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Preface

This is a free textbook for an undergraduate course on Discrete Structures for Computer Science students, which I have been teaching at Carleton University since the fall term of 2013. The material is offered as the second-year course COMP 2804 (Discrete Structures II). Students are assumed to have taken COMP 1805 (Discrete Structures I), which covers mathematical reasoning, basic proof techniques, sets, functions, relations, basic graph theory, asymptotic notation, and countability.

During a 12-week term with three hours of classes per week, I cover most of the material in this book, except for Chapter 2 which has been included so that students can review the material from COMP 1805.

Chapter 2 is largely taken from the free textbook *Introduction to Theory of Computation* by Anil Maheshwari and Michiel Smid, which is available at http://cg.scs.carleton.ca/~michiel/TheoryOfComputation/

Please let me know if you find errors, typos, simpler proofs, comments, omissions, or if you think that some parts of the book “need improvement”.
Chapter 1

Introduction

In this chapter, we introduce some problems that will be solved later in this book. Along the way, we recall some notions from discrete mathematics that you are assumed to be familiar with. These notions are reviewed in more detail in Chapter 2.

1.1 Ramsey Theory

Ramsey Theory studies problems of the following form: How many elements of a given type must there be so that we can guarantee that some property holds? In this section, we consider the case when the elements are people and the property is “there is a large group of friends or there is a large group of strangers”.

**Theorem 1.1.1** In any group of six people, there are

- three friends, i.e., three people who know each other,
- or three strangers, i.e., three people, none of which knows the other two.

In order to prove this theorem, we denote the six people by $P_1, P_2, \ldots, P_6$. Any two people $P_i$ and $P_j$ are either friends or strangers. We define the complete graph $G = (V, E)$ with vertex set

$$V = \{P_i : 1 \leq i \leq 6\}$$

and edge set

$$E = \{P_iP_j : 1 \leq i < j \leq 6\}.$$
Observe that $|V| = 6$ and $|E| = 15$. We draw each edge $P_iP_j$ as a straight-line segment according to the following rule:

- If $P_i$ and $P_j$ are friends, then the edge $P_iP_j$ is drawn as a *solid* segment.
- If $P_i$ and $P_j$ are strangers, then the edge $P_iP_j$ is drawn as a *dashed* segment.

In the example below, $P_3$ and $P_5$ are friends, whereas $P_1$ and $P_3$ are strangers.

Observe that a group of three friends corresponds to a solid triangle in the graph $G$, whereas a group of three strangers corresponds to a dashed triangle. In the example above, there is no solid triangle and, thus, there is no group of three friends. Since the triangle $P_2P_4P_5$ is dashed, there is a group of three strangers.

Using this terminology, Theorem 1.1.1 is equivalent to the following:

**Theorem 1.1.2** Consider a complete graph on six vertices, in which each edge is either solid or dashed. Then there is a solid triangle or a dashed triangle.

**Proof.** As before, we denote the vertices by $P_1, \ldots, P_6$. Consider the five edges that are incident on $P_1$. Using a proof by contradiction, it can easily be shown that one of the following two claims must hold:

- At least three of these five edges are solid.
- At least three of these five edges are dashed.
We may assume without loss of generality that the first claim holds. (Do you see why?) Consider three edges incident on $P_1$ that are solid and denote them by $P_1A$, $P_1B$, and $P_1C$.

If at least one of the edges $AB$, $AC$, and $BC$ is solid, then there is a solid triangle. In the left figure below, $AB$ is solid and we obtain the solid triangle $P_1AB$.

Otherwise, all edges $AB$, $AC$, and $BC$ are dashed, in which case we obtain the dashed triangle $ABC$; see the right figure above.

You should convince yourself that Theorem 1.1.2 also holds for complete graphs with more than six vertices. The example below shows an example of a complete graph with five vertices without any solid triangle and without any dashed triangle. Thus, Theorem 1.1.2 does not hold for complete graphs with five vertices. Equivalently, Theorem 1.1.1 does not hold for groups of five people.

What about larger groups of friends/strangers? Let $k \geq 3$ be an integer. The following theorem states that even if we take $\lfloor 2^{k/2} \rfloor$ people, we are not guaranteed that there is a group of $k$ friends or a group of $k$ strangers.

A $k$-clique in a graph is a set of $k$ vertices, any two of which are connected by an edge. For example, a 3-clique is a triangle.
Chapter 1. Introduction

Theorem 1.1.3 Let $k \geq 3$ and $n \geq 3$ be integers with $n \leq \lfloor 2^{k/2} \rfloor$. There exists a complete graph with $n$ vertices, in which each edge is either solid or dashed, such that this graph does not contain a solid $k$-clique and does not contain a dashed $k$-clique.

We will prove this theorem in Section 7.2 using elementary counting techniques and probability theory. This probably sounds surprising to you, because Theorem 1.1.3 does not have anything to do with probability. In fact, in Section 7.2, we will prove the following claim: Take $k = 20$ and $n = 1024$. If you go to the ByWard Market in downtown Ottawa and take a random group of $n$ people, then with very high probability, this group does not contain a subgroup of $k$ friends and does not contain a subgroup of $k$ strangers. In other words, with very high probability, every subgroup of $k$ people contains two friends and two strangers.

1.2 Sperner’s Theorem

Consider a set $S$ with five elements, say, $S = \{1, 2, 3, 4, 5\}$. Let $S_1, S_2, \ldots, S_m$ be a sequence of $m$ subsets of $S$, such that for all $i$ and $j$ with $i \neq j$,

$$S_i \not\subseteq S_j \text{ and } S_j \not\subseteq S_i,$$

i.e., $S_i$ is not a subset of $S_j$ and $S_j$ is not a subset of $S_i$. How large can $m$ be? The following example shows that $m$ can be as large as 10:

$$\{1, 2\}, \{1, 3\}, \{1, 4\}, \{1, 5\}, \{2, 3\}, \{2, 4\}, \{2, 5\}, \{3, 4\}, \{3, 5\}, \{4, 5\}$$

(Note that these are all 10 subsets of $S$ having size two.) Can there be such a sequence of more than 10 subsets? The following theorem states that the answer is “no”.

Theorem 1.2.1 (Sperner) Let $n \geq 1$ be an integer and let $S$ be a set with $n$ elements. Let $S_1, S_2, \ldots, S_m$ be a sequence of $m$ subsets of $S$, such that for all $i$ and $j$ with $i \neq j$,

$$S_i \not\subseteq S_j \text{ and } S_j \not\subseteq S_i.$$

Then

$$m \leq \binom{n}{\lfloor n/2 \rfloor}.$$
1.3. The Quick-Sort Algorithm

The right-hand side of the last line is a binomial coefficient, which we will define in Section 3.6. Its value is equal to the number of subsets of $S$ having size $\lfloor n/2 \rfloor$. Observe that these subsets satisfy the property in Theorem 1.2.1.

We will prove Theorem 1.2.1 in Section 7.3 using elementary counting techniques and probability theory. Again, this probably sounds surprising to you, because Theorem 1.2.1 does not have anything to do with probability.

1.3 The Quick-Sort Algorithm

You are probably familiar with the QuickSort algorithm. This algorithm sorts any sequence $S$ of $n \geq 0$ pairwise distinct numbers in the following way:

- If $n = 0$ or $n = 1$, then there is nothing to do.

- If $n \geq 2$, then the algorithm picks one of the numbers in $S$, let us call it $p$ (which stands for pivot), scans the sequence $S$, and splits it into three subsequences: one subsequence $S_1$ contains all elements in $S$ that are less than $p$, one subsequence only consists of the element $p$, and the third subsequence $S_2$ contains all elements in $S$ that are larger than $p$; see the figure below.

\[
\begin{array}{c|c|c}
<&p & p & >p \\
S_1 & & S_2 \\
\end{array}
\]

The algorithm then recursively runs QuickSort on the subsequence $S_1$, after which it, again recursively, runs QuickSort on the subsequence $S_2$.

Running QuickSort recursively on the subsequence $S_1$ means that we first check if $S_1$ has size at most one; if this is the case, nothing needs to be done, because $S_1$ is sorted already. If $S_1$ has size at least two, then we choose a pivot $p_1$ in $S_1$, use $p_1$ to split $S_1$ into three subsequences, and recursively run QuickSort on the subsequence of $S_1$ consisting of all elements that are less than $p_1$, and finally run QuickSort on the subsequence of $S_1$ consisting of all elements that are larger than $p_1$. (We will see recursive algorithms in more detail in Chapter 4.)

Algorithm QuickSort correctly sorts any sequence of numbers, no matter how we choose the pivot element. It turns out, however, that the running
time of the algorithm heavily depends on the pivots that are chosen in the recursive calls.

For example, assume that in each (recursive) call to the algorithm, the pivot happens to be the largest element in the sequence. Then, in each call, the subsequence of elements that are larger than the pivot is empty. Let us see what happens in this case:

- We start with a sequence $S$ of size $n$. The first pivot $p$ is the largest element in $S$. Thus, using the notation given above, the subsequence $S_1$ contains $n - 1$ elements (all elements of $S$ except for $p$), whereas the subsequence $S_2$ is empty. Computing these subsequences can be done in $n$ “steps”, after which we are in the following situation:

  \[
  \begin{array}{c}
  < p \\
  n - 1 \text{ elements}
  \end{array}
  \begin{array}{c}
  p
  \end{array}
  \]

- We now run \texttt{QuickSort} on a sequence of $n - 1$ elements. Again, the pivot $p_1$ is the largest element. In $n - 1$ steps, we obtain a subsequence of $n - 2$ elements that are less than $p_1$, and an empty subsequence of elements that are larger than $p_1$; see the figure below.

  \[
  \begin{array}{c}
  < p_1 \\
  n - 2 \text{ elements}
  \end{array}
  \begin{array}{c}
  p_1
  \end{array}
  \]

- Next we run \texttt{QuickSort} on a sequence of $n - 2$ elements. As before, the pivot $p_2$ is the largest element. In $n - 2$ steps, we obtain a subsequence of $n - 3$ elements that are less than $p_2$, and an empty subsequence of elements that are larger than $p_2$; see the figure below.

  \[
  \begin{array}{c}
  < p_2 \\
  n - 3 \text{ elements}
  \end{array}
  \begin{array}{c}
  p_2
  \end{array}
  \]

You probably see the pattern. The total running time of the algorithm is proportional to

\[n + (n - 1) + (n - 2) + (n - 3) + \cdots + 3 + 2 + 1,\]
which, by Theorem 2.2.10, is equal to
\[
\frac{1}{2} n(n + 1) = \frac{1}{2} n^2 + \frac{1}{2} n,
\]
which, using the Big-Theta notation (see Section 2.3) is \(\Theta(n^2)\), i.e., quadratic in \(n\). It can be shown that this is, in fact, the worst-case running time of QuickSort.

What would be a good choice for the pivot elements? Intuitively, a pivot is good if the sequences \(S_1\) and \(S_2\) have (roughly) the same size. Thus, after the first call, we are in the following situation:

|  \(< p\) |  \(p\) |  \(> p\) |
|  \((n - 1)/2\) |  \(p\) |  \((n - 1)/2\) |

In Section 4.5, we will prove that, if this happens in each recursive call, the running time of the algorithm is only \(O(n \log n)\). Obviously, it is not clear at all how we can guarantee that we always choose a good pivot. It turns out that there is an easy strategy: In each call, choose the pivot *randomly!* That is, among all elements involved in the recursive call, pick one uniformly at random so that each element has the same probability of being chosen. In Section 6.10, we will prove that this leads to an *expected* running time of \(O(n \log n)\).
Chapter 1. Introduction
Chapter 2
Mathematical Preliminaries

2.1 Basic Concepts

Throughout this book, we will assume that you know the following mathematical concepts:

1. A set is a collection of well-defined objects. Examples are (i) the set of all Dutch Olympic Gold Medallists, (ii) the set of all pubs in Ottawa, and (iii) the set of all even natural numbers.

2. The set of natural numbers is \( \mathbb{N} = \{0, 1, 2, 3, \ldots\} \).

3. The set of integers is \( \mathbb{Z} = \{\ldots, -3, -2, -1, 0, 1, 2, 3, \ldots\} \).

4. The set of rational numbers is \( \mathbb{Q} = \{m/n : m \in \mathbb{Z}, n \in \mathbb{Z}, n \neq 0\} \).

5. The set of real numbers is denoted by \( \mathbb{R} \).

6. The empty set is the set that does not contain any element. This set is denoted by \( \emptyset \).

7. If \( A \) and \( B \) are sets, then \( A \) is a subset of \( B \), written as \( A \subseteq B \), if every element of \( A \) is also an element of \( B \). For example, the set of even natural numbers is a subset of the set of all natural numbers. Every set \( A \) is a subset of itself, i.e., \( A \subseteq A \). The empty set is a subset of every set \( A \), i.e., \( \emptyset \subseteq A \). We say that \( A \) is a proper subset of \( B \), written as \( A \subset B \), if \( A \subseteq B \) and \( A \neq B \).
8. If $B$ is a set, then the power set $\mathcal{P}(B)$ of $B$ is defined to be the set of all subsets of $B$:

$$\mathcal{P}(B) = \{A : A \subseteq B\}.$$ 

Observe that $\emptyset \in \mathcal{P}(B)$ and $B \in \mathcal{P}(B)$.

9. If $A$ and $B$ are two sets, then

(a) their union is defined as

$$A \cup B = \{x : x \in A \text{ or } x \in B\},$$

(b) their intersection is defined as

$$A \cap B = \{x : x \in A \text{ and } x \in B\},$$

(c) their difference is defined as

$$A \setminus B = \{x : x \in A \text{ and } x \notin B\},$$

(d) the Cartesian product of $A$ and $B$ is defined as

$$A \times B = \{(x,y) : x \in A \text{ and } y \in B\},$$

(e) the complement of $A$ is defined as

$$\overline{A} = \{x : x \notin A\}.$$ 

10. A binary relation on two sets $A$ and $B$ is a subset of $A \times B$.

11. A function $f$ from $A$ to $B$, denoted by $f : A \rightarrow B$, is a binary relation $R$, having the property that for each element $a \in A$, there is exactly one ordered pair in $R$, whose first component is $a$. We will also say that $f(a) = b$, or $f$ maps $a$ to $b$, or the image of $a$ under $f$ is $b$. The set $A$ is called the domain of $f$, and the set

$$\{b \in B : \text{ there is an } a \in A \text{ with } f(a) = b\}$$

is called the range of $f$. 
12. A function \( f : A \to B \) is one-to-one (or injective), if for any two distinct elements \( a \) and \( a' \) in \( A \), we have \( f(a) \neq f(a') \). The function \( f \) is onto (or surjective), if for each element \( b \in B \), there exists an element \( a \in A \), such that \( f(a) = b \); in other words, the range of \( f \) is equal to the set \( B \). A function \( f \) is a bijection, if \( f \) is both injective and surjective.

13. A set \( A \) is countable, if \( A \) is finite or there is a bijection \( f : \mathbb{N} \to A \). The sets \( \mathbb{N} \), \( \mathbb{Z} \), and \( \mathbb{Q} \) are countable, whereas \( \mathbb{R} \) is not.

14. A graph \( G = (V, E) \) is a pair consisting of a set \( V \), whose elements are called vertices, and a set \( E \), where each element of \( E \) is a pair of distinct vertices. The elements of \( E \) are called edges.

15. The Boolean values are 1 and 0, that represent true and false, respectively. The basic Boolean operations include

- (a) negation (or NOT), represented by \( \neg \),
- (b) conjunction (or AND), represented by \( \land \),
- (c) disjunction (or OR), represented by \( \lor \),
- (d) exclusive-or (or XOR), represented by \( \oplus \),
- (e) equivalence, represented by \( \leftrightarrow \) or \( \iff \),
- (f) implication, represented by \( \rightarrow \) or \( \Rightarrow \).

The following table explains the meanings of these operations.

<table>
<thead>
<tr>
<th>NOT</th>
<th>AND</th>
<th>OR</th>
<th>XOR</th>
<th>equivalence</th>
<th>implication</th>
</tr>
</thead>
<tbody>
<tr>
<td>\neg 0 = 1</td>
<td>0 \land 0 = 0</td>
<td>0 \lor 0 = 0</td>
<td>0 \oplus 0 = 0</td>
<td>0 \leftrightarrow 0 = 1</td>
<td>0 \rightarrow 0 = 1</td>
</tr>
<tr>
<td>\neg 1 = 0</td>
<td>0 \land 1 = 0</td>
<td>0 \lor 1 = 1</td>
<td>0 \oplus 1 = 1</td>
<td>0 \leftrightarrow 1 = 0</td>
<td>0 \rightarrow 1 = 1</td>
</tr>
<tr>
<td>1 \land 0 = 0</td>
<td>1 \lor 0 = 1</td>
<td>1 \oplus 0 = 1</td>
<td>1 \leftrightarrow 0 = 0</td>
<td>1 \rightarrow 0 = 0</td>
<td></td>
</tr>
<tr>
<td>1 \land 1 = 1</td>
<td>1 \lor 1 = 1</td>
<td>1 \oplus 1 = 0</td>
<td>1 \leftrightarrow 1 = 1</td>
<td>1 \rightarrow 1 = 1</td>
<td></td>
</tr>
</tbody>
</table>

2.2 Proof Techniques

A proof is a proof. What kind of a proof? It’s a proof. A proof is a proof. And when you have a good proof, it’s because it’s proven.

— Jean Chrétien, Prime Minister of Canada (1993–2003)
In mathematics, a theorem is a statement that is true. A proof is a sequence of mathematical statements that form an argument to show that a theorem is true. The statements in the proof of a theorem include axioms (assumptions about the underlying mathematical structures), hypotheses of the theorem to be proved, and previously proved theorems. The main question is “How do we go about proving theorems?” This question is similar to the question of how to solve a given problem. Of course, the answer is that finding proofs, or solving problems, is not easy; otherwise life would be dull! There is no specified way of coming up with a proof, but there are some generic strategies that could be of help. In this section, we review some of these strategies, that will be sufficient for this course. The best way to get a feeling of how to come up with a proof is by solving a large number of problems. Here are some useful tips. (You may take a look at the book *How to Solve It*, by George Pólya).

