distribute cookies, we can represent that with a stars and bars chart. For example, the distribution in which kid A gets 6 cookies and kid B gets 1 cookie has the following chart:

```
** | * || **
```

Now after all that work we are finally ready to count. Each way to distribute cookies corresponds to a stars and bars chart with 7 stars and 3 bars. So there are 10 symbols, and we must choose 3 of them to be bars. Thus:

\[
\text{There are } \binom{10}{3} \text{ ways to distribute 7 cookies to 4 kids.}
\]

While we are at it, we can also answer a related question: how many ways are there to distribute 7 cookies to 4 kids so that each kid gets at least one cookie? What can you say about the corresponding stars and bars charts? The charts must start and end with at least one star (so that kids A and D) get cookies, and also no two bars can be adjacent (so that kids B and C are not skipped). One way to assure this is to only place bars in the spaces
between the stars. With 7 stars, there are 6 spots between the stars, so we must choose 3 of those 6 spots to fill with bars. Thus there are \( \binom{6}{3} \) ways to distribute 7 cookies to 4 kids giving at least one cookie to each kid.

Another (and more general) way to approach this modified problem is to first give each kid one cookie. Now the remaining 3 cookies can be distributed to the 4 kids without restrictions. So we have 3 stars and 3 bars for a total of 6 symbols, 3 of which must be bars. So again we see that there are \( \binom{6}{3} \) ways to distribute the cookies.

Stars and bars can be used in counting problems other than kids and cookies. Here are a few examples.

**Example:** Your favorite mathematical pizza chain offers 10 toppings. How many pizzas can you make if you are allowed 6 toppings? The order of toppings does not matter, but now you are allowed repeats - so one possible pizza is: triple sausage, double pineapple, and onions.

**Solution:** We get 6 toppings (counting possible repeats). Represent each of these toppings as a star. There are 10 toppings to choose from, so we must switch from considering one topping to the next 9 times - these are the bars. Think of going down the menu one topping at a time: you see anchovies first, and skip to the next, sausage. You say yes to sausage 3 times then switch to the next topping on the list. You keep skipping until you get to pineapple, which you say yes to twice. Another switch and you are at onions - you say yes once. Then you keep switching until you get to the last topping, never saying yes again (since you already have said yes 6 times.

Now that we are confident that we have the right number of stars and bars, we answer the question simply: there are 6 stars and 9 bars, so 15 symbols. We need to pick 9 of them to be bars, so there number of pizzas possible is

\[
\binom{15}{9}
\]

**Example:** How many 7 digit phone numbers are there in which the digits are non-increasing? That is, every digit is less than or equal to the previous one.

**Solution:** We need to decide on 7 digits - so we will use 7 stars. The bars will represent a switch from each possible single digit number down the next smaller one. So the phone number 866-5221 is represented by the stars and bars chart

```
| * | * | * * | * | * | * |
```

There are 10 choices for each digit (0-9) so we must switch between choices 9 times. We have 7 stars and 9 bars, so the total number of phone numbers is

\[
\binom{16}{9}
\]
**Example:** How many integer solutions\(^4\) are there to the equation

\[ x_1 + x_2 + x_3 + x_4 + x_5 = 13 \]

1. where \(x_i \geq 0\) for each \(x_i\)?
2. where \(x_i > 0\) for each \(x_i\)?
3. where \(x_i \geq 2\) for each \(x_i\)?

**Solution:** This problem is just like giving 13 cookies to 5 kids - we need to say how many of the 13 units go to each of the 5 variables. In other words, we have 13 stars and 4 bars (the bars are like the “+” signs in the equation).

1. If \(x_i\) can be 0 or greater, we are in the standard case with no restrictions. So 13 stars and 4 bars can be arranged in \(\binom{17}{4}\) ways.
2. Now each variable must be at least 1. So give one unit to each variable to satisfy that restriction. Now there are 8 stars left, and still 4 bars, so the number of solutions is \(\binom{12}{4}\).
3. Now each variable must be 2 or greater. So before any counting, give each variable 2 free units. We now have 3 remaining stars and 4 bars, so there are \(\binom{7}{4}\) solutions.

### 3.9 Functions

Have we now discovered all the rules needed in counting? Probably not. To see what we can do and what we still must figure out, let’s shift focus and think about counting in another way. Consider again the problem if giving cookies to kids. One way to think about this problem is to say that we are assigning each cookie to a kid. Now in the case we have looked at so far, the cookies were identical, so it did not matter which cookie went to which kid. But what if it did? Suppose we have 7 different cookies and 4 different kids - how many ways are there to give out the cookies?

Perhaps you already have a solution in mind. However, for the time being, we will consider this question more generally. What are we really doing here? We are mapping each cookie to a kid. This sounds like a function! It is. To help us completely characterize counting problems, we should look at them as problems about how to count functions. Before we proceed, here is a quick review of functions.

A function is simply a rule that assigns each input exactly one output. The set of all inputs for a function is called the *domain*. The set of all allowable outputs is called the *codomain*. For example, a function might assign each natural number a natural number from 1 to 5. In that case, the domain is the natural numbers and the codomain is the set of natural numbers from 1 to 5. Now it could be that this particular function we are thinking about assigns each even natural number to the number 2 and each odd natural number to the number 1. In this case, not all of the codomain is actually used. We would say that

\(^4\)An integer solution to an equation is a solution in which the unknown must have an integer value.
the set \( \{1, 2\} \) is the range of the function - these are the element in the codomain (allowable outputs) which are actually outputs for some input.

The key thing that makes a rule assigning inputs to outputs a function is that there is only one output for an input. In other words, it is important that the rule be a good rule. What output do we assign to the input 7? There can only be one answer for a particular function - otherwise where does 7 go to?

To specify the name of the function, as well as the domain and codomain, we write \( f : X \rightarrow Y \). The function is called \( f \), the domain is the set \( X \) and the codomain is the set \( Y \). This however does not describe the rule. To do that, we say something like this:

The function \( f : X \rightarrow Y \) is defined by \( f(x) = x^2 + 3 \).

In this case, the function takes an input \( x \) and computes the output by squaring \( x \) and then adding 3. In this case, you better hope that \( X \) is a set of numbers and \( Y \) is a set of number which can be 3 more than squares of numbers from \( X \). It would not work for \( Y \) to be negative numbers here - that would not be a valid function.

The description of the rule can vary greatly. If \( X \) is a finite set, we might just give a list of each output for each input. You could also describe the function with at able or a graph or in words.

**Example:** The following are all examples of functions:

1. \( f : \mathbb{Z} \rightarrow \mathbb{Z} \) defined by \( f(n) = 3n \). The domain and codomain are both the set of integers. However, the range is only the set of integer multiples of 3.

2. \( f : \{1, 2, 3\} \rightarrow \{a, b, c\} \) defined by \( f(1) = c, f(2) = a \) and \( f(3) = a \). The domain is the set \( \{1, 2, 3\} \), the codomain is the set \( \{a, b, c\} \) and the range is the set \( \{a, c\} \). Note that \( f(2) \) and \( f(3) \) are the same element of codomain. This is not a problem - each element in the domain still has only one output (although each output does not have a unique input).

3. \( f : \{1, 2, 3\} \rightarrow \{1, 2, 3\} \) defined as follows:

\[
\begin{array}{ccc}
1 & 2 & 3 \\
\hline
1 & 2 & 3
\end{array}
\]

The arrow diagram used to define the function above can be very helpful in visualizing functions. We will often be working with functions on finite sets so this kind of picture is often more useful than a traditional graph of a function. A graph of the function in example 3 above would look like this:
It is important to know how to recognize a function from a rule which is not a function. The arrow diagram can help.

**Example:** Which of the following diagrams represent a function. Let $X = \{1, 2, 3, 4\}$ and $Y = \{a, b, c, d\}$

<p>| | | | |</p>
<table>
<thead>
<tr>
<th></th>
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</thead>
<tbody>
<tr>
<td></td>
<td></td>
<td></td>
<td></td>
</tr>
<tr>
<td>$f: X \rightarrow Y$</td>
<td>$g: X \rightarrow Y$</td>
<td>$h: X \rightarrow Y$</td>
<td></td>
</tr>
<tr>
<td>1</td>
<td>2</td>
<td>3</td>
<td>4</td>
</tr>
<tr>
<td>a</td>
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</tr>
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<td>a</td>
<td>b</td>
<td>c</td>
<td>d</td>
</tr>
</tbody>
</table>

**Solution:** $f$ is a function. So is $g$. There is no problem with an element of the codomain not being the output for any input, and there is no problem with $a$ from the codomain being the output of both 2 and 3 from the input.

However, $h$ is not a function - in fact, for two reasons. First, the element 1 from the domain has not been mapped to any element from the codomain. Second, the element 2 from the domain has been mapped to more than one element from the codomain ($a$ and $c$). Note that either one of these problems is enough to make a rule not a function. Neither of these mappings are functions:

Not functions.

### 3.9.1 Surjections, Injections, and Bijections

We now turn to investigating special properties functions might or might not possess.

In the examples above, you may have noticed that sometimes there are elements of the codomain which are not in the range - which are not actual outputs for any input. When this sort of the thing does not happen, (that is, when everything in the codomain is in the range) we say the function is onto or that the function maps the domain onto the codomain. This terminology should make sense: the function puts the domain (entirely) on top the codomain. The fancy math term for an onto function is a surjection, and we say that an onto function is a surjective function.

In pictures:
Example: Which functions are surjective (i.e., onto)?

1. \( f : \mathbb{Z} \rightarrow \mathbb{Z} \) defined by \( f(n) = 3n \).
2. \( g : \{1, 2, 3\} \rightarrow \{a, b, c\} \) defined by \( g(1) = c \), \( g(2) = a \) and \( g(3) = a \).
3. \( h : \{1, 2, 3\} \rightarrow \{1, 2, 3\} \) defined as follows:

\[
\begin{array}{ccc}
1 & 2 & 3 \\
1 & 2 & 3
\end{array}
\]

Solution:

1. \( f \) is not surjective. There are elements in the codomain which are not in the range. For example, no \( n \in \mathbb{Z} \) gets mapped to the number 1 (the rule would say that \( \frac{1}{3} \) would be sent to 1, but \( \frac{1}{3} \) is not in the domain). In fact, the range of the function is \( 3\mathbb{Z} \) (the integer multiples of 3), which is not equal to \( \mathbb{Z} \).

2. \( g \) is not surjective. There is no \( x \in \{1, 2, 3\} \) (the domain) for which \( g(x) = b \). So \( b \), which is in the codomain, is not in the range, so once again the function is not onto.

3. \( h \) is surjective. Every element of the codomain is also in the range. Nothing is missed.

To be a function, a map cannot assign a single element of the domain to two or more different elements of the codomain. However, we have seen that the reverse is permissible. That is, a function might assign the same element of the codomain to two or more different elements of the domain. When this does not occur - that is, when each element of the codomain is assigned to at most one element of the domain - then we say the function is one-to-one. Again, this terminology makes sense: we are sending at most one element from the domain to one element from the codomain. One input to one output. The fancy math term for a one-to-one function is a injection. We call one-to-one functions injective functions.

In pictures:

\[
\begin{array}{c}
\bullet \bullet \bullet \\
\circ \circ \circ \\
\end{array}
\quad
\begin{array}{c}
\bullet \bullet \bullet \\
\circ \circ \circ \\
\end{array}
\]

\[
\begin{array}{c}
\bullet \bullet \bullet \\
\circ \circ \circ \circ \circ \\
\end{array}
\quad
\begin{array}{c}
\bullet \bullet \bullet \\
\circ \circ \circ \circ \circ \\
\end{array}
\]
**Example:** Which functions are injective (i.e., one-to-one)?

1. \( f : \mathbb{Z} \rightarrow \mathbb{Z} \) defined by \( f(n) = 3n \).
2. \( g : \{1, 2, 3\} \rightarrow \{a, b, c\} \) defined by \( g(1) = c \), \( g(2) = a \) and \( g(3) = a \).
3. \( h : \{1, 2, 3\} \rightarrow \{1, 2, 3\} \) defined as follows:

\[
\begin{array}{ccc}
1 & 2 & 3 \\
\hline
1 & 2 & 3 \\
\end{array}
\]

**Solution:**

1. \( f \) is injective. Each element in the codomain is assigned to at *most* one element from the domain - if \( x \) is a multiple of three, then only \( x/3 \) is mapped to \( x \). If \( x \) is not a multiple of 3, then there is no input corresponding to the output \( x \).
2. \( g \) is not injective. Both inputs 2 and 3 are assigned the output \( a \).
3. \( h \) is injective. Each output is only an output once.

From the examples above, it should be clear that there are functions which are surjective, injective, both or neither. In the case when a function is both one-to-one and onto (a injection and surjection) we say the function is a *bijection*, or that the function is a *bijective* function.

---

**Function Definitions**

- A *function* is a rule that assigns each element of a set, called the *domain*, to exactly one element of a second set, called the *codomain*.

- Notation: \( f : X \rightarrow Y \) is our way of saying that the function is called \( f \), the domain is the set \( X \) and the codomain is the set \( Y \).

- \( f(x) = y \) means the element \( x \) of the domain (input) is assigned to the element \( y \) of the codomain. We say \( y \) is an output. Alternatively, we call \( y \) the *image of* \( x \) *under* \( f \).

- The *range* is a subset of the codomain. It is the set of all elements which are assigned to at least one element of the domain by the function. That is, the range is the set of all outputs.

- A function is *one-to-one* (also called an *injection*) if every element of the codomain is the output for *at most* one element from the domain.

- A function is *onto* (also called a *surjection*) if every element of the codomain is the output of *at least* one element of the domain.

- A *bijection* is a function which is both an injection and surjection - so one-to-one and onto.
3.9.2 Counting Functions

Now back to counting. Suppose (yet again) that you have 7 cookies to give to 4 kids. Now the cookies are all different, so which cookie goes to which kid matters. How many ways are there to do this? What we really have is a function with domain the set of cookies and codomain the set of kids. So we might as well ask:

**Example:** How many functions are there $f : \{1, 2, 3, 4, 5, 6, 7\} \rightarrow \{a, b, c, d\}$?

**Solution:** We must assign $f(x)$ for each $x$ in the domain. We have four choices for $f(1)$ - we could have $f(1) = a$, $f(1) = b$, $f(1) = c$ or $f(1) = d$. Similarly, we have four choices for $f(2)$, and four choices for $f(3)$ and so on. In fact, for each of the seven elements of the domain, we have four choices, so the number of functions here is $4^7$.

We might also wonder how many of those functions are injective and how many are surjective. What would these mean in terms of cookies? The function would be injective if each kid got no more than one cookie. There is no way to do this, because we have more cookies than kids. The number of injective functions with this domain and codomain is 0. What if we switch the domain and codomain? So now we are mapping kids to cookies. Note that this is not the same, although it is a perfectly fine counting question (the answer is $P(7, 4)$, do you see why?).

On the other hand, a surjective function would be one in which each kid got at least one cookie. Can we count those? It turns out that counting surjective functions is hard. We will do it, but first back up and consider another simpler example.

**Example:** How many functions are there $f : \{1, 2, 3, 4, 5\} \rightarrow \{a, b, c, d, e\}$?

How many of those functions are injective?

**Solution:** First we count all the functions with domain \(\{1, 2, 3, 4, 5\}\) and codomain \(\{a, b, c, d, e\}\). For $f$ to be a function, it must assign each element of the domain to exactly one element of the codomain. In other words, we need to assign one of the elements of \(\{a, b, c, d, e\}\) to $f(1)$. There are 5 choices for this. Similarly, there are five choices for $f(2)$. In fact, there are 5 choices for $f(x)$ for all five possible values of $x$ $\in \{1, 2, 3, 4, 5\}$. Thus the total number of function is $5 \cdot 5 \cdot 5 \cdot 5 \cdot 5 = 5^5$.

What about the injective (one-to-one) functions? Again, we have 5 choices for $f(1)$. However, once we assign a value to $f(1)$, we cannot assign that value to $f(2)$, or any other $f(x)$.$^5$ So for each of the 5 choices for $f(1)$, we only have 4 choices for $f(2)$, and then only 3 choices for $f(3)$ and only 2 choices for $f(4)$, leaving only 1 choice (the last element of the codomain) for $f(5)$. Therefore the number of one-to-one functions is $5!$.

Using the same counting techniques, you should be able to count the total number of functions from any finite size domain to any finite size codomain, and also count the one-to-

$^5$Remember, a function is injective if and only if different elements from the domain go to *different* elements of the codomain.
one functions (as long as the size of the codomain is at least as large as the domain, otherwise there would be no one-to-one functions).

So what about surjective functions? Well in the previous example, there are exactly $5! = 120$ surjections. We know this because whenever the domain and codomain of a function are the same finite size, a function is injective if and only if it is surjective. In general, if the domain and codomain both contain $n$ elements, then the number of surjective functions is $n!$, since this is the number of injective functions. Also, if the codomain is strictly larger than the domain, then the surjective functions are easy to count - there are 0 of them (since there will always be elements of the codomain left out of the range). But what if the size of the codomain is smaller than the size of the domain?

The idea is to count the functions which are not surjective, and then subtract that from the total number of functions. This works very well when the codomain has two elements in it.

**Example:** How many functions $f : \{1, 2, 3, 4, 5\} \to \{a, b\}$ are surjective?

**Solution:** There are $2^5$ functions all together - there are two choices for where to send each of the 5 elements of the domain. Now of these, the functions which are not surjective must exclude one or more elements of the codomain from the range. So first, consider functions for which $a$ is not in the range. This can only happen one way - everything gets sent to $b$. Alternatively, we could exclude $b$ from the range. Then everything gets sent to $a$, so there is only one function like this. These are the only ways in which a function could not be onto (no function excludes both $a$ and $b$ from the range) so there are exactly $2^5 - 2$ surjective functions.

When there are three elements in the codomain, there are now three choices for a single element to exclude from the range. Additionally, we could pick pairs of two elements to exclude from the range, and we must make sure we don’t over count these.

**Example:** How many functions $f : \{1, 2, 3, 4, 5\} \to \{a, b, c\}$ are surjective?

**Solution:** Again start with the total number of functions: $3^5$ (as each of the five elements of the domain can go to any of three elements of the codomain). Now we count the functions which are not surjective.

Start by excluding $a$ from the range. Then we have two choices ($b$ or $c$) for where to send each of the five elements of the domain. Thus there are $2^5$ functions which exclude $a$ from the range. Similarly, there are $2^5$ functions which exclude $b$, and another $2^5$ which exclude $c$. Now have we counted all functions which are not surjective? Yes, but in fact, we have counted some multiple times. For example, the function which sends everything to $c$ was one of the $2^5$ functions we counted when we excluded $a$ from the range, and also one of the $2^5$ functions we counted when we excluded $b$ from the range. We must subtract out all the functions which specifically exclude two elements from the range. There is 1 function when we exclude $a$ and $b$, one function when we exclude $a$ and $c$, and one function when we exclude $b$ and $c$.
We are using PIE: to count the functions which are not surjective, we added up the functions which exclude \(a\), \(b\), and \(c\) separately, then subtracted the functions which exclude pairs of elements. We would then add back in the functions which exclude groups of three elements, except that there are no such functions. We find that the number of functions which are not surjective is

\[2^5 + 2^5 + 2^5 - 1 - 1 - 1 + 0\]

Perhaps a more descriptive way to write this is

\[\binom{3}{1}2^5 - \binom{3}{2}1^5 + \binom{3}{3}0^5\]

since each of the \(2^5\)'s was the result of choosing 1 of the 3 elements of the codomain to exclude from the range, each of the three \(1^5\)'s was the result of choosing 2 of the 3 elements of the codomain to exclude. Writing \(1^5\) instead of 1 makes sense too: we have 1 choice of were to send each of the 5 elements of the domain.

Now we can finally count the number of surjective functions:

\[3^5 - \left[\binom{3}{1}2^5 - \binom{3}{2}1^5\right] = 150\]

You might worry that to count surjective functions when the codomain is larger than 3 elements would be too tedious - we need to use PIE but with more than 3 sets the formula for PIE is very long. However, we have lucked out. As we saw in the example above, the number of functions which exclude a single element from the range is the same no matter which single element is excluded. Similarly, the number of functions which exclude a pair of elements will be the same for every pair. With larger codomains, we will see the same behavior with groups of 3, 4, and more elements excluded. So instead of adding/subtracting each of these, we can simply add or subtract all of them at once, if you know how many there are. Here’s what happens with 4 and 5 elements in the codomain.

**Example:**

1. How many functions \(f : \{1, 2, 3, 4, 5\} \rightarrow \{a, b, c, d\}\) are surjective?
2. How many functions \(f : \{1, 2, 3, 4, 5\} \rightarrow \{a, b, c, d, e\}\) are surjective?

**Solution:**

1. There are \(4^5\) functions all together - we will subtract the functions which are not surjective. We could exclude any one of the four elements of the codomain, and doing so will leave us with \(3^5\) functions. This counts too many so we subtract the functions which exclude two of the four elements of the codomain, each pair giving \(2^5\) functions. But this excludes too many, so we add back in the functions which exclude three of the four
elements of the codomain, each triple giving a function. There are \( \binom{5}{1} \) groups of functions excluding a single element, \( \binom{5}{2} \) groups of functions excluding a pair of elements, and \( \binom{5}{3} \) groups of functions excluding a triple of elements. This means that the number of functions which are not surjective is:

\[
\binom{4}{1}3^5 - \binom{4}{2}2^5 + \binom{4}{3}1^5
\]

We can now say that the number of functions which are surjective is:

\[
4^5 - \left( \binom{4}{1}3^5 - \binom{4}{2}2^5 + \binom{4}{3}1^5 \right)
\]

2. The number of surjective functions is:

\[
5^5 - \left[ \binom{5}{1}4^5 - \binom{5}{2}3^5 + \binom{5}{3}2^5 - \binom{5}{4}1^5 \right]
\]

We took the total number of functions 5 and subtracted all that were not surjective. There were \( \binom{5}{1} \) ways to select a single element from the codomain to exclude from the range, and for each there were 4 functions. But this double counts, so we use PIE and subtract functions excluding two elements from the range: there are \( \binom{5}{2} \) choices for the two elements to exclude, and for each pair, 3 functions. This takes out too many functions, so we add back in functions which exclude 3 elements from the range: \( \binom{5}{3} \) choices for which three to exclude, and then 2 functions for each choice of elements. Finally we take back out the 1 function which excludes 4 elements for each of the \( \binom{5}{4} \) choices of 4 elements.

### 3.10 Advanced Counting using PIE

Counting surjections used the Principle of Inclusion/Exclusion (PIE), which gives a method for finding the cardinality of the union of not necessarily disjoint sets. The idea is that to find how many things are in one or more of the sets \( A, B, \) and \( C \) we should just add up the number of things in each of these sets. However, if there is any overlap among the sets, those elements are counted multiple times. So we subtract the things in each intersection of a pair of sets. But doing this removes elements which are in all three sets once too often, so we need to add it back in.

For four or more sets, we do not write down a formula for PIE. Instead, we just think of the principle: add up all the elements in single sets, then subtract out things you counted twice (elements in the intersection of a pair of sets), then add back in elements you removed too often (elements in the intersection of groups of three sets), then take back out elements you added back in too often (elements in the intersection of groups of four sets), then add back in, take back out, add back in, etc.

Here are some additional examples of how it can be used.
3.10.1 Counting Derangements

A *derangement* of $n$ elements $\{1, 2, 3, \ldots, n\}$ is a permutation in which no element is fixed. For example, there are 6 permutations of the three elements $\{1, 2, 3\}$:

$$123 \quad 132 \quad 213 \quad 231 \quad 312 \quad 321$$

but most of these have one or more elements fixed: 123 has all three elements fixed, 132 has the first element fixed (1 is in its original first position), and so on. In fact, the only derangements of three elements are

$$231 \text{ and } 312$$

If we go up to 4 elements, there are 24 permutations (because we have 4 choices for the first element, 3 choices for the second, 2 choices for the third leaving only 1 choice for the last - $4! = 24$). How many of these are derangements? If you list out all 24 permutations and eliminate those which are not derangements, you will be left with just 9 derangements. Let’s see how we can get that number using PIE.

**Example:** How many derangements are there of 4 elements?

**Solution:** We count all permutations, and subtract those which are not derangements. There are $4! = 24$ permutations of 4 elements. Now for a permutation to not be a derangement, at least one of the 4 elements must be fixed. There are $\binom{4}{1}$ choices for which single element we fix. Once fixed, we need to find a permutation of the other three elements - there are $3!$ permutations on 3 elements. But now we have counted too many non-derangements, so we must subtract those permutations which fix two elements. There are $\binom{4}{2}$ choices for which two elements we fix, and then for each pair, 2! permutations of the remaining elements. But this subtracts too many, so add back in permutations which fix 3 elements, all $\binom{4}{3}!$ of them. Finally subtract the $\binom{4}{4}0!$ permutations (recall $0! = 1$) which fix all four elements. All together we get that the number of derangements of 4 elements is:

$$4! - \left[ \binom{4}{1}3! - \binom{4}{2}2! + \binom{4}{3}1! - \binom{4}{4}0! \right] = 24 - 15 = 9$$

Of course we can use a similar formula to count the derangements on any number of elements. However, the more elements we have, the longer the formula gets. It turns out that there is no closed formula for counting derangements or surjective functions for that matter. We need to use PIE.

Before moving on to another advanced counting technique, here is another example.

**Example:** Five gentlemen attend a party, leaving their hats at the door. At the end of the party, they hastily grab hats on their way out. How many different ways could this happen so that none of the gentlemen leave with their own hat?
**Solution:** We are counting derangements on 5 elements. There are $5!$ ways for the gentlemen to grab hats in any order - but many of these permutations will result in someone getting their own hat. So we subtract all the ways in which one or more of the men get their own hat. In other words, we subtract the non-derangements. Doing so requires PIE. Thus the answer is:

$$5! - \left[ \binom{5}{1} 4! - \binom{5}{2} 3! + \binom{5}{3} 2! - \binom{5}{4} 1! + \binom{5}{5} 0! \right]$$

### 3.10.2 Stars, Bars, and Pie

Another time PIE crops up is in stars and bars problems. Let’s go back to our 7 cookies and 4 kids problem. Recall we want to distribute the 7 identical cookies to 4 (non-identical) kids. Now we place a restriction on it.

**Example:** How many ways can you distribute 7 cookies to 4 kids so that no kid gets more than 2 cookies?

**Solution:** To answer this, we will subtract all the outcomes in which a kid gets 3 or more cookies. How many outcomes are there like that? We can force kid A to eat 3 or more cookies by giving him 3 cookies before we start. Doing so reduces the problem to one in which we have 4 cookies to give to 4 kids without any restrictions. In that case, we have 4 stars (the 4 remaining cookies) and 3 bars (one less than the number of kids) so we can distribute the cookies in $\binom{7}{3}$ ways. Of course we could choose any one of the 4 kids to give too many cookies, so it would appear that there are $\binom{4}{1} \binom{7}{3}$ ways to distribute the cookies giving too many to one kid. But in fact, we have over counted. We must get rid of the outcomes in which two kids have too many cookies. There are $\binom{4}{2}$ ways to select to 2 kids to over-feed. It takes 6 cookies to do this, leaving only 1 cookie. So we have 1 star and still 3 bars. The remaining 1 cookie can thus be distributed in $\binom{4}{1}$ ways (for each of the $\binom{4}{2}$ choices of which 2 kids to over-feed). We could continue in this fashion (using PIE) but there are no ways to give too many cookies to 3 kids - you don’t have enough cookies.

All together we get that the number of ways to distribute 7 cookies to 4 kids without giving any kid more than 2 cookies is:

$$\binom{10}{3} - \left[ \binom{4}{1} \binom{7}{3} - \binom{4}{2} \binom{4}{3} \right] = 120 - [140 - 24] = 4$$

This makes sense: the only way to distribute cookies with this restriction is for 3 kids to get 2 cookies and one kid to get 1. There are 4 choices for which kid gets the 1 cookie.

**Example:** Earlier we counted the number of solutions to the equation

$$x_1 + x_2 + x_3 + x_4 + x_5 = 13.$$
How many of those solutions have $0 \leq x_i \leq 3$ for each $x_i$?

**Solution:** We must subtract off the number of solutions in which one or more of the variables has a value greater than 3. We will need to use PIE because counting the number of solutions for which each of the five variables separately are greater than 3 counts solutions multiple times. Here is what we get:

- Total solutions: \( \binom{17}{4} \)
- Solutions where $x_1 > 3$: \( \binom{13}{4} \) - give $x_1$ 4 units first, then distribute the remaining 9 units to the 5 variables.
- Solutions where $x_1 > 3$ and $x_2 > 3$: \( \binom{9}{4} \) - after you give 4 units to $x_1$ and another 4 to $x_2$, you only have 5 units left to distribute.
- Solutions where $x_1 > 3$, $x_2 > 3$ and $x_3 > 3$: \( \binom{5}{4} \).
- Solutions where $x_1 > 3$, $x_2 > 3$, $x_3 > 3$, and $x_4 > 3$: 0.