1. Read and completely understand the statement of the theorem to be proved. Most often this is the hardest part.

2. Sometimes, theorems contain theorems inside them. For example, “Property $A$ if and only if property $B$”, requires showing two statements:

   (a) If property $A$ is true, then property $B$ is true ($A \Rightarrow B$).

   (b) If property $B$ is true, then property $A$ is true ($B \Rightarrow A$).

Another example is the theorem “Set $A$ equals set $B.” To prove this, we need to prove that $A \subseteq B$ and $B \subseteq A$. That is, we need to show that each element of set $A$ is in set $B$, and each element of set $B$ is in set $A$.

3. Try to work out a few simple cases of the theorem just to get a grip on it (i.e., crack a few simple cases first).

4. Try to write down the proof once you think you have it. This is to ensure the correctness of your proof. Often, mistakes are found at the time of writing.

5. Finding proofs takes time, we do not come prewired to produce proofs. Be patient, think, express and write clearly, and try to be precise as much as possible.
In the next sections, we will go through some of the proof strategies.

### 2.2.1 Direct proofs

As the name suggests, in a direct proof of a theorem, we just approach the theorem directly.

**Theorem 2.2.1** *If $n$ is an odd positive integer, then $n^2$ is odd as well.*

**Proof.** An odd positive integer $n$ can be written as $n = 2k + 1$, for some integer $k \geq 0$. Then

\[
    n^2 = (2k + 1)^2 = 4k^2 + 4k + 1 = 2(2k^2 + 2k) + 1.
\]

Since $2(2k^2 + 2k)$ is even, and “even plus one is odd”, we can conclude that $n^2$ is odd.

For a graph $G = (V,E)$, the *degree* of a vertex $v$, denoted by $\deg(v)$, is defined to be the number of edges that are incident on $v$.

**Theorem 2.2.2** *Let $G = (V,E)$ be a graph. Then the sum of the degrees of all vertices is an even integer, i.e.,

\[
    \sum_{v \in V} \deg(v)
\]

is even.

**Proof.** If you do not see the meaning of this statement, then first try it out for a few graphs. The reason why the statement holds is very simple: Each edge contributes 2 to the summation (because an edge is incident on exactly two distinct vertices).

Actually, the proof above proves the following theorem.

**Theorem 2.2.3** *Let $G = (V,E)$ be a graph. Then the sum of the degrees of all vertices is equal to twice the number of edges, i.e.,

\[
    \sum_{v \in V} \deg(v) = 2|E|.
\]
2.2.2 Constructive proofs

This technique not only shows the existence of a certain object, it actually gives a method of creating it:

**Theorem 2.2.4** There exists an object with property $\mathcal{P}$.

**Proof.** Here is the object: $[\ldots]$

And here is the proof that the object satisfies property $\mathcal{P}$: $[\ldots]$  

A graph is called 3-regular, if each vertex has degree three. We prove the following theorem using a constructive proof.

**Theorem 2.2.5** For every even integer $n \geq 4$, there exists a 3-regular graph with $n$ vertices.

**Proof.** Define

$$V = \{0, 1, 2, \ldots, n-1\},$$

and

$$E = \{(i, i+1) : 0 \leq i \leq n-2\} \cup \{(n-1, 0)\} \cup \{(i, i+n/2) : 0 \leq i \leq n/2-1\}.$$  

Then the graph $G = (V, E)$ is 3-regular.

Convince yourself that this graph is indeed 3-regular. It may help to draw the graph for, say, $n = 8$. ■

2.2.3 Nonconstructive proofs

In a nonconstructive proof, we show that a certain object exists, without actually creating it. Here is an example of such a proof:

**Theorem 2.2.6** There exist irrational numbers $x$ and $y$ such that $x^y$ is rational.

**Proof.** There are two possible cases.

**Case 1:** $\sqrt{2}^{\sqrt{2}} \in \mathbb{Q}$.

In this case, we take $x = y = \sqrt{2}$. In Theorem 2.2.9 below, we will prove that $\sqrt{2}$ is irrational.
Case 2: $\sqrt{2} \sqrt{2} \notin \mathbb{Q}$.

In this case, we take $x = \sqrt{2} \sqrt{2}$ and $y = \sqrt{2}$. Since

$$x^y = \left( \sqrt{2} \sqrt{2} \right)^{\sqrt{2}} = \sqrt{2}^2 = 2,$$

the claim in the theorem follows.

Observe that this proof indeed proves the theorem, but it does not give an example of a pair of irrational numbers $x$ and $y$ such that $x^y$ is rational.

### 2.2.4 Proofs by contradiction

This is how a proof by contradiction looks like:

**Theorem 2.2.7** Statement $S$ is true.

**Proof.** Assume that statement $S$ is false. Then, derive a contradiction (such as $1 + 1 = 3$).

In other words, show that the statement “$\neg S \Rightarrow \text{false}$” is true. This is sufficient, because the contrapositive of the statement “$\neg S \Rightarrow \text{false}$” is the statement “$\text{true} \Rightarrow S$”. The latter logical formula is equivalent to $S$, and that is what we wanted to show.

Below, we give two examples of proofs by contradiction.

**Theorem 2.2.8** Let $n$ be a positive integer. If $n^2$ is even, then $n$ is even.

**Proof.** We will prove the theorem by contradiction. So we assume that $n^2$ is even, but $n$ is odd. Since $n$ is odd, we know from Theorem 2.2.1 that $n^2$ is odd. This is a contradiction, because we assumed that $n^2$ is even.

**Theorem 2.2.9** $\sqrt{2}$ is irrational, i.e., $\sqrt{2}$ cannot be written as a fraction of two integers.

**Proof.** We will prove the theorem by contradiction. So we assume that $\sqrt{2}$ is rational. Then $\sqrt{2}$ can be written as a fraction of two integers $m \geq 1$ and $n \geq 1$, i.e., $\sqrt{2} = m/n$. We may assume that $m$ and $n$ are not both
even, because otherwise, we can get rid of the common factors. By squaring
\( \sqrt{2} = \frac{m}{n} \), we get \( 2n^2 = m^2 \). This implies that \( m^2 \) is even. Then, by
Theorem 2.2.8, \( m \) is even, which means that we can write \( m \) as \( m = 2k \); for
some positive integer \( k \). It follows that \( 2n^2 = m^2 = 4k^2 \), which implies that
\( n^2 = 2k^2 \). Hence, \( n^2 \) is even. Again by Theorem 2.2.8 it follows that \( n \) is
even.

We have shown that \( m \) and \( n \) are both even. But we know that \( m \) and
\( n \) are not both even. Hence, we have a contradiction. Our assumption that
\( \sqrt{2} \) is rational is wrong. Thus, we can conclude that \( \sqrt{2} \) is irrational.

There is a nice discussion of this proof in the book *My Brain is Open: The
Mathematical Journeys of Paul Erdős* by Bruce Schechter.

### 2.2.5 Proofs by induction

This is a very powerful and important technique for proving theorems.

For each positive integer \( n \), let \( P(n) \) be a mathematical statement that
depends on \( n \). Assume we wish to prove that \( P(n) \) is true for all positive
integers \( n \). A proof by induction of such a statement is carried out as follows:

**Basis:** Prove that \( P(1) \) is true.

**Induction step:** Prove that for all \( n \geq 1 \), the following holds: If \( P(n) \) is
true, then \( P(n + 1) \) is also true.

In the induction step, we choose an arbitrary integer \( n \geq 1 \) and assume
that \( P(n) \) is true; this is called the *induction hypothesis*. Then we prove that
\( P(n + 1) \) is also true.

**Theorem 2.2.10** For all positive integers \( n \), we have

\[
1 + 2 + 3 + \cdots + n = \frac{n(n + 1)}{2}.
\]

**Proof.** We start with the basis of the induction. If \( n = 1 \), then the left-hand
side is equal to 1, and so is the right-hand side. So the theorem is true for
\( n = 1 \).

For the induction step, let \( n \geq 1 \) and assume that the theorem is true for
\( n \), i.e., assume that

\[
1 + 2 + 3 + \cdots + n = \frac{n(n + 1)}{2}.
\]
2.2. Proof Techniques

We have to prove that the theorem is true for \( n + 1 \), i.e., we have to prove that

\[
1 + 2 + 3 + \cdots + (n + 1) = \frac{(n + 1)(n + 2)}{2}.
\]

Here is the proof:

\[
1 + 2 + 3 + \cdots + (n + 1) = \underbrace{1 + 2 + 3 + \cdots + n + (n + 1)}_{\text{n terms}}
\]

\[
= \frac{n(n + 1)}{2} + (n + 1)
\]

\[
= \frac{(n + 1)(n + 2)}{2}.
\]

By the way, here is an alternative proof of the theorem above: Let \( S = 1 + 2 + 3 + \cdots + n \). Then,

\[
\frac{S}{2} = \frac{1}{n} + \frac{2}{n+1} + \frac{3}{n+2} + \cdots + \frac{2}{n+1} + \frac{3}{n+2} + \cdots + \frac{n}{n+1}
\]

Since there are \( n \) terms on the right-hand side, we have \( 2S = n(n + 1) \). This implies that \( S = n(n + 1)/2 \).

**Theorem 2.2.11** For every positive integer \( n \), \( a - b \) is a factor of \( a^n - b^n \).

**Proof.** A direct proof can be given by providing a factorization of \( a^n - b^n \):

\[
a^n - b^n = (a - b)(a^{n-1} + a^{n-2}b + a^{n-3}b^2 + \cdots + ab^{n-2} + b^{n-1}).
\]

We now prove the theorem by induction. For the basis, let \( n = 1 \). The claim in the theorem is “\( a - b \) is a factor of \( a - b \)”, which is obviously true.

Let \( n \geq 1 \) and assume that \( a - b \) is a factor of \( a^n - b^n \). We have to prove that \( a - b \) is a factor of \( a^{n+1} - b^{n+1} \). We have

\[
a^{n+1} - b^{n+1} = a^n(a - b) + (a^n - b^n)b.
\]

The first term on the right-hand side is divisible by \( a - b \). By the induction hypothesis, the second term on the right-hand side is divisible by \( a - b \) as well. Therefore, the entire right-hand side is divisible by \( a - b \). Since the right-hand side is equal to \( a^{n+1} - b^{n+1} \), it follows that \( a - b \) is a factor of \( a^{n+1} - b^{n+1} \).

We now give an alternative proof of Theorem 2.2.3.
Theorem 2.2.12 Let $G = (V, E)$ be a graph with $m$ edges. Then the sum of the degrees of all vertices is equal to twice the number of edges, i.e.,

$$\sum_{v \in V} \deg(v) = 2m.$$ 

Proof. The proof is by induction on the number $m$ of edges. For the basis of the induction, assume that $m = 0$. Then the graph $G$ does not contain any edges and, therefore, $\sum_{v \in V} \deg(v) = 0$. Thus, the theorem is true if $m = 0$.

Let $m \geq 0$ and assume that the theorem is true for every graph with $m$ edges. Let $G$ be an arbitrary graph with $m + 1$ edges. We have to prove that $\sum_{v \in V} \deg(v) = 2(m + 1)$.

Let $\{a, b\}$ be an arbitrary edge in $G$, and let $G'$ be the graph obtained from $G$ by removing the edge $\{a, b\}$. Since $G'$ has $m$ edges, we know from the induction hypothesis that the sum of the degrees of all vertices in $G'$ is equal to $2m$. Using this, we obtain

$$\sum_{v \in G} \deg(v) = \sum_{v \in G'} \deg(v) + 2 = 2m + 2 = 2(m + 1).$$

2.2.6 More examples of proofs

Recall Theorem 2.2.5, which states that for every even integer $n \geq 4$, there exists a 3-regular graph with $n$ vertices. The following theorem explains why we stated this theorem for even values of $n$.

Theorem 2.2.13 Let $n \geq 5$ be an odd integer. There is no 3-regular graph with $n$ vertices.

Proof. The proof is by contradiction. So we assume that there exists a graph $G = (V, E)$ with $n$ vertices that is 3-regular. Let $m$ be the number of edges in $G$. Since $\deg(v) = 3$ for every vertex, we have

$$\sum_{v \in V} \deg(v) = 3n.$$ 

On the other hand, by Theorem 2.2.3, we have

$$\sum_{v \in V} \deg(v) = 2m.$$
2.2. Proof Techniques

It follows that $3n = 2m$, which can be rewritten as $m = 3n/2$. Since $m$ is an integer, and since $\gcd(2, 3) = 1$, $n/2$ must be an integer. Hence, $n$ is even, which is a contradiction.

Let $K_n$ be the complete graph on $n$ vertices. This graph has a vertex set of size $n$, and every pair of distinct vertices is joined by an edge.

If $G = (V, E)$ is a graph with $n$ vertices, then the complement $\overline{G}$ of $G$ is the graph with vertex set $V$ that consists of those edges of $K_n$ that are not present in $G$.

**Theorem 2.2.14** Let $n \geq 2$ and let $G$ be a graph on $n$ vertices. Then $G$ is connected or $\overline{G}$ is connected.

**Proof.** We prove the theorem by induction on the number $n$ of vertices. For the basis, assume that $n = 2$. There are two possibilities for the graph $G$:

1. $G$ contains one edge. In this case, $G$ is connected.

2. $G$ does not contain an edge. In this case, the complement $\overline{G}$ contains one edge and, therefore, $\overline{G}$ is connected.

Thus, for $n = 2$, the theorem is true.

Let $n \geq 2$ and assume that the theorem is true for every graph with $n$ vertices. Let $G$ be graph with $n + 1$ vertices. We have to prove that $G$ is connected or $\overline{G}$ is connected. We consider three cases.

**Case 1:** There is a vertex $v$ whose degree in $G$ is equal to $n$.

Since $G$ has $n + 1$ vertices, $v$ is connected by an edge to every other vertex of $G$. Therefore, $G$ is connected.

**Case 2:** There is a vertex $v$ whose degree in $G$ is equal to 0.

In this case, the degree of $v$ in the graph $\overline{G}$ is equal to $n$. Since $\overline{G}$ has $n + 1$ vertices, $v$ is connected by an edge to every other vertex of $\overline{G}$. Therefore, $\overline{G}$ is connected.

**Case 3:** For every vertex $v$, the degree of $v$ in $G$ is in $\{1, 2, \ldots, n - 1\}$.

Let $v$ be an arbitrary vertex of $G$. Let $G'$ be the graph obtained by deleting from $G$ the vertex $v$, together with all edges that are incident on $v$. Since $G'$ has $n$ vertices, we know from the induction hypothesis that $G'$ is connected or $\overline{G'}$ is connected.
Let us first assume that $G'$ is connected. Then the graph $G$ is connected as well, because there is at least one edge in $G$ between $v$ and some vertex of $G'$.

If $G'$ is not connected, then $\overline{G}$ must be connected. Since we are in Case 3, we know that the degree of $v$ in $G$ is in the set $\{1, 2, \ldots, n - 1\}$. It follows that the degree of $v$ in the graph $\overline{G}$ is in this set as well. Hence, there is at least one edge in $\overline{C}$ between $v$ and some vertex in $\overline{G'}$. This implies that $\overline{C}$ is connected.

The previous theorem can be rephrased as follows:

**Theorem 2.2.15** Let $n \geq 2$ and consider the complete graph $K_n$ on $n$ vertices. Color each edge of this graph as either red or blue. Let $R$ be the graph consisting of all the red edges, and let $B$ be the graph consisting of all the blue edges. Then $R$ is connected or $B$ is connected.

If you like surprising proofs of various mathematical results, you should read the book *Proofs from THE BOOK* by Martin Aigner and Günter Ziegler.

### 2.3 Asymptotic Notation

Let $f : \mathbb{N} \to \mathbb{R}$ and $g : \mathbb{N} \to \mathbb{R}$ be functions such that $f(n) > 0$ and $g(n) > 0$ for all $n \in \mathbb{N}$.

- We say that $f(n) = O(g(n))$ if the following is true: There exist constants $c > 0$ and $k > 0$ such that for all $n \geq k$,

  $$f(n) \leq c \cdot g(n).$$

- We say that $f(n) = \Omega(g(n))$ if the following is true: There exist constants $c > 0$ and $k > 0$ such that for all $n \geq k$,

  $$f(n) \geq c \cdot g(n).$$

- We say that $f(n) = \Theta(g(n))$ if $f(n) = O(g(n))$ and $f(n) = \Omega(g(n))$. Thus, there exist constants $c > 0, c' > 0$, and $k > 0$ such that for all $n \geq k$,

  $$c \cdot g(n) \leq f(n) \leq c' \cdot g(n).$$
For example, for all \( n \geq 1 \), we have
\[
13 + 7n - 5n^2 + 8n^3 \leq 13 + 7n + 8n^3 \\
\leq 13n^3 + 7n^3 + 8n^3 \\
= 28n^3.
\]

Thus, by taking \( c = 28 \) and \( k = 1 \), it follows that
\[
13 + 7n - 5n^2 + 8n^3 = O(n^3).
\]
(2.1)

We also have
\[
13 + 7n - 5n^2 + 8n^3 \geq -5n^2 + 8n^3.
\]

Since \( n^3 \geq 5n^2 \) for all \( n \geq 5 \), it follows that, again for all \( n \geq 5 \),
\[
13 + 7n - 5n^2 + 8n^3 \geq -5n^2 + 8n^3 \\
\geq -n^3 + 8n^3 \\
= 7n^3.
\]

Hence, by taking \( c = 7 \) and \( k = 5 \), we have shown that
\[
13 + 7n - 5n^2 + 8n^3 = \Omega(n^3).
\]
(2.2)

It follows from (2.1) and (2.2) that
\[
13 + 7n - 5n^2 + 8n^3 = \Theta(n^3).
\]

2.4 Logarithms

If \( b \) and \( x \) are real numbers with \( b > 1 \) and \( x > 0 \), then \( \log_b x \) denotes the logarithm of \( x \) with base \( b \). Note that
\[
\log_b x = y \text{ if and only if } b^y = x.
\]

If \( b = 2 \), then we write \( \log x \) instead of \( \log_2 x \). We write \( \ln x \) to refer to the natural logarithm of \( x \) with base \( e \).

**Lemma 2.4.1** If \( b > 1 \) and \( x > 0 \), then
\[
b^{\log_b x} = x.
\]
Proof. We have seen above that \( y = \log_b x \) if and only if \( b^y = x \). Thus, if we write \( y = \log_b x \), then \( b^{\log_b x} = b^y = x \). □

For example, if \( x > 0 \), then
\[
2^{\log x} = x.
\]

Lemma 2.4.2 If \( b > 1 \), \( x > 0 \), and \( a \) is a real number, then
\[
\log_b (x^a) = a \log_b x.
\]

Proof. Define \( y = \log_b x \). Then
\[
a \log_b x = ay.
\]
Since \( y = \log_b x \), we have \( b^y = x \) and, thus,
\[
x^a = (b^y)^a = b^{ay}.
\]
Taking logarithms (with base \( b \)) on both sides gives
\[
\log_b (x^a) = \log_b (b^{ay}) = ay = a \log_b x.
\]
□

For example, for \( x > 1 \), we get
\[
2 \log \log x = \log (\log^2 x)
\]
and
\[
2^{2\log \log x} = 2^{\log (\log^2 x)} = \log^2 x.
\]

Lemma 2.4.3 If \( b > 1 \), \( c > 1 \), and \( x > 0 \), then
\[
\log_b x = \frac{\log_c x}{\log_c b}.
\]

Proof. Define \( y = \log_b x \). Then \( b^y = x \), and we get
\[
\log_c x = \log_c (b^y) = y \log_c b = \log_b x \log_c b.
\]
□

For example, if \( x > 0 \), then
\[
\log x = \frac{\ln x}{\ln 2}.
\]
2.5 Exercises

2.1 Prove that $\sqrt{p}$ is irrational for every prime number $p$.

2.2 Let $n$ be a positive integer that is not a perfect square. Prove that $\sqrt{n}$ is irrational.

2.3 Use induction to prove that every positive integer can be written as a product of prime numbers.

2.4 Prove by induction that $n^4 - 4n^2$ is divisible by 3, for all integers $n \geq 1$.

2.5 Prove that
\[ \sum_{i=1}^{n} \frac{1}{i^2} < 2 - \frac{1}{n}, \]
for every integer $n \geq 2$.

2.6 Prove that 9 divides $n^3 + (n + 1)^3 + (n + 2)^3$, for every integer $n \geq 0$.

2.7 The Fermat numbers $F_0, F_1, F_2, \ldots$ are defined by $F_n = 2^{2^n} + 1$ for $n \geq 0$.
   - Prove by induction that
     \[ F_0F_1F_2\cdots F_{n-1} = F_n - 2 \]
     for $n \geq 1$.
   - Prove that for any two distinct integers $n \geq 0$ and $m \geq 0$, the greatest common divisor of $F_n$ and $F_m$ is equal to 1.
   - Conclude that there are infinitely many prime numbers.

2.8 Prove by induction that $n! > 2^{1+n/2}$ for all integers $n \geq 3$. 
Chapter 3
Counting

There are three types of people, those who can count and those who cannot count.

Given a set with 23 elements, how many subsets of size 17 are there? How many solutions are there to the equation

\[ x_1 + x_2 + \cdots + x_{12} = 873, \]

where \( x_1 \geq 0, x_2 \geq 0, \ldots, x_{12} \geq 0 \) are integers? In this chapter, we will introduce some general techniques that can be used to answer questions of these types.

3.1 The Product Rule

How many strings of two characters are there, if the first character must be an uppercase letter and the second character must be a digit? Examples of such strings are A0, K7, and Z9. The answer is obviously \( 26 \cdot 10 = 260 \), because there are 26 choices for the first character and, no matter which letter we choose for being the first character, there are 10 choices for the second character. We can look at this in the following way: Consider the “procedure” of writing down a string of two characters, the first one being an uppercase letter, and the second one being a digit. Then our original question becomes “how many ways are there to perform this procedure?” Observe that the procedure consists of two “tasks”, the first one being writing
down the first character, and the second one being writing down the second
character. Obviously, there are 26 ways to do the first task. Next, observe
that, regardless of how we do the first task, there are 10 ways to do the second
task. The Product Rule states that the total number of ways to perform the
entire procedure is $26 \cdot 10 = 260$.

**Product Rule:** Assume a procedure consists of performing a se-
quence of $m$ tasks in order. Furthermore, assume that for each
$i = 1, 2, \ldots, m$, there are $N_i$ ways to do the $i$-th task, regardless
of how the first $i - 1$ tasks were done. Then, there are $N_1 N_2 \cdots N_m$
ways to do the entire procedure.

In the example above, we have $m = 2$, $N_1 = 26$, and $N_2 = 10$.

### 3.1.1 Counting Bitstrings of Length $n$

Let $n \geq 1$ be an integer. A *bitstring* of length $n$ is a sequence of 0’s and 1’s.
How many bitstrings of length $n$ are there? To apply the Product Rule, we
have to specify the “procedure” and the “tasks”:

- The procedure is “write down a bitstring of length $n$”.
- For $i = 1, 2, \ldots, n$, the $i$-th task is “write down one bit”.

There are two ways to do the $i$-th task, regardless of how we did the first
$i - 1$ tasks. Therefore, we can apply the Product Rule with $N_i = 2$ for
$i = 1, 2, \ldots, n$, and conclude that there are $N_1 N_2 \cdots N_n = 2^n$ ways to do the
entire procedure. As a result, the number of bitstrings of length $n$ is equal
to $2^n$.

**Theorem 3.1.1** For any integer $n \geq 1$, the number of bitstrings of length $n$
is equal to $2^n$.

### 3.1.2 Counting Functions

Let $m \geq 1$ and $n \geq 1$ be integers, let $A$ be a set of size $m$, and let $B$ be a set
of size $n$. How many functions $f : A \rightarrow B$ are there?

Write the set $A$ as $A = \{a_1, a_2, \ldots, a_m\}$. Any function $f : A \rightarrow B$ is
completely specified by the values $f(a_1), f(a_2), \ldots, f(a_m)$, where each such
value can be any element of $B$. Again, we are going to apply the Product Rule. Thus, we have to specify the “procedure” and the “tasks”:

- The procedure is “specify the values $f(a_1), f(a_2), \ldots, f(a_m)$”.
- For $i = 1, 2, \ldots, m$, the $i$-th task is “specify the value $f(a_i)$”.

For each $i$, $f(a_i)$ can be any of the $n$ elements of $B$. As a result, there are $N_i = n$ ways to do the $i$-th task, regardless of how we did the first $i - 1$ tasks. By the Product Rule, there are $N_1N_2 \cdots N_m = n^m$ ways to do the entire procedure and, hence, this many functions $f : A \to B$. We have proved the following result:

**Theorem 3.1.2** Let $m \geq 1$ and $n \geq 1$ be integers, let $A$ be a set of size $m$, and let $B$ be a set of size $n$. The number of functions $f : A \to B$ is equal to $n^m$.

Recall that a function $f : A \to B$ is one-to-one if for any $i$ and $j$ with $i \neq j$, we have $f(a_i) \neq f(a_j)$. How many one-to-one functions $f : A \to B$ are there?

If $m > n$, then there is no such function. (Do you see why?) Assume that $m \leq n$. To determine the number of one-to-one functions, we use the same procedure and tasks as before.

- Since $f(a_1)$ can be any of the $n$ elements of $B$, there are $N_1 = n$ ways to do the first task.

- In the second task, we have to specify the value $f(a_2)$. Since the function $f$ is one-to-one and since we have already specified $f(a_1)$, we can choose $f(a_2)$ to be any of the $n - 1$ elements in the set $B \setminus \{f(a_1)\}$. As a result, there are $N_2 = n - 1$ ways to do the second task. Note that this is true, regardless of how we did the first task.
In general, in the $i$-th task, we have to specify the value $f(a_i)$. Since we have already specified $f(a_1), f(a_2), \ldots, f(a_{i-1})$, we can choose $f(a_i)$ to be any of the $n - i + 1$ elements in the set

$$B \setminus \{f(a_1), f(a_2), \ldots, f(a_{i-1})\}.$$ 

As a result, there are $N_i = n - i + 1$ ways to do the $i$-th task. Note that this is true, regardless of how we did the first $i - 1$ tasks.

By the Product Rule, there are

$$N_1 N_2 \cdots N_m = n(n - 1)(n - 2) \cdots (n - m + 1)$$

ways to do the entire procedure, which is also the number of one-to-one functions $f : A \to B$.

Recall the factorial function

$$k! = \begin{cases} 
1 & \text{if } k = 0, \\
1 \cdot 2 \cdot 3 \cdots k & \text{if } k \geq 1.
\end{cases}$$

We can simplify the product

$$n(n - 1)(n - 2) \cdots (n - m + 1)$$

by observing that it is “almost” a factorial:

$$n(n - 1)(n - 2) \cdots (n - m + 1)$$

$$= \frac{n(n - 1)(n - 2) \cdots (n - m + 1)}{(n - m)(n - m - 1) \cdots 1} \cdot \frac{(n - m)(n - m - 1) \cdots 1}{(n - m)(n - m - 1) \cdots 1}$$

$$= \frac{n!}{(n - m)!}.$$ 

We have proved the following result:

**Theorem 3.1.3** Let $m \geq 1$ and $n \geq 1$ be integers, let $A$ be a set of size $m$, and let $B$ be a set of size $n$.

1. If $m > n$, then there is no one-to-one function $f : A \to B$.

2. If $m \leq n$, then the number of one-to-one functions $f : A \to B$ is equal to

$$\frac{n!}{(n - m)!}.$$
3.1.3 Placing Books on Shelves

Let \( m \geq 1 \) and \( n \geq 1 \) be integers, and consider \( m \) books \( B_1, B_2, \ldots, B_m \) and \( n \) book shelves \( S_1, S_2, \ldots, S_n \). How many ways are there to place the books on the shelves? Placing the books on the shelves means that

- we specify for each book the shelf at which this book is placed, and
- we specify for each shelf the left-to-right order of the books that are placed on that shelf.

Some book shelves may be empty. We assume that each shelf is large enough to fit all books. In the figure below, you see two different placements.

We are again going to use the Product Rule to determine the number of placements.

- The procedure is “place the \( m \) books on the \( n \) shelves”.

- For \( i = 1, 2, \ldots, m \), the \( i \)-th task is “place book \( B_i \) on the shelves”. When placing book \( B_i \), we can place it on the far left or far right of any shelf or between any two of the books \( B_1, \ldots, B_{i-1} \) that have already been placed.

Let us see how many ways there are to do each task.

- Just before we place book \( B_1 \), all shelves are empty. Therefore, there are \( N_1 = n \) ways to do the first task.

- In the second task, we have to place book \( B_2 \). Since \( B_1 \) has already been placed, we have the following possibilities for placing \( B_2 \):
  - We place \( B_2 \) on the far left of any of the \( n \) shelves.
  - We place \( B_2 \) immediately to the right of \( B_1 \).
As a result, there are \( N_2 = n + 1 \) ways to do the second task. Note that this is true, regardless of how we did the first task.

- In general, in the \( i \)-th task, we have to place book \( B_i \). Since \( B_1, B_2, \ldots, B_{i-1} \) have already been placed, we have the following possibilities for placing \( B_i \):
  - We place \( B_i \) on the far left of any of the \( n \) shelves.
  - We place \( B_i \) immediately to the right of one of \( B_1, B_2, \ldots, B_{i-1} \).

As a result, there are \( N_i = n + i - 1 \) ways to do the \( i \)-th task. Note that this is true, regardless of how we did the first \( i - 1 \) tasks.

Thus, by the Product Rule, there are

\[
N_1 N_2 \cdots N_m = n(n + 1)(n + 2) \cdots (n + m - 1)
\]

ways to do the entire procedure, which is also the number of placements of the \( m \) books on the \( n \) shelves. As before, we use factorials to simplify this product:

\[
n(n + 1)(n + 2) \cdots (n + m - 1) = \frac{1 \cdot 2 \cdot 3 \cdots (n - 1)}{1 \cdot 2 \cdot 3 \cdots (n - 1)} \cdot n(n + 1)(n + 2) \cdots (n + m - 1) = \frac{(n + m - 1)!}{(n - 1)!}.
\]

We have proved the following result:

**Theorem 3.1.4** Let \( m \geq 1 \) and \( n \geq 1 \) be integers. The number of ways to place \( m \) books on \( n \) book shelves is equal to

\[
\frac{(n + m - 1)!}{(n - 1)!}.
\]

### 3.2 The Bijection Rule

Let \( n \geq 0 \) be an integer and let \( S \) be a set with \( n \) elements. How many subsets does \( S \) have? If \( n = 0 \), then \( S = \emptyset \) and there is only one subset of \( S \), namely \( S \) itself. Assume from now on that \( n \geq 1 \). As we will see below,
3.2. The Bijection Rule

asking for the number of subsets of $S$ is exactly the same as asking for the number of bitstrings of length $n$.

Let $A$ and $B$ be finite sets. Recall that a function $f : A \to B$ is a bijection if

- $f$ is one-to-one, i.e., if $a \neq a'$ then $f(a) \neq f(a')$, and
- $f$ is onto, i.e., for each $b$ in $B$, there is an $a$ in $A$ such that $f(a) = b$.

This means that

- each element of $A$ corresponds to a unique element of $B$ and
- each element of $B$ corresponds to a unique element of $A$.

It should be clear that this means that $A$ and $B$ contain the same number of elements.

**Bijection Rule:** Let $A$ and $B$ be finite sets. If there exists a bijection $f : A \to B$, then $|A| = |B|$, i.e., $A$ and $B$ have the same size.

Let us see how we can apply this rule to the subset problem. We define the following two sets $A$ and $B$:

- $A = \mathcal{P}(S)$, i.e., the power set of $S$, which is the set of all subsets of $S$:
  \[
  \mathcal{P}(S) = \{T : T \subseteq S\}. 
  \]

- $B$ is the set of all bitstrings of length $n$.

We have seen in Theorem 3.1.1 that the set $B$ has size $2^n$. Therefore, if we can show that there exists a bijection $f : A \to B$, then, according to the Bijection Rule, we have $|A| = |B|$ and, thus, the number of subsets of $S$ is equal to $2^n$.

Write the set $S$ as $S = \{s_1, s_2, \ldots, s_n\}$. We define the function $f : A \to B$ in the following way:

- For any $T \in A$ (i.e., $T \subseteq S$), $f(T)$ is the bitstring $b_1b_2\ldots b_n$, where
  \[
  b_i = \begin{cases} 
  1 & \text{if } s_i \in T, \\
  0 & \text{if } s_i \notin T.
  \end{cases}
  \]
For example, assume that $n = 5$.

- If $T = \{s_1, s_3, s_4\}$, then $f(T) = 10110$.
- If $T = S = \{s_1, s_2, s_3, s_4, s_5\}$, then $f(T) = 11111$.
- If $T = \emptyset$, then $f(T) = 00000$.

It is not difficult to see that each subset of $S$ corresponds to a unique bitstring of length $n$, and each bitstring of length $n$ corresponds to a unique subset of $S$. Therefore, this function $f$ is a bijection between $A$ and $B$.

Thus, we have shown that there exists a bijection $f : A \rightarrow B$. This, together with Theorem 3.1.1 and the Bijection Rule, implies the following result:

**Theorem 3.2.1** For any integer $n \geq 0$, the number of subsets of a set of size $n$ is equal to $2^n$.

You will probably have noticed that we could have proved this result directly using the Product Rule: The procedure “specify a subset of $S = \{s_1, s_2, \ldots, s_n\}$” can be carried out by specifying, for $i = 1, 2, \ldots, n$, whether or not $s_i$ is contained in the subset. For each $i$, there are two choices. As a result, there are $2^n$ ways to do the procedure.

### 3.3 The Complement Rule

Consider strings consisting of 8 characters, each character being a lowercase letter or a digit. Such a string is called a valid password if it contains at least one digit. How many valid passwords are there? One way to answer this question is to first count the valid passwords with exactly one digit, then count the valid passwords with exactly two digits, then count the valid passwords with exactly three digits, etc. As we will see below, it is easier to first count the strings that do not contain any digit.

Recall that the difference $U \setminus A$ of the two sets $U$ and $A$ is defined as $U \setminus A = \{x : x \in U \text{ and } x \notin A\}$.

**Complement Rule:** Let $U$ be a finite set and let $A$ be a subset of $U$. Then

$$|A| = |U| - |U \setminus A|.$$
This rule follows easily from the fact that $|U| = |A| + |U \setminus A|$, which holds because each element in $U$ is either in $A$ or in $U \setminus A$.

To apply the Complement Rule to the password problem, let $U$ be the set of all strings consisting of 8 characters, each character being a lowercase letter or a digit, and let $A$ be the set of all valid passwords, i.e., all strings in $U$ that contain at least one digit. Note that $U \setminus A$ is the set of all strings of 8 characters, each character being a lowercase letter or a digit, that do not contain any digit. In other words, $U \setminus A$ is the set of all strings of 8 characters, where each character is a lowercase letter.

By the Product Rule, the set $U$ has size $36^8$, because each string in $U$ has 8 characters, and there are $26 + 10 = 36$ choices for each character. Similarly, the set $U \setminus A$ has size $26^8$, because there are 26 choices for each of the 8 characters. Then, by the Complement Rule, the number of valid passwords is equal to

$$|A| = |U| - |U \setminus A| = 36^8 - 26^8 = 2,612,282,842,880.$$  

### 3.4 The Sum Rule

If $A$ and $B$ are two finite sets that are disjoint, i.e., $A \cap B = \emptyset$, then it is obvious that the size of $A \cup B$ is equal to the sum of the sizes of $A$ and $B$.