We also need to account for the fact that we could choose any of the five variables in the place of $x_1$ above, any pair of variables in the place of $x_1$ and $x_2$ and so on. It is because of this that the double counting occurs, so we need to use PIE. All together we have that the number of solutions with $0 \leq x_i \leq 3$ is

\[
\binom{17}{4} - \left[ \binom{5}{1} \binom{13}{4} - \binom{5}{2} \binom{9}{4} + \binom{5}{3} \binom{5}{4} \right]
\]

**Exercises**

1. Your wardrobe consists of 5 shirts, 3 pairs of pants, and 17 bow ties. How many different outfits can you make?

2. For your college interview, you must wear a tie. You own 3 regular (boring) ties and 5 (cool) bow ties. How many choices do you have for your neck-wear?

3. You realize that the interview is for clown-college, so you should probably wear both a regular tie and a bow tie. How many choices do you have now?

4. You realize that it would also be okay to wear more than two ties.
   
   (a) You must select some of your ties to wear - everything is okay, from no ties up to all ties. How many choices do you have?
   
   (b) If you want to wear at least one regular tie and one bow tie, but are willing to wear up to all your ties, how many choices do you have for which ties to wear?
   
   (c) How many choices do you have if you wear exactly 2 of the 3 regular ties and 3 of the 5 bow ties?
   
   (d) Once you have selected 2 regular and 3 bow ties, in how many orders could you put the ties on, assuming you must have one of the three bow ties on top?
5. Your Blu-ray collection consists of 9 comedies and 7 horror movies. Give an example of a question for which the answer is:

(a) 16.
(b) 63.

6. If \(|A| = 10\) and \(|B| = 15\), what is the largest possible value for \(|A \cap B|\)? What is the smallest? What are the possible values for \(|A \cup B|\)?

7. If \(|A| = 8\) and \(|B| = 5\), what is \(|A \cup B| + |A \cap B|\)?

8. A group of college students were asked about their TV watching habits. Of those surveyed, 28 students watch *Elementary*, 19 watch *Castle* and 24 watch *The Mentalist*. Additionally, 16 watch *Elementary* and *Castle*, 14 watch *elementary* and *The Mentalist* and 10 watch *Castle* and *The Mentalist*. There are 8 students who watch all three shows. How many students surveyed watched at least one of the shows?

9. Find \(|(A \cup C) \cap \overline{B}|\) provided \(|A| = 50\), \(|B| = 45\), \(|C| = 40\), \(|A \cap B| = 20\), \(|A \cap C| = 15\), \(|B \cap C| = 23\) and \(|A \cap B \cap C| = 12\).

10. Using the same data as the previous question, describe a set with cardinality 26.

11. Consider all 5 letter “words” made from the letters a through h.

(a) How many of these words are there total?
(b) How many of these words contain no repeated letters?
(c) How many of these words (repetitions allowed) start with the sub-word “aha”?
(d) How many of these words (repetitions allowed) either start with “aha” or end with “bah” or both?
(e) How many of the words containing no repeats also do not contain the sub-word “bad” (in consecutive letters)?

12. How many 10-bit strings contain 6 or more 1’s?

13. What is the coefficient of \(x^9\) in the expansion of \((x + 1)^{14} + x^3(x + 2)^{15}\)?

14. Let \(S = \{1, 2, 3, 4, 5, 6\}\)

(a) How many subsets are there total?
(b) How many subsets contain \(\{2, 3, 5\}\) as a subset?
(c) How many subsets of \(S\) contain no prime numbers?
(d) How many subsets contain at least one odd number?
(e) How many doubletons (i.e., subsets of two elements) contain exactly one even number?

15. How many shortest lattice paths start at (3,3) and
(a) end at (10,10)?
(b) end at (10,10) and pass through (5,7)?
(c) end at (10,10) and avoid (5,7)?

16. A pizza parlor offers 10 toppings.
(a) How many 3-topping pizzas could they put on their menu? Assume double toppings are not allowed.
(b) How many total pizzas are possible, with between zero and ten toppings (but not double toppings) allowed?
(c) The pizza parlor will list the 10 toppings in two columns on their menu. How many ways can they arrange the toppings in the left column?

17. How many quadrilaterals can you draw using the dots below as vertices (corners)?

\[
\begin{array}{cccccccc}
& & & & & & & \\
& & & & & & & \\
& & & & & & & \\
& & & & & & & \\
& & & & & & & \\
& & & & & & & \\
& & & & & & & \\
& & & & & & & \\
& & & & & & & \\
& & & & & & & \\
\end{array}
\]

18. How many of the quadrilaterals possible in the previous problem are:
(a) Squares?  (b) Rectangles?  (c) Parallelograms?  (d) Trapezoids?

19. On a business retreat, your company of 20 businessmen go golfing.
(a) You need to divide up into foursomes (groups of 4 people): a first foursome, a second foursome, and so on. How many ways can you do this?
(b) After all your hard work, you realize that in fact, you want each foursome to include one of the five CEO’s. How many ways can you do this?

20. How many different seating arrangements are possible for King Arthur and his 9 knights around their round table?

21. Give a combinatorial proof for the identity \(1 + 2 + 3 + \cdots + n = \binom{n+1}{2}\).

22. A woman is getting married. She has 15 best friends but can only select 6 of them to be her bride's maids, one of which needs to be her maid of honor. How many ways can she do this?
(a) What if she first selects the 6 bride's maids, and then selects one of them to be the maid of honor?
(b) What if she first selects her maid of honor, and then 5 other bride's maids?
(c) Explain why \(6 \binom{15}{6} = 15 \binom{14}{5}\).
23. Consider the bit strings in $B_2^6$ (bit strings of length 6 and weight 2).
   (a) How many of those bit strings start with 1?
   (b) How many of those bit strings start with 01?
   (c) How many of those bit strings start with 001?
   (d) Are there any other strings we have not counted yet? Which ones, and how many are there?
   (e) How many bit strings are there total in $B_2^6$?
   (f) What binomial identity have you just given a combinatorial proof for?

24. Let’s count ternary digit strings - strings in which each digit can be 0, 1, or 2.
   (a) How many ternary digit strings contain exactly $n$ digits?
   (b) How many ternary digit strings contain exactly $n$ digits and $n$ 2’s.
   (c) How many ternary digit strings contain exactly $n$ digits and $n - 1$ 2’s. (Hint: where can you put the non-2 digit, and then what could it be?)
   (d) How many ternary digit strings contain exactly $n$ digits and $n - 2$ 2’s. (Hint: see previous hint)
   (e) How many ternary digit strings contain exactly $n$ digits and $n - k$ 2’s.
   (f) How many ternary digit strings contain exactly $n$ digits and no 2’s. (Hint: what kind of a string is this?)
   (g) Use the above parts to give a combinatorial proof for the identity
   \[ \binom{n}{0} + 2\binom{n}{1} + 2^2\binom{n}{2} + 2^3\binom{n}{3} + \cdots + 2^n\binom{n}{n} = 3^n \]

25. Give a combinatorial proof for the identity $P(n, k) = \binom{n}{k}k!$

26. After gym class you are tasked with putting the 14 identical dodge-balls away into 5 bins.
   (a) How many ways can you do this if there are no restrictions?
   (b) How many ways can you do this if each bin must contain at least one dodge-ball?
   (c) How many ways can you do this if no bin can hold more than 6 balls?

27. How many integer solutions are there to the equation $x + y + z = 8$ for which
   (a) $x$, $y$, and $z$ are all positive?
   (b) $x$, $y$, and $z$ are all non-negative?
   (c) $x$, $y$, and $z$ are all greater than $-3$.

28. When playing Yahtzee, you roll five regular 6-sided dice. How many different outcomes are possible from a single roll? The order of the dice does not matter.
29. How many integer solutions to \( x_1 + x_2 + x_3 + x_4 = 25 \) are there for which \( x_1 \geq 1, x_2 \geq 2, x_3 \geq 3 \) and \( x_4 \geq 4 \)?

30. Write out all functions \( f : \{1,2,3\} \to \{a,b\} \). How many are there? How many are injective? How many are surjective? How many are both?

31. Write out all functions \( f : \{1,2\} \to \{a,b,c\} \). How many are there? How many are injective? How many are surjective? How many are both?

32. Consider the function \( f : \{1,2,3,4,5\} \to \{1,2,3,4\} \) given by the table below:

<table>
<thead>
<tr>
<th>( x )</th>
<th>1</th>
<th>2</th>
<th>3</th>
<th>4</th>
<th>5</th>
</tr>
</thead>
<tbody>
<tr>
<td>( f(x) )</td>
<td>3</td>
<td>2</td>
<td>4</td>
<td>1</td>
<td>2</td>
</tr>
</tbody>
</table>

(a) Is \( f \) injective? Explain.
(b) Is \( f \) surjective? Explain.

33. Consider the function \( f : \{1,2,3,4\} \to \{1,2,3,4\} \) given by the graph below.

(a) Is \( f \) injective? Explain.
(b) Is \( f \) surjective? Explain.

34. For each function given below, determine whether or not the function is injective and whether or not the function is surjective.

(a) \( f : \mathbb{N} \to \mathbb{N} \) given by \( f(n) = n + 4 \).
(b) \( f : \mathbb{Z} \to \mathbb{Z} \) given by \( f(n) = n + 4 \).
(c) \( f : \mathbb{Z} \to \mathbb{Z} \) given by \( f(n) = 5n - 8 \).
(d) \( f : \mathbb{Z} \to \mathbb{Z} \) given by \( f(n) = \begin{cases} n/2 & \text{if } n \text{ is even} \\ (n+1)/2 & \text{if } n \text{ is odd} \end{cases} \).

35. What can you deduce about the sets \( X \) and \( Y \) if you know,

(a) there is a injective function \( f : X \to Y \). Explain.
(b) there is a surjective function \( f : X \to Y \). Explain.
(c) there is a bijection \( f : X \to Y \). Explain.
36. Suppose $f : X \to Y$ is a function. Which of the following are possible? Explain.
   (a) $f$ is injective but not surjective.
   (b) $f$ is surjective but not injective.
   (c) $|X| = |Y|$ and $f$ is injective but not surjective.
   (d) $|X| = |Y|$ and $f$ is surjective but not injective.
   (e) $|X| = |Y|$, $X$ and $Y$ are finite, and $f$ is injective but not surjective.
   (f) $|X| = |Y|$, $X$ and $Y$ are finite, and $f$ is surjective but not injective.

37. Consider functions $f : \{1, 2, 3, 4\} \to \{a, b, c, d, e, f\}$.
   (a) How many functions are there total?
   (b) How many functions are injective?
   (c) How many functions are surjective?
   (d) How many functions have the property that $f(1) \neq a$ or $f(2) \neq b$, or both?

38. Consider sets $A$ and $B$ with $|A| = 10$ and $|B| = 17$.
   (a) How many functions $f : A \to B$ are there?
   (b) How many functions $f : A \to B$ are injective?

39. Consider sets $A$ and $B$ with $|A| = 10$ and $|B| = 5$. How many functions $f : A \to B$ are surjective?

40. Let $A = \{1, 2, 3, 4, 5\}$. How many injective functions $f : A \to A$ have the property that for each $x \in A$, $f(x) \neq x$?

41. Ten ladies of a certain age drop off their red hats at the hat check of a museum. As they are leaving, the hat check attendant gives the hats back randomly. In how many ways can exactly six of the ladies receive their own hat (and the other four not)?

42. You have 9 presents to give to your 4 kids. How many ways can this be done if
   (a) The presents are identical, and each kid gets at least one present?
   (b) The presents are identical, and some kids might get no presents?
   (c) The presents are unique, and some kids might get no presents?
   (d) the presents are unique and each kid gets at least one present?
43. For each of the following counting problems, say whether the answer is

A \( \binom{10}{4} \)  \hspace{1cm} B \ P(10, 4) \hspace{1cm} C \text{ Neither}

If you answer is “Neither,” say what the answer should be instead.

(a) How many shortest lattice paths are there from (0, 0) to (10, 4)?
(b) If you have 10 bow ties, and you want to select 4 of them for next week, how many choices do you have?
(c) Suppose you have 10 bow ties and you will wear one on each of the next 4 days. How many choices do you have?
(d) If you want to wear 4 of your 10 bow ties next week (Monday through Sunday), how many ways can this be accomplished?
(e) Out of a group of 10 classmates, how many ways can you rank your top 4 friends?
(f) If 10 students come to their professor’s office but only 4 can fit at a time, how different combinations of 4 students can see the prof first?
(g) How many 4 letter words can be made from the first 10 letters of the alphabet?
(h) How many ways can you make the word “cake” from the first 10 letters of the alphabet?
(i) How many ways are there to distribute 10 apples among 4 children?
(j) If you have 10 kids (and live in a shoe) and 4 types of cereal, how many ways can your kids eat breakfast?
(k) How many ways can you arrange exactly 4 ones in a string of 10 binary digits?
(l) You want to select 4 single digit numbers as your lotto picks. How many choices do you have?
(m) 10 kids want ice-cream. You have 4 varieties. How many ways are there to give the kids as much ice-cream as they want?
(n) How many 1-1 functions are there from \{1, 2, \ldots, 10\} to \{a, b, c, d\}?
(o) How many surjective functions are there from \{1, 2, \ldots, 10\} to \{a, b, c, d\}?
(p) Each of your 10 bow ties match 4 pairs of suspenders. How many outfits can you make?
(q) After the party, the 10 kids each choose one of 4 party-favors. How many outcomes?
(r) How many 6-elements subsets are there of the set \{1, 2, \ldots, 10\}
(s) How many ways can you split up 11 kids into 5 teams?
(t) How many solutions are there to \[ x_1 + x_2 + \cdots + x_5 = 6 \] where each \( x_i \) is non-negative?
(u) Your band goes on tour. There are 10 cities within driving distance, but only enough time to play 4 of them. How many choices do you have for the cities on your tour?
(v) In how many different ways can you play the 4 cities you choose?

(w) Out of the 10 breakfast cereals available, you want to have 4 bowls. How many ways can you do this?

(x) There are 10 types of cookies available. You want to make a 4 cookie stack. How many different stacks can you make?

(y) From you home at (0,0) you want to go to either the donut shop at (5,4) or the one at (3,6). How many paths could you take?

(z) How many 10-digit numbers do not contain a sub-string of 4 repeated digits?
Chapter 4

Sequences

As the legend goes, there is a monastery in Hanoi with a great hall containing 3 tall pillars. Resting on the first pillar are 64 giant disks (or washers) - all different sizes, stacked from largest to smallest. The monks are charged with the following task: they must move the entire stack of disks to the third pillar. However, due to the size of the disks, the monks cannot move more than one at a time. Each disk must be placed on one of the pillars before the next disk is moved. And because the disks are so heavy and fragile, the monks may never place a larger disk on top of a smaller disk.

When the monks finally complete their task, the world shall come to an end. Your task: figure out how long before we need to start worrying about the end of the world.

This puzzle is called the Tower of Hanoi. You are tasked with finding the minimum number of moves to complete the puzzle. This certainly sounds like a counting problem. Perhaps you have an answer? If not, what else could we try? The answer depends on the number of disks you need to move. In fact, we could answer the puzzle first for 1 disk, then 2, then 3 and so on. If we list out all of the answers for each number of disks, we will get a sequence of numbers. The $n$th term in the sequence is the answer to the question, “what is the smallest number of moves required to complete the Tower of Hanoi puzzle with $n$ disks?” You might wonder why we would create such a sequence, instead of just answering the question. The reason is, but looking at how the sequence of numbers grows, we gain insight into the problem. It is easy to count the number of moves required for small numbers of disks. We can then look for a pattern among the first few terms of the sequence. Hopefully this will suggest a method for finding the $n$th term - the answer to our question. Of course we will also need to verify that our suspected pattern is correct, and that this correct pattern really does give us the $n$th term we think it does, but it is impossible to prove that your formula is correct without having a formula to start with.

Of course sequences are also interesting mathematical objects to study in their own right. Let’s see why.

4.1 Defining Sequences

A sequence is simply an ordered list of numbers. For example, here is a sequence: 0, 1, 2, 3, 4, 5, . . .. This is different from the set $\mathbb{N}$, because while the sequence is a complete list of every element in the set of natural numbers, in the sequence, we very much care what order the numbers come in. For this reason, when we use variables to represent terms in a sequence, they will look like this:

$$a_0, a_1, a_2, a_3, \ldots$$

We might replace the $a$ with another letter, and sometimes we omit $a_0$, starting with $a_1$ instead. The numbers in the subscripts are called indices (the plural of index). We can think
of the terms in the sequence as the outputs of a function with domain \( \mathbb{N} \): the \( n \)th term in the sequence is \( f(n) = a_n \).

**Example:** Can you find the next term in the following sequences?

1. 7, 7, 7, 7, 7, …
2. 3, -3, 3, -3, 3, …
3. 1, 5, 2, 10, 3, 15, …
4. 1, 2, 4, 8, 16, 32, …
5. 1, 4, 9, 16, 25, 36, …
6. 1, 2, 3, 5, 8, 13, 21, …
7. 1, 3, 6, 10, 15, 21, …
8. 2, 3, 5, 7, 11, 13, …
9. 3, 2, 1, 0, -1, …
10. 1, 1, 2, 6, …

**Solution:** No you cannot. You might guess that the next terms are:

1. 7
2. -3
3. 4
4. 64
5. 49
6. 28
7. 24
8. 17
9. -2
10. 24

In fact, those are the next terms of the sequences I had in mind when I made up the example. But there is no way to be sure they are correct.

That said, we will often do this. Given the first few terms of a sequence, we can ask what the pattern in the sequence suggests the next terms are.

Given that no number of initial terms in a sequence is enough to say for certain which sequence we are dealing with, we need to find another way to specify a sequence. There are two ways to do this.

**Closed formula**
A *closed formula* for a sequence \( a_0, a_1, a_2, \ldots \) is a formula for \( a_n \) using a fixed finite number of operations on \( n \). This is what you normally think of as a formula in \( n \).

**Recursive definition**
A *recursive definition* (sometimes called an *inductive definition*) for a sequence \( a_0, a_1, a_2, \ldots \) consists of a *recurrence relation*: an equation relating a term of the sequence to previous terms (terms with smaller index) and an *initial condition*: a list of a few terms of the sequence (one less than the number of terms in the recurrence relation).

It is easier to understand what is going on here with an example:

**Example:** Here are a few closed formulas for sequences:

- \( a_n = n^2 \)
Sequences 4.1. Defining Sequences

- \( a_n = \frac{n(n+1)}{2} \)
- \( a_n = \frac{(\frac{1+\sqrt{5}}{2})^n - 1}{\frac{1+\sqrt{5}}{2}} \).

Note in each case, if you are given \( n \), you can calculate \( a_n \) directly - just plug in \( n \).

Here are a few recursive definitions for sequences:

- \( a_n = 2a_{n-1} \) with \( a_0 = 1 \).
- \( a_n = 2a_{n-1} \) with \( a_0 = 27 \).
- \( a_n = a_{n-1} + a_{n-2} \) with \( a_0 = 0 \) and \( a_1 = 1 \).

In these cases, if you are given \( n \), you cannot calculate \( a_n \) directly, you first need to find \( a_{n-1} \) or \( a_{n-1} \) and \( a_{n-2} \).

You might wonder why we would bother with recursive definitions for sequences - after all, it is harder to find \( a_n \) with a recursive definition than with a closed formula. This is true, but it is also harder to find a closed formula for a sequence than it is to find a recursive definition. So to find a useful closed formula, we might first find the recursive definition, then use that to find the closed formula.

This is not to say that recursive definitions aren’t useful in finding \( a_n \). You can always calculate \( a_n \) given a recursive definition - it might just take a while.

**Example:** Find \( a_6 \) in the sequence defined by \( a_n = 2a_{n-1} - a_{n-2} \) with \( a_0 = 3 \) and \( a_1 = 4 \).

**Solution:** We know that \( a_6 = 2a_5 - a_4 \). So to find \( a_6 \) we need to find \( a_5 \) and \( a_4 \). Well

\[
a_5 = 2a_4 - a_3 \quad \text{and} \quad a_4 = 2a_3 - a_2
\]

so if we can only find \( a_3 \) and \( a_2 \) we would be set. Of course

\[
a_3 = 2a_2 - a_1 \quad \text{and} \quad a_2 = 2a_1 - a_0
\]

so we only need to find \( a_1 \) and \( a_0 \). But we are given these. Thus

\[
a_0 = 3 \\
a_1 = 4 \\
a_2 = 2 \cdot 4 - 3 = 5 \\
a_3 = 2 \cdot 5 - 4 = 6 \\
a_4 = 2 \cdot 6 - 5 = 7 \\
a_5 = 2 \cdot 7 - 6 = 8 \\
a_6 = 2 \cdot 8 - 7 = 9.
\]

Note that now we can guess a closed formula for the \( n \)th term of the sequence: \( a_n = n + 3 \). To be sure this will always work, we could plug in this formula into the recurrence relation:

\[
2a_{n-1} - a_{n-2} = 2((n - 1) + 3) - ((n - 2) + 3) = 2n + 4 - n - 1 = n + 3 = a_n
\]
Since \(a_0 = 0 + 3 = 3\) and \(a_1 = 1 + 3 = 4\) are the correct initial conditions, we have lucked upon the correct closed formula.

Finding closed formulas, or even recursive definitions, for sequences is not trivial - there is no one method for doing this. Just like in evaluating integrals or solving differential equations, it is useful to have a bag of tricks you can apply, but sometimes there is no easy answer.

One useful trick to keep in your bag is relating a given sequence to another sequence for which we already know the closed formula.

**Example:** Use the formulas \(T_n = \frac{n(n+1)}{2}\) and \(a_n = 2^n\) to find closed formulas for the following sequences. Assume the first term has index 1 (not 0).

1. 2, 4, 7, 11, 16, 22, . . .
2. 2, 3, 5, 9, 17, 33, . . .
3. 2, 6, 12, 20, 30, 42, . . .
4. 6, 10, 15, 21, 28, . . .
5. 1, 3, 7, 15, 31, . . .
6. 3, 6, 12, 24, 48, . . .
7. 6, 10, 18, 34, 66, . . .
8. 15, 33, 57, 87, 123, . . .

**Solution:** Before you say this is impossible, what we are asking for is simply to find a closed formula which agrees with all of the initial terms of the sequences. Of course there is no way to read into the mind of the person who wrote the numbers down, but we can at least do this.

Now the first few terms of \(T_n\), starting with \(T_1\) are 1, 3, 6, 10, 15, 21, . . . (these are the triangular numbers). The first few terms of \(a_n\) (starting this time with \(a_0\)) are 1, 2, 4, 8, 16, . . . Now let’s try to find formulas for the given sequences.

1. 2, 4, 7, 11, 16, 22, . . . - Note that if subtract 1 from each term, we get the sequence \(T_n\). So this sequence is \(T_n + 1\). Therefore a closed formula is \(\frac{n(n+1)}{2} + 1\). A quick check of the first few \(n\) confirms we have it right.
2. 2, 3, 5, 9, 17, 33, . . . - This sequence is the result of adding 1 to each term in \(a_n\). So we might guess the closed formula \(2^n + 1\). If we try this though, we get the first terms \(2^1 + 1 = 3\) and \(2^2 + 1 = 5\). We are off, because \(a_n\) started with \(n = 0\), and now we are starting with \(n = 1\). So we shift: the closed formula is \(2^{n-1} + 1\).
3. 2, 6, 12, 20, 30, 42, . . . - Notice that all these terms are even. What happens if we factor out a 2? We get \(T_n\)! So this sequence has closed formula \(n(n + 1)\).
4.6, 10, 15, 21, 28, . . . - These are all triangular numbers. However, we are starting with 6 as our first term instead of as our third term. So if we could plug in 3 instead of 1 into the formula for $T_n$, we would be set. Therefore the closed formula is $\frac{(n+2)(n+3)}{2}$ (where $n + 3$ came from $(n + 2) + 1$).

5. 1, 3, 7, 15, 31, . . . - Try adding one to each term and we get powers of 2. You might guess this because each term is approximately twice the previous term. Closed formula: $2^n - 1$.

6. 3, 6, 12, 24, 48, . . . - These numbers are all multiples of 3. Let’s try dividing each by 3. Doing so gives 1, 2, 4, 8, . . . Aha. We get the closed formula $3 \cdot 2^{n-1}$.

7. 6, 10, 18, 34, 66, . . . - To get from one term to the next, we almost double each term. So maybe we can relate this back to $2^n$. Yes, each term is 2 more than a power of 2. So we get $2^{n+1} + 2$ (the $n + 1$ is because the first term is 2 more than $2^1$, not $2^1$). Alternatively, we could have related this sequence to the second sequence in this example: starting with 3, 5, 9, 17 . . . we see that this sequence is twice the terms from that sequence (omitting the 2). That sequence had closed formula $2^{n-1} + 1$. To make the 3 first, we would write $2^n + 1$. Our sequence would be twice this, so $2(2^n + 1)$, which of course is the same as we got before.

8. 15, 33, 57, 87, 123, . . . - Try dividing each term by 3. That gives the sequence 5, 11, 19, 29, 41, . . . Now add one: 12, 20, 30, 42, . . . - which is sequence 3 in this example, except starting with 6 instead of 2. So let’s start with the formula for sequence 3: $n(n + 1)$. To start with the 6, we shift: $(n + 1)(n + 2)$. But this is one too many, so subtract 1: $(n + 1)(n + 2) - 1$. That gives us our sequence, but divided by 3. So we want $3((n + 1)(n + 2) - 1)$.

4.2 Arithmetic and Geometric Sequences

We now turn to the question of finding closed formulas for particular types of sequences.

**Arithmetic Sequences**

If the terms of a sequence differ by a constant, we say the sequence is **arithmetic**.

If the first term ($a_1$) of the sequence is $a$ and the common difference is $d$, then we have,

- Recursive definition: $a_n = a_{n-1} + d$ with $a_1 = a$.
- Closed formula: $a_n = a + d(n - 1)$.

How do we know this? For the recursive definition, we need to specify $a_1$. Then we need to express $a_n$ in terms of $a_{n-1}$. If we call the first term $a$, then $a_1 = a$. For the recurrence
relation, by the definition of an arithmetic sequence, the difference between successive terms is some constant, say $d$. So $a_n - a_{n-1} = d$, or in other words,

$$a_1 = a \quad a_n = a_{n-1} + d$$

To find a closed formula, first try writing out the first write out the sequence using the recursive definition without simplifying:

$$a, a+d, a+d+d, a+d+d+d, \ldots$$

We see that to find the $n$th term, we need to start with $a$ and then add $d$ a bunch of times. In fact, add it $n-1$ times. Thus $a_n = a + d(n-1)$.

**Example:** Find recursive definitions and closed formulas for the sequences:

1. 2, 5, 8, 11, 14, ...
2. 50, 43, 36, 29, ...

**Solution:** First we should check that these sequences really are arithmetic. Doing so will reveal the common difference $d$.

1. $5 - 2 = 3, \ 8 - 5 = 3$, etc. To get from each term to the next, we add three, so $d = 3$. The recursive definition is therefore $a_n = a_{n-1} + 3$ with $a_1 = 2$. The closed formula is $a_n = 2 + 3(n-1)$.

2. Here the common difference is $-7$, since we add $-7$ to 50 to get 43, and so on. Thus we have a recursive definition of $a_n = a_{n-1} - 7$ with $a_1 = 50$. The closed formula is $a_n = 50 - 7(n-1)$.

What about sequences like 2, 6, 18, 54, ...? This is not arithmetic, because the difference between terms is not constant. However, the ratio between successive terms is constant. We call such sequences geometric.

The recursive definition for the geometric progression with first term $a$ and common ratio $r$ is

$$a_1 = a \quad a_n = a_{n-1} \cdot r$$

This makes sense. To get the next term we multiply the previous term by $r$. We can find the closed formula like we did for the arithmetic progression. Write $a_1 = a, \ a_2 = a \cdot r, \ a_3 = a_2 \cdot r = a \cdot r \cdot r$ and so on. We must multiply the first term $a$ by $r$ a number of times - $n-1$ to be precise. We get $a_n = a \cdot r^{n-1}$.

**Geometric Sequences**

A sequence is called geometric if the ratio between successive terms is constant. Suppose the first term $a_1$ is $a$ and the common ratio is $r$. Then we have,

- **Recursive definition:** $a_n = r a_{n-1}$ with $a_1 = a$.
- **Closed formula:** $a_n = a \cdot r^{n-1}$. 

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Example: Find the recursive and closed formula for the sequences:

1. 3, 6, 12, 24, 48, \ldots
2. 27, 9, 3, 1, 1/3, \ldots

Solution: Again, we should first check that these sequences really are geometric - but doing so will tell us what \( r \) is.

1. \( 6/3 = 2, 12/6 = 2, 24/12 = 2, \) etc. Yes, to get from any term to the next, we multiply by \( r = 2 \). So the recursive definition is \( a_n = 2a_{n-1} \) with \( a_1 = 3 \). The closed formula is \( a_n = 3 \cdot 2^{n-1} \).

2. \( 9/27 = 1/3 \) and \( 3/9 = 1/3 \) and so on. The common ratio is \( r = 1/3 \). So the sequence has recursive definition \( a_n = \frac{1}{3}a_{n-1} \) with \( a_1 = 27 \) and closed formula \( a_n = 27 \cdot \frac{1}{3}^{n-1} \).