**Sum Rule:** Let $A_1, A_2, \ldots, A_m$ be a sequence of finite and pairwise disjoint sets. Then

$$|A_1 \cup A_2 \cup \cdots \cup A_m| = |A_1| + |A_2| + \cdots + |A_m|.$$  

Note that we already used this rule in Section 3.3 when we argued why the Complement Rule is correct!

To give an example, consider strings consisting of 6, 7, or 8 characters, each character being a lowercase letter or a digit. Such a string is called a *valid password* if it contains at least one digit. Let $A$ be the set of all valid passwords. What is the size of $A$?

For $i = 6, 7, 8$, let $A_i$ be the set of all valid passwords of length $i$. It is obvious that $A = A_6 \cup A_7 \cup A_8$. Since the three sets $A_6$, $A_7$, and $A_8$ are pairwise disjoint, we have, by the Sum Rule,

$$|A| = |A_6| + |A_7| + |A_8|.$$
We have seen in Section 3.3 that $|A_8| = 36^8 - 26^8$. By the same arguments, we have $|A_6| = 36^6 - 26^6$ and $|A_7| = 36^7 - 26^7$. Thus, the number of valid passwords is equal to

$$|A| = (36^6 - 26^6) + (36^7 - 26^7) + (36^8 - 26^8) = 2,684,483,063,360.$$ 

### 3.5 The Principle of Inclusion and Exclusion

The Sum Rule holds only for sets that are pairwise disjoint. Consider two finite sets $A$ and $B$ that are not necessarily disjoint. How can we determine the size of the union $A \cup B$? We can start with the sum $|A| + |B|$, i.e., we include both $A$ and $B$. In the figure below,

- $x$ is in $A$ but not in $B$; it is counted exactly once in $|A| + |B|$,
- $z$ is in $B$ but not in $A$; it is counted exactly once in $|A| + |B|$,
- $y$ is in $A$ and in $B$; it is counted exactly twice in $|A| + |B|$.

Based on this, if we subtract the size of the intersection $A \cap B$, i.e., we exclude $A \cap B$, then we have counted every element of $A \cup B$ exactly once.

**Inclusion-Exclusion:** Let $A$ and $B$ be finite sets. Then

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

To give an example, let us count the bitstrings of length 17 that start with 010 or end with 11. Let $S$ be the set of all such bitstrings. Define $A$ to be the set of all bitstrings of length 17 that start with 010, and define $B$ to be the set of all bitstrings of length 17 that end with 11. Then $S = A \cup B$ and, thus, we have to determine the size of $A \cup B$. 

\[
\begin{array}{c}
\text{Inclusion-Exclusion: Let } A \text{ and } B \text{ be finite sets. Then} \\
|A \cup B| = |A| + |B| - |A \cap B|.
\end{array}
\]
3.5. Inclusion-Exclusion

- The size of $A$ is equal to the number of bitstrings of length 14, because the first three bits of every string in $A$ are fixed. Therefore, by the Product Rule, we have $|A| = 2^{14}$.

- The size of $B$ is equal to the number of bitstrings of length 15, because the last two bits of every string in $B$ are fixed. Therefore, by the Product Rule, we have $|B| = 2^{15}$.

- Each string in $A \cap B$ starts with 010 and ends with 11. Thus, five bits are fixed for every string in $A \cap B$. It follows that the size of $A \cap B$ is equal to the number of bitstrings of length 12. Therefore, by the Product Rule, we have $|A \cap B| = 2^{12}$.

By applying the Inclusion-Exclusion formula, it follows that

$$|S| = |A \cup B| = |A| + |B| - |A \cap B| = 2^{14} + 2^{15} - 2^{12} = 45,056.$$ 

The Inclusion-Exclusion formula can be generalized to more than two sets. You are encouraged to verify, using the figure below, that the following formula is the correct one for three sets.

**Inclusion-Exclusion:** Let $A$, $B$, and $C$ be finite sets. Then

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

To give an example, how many bitstrings of length 17 are there that start with 010, or end with 11, or have 10 at positions 7 and 8? Let $S$ be the set of all such bitstrings. Define $A$ to be the set of all bitstrings of length 17 that
start with 010, define $B$ to be the set of all bitstrings of length 17 that end with 11, and define $C$ to be the set of all bitstrings of length 17 that have 10 at positions 7 and 8. Then $S = A \cup B \cup C$ and, thus, we have to determine the size of $A \cup B \cup C$.

- We have seen before that $|A| = 2^{14}$, $|B| = 2^{15}$, and $|A \cap B| = 2^{12}$.
- We have $|C| = 2^{15}$, because the bits at positions 7 and 8 are fixed for every string in $C$.
- We have $|A \cap C| = 2^{12}$, because 5 bits are fixed for every string in $A \cap C$.
- We have $|B \cap C| = 2^{13}$, because 4 bits are fixed for every string in $B \cap C$.
- We have $|A \cap B \cap C| = 2^{10}$, because 7 bits are fixed for every string in $A \cap B \cap C$.

By applying the Inclusion-Exclusion formula, it follows that

$$|S| = |A \cup B \cup C|$$

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

$$= 2^{14} + 2^{15} + 2^{15} - 2^{12} - 2^{12} - 2^{13} + 2^{10}$$

$$= 66,560.$$

### 3.6 Permutations and Binomial Coefficients

A permutation of a finite set $S$ is an ordered sequence of the elements of $S$, in which each element occurs exactly once. For example, the set $S = \{a, b, c\}$ has exactly six permutations:

$$\text{abc, acb, bac, bca, cab, cba}$$

**Theorem 3.6.1** Let $n \geq 0$ be an integer and let $S$ be a set with $n$ elements. There are exactly $n!$ many permutations of $S$.

**Proof.** If $n = 0$, then $S = \emptyset$ and the only permutation of $S$ is the empty sequence. Since $0! = 1$, the claim holds for $n = 0$. Assume that $n \geq 1$ and denote the elements of $S$ by $s_1, s_2, \ldots, s_n$. Consider the procedure “write
down a permutation of $S$ and, for $i = 1, 2, \ldots, n$, the task “write down the $i$-th element in the permutation”. When we do the $i$-th task, we have already written down $i - 1$ elements of the permutation; we cannot take any of these elements for the $i$-th task. Therefore, there are $n - (i - 1) = n - i + 1$ ways to do the $i$-th task. By the Product Rule, there are 

$$n \cdot (n - 1) \cdot (n - 2) \cdots 2 \cdot 1 = n!$$

ways to do the procedure. This number is equal to the number of permutations of $S$.

Note that we could also have used Theorem 3.1.3 to prove Theorem 3.6.1: A permutation of $S$ is a one-to-one function $f : S \to S$. Therefore, by applying Theorem 3.1.3 with $A = S$, $B = S$ and, thus, $m = n$, we obtain Theorem 3.6.1.

Consider the set $S = \{a, b, c, d, e\}$. How many 3-element subsets does $S$ have? Recall that in a set, the order of the elements does not matter. Here is a list of all 10 subsets of $S$ having size 3:

$$\{a, b, c\}, \{a, b, d\}, \{a, b, e\}, \{a, c, d\}, \{a, c, e\}, \{a, d, e\}, \{b, c, d\}, \{b, c, e\}, \{b, d, e\}, \{c, d, e\}$$

**Definition 3.6.2** Let $n$ and $k$ be integers. We define the binomial coefficient $\binom{n}{k}$ to be the number of $k$-element subsets of an $n$-element set.

The symbol $\binom{n}{k}$ is pronounced as “$n$ choose $k$”.

The example above shows that $\binom{5}{3} = 10$. Since the empty set has exactly one subset of size zero (the empty set itself), we have $\binom{0}{0} = 1$. Note that $\binom{n}{k} = 0$ if $k < 0$ or $k > n$. Below, we derive a formula for the value of $\binom{n}{k}$ if $0 \leq k \leq n$.

Let $S$ be a set with $n$ elements and let $A$ be the set of all ordered sequences consisting of exactly $k$ pairwise distinct elements of $S$. We are going to count the number of elements of $A$ in two different ways.

The first way is by using the Product Rule. This gives

$$|A| = n(n - 1)(n - 2) \cdots (n - k + 1) = \frac{n!}{(n - k)!}. \quad (3.1)$$

In the second way, we do the following:
• Write down all $\binom{n}{k}$ subsets of $S$ having size $k$.

• For each of these subsets, write down a list of all $k!$ permutations of this subset.

If we put all these lists together, then we obtain a big list in which each ordered sequence of $k$ pairwise distinct elements of $S$ appears exactly once. In other words, the big list contains each element of $A$ exactly once. Since the big list has size $\binom{n}{k}k!$, it follows that

$$|A| = \binom{n}{k}k!$$  \hspace{1cm} (3.2)

By equating the right-hand sides of (3.1) and (3.2), we obtain the following result:

**Theorem 3.6.3** Let $n$ and $k$ be integers with $0 \leq k \leq n$. Then

$$\binom{n}{k} = \frac{n!}{k!(n-k)!}.$$ 

For example,

$$\binom{5}{3} = \frac{5!}{3!(5-3)!} = \frac{5!}{3!2!} = \frac{1 \cdot 2 \cdot 3 \cdot 4 \cdot 5}{1 \cdot 2 \cdot 3 \cdot 1 \cdot 2} = 10$$

and

$$\binom{0}{0} = \frac{0!}{0!0!} = \frac{1}{1 \cdot 1} = 1;$$

recall that we defined $0!$ to be equal to 1.

### 3.6.1 Some Examples

Consider a standard deck of 52 cards. How many hands of 5 cards are there? Any such hand is a 5-element subset of the set of 52 cards and, therefore, the number of hands of 5 cards is equal to

$$\binom{52}{5} = \frac{52!}{5!47!} = \frac{52 \cdot 51 \cdot 50 \cdot 49 \cdot 48}{5 \cdot 4 \cdot 3 \cdot 2 \cdot 1} = 2,598,960.$$ 

Let $n$ and $k$ be integers with $0 \leq k \leq n$. How many bitstrings of length $n$ have exactly $k$ many 1s? We can answer this question using the Product Rule:
3.6. Permutations and Binomial Coefficients

- The procedure is “write down a bitstring of length \( n \) having exactly \( k \) many 1s”.
- Task 1: Choose \( k \) positions out of 1, 2, \ldots, \( n \).
- Task 2: Write a 1 in each of the \( k \) chosen positions.
- Task 3: Write a 0 in each of the \( n - k \) remaining positions.

There are \( \binom{n}{k} \) ways to do the first task, there is only one way to do the second task, and there is only one way to do the third task. Thus, by the Product Rule, the number of ways to do the procedure and, therefore, the number of bitstrings of length \( n \) having exactly \( k \) many 1s, is equal to

\[
\binom{n}{k} \cdot 1 \cdot 1 = \binom{n}{k}.
\]

We can also use the Bijection Rule, by observing that there is a bijection between

- the set of all bitstrings of length \( n \) having exactly \( k \) many 1s, and
- the set of all \( k \)-element subsets of an \( n \)-element set.

Since the latter set has size \( \binom{n}{k} \), the former set has size \( \binom{n}{k} \) as well.

**Theorem 3.6.4** Let \( n \) and \( k \) be integers with \( 0 \leq k \leq n \). The number of bitstrings of length \( n \) having exactly \( k \) many 1s is equal to \( \binom{n}{k} \).

### 3.6.2 Newton’s Binomial Theorem

Consider

\[
(x + y)^5 = (x + y)(x + y)(x + y)(x + y)(x + y).
\]

If we expand the expression on the right-hand side, we get terms

\[
x^5, x^4y, x^3y^2, x^2y^3, xy^4, y^5,
\]

each with some coefficient. What is the coefficient of \( x^2y^3 \)? We obtain a term \( x^2y^3 \), by

- choosing 3 of the 5 terms \( x + y \),
• taking \( y \) in each of the 3 chosen terms \( x + y \), and
• taking \( x \) in each of the other 2 terms \( x + y \).

Since there are \( \binom{3}{3} \) ways to do this, the coefficient of \( x^2 y^3 \) is equal to \( \binom{5}{3} = 10 \).

**Theorem 3.6.5 (Newton’s Binomial Theorem)** For any integer \( n \geq 0 \), we have

\[
(x + y)^n = \sum_{k=0}^{n} \binom{n}{k} x^{n-k} y^k.
\]

**Proof.** The expression \((x + y)^n\) is a product of \( n \) terms \( x + y \). By expanding this product, we get a term \( x^{n-k} y^k \) for each \( k = 0, 1, \ldots, n \). We get a term \( x^{n-k} y^k \) by

• choosing \( k \) of the \( n \) terms \( x + y \),
• taking \( y \) in each of the \( k \) chosen terms \( x + y \), and
• taking \( x \) in each of the other \( n - k \) terms \( x + y \).

Since there are \( \binom{n}{k} \) ways to do this, the coefficient of \( x^{n-k} y^k \) is equal to \( \binom{n}{k} \).

For example, we have

\[
(x + y)^3 = \binom{3}{0} x^3 + \binom{3}{1} x^2 y + \binom{3}{2} xy^2 + \binom{3}{3} y^3
= x^3 + 3x^2 y + 3xy^2 + y^3.
\]

To determine the coefficient of \( x^{12} y^{13} \) in \((x + y)^{25}\), we take \( n = 25 \) and \( k = 13 \) in Newton’s Binomial Theorem, and get \( \binom{25}{13} \).

What is the coefficient of \( x^{12} y^{13} \) in \((2x - 5y)^{25}\)? Observe that

\[
(2x - 5y)^{25} = ((2x) + (-5y))^{25}.
\]

By replacing \( x \) by \( 2x \), and \( y \) by \( -5y \) in Newton’s Binomial Theorem, we get

\[
(2x - 5y)^{25} = \sum_{k=0}^{25} \binom{25}{k} (2x)^{25-k} (-5y)^k.
\]
By taking \( k = 13 \), we obtain the coefficient of \( x^{12}y^{13} \):

\[
\binom{25}{13} \cdot 2^{25-13} \cdot (-5)^{13} = -\binom{25}{13} \cdot 2^{12} \cdot 5^{13}.
\]

Newton’s Binomial Theorem leads to identities for summations involving binomial coefficients:

**Theorem 3.6.6** For any integer \( n \geq 0 \), we have

\[
\sum_{k=0}^{n} \binom{n}{k} = 2^n.
\]

**Proof.** Take \( x = y = 1 \) in Newton’s Binomial Theorem. □

In Section 3.7, we will see a proof of Theorem 3.6.6 that does not use Newton’s Binomial Theorem.

**Theorem 3.6.7** For any integer \( n \geq 1 \), we have

\[
\sum_{k=0}^{n} (-1)^k \binom{n}{k} = 0.
\]

**Proof.** Take \( x = 1 \) and \( y = -1 \) in Newton’s Binomial Theorem. □

### 3.7 Combinatorial Proofs

In a combinatorial proof, we show the validity of an identity, such as the one in Theorem 3.6.6, by interpreting it as a counting problem. In this section, we will give several examples.

**Theorem 3.7.1** For integers \( n \) and \( k \) with \( 0 \leq k \leq n \), we have

\[
\binom{n}{k} = \binom{n}{n-k}.
\]

**Proof.** The claim can be proved using Theorem 3.6.3. To obtain a combinatorial proof, let \( S \) be a set with \( n \) elements. Recall that
• \( \binom{n}{k} \) is the number of ways to choose \( k \) elements from the set \( S \), which is the same as
  
  • the number of ways to not choose \( n-k \) elements from the set \( S \).

The latter number is equal to \( \binom{n}{n-k} \).

We can also prove the claim using Theorem 3.6.4:

• The number of bitstrings of length \( n \) with exactly \( k \) many 1s is equal to \( \binom{n}{k} \).

• The number of bitstrings of length \( n \) with exactly \( n-k \) many 0s is equal to \( \binom{n}{n-k} \).

Since these two quantities are equal, the theorem follows.

**Theorem 3.7.2 (Pascal’s Identity)** For integers \( n \) and \( k \) with \( n \geq k \geq 1 \), we have

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

**Proof.** As in the previous theorem, the claim can be proved using Theorem 3.6.3. To obtain a combinatorial proof, let \( S \) be a set with \( n+1 \) elements. We are going to count the number of \( k \)-element subsets of \( S \) in two different ways.

First, by definition, the number of \( k \)-element subsets of \( S \) is equal to

\[
\binom{n+1}{k}.
\]

(3.3)

For the second way, we choose an element \( x \) in \( S \) and consider the set \( T = S \setminus \{x\} \), i.e., the set obtained by removing \( x \) from \( S \). Any \( k \)-element subset of \( S \) is of one of the following two types:

• The \( k \)-element subset of \( S \) does not contain \( x \).
  
  – Any such subset is a \( k \)-element subset of \( T \). Since \( T \) has size \( n \), there are \( \binom{n}{k} \) many \( k \)-element subsets of \( S \) that do not contain \( x \).

• The \( k \)-element subset of \( S \) contains \( x \).
3.7. Combinatorial Proofs

- If \( A \) is any such subset, then \( B = A \setminus \{x\} \) is a \((k - 1)\)-element subset of \( T \).
- Conversely, for any \((k - 1)\)-element subset \( B \) of \( T \), the set \( A = B \cup \{x\} \) is a \( k \)-element subset of \( S \) that contains \( x \).
- It follows that the number of \( k \)-element subsets of \( S \) containing \( x \) is equal to the number of \((k - 1)\)-element subsets of \( T \). The latter number is equal to \( \binom{n}{k-1} \).

Thus, the second way of counting shows that the number of \( k \)-element subsets of \( S \) is equal to
\[
\binom{n}{k} + \binom{n}{k-1}.
\]

Since the expressions in (3.3) and (3.4) count the same objects, they must be equal. Therefore, the proof is complete.

Theorem 3.7.3 For any integer \( n \geq 0 \), we have
\[
\sum_{k=0}^{n} \binom{n}{k} = 2^n.
\]

Proof. We have seen in Theorem 3.6.6 that this identity follows from Newton’s Binomial Theorem. Below, we give a combinatorial proof.