### 4.2.1 Sums of Arithmetic and Geometric Sequences

Look at the sequence 1, 3, 6, 10, 15, \ldots. We have a formula for this already (they are the triangular numbers) but let’s see if we can derive the formula anew. First, is this sequence arithmetic? No, since \( 3 - 1 = 2 \) and \( 6 - 3 = 3 \neq 2 \), so there is no common difference. Is the sequence geometric? No. \( 3/1 = 3 \) but \( 6/3 = 2 \), so there is no common ratio. What to do?

Notice thought that the differences between terms are a arithmetic sequence: 2, 3, 4, 5, 6, \ldots. This says that the \( n \)th term of the sequence 1, 3, 5, 10, 15, \ldots is the sum of the first \( n \) terms in the sequence 1, 2, 3, 4, 5, \ldots. If we know how to add up the terms of an arithmetic sequence, we could use this to find a closed formula for a sequence whose differences are the terms of that arithmetic sequence. In fact, this is how we found the formula \( \frac{n(n+1)}{2} \) for the triangular numbers.

We could use a similar technique to find a closed formula for the sequence 2, 3, 5, 9, 17, \ldots - the differences are 1, 2, 4, 8, \ldots, a geometric sequence. If we had a method for summing geometric sequences, we could get a formula for sequences with geometric differences.

**Summing Arithmetic Sequences: Reverse and Add**

Here is a trick that allows us to quickly find the sum of an arithmetic sequence.

**Example:** Find the sum: 2 + 5 + 8 + 11 + 14 + \cdots + 470.

**Solution:** The idea is to mimic the trick we used to find the formula for triangular numbers. If we add the first and last, we get 472. The second term and second-to-last term also add up to 472. To make this clearer, we might express this as follows. Call the sum \( S \). Then,

\[
S = 2 + 5 + 8 + \cdots + 467 + 470
\]

\[
+ S = 470 + 467 + 464 + \cdots + 5 + 2
\]

\[
2S = 472 + 472 + 472 + \cdots + 472 + 472
\]

To find \( 2S \) then we add 472 to itself a number of times. What number? We need to decide how many terms (summands) are in the sum. The \( n \)th term in
the sum can be expressed as $2 + 3(n - 1)$ - this is an arithmetic sequence after all. We want to find $n$ when $2 + 3(n - 1) = 470$. Solving for $n$ gives $n = 157$. So $n$ ranges from 1 to 157, giving 157 terms in the sum. This is the number of 472’s in the sum for $2S$. Thus

$$2S = 157 \cdot 472 = 74104$$

It is now easy to find $S$:

$$S = \frac{74104}{2} = 37052$$

This will work in general, for any sum of arithmetic sequences. Call the sum $S$. Reverse and add. This produces a single number added to itself many times. Find the number of times. Multiply. Divided by 2. Done.

**Example:** Find a closed formula for $6 + 10 + 14 + \cdots + (4n - 2)$.

**Solution:** Again, we have a sum of an arithmetic sequence. We need to know how many terms are in the sequence. Clearly each term in the sequence has the form $4k - 2$ (as evidenced by the last term). For which values of $k$ though? To get 6, $k = 2$. To get $4n - 2$ take $k = n$. So to find the number of terms, we need to know how many integers are in the range 2, 3, \ldots, $n$. The answer is $n - 1$. (There are $n$ numbers from 1 to $n$, so one less if we start with 2.)

Now use the reverse and add trick:

$$S = \begin{array}{ccccccc}
6 & + & 10 & + & 14 & + & \cdots & + & 4n - 6 & + & 4n - 2 \\
+ & S = & 4n - 2 & + & 4n - 6 & + & 4n - 10 & + & \cdots & + & 10 & + & 6
\end{array}$$

Since there are $n - 2$ terms, we get

$$2S = (n - 2)(4n + 4)$$ 

so

$$S = \frac{(n - 2)(4n + 4)}{2}$$

**Summing Geometric Sequences: Multiply, Shift and Subtract**

To find the sum of a geometric sequence, we cannot use the previous trick. Do you see why? The reason we got the same term added to itself many times is because there was a constant difference. So as we added that difference in one direction, we subtracted the difference going the other way, leaving a constant total. For geometric sums, we have a different trick.

**Example:** What is $3 + 6 + 12 + 24 + \cdots + 12288$?

**Solution:** Multiply each term by the common ratio (2). You get $2S = 6 + 12 + 24 + \cdots + 24576$. Now subtract: $2S - S = -3 + 24576 = 24573$. Since $2S - S = S$, we have our answer.

To see what happened in the above example, try writing it this way:
Sequences 4.2. Arithmetic and Geometric Sequences

\[
S = 3 + 6 + 12 + 24 + \cdots + 12288 \\
- 2S = 6 + 12 + 24 + \cdots + 12288 + 24576 \\
-S = 3 + 0 + 0 + 0 + \cdots + 0 - 24576
\]

Then divide both sides by \(-1\) and we have the same result for \(S\). The idea is, by multiplying the sum by the common ratio, each term becomes the next term. We shift over the sum to get the subtraction to mostly cancel out, leaving just the first term and new last term.

**Example:** Find a closed formula for \(S(n) = 2 + 10 + 50 + \cdots + 2 \cdot 5^n\).

**Solution:** The common ratio is 5. So we have

\[
S = 2 + 10 + 50 + \cdots + 2 \cdot 5^n \\
- 5S = 10 + 50 + \cdots + 2 \cdot 5^n + 2 \cdot 5^{n+1} \\
-4S = 2 - 2 \cdot 5^{n+1} \\
\]

Thus \(S = \frac{2 - 2 \cdot 5^{n+1}}{-4}\)

Even though this might seem like a new technique, you have probably used it before.

**Example:** Express 0.464646\ldots as a fraction.

**Solution:** Let \(N = 0.46464646\ldots\). Consider 100\(N\). We get:

\[
N = 0.46464646\ldots \\
-0.01N = 0.00464646\ldots \\
0.99N = 0.46 \\
\]

So \(N = \frac{46}{99}\). What have we done? We viewed the repeating decimal 0.464646\ldots as a sum of the geometric sequence 0.46, 0.0046, 0.000046,\ldots. The common ratio is .01. The only real difference is that we are now computing an *infinite* geometric sum, we we do not have the extra term on the end to consider.

**\(\sum\) and \(\prod\) notation**

To simplify writing out sums, we will use notation like \(\sum_{k=1}^{n} a_k\). This means add up the \(a_k\)s where \(k\) changes from 1 to \(n\).

**Example:** Use \(\sum\) notation to rewrite the sums:

1. \(1 + 2 + 3 + 4 + \cdots + 100\)
2. \(1 + 2 + 4 + 8 + \cdots + 2^{50}\)
3. \(6 + 10 + 14 + \cdots + (4n - 2)\)

**Solution:**
If we want to multiply the $a_k$ instead, we would write $\prod_{k=1}^{n} a_k$. For example, $\prod_{k=1}^{n} k = n!$.

### 4.3 Polynomial Fitting

So far we have seen methods for finding the closed formulas for certain types of sequences - arithmetic and geometric. Since we know how to compute the sum of the first $n$ terms of arithmetic and geometric sequences, we can compute the closed formulas for sequences which have an arithmetic (or geometric) sequence of differences between terms. But what if we consider a sequence which is the sum of the first $n$ terms of a sequence which is the sum of an arithmetic sequence!

Before we get too carried away, let’s consider an example: How many squares (of all sizes) are there on a chessboard? A chessboard consists of 64 squares, but we also want to consider squares of longer side length. Even though we are only considering an $8 \times 8$ board, there is already a lot to count. So instead, let us build a sequence: the first term will be the number of squares on a $1 \times 1$ board, the second term will be the number of squares on a $2 \times 2$ board, and so on. After a little thought, we arrive at the sequence

$$1, 5, 14, 30, 55, \ldots$$

This sequence is not arithmetic (or geometric for that matter), but perhaps it’s sequence of differences is. For differences we get

$$4, 9, 16, 25, \ldots$$

Not a huge surprise: one way to count the number of squares in a $4 \times 4$ chessboard is to notice that there are 16 squares with side length 1, 9 with side length 2, 4 with side length 3 and 1 with side length 4. So the original sequence is just the sum of squares. Now this sequence of differences is not arithmetic since it’s sequence of differences (the differences of the differences of the original sequence) is not constant. In fact, this sequence of *second differences* is

$$5, 7, 9, \ldots$$

which is an arithmetic sequence - *its* sequence of differences is constant (2, 2, 2, \ldots). Notice that our original sequence had *third differences* (that is, differences of differences of differences of the original) constant. We will call such a sequence $\Delta^3$. The sequence 1, 4, 9, 16, \ldots has second differences constant, so it will be a $\Delta^2$ sequence. In general, we will say a sequence is a $\Delta^k$ sequence if the $k$th differences are constant.

Now $\Delta^0$ sequences are constant, so a closed formula for them is easy to compute (it’s just the constant). $\Delta^1$ sequences are arithmetic - we have a method for finding closed formulas for them as well. $\Delta^2$ sequences are sums of arithmetic sequences - we can find formulas for these as well. But notice that the format of the closed formula for a $\Delta^2$ sequence is always
quadratic - the squares are $\Delta^2$ with closed formula $a_n = n^2$, the triangular numbers (also $\Delta^2$) have closed formula $a_n = \frac{n(n+1)}{2}$, which when multiplied out gives you an $n^2$ term as well. It appears that every time we increase the complexity of the sequence - that is, increase the number of differences before we get constants - we also increase the degree of the polynomial used for the closed formula. We go from constant to linear to quadratic. This makes sense. The sequence of differences between terms tells us something about the rate of growth of the sequence. If a sequence is growing at a constant rate, then the formula for the sequence will be linear. If the sequence is growing at a rate which itself is growing at a constant rate, then the formula is quadratic. You have seen this elsewhere - if a function has a constant second derivative (rate of change) then the function must be quadratic.

This works in general:

**Finite Differences**
The closed formula for a sequence will be a degree $k$ polynomial if and only if the sequence is $\Delta^k$ - that is, the $k$th sequence of differences is constant.

This tells us that the sequence 1, 5, 14, 30, 55, ... will have a cubic (degree 3 polynomial) for its closed formula.

Now once we know what format the closed formula for a sequence will take, it is much easier to actually find the closed formula. In the case that the closed formula is a degree $k$ polynomial, we just need $k + 1$ data points to “fit” the polynomial to the data.

**Example:** Find a formula for the sequence 3, 7, 14, 24, .... Assume $a_1 = 3$.

**Solution:** First, check to see if the formula has constant differences at some level. The sequence of first differences is 4, 7, 10, ... which is arithmetic, so the sequence of second differences is constant. The sequence is $\Delta^2$, so the formula for $a_n$ will be a degree 2 polynomial. That is, we know that for some constants $a$, $b$, and $c$,

$$a_n = an^2 + bn + c$$

Now to find $a$, $b$, and $c$. First, it would be nice to know what $a_0$ is, since plugging in $n = 0$ simplifies the above formula greatly. In this case, $a_0 = 2$ (work backwards from the sequence of constant differences). Thus

$$a_0 = 2 = a \cdot 0^2 + b \cdot 0 + c$$

so $c = 2$. Now plug in $n = 1$ and $n = 2$. We get

$$a_1 = 3 = a + b + 2$$

$$a_2 = 7 = a4 + b2 + 2.$$ 

At this point we have two (linear) equations and two unknowns, so we can solve the system for $a$ and $b$. We find $a = \frac{3}{2}$ and $b = \frac{-1}{2}$, so $a_n = \frac{3}{2}n^2 - \frac{1}{2}n + 2.$
Example: Find a closed formula for the number of squares on an $n \times n$ chessboard.

Solution: We have seen that the sequence $1, 5, 14, 30, 55, \ldots$ is $\Delta^3$, so we are looking for a degree 3 polynomial. That is,

$$a_n = an^3 + bn^2 + cn + d$$

We can find $d$ if we know what $a_0$ is. Working backwards from the third differences, we find $a_0 = 0$ (which makes sense - there are no squares on a $0 \times 0$ chessboard). Thus $d = 0$. Now plug in $n = 1, n = 2,$ and $n = 3$:

$$1 = a + b + c$$
$$5 = 8a + 4b + 2c$$
$$14 = 27a + 9b + 3c$$

If we solve this system of equations (using elimination, or an augmented matrix, or a computer) we get $a = \frac{1}{3}, b = \frac{1}{2}$ and $c = \frac{1}{6}$. Therefore the number of squares on an $n \times n$ chessboard is $a_n = \frac{1}{3}n^3 + \frac{1}{2}n^2 + \frac{1}{6}n$. Note: since the squares on a chessboard problem really is asking for the sum of squares, we now have a nice formula for $\sum_{k=1}^{n} k^2$.

Not all sequences will have polynomials as their closed formula. That is, this process doesn’t always work.

Example: Determine whether the following sequences can be described by a polynomial, and if so, of what degree.

1. $1, 2, 4, 8, 16, \ldots$
2. $0, 7, 50, 183, 484, 1055, \ldots$
3. $1, 1, 2, 3, 5, 8, 13, \ldots$

Solution:

1. The sequence of first differences is $1, 2, 4, 8, 16, \ldots$. This is identical to the original sequence, so taking any additional differences will not give anything different either. So there is no number of differences you could take to get constants, so the sequence is not $\Delta^k$ for any $k$. Therefore the closed formula for the sequence is not a polynomial. (In fact, we know the closed formula is $a_n = 2^{n-1}$, not a polynomial.)

2. First differences: $7, 43, 133, 301, 571, \ldots$. Second differences: $36, 90, 168, 270, \ldots$.
   Third difference: $54, 78, 102, \ldots$. Fourth differences: $24, 24, \ldots$ - constant.
   Thus the sequence is $\Delta^4$, so the closed formula is a degree 4 polynomial.

3. This is the Fibonacci sequence. The first differences are $0, 1, 1, 2, 3, 5, 8, \ldots$; the second differences are $1, 0, 1, 1, 2, 3, 5 \ldots$ - we notice that after the first
Sequences 4.4. Solving Recurrence Relations

few terms, we get the original sequence back. So there will never be constant differences, so the closed formula for the Fibonacci sequence is not a polynomial.

4.4 Solving Recurrence Relations

We have seen that it is often easier to find recursive definitions than closed formulas. Lucky for us, there are a few techniques for converting recursive definitions to closed formulas. Doing so is called solving a recurrence relation. Recall that the recurrence relation is the part of a recursive definition besides the initial conditions. For example, the recurrence relation for the Fibonacci sequence is $F_n = F_{n-1} + F_{n-2}$. (This, together with the initial conditions $F_0 = 0$ and $F_1 = 1$ give the entire recursive definition for the sequence.)

Example: Find a recurrence relation and initial conditions for the sequence $1, 5, 17, 53, 161, 485, \ldots$

Solution: Finding the recurrence relation would be easier if we had some context for the problem (like the Tower of Hanoi, for example). Alas, we have only the sequence. Remember, the recurrence relation tells you how to get from previous terms to future terms. What is going on here? We could look at the differences between terms: $4, 12, 36, 108, \ldots$ Notice that these are growing by a factor of 3. Is the original sequence as well? $1 \cdot 3 = 3, 5 \cdot 3 = 15, 17 \cdot 3 = 51$ and so on. It appears that we always end up with 2 less than the next term. Aha!

So $a_n = 3a_{n-1} + 2$ is our recurrence relation. The initial condition is $a_1 = 1$.

We are going to try to solve the recurrence relations. By this we mean something very similar to solving differential equations: we want to find a function of $n$ (a closed formula) which satisfies the recurrence relation, as well as the initial condition. (Note: recurrence relations are sometimes called difference equations since they can describe the difference between terms - this highlights the relation to differential equations further.) Just like for differential equations, finding a solution might be tricky, but checking that the solution is correct is easy - just plug it in.

Example: Check that $a_n = 2^n + 1$ is a solution to the recurrence relation $a_n = 2a_{n-1} - 1$ with $a_1 = 3$.

Solution: First, it is easy to check the initial condition: $a_1$ should be $2^1 + 1$ according to our closed formula. But $2^1 + 1 = 3$, which is what we want. To check that our proposed solution satisfies the the recurrence relation, try plugging it in.

$$2a_{n-1} - 1 = 2(2^{n-1} + 1) - 1 = 2^n + 2 - 1 = 2^n + 1 = a_n$$

That’s what our recurrence relation says! We have a solution.
We will consider three techniques for solving recurrence relations: Telescoping, Iteration and the Characteristic Root Technique.

4.4.1 Telescoping

Telescoping refers to the phenomenon when many terms in a large sum cancel out - so the sum “telescopes.” For example:

\[(2 - 1) + (3 - 2) + (3 - 4) + \cdots + (100 - 99) + (101 - 100) = -1 + 101\]

because every third term looks like: \(2 + -2 = 0\), the \(3 + -3 = 0\) and so on.

We can use this behavior to solve recurrence relations. Here is an example.

**Example:** Solve the recurrence relation \(a_n = a_{n-1} + n\) with initial term \(a_0 = 4\).

**Solution:** To get a feel for the recurrence relation, write out the first few terms of the sequence: \(4, 5, 7, 10, 14, 19, \ldots\). Look at the difference between terms. \(a_1 - a_0 = 1\) and \(a_2 - a_1 = 2\) and so on. The key thing here is that the difference between terms is \(n\). We can write this explicitly: \(a_n - a_{n-1} = n\). Of course, we could have arrived at this conclusion directly from the recurrence relation - just subtract \(a_{n-1}\) from both sides.

Now use this equation over and over again, changing \(n\) each time:

\[
\begin{align*}
a_1 - a_0 &= 1 \\
a_2 - a_1 &= 2 \\
a_3 - a_2 &= 3 \\
& \vdots \\
a_n - a_{n-1} &= n
\end{align*}
\]

Now add all these equations together. On the right hand side, we get the sum \(1 + 2 + 3 + \cdots + n\). We already know that that can be simplified to \(\frac{n(n+1)}{2}\).

What happens on the left hand side? We get

\[
(a_1 - a_0) + (a_2 - a_1) + (a_3 - a_2) + \cdots (a_{n-1} - a_{n-2}) + (a_n - a_{n-1})
\]

This sum telescopes. We are left with only the \(-a_0\) from the first equation and the \(a_n\) from the last equation. Putting this all together we have \(-a_0 + a_n = \frac{n(n+1)}{2}\) or \(a_n = \frac{n(n+1)}{2} + a_0\). But we know that \(a_0 = 4\). So the solution to the recurrence relation, subject to the initial condition is

\[
a_n = \frac{n(n+1)}{2} + 4
\]

(Now that we know that, we should notice that yes, the sequence is indeed the result of adding 4 to each of the triangular numbers.)
The above example shows a way to solve recurrence relations of the form $a_n = a_{n-1} + f(n)$ where $\sum_{k=1}^{n} f(k)$ has a known closed formula. If you rewrite the recurrence relation as $a_n - a_{n-1} = f(n)$, and then add up all the different equations with $n$ ranging between 1 and $n$, the left hand side will always give you $a_n - a_0$. The right hand side will be $\sum_{k=1}^{n} f(k)$, which is why we need to know the closed formula for that sum.

However, telescoping will not help us with a recursion such as $a_n = 3a_{n-1} + 2$ - the left hand side will not telescope, since you will have $-3a_{n-1}$’s but only one $a_{n-1}$. We need another method for this.

### 4.4.2 Iteration

We have already seen an example of iteration when we found the closed formula for arithmetic and geometric sequences. The idea is, we *iterate* the process of finding the next term, starting with the known initial condition, up until we have $a_n$. Then we simplify. In the arithmetic sequence example, we simplified by multiplying $d$ by the number of times we add it to $a$ when we get to $a_n$, to get from $a_n = a + d + d + d + \cdots + d$ to $a_n = a + d(n-1)$.

To see how this works, let’s go through the same example we used for telescoping, but this time use iteration.

**Example:** Use iteration to solve the recurrence relation $a_n = a_{n-1} + n$ with $a_0 = 4$.

**Solution:** Again, start by writing down the recurrence relation when $n = 1$. This time, don’t subtract the $a_{n-1}$ terms to the other side.

$$a_1 = a_0 + 1$$

Now $a_2 = a_1 + 2$, but we know what $a_1$ is. We get

$$a_2 = (a_0 + 1) + 2$$

Now go to $a_3 = a_2 + 3$, using our known value of $a_2$:

$$a_3 = ((a_0 + 1) + 2) + 3$$

We notice a pattern. Each time, we take the previous term and add the current index. So

$$a_n = (((a_0 + 1) + 2) + 3) + \cdots + n - 1 + n$$

Regrouping terms, we notice that $a_n$ is just $a_0$ plus the sum of the integers from 1 to $n$. So, since $a_0 = 4$,

$$a_n = 4 + \frac{n(n + 1)}{2}$$

Of course in this case we still needed to know formula for the sum of 1, $\ldots$, $n$. Let’s try iteration with a sequence for which telescoping doesn’t work.
Example: Solve the recurrence relation $a_n = 3a_{n-1} + 2$ subject to $a_0 = 1$.

Solution: Again, we iterate the recurrence relation, building up to the index $n$.

\[
\begin{align*}
a_1 &= 3a_0 + 2 \\
a_2 &= 3(a_1) + 2 = 3(3a_0 + 2) + 2 \\
a_3 &= 3[a_2] + 2 = 3[3(3a_0 + 2) + 2] + 2 \\
&\quad \vdots \quad \vdots \\
a_n &= 3(a_{n-1}) + 2 = 3(3(3 \cdots (3a_0 + 2) + 2) \cdots + 2) + 2
\end{align*}
\]

It is difficult to see what is happening here because we have a lot of distributing the 3’s to do. Let’s try again, this time simplifying a bit as we go.

\[
\begin{align*}
a_1 &= 3a_0 + 2 \\
a_2 &= 3(a_1) + 2 = 3(3a_0 + 2) + 2 = 3^2a_0 + 2 \cdot 3 + 2 \\
a_3 &= 3[a_2] + 2 = 3[3^2a_0 + 2 \cdot 3 + 2] + 2 = 3^3a_0 + 2 \cdot 3^2 + 2 \cdot 3 + 2 \\
&\quad \vdots \quad \vdots \\
a_n &= 3(a_{n-1}) + 2 = 3(3^{n-1}a_0 + 2 \cdot 3^{n-2} + \cdots + 2) + 2 = 3^n a_0 + 2 \cdot 3^{n-1} + 2 \cdot 3^{n-2} + \cdots + 2 \cdot 3 + 2
\end{align*}
\]

Now we simplify. $a_0 = 1$, so we have $3^n + (\cdots)$. Note that every other term has a 2 in it - in fact, we have a geometric sum with first term 2 and common ratio 3. We have seen how to simplify $2 + 2 \cdot 3 + 2 \cdot 3^2 + \cdots + 2 \cdot 3^{n-1}$. We get \(\frac{2 - 2 \cdot 3^n}{-2}\) which simplifies to $3^n - 1$. Putting this together with the first $3^n$ term gives our closed formula:

\[a_n = 2 \cdot 3^n - 1\]

Iteration can be messy, but when the recurrence relation only refers to one previous term (and maybe some function of $n$) it can work well. However, trying to iterate a recurrence relation such as $a_n = 2a_{n-1} + 3a_{n-2}$ will be way to complicated - we would need to keep track of two sets of previous terms, each of which were expressed by two previous terms, and so on. The length of the formula would grow exponentially (double each time, in fact). And anyway, there happens to be a third method for solving recurrence relations which works very well on relations like this.

4.4.3 The Characteristic Root Technique

Suppose we want to solve a recurrence relation expressed as a combination of the two previous terms, such as $a_n = a_{n-1} + 6a_{n-2}$. In other words, we want to find a function of $n$ which satisfies $a_n - a_{n-1} - 6a_{n-2} = 0$. Now iteration is too complicated, but think just for a second what would happen if we did iterate. In each step, we would, among other things, multiply a previous iteration by 6. So our closed formula would have include 6 multiplied some number of times. Thus it is reasonable to guess the solution will contain parts that look geometric. So perhaps the solution will take the form $r^n$, for some constant $r$. 

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The nice thing is, we know how to check whether a formula is actually a solution to a recurrence relation - plug it in. What happens if we plug in \( r^n \) into the recursion above? We get
\[
r^n - r^{n-1} - 6r^{n-2} = 0
\]
Now solve for \( r \):
\[
r^{n-2}(r^2 - r - 6) = 0
\]
so by factoring, \( r = -2 \) or \( r = 3 \). This tells us that \( a_n = (-2)^n \) is a solution to the recurrence relation, as is \( a_n = 3^n \). Which one is correct? They both are, unless we specify initial conditions. Notice we could also have \( a_n = (-2)^n + 3^n \). Or \( a_n = 7(-2)^n + 4 \cdot 3^n \). In fact, for any \( a \) and \( b \), \( a_n = a(-2)^n + b3^n \) is a solution - try plugging that into the recurrence relation. To find the values of \( a \) and \( b \), use the initial conditions.

The points us in the direction of a more general technique for solving recurrence relations. Notice will will always be able to factor out the \( r^{n-2} \) as we did above. So we really only care about the other part. We call this other part the characteristic equation for the recurrence relation. We are interested in finding the roots of the characteristic equation, which are called (surprise) the characteristic roots.

### Characteristic Roots
Given a recurrence relation \( a_n + \alpha a_{n-1} + \beta a_{n-2} = 0 \), the characteristic polynomial is
\[
x^2 + \alpha x + \beta
\]
giving the characteristic equation:
\[
x^2 + \alpha x + \beta = 0
\]
If \( r_1 \) and \( r_2 \) are two distinct roots of the characteristic polynomial (i.e., solutions to the characteristic equation), then the solution to the recurrence relation is
\[
a_n = ar_1^n + br_2^n
\]
where \( a \) and \( b \) are constants determined by the initial conditions.

**Example:** Solve the recurrence relation \( a_n = 7a_{n-1} - 10a_{n-2} \) with \( a_0 = 2 \) and \( a_1 = 3 \).

**Solution:** Rewrite the recurrence relation \( a_n - 7a_{n-1} + 10a_{n-2} = 0 \). Now form the characteristic equation:
\[
x^2 - 7x + 10 = 0
\]
and solve for \( x \):
\[
(x - 2)(x - 5) = 0
\]
so $x = 2$ and $x = 5$ are the characteristic roots. We therefore know that the solution to the recurrence relation will have the form

$$a_n = a2^n + b5^n$$

To find $a$ and $b$, plug in $n = 0$ and $n = 1$ to get a system of two equations with two unknowns:

$$
\begin{align*}
2 &= a2^0 + b5^0 = a + b \\
3 &= a2^1 + b5^1 = 2a + 5b
\end{align*}
$$

Solving this system gives $a = \frac{7}{3}$ and $b = -\frac{1}{3}$ so the solution to the recurrence relation is

$$a_n = \frac{7}{3}2^n - \frac{1}{3}3^n$$

Perhaps the most famous recurrence relation is $F_n = F_{n-1} + F_{n-2}$, which together with the initial conditions $F_0 = 0$ and $F_1 = 1$ defines the Fibonacci sequence. But notice that this is precisely the type of recurrence relation on which we can use the characteristic root technique. When you do, the only thing that changes is that the characteristic equation does not factor - you need to use the quadratic formula to find the characteristic roots. In fact, doing so gives the third most famous irrational number - $\varphi$, the golden ratio.

Before leaving the characteristic root technique, we should think about what all might happen when you solve the characteristic equation. Above we have an example in which the characteristic polynomial has two distinct roots. These roots can be integers, or perhaps irrational numbers (requiring the quadratic formula to find them). In these cases, we know what the solution to the recurrence relation looks like.

However, it is possible for the characteristic polynomial to only have one root - say the characteristic polynomial factors as $(x - r)^2$. It is still the case that $r^n$ would be a solution to the recurrence relation, but we won’t be able to find solutions for all initial conditions using the general form $a_n = ar_1^n + br_2^n$, since we can’t distinguish between $r_1^n$ and $r_2^n$ - they are the same (repeated) root. We are in luck though:

### Characteristic Root Technique for Repeated Roots

Suppose the recurrence relation $a_n = \alpha a_{n-1} + \beta a_{n-2}$ has a characteristic polynomial with only one root $r$. Then the solution to the recurrence relation is

$$a_n = ar^n + bnr^n$$

where $a$ and $b$ are constants determined by the initial conditions.

Notice the extra $n$ in $bnr^n$. This allows us to solve for the constants $a$ and $b$ from the initial conditions.
Example: Solve the recurrence relation \( a_n = 6a_{n-1} - 9a_{n-2} \) with initial conditions \( a_0 = 1 \) and \( a_1 = 4 \).

Solution: The characteristic polynomial is \( x^2 - 6x + 9 = 0 \). We solve the characteristic equation

\[
(x - 3)^2 = 0
\]

so \( x = 3 \) is the only characteristic root. Therefore we know that the solution to the recurrence relation has the form

\[
a_n = a3^n + bn3^n
\]

for some constants \( a \) and \( b \). Now use the initial conditions:

\[
a_0 = 1 = a3^0 + b \cdot 0 \cdot 3^0 = a
\]
\[
a_1 = 4 = a \cdot 3 + b \cdot 1 \cdot 3 = 3a + 3b
\]

Since \( a = 1 \), we find that \( b = \frac{1}{3} \). Therefore the solution to the recurrence relation is

\[
a_n = 3^n + \frac{1}{3}n3^n
\]

Although we will not consider examples more complicated than these, this characteristic root technique can be applied to much more complicated recurrence relations. For example, \( a_n = 2a_{n-1} + a_{n-2} - 3a_{n-3} \) has characteristic polynomial \( x^3 - 2x^2 - x + 3 \). Assuming you see how to factor such a degree 3 (or more) polynomial you can easily find the characteristic roots and as such solve the recurrence relation (the solution would look like \( a_n = ar_1^n + br_2^n + cr_3^n \) if there were 3 distinct roots). It is also possible to solve recurrence relations of the form \( a_n = \alpha a_{n-1} + \beta a_{n-2} + C \) for some constant \( C \). Additionally, if the characteristic roots could be complex numbers - this is acceptable as well (and there are tricks to deal with this).

4.5 Generating Functions

We complete our discussion of sequences with a powerful tool of discrete mathematics. A method of manipulating sequences: the generating function. The idea is this: instead of an infinite sequence (for example: 2, 3, 5, 8, 12, ...) we look at a single function which encodes the sequence. But not a function which gives the \( n \)th term as output. Instead, a function whose power series (like from calculus 2) “displays” the terms of the sequence. So for example, we would look at the power series \( 2 + 3x + 5x^2 + 8x^3 + 12x^4 + \cdots \) which displays (as coefficients) the sequence 2, 3, 5, 8, 12, ...
An infinite power series is simply an infinite sum of terms of the form \(c_n x^n\) were \(c_n\) is some constant. So we might write a power series like this:

\[
\sum_{k=0}^{\infty} c_k x^k
\]

or expanded like this

\[
c_0 + c_1 x + c_2 x^2 + c_3 x^3 + c_4 x^4 + c_5 x^5 + \cdots
\]

When viewed in the context of generating functions, we call such a power series a \textit{generating series}. The generating series generates the sequence

\[
c_0, c_1, c_2, c_3, c_4, c_5, \ldots
\]

that is, the sequence generated by a generating series is simply the sequence of \textit{coefficients} of the infinite polynomial.

\textbf{Example:} What sequence is represented by the generating series \(3 + 8x^2 + x^3 + \frac{x^5}{7} + 100x^6 + \cdots\)?

\textit{Solution:} We just read off the coefficients of each \(x^n\) term. So \(a_0 = 3\) since the coefficient of \(x^0\) is \(3\) (\(x^0 = 1\) so this is the constant term). What is \(a_1\)? It is NOT 8, since 8 is the coefficient of \(x^2\), so 8 is the term \(a_2\) of the sequence.

To find \(a_1\) we need to look for the coefficient of \(x^1\) which in this case is 0. So \(a_1 = 0\). We have \(a_3 = 1\), \(a_4 = 0\), and \(a_5 = \frac{1}{7}\). So we have the sequence

\[
3, 0, 8, 1, \frac{1}{7}, 100, \ldots
\]

Note, when discussing generating functions, we always start our sequence with \(a_0\).

Now you might very naturally ask why we would do such a thing. One reason is that encoding a sequence with a power series helps us keep track of which term is which in the sequence. For example, if we write the sequence \(1, 3, 4, 6, 9, \ldots, 24, 41, \ldots\) it is impossible to which term 24 is (even if we agreed that the first term was supposed to be \(a_0\)). However, if wrote the generating series instead we would have \(1 + 3x + 4x^2 + 6x^3 + 9x^4 + \cdots + 24x^{17} + 41x^{18} + \cdots\). Now it is clear that 24 is the 17th term of the sequence (that is, \(a_{17} = 24\). Of course to get this benefit we could have displayed our sequence in any number of ways - perhaps \([1, 3, 4, 6, 9, \ldots, 24, 41, \ldots]\) - but we do not do this. The reason is that the generating series looks like an ordinary power series (although we are interpreting it differently) so we can do things with it that we ordinarily do with power series - such as write down what it converges to.

For example, from calculus we know that the power series \(1 + x + \frac{x^2}{2} + \frac{x^3}{6} + \frac{x^4}{24} + \cdots + \frac{x^n}{n!} + \cdots\) converges to the function \(e^x\). So we can use \(e^x\) as a way of talking about the sequence of coefficients of the power series for \(e^x\). When we write down a nice compact function which has an infinite power series that we view as a generating series, then we call that function a \textit{generating function}. In this example we would say,

\[
1, 1, \frac{1}{2}, \frac{1}{6}, \frac{1}{24}, \ldots, \frac{1}{n!}, \ldots
\]

has generating function \(e^x\)
4.5.1 Building Generating Functions

The \( e^x \) example is very specific - we have that one rather odd sequence, and the only reason we know its generating function is because we happen to know the Taylor series for \( e^x \). Our goal now is to gather some tricks to build the generating function when given a particular sequence.

Let’s see what the generating functions are for some very simple sequences. The simplest of all: 1, 1, 1, 1, 1, \ldots. What does the generating series look like? It is simply \( 1 + x + x^2 + x^3 + x^4 + \cdots \). Now, can we find a closed formula for this power series? Yes! This particular series is really just a geometric series with common ratio \( x \). So if we use our “multiply, shift and subtract” technique, we have

\[
S = 1 + x + x^2 + x^3 + \cdots \\
-xS = x + x^2 + x^3 + x^4 + \cdots \\
(1 - x)S = 1
\]

Therefore we see that

\[
1 + x + x^2 + x^3 \cdots = \frac{1}{1-x}
\]

You might remember from calculus that this is only true on the interval of convergence for the power series - when \(|x| < 1\). That is true for us, but we don’t care - we are never going to plug anything in for \( x \), so as long as there is some value of \( x \) for which the generating function and generating series agree, we are happy. And in this case we are happy:

| The generating function for 1,1,1,1,1,\ldots is \( \frac{1}{1-x} \) |

Now let’s use this basic generating function to find generating functions for more sequences. What if we replace \( x \) by \( -x \). We get

\[
\left( \frac{1}{1-x} \right) = 1 - x + x^2 - x^3 + \cdots \text{ which generates } 1, -1, 1, -1, \ldots \\
\]

If we replace \( x \) by \( 3x \) we get

\[
\left( \frac{1}{1-3x} \right) = 1 + 3x + 9x^2 + 27x^3 + \cdots \text{ which generates } 1, 3, 9, 27, \ldots \\
\]

So by replacing the \( x \) in \( \frac{1}{1-x} \) we can get generating functions for a variety of sequences. But not all. For example, you cannot plug in anything for \( x \) to get the generating function for 2, 2, 2, 2, \ldots. However, we are not lost yet. Notice that each term of 2, 2, 2, 2, \ldots is the result of multiplying the terms of 1, 1, 1, 1, \ldots by a constant (2). So let’s multiply the generating function by 2 as well.

\[
\left( \frac{2}{1-x} \right) = 2 + 2x + 2x^2 + 2x^3 + \cdots \text{ which generates } 2, 2, 2, 2, \ldots \\
\]
Similarly, to find the generating function for the sequence 3, 9, 27, 81, ... we note that this sequence is the result of multiplying each term of 1, 3, 9, 27, ... by 3. Since we have the generating function for 1, 3, 9, 27, ... we can say
\[
\frac{3}{1 - 3x} = 3 \cdot 1 + 3 \cdot 3x + 3 \cdot 9x^2 + 3 \cdot 27x^3 + \cdots \text{ which generates } 3, 9, 27, 81, \ldots
\]

What about the sequence 2, 4, 10, 28, 82, ...? Here the terms are always 1 more than powers of 3. We have added the sequences 1, 1, 1, 1, ... and 1, 3, 9, 27, ... term by term. Therefore we can get a generating function by adding the respective generating functions:
\[
2 + 4x + 10x^2 + 28x^3 + \cdots = (1 + 1)(1 + 3)x + (1 + 9)x^2 + (1 + 27)x^3 + \cdots = \frac{1}{1 - x} + \frac{1}{1 - 3x}
\]

The fun does not stop there: if we replace \(x\) in our original generating function by \(x^2\) we get
\[
\frac{1}{1 - x^2} = 1 + x^2 + x^4 + x^6 \cdots \text{ which generates } 1, 0, 1, 0, 1, 0, \ldots,
\]

How could we get 0, 1, 0, 1, 0, 1, ...? Start with the previous sequence and shift it over by 1. But how do you do this? To see how shifting works, let’s first try to get the generating function for the sequence 0, 1, 3, 9, 27, ... We know that \(\frac{1}{1 - 3x} = 1 + 3x + 9x^2 + 27x^3 + \cdots\). To get the zero out front, we need the generating series to look like \(x + 3x^2 + 9x^3 + 27x^4 + \cdots\) (so there is no constant term). Multiplying by \(x\) has this effect. So the generating function for 0, 1, 3, 9, 27, ... is \(\frac{x}{1 - 3x}\). This will also work to get the generating function for 0, 1, 0, 1, 0, 1, 0, ...:
\[
\frac{x}{1 - x^2} = x + x^3 + x^5 + \cdots \text{ which generates } 0, 1, 0, 1, 0, 1, 0, \ldots
\]

What would happen if we add (term by term) the sequences 1, 0, 1, 0, 1, ... and 0, 1, 0, 1, 0, 1, ...? We should get 1, 1, 1, 1, 1, 1, 1, ... What happens when we add the generating functions? It works (try it)!

\[
\frac{1}{1 - x^2} + \frac{x}{1 - x^2} = \frac{1}{1 - x}
\]

Here’s a tricky one: what happens if you take the derivative of \(\frac{1}{1 - x}\)? We simply get \(\frac{1}{(1 - x)^2}\). But if we differentiate term by term in the power series, we get \((1 + x + x^2 + x^3 + \cdots)' = 1 + 2x + 3x^2 + 4x^3 + \cdots\) which is the generating series for 1, 2, 3, 4, ... This says

The generating function for 1, 2, 3, 4, 5, ... is \(\frac{1}{(1 - x)^2}\)

What happens if we take a second derivative: \(\frac{2}{(1 - x)^3} = 2 + 6x + 12x^2 + 20x^3 + \cdots\). So \(\frac{1}{(1 - x)^3} = 1 + 3x + 6x^2 + 10x^3 + \cdots\) is a generating function for the triangular numbers.
4.5.2 Differencing

We have seen how to find generating functions from \( \frac{1}{1-x} \) using multiplication (by a constant or by \( x \)), substitution, addition, and differentiation. To use each of these, you must notice a way to transform the sequence 1, 1, 1, 1, 1... into your desired sequence. This is not always easy. It is also not really the way we have been looking at analyzing sequences. One thing we have considered often is the sequence of differences between terms of a sequence. This will turn out to be helpful in finding generating functions as well. The idea is that the sequence of differences is often simpler than the original sequence. So if we know a generating function for the differences, we would like to use this to find a generating function for the original sequence. First, we must figure out how to relate a generating series to the generating series for the sequence of differences.

For example, consider the sequence 2, 4, 10, 28, 82,... How could we move to the sequence of first differences: 2, 6, 18, 54,...? We want to subtract 2 from the 4, 4 from the 10, 10 from the 28, and so on. So if we subtract (term by term) the sequence 0, 2, 4, 10, 28,... from 2, 4, 10, 28,..., we will be set. Of course it is easy to find the generating function for 0, 2, 4, 10, 28,... (multiply the generating function for 2, 4, 10, 28,... by \( x \)) - then just subtract. Use \( A \) to represent the generating function for 2, 4, 10, 28, 82,... Then:

\[
A = 2 + 4x + 10x^2 + 28x^3 + 82x^4 + \cdots \\
-xA = 0 + 2x + 4x^2 + 10x^3 + 28x^4 + 82x^5 + \cdots \\
(1 - x)A = 2 + 2x + 6x^2 + 18x^3 + 54x^4 + \cdots
\]

Now we don’t get exactly the sequence of differences - but we get something close. In this particular case, we already know the generating function \( A \) (we found it in the previous section) but most of the time we will use this differencing technique to find \( A \): if we have the generating function for the sequence of differences, we can then solve for \( A \). Here is an example.

**Example:** Find a generating function for 1, 3, 5, 7, 9,...

**Solution:** We notice that the sequence of differences is constant, and we know how to find the generating function for any constant sequence. So call the generating function for 1, 3, 5, 7, 9,... simply \( A \). We have

\[
A = 1 + 3x + 5x^2 + 7x^3 + 9x^4 + \cdots \\
-xA = 0 + x + 3x^2 + 5x^3 + 7x^4 + 9x^5 + \cdots \\
(1 - x)A = 1 + 2x + 2x^2 + 2x^3 + 2x^4 + \cdots
\]

Now we know that \( 2x + 2x^2 + 2x^3 + 2x^4 + \cdots = \frac{2x}{1 - x}. \) Thus

\[
(1 - x)A = 1 + \frac{2x}{1 - x}
\]

Now solve for \( A \):

\[
A = \frac{1}{1 - x} + \frac{2x}{(1 - x)^2} = \frac{1 + x}{(1 - x)^2}
\]
Does this make sense? Before we simplified the two fractions into one, we were adding the generating function for the sequence 1, 1, 1, 1, ... to the generating function for the sequence 0, 2, 4, 6, 8, 10, ... (remember \( \frac{1}{(1-x)^2} \) generates 1, 2, 3, 4, 5, ... - multiplying by 2x shifts it over, putting the zero out front, and doubles each term). If we add these term by term, we get the correct sequence 1, 3, 5, 7, 9, ....

Now that we have a generating function for the odd numbers, we can use that to find the generating function for the squares.

**Example:** Find the generating function for 1, 4, 9, 16, ...

**Solution:** Again we call the generating function for the sequence \( A \). Use differencing:

\[
A = 1 + 4x + 9x^2 + 16x^3 + \cdots \\
-xA = 0 + x + 4x^2 + 9x^3 + 16x^4 + \cdots \\
(1 - x)A = 1 + 3x + 5x^2 + 7x^3 + \cdots
\]

Since \( 1 + 3x + 5x^2 + 7x^3 + \cdots = \frac{1 + x}{(1 - x)^2} \) we have \( A = \frac{1 + x}{(1 - x)^3} \)

In each of the examples above, found the difference between consecutive terms which gave us a sequence of differences we knew a generating function for. We can generalize this to more complicated relationships between terms of the sequence. For example, what if we know that the sequence satisfies the recurrence relation \( a_n = 3a_{n-1} - 2a_{n-2} \)? In other words, if we take a term of the sequence and subtracted 3 times the previous term and then added 2 times the term before that, we would get 0 (since \( a_n - 3a_{n-1} + 2a_{n-2} = 0 \)). That will work for all but the first two terms of the sequence. So after the first two terms, the sequence of results of these calculations would be a sequence of 0’s, which we definitely know a generating function for.

**Example:** The sequence 1, 3, 7, 15, 31, 63, ... satisfies the recurrence relation \( a_n = 3a_{n-1} - 2a_{n-2} \). Find the generating function for the sequence.

**Solution:** Call the generating function for the sequence \( A \). We have

\[
A = 1 + 3x + 7x^2 + 15x^3 + 31x^4 + \cdots + a_nx^n + \cdots \\
-xA = 0 - 3x - 9x^2 - 21x^3 - 45x^4 - \cdots - 3a_{n-1}x^n - \cdots \\
+ \quad 2x^2A = 0 + 0x + 2x^2 + 6x^3 + 14x^4 + \cdots + 2a_{n-2}x^n + \cdots \\
(1 - 3x + 2x^2)A = 1
\]

Let us see what happened there - we multiplied \( A \) by \(-3x\) which shifts every term over one spot and multiplies them by \(-3\). On the third line, we multiplied \( x \) by \( 2x^2 \), which shifted every term over two spots and multiplied them by 2. When we add the corresponding terms up, we are taking each term, subtracting
3 times the previous term, and adding 2 times the term before that. You can see that for the initial terms this does indeed give \(0x^n\). This will happen for each term because \(a_n - 3a_{n-1} + 2a_{n-2} = 0\). In general, we might have two terms from the beginning of the generating series, although in this case the second term happens to be 0 as well.

Now we just need to solve for \(A\):

\[
A = \frac{1}{1 - 3x + 2x^2}
\]

### 4.5.3 Multiplication - Partial Sums

What happens to the sequences when you multiply two generating functions? Let’s see:

\[ A = a_0 + a_1x + a_2x^2 + \cdots \quad \text{and} \quad B = b_0 + b_1x + b_2x^2 + \cdots \]

To multiply \(A\) and \(B\), we need to do a lot of distributing (infinite FOIL?) but keep in mind we will regroup and only need to write down the first few terms to see the pattern. What is the constant term? \(a_0b_0\). What is the coefficient of \(x\)? \(a_0b_1 + a_1b_0\). And so on. We get:

\[
AB = a_0b_0 + (a_0b_1 + a_1b_0)x + (a_0b_2 + a_1b_1 + a_2b_0)x^2 + (a_0b_3 + a_1b_2 + a_2b_1 + a_3b_0)x^3 + \cdots
\]

**Example:** “Multiply” the sequence 1, 2, 3, 4, \ldots by the sequence 1, 2, 4, 8, 16, \ldots.

**Solution:** The new constant term is just 1·1. The next term will be 1·2+2·1 = 4.

The next term: 1·4 + 2·2 + 3·1 = 11. One more: 1·8 + 2·4 + 3·2 + 4·1 = 28.

The resulting sequence is

1, 4, 11, 28, 57, \ldots

Since the generating function for 1, 2, 3, 4, \ldots is \(\frac{1}{(1-x)^2}\) and the generating function for 1, 2, 4, 8, 16, \ldots is \(\frac{1}{1-2x}\), we have that the generating function for 1, 4, 11, 28, 57, \ldots is \(\frac{1}{(1-x)^3(1-2x)}\)

Now what happens when you multiply a sequence by 1, 1, 1, \ldots? Try it with 1, 2, 3, 4, 5, \ldots.

The first term is 1·1 = 1. Then 1·2 + 1·1 = 3. Then 1·3 + 1·2 + 1·1 = 6. The next term will be 10. We are getting the triangular numbers. More precisely we get the sequence of partial sums of 1, 2, 3, 4, 5, \ldots. In terms of generating functions, we take \(\frac{1}{1-x}\) (generating 1, 1, 1, 1, 1, \ldots) and multiply it by \(\frac{1}{(1-x)^2}\) (generating 1, 2, 3, 4, 5, \ldots) and this gives \(\frac{1}{(1-x)^3}\). This should not be a surprise - we found the same generating function for the triangular numbers earlier.

The point is, if you need to find a generating function for the sum of the first \(n\) terms of a particular sequence, and you know the generating function for that sequence, you can multiply it by \(\frac{1}{1-x}\). This makes sense - to go back from the sequence of partial sums to the original sequence, you look at the sequence of differences. When you get the sequence of differences you end up multiplying by \(1-x\) - that is, dividing by \(\frac{1}{1-x}\). Multiplying by \(\frac{1}{1-x}\) gives partial sums, dividing by \(11-x\) gives differences.
4.5.4 Solving Recurrence Relations with Generating Functions

We end with an example of one of the many reasons studying generating functions is helpful - we can use generating functions to solve recurrence relations.

**Example:** Solve the recurrence relation \(a_n = 3a_{n-1} - 2a_{n-2}\) with initial conditions \(a_0 = 1\) and \(a_1 = 3\).

**Solution:** We saw in an example above that this recurrence relation gives the sequence 1, 3, 7, 15, 31, 63, \ldots which has generating function \(\frac{1}{1 - 3x + 2x^2}\). We did this by calling the generating function \(A\) and then computing \(A - 3xA + 2x^2A\) which was just 1, since every other term canceled out.

But how does knowing the generating function help us? Well, we must first break up the generating function into two simpler ones. For this, we need partial fraction decomposition. Start by factoring the denominator:

\[
\frac{1}{1 - 3x + 2x^2} = \frac{1}{(1 - x)(1 - 2x)}
\]

Now partial fraction decomposition tells us that we can write this faction as the sum of two fractions (we decompose the given fraction):

\[
\frac{1}{(1 - x)(1 - 2x)} = \frac{a}{1 - x} + \frac{b}{1 - 2x}
\]

for some constants \(a\) and \(b\)

To find \(a\) and \(b\) we add the two decomposed fractions using a common denominator. This gives

\[
\frac{1}{(1 - x)(1 - 2x)} = \frac{a(1 - 2x) + b(1 - x)}{(1 - x)(1 - 2x)}
\]

so

\[
1 = a(1 - 2x) + b(1 - x)
\]

This must be true for all values of \(x\). If \(x = 1\), then the equation becomes \(1 = -a\) so \(a = -1\). When \(x = \frac{1}{2}\) we get \(1 = b/2\) so \(b = 2\). This tells us that we can decompose the fraction like this:

\[
\frac{1}{(1 - x)(1 - 2x)} = \frac{-1}{1 - x} + \frac{2}{1 - 2x}
\]

This completes the partial fraction decomposition. But now notice that these two fractions look like generating functions we know. In fact, we should be able to expand each of them.

\[
\frac{-1}{1 - x} = -1 - x - x^2 - x^3 - x^4 - \cdots \text{ which generates } -1, -1, -1, -1, -1, \ldots
\]

\[
\frac{2}{1 - 2x} = 2 + 4x + 8x^2 + 16x^3 + 32x^4 + \cdots \text{ which generates } 2, 4, 8, 16, 32, \ldots
\]
We can in fact give a closed formula for the \( n \)th term of each of these sequences. The first is just \( a_n = -1 \). The second is \( a_n = 2^{n+1} \). The sequence we are interested in is just the sum of these, so the solution to the recurrence relation is

\[
a_n = 2^{n+1} - 1
\]

So now we can add generating functions to our list of methods for solving recurrence relations - although we do need to know how to do partial fraction decomposition.

**Exercises**

1. Find the closed formula for each of the following sequences by relating them to a well know sequence. Assume the first term given is \( a_1 \).
   (a) 2, 5, 10, 17, 26, . . .
   (b) 0, 2, 5, 9, 14, 20, . . .
   (c) 8, 12, 17, 23, 30, . . .
   (d) 1, 5, 23, 119, 719, . . .

2. The Fibonacci sequence is 0, 1, 1, 2, 3, 5, 8, 13, . . . (where \( F_0 = 0 \)).
   (a) Give the recursive definition for the sequence.
   (b) Write out the first few terms of the sequence of partial sums.
   (c) Give a closed formula for the sequence of partial sums in terms of \( F_n \) (for example, you might say \( F_0 + F_1 + \cdots + F_n = 3F^2_{n-1} + n \), although that is definitely not correct).

3. Write out the first few terms of the sequence given by \( a_1 = 3; a_n = 2a_{n-1} + 4 \). Then find a recursive definition for the sequence 10, 24, 52, 108, . . .

4. Write out the first few terms of the sequence given by \( a_n = n^2 - 3n + 1 \). Then find a closed formula for the sequence (starting with \( a_1 \)) 0, 2, 6, 12, 20, . . .

5. Consider the sequence 8, 14, 20, 26, . . .
   (a) What is the next term in the sequence?
   (b) Find a formula for the \( n \)th term of this sequence, assuming \( a_1 = 8 \).
   (c) Find the sum of the first 100 terms of the sequence: \( \sum_{k=1}^{100} a_k \).

6. Consider the sequence 1, 7, 13, 19, . . ., \( 6n + 7 \).
   (a) How many terms are there in the sequence?
   (b) What is the second-to-last term?
(c) Find the sum of all the terms in the sequence.

7. Find \(5 + 7 + 9 + 11 + \cdots + 521\).

8. Find \(5 + 15 + 45 + \cdots + 5 \cdot 3^{20}\)

9. Find \(1 - \frac{2}{3} + \frac{4}{5} - \cdots + \frac{2^{30}}{3^{29}}\)

10. Find \(x\) and \(y\) such that \(27, x, y, 1\) is part of an arithmetic sequence. Then find \(x\) and \(y\) so that the sequence is part of a geometric sequence. (\(x\) and \(y\) might not be integers.)

11. Use summation (\(\sum\)) or product (\(\prod\)) notation to rewrite the following.
   (a) \(2 + 4 + 6 + 8 + \cdots + 2n\)
   (b) \(1 + 5 + 9 + 13 + \cdots + 425\)
   (c) \(1 + \frac{1}{2} + \frac{1}{3} + \frac{1}{4} + \cdots + \frac{1}{50}\)
   (d) \(2 \cdot 4 \cdot 6 \cdots \cdot 2n\)
   (e) \((\frac{1}{2}) \cdot (\frac{2}{3}) \cdot (\frac{3}{4}) \cdots (\frac{100}{101})\)

12. Expand the following sums and products. That is, write them out the long way.
   (a) \(\sum_{k=1}^{100} (3 + 4k)\)
   (b) \(\sum_{k=0}^{n} 2^k\)
   (c) \(\sum_{k=2}^{50} \frac{1}{(k^2 - 1)}\)
   (d) \(\prod_{k=2}^{100} \frac{k^2}{(k^2 - 1)}\)
   (e) \(\prod_{k=0}^{n} (2 + 3k)\)

13. Use polynomial fitting to find the formula for the \(n\)th term of the following sequences:
   (a) \(2, 5, 11, 21, 36, \ldots\)
   (b) \(0, 2, 6, 12, 20, \ldots\)

14. Can you use polynomial fitting to find the formula for the \(n\)th term of the sequence \(4, 7, 11, 18, 29, 47, \ldots\)? Explain why or why not.

15. Find the next 2 terms in the sequence \(3, 5, 11, 21, 43, 85, \ldots\). Then give a recursive definition for the sequence. Finally, use the characteristic root technique to find a closed formula for the sequence.
16. Solve the recurrence relation \( a_n = a_{n-1} + 2^n \) with \( a_0 = 5 \).

17. Show that \( 4^n \) is a solution to the recurrence relation \( a_n = 3a_{n-1} + 4a_{n-2} \).

18. Find the solution to the recurrence relation \( a_n = 3a_{n-1} + 4a_{n-2} \) with initial terms \( a_0 = 2 \) and \( a_1 = 3 \).

19. Find the solution to the recurrence relation \( a_n = 3a_{n-1} + 4a_{n-2} \) with initial terms \( a_0 = 5 \) and \( a_1 = 8 \).

20. Solve the recurrence relation \( a_n = 2a_{n-1} - a_{n-2} \).
   
   (a) What is the solution if the initial terms are \( a_0 = 1 \) and \( a_1 = 2 \)?
   
   (b) What do the initial terms need to be in order for \( a_9 = 30 \)?
   
   (c) For which \( x \) are there initial terms which make \( a_9 = x \)?

21. Solve the recurrence relation \( a_n = 3a_{n-1} + 10a_{n-2} \) with initial terms \( a_0 = 4 \) and \( a_1 = 1 \).

22. Find the generating function for each of the following sequences by relating them back to a sequence with known generating function.

   (a) \( 4, 4, 4, 4, 4, \ldots \)
   
   (b) \( 2, 4, 6, 8, 10, \ldots \)
   
   (c) \( 0, 0, 0, 2, 4, 6, 8, 10, \ldots \)
   
   (d) \( 1, 5, 25, 125, \ldots \)
   
   (e) \( 1, -3, 9, -27, 81, \ldots \)
   
   (f) \( 1, 0, 5, 0, 25, 0, 125, 0, \ldots \)
   
   (g) \( 0, 1, 0, 2, 0, 0, 3, 0, 0, 4, 0, 0, 5, \ldots \)

23. Find the sequence generated by the following generating functions:

   (a) \( \frac{4x}{1 - x} \)
   
   (b) \( \frac{1}{1 - 4x} \)
   
   (c) \( \frac{x}{1 + x} \)
   
   (d) \( \frac{3x}{(1 + x)^2} \)
   
   (e) \( \frac{1 + x + x^2}{(1 - x)^2} \) (Hint: multiplication)

24. Show how you can get the generating function for the triangular numbers in three different ways:

   (a) Take two derivatives of the generating function for \( 1, 1, 1, 1, 1, \ldots \)
(b) Use differencing.
(c) Multiply two known generating functions.

25. Use differencing to find the generating function for $4, 5, 7, 10, 14, 19, 25, \ldots$

26. Find a generating function for the sequence with recurrence relation $a_n = 3a_{n-1} - a_{n-2}$ with initial terms $a_0 = 1$ and $a_1 = 5$.

27. Use the recurrence relation for the Fibonacci numbers to find the generating function for the Fibonacci sequence.

28. Use multiplication to find the generating function for the sequence of partial sums of Fibonacci numbers, $S_0, S_1, S_2, \ldots$ where $S_0 = F_0$, $S_1 = F_0 + F_1$, $S_2 = F_0 + F_1 + F_2$, $S_3 = F_0 + F_1 + F_2 + F_3$ and so on.

29. Find the generating function for the sequence with closed formula $a_n = 2(5^n) + 7(-3)^n$.

30. Find a closed formula for the $n$th term of the sequence with generating function $\frac{3x}{1 - 4x} + \frac{1}{1 - x}$.