Consider a set \( S \) with \( n \) elements. According to Theorem 3.2.1, this set has \( 2^n \) many subsets. A different way to count the subsets of \( S \) is by dividing them into groups according to their sizes. For each \( k \) with \( 0 \leq k \leq n \), consider all \( k \)-element subsets of \( S \). The number of such subsets is equal to \( \binom{n}{k} \). If we take the sum of all these binomial coefficients, then we have counted each subset of \( S \) exactly once. Thus,
\[
\sum_{k=0}^{n} \binom{n}{k}
\]
is equal to the total number of subsets of \( S \).

Theorem 3.7.4 (Vandermonde’s Identity) For integers \( m \geq 0 \), \( n \geq 0 \), and \( r \geq 0 \) with \( r \leq m \) and \( r \leq n \), we have
\[
\sum_{k=0}^{r} \binom{m}{k} \binom{n}{r-k} = \binom{m+n}{r}.
\]
Proof. Consider a set $S$ with $m + n$ elements. We are going to count the $r$-element subsets of $S$ in two different ways.

First, by using the definition of binomial coefficients, the number of $r$-element subsets of $S$ is equal to $\binom{m+n}{r}$.

For the second way, we partition the set $S$ into two subsets $A$ and $B$, where $A$ has size $m$ and $B$ has size $n$. Observe that any $r$-element subset of $S$ contains

- some elements of $A$, say $k$ many, and
- $r - k$ elements of $B$.

Let $k$ be any integer with $0 \leq k \leq r$, and let $N_k$ be the number of $r$-element subsets of $S$ that contain exactly $k$ elements of $A$ (and, thus, $r - k$ elements of $B$). Then,

$$\sum_{k=0}^{r} N_k = \binom{m+n}{r}.$$

To determine $N_k$, we use the Product Rule: We obtain any subset that is counted in $N_k$, by

- choosing $k$ elements in $A$ (there are $\binom{m}{k}$ ways to do this) and
- choosing $r - k$ elements in $B$ (there are $\binom{n}{r-k}$ ways to do this).

It follows that

$$N_k = \binom{m}{k} \binom{n}{r-k}.$$

\[\square\]

**Corollary 3.7.5** For any integer $n \geq 0$, we have

$$\sum_{k=0}^{n} \binom{n}{k}^2 = \binom{2n}{n}.$$

**Proof.** By taking $m = n = r$ in Vandermonde’s Identity, we get

$$\sum_{k=0}^{n} \binom{n}{k} \binom{n}{n-k} = \binom{2n}{n}.$$
Using Theorem 3.7.1, we get
\[
\binom{n}{k} \binom{n}{n-k} = \binom{n}{k} \binom{n-k}{k} = \binom{n}{k}^2.
\]

3.8 Pascal’s Triangle

The computational method at the heart of Pascal’s work was actually discovered by a Chinese mathematician named Jia Xian around 1050, published by another Chinese mathematician, Zhu Shijie, in 1303, discussed in a work by Cardano in 1570, and plugged into the greater whole of probability theory by Pascal, who ended up getting most of the credit.


We have seen that

• \( \binom{n}{0} = 1 \) for all integers \( n \geq 0 \),

• \( \binom{n}{n} = 1 \) for all integers \( n \geq 0 \),

• \( \binom{n}{k} = \binom{n-1}{k-1} + \binom{n-1}{k} \) for all integers \( n \geq 2 \) and \( k \) with \( 1 \leq k \leq n-1 \); see Theorem 3.7.2.

These relations lead to an algorithm for generating binomial coefficients:

**Algorithm** GENERATEBINOMCOEFF:

\[
BCoeff(0, 0) = 1;
\]

for \( n = 1, 2, 3, \ldots \)

do \( BCoeff(n, 0) = 1; \)

for \( k = 1 \) to \( n - 1 \)

do \( BCoeff(n, k) = BCoeff(n-1, k-1) + BCoeff(n-1, k) \)
endfor;

\( BCoeff(n, n) = 1 \)
endfor
The values $BC_{\text{Coeff}}(n,k)$ that are computed by this (non-terminating) algorithm satisfy

$$BC_{\text{Coeff}}(n,k) = \binom{n}{k} \text{ for } 0 \leq k \leq n.$$ 

The triangle obtained by arranging these binomial coefficients, with the $n$-th row containing all values $\binom{n}{0}, \binom{n}{1}, \ldots, \binom{n}{n}$, is called Pascal’s Triangle. The figure below shows rows 0, 1, \ldots, 6:

\[
\begin{array}{cccccccc}
\binom{0}{0} \\
\binom{1}{0} & \binom{1}{1} \\
\binom{2}{0} & \binom{2}{1} & \binom{2}{2} \\
\binom{3}{0} & \binom{3}{1} & \binom{3}{2} & \binom{3}{3} \\
\binom{4}{0} & \binom{4}{1} & \binom{4}{2} & \binom{4}{3} & \binom{4}{4} \\
\binom{5}{0} & \binom{5}{1} & \binom{5}{2} & \binom{5}{3} & \binom{5}{4} & \binom{5}{5} \\
\binom{6}{0} & \binom{6}{1} & \binom{6}{2} & \binom{6}{3} & \binom{6}{4} & \binom{6}{5} & \binom{6}{6} \\
\end{array}
\]

We obtain the values for the binomial coefficients by using the following rules:

- Each value along the boundary is equal to 1.
- Each value in the interior is equal to the sum of the two values above it.

In Figure 3.1, you see rows 0, 1, \ldots, 12.
Below, we state some of our earlier results using Pascal’s Triangle.

- The values in the $n$-th row are equal to the coefficients in Newton’s Binomial Theorem. For example, the coefficients in the expansion of $(x + y)^5$ are given in the 5-th row:

$$\text{\(x + y\)^5} = x^5 + 5x^4y + 10x^3y^2 + 10x^2y^3 + 5xy^4 + y^5.$$
3.8. Pascal’s Triangle

Figure 3.1: Rows 0, 1, . . . , 12 of Pascal’s Triangle.
• Theorem 3.6.6 states that the sum of all values in the \(n\)-th row is equal to \(2^n\).

• Theorem 3.7.1 states that reading the \(n\)-th row from left to right gives the same sequence as reading this row from right to left.

• Corollary 3.7.5 states that the sum of the squares of all values in the \(n\)-th row is equal to the middle element in the \(2n\)-th row.

3.9 More Counting Problems

3.9.1 Reordering the Letters of a Word

How many different strings can be made by reordering the letters of the 7-letter word

SUCCESS.

It should be clear that the answer is not \(7!\): If we swap, for example, the two occurrences of C, then we obtain the same string.

The correct answer can be obtained by applying the Product Rule. We start by counting the frequencies of each letter:

• The letter S occurs 3 times.
• The letter C occurs 2 times.
• The letter U occurs 1 time.
• The letter E occurs 1 time.

To apply the Product Rule, we have to specify the procedure and the tasks:

• The procedure is “write down the letters occurring in the word SUCCESS”.

• The first task is “choose 3 positions out of 7, and write the letter S in each chosen position”.

• The second task is “choose 2 positions out of the remaining 4, and write the letter C in each chosen position”.
• The third task is “choose 1 position out of the remaining 2, and write the letter $U$ in the chosen position”.

• The fourth task is “choose 1 position out of the remaining 1, and write the letter $E$ in the chosen position”.

Since there are \( \binom{7}{3} \) ways to do the first task, \( \binom{4}{2} \) ways to do the second task, \( \binom{2}{1} \) ways to do the third task, and \( \binom{1}{1} \) way to do the fourth task, it follows that the total number of different strings that can be made by reordering the letters of the word SUCCESS is equal to

\[
\binom{7}{3} \binom{4}{2} \binom{2}{1} \binom{1}{1} = 420.
\]

In the four tasks above, we first chose the positions for the letter $S$, then the positions for the letter $C$, then the position for the letter $U$, and finally the position for the letter $E$. If we change the order, then we obtain the same answer. For example, if we choose the positions for the letters in the order $C, E, U, S$, then we obtain

\[
\binom{7}{2} \binom{5}{1} \binom{4}{3} \binom{3}{3},
\]

which is indeed equal to 420.

### 3.9.2 Counting Solutions of Linear Equations

Consider the equation

\[x_1 + x_2 + x_3 = 11.\]

We are interested in the number of solutions \((x_1, x_2, x_3)\), where \(x_1 \geq 0, x_2 \geq 0, x_3 \geq 0\) are integers. Examples of solutions are

\[(2, 3, 6), (3, 2, 6), (0, 11, 0), (2, 0, 9).\]

Observe that we consider \((2, 3, 6)\) and \((3, 2, 6)\) to be different solutions.

We are going to use the Bijection Rule to determine the number of solutions. For this, we define \(A\) to be the set of all solutions, i.e.,

\[A = \{(x_1, x_2, x_3) : x_1 \geq 0, x_2 \geq 0, x_3 \geq 0 \text{ are integers, } x_1 + x_2 + x_3 = 11\}.\]
To apply the Bijection Rule, we need a set $B$ and a bijection $f : A \to B$, such that it is easy to determine the size of $B$. This set $B$ should be chosen such that its elements “encode” the elements of $A$ in a unique way. Consider the following set $B$:

- $B$ is the set of all binary strings of length 13 that contain exactly 2 many 1s (and, thus, exactly 11 many 0s).

The function $f : A \to B$ is defined as follows: If $(x_1, x_2, x_3)$ is an element of $A$, then $f(x_1, x_2, x_3)$ is the binary string

- that starts with $x_1$ many 0s,
- is followed by one 1,
- is followed by $x_2$ many 0s,
- is followed by one 1,
- and ends with $x_3$ many 0s.

For example, we have

\[
f(2, 3, 6) = 0010001000000,
\]

\[
f(3, 2, 6) = 0001001000000,
\]

\[
f(0, 11, 0) = 1000000000001,
\]

and

\[
f(2, 0, 9) = 0011000000000.
\]

To show that this function $f$ maps elements of $A$ to elements of $B$, we have to verify that the string $f(x_1, x_2, x_3)$ belongs to the set $B$. This follows from the following observations:

- The string $f(x_1, x_2, x_3)$ contains exactly 2 many 1s.
- The number of 0s in the string $f(x_1, x_2, x_3)$ is equal to $x_1 + x_2 + x_3$, which is equal to 11, because $(x_1, x_2, x_3)$ belongs to the set $A$. 

It should be clear that this function $f$ is one-to-one: If we take two different elements $(x_1, x_2, x_3)$ in $A$, then $f$ gives us two different bitstrings $f(x_1, x_2, x_3)$.

To prove that $f$ is onto, we have to show that for every bitstring $b$ in the set $B$, there is an element $(x_1, x_2, x_3)$ in $A$ such that $f(x_1, x_2, x_3) = b$. This element of $A$ is obtained by taking

- $x_1$ to be the number of 0s to the left of the first 1,
- $x_2$ to be the number of 0s between the two 1s, and
- $x_3$ to be the number of 0s to the right of the second 1.

For example, if $b = 000011000000$, then $x_1 = 4$, $x_2 = 0$, and $x_3 = 7$. Note that, since $b$ has length 13 and contains exactly 2 many 1s, we have $x_1 + x_2 + x_3 = 11$ and, therefore, $(x_1, x_2, x_3) \in A$.

Thus, we have shown that $f : A \rightarrow B$ is indeed a bijection. We know from Theorem 3.6.4 that $B$ has size $\binom{13}{2}$. Therefore, it follows from the Bijection Rule that

$$|A| = |B| = \binom{13}{2} = 78.$$

The following theorem states this result for general linear equations. You are encouraged to come up with the proof.

**Theorem 3.9.1** Let $k \geq 1$ and $n \geq 0$ be integers. The number of solutions to the equation

$$x_1 + x_2 + \cdots + x_k = n,$$

where $x_1 \geq 0$, $x_2 \geq 0$, $\ldots$, $x_k \geq 0$ are integers, is equal to

$$\binom{n + k - 1}{k - 1}.$$

Let us now consider inequalities instead of equations. For example, consider the inequality

$$x_1 + x_2 + x_3 \leq 11.$$

Again, we are interested in the number of solutions $(x_1, x_2, x_3)$, where $x_1 \geq 0$, $x_2 \geq 0$, $x_3 \geq 0$ are integers. This inequality contains the same solutions as before, but it has additional solutions such as

$$(2, 3, 5), (3, 2, 5), (0, 1, 0), (0, 0, 0).$$
As before, we are going to apply the Bijection Rule. We define

\[ A = \{(x_1, x_2, x_3) : x_1 \geq 0, x_2 \geq 0, x_3 \geq 0 \text{ are integers, } x_1 + x_2 + x_3 \leq 11\} \]

and \( B \) to be the set of all binary strings of length 14 that contain exactly 3 many 1s (and, thus, exactly 11 many 0s).

The function \( f : A \rightarrow B \) is defined as follows: If \((x_1, x_2, x_3)\) is an element of \( A \), then \( f(x_1, x_2, x_3)\) is the binary string

- that starts with \( x_1 \) many 0s,
- is followed by one 1,
- is followed by \( x_2 \) many 0s,
- is followed by one 1,
- is followed by \( x_3 \) many 0s,
- is followed by one 1,
- and ends with \( 14 - (x_1 + x_2 + x_3 + 3) \) many 0s.

For example, we have

\[ f(2, 3, 6) = 00100010000001, \]
\[ f(2, 3, 5) = 00100010000010, \]
\[ f(0, 1, 0) = 10110000000000, \]

and

\[ f(0, 0, 0) = 11100000000000. \]

As before, it can be verified that the string \( f(x_1, x_2, x_3) \) belongs to the set \( B \) and the function \( f \) is a bijection. It then follows from the Bijection Rule that

\[ |A| = |B| = \binom{14}{3} = 364. \]

The next theorem gives the answer for the general case. As before, you are encouraged to give a proof.
3.10. The Pigeonhole Principle

**Theorem 3.9.2** Let \( k \geq 1 \) and \( n \geq 0 \) be integers. The number of solutions to the inequality

\[ x_1 + x_2 + \cdots + x_k \leq n, \]

where \( x_1 \geq 0, x_2 \geq 0, \ldots, x_k \geq 0 \) are integers, is equal to

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

3.10 The Pigeonhole Principle

In any group of 366 people, there must be two people having the same birthday: Since there are 365 days in a year (ignoring leap years), it is not possible that the birthdays of 366 people are all distinct.

**Pigeonhole Principle:** Let \( k \geq 1 \) be an integer. If \( k+1 \) or more objects are placed into \( k \) boxes, then there is at least one box containing two or more objects.

Equivalently, if \( A \) and \( B \) are two finite sets such that \(|A| > |B|\), then there is no one-to-one function from \( A \) to \( B \).

3.10.1 India Pale Ale

Simon Pratt\(^1\) loves to drink India Pale Ale (IPA\(^2\)). During each day of the month of April (which has 30 days), Simon drinks at least one bottle of IPA. Also, during this entire month, he drinks exactly 45 bottles of IPA. The claim is that there must be a sequence of consecutive days in April, during which Simon drinks exactly 14 bottles of IPA.

To prove this, let \( b_i \) be the number of bottles that Simon drinks on April \( i \), for \( i = 1, 2, \ldots, 30 \). We are given that each \( b_i \) is a positive integer (i.e., \( b_i \geq 1 \)) and

\[ b_1 + b_2 + \cdots + b_{30} = 45. \]

Define, for \( i = 1, 2, \ldots, 30 \),

\[ a_i = b_1 + b_2 + \cdots + b_i. \]

---

\(^1\)President of the Carleton Computer Science Society (2013–2014).

\(^2\)Lindsay Bangs, President of the CCSS (2014–2015), prefers wheat beer, whereas Connor Hillen, President of the CCSS (2015–2016), prefers Black IPA.
i.e., $a_i$ is the total number of bottles of IPA that Simon drinks during the first $i$ days of April. Consider the sequence of 60 numbers

$$a_1, a_2, \ldots, a_{30}, a_1 + 14, a_2 + 14, \ldots, a_{30} + 14.$$  

Each number in this sequence is an integer that belongs to the set

$$\{1, 2, \ldots, 59\}.$$  

Therefore, by the Pigeonhole Principle, these 60 numbers cannot all be distinct. Observe that there are no duplicates in the sequence $a_1, a_2, \ldots, a_{30}$, because all $b_i$ are at least one. Similarly, there are no duplicates in the sequence $a_1 + 14, a_2 + 14, \ldots, a_{30} + 14$. It follows that there are two indices $i$ and $j$ such that

$$a_i = a_j + 14.$$  

Observe that $j$ must be less than $i$ and

$$14 = a_i - a_j = b_{j+1} + b_{j+2} + \cdots + b_i.$$  

Thus, in the period from April $j + 1$ until April $i$, Simon drinks exactly 14 bottles of IPA.

### 3.10.2 Sequences Containing Divisible Numbers

Let $A = \{1, 2, \ldots, 2n\}$ and consider the sequence $n + 1, n + 2, \ldots, 2n$ of elements in $A$. This sequence has the property that none of its elements divides any other element in the sequence. Note that the sequence has length $n$. The following theorem states that such a sequence of length $n + 1$ does not exist.

**Theorem 3.10.1** Let $n \geq 1$ and consider a sequence $a_1, a_2, \ldots, a_{n+1}$ of $n+1$ elements from the set $\{1, 2, \ldots, 2n\}$. Then there are two distinct indices $i$ and $j$ such that $a_i$ divides $a_j$ or $a_j$ divides $a_i$.

**Proof.** For each $i$ with $1 \leq i \leq n + 1$, write

$$a_i = 2^{k_i} \cdot q_i,$$

where $k_i \geq 0$ is an integer and $q_i$ is an odd integer. For example,

- if $a_i = 48$, then $k_i = 4$ and $q_i = 3$, because $48 = 2^4 \cdot 3$,
3.10. The Pigeonhole Principle

• if \( a_i = 1 \), then \( k_i = 0 \) and \( q_i = 1 \), because \( 1 = 2^0 \cdot 1 \),

• if \( a_i = 7 \), then \( k_i = 0 \) and \( q_i = 7 \), because \( 7 = 2^0 \cdot 7 \).

Consider the sequence \( q_1, q_2, \ldots, q_{n+1} \) of \( n+1 \) integers. Each of these numbers is an odd integer that belongs to the set

\[ \{1, 3, 5, \ldots, 2n - 1\} = \{2m - 1 : m = 1, 2, \ldots, n\}. \]

Since this set has size \( n \), the Pigeonhole Principle implies that there must be two numbers in the sequence \( q_1, q_2, \ldots, q_{n+1} \) that are equal. In other words, there are two distinct indices \( i \) and \( j \) such that \( q_i = q_j \). It follows that

\[ \frac{a_i}{a_j} = \frac{2^{k_i} \cdot q_i}{2^{k_j} \cdot q_j} = 2^{k_i-k_j}. \]

Thus, if \( k_i \geq k_j \), then \( a_j \) divides \( a_i \). Otherwise, \( k_i < k_j \), and \( a_i \) divides \( a_j \). □

3.10.3 Long Sequences Contain Long Monotone Subsequences

Let \( n = 3 \), and consider the sequence \( 20, 10, 9, 7, 11, 2, 21, 1, 20, 31 \) of \( 10 = n^2 + 1 \) numbers. This sequence contains an increasing subsequence of length \( 4 = n + 1 \), namely \( 10, 11, 21, 31 \). The following theorem states this result for arbitrary values of \( n \).

**Theorem 3.10.2** Let \( n \geq 1 \) be an integer. Every sequence of \( n^2 + 1 \) distinct real numbers contains a subsequence of length \( n + 1 \) that is either increasing or decreasing.

**Proof.** Let \( a_1, a_2, \ldots, a_{n^2+1} \) be an arbitrary sequence of \( n^2 + 1 \) distinct real numbers. For each \( i \) with \( 1 \leq i \leq n^2 + 1 \), let \( inc_i \) denote the length of the longest increasing subsequence that starts at \( a_i \), and let \( dec_i \) denote the length of the longest decreasing subsequence that starts at \( a_i \).

Using this notation, the claim in the theorem can be formulated as follows: There is an index \( i \) such that \( inc_i \geq n + 1 \) or \( dec_i \geq n + 1 \).

We will prove the claim by contradiction. Thus, we assume that \( inc_i \leq n \) and \( dec_i \leq n \) for all \( i \) with \( 1 \leq i \leq n^2 + 1 \).
Consider the set
\[ B = \{(b, c) : 1 \leq b \leq n, 1 \leq c \leq n\}, \]
and think of the elements of \( B \) as being boxes. For each \( i \) with \( 1 \leq i \leq n^2 + 1 \), the pair \((inc_i, dec_i)\) is an element of \( B \). Thus, we have \( n^2 + 1 \) elements \((inc_i, dec_i)\), which are placed in the \( n^2 \) boxes of \( B \). By the Pigeonhole Principle, there must be a box that contains two (or more) elements. In other words, there exist two integers \( i \) and \( j \) such that \( i < j \) and
\[ (inc_i, dec_i) = (inc_j, dec_j). \]
Recall that the elements in the sequence are distinct. Hence, \( a_i \neq a_j \). We consider two cases.

First assume that \( a_i < a_j \). Then the length of the longest increasing subsequence starting at \( a_i \) must be at least \( 1 + inc_j \), because we can append \( a_i \) to the longest increasing subsequence starting at \( a_j \). Therefore, \( inc_i \neq inc_j \), which is a contradiction.

The second case is when \( a_i > a_j \). Then the length of the longest decreasing subsequence starting at \( a_i \) must be at least \( 1 + dec_j \), because we can append \( a_i \) to the longest decreasing subsequence starting at \( a_j \). Therefore, \( dec_i \neq dec_j \), which is again a contradiction.

### 3.10.4 There are Infinitely Many Primes

As a final application of the Pigeonhole Principle, we prove the following result:

**Theorem 3.10.3** There are infinitely many prime numbers.

**Proof.** The proof is by contradiction. Thus, we assume that there are, say, \( k \) prime numbers, and denote them by
\[ 2 = p_1 < p_2 < \cdots < p_k. \]
Note that \( k \) is a fixed integer. Since
\[ \lim_{n \to \infty} \frac{2^n}{(n + 1)^k} = \infty, \]
we can choose an integer \( n \) such that

\[
2^n > (n + 1)^k.
\]

Define the function

\[
f : \{1, 2, \ldots, 2^n\} \to \mathbb{N}^k
\]

as follows: For any integer \( x \) with \( 1 \leq x \leq 2^n \), consider its prime factorization

\[
x = p_1^{m_1} \cdot p_2^{m_2} \cdots p_k^{m_k}.
\]

We define

\[
f(x) = (m_1, m_2, \ldots, m_k).
\]

Since

\[
m_i \leq m_1 + m_2 + \cdots + m_k
\leq m_1 \log p_1 + m_2 \log p_2 + \cdots + m_k \log p_k
= \log (p_1^{m_1} \cdot p_2^{m_2} \cdots p_k^{m_k})
= \log x
\leq n,
\]

it follows that

\[
f(x) \in \{0, 1, 2, \ldots, n\}^k.
\]

Thus, \( f \) is a function

\[
f : \{1, 2, \ldots, 2^n\} \to \{0, 1, 2, \ldots, n\}^k.
\]

It is easy to see that this function is one-to-one. The set on the left-hand side has size \( 2^n \), whereas the set on the right-hand side has size \( (n + 1)^k \). It then follows from the Pigeonhole Principle that

\[
(n + 1)^k \geq 2^n,
\]

which contradicts our choice for \( n \).
3.11 Exercises

3.1 A licence plate number consists of a sequence of four uppercase letters followed by three digits. How many licence plate numbers are there?

3.2 A multiple-choice exam consists of 100 questions. Each question has four possible answers $a$, $b$, $c$, and $d$. How many ways are there to answer the 100 questions (assuming that each question is answered)?

3.3 For each of the following 7 cases, determine how many strings of 8 uppercase letters there are.

- Letters can be repeated.

- No letter can be repeated.

- The strings start with PQ (in this order) and letters can be repeated.

- The strings start with PQ (in this order) and no letter can be repeated.

- The strings start and end with PQ (in this order) and letters can be repeated.

- The strings start with XYZ (in this order), end with QP (in this order), and letters can be repeated.

- The strings start with XYZ (in this order) or end with QP (in this order), and letters can be repeated.

3.4 The Carleton Computer Science Society has a Board of Directors consisting of one president, one vice-president, one secretary, one treasurer, and a three-person party committee (whose main responsibility is to buy beer for the other four board members). The entire board consists of seven distinct students. If there are $n \geq 7$ students in Carleton’s Computer Science program, how many ways are there to choose a Board of Directors?

3.5 The Carleton Computer Science Society has an Academic Events Committee (AEC) consisting of five students and a Beer Committee (BC) consisting of six students (whose responsibility is to buy beer for the AEC).
3.11. Exercises

• Assume there are \( n \geq 6 \) students in Carleton’s Computer Science program. Also, assume that a student can be both on the AEC and on the BC. What is the total number of ways in which these two committees can be chosen?

• Assume there are \( n \geq 11 \) students in Carleton’s Computer Science program. Also, assume that a student cannot be both on the AEC and on the BC. What is the total number of ways in which these two committees can be chosen?

3.6 Let \( m \) and \( n \) be integers with \( 0 \leq m \leq n \). There are \( n + 1 \) students in Carleton’s Computer Science program. The Carleton Computer Science Society has a Board of Directors, consisting of one president and \( m \) vice-presidents. The president cannot be vice-president. Prove that

\[
(n + 1) \binom{n}{m} = (n + 1 - m) \binom{n + 1}{m},
\]

by counting, in two different ways, the number of ways to choose a Board of Directors.

3.7 In Tic-Tac-Toe, we are given a 3 \( \times \) 3 grid, consisting of unmarked cells. Two players, Xavier and Olivia, take turns marking the cells of this grid. When it is Xavier’s turn, he chooses an unmarked cell and marks it with the letter \( X \). Similarly, when it is Olivia’s turn, she chooses an unmarked cell and marks it with the letter \( O \). The first turn is by Xavier. The players continue making turns until all cells have been marked. Below, you see an example of a completely marked grid.

<table>
<thead>
<tr>
<th></th>
<th>O</th>
<th>O</th>
<th>X</th>
</tr>
</thead>
<tbody>
<tr>
<td>X</td>
<td>X</td>
<td>O</td>
<td></td>
</tr>
<tr>
<td>X</td>
<td>X</td>
<td>O</td>
<td></td>
</tr>
</tbody>
</table>

• What is the number of completely marked grids?

• What is the number of different ways (i.e., ordered sequences) in which the grid can be completely marked, when following the rules given above?
3.8 In how many ways can you paint 200 chairs, if 33 of them must be painted red, 66 of them must be painted blue, and 101 of them must be painted green?

3.9 Let $A$ be a set of size $m$, let $B$ be a set of size $n$, and assume that $n \geq m \geq 1$. How many functions $f : A \to B$ are there that are not one-to-one?

3.10 How many bitstrings of length 8 are there that contain at least 4 consecutive 0s or at least 4 consecutive 1s?

3.11 How many bitstrings of length 77 are there that start with 010 (i.e., have 010 at positions 1, 2, and 3) or have 101 at positions 2, 3, and 4, or have 010 at positions 3, 4, and 5?

3.12 Let $m$ and $n$ be integers with $m \geq n \geq 1$. How many ways are there to place $m$ books on $n$ shelves, if there must be at least one book on each shelf? As in Section 3.1.3 the order on each shelf matters.

3.13 You are given $m$ distinct books $B_1, B_2, \ldots, B_m$ and $n$ identical blocks of wood. How many ways are there to arrange these books and blocks in a straight line?

For example, if $m = 5$ and $n = 3$, then three possible arrangements are

$(W$ stands for a block of wood)$

\begin{align*}
WB_3B_1WB_5B_4WB_2, \\
WB_1B_3WB_5B_4WB_2,
\end{align*}

and

\begin{align*}
B_3WB_3B_1WWB_2B_4.
\end{align*}

3.14 Let $n \geq 1$ be an integer and consider $n$ boys and $n$ girls. For each of the following three cases, determine how many ways there are to arrange these $2n$ people on a straight line (the order on the line matters):

- All boys stand next to each other and all girls stand next to each other.
- All girls stand next to each other.
- Boys and girls alternate.
3.15 Let \( n \geq 1 \) be an integer. Consider a tennis tournament with \( 2n \) participants. In the first round of this tournament, \( n \) games will be played and, thus, the \( 2n \) people have to be divided into \( n \) pairs. What is the number of ways in which this can be done?

3.16 The Beer Committee of the Carleton Computer Science Society has bought large quantities of 10 different types of beer. In order to test which beer students prefer, the committee does the following experiment:

- Out of the \( n \geq 10 \) students in Carleton’s Computer Science program, 10 students are chosen.
- Each of the 10 students chosen drinks one of the 10 beers; no two students drink the same beer.

What is the number of ways in which this experiment can be done?

3.17 Consider strings consisting of 12 characters, each character being \( a, b, \) or \( c \). Such a string is called valid if at least one of the characters is missing. For example, \( abababababab \) is a valid string, whereas \( ababacabab \) is not a valid string. How many valid strings are there?

3.18 A password consists of 100 characters, each character being a digit or a lowercase letter. A password must contain at least two digits. How many passwords are there?

3.19 A password consists of 100 characters, each character being a digit, a lowercase letter, or an uppercase letter. A password must contain at least one digit, at least one lowercase letter, and at least one uppercase letter. How many passwords are there?

Hint: Recall De Morgan’s Law

\[ A \cap B \cap C = \overline{A \cup B \cup C}. \]

3.20 A password is a string of 100 characters, where each character is a digit or a lowercase letter. A password is called valid if

- it does not start with \( abc \), and
- it does not end with \( xyz \), and
• it does not start with 3456.

Determine the number of valid passwords.

3.21 Consider permutations of the 26 lowercase letters $a$, $b$, $c$, \ldots, $z$.

• How many such permutations contain the string $wine$?

• How many such permutations do not contain any of the strings $wine$, $vodka$, or $coke$?

3.22 Let $n$ and $k$ be integers with $2 \leq k \leq n$, and consider the set $S = \{1, 2, 3, \ldots, 2n\}$. An ordered sequence of $k$ elements of $S$ is called valid if

• this sequence is strictly increasing, or

• this sequence is strictly decreasing, or

• this sequence contains only even numbers (and duplicate elements are allowed).

Determine the number of valid sequences.

3.23 In a group of 20 people,

• 6 are blond,

• 7 have green eyes,

• 11 are not blond and do not have green eyes.

How many people in this group are blond and have green eyes?

3.24 Use Pascal’s Identity (Theorem 3.7.2) to prove Newton’s Binomial Theorem (Theorem 3.6.5) by induction.

3.25 Determine the coefficient of $x^{111}y^{444}$ in the expansion of

$$(-17x + 71y)^{555}.$$
3.11. Exercises

3.26 Let \( n \geq 1 \) be an integer. Use Newton’s Binomial Theorem (Theorem 3.6.5) to prove that
\[
\sum_{k=1}^{n} \binom{n}{k} 10^k \cdot 26^{n-k} = 36^n - 26^n. \tag{3.5}
\]

In the rest of this exercise, you will give a combinatorial proof of this identity.

Consider passwords consisting of \( n \) characters, each character being a digit or a lowercase letter. A password must contain at least one digit.

- Use the Complement Rule of Section 3.3 to show that the number of passwords is equal to \( 36^n - 26^n \).
- Let \( k \) be an integer with \( 1 \leq k \leq n \). Prove that the number of passwords with exactly \( k \) digits is equal to \( \binom{n}{k} 10^k \cdot 26^{n-k} \).
- Explain why the above two parts imply the identity in (3.5).

3.27 Use Newton’s Binomial Theorem (Theorem 3.6.5) to prove that for every integer \( n \geq 1 \),
\[
\sum_{k=0}^{n} \binom{n}{k} 2^k = 3^n.
\]

In the rest of this exercise, you will give a combinatorial proof of this identity.

Let \( A = \{1, 2, 3, \ldots, n\} \) and \( B = \{a, b, c\} \). According to Theorem 3.1.2, the number of functions \( f : A \rightarrow B \) is equal to \( 3^n \).

- Consider a fixed integer \( k \) with \( 0 \leq k \leq n \) and a fixed subset \( S \) of \( A \) having size \( k \). How many functions \( f : A \rightarrow B \) are there for which
  - for all \( x \in S \), \( f(x) \in \{a, b\} \), and
  - for all \( x \in A \setminus S \), \( f(x) = c \).

- Argue that the summation \( \sum_{k=0}^{n} \binom{n}{k} 2^k \) counts all functions \( f : A \rightarrow B \).

3.28 Let \( n \geq 66 \) be an integer and consider the set \( S = \{1, 2, \ldots, n\} \).

- Let \( k \) be an integer with \( 66 \leq k \leq n \). How many 66-element subsets of \( S \) are there whose largest element is equal to \( k \)?
• Use the result in the first part to prove that
\[ \sum_{k=65}^{n} \binom{k-1}{65} = \binom{n}{66}. \]

3.29 Let \( n \geq 0 \) and \( k \geq 0 \) be integers.

• How many bitstrings of length \( n+1 \) have exactly \( k+1 \) many 1s?

• Let \( i \) be an integer with \( k \leq i \leq n \). What is the number of bitstrings of length \( n+1 \) that have exactly \( k+1 \) many 1s and in which the rightmost 1 is at position \( i+1 \)?

• Use the above two results to prove that
\[ \sum_{i=k}^{n} \binom{i}{k} = \binom{n+1}{k+1}. \]

3.30 Let \( k, m, \) and \( n \) be integers with \( 0 \leq k \leq m \leq n \), and let \( S \) be a set of size \( n \). Prove that
\[ \binom{n}{k} \binom{n-k}{m-k} = \binom{n}{m} \binom{m}{k}, \]
by counting, in two different ways, the number of ordered pairs \((A, B)\) with \( A \subseteq S, B \subseteq S, A \subseteq B, |A| = k, \) and \( |B| = m \).

3.31 Let \( m \) and \( n \) be integers with \( 0 \leq m \leq n \).

• How many bitstrings of length \( n+1 \) have exactly \( m \) many 1s?

• Let \( k \) be an integer with \( 0 \leq k \leq m \). What is the number of bitstrings of length \( n+1 \) that have exactly \( m \) many 1s and that start with \( \underbrace{1\ldots1}_{k}0 \)?

• Use the above two results to prove that
\[ \sum_{k=0}^{m} \binom{n-k}{m-k} = \binom{n+1}{m}. \]
3.32 Let $m$ and $n$ be integers with $0 \leq m \leq n$. Use Exercises 3.6, 3.30, and 3.31 to prove that

$$\sum_{k=0}^{m} \binom{m}{k} = \frac{n+1}{n+1-m}.$$ 

3.33 Let $m \geq 1$ and $n \geq 1$ be integers. Consider a rectangle whose horizontal side has length $m$ and whose vertical side has length $n$. A path from the bottom-left corner to the top-right corner is called valid, if in each step, it either goes one unit to the right or one unit upwards. In the example below, you see a valid path for the case when $m = 5$ and $n = 3$.

How many valid paths are there?

3.34 Let $n \geq 1$ be an integer. Prove that

$$\sum_{k=1}^{n} k \binom{n}{k} = n \cdot 2^{n-1}. $$

*Hint:* Take the derivative of $(1 + x)^n$.

3.35 How many different strings can be obtained by reordering the letters of the word MississippiMills. (This is a town close to Ottawa. James Naismith, the inventor of basketball, was born there.)