31. Find $a_7$ for the sequence with generating function $\frac{2}{(1-x)^2} \cdot \frac{x}{1-x-x^2}$
Chapter 5
Mathematical Induction

Mathematical induction is a proof technique, not unlike direct proof or proof by contradiction. However, while both direct proof and proof by contradiction are somewhat natural styles of proof - they arise from the way we think of the statements we want to prove - proof by mathematical induction is different and a little bit strange. Let’s start with an example.

5.1 Stamps

Here is a question about making different amounts of postage:

You need to mail a package, but don’t yet know how much postage you will need.
You have a large supply of 8-cent stamps and 5-cent stamps. Which amounts of postage can you make exactly using these stamps? Which amounts are impossible to make?

Try this problem on your own first. Perhaps in investigating the problem you picked some amounts of postage, and then figured out whether you could make that amount using just 8-cent and 5-cent stamps. Perhaps you did this in order - can you make 1 cents of postage? Can you make 2 cents? 3 cents? And so on. If this is what you did, you were actually answering a sequence of questions. We have methods for dealing with sequences. Let’s see if that helps.

Actually, we will not make a sequence of questions, but rather a sequence of statements. Let \( P(n) \) be the statement “you can make \( n \) cents of postage using just 8-cent and 5-cent stamps.” Since each \( P(n) \) is a statement, it is either true or false. So if we form the sequence of statements

\[
P(1), P(2), P(3), P(4), \ldots
\]

the sequence will consist of T’s (for true) and F’s (for false). In our particular case the sequence starts

\[
\]

because \( P(1), P(2), P(3), P(4) \) are all false (you cannot make 1, 2, 3, or 4 cents of postage) but \( P(5) \) is true (use one 5-cent stamp), and so on.

Let’s think a bit about how we could find the value of \( P(n) \) for some specific value of \( n \) (the “value” will be either \( T \) or \( F \)). How did we find the value of the \( n \)th term of a sequence of numbers? How did we find \( a_n \)? There were two ways we could do this: either there was a closed formula for \( a_n \), so we could plug in \( n \) into the formula and get our output value, or we had a recursive definition for the sequence, so we could use the previous terms of the sequence to compute the \( n \)th term. When dealing with sequences of statements, we could use either of these techniques as well. Maybe there is a way to use \( n \) itself to determine
whether we can make \( n \) cents of postage. That would be something like a closed formula. Or instead we could use the previous terms in the sequence (of statements) to determine whether we can make \( n \) cents of postage. That is, if we know the value of \( P(n-1) \), can we get from that to the value of \( P(n) \)? That would be something like a recursive definition for the sequence. Remember, finding recursive definitions for sequences was often easier than finding closed formulas. The same is true here.

Suppose I told you that \( P(43) \) was true (it is). Can you determine from this fact the value of \( P(44) \) (whether it true or false)? Yes you can. Even if we don’t know how exactly we made 43 cents out of the 5-cent and 8-cent stamps, we do know that there was some way to do it. What if that way used at least three 5-cent stamps (making 15 cents)? We could replace those three 5-cent stamps with two 8-cent stamps (making 16 cents). The total postage has gone up by 1, so we have a way to make 44 cents, so \( P(44) \) is true. Of course, we assumed that we had at least three 5-cent stamps. What if we didn’t? Then we must have at least three 8-cent stamps (making 24 cents). If we replace those three 8-cent stamps with five 5-cent stamps (making 25 cents) then again we have bumped up our total by 1 cent so we can make 44 cents, so \( P(44) \) is true.

Notice that we have not said how to make 44 cents, just that we can, on the basis that we can make 43 cents. How do we know we can make 43 cents? Perhaps because we know we can make 42 cents, which we know we can do because we know we can make 41 cents, and so on. It’s a recursion! As with a recursive definition of a numerical sequence, we must specify our initial value. In this case, the initial value is “\( P(1) \) is false.” That’s not good, since our recurrence relation just says that \( P(k+1) \) is true if \( P(k) \) is also true. We need to start the process with a true \( P(k) \). So instead, we might want to use “\( P(31) \) is true” as the initial condition.

Putting this all together we arrive at the following fact: it is possible to (exactly) make any amount of postage greater than 27 cents using just 5-cent and 8-cent stamps.\(^1\) In other words, \( P(k) \) is true for any \( k \geq 28 \). To prove this, we could do the following:

1. Demonstrate that \( P(28) \) is true.

2. Prove that if \( P(k) \) is true, then \( P(k+1) \) is true (for any \( k \geq 28 \)).

Suppose we have done this. Then we know that the 28th term of the sequence above is \( a \) (using (1) - the initial condition or base case), and that every term after the 28th is \( T \) also (using (2) - the recursive part or inductive case). Here is what the proof would actually look like.

**Proof.** Let \( P(n) \) be the statement “it is possible to make exactly \( n \) cents of postage using 5-cent and 8-cent stamps.” We will show \( P(n) \) is true for all \( n \geq 28 \).

First, we show that \( P(28) \) is true: \( 28 = 4 \cdot 5 + 1 \cdot 8 \), so we can make 28 cents using four 5-cent stamps and one 8-cent stamp.

Now suppose \( P(k) \) is true for some arbitrary \( k \geq 28 \). Then it is possible to make \( k \) cents using 5-cent and 8-cent stamps. Note that since \( k \geq 28 \), it cannot be that we use less than three 5-cent stamps and all 8-cent stamps: using two of each would give only 26

\(^1\) This is not claiming that there are no amounts less than 27 cents which can also be made - there are.
cents. Now if we have made \( k \) cents using at least three 5-cent stamps, replace three 5-cent stamps by two 8-cent stamps. This replaces 15 cents of postage with 16 cents, moving from a total of \( k \) cents to \( k + 1 \) cents - so \( P(k+1) \) is true. On the other hand, if we have made \( k \) cents using at least three 8-cent stamps, then we can replace three 8-cent stamps with five 5-cent stamps, moving from 24 cents to 25 cents, giving a total of \( k + 1 \) cents of postage. So in this case as well \( P(k+1) \) is true.

Therefore, by the principle of mathematical induction, \( P(n) \) is true for all \( n \geq 28 \). QED

### 5.2 Formalizing proofs

What we did in the stamp example above works for many types of problems. Proof by induction is useful when trying to prove statements about all natural numbers, or all natural numbers greater than some fixed first case (like 28 in the example above), and in some other situations besides. In particular, induction should be used when there is some way to go from one case to the next - when you can see how to always “do one more.”

This is a big idea. Thinking about a problem *inductively* can give new insight into the problem. For example, to really understand the stamp problem, you should think about how any amount of postage (greater than 28 cents) can be made (this in non-inductive reasoning) and also how the ways in which postage can be made *changes* as the amount increases (inductive reasoning). When you are asked to provide a proof by induction, you are being asked to think about the problem *dynamically*: how does increasing \( n \) change the problem?

But there is another side to proofs by induction as well. In mathematics, it is not enough to understand a problem, you must also be able to communicate the problem to others. Like any discipline, mathematics has standard language and style, allowing mathematicians to share their ideas efficiently. Proofs by induction have a certain formal style, and being able to write in this style is important. It allows us to keep our ideas organized and might even help us with formulating a proof.

Here is the general structure of a proof by mathematical induction:

---

**Induction Proof Structure**

Start by saying what the statement is which you want to prove: “Let \( P(n) \) be the statement…” To prove that \( P(n) \) is true for all \( n \geq 0 \), you must prove two facts:

1. **Base case:** Prove that \( P(0) \) is true. You do this directly. This is often easy.

2. **Inductive case:** Prove that \( P(k) \rightarrow P(k+1) \) for all \( k \geq 0 \). That is, prove that for any \( k \geq 0 \) if \( P(k) \) is true, then \( P(k+1) \) is true as well. This is the proof of an if...then...statement, so you can assume \( P(k) \) is true (\( P(k) \) is called the *inductive hypothesis*). You must then explain why \( P(k+1) \) is also true, given that assumption.

Assuming you are successful on both parts above, you can conclude, “Therefore by the principle of mathematical induction, \( P(n) \) is true for all \( n \geq 0 \).”

---
Sometimes the statement \( P(n) \) will only be true for values of \( n \geq 4 \), for example, or some other value. In such cases, replace all the 0’s above with 4’s (or the other value).

The other advantage of formalizing inductive proofs is it allows us to verify that the logic behind this style of argument is valid. Why does induction work? Think of a row of dominoes set up standing on their edges. We want to argue that in a minute, all the dominoes will have fallen down. For this to happen, you will need to push the first domino - that is the base case. It will also have to be that the dominoes are close enough together that when any particular domino falls, it will cause the next domino to fall - that is the inductive case. If both of these conditions are met - you push the first domino over and each domino will cause the next to fall, then all the dominoes will fall.

Induction is powerful! Think how much easier it is to knock over dominoes when you don’t have to push over each domino yourself. You just start the chain reaction, and the rely on the relative nearness of the dominoes to take care of the rest.

Think about our study of sequences. It is easier to find recursive definitions for sequences than closed formulas - going from one case to the next is easier than going directly to a particular case. That is what is so great about induction. Instead of going directly to the (arbitrary) case for \( n \), we just need to say how to get from one case to the next.

When you are asked to prove a statement by mathematical induction, you should first think about why the statement is true, using inductive reasoning. Explain why induction is the right thing to do, and roughly why the inductive case will work. Then, sit down and write out a careful, formal proof using the structure above.

### 5.3 Examples

Here are some examples of proof my mathematical induction.

**Example:** Prove for each natural number \( n \geq 1 \) that \( 1+2+3 + \cdots + n = \frac{n(n+1)}{2} \).

**Solution:** First, let’s think inductively about this equation. In fact, we know this is true for other reasons (reverse and add comes to mind). But why might induction be applicable? The left hand side is adding up the numbers from 1 to \( n \). If we know how to do that, adding just one more term \( (n+1) \) would not be that hard. For example, if \( n = 100 \), suppose we know that the sum of the first 100 numbers is 5050 (so \( 1 + 2 + 3 + \cdots + 100 = 5050 \), which is true). Now to find the sum of the first 101 numbers, it makes more sense to just add 101 to 5050, instead of computing the entire sum again. We would have \( 1 + 2 + 3 + \cdots + 100 + 101 = 5050 + 101 = 5151 \). In fact, it would always be easy to add just one more term. This is why we should use induction.

Now the formal proof.

**Proof.** Let \( P(n) \) be the statement \( 1+2+3 + \cdots + n = \frac{n(n+1)}{2} \). We will show that \( P(n) \) is true for all natural numbers \( n \geq 1 \).

**Base case:** \( P(1) \) is the statement \( 1 = \frac{1(1+1)}{2} \) which is clearly true.

**Inductive case:** Let \( k \geq 1 \) be a natural number. Assume (for induction) that \( P(k) \) is true. That means \( 1+2+3+\cdots+k = \frac{k(k+1)}{2} \). We will prove that \( P(k+1) \)
is true as well. That is, we must prove that $1+2+3+\cdots+k+(k+1) = \frac{(k+1)(k+2)}{2}$.

To prove this equation, start with the left hand side:

$$1 + 2 + 3 + \cdots + k + (k + 1) = \frac{k(k + 1)}{2} + (k + 1)$$

by the induction hypothesis. Now simplifying:

$$\frac{k(k + 1)}{2} + (k + 1) = \frac{k(k + 1) + 2(k + 1)}{2} = \frac{(k + 2)(k + 1)}{2}$$

Thus $P(k + 1)$ is true, so by the principle of mathematical induction $P(n)$ is true for all natural numbers $n \geq 1$. QED

Note that in the part of the proof in which we proved $P(k + 1)$ from $P(k)$, we used the equation $P(k)$ - this was the inductive hypothesis. This is very easy to do with proving a fact about a sum like this. Sometimes it is not as easy to see where to use the inductive hypothesis.

**Example:** Prove that for all $n \in \mathbb{N}$, $6^n - 1$ is a multiple of 5.

**Solution:** Again, start by understanding the dynamics of the problem. What does increasing $n$ do? Let’s try with a few examples. If $n = 1$, then yes, $6^1 - 1 = 5$ is a multiple of 5. What does incrementing $n$ to 2 look like? We get $6^2 - 1 = 35$, which again is a multiple of 5. Next, $n = 3$: but instead of just finding $6^3 - 1$, what did the increase in $n$ do? We will still subtract 1, but now we are multiplying by another 6 first. Viewed another way, we are multiplying a number which is one more than a multiple of 5 by 6 (because $6^2 - 1$ is a multiple of 5, so $6^2$ is one more than a multiple of 5). What do numbers which are one more than a multiple of 5 look like? The must have last digit 1 or 6. What happens when you multiply such a number by 6? Depends on the number, but in any case, the last digit of the new number must be a 6. And then if you subtract 1, you get last digit 5, so a multiple of 5.

The point is, every time we multiply by just one more six, we still get a number with last digit 6, so subtracting 1 gives us a multiple of 5. Now the formal proof:

**Proof.** Let $P(n)$ be the statement, “$6^n - 1$ is a multiple of 5.” We will prove that $P(n)$ is true for all $n \in \mathbb{N}$.

Base case: $P(0)$ is true: $6^0 - 1 = 0$ which is a multiple of 5.

Inductive case: Let $k$ be an arbitrary natural number. Assume, for induction, that $P(k)$ is true. That is, $6^k - 1$ is a multiple of 5. Then $6^k - 1 = 5j$ for some integer $j$. This means that $6^k = 5j + 1$. Multiply both sides by 6:

$$6^{k+1} = 6(5j + 1) = 30j + 6$$

But we want to know about $6^{k+1} - 1$, so subtract 1 from both sides:

$$6^{k+1} - 1 = 30j + 5$$
Of course $30j + 5 = 5(6j + 1)$, so is a multiple of 5.

Therefore $6^{k+1} - 1$ is a multiple of 5, or in other words, $P(k+1)$ is true. Thus, by the principle of mathematical induction $P(n)$ is true for all $n \in \mathbb{N}$. QED

We had to be a little bit clever (i.e., use some algebra) to locate the $6^k - 1$ inside of $6^{k+1} - 1$ before we could apply the inductive hypothesis. This is what can make inductive proofs challenging.

In the two examples above, we started with $n = 1$ or $n = 0$. We can start later if we need to.

**Example:** Prove that $n^2 < 2^n$ for all integers $n \geq 5$.

**Solution:** First, the idea of the argument. What happens when we increase $n$ by 1? On the left hand side, we increase the base of the square - we go to the next square number. On the right hand side, we increase the power of 2. This means we double the number. So the question is, how does doubling a number relate to increasing to the next square? Think about what the difference of two consecutive squares looks like. We have $(n + 1)^2 - n^2$. This factors:

$$(n + 1)^2 - n^2 = (n + 1 - n)(n + 1 + n) = 2n + 1$$

(this should not be a surprise - we already know that the sum of consecutive odd numbers gives the perfect squares). But doubling the right hand side increases it by $2^n$, since $2^{n+1} = 2^n + 2^n$. Of course when $n$ is large enough, $2^n > 2n + 1$.

What we are saying here is that each time $n$ increases, the left hand side grows by less than the right hand side. So if the left hand side starts smaller (as it does when $n = 5$), it will never catch up. Now the formal proof:

**Proof.** Let $P(n)$ be the statement $n^2 < 2^n$. We will prove $P(n)$ is true for all integers $n \geq 5$.

Base case: $P(5)$ is the statement $5^2 < 2^5$. Since $5^2 = 25$ and $2^5 = 32$, we see that $P(5)$ is indeed true.

Inductive case: Let $k \geq 5$ be an arbitrary integer. Assume, for induction, that $P(k)$ is true. That is, assume $k^2 < 2^k$. We will prove that $P(k+1)$ is true, i.e., $(k+1)^2 < 2^{k+1}$. To prove such an inequality, start with the left hand side and work towards the right hand side:

$$(k + 1)^2 = k^2 + 2k + 1$$

$$(k + 1)^2 < 2^k + 2k + 1$$

$$(k + 1)^2 < 2^k + 2k$$

$$(k + 1)^2 = 2^{k+1}$$

$\ldots$ by the inductive hypothesis

$\ldots$ since $2k + 1 < 2^k$ for $k \geq 5$

$2P(n)$ is not the statement “$n^2 < 2^n$ for all integers $n \geq 5$” - you do not include the quantifier on $n$ in the statement of $P(n)$.  

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Following the equalities and inequalities through, we get $(k + 1)^2 < 2^{k+1}$, in other words, $P(k + 1)$. Therefore by the principle of mathematical induction, $P(n)$ is true for all $n \geq 5$. QED

The previous example might remind you of the *racetrack principle* from calculus, which says that if $f(a) < g(a)$, and $f'(x) < g'(x)$ for $x > a$, then $f(x) < g(x)$ for $x > a$. Same idea: the larger function is increasing at a faster rate than the smaller function, so the larger function will stay larger. In discrete math, we don’t have derivatives, so we look at differences - when $n$ increase by 1. Thus induction is the way to go.

**Warning:**

With great power, comes great responsibility. Induction isn’t magic. It seems very powerful to be able to assume $P(k)$ is true - after all, we are trying to prove $P(n)$ is true (and the only difference is in the variable: $k$ vs. $n$). Are we assuming that what we want to prove is true? Not really. We assume $P(k)$ is true only for the sake of proving that $P(k+1)$ is true.

Still you might start to believe that you can prove anything with induction. Consider this incorrect “proof” that every Canadian has the same eye color: Let $P(n)$ be the statement that any $n$ Canadians have the same eye color. $P(1)$ is true, since everyone has the same eye color as themselves. Now assume $P(k)$ is true. That is, assume that in any group of $k$ Canadians, everyone has the same eye color. Now consider an arbitrary group of $k+1$ Canadians. The first $k$ of these must all have the same eye color, since $P(k)$ is true. Also, the last $k$ of these must have the same eye color, since $P(k)$ is true. So in fact, everyone the group must have the same eye color. Thus $P(k+1)$ is true. So by the principle of mathematical induction, $P(n)$ is true for all $n$.

Clearly something went wrong. The problem is that the proof that $P(k)$ implies $P(k+1)$ assumes that $k \geq 2$. We have only shown $P(1)$ is true. In fact, $P(2)$ is false.

### 5.4 Strong Induction

Consider the following classic puzzle:

You have 9 coins which all look identical. However, you know that one of the coins is counterfeit, and weighs slightly less than the other coins (which all have the same weight). The weight difference is not great enough to tell by holding the coins, but luckily you have an old-fashioned balance scale (like the one held by Lady Justice). What is the smallest number of weighings needed to find the counterfeit coin?

If you have never tried to solve this puzzle, you should try now. Don’t read on until you have given it some thought.

Need a hint? The answer is 2. Why?

Here’s what you do. Put 3 coins on each side of the balance scale, leaving 3 off. Now if the scales are balanced, then the light coin must be one of the 3 you didn’t weigh. If the scale
Induction 5.4. Strong Induction

tips, then the light coin is one of the three on the higher side. Now you have eliminated all but 3 coins using just one weighing. For the second weighing, put one of the three suspect coins on each side of the scales (leaving one off). Same idea: if the scales balance, then the unweighed coin is the counterfeit. Otherwise, the side of the scale that tipped up contains the light coin. Neat huh?

Let’s generalize? If you had more than 9 coins to start with, say you had \( n \) coins, how many weighings would you need? A reasonable guess would be that each time you weigh the coins, you can eliminate \( 2/3 \) of them. So using 3 weighings, you could find the counterfeit coin among up to 27 coins. A forth weighing would be required for 28 to 81 coins.

**Conjecture.** The minimum number of weighings required to find the counterfeit coin among \( n \) coins is \( \log_3(n) \), rounded up. In other words, if \( x \) is the smallest integer such that \( n \leq 3^x \), then the minimum number of weighings is \( x \).

How do we know this is true? If you have \( n \) coins, divide them into 3 equal piles (if this is not possible, make two piles the same size, and the third pile one more or less as needed). Put one pile on each side of the scale, leaving one pile off (if the piles were not equal, weigh the piles with an equal number of coins against each other). If the scales balance, the counterfeit coin is in the unweighed pile. If the scales tip, the counterfeit coin is in the light pile. In any event, we have now reduced the number of coins to roughly \( n/3 \).

What now? Well obviously we do this again. And again, until we find our bad coin. Here is where induction comes in: instead of going through the process over and over again, we can just say that we have reduced the problem to an earlier case, for which we assume our conjecture holds. For example if \( n = 81 \), we do the weighing as described, leaving us with \( n/3 = 27 \) coins. By induction, we already know how to complete the problem for 27 coins, so we are done.

This “reducing to a simpler case” is a type of induction. But not exactly the same sort of induction we have encountered so far. In ordinary induction, we assume our result holds for \( k \), and prove that it holds for the next higher case: \( k + 1 \). There is not much difference in showing that the result holds for the case \( k \) assuming that it holds for the previous case \( (k - 1) \). But what we are doing here, is showing that the result holds for the case \( k \), assuming it works for all smaller cases. Is this okay? In fact, it is. Think of the dominoes: to prove that they all fall down, it is enough to prove that the first one falls, and that if all the dominoes up to the \( k \)th one have fallen, then the \( k \)th one will fall too.

The advantage is that we now have a stronger inductive hypothesis. We can assume that \( P(1), P(2), P(3), \ldots P(k) \) is true, just to show that \( P(k + 1) \) is true. Previously, we just assumed \( P(k) \) for this purpose.

It is slightly easier to change our variables for strong induction. Here is what the formal proof would look like:
Induction 5.4. Strong Induction

**Strong Induction Proof Structure**

Again, start by saying what you want to prove: “Let \( P(n) \) be the statement...” Then establish two facts:

1. **Base case**: Prove that \( P(0) \) is true.
2. **Inductive case**: Assume \( P(k) \) is true for all \( k < n \). Prove that \( P(n) \) is true.

Conclude, “therefore, by strong induction, \( P(n) \) is true for all \( n > 0 \).”

Of course, it is acceptable to replace 0 with a larger base case if needed.\(^3\)

Now let’s prove our conjecture about the coin weighing puzzle.

**Proof.** Let \( P(n) \) be the statement, “if \( x \) is the smallest natural number such that \( n \leq 3^x \), then it is possible to find the light coin in no more than \( x \) weighings.”

Base case: Consider \( P(2) \). Here \( x = 1 \), since \( 2 \leq 3^1 \). It is possible to find the heavy coin using just 1 weighing - put the coins on the scale.

Inductive case: Suppose we have \( n \) coins, and assume \( P(k) \) is true for all \( k < n \). Now divide the \( n \) coins into three piles - equal if possible, otherwise two equal with the third either one more or one less than the other two. Put two equal piles on each side of the scales. If the scales balance, then the counterfeit coin is in the unweighed pile. If the scales tip, the coin is in the pile on the lighter side of the scale. In any event, we now know the counterfeit coin is one of \( k = n/3 \) or \( k = n/3 \pm 1 \) coins. What’s more, if \( x \) is the smallest number such that \( n \leq 3^x \), then \( k \leq 3^{x-1} \). By the inductive hypothesis, it is now possible to find the counterfeit coin in \( x-1 \) more weighings. Thus \( P(n) \) is true.

Therefore, by strong induction, \( P(n) \) is true for all \( n \geq 2 \). QED

Here is a more mathematically relevant example.

**Example:** Prove that any natural number greater than 1 is either prime or can be written as the product of primes.

**Solution:** First, the idea: if we take some number \( n \), maybe it is prime. If so, we are done. If not, then it is composite, so it is the product of two smaller numbers. Each of these factors is smaller than \( n \) (but at least 2), so we can repeat the argument with these numbers. We have reduced to a smaller case.

Now the formal proof:

**Proof.** Let \( P(n) \) be the statement, “\( n \) is either prime or can be written as the product of primes.” We will prove \( P(n) \) is true for all \( n \geq 2 \).

Base case: \( P(2) \) is true because 2 is indeed prime.

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\(^3\)Technically, strong induction does not require you to prove a separate base case. This is because when proving the inductive case, you must show that \( P(0) \) is true, assuming \( P(k) \) is true for all \( k < 0 \). But this is not any help - so you end up proving \( P(0) \) anyway. To be on the safe side, we will always include the base case separately.
Inductive case: assume \( P(k) \) is true for all \( k < n \). We want to show that \( P(n) \) is true. That is, we want to show that \( n \) is either prime or is the product of primes. If \( n \) is prime, we are done. If not, then \( n \) has more than 2 divisors, so we can write \( n = m_1 \cdot m_2 \), with \( m_1 \) and \( m_2 \) less than \( n \) (and greater than 1). By the inductive hypothesis, \( m_1 \) and \( m_2 \) are each either prime or can be written as the product of primes. In either case, we have that \( n \) is written as the product of primes.

Thus by the strong induction, \( P(n) \) is true for all \( n \geq 2 \). QED

Exercises

1. Use induction to prove for all \( n \in \mathbb{N} \) that \( \sum_{k=0}^{n} 2^k = 2^{n+1} - 1 \).

2. Prove that \( 7^n - 1 \) is a multiple of 6 for all \( n \in \mathbb{N} \).

3. Prove that \( 1 + 3 + 5 + \cdots + (2n - 1) = n^2 \) for all \( n \geq 1 \).

4. Prove that \( F_0 + F_2 + F_4 + \cdots + F_{2n} = F_{2n+1} - 1 \) where \( F_n \) is the \( n \)th Fibonacci number.

5. Prove that \( 2^n < n! \) for all \( n \geq 4 \). (Recall, \( n! = 1 \cdot 2 \cdot 3 \cdots \cdot n \).)

6. What is wrong with the following “proof” of the “fact” that \( n + 3 = n + 7 \) for all values of \( n \) (besides of course that the thing it is claiming to prove is false)?

   \[ \text{Proof. Let } P(n) \text{ be the statement that } n + 3 = n + 7. \text{ We will prove that } P(n) \text{ is true for all } n \in \mathbb{N}. \text{ Assume, for induction that } P(k) \text{ is true. That is, } k + 3 = k + 7. \text{ We must show that } P(k + 1) \text{ is true. Now since } k + 3 = k + 7, \text{ add 1 to both sides. This gives } k + 3 + 1 = k + 7 + 1. \text{ Regrouping } (k + 1) + 3 = (k + 1) + 7. \text{ But this is simply } P(k + 1). \text{ Thus by the principle of mathematical induction } P(n) \text{ is true for all } n \in \mathbb{N}. \quad \text{QED} \]

7. The proof in the previous problem does not work. But if we modify the “fact,” we can get a working proof. Prove that \( n + 3 < n + 7 \) for all values of \( n \in \mathbb{N} \). You can do this proof with algebra (without induction), but the goal of this exercise is to write out a valid induction proof.

8. Find the flaw in the following “proof” of the “fact” that \( n < 100 \) for every \( n \in \mathbb{N} \).

   \[ \text{Proof. Let } P(n) \text{ be the statement } n < 100. \text{ We will prove } P(n) \text{ is true for all } n \in \mathbb{N}. \text{ First we establish the base case: when } n = 0, \text{ } P(n) \text{ is true, because } 0 < 100. \text{ Now for the inductive step, assume } P(k) \text{ is true. That is, } k < 100. \text{ Now if } k < 100, \text{ then } k \text{ is some number, like 80. Of course } 80 + 1 = 81 \text{ which is still less than 100. So } k + 1 < 100 \text{ as well. But this is what } P(k + 1) \text{ claims, so we have shown that } P(k) \rightarrow P(k + 1). \text{ Thus by the principle of mathematical induction, } P(n) \text{ is true for all } n \in \mathbb{N}. \quad \text{QED} \]
9. While the above proof does not work (it better not - the statement it is trying to prove is false!) we can prove something similar. Prove that there is a strictly increasing sequence \(a_1, a_2, a_3, \ldots\) of numbers (not necessarily integers) such that \(a_n < 100\) for all \(n \in \mathbb{N}\). (By strictly increasing we mean \(a_n < a_{n+1}\) for all \(n\) - so each term must be larger than the last.)

10. What is wrong with the following “proof” of the “fact” that for all \(n \in \mathbb{N}\), the number \(n^2 + n\) is odd?

\[\text{Proof.} \quad \text{Let } P(n) \text{ be the statement “} n^2 + n \text{ is odd.” We will prove that } P(n) \text{ is true for all } n \in \mathbb{N}. \text{ Suppose for induction that } P(k) \text{ is true, that is, that } k^2 + k \text{ is odd. Now consider the statement } P(k+1). \text{ Now } (k+1)^2 + (k+1) = k^2 + 2k + 1 + k + 1 = k^2 + k + 2k + 2. \text{ By the inductive hypothesis, } k^2 + k \text{ is odd, and of course } 2k + 2 \text{ is even. An odd plus an even is always odd, so therefore } (k+1)^2 + (k+1) \text{ is odd. Therefore by the principle of mathematical induction, } P(n) \text{ is true for all } n \in \mathbb{N}. \quad \text{QED}\]

11. Now give a valid proof (by induction - even though you might be able to do so without using induction) of the statement, “for all \(n \in \mathbb{N}\), the number \(n^2 + n\) is even.”

12. Prove that there is a sequence of positive real numbers \(a_1, a_2, a_3, \ldots\) such that the partial sum \(a_1 + a_2 + a_3 + \cdots + a_n\) is strictly less than 2 for all \(n \in \mathbb{N}\). Hint: think about how you could define what \(a_{k+1}\) is to make the induction argument work.

13. Prove that every natural number is either a power of 2, or can be written as the sum of distinct powers of 2.

14. Use induction to prove that if \(n\) people all shake hands with each other, that the total number of handshakes is \(\frac{n(n-1)}{2}\).

15. Prove that if \(G\) is a connected planar graph with \(V\) vertices, \(E\) edges, and \(F\) faces, then \(V - E + F = 2\). Hint: use induction on the number of edges in the graph.

16. Use induction to prove that \(\sum_{k=0}^{n} \binom{n}{k} = 2^n\). That is, the sum of the \(n\)th row of Pascal’s Triangle is \(2^n\).

17. Use induction to prove \(\binom{4}{0} + \binom{5}{1} + \binom{6}{2} + \cdots + \binom{4+n}{n} = \binom{5+n}{n}\). (This is an example of the hockey stick theorem.)

18. Use the product rule for logarithms \((\log(ab) = \log(a) + \log(b))\) to prove, by induction on \(n\), that \(\log(a^n) = n \log(a)\), for all natural numbers \(n \geq 2\).

19. Let \(f_1, f_2, \ldots, f_n\) be differentiable functions. Prove, using induction, that

\[(f_1 + f_2 + \cdots + f_n)' = f_1' + f_2' + \cdots + f_n'\]

You may assume \((f + g)' = f' + g'\) for any differentiable functions \(f\) and \(g\).
20. Suppose $f_1, f_2, \ldots, f_n$ are differentiable functions. Use mathematical induction to prove the generalized product rule:

$$(f_1 f_2 f_3 \cdots f_n)' = f_1' f_2 f_3 \cdots f_n + f_1 f_2' f_3 \cdots f_n + \cdots + f_1 f_2 f_3 \cdots f_n'$$

You may assume the product rule for two functions is true.
Chapter 6

Introduction to Number Theory

So far in this course, we have used the natural numbers to solve problems. This was the right set of numbers to work with in discrete mathematics because we always dealt with a whole number of things. The natural numbers have been a tool. Now let’s take a moment to inspect that tool. What mathematical discoveries can we make about the natural numbers themselves?

This is the main question of number theory, a huge, ancient, complex, and above all, beautiful branch of mathematics. Historically, number theory was known as the Queen of Mathematics - it was very much a branch of pure mathematics, studied for its own sake instead of as a means to understanding real world applications. However, this has changed in recent years, as applications of number theory have been unearthed. Probably the most well known example of this is RSA cryptography, one of the methods used in encrypt data on the internet. It is number theory that makes this possible.

So what sorts of questions belong to the realm of number theory? Here is a motivating example. Recall in our study of induction, we asked:

Which amounts of postage can be made exactly using just 5-cent and 8-cent stamps?

We were able to prove that any amount greater than 27 cents could be made. You might wonder what would happen if we changed the denomination of the stamps. What if we instead had 4- and 9-cent stamps? Would there be some amount after which all amounts would be possible? Well, again, we could replace two 4-cent stamps with a 9-cent stamp, or three 9-cent stamps with seven 4-cent stamps. In each case we can create one more cent of postage. Using this as the inductive case would allow us to prove that any amount of postage greater than 23 cents can be made.

What if we had 2-cent and 4-cent stamps. Here it looks less promising. If we take some number of 2-cent stamps and some number of 4-cent stamps, what can we say about the total? Could it ever be odd? Doesn’t look like it.

So why does 5 and 8 work, 4 and 9 work, but 2 and 4 not work? What is it about these numbers? If I gave you a pair of numbers, could you tell me right away if they would work or not? We will answer these questions, and more, after first investigating some simpler properties of numbers themselves.

6.1 Divisibility

It is easy to add and multiply natural numbers. If we extend our focus to all integers, then subtraction is also easy (we need the negative numbers so we can subtract any number from any other numbers, even larger from smaller). Division is the first operation that presents
a challenge. If we wanted to extend our set of numbers so any division would be possible (maybe excluding division by 0) we would need to look at the rational numbers - the set of all numbers which can be written as fractions. This would be going too far, so we will refuse this option.

In fact, it is a good thing that not every number can be divided by other numbers. This helps us understand the structure of the natural numbers and opens the door to many interesting questions and applications.

If given numbers \( a \) and \( b \), it is possible that \( a \div b \) gives a whole number. In this case, we say that \( b \) divides \( a \), in symbols, we write \( b \mid a \). If this holds, then \( b \) is a divisor or factor of \( a \), and \( a \) is a multiple of \( b \). This also means that if \( b \mid a \), then \( a = bk \) for some integer \( k \) (this is saying \( a \) is some multiple of \( b \)).

### The Divisibility Relation

Given integers \( m \) and \( n \), we say,

\[
m \mid n \quad \text{“} m \text{ divides } n \text{”}
\]

provided \( n \div m \) is an integer. Thus the following assertions mean the same thing:

1. \( m \mid n \)
2. \( n = mk \) for some integer \( k \)
3. \( m \) is a factor (or divisor) of \( n \)
4. \( n \) is a multiple of \( m \).

Notice that \( m \mid n \) is a statement - it is either true or false. On the other hand, \( n \div m = n/m \) is some number. If we want to claim that \( n/m \) is not an integer, so \( m \) does not divide \( n \), then we can write \( m \nmid n \).

**Example:** Decide whether each of the statements below are true or false.

1. \( 4 \mid 20 \)
2. \( 20 \mid 4 \)
3. \( 0 \mid 5 \)
4. \( 5 \mid 0 \)
5. \( 7 \mid 7 \)
6. \( 1 \mid 37 \)
7. \( -3 \mid 12 \)
8. \( 8 \mid 12 \)
9. \( 1642 \mid 136299 \)

**Solution:**

1. True. \( 4 \) “goes into” \( 20 \) five times without remainder. In other words, \( 20 \div 4 = 5 \), an integer. We could also justify this by saying that \( 20 \) is a multiple of \( 4 \): \( 20 = 4 \cdot 5 \).
2. False. While \( 20 \) is a multiple of \( 4 \), it is false that \( 4 \) is a multiple of \( 20 \).
3. False. \( 5 \div 0 \) is not even defined, let alone an integer.
4. True. In fact, \( x \mid 0 \) is true for all \( x \). This is because 0 is a multiple of every number: \( 0 = x \cdot 0 \).

5. True. In fact, \( x \mid x \) is true for all \( x \).

6. True. 1 divides every number (other than 0).

7. True. Negative numbers work just fine for the divisibility relation. Here \( 12 = -3 \cdot 4 \). It is also true that \( 3 \mid -12 \) and that \( -3 \mid -12 \).

8. False. Both 8 and 12 are divisible by 4, but this does not mean that 12 is divisible by 8.


This last example raises a question: how might one decide whether \( m \mid n \)? Of course, if you had a trusted calculator, you could ask it for the value of \( n \div m \) - if it spits out anything other than an integer, you know \( m \nmid n \). This seems a little like cheating though: we don’t have division, so should we really use division to check divisibility?

While we don’t really know how to divide, we do know how to multiply. So we might try multiplying \( m \) by larger and larger numbers until we get close to \( n \). How close? Well, we want to be sure that if we multiply \( m \) by the next larger integer, we go over \( n \).

For example, let’s try this to decide whether \( 1642 \mid 136299 \). Start finding multiples of 1642:

\[
\begin{align*}
1642 \cdot 2 &= 3284 \\
1642 \cdot 3 &= 4926 \\
1642 \cdot 4 &= 6568 \\
\vdots
\end{align*}
\]

All of these are well less than 136299. I suppose we can jump ahead a bit:

\[
\begin{align*}
1642 \cdot 50 &= 82100 \\
1642 \cdot 80 &= 131360 \\
1642 \cdot 85 &= 139570
\end{align*}
\]

Ah, so we need to look somewhere between 80 and 85. Try 83:

\[
1642 \cdot 83 = 136286
\]

Is this the best we can do? How far are we from our desired 136299? If we subtract, we get \( 136299 - 136286 = 13 \). So we know we cannot go up to 84, that will be too much. In other words, we have found that

\[
136299 = 83 \cdot 1642 + 13
\]

Since \( 13 < 1642 \), we can now safely say that \( 1642 \nmid 136299 \).

It turns out that the process we went through above can be repeated for any pair of numbers. We can always write the number \( a \) as some multiple of the number \( b \) plus some remainder. We know this because we know about division with remainder from elementary school. This is just a way of saying it using multiplication. Due to the procedural nature that can be used to find the remainder, this fact is usually called the division algorithm:

<table>
<thead>
<tr>
<th>The Division Algorithm</th>
</tr>
</thead>
<tbody>
<tr>
<td>Given any two integers ( a ) and ( b ), we can always find an integer ( q ) such that</td>
</tr>
<tr>
<td>( a = qb + r )</td>
</tr>
<tr>
<td>where ( r ) is an integer satisfying ( 0 \leq r &lt;</td>
</tr>
</tbody>
</table>
The idea is we can always take a large enough multiple of \( b \) so that the remainder \( r \) is as small as possible. Note that we do allow the possibility of \( r = 0 \). In this case, we get \( b \div a \).

6.2 Remainder Classes

The division algorithm tells us that there are only \( b \) possible remainders when dividing by \( b \). If we fix this divisor, we can group integers by the remainder. Each group is called a remainder class modulo \( b \) (or sometimes residue class).

**Example:** Describe the remainder classes modulo 5.

**Solution:** We want to classify numbers by what their remainder would be when divided by 5. From the division algorithm, we know there will be exactly 5 remainder classes, because there are only 5 choices for what \( r \) could be \((0 \leq r < 5)\).

First consider \( r = 0 \). Here we are looking for all the numbers divisible by 5. In other words, the multiples of 5. We get the infinite set:

\[
\{ \ldots, -15, -10, -5, 0, 5, 10, 15, 20, \ldots \}
\]

Notice we also include negative integers.

Next consider \( r = 1 \). Which integers, when divided by 5, have remainder 1? Well, certainly 1, does, as does 6, and 11. Negatives? Here we must be careful: -6 does NOT have remainder 1. We can write -6 = -2 \cdot 5 + 4 or -6 = -1 \cdot 5 - 1, but only one of these is a “correct” instance of the division algorithm: \( r = 4 \) since we need \( r \) to be non-negative. So in fact, to get \( r = 1 \), we would could have -4, or -9, etc. Thus we get the remainder class:

\[
\{ \ldots, -14, -9, -4, 1, 6, 11, 16, 21, \ldots \}
\]

There are three more to go. The remainder classes for 2, 3, and 4 are, respectively:

\[
\{ \ldots, -13, -8, -3, 2, 7, 12, 17, 22, \ldots \}
\]
\[
\{ \ldots, -12, -7, -2, 3, 8, 13, 18, 23, \ldots \}
\]
\[
\{ \ldots, -11, -6, -1, 4, 9, 14, 19, 24, \ldots \}
\]

Note that in the example above, *every* integer is in exactly one remainder class. The technical way to say this is that the remainder classes modulo \( b \) form a partition of the integers.\(^1\) The most important fact about partitions, is that it is possible to define an equivalence relation from a partition: this is a relationship between pairs of numbers which acts in all the important ways like the “equals” relationship.\(^2\)

---

\(^1\)It is possible to develop a mathematical theory of partitions, prove statements about all partitions in general and then apply those observations to our case here.

\(^2\)Again, there is a mathematical theory of equivalence relations which applies in many more instances than the one we look at here.
All fun technical language aside, the idea is really simple. If two numbers belong to the same remainder class, then in some way, they are the same. That is, they are the same up to division by \( b \). In the case where \( b = 5 \) above, the numbers 8 and 23, while not the same number, are the same when it comes to dividing by 5, because both have remainder 3.

It matters what the divisor is: 8 and 23 are the same up to division by 5, but not up to division by 7, since 8 has remainder of 1 when divided by 7 while 23 has a remainder of 2.

With all this in mind, let’s introduce some notation. We want to say that 8 and 23 are basically the same, even though they are not equal. So it would be wrong to say \( 8 = 23 \). Instead, we write \( 8 \equiv 23 \). But this is not always true - it works if we are thinking division by 5, so that needs to be there as well. So what we will actually write is this:

\[
8 \equiv 23 \pmod{5}
\]

which is read, “8 is congruent to 23 modulo 5.” Of course then we could observe that

\[
8 \not\equiv 23 \pmod{7}
\]

**Congruence Modulo \( n \)**

We say \( a \) is congruent to \( b \) modulo \( n \), and write,

\[
a \equiv b \pmod{n}
\]

provided \( a \) and \( b \) have the same remainder when divided by \( n \). In other words, provided \( a \) and \( b \) belong to the same remainder class modulo \( n \).

Many books define congruence modulo \( n \) slightly differently. They say that \( a \equiv b \pmod{n} \) if and only if \( n \mid a - b \). In other words, two numbers are congruent modulo \( n \), if their difference is a multiple of \( n \). So which definition is correct? Turns out, it doesn’t matter - they are equivalent.

To see why, consider two numbers \( a \) and \( b \) which are congruent modulo \( n \). Then \( a \) and \( b \) have the same remainder when divided by \( n \). So we have:

\[
a = q_1n + r \quad \quad b = q_2n + r
\]

Here the two \( r \)'s really are the same. Consider what we get when we take the difference of \( a \) and \( b \):

\[
a - b = q_1n + r - (q_2n + r) = q_1n - q_2n = (q_1 - q_2)n
\]

So \( a - b \) is a multiple of \( n \), or equivalently, \( n \mid a - b \).

On the other hand, if we assume first that \( n \mid a - b \), so \( a - b = kn \), then consider what happens if we divide each term by \( n \). Dividing \( a \) by \( n \) will leave some remainder, as will dividing \( b \) by \( n \). However, dividing \( kn \) by \( n \) will leave 0 remainder. So the remainders on the left hand side must cancel out. That is, the remainders must be the same.

Thus we have:
Congruence and Divisibility
For any integers \(a, b,\) and \(n,\) we have
\[a \equiv b \pmod{n}\]
if and only if
\[n \mid a - b\]

It will also be useful to switch back and forth between congruences and regular equations. The above fact helps with this. We know that \(a \equiv b \pmod{n}\) if and only if \(n \mid a - b,\) if and only if \(a - b = kn\) for some integer \(k.\) Rearranging that equation, we get \(a = b + kn.\) In other words, if \(a\) and \(b\) are congruent modulo \(n,\) then \(a\) is \(b\) more than some multiple of \(n.\) This conforms with our earlier observation that all the numbers in a particular remainder class are the same amount larger than the multiples of \(n.\)

Congruence and Equality
For any integers \(a, b,\) and \(n,\) we have
\[a \equiv b \pmod{n}\]
if and only if
\[a = b + kn\] for some integer \(k\)

6.3 Properties of Congruence

We said earlier that congruence modulo \(n\) behaves in many important ways the same way equality does. Specifically, we could prove that congruence modulo \(n\) is an equivalence relation, which would require checking the following three facts:

Congruence Modulo \(n\) is an Equivalence Relation
Given any integers \(a, b, c,\) and \(n,\) the following hold:
1. \(a \equiv a \pmod{n}\)
2. If \(a \equiv b \pmod{n}\) then \(b \equiv a \pmod{n}\)
3. If \(a \equiv b \pmod{n}\) and \(b \equiv c \pmod{n},\) then \(a \equiv c \pmod{n}\)

In other words, congruence modulo \(n\) is reflexive, symmetric, and transitive, so is an equivalence relation.

You should take a minute to convince yourself that each of the properties above actually hold of congruence. Try explaining each using both the remainder and divisibility definitions.

Next let’s consider how congruence behaves when doing basic arithmetic. We already know that if you subtract two congruent numbers, the result will be congruent to \(0\) (be a multiple of \(n\)). What if we add something congruent to \(1\) to something congruent to \(2?\) Will we get something congruent to \(3?\)
**Congruence and Arithmetic**

Suppose \( a \equiv b \pmod{n} \) and \( c \equiv d \pmod{n} \). Then the following hold:

1. \( a + c \equiv b + d \pmod{n} \)
2. \( a - c \equiv b - d \pmod{n} \)
3. \( ac \equiv bd \pmod{n} \)

The above facts might be written a little strangely, but the idea is simple. If we have a true congruence, and we add the same thing to both sides, the result is still a true congruence. This sounds like we are saying:

If \( a \equiv b \pmod{n} \) then \( a + c \equiv b + c \pmod{n} \).

Of course that is true as well - it is the special case where \( c = d \). But what we have works in more generality. Think of congruence as being “basically equal.” If we have two numbers which are basically equal, and we add basically the same thing to both sides, the result will basically be equal.

This seems reasonable. Is it really true? Let’s prove the first fact.

**Proof.** Suppose \( a \equiv b \pmod{n} \) and \( c \equiv d \pmod{n} \). That means \( a = b + kn \) and \( c = d + jn \) for integers \( k \) and \( j \). Add these equations:

\[
    a + c = b + d + kn + jn.
\]

But \( kn + jn = (k + j)n \), which is just a multiple of \( n \). So \( a + c = b + d + (j + k)n \), or in other words, \( a + c \equiv b + d \pmod{n} \)

The other two facts can be proved in a similar way.

One of the important consequences of these facts about congruences, is that we can basically replace any number in a congruence with any other number it is congruent to. Here are some examples to see how (and why) that works:

**Example:** Find the remainder of 3491 divided by 9.

**Solution:** We could do long division, but there is another way. We want to find \( x \) such that \( x \equiv 3491 \pmod{9} \). Now 3491 = 3000 + 400 + 90 + 1. Of course 90 \( \equiv 0 \pmod{9} \), so we can replace the 90 in the sum with 0. Why is this okay? We are actually subtracting the “same” thing from both sides:

\[
    x \equiv 3000 + 400 + 90 + 1 \pmod{9}
    \]

\[
    - 0 \equiv 90 \pmod{9}
    \]

\[
    x \equiv 3000 + 400 + 0 + 1 \pmod{9}
    \]

Next, note that 400 = 4 \cdot 100, and 100 \( \equiv 1 \pmod{9} \) (since 9 \( | 99 \)). So we can in fact replace the 400 with simply a 4. Again, we are appealing to our
claim that we can replace congruent elements, but we are really appealing to property 3 about the arithmetic of congruence: we know \(100 \equiv 1 \pmod{9}\), so if we multiply both sides by 4, we get \(400 \equiv 4 \pmod{9}\).

Similarly, we can replace 3000 with 3, since \(1000 = 1 + 999 \equiv 1 \pmod{9}\). So our original congruence becomes

\[
x \equiv 3 + 4 + 0 + 1 \pmod{9}
\]

\[
x \equiv 8 \pmod{9}
\]

Therefore 3491 divided by 9 has remainder 8.

The above example should convince you that the well known divisibility test for 9 is true: the sum of the digits of a number is divisible by 9 if and only if the original number is divisible by 9. In fact, we now know something more: any number is congruent to the sum of its digits, modulo 9\(^3\).

Let’s try another.

**Example:** Find the remainder when \(3^{123}\) is divided by 7.

**Solution:** Of course, we are working with congruence because we want to find the smallest positive \(x\) such that \(x \equiv 3^{123} \pmod{7}\). Now first write \(3^{123} = (3^3)^{41}\). We have:

\[
3^{123} = 27^{41} \equiv 6^{41} \pmod{7}
\]

since \(27 \equiv 6 \pmod{7}\). Notice further that \(6^2 = 36\) is congruent to 1 modulo 7. Thus we can simplify further:

\[
6^{41} = 6 \cdot (6^2)^{20} \equiv 6 \cdot 1^{20} \pmod{7}
\]

But \(1^{20} = 1\), so we are done:

\[
3^{123} \equiv 6 \pmod{7}
\]

The above example works. We are using the fact that if \(a \equiv b \pmod{n}\), then \(a^p \equiv b^p \pmod{n}\). This is just applying property 3 a bunch of times.

So far we have seen how to add, subtract and multiply with congruences. What about division? There is a reason we have waited to discuss it. It turns out that we cannot simply divide. In other words, even if \(ad \equiv bd \pmod{n}\), we do not know that \(a \equiv b \pmod{n}\). Consider, for example:

\[
18 \equiv 42 \pmod{8}
\]

This is true. Now 18 and 42 are both divisible by 6. However,

\[
3 \not\equiv 7 \pmod{8}
\]

While this doesn’t work, note that \(3 \equiv 7 \pmod{4}\). We cannot divide 8 by 6, but we can divide 8 by the greatest common factor of 8 and 6. Will this always work?

\(^3\)This works for 3 as well, but definitely not for any modulus in general.
Suppose \( ad \equiv bd \pmod{n} \) In other words, we have \( ad = bd + kn \) for some integer \( k \). Of course \( ad \) is divisible by \( d \), as is \( bd \). So \( kn \) must also be divisible by \( d \). Now if \( n \) and \( d \) have no common factors (other than 1), then we must have \( d \mid k \). But in general, if we try to divide \( kn \) by \( d \), we don’t know that we will get an integer multiple of \( n \). Some of the \( n \) might get divided as well. To be safe, let’s divide as much of \( n \) as we can. Take the largest factor of both \( d \) and \( n \), and cancel that out from \( n \). The rest of the factors of \( d \) will come from \( k \), no problem.

We will call the largest factor of both \( d \) and \( n \), the \( \gcd(d,n) \), for greatest common divisor. In our example above, \( \gcd(6,8) = 2 \) since the greatest divisor common to 6 and 8 is 2.

### Congruence and Division

Suppose \( ad \equiv bd \pmod{n} \). Then \( a \equiv b \pmod{\frac{n}{\gcd(d,n)}} \)

If \( d \) and \( n \) have no common factors then \( \gcd(d,n) = 1 \), so \( a \equiv b \pmod{n} \).

**Example:** Simplify the following congruences using division: (a) \( 24 \equiv 39 \pmod{5} \) and (b) \( 24 \equiv 39 \pmod{15} \).

**Solution:** (a) Both 24 and 39 are divisible by 3, and 3 and 5 have no common factors, so we get

\[
8 \equiv 13 \pmod{5}
\]

(b) Again, we can divide by 3. However, if doing so blindly gives us \( 8 \equiv 13 \pmod{15} \) which is no longer true. Instead, we must also divide the modulus 15 by the greatest common factor of 3 and 15, which is 3. Again we get

\[
8 \equiv 13 \pmod{5}
\]

### 6.4 Solving Congruences

Now that we have some algebraic rules to govern congruence relations, we can attempt to solve for an unknown in a congruence. For example, is there a value of \( x \) that satisfies,

\[
3x + 2 \equiv 4 \pmod{5},
\]

and if so, what is it?

In this example, since the modulus is small, we could simply try every possible value for \( x \). There are really one 5 to consider, since any integer that satisfied the congruence could be replaced with any other integer it was congruent to modulo 5. Here, \( x = 4 \) gives 14 which is indeed congruent to 4 modulo 5. This means that \( x = 9 \) and \( x = 14 \) and \( x = 19 \) and so on will each also be a solution because as we saw above, replacing any number in a congruence with a congruent number does not change the truth of the congruence.

So in this example, simply compute \( 3x + 2 \) for values of \( x \in \{0, 1, 2, 3, 4\} \). This gives 2, 5, 8, 11, and 14 respectively, for which only 14 is congruent to 4.
Let’s also see how you could solve this using our rules for algebra of congruences. Such an approach would be much simpler that the trail and error tactic if the modulus was larger. First, we know we can subtract 2 from both sides:

\[ 3x \equiv 2 \pmod{5}. \]

Then to divide both sides by 3, we first add 0 to both sides. Of course, on the right hand side, we want that 0 to be a 10 (yes, 10 really is 0 - they are congruent modulo 5). This gives,

\[ 3x \equiv 12 \pmod{5}. \]

Now divide both sides by 3. Since \( \gcd(3, 5) = 1 \), we do not need to change the modulus:

\[ x \equiv 4 \pmod{5}. \]

Notice that this in fact gives the general solution: not only can \( x = 4 \), but \( x \) can be any number which is congruent to 4. We can leave it like this, or write “\( x = 4 + 5k \) for any integer \( k \).”

Here are a few of more examples of this process.

**Example:** Solve the following congruences for \( x \).

1. \( 7x \equiv 12 \pmod{13} \)
2. \( 84x - 38 \equiv 79 \pmod{15} \)
3. \( 20x \equiv 23 \pmod{14} \)

**Solution:**

1. All we need to do here is divide both sides by 7. We add 13 to the right hand side repeatedly until we get a multiple of 7 (adding 13 is the same as adding 0, so this is legal). We get 25, 38, 51, 64, 77 - got it. So we have:

\[ 7x \equiv 12 \pmod{13} \]
\[ 7x \equiv 77 \pmod{13} \]
\[ x \equiv 11 \pmod{13} \]

2. Here, since we have numbers larger than the modulus, we can reduce them prior to applying any algebra. We have \( 84 \equiv 9, 38 \equiv 8 \) and \( 79 \equiv 4 \). Thus,

\[ 84x - 38 \equiv 79 \pmod{15} \]
\[ 9x - 8 \equiv 4 \pmod{15} \]
\[ 9x \equiv 12 \pmod{15} \]
\[ 9x \equiv 72 \pmod{15} \]

We got the 72 by adding \( 0 \equiv 60 \pmod{15} \) to both sides of the congruence. Now divide both sides by 9. However, since \( \gcd(9, 15) = 3 \), we must divide the modulus by 3 as well:

\[ x \equiv 8 \pmod{5} \]
So the solutions are those values which are congruent to 8, or equivalently 3, modulo 5. This means that in some sense there are 3 solutions modulo 15: 3, 8, and 13. We can write the solution:

\[ x \equiv 3 \pmod{15}; \quad x \equiv 8 \pmod{15}; \quad x \equiv 13 \pmod{15} \]

3. First reduce modulo 14:

\[ 6x \equiv 9 \pmod{14} \]

We could now divide both sides by 3, or try to increase 9 by a multiple of 14 to get a multiple of 6. If we divide by 3, we get,

\[ 2x \equiv 3 \pmod{14} \]

Now try adding multiples of 14 to 3, in hopes of getting a number we can divide by 2. This will not work! Every time we add 14 to the right side, the result will still be odd. We will never get an even number, so we will never be able to divide by 2. Thus there are no solutions to the congruence.

The last congruence above illustrates the way in which congruences might not have solutions. We could have seen this immediately in fact. Look at the original congruence:

\[ 20x \equiv 23 \pmod{14} \]

If we write this as an equation, we get

\[ 20x = 23 + 14k \]

or equivalently \( 20x - 14k = 23 \). This is a linear Diophantine equation, and one for which we can easily see there will be no solution. The left hand side will always be even, but the right hand side is odd. A similar problem would occur if the right hand side was divisible by any number the left hand side was not.

So in general, given the congruence

\[ ax \equiv b \pmod{n}, \]

if \( a \) and \( n \) are divisible by a number which \( b \) is not divisible by, then there will be no solutions. In fact, we really only need to check one divisor of \( a \) and \( n \): the greatest common divisor. Thus, a more compact way to say this is:

<table>
<thead>
<tr>
<th>Congruences with no solutions</th>
</tr>
</thead>
<tbody>
<tr>
<td>If ( \gcd(a, n) \nmid b ), then ( ax \equiv b \pmod{n} ) has no solutions.</td>
</tr>
</tbody>
</table>
6.5 Solving Linear Diophantine Equations

We now have the tools need to solve linear Diophantine equations such as

\[ 51x + 87y = 123. \]

The general strategy will be to convert the equation to a congruence, then solve that congruence.\(^5\) Let’s work this particular example to see how this might go.

First, check if perhaps there are no solutions because a divisor of 51 and 87 is not a divisor of 123. Really, we just need to check whether \( \gcd(51, 87) \mid 123 \). This greatest common divisor is 3, and yes 3 \( \mid 123 \). At this point, we might as well factor out this greatest common divisor. So instead, we will solve:

\[ 17x + 29y = 41 \]

Now observe that if there are going to be solutions, then for those values of \( x \) and \( y \), the two sides of the equation must have the same remainder as each other, no matter what we divide by. In particular, if we divide both sides by 17, we must get the same remainder. Thus we can safely write,

\[ 17x + 29y \equiv 41 \pmod{17} \]

We choose 17 because 17\( x \) will have remainder 0. This will allow us to reduce the congruence to just one variable. We could have also moved to a congruence modulo 29, although there is usually a good reason to select the smaller choice, as this will allow us to reduce the other coefficient. In our case, we reduce the congruence as follows:

\[
\begin{align*}
17x + 29y & \equiv 41 \pmod{17} \\
0x + 12y & \equiv 7 \pmod{17} \\
12y & \equiv 24 \pmod{17} \\
y & \equiv 2 \pmod{17}
\end{align*}
\]

Now at this point we know \( y = 2 + 17k \) will work for any integer \( k \). If we haven’t made a mistake, we should be able to plug this back into our original Diophantine equation to find \( x \):

\[
17x + 29(2 + 17k) = 41
\]

\[
17x = -17 - 29 \cdot 17k
\]

\[
x = -1 - 29k
\]

We have now found all solutions to the Diophantine equation. For each \( k \), \( x = -1 - 29k \) and \( y = 2 + 17k \) will satisfy the equation. We could check this for a few cases. If \( k = 0 \), the solution is \((-1, 2)\), and yes, \(-17 + 2 \cdot 29 = 41\). If \( k = 3 \), the solution is \((-88, 53)\). If \( k = -2 \), we get \((57, -32)\).