3.36 In Theorems 3.9.1 and 3.9.2, we have seen how many solutions (in nonnegative integers) there are for equations of the type

$$x_1 + x_2 + \cdots + x_k = n$$
and inequalities of the type

\[ x_1 + x_2 + \cdots + x_k \leq n. \]

Use this to prove the following identity:

\[ \sum_{i=0}^{n} \binom{i+k-1}{k-1} = \binom{n+k}{k}. \]

3.37 Let \( n \) and \( k \) be integers with \( n \geq k \geq 1 \). How many solutions are there to the equation

\[ x_1 + x_2 + \cdots + x_k = n, \]

where \( x_1 \geq 1, x_2 \geq 1, \ldots, x_k \geq 1 \) are integers?

*Hint:* In Theorem 3.9.1 we have seen the answer if \( x_1 \geq 0, x_2 \geq 0, \ldots, x_k \geq 0 \).

3.38 Let \( n \geq 1 \) be an integer. Use the Pigeonhole Principle to prove that in any set of \( n+1 \) integers from \( \{1, 2, \ldots, 2n\} \), there are two integers that are consecutive (i.e., differ by one).

3.39 Consider five points in a square with sides of length one. Use the Pigeonhole Principle to prove that there are two of these points having distance at most \( 1/\sqrt{2} \).

3.40 Let \( S \) be a set of 90 positive integers, each one having at most 25 digits in decimal notation. Use the Pigeonhole Principle to prove that there are two different subsets \( A \) and \( B \) of \( S \) that have the same sum, i.e.,

\[ \sum_{x \in A} x = \sum_{x \in B} x. \]

3.41 Let \( n \geq 1 \) be an integer and consider the set \( S = \{1, 2, \ldots, 2n\} \). Let \( T \) be an arbitrary subset of \( S \) having size \( n+1 \). Prove that this subset \( T \) contains two elements whose sum is equal to \( 2n + 1 \).

*Hint:* Consider the pairs \( (1, 2n), (2, 2n - 1), (3, 2n - 2), \ldots, (n, n + 1) \) and use the Pigeonhole Principle.

3.42 In this exercise, we consider the sequence

\[ 3^0, 3^1, 3^2, \ldots, 3^{1000} \]

of integers.
3.11. Exercises

- Prove that this sequence contains two distinct elements whose difference is divisible by 1000. That is, prove that there exist two integers \( m \) and \( n \) with \( 0 \leq m < n \leq 1000 \), such that \( 3^n - 3^m \) is divisible by 1000.

  *Hint:* Consider each element in the sequence modulo 1000 and use the Pigeonhole Principle.

- Use the first part to prove that the sequence

\[ 3^1, 3^2, \ldots, 3^{1000} \]

contains an element whose decimal representation ends with 001. In other words, the last three digits in the decimal representation are 001.
Chapter 4

Recursion

In order to understand recursion, you must first understand recursion.

Recursion is the concept where an object (such as a function, a set, or an algorithm) is defined in the following way:

- There are one or more base cases.
- There are one or more rules that define an object in terms of “smaller” objects that have already been defined.

In this chapter, we will see several examples of such recursive definitions and how to use them to solve counting problems.

4.1 Recursive Functions

Recall that $\mathbb{N} = \{0, 1, 2, \ldots\}$ denotes the set of natural numbers. Consider the following recursive definition of a function $f : \mathbb{N} \to \mathbb{N}$:

\[
\begin{align*}
    f(0) &= 3, \\
    f(n) &= 2 \cdot f(n - 1) + 3, \text{ if } n \geq 1.
\end{align*}
\]

These two rules indeed define a function, because $f(0)$ is uniquely defined and for any integer $n \geq 1$, if $f(n - 1)$ is uniquely defined, then $f(n)$ is also uniquely defined, because it is equal to $2 \cdot f(n - 1) + 3$. Therefore, by induction, for any natural number $n$, the function value $f(n)$ is uniquely defined. We can obtain the values $f(n)$ in the following way:
• We are given that $f(0) = 3$.

• If we apply the recursive rule with $n = 1$, then we get
  
  $$f(1) = 2 \cdot f(0) + 3 = 2 \cdot 3 + 3 = 9.$$  

• If we apply the recursive rule with $n = 2$, then we get
  
  $$f(2) = 2 \cdot f(1) + 3 = 2 \cdot 9 + 3 = 21.$$  

• If we apply the recursive rule with $n = 3$, then we get
  
  $$f(3) = 2 \cdot f(2) + 3 = 2 \cdot 21 + 3 = 45.$$  

• If we apply the recursive rule with $n = 4$, then we get
  
  $$f(4) = 2 \cdot f(3) + 3 = 2 \cdot 45 + 3 = 93.$$  

Can we “solve” this recurrence? That is, can we express $f(n)$ in terms of $n$ only? By looking at these values, you may see a pattern, i.e., you may guess that for each $n \geq 0$,  

$$f(n) = 3 \cdot 2^{n+1} - 3.$$  

(4.1)

We prove by induction that this is correct: If $n = 0$, then $f(n) = f(0) = 3$ and $3 \cdot 2^{n+1} - 3 = 3 \cdot 2^1 - 3 = 3$. Thus, (4.1) is true for $n = 0$. Let $n \geq 1$ and assume that (4.1) is true for $n - 1$, i.e., assume that

$$f(n - 1) = 3 \cdot 2^n - 3.$$  

Then

$$f(n) = 2 \cdot f(n - 1) + 3$$

$$= 2 (3 \cdot 2^n - 3) + 3$$

$$= 3 \cdot 2^{n+1} - 3.$$  

Thus, we have proved by induction that (4.1) holds for all integers $n \geq 0$.  

A recursive definition of factorials: Consider the following recursive definition of a function $g : \mathbb{N} \to \mathbb{N}$:

$$
\begin{align*}
g(0) &= 1, \\
g(n) &= n \cdot g(n-1), \text{ if } n \geq 1.
\end{align*}
$$

As in the previous example, a simple induction proof shows that these rules uniquely define the value $g(n)$ for each $n \geq 0$. We leave it to the reader to verify that $g$ is the factorial function, i.e., $g(n) = n!$ for each $n \geq 0$.

A recursive definition of binomial coefficients: Consider the following recursive definition of a function $B : \mathbb{N} \times \mathbb{N} \to \mathbb{N}$ with two variables:

$$
\begin{align*}
B(n, 0) &= 1, \text{ if } n \geq 0, \\
B(n, n) &= 1, \text{ if } n \geq 0, \\
B(n, k) &= B(n-1, k-1) + B(n-1, k), \text{ if } n \geq 2 \text{ and } 1 \leq k \leq n-1.
\end{align*}
$$

The recursive rule has the same form as Pascal’s Identity in Theorem 3.7.2. Also, the first base case shows that $B(n, 0) = 1 = \binom{n}{0}$, whereas the second base case shows that $B(n, n) = 1 = \binom{n}{n}$. From this, it can be shown by induction that $B(n, k) = \binom{n}{k}$ for all $n$ and $k$ with $0 \leq k \leq n$.

### 4.2 Fibonacci Numbers

I’ll have an order of the Fibonachos.

The Fibonacci numbers are defined using the following rules:

$$
\begin{align*}
f_0 &= 0, \\
f_1 &= 1, \\
f_n &= f_{n-1} + f_{n-2}, \text{ if } n \geq 2.
\end{align*}
$$

In words, there are two base cases (i.e., 0 and 1) and each next element in the sequence is the sum of the previous two elements. This gives the sequence

$$
0, 1, 1, 2, 3, 5, 8, 13, 21, 34, 55, 89, 144, \ldots
$$

The following theorem states that we can “solve” this recurrence. That is, we can express the $n$-th Fibonacci number $f_n$ in a non-recursive way, i.e., without using any other Fibonacci numbers.
Theorem 4.2.1 Let \( \varphi = \frac{1 + \sqrt{5}}{2} \) and \( \psi = \frac{1 - \sqrt{5}}{2} \) be the two solutions of the quadratic equation \( x^2 = x + 1 \). Then, for all \( n \geq 0 \), we have
\[
 f_n = \frac{\varphi^n - \psi^n}{\sqrt{5}}.
\]

Proof. We prove the claim by induction on \( n \). There are two base cases:\footnote{Do you see why there are two base cases?}

- Both \( f_0 \) and \( \frac{\varphi^0 - \psi^0}{\sqrt{5}} \) are equal to 0.
- Both \( f_1 \) and \( \frac{\varphi^1 - \psi^1}{\sqrt{5}} \) are equal to 1.

Let \( n \geq 2 \) and assume that the claim is true for \( n - 2 \) and \( n - 1 \). In other words, assume that
\[
 f_{n-2} = \frac{\varphi^{n-2} - \psi^{n-2}}{\sqrt{5}}
\]
and
\[
 f_{n-1} = \frac{\varphi^{n-1} - \psi^{n-1}}{\sqrt{5}}.
\]

We have to prove that the claim is true for \( n \) as well. Using the definition of \( f_n \), the two assumptions, and the identities \( \varphi^2 = \varphi + 1 \) and \( \psi^2 = \psi + 1 \), we get
\[
 f_n = f_{n-1} + f_{n-2} = \frac{\varphi^{n-1} - \psi^{n-1}}{\sqrt{5}} + \frac{\varphi^{n-2} - \psi^{n-2}}{\sqrt{5}} = \frac{\varphi^{n-2}(\varphi + 1) - \psi^{n-2}(\psi + 1)}{\sqrt{5}} = \frac{\varphi^{n-2} \cdot \varphi^2 - \psi^{n-2} \cdot \psi^2}{\sqrt{5}} = \frac{\varphi^n - \psi^n}{\sqrt{5}}.
\]
4.2.1 Counting 00-Free Bitstrings

Consider finite bitstrings that do not contain 00. Examples of such bitstrings are 10, 010, 01010101, and 11110111. For any integer \( n \geq 2 \), how many such strings are there having length \( n \)? Since we do not know the answer yet, we introduce a variable \( B_n \), one for each \( n \geq 2 \), for the number of such strings. Thus,

- \( B_n \) denotes the number of bitstrings of length \( n \) that do not contain 00.

There are four bitstrings of length 2:

\[
00, 10, 01, 11.
\]

Since three of them do not contain 00, it follows that \( B_2 = 3 \). Similarly, there are eight bitstrings of length 3:

\[
000, 001, 010, 100, 011, 101, 110, 111.
\]

Since five of them do not contain 00, we have \( B_3 = 5 \).

Let \( n \geq 4 \). We are going to express \( B_n \) in terms of the previous two values \( B_{n-1} \) and \( B_{n-2} \). This, together with the two base cases \( B_2 = 3 \) and \( B_3 = 5 \), will give a recurrence for the entire sequence.

Consider a matrix that contains all bitstrings of length \( n \) that do not contain 00, one string per row. Since there are \( B_n \) many such strings, the matrix has \( B_n \) many rows. Also, the matrix has \( n \) columns, because the strings have length \( n \).

We rearrange the rows of the matrix such that all strings in the top part start with 1 and all strings in the bottom part start with 0.

- How many rows are there in the top part? Any string in the top part starts with 1 and is followed by a bitstring of length \( n - 1 \) that does not contain 00. Thus, if we take the rows in the top part and delete the first bit from each row, then we obtain all bitstrings of length \( n - 1 \) that do not contain 00. Since there are \( B_{n-1} \) many such strings of length \( n - 1 \), it follows that the top part of the matrix consists of \( B_{n-1} \) many rows.

- How many rows are there in the bottom part? Any string in the bottom part starts with 0. Since the string does not contain 00, the second bit
must be 1. After these first two bits, we have a bitstring of length \( n - 2 \) that does not contain 00. Thus, if we take the rows in the bottom part and delete the first two bits from each row, then we obtain all bitstrings of length \( n - 2 \) that do not contain 00. Since there are \( B_{n-2} \) many such strings of length \( n - 2 \), it follows that the bottom part of the matrix consists of \( B_{n-2} \) many rows.

Thus, on the one hand, the matrix has \( B_n \) many rows. On the other hand, this matrix has \( B_{n-1} + B_{n-2} \) many rows. Therefore, we have \( B_n = B_{n-1} + B_{n-2} \).

To summarize, we have proved that the values \( B_n \), for \( n \geq 2 \), satisfy the following recurrence:

\[
B_2 = 3,
B_3 = 5,
B_n = B_{n-1} + B_{n-2}, \text{ if } n \geq 4.
\]

This recurrence is the same as the one for the Fibonacci numbers, except that the two base cases are different. The sequence \( B_n \), \( n \geq 2 \), consists of the integers

\[3, 5, 8, 13, 21, 34, 55, 89, 144, \ldots\]

We obtain this sequence by removing the first four elements (i.e., \( f_0 \), \( f_1 \), \( f_2 \), and \( f_3 \)) from the Fibonacci sequence. We leave it to the reader to verify (using induction) that for all \( n \geq 2 \),

\[B_n = f_{n+2}.
\]

## 4.3 A Recursively Defined Set

Consider the set \( S \) that is defined by the following two rules:

- 5 is an element of the set \( S \).
- If \( x \) and \( y \) are elements of the set \( S \), then \( x - y \) is also an element of the set \( S \).

Thus, if we already know that \( x \) and \( y \) belong to the set \( S \), then the second rule gives us a new element, i.e., \( x - y \), that also belongs to \( S \).

Can we give a simple description of the set \( S \)? We are going to use the rules to obtain some elements of \( S \). From these examples, we then hope to see a pattern from which we guess the simple description of \( S \). The final step consists of proving that our guess is correct.
• We are given that 5 is an element of $S$.

• Applying the rule with $x = 5$ and $y = 5$ implies that $x - y = 0$ is also an element of $S$.

• Applying the rule with $x = 0$ and $y = 5$ implies that $x - y = -5$ is also an element of $S$.

• Applying the rule with $x = 5$ and $y = -5$ implies that $x - y = 10$ is also an element of $S$.

• Applying the rule with $x = 0$ and $y = 10$ implies that $x - y = -10$ is also an element of $S$.

• Applying the rule with $x = 5$ and $y = -10$ implies that $x - y = 15$ is also an element of $S$.

• Applying the rule with $x = 0$ and $y = 15$ implies that $x - y = -15$ is also an element of $S$.

Thus, we have obtained the following elements of $S$:

\[-15, -10, -5, 0, 5, 10, 15\]

Since there is clearly a pattern, it is natural to guess that

\[S = \{5n : n \in \mathbb{Z}\},\]  \hspace{1cm} (4.2)

where $\mathbb{Z}$ is the set of all (positive and negative) integers, including 0. To prove that this is correct, we will first prove that the set on the left-hand side is a subset of the set on the right-hand side. Then we prove that the set on the right-hand side is a subset of the set on the left-hand side.

We start by proving that

\[S \subseteq \{5n : n \in \mathbb{Z}\},\]

which is equivalent to proving that

\[\text{every element of } S \text{ is a multiple of } 5.\]  \hspace{1cm} (4.3)

How do we prove this? The set $S$ is defined using a base case and a recursive rule. The only way to obtain an element of $S$ is by starting with the base case and then applying the recursive rule a finite number of times. Therefore, the following will prove that (4.3) holds:
• The element in the base case, i.e., 5, is a multiple of 5.

• Let \( x \) and \( y \) be two elements of \( S \) and assume that they are both multiples of 5. Then \( x - y \) (which is the “next” element of \( S \)) is also a multiple of 5.

Next we prove that
\[
\{5n : n \in \mathbb{Z}\} \subseteq S.
\]
We will do this by proving that for all \( n \geq 0 \),
\[
5n \in S \quad \text{and} \quad -5n \in S. \tag{4.4}
\]
The proof is by induction on \( n \). For the base case, i.e., when \( n = 0 \), we observe that, from the definition of \( S \), \( x = 5 \) and \( y = 5 \) are in \( S \) and, therefore, \( x - y = 0 \) is also in \( S \). Therefore, (4.4) is true for \( n = 0 \).

Let \( n \geq 0 \) and assume that (4.4) is true for \( n \), i.e., assume that
\[
5n \in S \quad \text{and} \quad -5n \in S.
\]
We have to show that (4.4) is also true for \( n + 1 \), i.e.,
\[
5(n + 1) \in S \quad \text{and} \quad -5(n + 1) \in S.
\]
• It follows from the definition of \( S \) and our assumption that both \( x = 5 \) and \( y = -5n \) are in \( S \). Therefore, \( x - y = 5(n + 1) \) is also in \( S \).

• It follows from the definition of \( S \) and our assumption that both \( x = -5n \) and \( y = 5 \) are in \( S \). Therefore, \( x - y = -5(n + 1) \) is also in \( S \).

Thus, we have shown by induction that (4.4) holds for all \( n \geq 0 \).

Since we have shown that both (4.3) and (4.4) hold, we conclude that (4.2) holds as well. In other words, we have indeed obtained a simple description of the set \( S \): It is the set of all multiples of 5.

### 4.4 A Gossip Problem

Let \( n \geq 4 \) be an integer and consider a group \( P_1, P_2, \ldots, P_n \) of \( n \) people. Assume that each person \( P_i \) knows some scandal \( S_i \) that nobody else knows. For any \( i \) and \( j \), if person \( P_i \) makes a phone call with person \( P_j \), they exchange
the scandals they know at that moment, i.e., $P_i$ tells all scandals he knows to $P_j$, and $P_j$ tells all scandals he knows to $P_i$. How many phone calls are needed until each of the $n$ people knows all $n$ scandals?

An obvious solution is that each pair of people in the group makes one phone call. At the end, each person knows all scandals. The number of phone calls is

$$\binom{n}{2} = \frac{n(n-1)}{2},$$

which is quadratic in the number $n$ of people. We will see below that only a linear number of phone calls are needed.