To summarize this process, to solve \( ax + by = c \), we,

\(^5\)This is certainly not the only way to proceed. A more common technique would be to apply the Euclidean algorithm. Our way can be a little faster, and is presented here primarily for variety.
1. Divide both sides of the equation by \( \gcd a, b \) (if this does not leave the right hand side as an integer, there are no solutions). Let’s assume that \( ax + by = c \) has already been reduced in this way.

2. Pick the smaller of \( a \) and \( b \) (here, assume it is \( b \)), and convert to a congruence modulo \( b \):

   \[
   ax + by \equiv c \pmod{b}
   \]

   This will reduce to a congruence with one variable, \( x \):

   \[
   ax \equiv c \pmod{b}
   \]

3. Solve the congruence as we did in the previous section. Write your solution as an equation, such as,

   \[
   x = n + kb
   \]

4. Plug this into the original Diophantine equation, and solve for \( y \).

5. If we want to know solutions in a particular range (for example, \( 0 \leq x, y \leq 20 \)), pick different values of \( k \) until you have all required solutions.

Here is another example.

**Example:** How can you make $6.37 using just 5-cent and 8-cent stamps? What is the smallest and largest number of stamps you could use?

**Solution:** First we need a Diophantine equation. We will work in numbers of cents. Let \( x \) be the number of 5-cent stamps, and \( y \) be the number of 8-cent stamps. We have:

\[
5x + 8y = 637.
\]

Convert to a congruence and solve:

\[
8y \equiv 367 \pmod{5}
\]

\[
8y \equiv 2 \pmod{5}
\]

\[
8y \equiv 32 \pmod{5}
\]

\[
y \equiv 4 \pmod{5}
\]

Thus \( y = 4 + 5k \). Then \( 5x + 8(4 + 5k) = 637 \), so \( x = 121 - 8k \).

This says that one way to make $6.37 is to take 121 of the 5-cent stamps and 4 of the 8-cent stamps. To find the smallest and largest number of stamps, try different values of \( k \).

<table>
<thead>
<tr>
<th>( k )</th>
<th>( (x, y) )</th>
<th>Stamps</th>
</tr>
</thead>
<tbody>
<tr>
<td>-1</td>
<td>(129, -1)</td>
<td>not possible</td>
</tr>
<tr>
<td>0</td>
<td>(121, 4)</td>
<td>125</td>
</tr>
<tr>
<td>1</td>
<td>(113, 9)</td>
<td>122</td>
</tr>
<tr>
<td>2</td>
<td>(105, 13)</td>
<td>119</td>
</tr>
<tr>
<td>...</td>
<td>...</td>
<td>...</td>
</tr>
</tbody>
</table>
Of course this is no surprise - having the most stamps means we have as many 5-cent stamps as possible, and to get the smallest number of stamps would require have the least number of 5-cent stamps. To minimize the number of 5-cent stamps, we want to pick \( k \) so that \( 121 - 8k \) is as small as possible (but still positive). When \( k = 15 \), we have \( x = 1 \) and \( y = 79 \).

Therefore, to make \$6.37, you can use as few as 80 stamps (1 5-cent stamp and 79 8-cent stamps) or as many as 125 stamps (121 5-cent stamps and 4 8-cent stamps).

Using this method, as long as you can solve linear congruences in one variable, you can solve linear Diophantine equations of two variables. There are times though that solving the linear congruence is a lot of work. For example, suppose you need to solve,

\[
13x \equiv 6 \pmod{51}
\]

You could keep adding 51 to the right side until you get a multiple of 13: You would get 57, 108, 159, 210, 261, 312, and 312 is the first of these that is divisible by 13. This works, but is really too much work. Instead we could convert back to a Diophantine equation:

\[
13x = 6 + 51k
\]

Now solve this like we have in this section. Write it as a congruence modulo 13:

\[
0 \equiv 6 + 51k \pmod{13}
\]

\[
-12k \equiv 6 \pmod{13}
\]

\[
2k \equiv -1 \pmod{13}
\]

\[
2k \equiv 12 \pmod{13}
\]

\[
k \equiv 6 \pmod{13}
\]

so \( k = 6 + 13j \). Now go back and figure out \( x \):

\[
13x = 6 + 51(6 + 13j)
\]

\[
x = 24 + 51j
\]

Of course you could do this switching back and forth between congruences and Diophantine equations as many times as you like. If you only used this technique, you would essentially replicate the Euclidean algorithm (a more standard way to solve Diophantine equations).

**Exercises**

1. Suppose \( a, b, \) and \( c \) are integers. Prove that if \( a \mid b \), then \( a \mid bc \).

2. Suppose \( a, b, \) and \( c \) are integers. Prove that if \( a \mid b \) and \( a \mid c \) then \( a \mid b+c \) and \( a \mid b-c \).

3. Write out the remainder classes for \( n = 4 \).
4. Let $a$, $b$, $c$, and $n$ be integers. Prove that if $a \equiv b \pmod{n}$ and $c \equiv d \pmod{n}$, then $a - c \equiv b - d \pmod{n}$.

5. Find the remainder of $3^{456}$ when divided by

(a) 2
(b) 5
(c) 7
(d) 9

6. Determine which of the following congruences have solutions, and find any solutions (between 0 and the modulus) by trial and error.

(a) $4x \equiv 5 \pmod{6}$
(b) $4x \equiv 5 \pmod{7}$
(c) $6x \equiv 3 \pmod{9}$
(d) $6x \equiv 4 \pmod{9}$
(e) $x^2 \equiv 2 \pmod{4}$
(f) $x^2 \equiv 2 \pmod{7}$

7. Solve the following congruences (describe the general solution).

(a) $5x + 8 \equiv 11 \pmod{22}$
(b) $6x \equiv 4 \pmod{10}$
(c) $4x \equiv 24 \pmod{30}$
(d) $341x \equiv 2941 \pmod{9}$

8. I’m thinking of a number. If you multiply my number by 7, add 5, and divide the result by 11, you will be left with a remainder of 2. What remainder would you get if you divided my original number by 11?

9. Solve the following linear Diophantine equations, using modular arithmetic (describe the general solutions).

(a) $6x + 10y = 32$
(b) $17x + 8y = 31$
(c) $35x + 47y = 1$

10. You have a 13 oz. bottle and a 20 oz. bottle, with which you wish to measure exactly 2 oz. However, you have a limited supply of water. If any water enters either bottle and then gets dumped out, it is gone forever. What is the least amount of water you can start with and still complete the task?
Appendix A

Solutions to Exercises
Solutions for Chapter 1

\[
\begin{array}{c|c|c}
P & Q & (P \lor Q) \rightarrow (P \land Q) \\
\hline
T & T & T \\
1.1. & T & F \\
& F & T \\
& F & F \\
\end{array}
\]

\[
\begin{array}{c|c|c}
P & Q & \neg P \land (Q \rightarrow P) \\
\hline
T & T & F \\
1.2. & T & F \\
& F & T \\
& F & T \\
\end{array}
\]

1.3. Hint: Like above, only now you will need 8 rows instead of just 4.

1.4. Make a truth table for each and compare. The statements are logically equivalent.

1.5. Again, make two truth tables. The statements are logically equivalent.

1.6. The argument is valid. To see this, make a truth table which contains \( P \lor Q \) and \( \neg P \) (and \( P \) and \( Q \) of course). Look at the truth value of \( Q \) in each of the rows that have \( P \lor Q \) and \( \neg P \) true.

1.7. The argument form is valid. Again, make a truth table containing the premises and conclusion - look at the rows for which the premises are true.

1.8. The argument is NOT valid. If you make a truth table containing the premises and conclusion, there will be a row with both premises true but the conclusion false. For example, if \( P \) and \( Q \) are false and \( R \) is true, then \( P \land Q \) is false, so \( (P \land Q) \rightarrow R \) is true. Also \( \neg P \) is true, so \( \neg P \lor \neg Q \) is true. However, \( \neg R \) is false.

1.9. (a) \( P \): it’s your birthday; \( Q \): there will be cake. \( (P \lor Q) \rightarrow Q \)
(b) Hint: you should get three T’s and one F.
(c) Only that there will be cake.
(d) It’s NOT your birthday!
(e) It’s your birthday, but the cake is a lie.

1.10. (a) \( P \land Q \)
(b) \( P \rightarrow \neg Q \)
(c) Jack passed math or Jill passed math (or both).
(d) If Jack and Jill did not both pass math, then Jill did.
(e) (a) Nothing else.
     (b) Jack did not pass math either.
1.11. (a) Three statements: \( P \lor S, S \rightarrow Q, (P \lor Q) \rightarrow R \). You could also connect the first two with a \( \land \).

(b) He cannot be lying about all three sentences, so he is telling the truth.

(c) No matter what, Geoff wants ricotta. If he doesn’t have quail, then he must have pepperoni but not sausage.

1.12. (a) If Oscar drinks milk, then he eats Chinese food.

(b) If Oscar does not drink milk, then he does not eat Chinese food.

(c) Yes. The original statement would be false too.

(d) Nothing. The converse need not be true.

(e) He does not eat Chinese food. The contrapositive would be true.

1.13. (a) \( P \land Q \)

(b) \( (P \lor Q) \lor (Q \land \neg R) \)

(c) F. Or \( (P \land Q) \lor (R \land \neg R) \)

(d) Either Sam is a woman and Chris is a man, or Chris is a woman.

1.14. The statements are equivalent to the...

(a) converse.

(b) implication.

(c) neither.

(d) implication.

(e) converse.

(f) converse.

(g) implication.

(h) converse.

(i) converse.

(j) converse (in fact, this IS the converse).

(k) implication (the statement is the contrapositive of the implication).

(l) neither.

1.15. Hint: of course there are many answers. It helps to assume that the statement is true and the converse is NOT true. Think about what that means in the real world and then start saying it in different ways. Some ideas: use necessary and sufficient language, use “only if,” consider negations, use “or else” language.

1.16. (a) \( \neg \exists x(E(x) \land O(x)) \)

(b) \( \forall x(E(x) \rightarrow O(x + 1)) \)
(c) \( \exists x (P(x) \land E(x)) \) (where \( P(x) \) means “\( x \) is prime”)

(d) \( \forall x \forall y \exists z (x < z < y \lor y < z < x) \)

(e) \( \forall x \neg \exists y (x < y < x + 1) \)

1.17. (a) Any even number plus 2 is an even number.

(b) For any \( x \) there is a \( y \) such that \( \sin(x) = y \). In other words, every number \( x \) is in the domain of sine.

(c) For every \( y \) there is an \( x \) such that \( \sin(x) = y \). In other words, every number \( y \) is in the range of sine (which is false).

(d) For any numbers, if the cubes of two numbers are equal, then the numbers are equal.

1.18. (a) \( \forall x \exists y (O(x) \land \neg E(y)) \)

(b) \( \exists x \forall y (x \geq y \lor \forall z (x \geq z \land y \geq z)) \)

(c) There is a number \( n \) for which every other number is strictly greater than \( n \).

(d) There is a number \( n \) which is not between any other two numbers.

1.19. (a) \( A \cap B = \{3, 4, 5\} \).

(b) \( A \cup B = \{1, 2, 3, 4, 5, 6, 7\} \).

(c) \( A \setminus B = \{1, 2\} \).

(d) Yes.

(e) No.

1.20. (a) \( A \cap B = \{4, 6, 8, 10, 12\} \)

(b) \( A \cup B = \{x \in \mathbb{N} : (3 \leq x \leq 13) \lor x \text{ is even}\} \). (the set of all natural numbers which are either even or between 3 and 13 inclusive).

(c) \( B \cap C = \emptyset \).

(d) \( B \cup C = \mathbb{N} \).

1.21. For example, \( A = \{2, 3, 5, 7, 8\} \) and \( B = \{3, 5\} \).

1.22. Let \( A = \{1, 2, 3\} \) and \( B = \{1, 2, 3, 4, 5, \{1, 2, 3\}\} \)

1.23. (a) No.

(b) No.

(c) \( 2\mathbb{Z} \cap 3\mathbb{Z} \) is the set of all integers which are multiples of both 2 and 3 (so multiples of 6). Therefore \( 2\mathbb{Z} \cap 3\mathbb{Z} = \{x \in \mathbb{Z} : \exists y \in \mathbb{Z}(x = 6y)\} \).

(d) \( 2\mathbb{Z} \cup 3\mathbb{Z} \).

1.24. The set of primes.
1.25. (a) $A \cup B$:

(b) $(A \cup B)^c$:

(c) $A \cap (B \cup C)$:

(d) $(A \cap B) \cup C$:

(e) $\overline{A} \cap B \cap C$:

(f) $(A \cup B) \setminus C$:
1.26. For example, $A \cup B \cap \overline{(A \cap B)}$. Note that $\overline{A \cap B}$ would almost work, but also contain the area outside of both circles.

1.27. (a) 34.
(b) 103.
(c) 8.

1.28. $\mathcal{P}(A) = \{\emptyset, \{a\}, \{b\}, \{c\}, \{a,b\}, \{a,c\}, \{b,c\}, \{a,b,c\}\}.$

1.29. There are 10 singletons. There are 45 doubletons (because $45 = 9 + 8 + 7 + \cdots + 2 + 1$).

1.30. $\{2, 3, 5\}, \{1, 2, 3, 5\}, \{2, 3, 4, 5\}, \{2, 3, 5, 6\}, \{1, 2, 3, 4, 5\}, \{1, 2, 3, 5, 6\}, \{2, 3, 4, 5, 6\},$ and $\{1, 2, 3, 4, 5, 6\}$.

1.31. For example $A = \{1, 2, 3, 4\}$ and $B = \{5, 6, 7, 8, 9\}$.

1.32. For example, $A = \{1, 2, 3\}$ and $B = \{2, 3, 4, 5\}$.

1.33. If $R$ is the set of red cards and $F$ is the set of face cards, we want to find $|R \cup F|$. This is not simply $|R| + |F|$ because there are 6 cards which are both red and a face card; $|R \cap F| = 6$. We find $|R \cup F| = 32$. 
Solutions for Chapter 2

2.1. This is asking for the number of edges in $K_{10}$. Each vertex (person) has degree (shook hands with) 9 (people). So the sum of the degrees is 90. However, the degrees count each edge (handshake) twice, so there are 45 edges in the graph. That is how many handshakes took place.

2.2. It is possible for everyone to be friends with exactly 2 people - you could arrange the 5 people in a circle and say that everyone is friends with the two people on either side of them (so you get the graph $C_5$). However, it is not possible for everyone to be friends with 3 people - that would lead to a graph with an odd number of odd degree vertices which is impossible - the sum of the degrees must be even.

2.3. This is a question about finding Euler paths. Draw a graph with a vertex in each state, and connect vertices if their states share a border. Exactly two vertices will have odd degree - the vertices for Nevada and Utah. Thus you must start your road trip at in one of those states and end it in the other.

2.4. The first and the third graphs are the same, but the middle graph is different.

2.5. The first (and third) graphs contain an Euler path. All the graphs are planar.

2.6. Yes. For example, both graphs below contain 6 vertices, 7 edges, and have degrees $(2,2,2,2,3,3)$.

![Graphs](image)

2.7. (a) $K_4$ does not have an Euler path or circuit.
(b) $K_5$ has an Euler circuit (so also an Euler path).
(c) $K_{5,7}$ does not have an Euler path or circuit.
(d) $K_{2,7}$ has an Euler path but not an Euler circuit.
(e) $C_7$ has an Euler circuit (it is a circuit graph!)
(f) $P_7$ has an Euler path but no Euler circuit.

2.8. When $n$ is odd, $K_n$ contains an Euler circuit. This is because every vertex has degree $n - 1$, so an odd $n$ results in all degrees being even.

2.9. If both $m$ and $n$ are even, then $K_{m,n}$ has an Euler circuit. When both are odd, there is no Euler path or circuit. If one is 2 and the other is odd, then there is an Euler path but not an Euler circuit.

2.10. Three of the graphs are bipartite. The one which is not is $C_7$ (second from the right).
2.11. $C_n$ is bipartite if and only if $n = 1$ or is even.

2.12. For example, $K_5$.

2.13. For example, $K_{3,3}$.


2.15. Yes. According to Euler’s formula it would have 2 faces. It does. The only such graph is $C_{10}$.

2.16. $G$ has 10 edges. It could be planar, and then it would have 6 faces.

2.17. 2, since the graph is bipartite. One color for the top set of vertices, another color for the bottom set of vertices.

2.18. For example, $K_6$. If the chromatic number is 6, then the graph is not planar - the 4-color theorem states that all planar graphs can be colored with 4 or fewer colors.

2.19. The chromatic numbers are 2, 3, 4, 5, and 3 respectively from left to right.

2.20. (a) Only if $n \geq 6$ and is even.

(b) None.

(c) 12. Such a graph would have $\frac{5n}{2}$ edges. If the graph is planar, then $n - \frac{5n}{2} + F = 2$ so there would be $\frac{4 + 3n}{2}$ faces. Also, we must have $3F \leq 2E$, since the graph is simple. So we must have $3 \frac{4 + 3n}{2} \leq 5n$. Solving for $n$ gives $n \geq 12$.  

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Chapter 3 Solutions

Solutions for Chapter 3

3.1. 255.

3.2. 8.

3.3. 15.

3.4. (a) $2^8 = 256$. You have two choices for each tie - wear it or don’t.

(b) You have 7 choices for regular ties (the 8 choices less the “no regular tie” option) and 31 choices for bow ties (32 total minus the “no bow tie” option). Thus total you have $7 \cdot 31 = 217$.

(c) $\binom{3}{2}\binom{5}{3} = 30$

(d) $5! = 120$

3.5. (a) 16 is the number of choices you have if you want to watch one movie, either a comedy or horror flick.

(b) 63 is the number of choices you have if you will watch two movies, first a comedy and then a horror.

3.6. $0 \leq |A \cap B| \leq 10$ and $15 \leq |A \cup B| \leq 25$.

3.7. $|A \cup B| + |A \cap B| = 13$

3.8. 39.

3.9. $|(A \cup C) \cap \bar{B}| = 44$. Use a Venn diagram.

3.10. One possibility: $(A \cup B) \cap C$.

3.11. (a) $8^5$, since you select from 8 letters 5 times.

(b) $P(8, 5) = 8 \cdot 7 \cdot 6 \cdot 5 \cdot 4$. After selecting a letter, you have fewer letters to select for the next one.

(c) 64 - you need to select the 4th and 5th letters.

(d) $64 + 64 - 0 = 128$. There are 64 words which start with “aha” and another 64 words that end with “bah.” Perhaps we over counted the words that both start with “aha” and end with “bah” but since the words are only 5 letters long, there are no such words.

(e) $P(8, 5) - 3 \cdot P(5, 2) = 6660$ - all the words minus the bad ones. The taboo word can be in any of three positions (starting with letter 1, 2, or 3) and for each position we must choose the other two letters (from the remaining 5 letters)

3.12. $\binom{10}{6} + \binom{10}{7} + \binom{10}{8} + \binom{10}{9} + \binom{10}{10} = 386$

3.13. Use the binomial theorem. $\binom{14}{9} + \binom{15}{6}2^9$.

3.14. (a) $2^6 = 64$
(b) \(2^3 = 8\). We need to select yes/no for each of the remaining three elements.
(c) \(2^3 = 8\). We need to decide yes/no for the three non-prime elements.
(d) \(2^6 - 2^3 = 56\). There are 8 subsets which do not contain any odd numbers.
(e) 9. We need to select one odd (3 choices) and one even (3 choices).

3.15. (a) \(\binom{14}{7}\)
(b) \(\binom{6}{2}\binom{8}{5}\)
(c) \(\binom{14}{7} - \binom{6}{2}\binom{8}{5}\)

3.16. (a) \(\binom{10}{3}\)
(b) \(2^{10}\)
(c) \(P(10, 5)\)

3.17. \(\binom{5}{2}\binom{7}{2}\)

3.18. (a) 5 (you need to skip one dot the top and the bottom).
(b) \(\binom{7}{2}\) - once you select the two dots on the top, the bottom two are determined.
(c) This is tricky - you need to worry about running out of space. One way to count: break into cases by the location of the top left corner. You get \(\binom{7}{2} + (\binom{7}{2} - 1) + (\binom{7}{2} - 3) + (\binom{7}{2} - 6) + (\binom{7}{2} - 10) + (\binom{7}{2} - 15)\)
(d) All of them

3.19. (a) \(\binom{20}{4}\binom{16}{4}\binom{12}{4}\binom{8}{4}\binom{4}{4}\)
(b) \(5!(\binom{15}{3}\binom{12}{3}\binom{9}{3}\binom{6}{3}\binom{3}{3})\)

3.20. 9! (there are 10 people seated around the table, but it does not matter where King Arthur sits, only who sits to his left, two seats to his left, and so on).

3.21. **Proof.** Question: How many subsets of \(A = 1, 2, 3, \ldots, n + 1\) contain exactly two elements?

**Answer 1:** We must choose 2 elements from \(n + 1\) choices, so there are \(\binom{n+1}{2}\) subsets.

**Answer 2:** We break this question down into cases, based on what the larger of the two elements in the subset is. The larger element can’t be 1, since we need at least one element smaller than it.

Larger element is 2: there is 1 choice for the smaller element.

Larger element is 3: there are 2 choices for the smaller element.

Larger element is 4: there are 3 choices for the smaller element.

And so on. When the larger element is \(n + 1\), there are \(n\) choices for the smaller element. Since each two element subset must be in exactly one of these cases, the total number of two element subsets is \(1 + 2 + 3 + \cdots + n\).
Answer 1 and answer 2 are both correct, so they must be equal. Therefore

\[ 1 + 2 + 3 + \cdots + n = \binom{n+1}{2} \]

QED

3.22. (a) She has \(\binom{15}{6}\) ways to select the 6 bride’s maids, and then for each way, has 6 choices for the maid of honor. Thus she has \(\binom{15}{6} \cdot 6\) choices.

(b) She has 15 choices for who will be her maid of honor. Then she needs to select 5 of the remaining 14 friends to be bride’s maids, which she can do in \(\binom{14}{5}\) ways. Thus she has \(15 \binom{14}{5}\) choices.

(c) We have answered the question (how many wedding parties can the bride choose from) in two ways. The first way gives the left hand side of the identity and the second way gives the right hand side of the identity. Therefore the identity holds.

3.23. (a) After the 1, we need to find a 5-bit string with one 1. There are \(\binom{5}{1}\) ways to do this.

(b) \(\binom{4}{1}\) (we need to pick 1 of the remaining 4 slots to be the second 1).

(c) \(\binom{3}{1}\)

(d) Yes. We still need strings starting with 0001 (there are \(\binom{2}{1}\) of these) and strings starting 00001 (there is only \(\binom{1}{1}\) = 1 of these).

(e) \(\binom{6}{2}\)

(f) An example of the Hockey Stick Theorem:

\[ \binom{1}{1} + \binom{2}{1} + \binom{3}{1} + \binom{4}{1} + \binom{5}{1} = \binom{6}{2} \]

3.24. (a) \(3^n\), since there are 3 choices for each of the \(n\) digits.

(b) 1, since all the digits need to be 2’s. However, we might write this as \(\binom{n}{0}\).

(c) There are \(\binom{n}{1}\) places to put the non-2 digit. That digit can be either a 0 or a 1, so there are \(2 \binom{n}{1}\) such strings.

(d) We must choose two slots to fill with 0’s or 1’s. There are \(\binom{n}{2}\) ways to do that. Once the slots are picked, we have two choices for the first slot (0 or 1) and two choices for the second slot (0 or 1). So there are a total of \(2^2 \binom{n}{2}\) such strings.

(e) There are \(\binom{n}{k}\) ways to pick which slots don’t have the 2’s. Then those slots can be filled in \(2^k\) ways (0 or 1 for each slot). So there are \(2^k \binom{n}{k}\) such strings.

(f) These strings contain just 0’s and 1’s - so they are bit strings. There are \(2^n\) bit strings. But keeping with the pattern above, we might write this as \(2^n \binom{n}{k}\).
(g) We answer the question of how many length \( n \) ternary digit strings there are in two ways. First, each digit can be one of three choices, so the total number of strings is \( 3^n \). On the other hand, we could break the question down into cases by how many of the digits are 2’s. If they are all 2’s, then there are \( \binom{n}{0} \) strings. If all but one is a 2, then there are \( 2 \binom{n}{1} \) strings. If all but 2 of the digits are 2’s, then there are \( 2^2 \binom{n}{2} \) strings - we choose 2 of the \( n \) digits to be non-2, and then there are 2 choices for each of those digits. And so on for every possible number of 2’s in the string.

3.25. **Proof.** Question: How many \( k \)-letter words can you make using \( n \) different letters without repeating any letter?

**Answer 1:** There are \( n \) choices for the first letter, \( n - 1 \) choices for the second letter, \( n - 2 \) choices for the third letter, and so on until \( n - (k - 1) \) choices for the \( k \)th letter (since \( k - 1 \) letters have already been assigned at that point). The product of these numbers can be written \( \frac{n!}{(n-k)!} \) which is \( P(n,k) \).

**Answer 2:** First pick \( k \) letters to be in the word from the \( n \) choices. This can be done in \( \binom{n}{k} \) ways. Now arrange those letters into a word - there are \( k \) choices for the first letter, \( k - 1 \) choices for the second, and so on, for a total of \( k! \) arrangements of the \( k \) letters. Thus the total number of words is \( \binom{n}{k} k! \). QED

3.26. (a) \( \binom{18}{4} \). Each outcome can be represented by a sequence of 14 stars and 4 bars.

(b) \( \binom{13}{4} \). First put one ball in each bin. This leaves 9 stars and 4 bars.

(c) \( \binom{18}{4} - \left[ \binom{5}{1} \binom{11}{4} - \binom{5}{2} \binom{4}{4} \right] \). Subtract all the distributions for which one or more bins contain 7 or more balls.

3.27. (a) \( \binom{7}{2} \). After each variable gets 1 star for free, we are left with 5 stars and 2 bars.

(b) \( \binom{10}{2} \). We have 8 stars and 2 bars.

(c) \( \binom{19}{2} \). This problem is equivalent to finding the number of solutions to \( x' + y' + z' = 17 \) where \( x', y' \) and \( z' \) are non-negative. (In fact, we really just do a substitution. Let \( x = x' - 3, y = y' - 3 \) and \( z = z' - 3 \).

3.28. \( \binom{10}{5} \). We have 5 stars (the five dice) and 5 bars (the five switches between the number 1-6).

3.29. \( \binom{18}{3} \). Distribute 10 units to the variables before finding all solutions to \( x_1' + x_2' + x_3' + x_4' = 15 \) in non-negative integers.

3.30. There are 8 different functions. For example, \( f(1) = a, f(2) = a, f(3) = a; \) or \( f(1) = a, f(2) = b, f(3) = a, \) and so on. None of the functions are injective. Exactly 6 of the functions are surjective. No functions are both (since no functions here are injective).

3.31. There are nine functions - you have a choice of three outputs for \( f(1) \), and for each, you have three choices for the output \( f(2) \). Of these functions, 6 are injective, 0 are surjective, and 0 are both.
3.32. (a) $f$ is not injective, since $f(2) = f(5)$ - two different inputs have the same output.  
(b) $f$ is surjective, since every element of the codomain is an element of the range.

3.33. (a) $f$ is not injective, since $f(1) = 3$ and $f(4) = 3$.  
(b) $f$ is not surjective, since there is no input which gives 2 as an output.

3.34. (a) $f$ is injective, but not surjective.  
(b) $f$ is injective and surjective.  
(c) $f$ is injective, but not surjective.  
(d) $f$ is not injective, but is surjective.

3.35. (a) $|X| \leq |Y|$ - otherwise two or more of the elements of $X$ would need to map to the same element of $Y$.  
(b) $|X| \geq |Y|$ - otherwise there would be one or more elements of $Y$ which were never an output.  
(c) $|X| = |Y|$. This is the only way for both of the above to occur.

3.36. (a) Yes. (Can you give an example?)  
(b) Yes.  
(c) Yes.  
(d) Yes.  
(e) No.  
(f) No.

3.37. (a) $6^4 = 1296$, since there are six choices of where to send each of the 4 elements of the domain.  
(b) $P(6, 4) = 6 \cdot 5 \cdot 4 \cdot 3 = 360$, since outputs cannot be repeated.  
(c) None.  
(d) There are $5 \cdot 6^3$ functions for which $f(1) \neq a$ and another $5 \cdot 6^3$ functions for which $f(2) \neq b$. There are $5^2 \cdot 6^2$ functions for which both $f(1) \neq a$ and $f(2) \neq b$. So the total number of functions for which $f(1) \neq a$ or $f(2) \neq b$ or both is  
$$5 \cdot 6^3 + 5 \cdot 6^3 - 5^2 \cdot 6^2 = 1260$$

3.38. (a) $17^{10}$  
(b) $P(17, 10)$

3.39. $5^{10} - \left[ {5 \choose 4} 4^{10} - {5 \choose 3} 3^{10} + {5 \choose 2} 2^{10} - {5 \choose 1} 1^{10} \right]$  

3.40. $5! - \left[ {5 \choose 1} 4! - {5 \choose 3} 3! + {5 \choose 2} 2! - {5 \choose 1} 1! + {5 \choose 0} 0! \right]$. This is a sneaky way to ask for the number of derangements on 5 elements.
3.41. \( \binom{10}{6} (4! - \left[ \binom{4}{3} 3! - \binom{4}{2} 2! + \binom{4}{1} 1! - \binom{4}{0} 0! \right]) \). We choose 6 of the 10 ladies to get their own hat, and the multiply by the number of ways the remaining hats can be deranged.

3.42. (a) \( \binom{8}{3} \), after giving one present to each kid, you are left with 5 presents (stars) which need to be divide among the 4 kids (giving 3 bars).

(b) \( \binom{12}{3} \). You have 9 stars and 3 bars.

(c) \( 4^9 \). You have 4 choices for whom to give each present. This is like making a function from the set of present to the set of kids.

(d) \( 4^9 - \left[ \binom{4}{3} 3^9 - \binom{4}{2} 2^9 + \binom{4}{1} 1^9 \right] \). Now the function from the set of present to the set of kids must be surjective.

3.43. (a) Neither. \( \binom{14}{4} \).

(b) \( \binom{10}{4} \)

(c) \( P(10, 4) \), since order is important.

(d) Neither. Assuming you will wear each of the 4 ties on just 4 of the 7 days, without repeats: \( \binom{10}{4} P(7, 4) \).

(e) \( P(10, 4) \)

(f) \( \binom{10}{4} \)

(g) Neither. Since you could repeat letters: \( 10^4 \). If no repeats are allowed, it would be \( P(10, 4) \).

(h) Neither. Actually, “k” is the 11th letter of the alphabet, so the answer is 0. If “k” was among the first 10 letters, there would only be 1 way - write it down.

(i) Neither. Either \( \binom{9}{3} \) (if every kid gets an apple) or \( \binom{13}{3} \) (if appleless kids are allowed).

(j) Neither. Note that this could not be \( \binom{10}{4} \) since the 10 things and 4 things are from different groups. \( 4^{10} \)

(k) \( \binom{10}{4} \) - don’t be fooled by the “arrange” in there - you are picking 4 out of 10 spots to put the 1’s.

(l) \( \binom{10}{4} \) (assuming order is irrelevant).

(m) Neither. \( 16^{10} \) (each kid chooses yes or no to 4 varieties).

(n) Neither. 0.

(o) Neither. \( 4^{10} - \left[ \binom{4}{1} 3^{10} - \binom{4}{2} 2^{10} + \binom{4}{3} 1^{10} \right] \)

(p) Neither. \( 10 \cdot 4 \).

(q) Neither. \( 4^{10} \).

(r) \( \binom{10}{4} \) (which is the same as \( \binom{10}{6} \)).

(s) Neither. If all the kids were identical, and you wanted no empty teams, it would be \( \binom{10}{4} \). Instead, this will be the same as the number of surjective functions from a set of size 11 to a set of size 5.
(t) \binom{10}{4}
(u) \binom{10}{4}
(v) Neither. 4!
(w) Neither. It’s \binom{10}{4} if you won’t repeat any choices. If repetition is allowed, then this becomes \(x_1 + x_2 + \cdots + x_{10} = 4\), which has \binom{13}{9} solutions in non-negative integers.
(x) Neither. Since repetition of cookie type is allowed, the answer is 10\(^4\). Without repetition, you would have \(P(10, 4)\).
(y) \binom{10}{4} since that is equal to \binom{9}{4} + \binom{9}{3}.
(z) Neither. It will be a complicated (possibly PIE) counting problem.
Solutions for Chapter 4

4.1. (a) \[ a_n = n^2 + 1 \]
(b) \[ a_n = \frac{n(n+1)}{2} - 1 \]
(c) \[ a_n = \frac{(n+2)(n+3)}{2} + 2 \]
(d) \[ a_n = (n + 1)! - 1 \text{ (where } n! = 1 \cdot 2 \cdot 3 \cdots n \text{)} \]

4.2. (a) \[ F_n = F_{n-1} + F_{n-2} \text{ with } F_0 = 0 \text{ and } F_1 = 1. \]
(b) 0, 1, 2, 4, 7, 12, 20, ...
(c) \[ F_0 + F_1 + \cdots + F_n = F_{n+2} - 1 \]

4.3. 3, 10, 24, 52, 108, .... The recursive definition for 10, 24, 52, ... is \[ a_n = 2a_{n-1} + 4 \text{ with } a_1 = 10. \]

4.4. \(-1, -1, 1, 5, 11, 19, \ldots\) Thus the sequence 0, 2, 6, 12, 20, ... has closed formula \[ a_n = (n + 1)^2 - 3(n + 1) + 2. \]

4.5. (a) \(32.\)
(b) \[ a_n = 8 + 6(n - 1) \]
(c) \(30500.\)

4.6. (a) \(n + 2 \text{ terms.}\)
(b) \(6n + 1.\)
(c) \[ \frac{6n+8(n+2)}{2} \]

4.7. 68117

4.8. \(\frac{5 - 5 \cdot 3^{21}}{-2}\)

4.9. \(\frac{1 + \frac{3^{21}}{5 \pi}}{5/3}\)

4.10. For arithmetic: \(x = 55/3, y = 29/3.\) For geometric: \(x = 9\) and \(y = 3.\)

4.11. (a) \[ \sum_{k=1}^{n} 2k \]
(b) \[ \sum_{k=1}^{107} (1 + 4(k - 1)) \]
(c) \[ \sum_{k=1}^{50} \frac{1}{k} \]
(d) \[ \prod_{k=1}^{n} 2k \]
(e) \[ \prod_{k=1}^{100} \frac{k}{k+1} \]

4.12. (a) \[ \sum_{k=1}^{100} (3 + 4k) = 7 + 11 + 15 + \cdots + 403 \]

(b) \[ \sum_{k=0}^{n} 2^k = 1 + 2 + 4 + 8 + \cdots + 2^n \]

(c) \[ \sum_{k=2}^{50} \frac{1}{(k^2 - 1)} = 1 + \frac{1}{3} + \frac{1}{8} + \frac{1}{15} + \cdots + \frac{1}{2499} \]

(d) \[ \prod_{k=2}^{100} \frac{k^2}{(k^2 - 1)} = \frac{4}{3} \cdot \frac{9}{8} \cdot \frac{16}{15} \cdots \frac{10000}{9999} \]

(e) \[ \prod_{k=0}^{n} (2 + 3k) = (2)(5)(8)(11)(14) \cdots (2 + 3n) \]

4.13. (a) Hint: third differences are constant, so \( a_n = an^3 + bn^2 + cn + d \). Use the terms of the sequence to solve for \( a, b, c, \) and \( d \).

(b) \( a_n = n^2 - n \)

4.14. No. The sequence of differences is the same as the original sequence so no differences will be constant.

4.15. 171 and 341. \( a_n = a_{n-1} + 2a_{n-2} \) with \( a_0 = 3 \) and \( a_1 = 5 \). Closed formula: \( a_n = \frac{8}{3}2^n + \frac{1}{3}(-1)^n \)

4.16. By telescoping or iteration. \( a_n = 3 + 2^{n+1} \)

4.17. We claim \( a_n = 4^n \) works. Plug it in: \( 4^n = 3(4^{n-1}) + 4(4^{n-2}) \). This works - just simplify the right hand side.

4.18. By the Characteristic Root Technique. \( a_n = 4^n + (-1)^n \).

4.19. \( a_n = \frac{13}{5}4^n + \frac{12}{5}(-1)^n \)

4.20. The general solution is \( a_n = a + bn \) where \( a \) and \( b \) depend on the initial conditions.

(a) \( a_n = 1 + n \)

(b) For example, we could have \( a_0 = 21 \) and \( a_1 = 22 \).

(c) For every \( x \) - take \( a_0 = x - 9 \) and \( a_1 = x - 8 \).

4.21. \( a_n = \frac{19}{7}(-2)^n + \frac{9}{7}5^n \)

4.22. (a) \[ \frac{4}{1 - x} \]
(b) \( \frac{2}{(1-x)^2} \)

(c) \( \frac{2x^3}{(1-x)^2} \)

(d) \( \frac{1}{1-5x} \)

(e) \( \frac{1}{1+3x} \)

(f) \( \frac{1}{1-5x^2} \)

(g) \( \frac{x}{(1-x)^3} \)

4.23. (a) 0, 4, 4, 4, 4, 4, ...

(b) 1, 4, 16, 64, 256, ...

(c) 0, 1, -1, 1, -1, 1, -1, ...

(d) 0, 3, -6, 9, -12, 15, -18, ...

(e) 1, 3, 6, 9, 12, 15, ...

4.24. (a) The second derivative of \( \frac{1}{1-x} \) is \( \frac{2}{(1-x)^3} \) which expands to \( 2 + 6x + 12x^2 + 20x^3 + 30x^4 + \cdots \). Dividing by 2 gives the generating function for the triangular numbers.

(b) Compute \( A - xA \) and you get \( 1 + 2x + 3x^2 + 4x^3 + \cdots \) which can be written as \( \frac{1}{(1-x)^2} \). Solving for \( A \) gives the correct generating function.

(c) The triangular numbers are the sum of the first \( n \) numbers \( 1, 2, 3, 4, \ldots \). To get the sequence of partial sums, we multiply by \( \frac{1}{1-x} \). so this gives the correct generating function again.

4.25. Call the generating function \( A \). Compute \( A - xA = 4 + x + 2x^2 + 3x^3 + 4x^4 + \cdots \). Thus \( A - xA = 4 + \frac{x}{1-x} \). Solving for \( A \) gives \( \frac{4}{1-x} + \frac{x}{(1-x)^3} \).

4.26. \( \frac{1 + 2x}{1 - 3x + x^2} \)

4.27. Compute \( A - xA - x^2A \) and the solve for \( A \). The generating function will be \( \frac{x}{1-x-x^2} \).

4.28. \( \frac{x}{(1-x)(1-x-x^2)} \)

4.29. \( \frac{2}{1-5x} + \frac{7}{1+3x} \).
4.30. \( a_n = 3 \cdot 4^{n-1} + 1 \)

4.31. Hint: you should “multiply” the two sequences. Answer: 158.
Solutions for Chapter 5

5.1. Proof. We must prove that \(1 + 2 + 2^2 + 2^3 + \cdots + 2^n = 2^{n+1} - 1\) for all \(n \in \mathbb{N}\). Thus let \(P(n)\) be the statement \(1 + 2 + 2^2 + \cdots + 2^n = 2^{n+1} - 1\). We will prove that \(P(n)\) is true for all \(n \in \mathbb{N}\).

First we establish the base case, \(P(0)\), which claims that \(1 = 2^{0+1} - 1\). Since \(2^1 - 1 = 2 - 1 = 1\), we see that \(P(0)\) is true.

Now for the inductive case. Assume that \(P(k)\) is true for an arbitrary \(k \in \mathbb{N}\). That is, \(1 + 2 + 2^2 + \cdots + 2^k \equiv 2^{k+1} - 1\). We must show that \(P(k+1)\) is true (i.e., that \(1 + 2 + 2^2 + \cdots + 2^{k+1} = 2^{k+2} - 1\)). To do this, we start with the left hand side of \(P(k+1)\) and work to the right hand side:

\[
1 + 2 + 2^2 + \cdots + 2^k + 2^{k+1} = 2^{k+1} - 1 + 2^{k+1} \quad \text{by the inductive hypothesis}
\]
\[
= 2 \cdot 2^{k+1} - 1
\]
\[
= 2^{k+2} - 1
\]

Thus \(P(k+1)\) is true so by the principle of mathematical induction, \(P(n)\) is true for all \(n \in \mathbb{N}\). QED

5.2. Proof. Let \(P(n)\) be the statement “\(7^n - 1\) is a multiple of 6.” We will show \(P(n)\) is true for all \(n \in \mathbb{N}\).

First we establish the base case, \(P(0)\). Since \(7^0 - 1 = 0\), and 0 is a multiple of 6, \(P(0)\) is true.

Now for the inductive case. Assume \(P(k)\) holds for an arbitrary \(k \in \mathbb{N}\). That is, \(7^k - 1 = 6j\) for some integer \(j\). Now consider \(7^{k+1} - 1\):

\[
7^{k+1} - 1 = 7^{k+1} - 7 + 6 \quad \text{by cleverness: } -1 = -7 + 6
\]
\[
= 7(7^k - 1) + 6 \quad \text{factor out a 7 from the first two terms}
\]
\[
= 7(6j) + 6 \quad \text{by the inductive hypothesis}
\]
\[
= 6(7j + 1) \quad \text{factor out a 6}
\]

Therefore \(7^{k+1} - 1\) is a multiple of 6, or in other words, \(P(k+1)\) is true. Therefore by the principle of mathematical induction, \(P(n)\) is true for all \(n \in \mathbb{N}\). QED

5.3. Proof. Let \(P(n)\) be the statement \(1 + 3 + 5 + \cdots + (2n - 1) = n^2\). We will prove that \(P(n)\) is true for all \(n \geq 1\).

First the base case, \(P(1)\). We have \(1 = 1^2\) which is true, so \(P(1)\) is established.

Now the inductive case. Assume that \(P(k)\) is true for some fixed arbitrary \(k \geq 1\). That is, \(1 + 3 + 5 + \cdots + (2k - 1) = k^2\). We will now prove that \(P(k+1)\) is also true (i.e., that \(1 + 3 + 5 + \cdots + (2k + 1) = (k+1)^2\)). We start with the left hand side of \(P(k+1)\) and work to the right hand side:

\[
1 + 3 + 5 + \cdots + (2k - 1) + (2k + 1) = k^2 + (2k + 1) \quad \text{by the induction hypothesis}
\]
\[
= (k + 1)^2 \quad \text{by factoring}
\]
Thus $P(k + 1)$ holds, so by the principle of mathematical induction, $P(n)$ is true for all $n \geq 1$. QED

5.4. \textit{Proof.} Let $P(n)$ be the statement $F_0 + F_2 + F_4 + \ldots + F_{2n} = F_{2n+1} - 1$. We will show that $P(n)$ is true for all $n \geq 0$. First the base case is easy because $F_0 = 0$ and $F_1 = 1$ so $F_0 = F_1 - 1$. Now consider the inductive case. Assume $P(k)$ is true, that is, assume $F_0 + F_2 + F_4 + \ldots + F_{2k} = F_{2k+1} - 1$. To establish $P(k+1)$ we work from left to right:

$$F_0 + F_2 + F_4 + \ldots + F_{2k} + F_{2k+2} = F_{2k+1} - 1 + F_{2k+2}$$

by the inductive hypothesis

$$= F_{2k+1} + F_{2k+2} - 1$$

by the recursive definition of the Fibonacci numbers

$$= F_{2k+3} - 1$$

Therefore $F_0 + F_2 + F_4 + \ldots + F_{2k+2} = F_{2k+3} - 1$, which is to say $P(k+1)$ holds. Therefore by the principle of mathematical induction, $P(n)$ is true for all $n \geq 0$. QED

5.5. \textit{Proof.} Let $P(n)$ be the statement $2^n < n!$. We will show $P(n)$ is true for all $n \geq 4$. First, we check the base case and see that yes, $2^4 < 4!$ (as $16 < 24$) so $P(4)$ is true. Now for the inductive case. Assume $P(k)$ is true for an arbitrary $k \geq 4$. That is, $2^k < k!$. Now consider $P(k + 1)$: $2^{k+1} < (k + 1)!$. To prove this, we start with the left side and work to the right side.

$$2^{k+1} = 2 \cdot 2^k$$

by the inductive hypothesis

$$< 2 \cdot k!$$

since $k + 1 > 2$

$$< (k + 1) \cdot k!$$

$$= (k + 1)!$$

Therefore $2^{k+1} < (k + 1)!$ so we have established $P(k + 1)$. Thus by the principle of mathematical induction $P(n)$ is true for all $n \geq 4$. QED

5.6. The only problem is that we never established the base case. Of course, when $n = 0$, $0 + 3 \neq 0 + 7$.

5.7. \textit{Proof.} Let $P(n)$ be the statement that $n + 3 < n + 7$. We will prove that $P(n)$ is true for all $n \in \mathbb{N}$. First, note that the base case holds: $0 + 3 < 0 + 7$. Now assume for induction that $P(k)$ is true. That is, $k + 3 < k + 7$. We must show that $P(k + 1)$ is true. Now since $k + 3 < k + 7$, add 1 to both sides. This gives $k + 3 + 1 < k + 7 + 1$. Regrouping $(k + 1) + 3 < (k + 1) + 7$. But this is simply $P(k + 1)$. Thus by the principle of mathematical induction $P(n)$ is true for all $n \in \mathbb{N}$. QED

5.8. The problem here is that while $P(0)$ is true, and while $P(k) \to P(k + 1)$ for some values of $k$, there is at least one value of $k$ (namely $k = 99$) when that implication fails. For a valid proof by induction, $P(k) \to P(k + 1)$ must be true for all values of $k$ greater than or equal to the base case.

5.9. \textit{Proof.} Let $P(n)$ be the statement “there is a strictly increasing sequence $a_1, a_2, a_3, \ldots, a_n$ with $a_n < 100$.” We will prove $P(n)$ is true for all $n \geq 1$. First we establish the base case: $P(1)$ says there is a single number $a_1$ with $a_1 < 100$. This is true - take $a_1 = 0$. 

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Chapter 5 Solutions

Now for the inductive step, assume \( P(k) \) is true. That is there exists a strictly increasing sequence \( a_1, a_2, a_3, \ldots, a_k \) with \( a_k < 100 \). Now consider this sequence, plus one more term, \( a_{k+1} \) which is greater than \( a_k \) but less than 100. Such a number exists - for example, the average between \( a_k \) and 100. So then \( P(k + 1) \) is true, so we have shown that \( P(k) \rightarrow P(k+1) \). Thus by the principle of mathematical induction, \( P(n) \) is true for all \( n \in \mathbb{N} \).

QED

5.10. We once again failed to establish the base case: when \( n = 0 \), \( n^2 + n = 0 \) which is even, not odd.

5.11. \( \text{Proof.} \) Let \( P(n) \) be the statement “\( n^2 + n \) is even.” We will prove that \( P(n) \) is true for all \( n \in \mathbb{N} \). First the base case: when \( n = 0 \), we have \( n^2 + n = 0 \) which is even, so \( P(0) \) is true. Now suppose for induction that \( P(k) \) is true, that is, that \( k^2 + k \) is even. Now consider the statement \( P(k+1) \). Now \( (k + 1)^2 + (k + 1) = k^2 + 2k + 1 + k + 1 = k^2 + k + 2k + 2 \). By the inductive hypothesis, \( k^2 + k \) is even, and of course \( 2k + 2 \) is even. An even plus an even is always even, so therefore \( (k + 1)^2 + (k + 1) \) is even. Therefore by the principle of mathematical induction, \( P(n) \) is true for all \( n \in \mathbb{N} \).

QED

5.12. Further hint: the idea is to define the sequence so that \( a_n \) is less than the distance between the previous partial sum and 2. That way when you add it into the next partial sum, the partial sum is still less than 2. You could do this ahead of time, or use a clever \( P(n) \) in the induction proof. Let \( P(n) \) be the statement, “there is a sequence of positive real numbers \( a_1, a_2, a_3, \ldots, a_n \) such that \( a_1 + a_2 + a_3 + \cdots + a_n < 2 \).” The base case should be easy (just pick \( a_1 < 2 \)). For the inductive case, you know that \( a_1 + a_2 + \cdots + a_k < 2 \) so you just need to argue that you can find some \( a_{k+1} \) small enough to have \( a_1 + a_2 + \cdots + a_k + a_{k+1} < 2 \).

5.13. The base case should be easy - 0 is a power of 2. For the inductive case, you actually want to use strong induction. Suppose \( k \) is either a power of 2 or can be written as the sum of distinct powers of 2, for any \( k < n \). Now if \( n \) is a power of 2, we are done. If not, subtract the largest power of 2 from \( n \) possible. You get \( n - 2^x \), which is a smaller number, in fact smaller than both \( n \) and \( 2^x \). Thus \( n - 2^x \) is either a power of 2 or can be written as the sum of distinct powers of 2, but none of them are going to be \( 2^x \), so the together with \( 2^x \) we have written \( n \) as the sum of distinct powers of 2.

5.14. If \( n = 2 \), this should work out (so theirs’ is your base case). If we assume it works for \( k \) people (that the number of handshakes is \( \frac{k(k-1)}{2} \), what happens if a \( k+1 \)st person shows up. How many new handshakes take place? Now make this into a formal induction argument.

Note, we have already proven this without using induction, but this is fun too.

5.15. With just one edge, it should be easy to verify the equation. If the equation holds with \( k \) edges, what will happen to \( V, E, \) and \( F \) if we add and edge? This can be done in two ways - either adding a new vertex, or adding a new face. Explain how this works, and what happens to the equation.
5.16. Here’s the idea: since every entry in Pascal’s Triangle is the sum of the two entries above it, we can get the \( k+1 \)st row by adding up all the pairs of entry from the \( k \)th row. But doing this uses each entry on the \( k \)th row twice. Thus each time we drop to the next row, we double the total. Of course, row 0 has sum 1 = 2^0 (the base case). Now try to make this precise with a formal induction proof. You will use the fact that \( \binom{n}{k} = \binom{n-1}{k-1} + \binom{n-1}{k} \) for the inductive case.

5.17. To see why this works, try it on a copy of Pascal’s triangle. We are adding up the entries along a diagonal, starting with the 1 on the left hand side of the 4th row. Suppose we add up the first 5 entries on this diagonal. The claim is that the sum is the entry below and to the left of the last of these 5 entries. Note that if this is true, and we instead add up the first 6 entries, we will need to add the entry one spot to the right of the previous sum. But these two together give the entry below them, which is below and left of the last of the 6 entries on the diagonal.

If you follow that, you can see what is going on. But it is not a great proof. A formal induction proof is needed:

\[P(n)\] be the statement \( \binom{4}{0} + \binom{5}{1} + \binom{6}{2} + \cdots + \binom{4+k}{k} = \binom{5+k+1}{k+1}\). For the base case, consider \( n = 0 \). This says \( \binom{4}{0} = \binom{5}{0} \). Since these are both 1, the base case is true. Now for the inductive case, suppose \( P(k) \) is true. That is, \( \binom{4}{0} + \binom{5}{1} + \binom{6}{2} + \cdots + \binom{4+k}{k} = \binom{5+k}{k} \). If we add \( \binom{4+k+1}{k+1} \) to both sides, we get

\[
\binom{4}{0} + \binom{5}{1} + \binom{6}{2} + \cdots + \binom{4+k}{k} + \binom{5+k}{k+1} = \binom{5+k}{k} + \binom{5+k}{k+1}
\]

But \( \binom{5+k}{k} + \binom{5+k}{k+1} = \binom{5+k+1}{k+1} \). In other words, we have

\[
\binom{4}{0} + \binom{5}{1} + \binom{6}{2} + \cdots + \binom{4+k}{k} + \binom{5+k}{k+1} = \binom{5+k+1}{k+1}
\]

which is to say that \( P(k+1) \) is true.

Therefore, by the principle of mathematical induction, \( P(n) \) is true for all \( n \geq 0 \). QED

5.18. The idea here is that if we take the logarithm of \( a^n \), we can increase \( n \) by 1 if we multiply by another \( a \) (inside the logarithm). This results in adding 1 more \( \log(a) \) to the total.

\[P(n)\] be the statement \( \log(a^n) = n \log(a) \). The base case, \( P(2) \) is true, because \( \log(a^2) = \log(a \cdot a) = \log(a) + \log(a) = 2 \log(a) \), by the product rule for logarithms.

Now assume, for induction, that \( P(k) \) is true. That is, \( \log(a^k) = k \log(a) \). Consider \( \log(a^{k+1}) \). We have

\[
\log(a^{k+1}) = \log(a^k \cdot a) = \log(a^k) + \log(a) = k \log(a) + \log(a)
\]
with the last equality due to the inductive hypothesis. But this simplifies to \((k + 1) \log(a)\), establishing \(P(k + 1)\).

Therefore by the principle of mathematical induction, \(P(n)\) is true for all \(n \geq 2\). QED

5.19. Hint: You are allowed to assume the base case. For the inductive case, group all but the last function together as one sum of functions, then apply the usual sum of derivatives rule, and then the inductive hypothesis.

5.20. Hint: for the inductive step, we know by the product rule for two functions that

\[
(f_1f_2f_3\cdots f_k f_{k+1})' = (f_1f_2f_3\cdots f_k)'f_{k+1} + (f_1f_2f_3\cdots f_k)f_{k+1}'
\]

Then use the inductive hypothesis on the first summand, and distribute.
Solutions for Chapter 6

6.1. Proof. Suppose \( a \mid b \). Then \( b \) is a multiple of \( a \), or in other words, \( b = ak \) for some \( k \). But then \( bc = akc \), and since \( kc \) is an integer, this says \( bc \) is a multiple of \( a \). In other words, \( a \mid bc \). QED

6.2. Proof. Assume \( a \mid b \) and \( a \mid c \). This means that \( b \) and \( c \) are both multiples of \( a \), so \( b = am \) and \( c = an \) for integers \( m \) and \( n \). Then \( b + c = am + an = a(m + n) \), so \( b + c \) is a multiple of \( a \), or equivalently, \( a \mid b + c \). Similarly, \( b - c = am - an = a(m - n) \), so \( b - c \) is a multiple of \( a \), which is to say \( a \mid b - c \). QED

6.3. \( \{ \ldots, -8, -4, 0, 4, 8, 12, \ldots \}, \{ \ldots, -7, -3, 1, 5, 9, 13, \ldots \}, \{ \ldots, -6, -2, 2, 6, 10, 14, \ldots \}, \) and \( \{ \ldots, -5, -1, 3, 7, 11, 15, \ldots \} \).

6.4. Proof. Assume \( a \equiv b \pmod{n} \) and \( c \equiv d \pmod{n} \). This means \( a = b + kn \) and \( c = d + jn \) for some integers \( k \) and \( j \). Consider \( a - c \). We have:

\[
-a - c = b + kn - (d + jn) = b - d + (k - j)n
\]

In other words, \( a - c \) is \( b - d \) more than some multiple of \( n \), so \( a - c \equiv b - d \pmod{n} \). QED

6.5. (a) \( 3^{456} \equiv 1^{456} = 1 \pmod{2} \).
(b) \( 3^{456} = 9^{228} \equiv (-1)^{228} = 1 \pmod{5} \)
(c) \( 3^{456} = 9^{228} \equiv 2^{228} = 8^{76} \equiv 1^{76} = 1 \pmod{7} \)
(d) \( 3^{456} = 9^{228} \equiv 0^{228} = 0 \pmod{9} \)

6.6. For all of these, just plug in all integers between 0 and the modulus to see which, if any, work.

(a) No solutions.
(b) \( x = 3 \).
(c) \( x = 2, x = 5, x = 8 \).
(d) No solutions.
(e) No solutions.
(f) \( x = 3 \).

6.7. (a) \( x = 5 + 22k \) for \( k \in \mathbb{Z} \).
(b) \( x = 4 + 5k \) for \( k \in \mathbb{Z} \).
(c) \( x = 6 + 15k \) for \( k \in \mathbb{Z} \).
(d) Hint: first reduce each number modulo 9, which can be done by adding up the digits of the numbers. Answer: \( x = 2 + 9k \) for \( k \in \mathbb{Z} \).

6.8. We must solve \( 7x + 5 \equiv 2 \pmod{11} \). This gives \( x \equiv 9 \pmod{11} \). In general, \( x = 9 + 11k \), but when you divide any such \( x \) by 11, the remainder will be 9.
6.9. (a) Divide through by 2: \(3x + 5y = 16\). Convert to a congruence, modulo 3: \(5y \equiv 16 \pmod{3}\), which reduces to \(2y \equiv 1 \pmod{3}\). So \(y \equiv 2 \pmod{3}\) or \(y = 2 + 3k\). Plug this back into \(3x + 5y = 16\) and solve for \(x\), to get \(x = 2 - 5k\). So the general solution is \(x = 2 - 5k\) and \(y = 2 + 3k\) for \(k \in \mathbb{Z}\).

(b) \(x = 7 + 8k\) and \(y = -11 - 17k\) for \(k \in \mathbb{Z}\).

(c) \(x = -4 - 47k\) and \(y = 3 + 35k\) for \(k \in \mathbb{Z}\).

6.10. First, solve the Diophantine equation \(13x + 20y = 2\). The general solution is \(x = -6 - 20k\) and \(y = 4 + 13k\). Now if \(k = 0\), this correspond to filling the 20 oz. bottle 4 times, and emptying the 13 oz. bottle 6 times, which would require 80 oz. of water. Increasing \(k\) would require considerably more water. Perhaps \(k = -1\) would be better? Then we would have \(x = -6 + 20 = 14\) and \(y = 4 - 13 = -11\), which describes the solution where we fill the 13 oz. bottle 14 times, and empty the 20 oz. bottle 11 times. This would require 182 oz. of water. Thus the most efficient procedure is to repeatedly fill the 20 oz bottle, emptying it into the 13 oz bottle, and discarding full 13 oz. bottles, which requires 80 oz. of water.