Let us first consider the case when $n = 4$. At the start, each person $P_i$ only knows the scandal $S_i$, which we visualize in the following table:

<table>
<thead>
<tr>
<th></th>
<th>$P_1$</th>
<th>$P_2$</th>
<th>$P_3$</th>
<th>$P_4$</th>
</tr>
</thead>
<tbody>
<tr>
<td>$S_1$</td>
<td>$S_1$</td>
<td>$S_1$</td>
<td>$S_1$</td>
<td>$S_1$</td>
</tr>
<tr>
<td>$S_2$</td>
<td>$S_2$</td>
<td>$S_2$</td>
<td>$S_2$</td>
<td>$S_2$</td>
</tr>
</tbody>
</table>

Consider the following sequence of phone calls:

1. $P_1$ calls $P_2$. After this phone call, the table looks as follows:

<table>
<thead>
<tr>
<th></th>
<th>$P_1$</th>
<th>$P_2$</th>
<th>$P_3$</th>
<th>$P_4$</th>
</tr>
</thead>
<tbody>
<tr>
<td>$S_1$</td>
<td>$S_1$</td>
<td>$S_1$</td>
<td>$S_1$</td>
<td>$S_1$</td>
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<tr>
<td>$S_2$</td>
<td>$S_2$</td>
<td>$S_2$</td>
<td>$S_2$</td>
<td>$S_2$</td>
</tr>
</tbody>
</table>

2. $P_3$ calls $P_4$. After this phone call, the table looks as follows:

<table>
<thead>
<tr>
<th></th>
<th>$P_1$</th>
<th>$P_2$</th>
<th>$P_3$</th>
<th>$P_4$</th>
</tr>
</thead>
<tbody>
<tr>
<td>$S_1$</td>
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<td>$S_1$</td>
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<tr>
<td>$S_2$</td>
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<td>$S_4$</td>
<td>$S_4$</td>
</tr>
</tbody>
</table>

3. $P_1$ calls $P_3$. After this phone call, the table looks as follows:

<table>
<thead>
<tr>
<th></th>
<th>$P_1$</th>
<th>$P_2$</th>
<th>$P_3$</th>
<th>$P_4$</th>
</tr>
</thead>
<tbody>
<tr>
<td>$S_1$</td>
<td>$S_1$</td>
<td>$S_1$</td>
<td>$S_1$</td>
<td>$S_1$</td>
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<tr>
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<td>$S_4$</td>
<td>$S_4$</td>
</tr>
</tbody>
</table>

4. $P_2$ calls $P_4$. After this phone call, the table looks as follows:

<table>
<thead>
<tr>
<th></th>
<th>$P_1$</th>
<th>$P_2$</th>
<th>$P_3$</th>
<th>$P_4$</th>
</tr>
</thead>
<tbody>
<tr>
<td>$S_1$</td>
<td>$S_1$</td>
<td>$S_1$</td>
<td>$S_1$</td>
<td>$S_1$</td>
</tr>
<tr>
<td>$S_2$</td>
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<td>$S_4$</td>
<td>$S_4$</td>
</tr>
</tbody>
</table>
We see that after four phone calls, each person knows all four scandals. Observe that the number of phone calls is \( \binom{4}{2} = 6 \) if we would have used the obvious solution mentioned above.

Thus, we now have an algorithm that schedules the phone calls for groups of four people. Below, we will extend this “base case” to a recursive algorithm that schedules the phone calls for any group of \( n \geq 4 \) people. The approach is as follows:

- We assume that we know how to schedule the phone calls for groups of \( n - 1 \) people.
- We use this assumption to schedule the phone calls for groups of \( n \) people.

Let us see how this is done.

- At the start, \( P_1 \) knows \( S_1 \), \( P_2 \) knows \( S_2 \), \ldots, \( P_n \) knows \( S_n \).
- \( P_{n-1} \) calls \( P_n \). After this phone call, \( P_1 \) knows \( S_1 \), \( P_2 \) knows \( S_2 \), \ldots, \( P_{n-2} \) knows \( S_{n-2} \), and both \( P_{n-1} \) and \( P_n \) know \( S_{n-1} \) and \( S_n \).
- Consider \( S_{n-1} \) and \( S_n \) to be one scandal \( S'_{n-1} \).
- Schedule the phone calls for the group \( P_1, P_2, \ldots, P_{n-1} \) of \( n - 1 \) people, using the scandals \( S_1, S_2, \ldots, S_{n-2}, S'_{n-1} \). (We have assumed that we know how to do this!) At the end, each of \( P_1, P_2, \ldots, P_{n-1} \) knows all scandals \( S_1, S_2, \ldots, S_n \).
- At this moment, \( P_n \) only knows \( S_{n-1} \) and \( S_n \). Therefore, \( P_{n-1} \) again calls \( P_n \) and tells him all scandals \( S_1, S_2, \ldots, S_{n-2} \).

Below, you see this recursive algorithm in pseudocode.
Algorithm GOSSIP($n$):

```plaintext
// n ≥ 4, this algorithm schedules phone calls for $P_1, P_2, \ldots, P_n$
if $n = 4$
  then $P_1$ calls $P_2$;
  $P_3$ calls $P_4$;
  $P_1$ calls $P_3$;
  $P_2$ calls $P_4$
else $P_{n-1}$ calls $P_n$;
  GOSSIP($n - 1$);
  $P_{n-1}$ calls $P_n$
endif
```

What is the number of phone calls made when running algorithm GOSSIP($n$)? Since we do not know the answer yet, we introduce a variable $C(n)$ to denote this number. It follows from the pseudocode that

$$C(4) = 4.$$

Let $n \geq 5$. Algorithm GOSSIP($n$) starts and ends with the same phone call: $P_{n-1}$ calls $P_n$. In between, it runs algorithm GOSSIP($n - 1$), during which, by definition, $C(n - 1)$ phone calls are made. It follows that

$$C(n) = 2 + C(n - 1) \text{ for } n \geq 5.$$

Thus, we have obtained a recurrence for the numbers $C(n)$. The first few numbers in the sequence are

$$
C(4) = 4, \\
C(5) = 2 + C(4) = 2 + 4 = 6, \\
C(6) = 2 + C(5) = 2 + 6 = 8, \\
C(7) = 2 + C(6) = 2 + 8 = 10.
$$

From this, we guess that

$$C(n) = 2n - 4 \text{ for } n \geq 4.$$

We can easily prove by induction that our guess is correct. Indeed, since both $C(4)$ and $2 \cdot 4 - 4$ are equal to 4, the claim is true for $n = 4$. If $n \geq 5$ and $C(n - 1) = 2(n - 1) - 4$, then

$$C(n) = 2 + C(n - 1) = 2 + (2(n - 1) - 4) = 2n - 4.$$
This shows that \( C(n) = 2n - 4 \) for all \( n \geq 4 \).

It can be shown that algorithm GOSSIP is optimal: Any algorithm that schedules phone calls for \( n \geq 4 \) people must make at least \( 2n - 4 \) phone calls.

You may wonder why the base case for algorithm \( \text{GOSSIP}(n) \) is when \( n = 4 \). You will find the reason in Exercise 4.27.

### 4.5 The Merge-Sort Algorithm

\( \text{MERGE} \text{SORT} \) is a recursive sorting algorithm that works as follows. To sort the sequence \( a_1, a_2, \ldots, a_n \) of numbers,

- it recursively sorts the sequence \( a_1, a_2, \ldots, a_m \), where \( m = \lfloor n/2 \rfloor \), and stores the sorted sequence in a list \( L_1 \),
- it recursively sorts the sequence \( a_{m+1}, a_{m+2}, \ldots, a_n \) and stores the sorted sequence in a list \( L_2 \),
- it merges the two sorted lists \( L_1 \) and \( L_2 \) into one sorted list.

Below, you see this recursive algorithm in pseudocode.

```plaintext
Algorithm \text{MERGE} \text{SORT}(L, n):

// \( L \) is a list of \( n \geq 0 \) numbers
if \( n \geq 2 \)
then \( m = \lfloor n/2 \rfloor \);
\( L_1 \) = list containing the first \( m \) elements of \( L \);
\( L_2 \) = list containing the last \( n - m \) elements of \( L \);
\( L_1 \) = \text{MERGE} \text{SORT}(L_1, m);
\( L_2 \) = \text{MERGE} \text{SORT}(L_2, n - m);
\( L \) = \text{MERGE}(L_1, L_2);
endif;
return \( L \)
```

We still have to specify algorithm \( \text{MERGE}(L_1, L_2) \). Of course, this algorithm uses the fact that both \( L_1 \) and \( L_2 \) are sorted lists. The task is to merge them into one sorted list. This is done in the following way. Initialize an empty list \( L \). (At the end, this list will contain the final sorted sequence.)

- Let \( x \) be the first element of \( L_1 \) and let \( y \) be the first element of \( L_2 \).
• If \( x \leq y \), then remove \( x \) from \( L_1 \) and append it to \( L \) (i.e., add \( x \) at the end of \( L \)).

• Otherwise (i.e., if \( x > y \)), remove \( y \) from \( L_2 \) and append it to \( L \).

Repeat these steps until one of \( L_1 \) and \( L_2 \) is empty. If \( L_1 \) is empty, then append \( L_2 \) to \( L \). Otherwise, append \( L_1 \) to \( L \). Here is the algorithm in pseudocode:

```plaintext
Algorithm Merge\((L_1, L_2)\):

// \( L_1 \) and \( L_2 \) are sorted lists
\( L = \) empty list;
while \( L_1 \) is not empty and \( L_2 \) is not empty
  do \( x = \) first element of \( L_1 \);
      \( y = \) first element of \( L_2 \);
      if \( x \leq y \)
        then remove \( x \) from \( L_1 \);
            append \( x \) to \( L \)
        else remove \( y \) from \( L_2 \);
            append \( y \) to \( L \)
      endif
  enddo
endwhile;
if \( L_1 \) is non-empty
  then append \( L_1 \) to \( L \)
else append \( L_2 \) to \( L \)
endif;
return \( L \)
```

### 4.5.1 Correctness of Algorithm MergeSort

I hope you are convinced that the output \( L \) of algorithm Merge\((L_1, L_2)\) is a sorted list that contains all elements of \( L_1 \) and \( L_2 \) (and no other elements). How do we prove that algorithm MergeSort\((L, n)\) is correct, i.e., correctly sorts the elements in any list \( L \) of \( n \) numbers? Since the algorithm is recursive, we prove this by induction.

The two base cases are when \( n = 0 \) or \( n = 1 \). It follows from the pseudocode for MergeSort\((L, n)\) that it simply returns the input list \( L \),
which is obviously sorted.

Let \( n \geq 2 \) and assume that for any integer \( k \) with \( 0 \leq k < n \) and for any list \( L' \) of \( k \) numbers, algorithm `MergeSort(L', k)` returns a list containing the elements of \( L' \) in sorted order. Let \( L \) be a list of \( n \) numbers. By going through the pseudocode for `MergeSort(L, n)`, we observe the following:

- The recursive call `MergeSort(L_1, m)` is on a list with less than \( n \) numbers. Therefore, by the induction hypothesis, its output, which is the list \( L_1 \), is sorted.

- The recursive call `MergeSort(L_2, n - m)` is on a list with less than \( n \) numbers. Again by the induction hypothesis, its output, which is the list \( L_2 \), is sorted.

- Algorithm `Merge(L_1, L_2)` gets as input the two sorted lists \( L_1 \) and \( L_2 \), and returns a list \( L \). Since algorithm `Merge` is correct, it then follows that \( L \) is a sorted list.

It follows that the final list \( L \), which is returned by algorithm `MergeSort`, is sorted.

This proves the correctness of algorithm `MergeSort(L, n)` for any integer \( n \geq 0 \) and any list \( L \) of \( n \) numbers.

### 4.5.2 Running Time of Algorithm `MergeSort`

We now analyze the running time of algorithm `MergeSort`. It follows from the pseudocode that, when running this algorithm together with its recursive calls, several calls are made to algorithm `Merge`. We are going to count the total number of comparisons that are made. That is, we will determine the total number of times that the line “if \( x \leq y \)” in algorithm `Merge` is executed when running algorithm `MergeSort(L, n)`.

We first observe that the number of comparisons made by algorithm `Merge(L_1, L_2)` is at most \( |L_1| + |L_2| \).

Let \( n \) be an integer and assume for simplicity that \( n \) is a power of two, i.e., \( n = 2^k \) for some integer \( k \geq 0 \). We define \( T(n) \) to be the maximum number of comparisons made when running algorithm `MergeSort(L, n)` on any input list \( L \) of \( n \) numbers. Note that we include in \( T(n) \) all comparisons that are made during all calls to `Merge` that are part of all recursive calls that are generated when running `MergeSort(L, n)`. 
Consider a list $L$ of $n$ numbers, where $n$ is a power of two. For $n = 1$, it follows from the pseudocode for $\text{MERGESORT}(L, n)$ that

$$T(1) = 0.$$ 

Assume that $n \geq 2$ and consider again the pseudocode for $\text{MERGESORT}(L, n)$. Which parts of the algorithm make comparisons between input elements?

- The call $\text{MERGESORT}(L_1, m)$ is a recursive call on a list of $m = n/2$ numbers. By definition, the total number of comparisons made in this call (together with all its recursive subcalls) is at most $T(n/2)$.

- The call $\text{MERGESORT}(L_2, n - m)$ is a recursive call on a list of $n - m = n/2$ numbers. By definition, the total number of comparisons made in this call (together with all its recursive subcalls) is at most $T(n/2)$.

- Finally, algorithm $\text{MERGESORT}(L, n)$ calls the non-recursive algorithm $\text{MERGE}(L_1, L_2)$. We have seen above that the number of comparisons made in this call is at most $|L_1| + |L_2| = n$.

By adding the number of comparisons, we get

$$T(n) \leq T(n/2) + T(n/2) + n = 2 \cdot T(n/2) + n.$$ 

Thus, we obtain the following recurrence:

$$

t(1) = 0,

T(n) \leq 2 \cdot T(n/2) + n, \text{ if } n \geq 2 \text{ and } n \text{ is a power of } 2. \quad (4.5)
$$

Our goal was to determine $T(n)$, but at this moment, we only have a recurrence for this function. We will solve the recurrence using a technique called *unfolding*:

Recall that we assume that $n = 2^k$ for some integer $k \geq 0$. Let $n$ be a large integer. We know from (4.5) that

$$T(n) \leq 2 \cdot T(n/2) + n.$$ 

If we replace $n$ by $n/2$ in (4.5), which is a valid thing to do, we get

$$T(n/2) \leq 2 \cdot T(n/2^2) + n/2.$$
By combining these two inequalities, we get

\[
T(n) \leq 2 \cdot T(n/2) + n \\
\leq 2 \left(2 \cdot T(n/2^2) + n/2\right) + n \\
= 2^2 \cdot T(n/2^2) + 2n.
\]

Let us repeat this: Replacing \(n\) by \(n/2^2\) in (4.5) gives

\[
T(n/2^2) \leq 2 \cdot T(n/2^3) + n/2^2.
\]

By substituting this into the inequality for \(T(n)\), we get

\[
T(n) \leq 2^2 \cdot T(n/2^2) + 2n \\
\leq 2^2 \left(2 \cdot T(n/2^3) + n/2^2\right) + 2n \\
= 2^3 \cdot T(n/2^3) + 3n.
\]

In the next step, we replace \(n\) by \(n/2^3\) in (4.5), which gives

\[
T(n/2^3) \leq 2 \cdot T(n/2^4) + n/2^3.
\]

By substituting this into the inequality for \(T(n)\), we get

\[
T(n) \leq 2^3 \cdot T(n/2^3) + 3n \\
\leq 2^3 \left(2 \cdot T(n/2^4) + n/2^3\right) + 3n \\
= 2^4 \cdot T(n/2^4) + 4n.
\]

At this moment, you will see the pattern and, at the end, we get the inequality

\[
T(n) \leq 2^k \cdot T(n/2^k) + kn.
\]

Since \(n = 2^k\), we have \(T(n/2^k) = T(1)\), which is 0 from the base case of the recurrence. Also, \(n = 2^k\) implies that \(k = \log n\). We conclude that

\[
T(n) \leq n \cdot T(1) + n \log n = n \log n.
\]

We thus have solved the recurrence. In case you have doubts about the validity of the unfolding method, we verify by induction that indeed

\[
T(n) \leq n \log n, \text{ for any integer } n \text{ that is a power of 2.}
\]
The base case is when $n = 1$. In this case, we have $T(1) = 0$ and $1 \log 1 = 1 \cdot 0 = 0$. Let $n \geq 2$ and assume that

$$T(n/2) \leq (n/2) \log(n/2).$$

From the recurrence, we get

$$T(n) \leq 2 \cdot T(n/2) + n.$$ 

By substituting the induction hypothesis into this inequality, we get

$$T(n) \leq 2 \cdot (n/2) \log(n/2) + n$$

$$= n \log(n/2) + n$$

$$= n \log n.$$

Thus, by induction, $T(n) \leq n \log n$ for any integer $n$ that is a power of 2.

Until now, we have only counted the number of comparisons made by algorithm MERGESORT. It follows from the pseudocode that the total running time, i.e., the total number of “elementary” steps, is within a constant factor of the total number of comparisons. Therefore, if $n$ is a power of 2, the running time of algorithm MERGESORT($L, n$) is $O(n \log n)$.

For general values of $n$, the recurrence for the number of comparisons becomes the following:

$$T(n) = 0, \text{ if } n = 0 \text{ or } n = 1,$$
$$T(n) \leq T(\lfloor n/2 \rfloor) + T(\lceil n/2 \rceil) + n, \text{ if } n \geq 2.$$ 

It can be shown by induction that this recurrence solves to $T(n) = O(n \log n)$.

We have proved the following result:

**Theorem 4.5.1** For any list $L$ of $n$ numbers, the running time of algorithm MERGESORT($L, n$) is $O(n \log n)$.

### 4.6 Counting Regions when Cutting a Circle

Take a circle, place $n$ points on it, and connect each pair of points by a straight-line segment. The points must be placed in such a way that no three segments pass through one point. These segments divide the circle into regions. Define $R_n$ to be the number of such regions. Can we determine $R_n$?
By looking at the figure above, we see that

\[ R_1 = 1, R_2 = 2, R_3 = 4, R_4 = 8, R_5 = 16. \]

There seems to be a clear pattern and it is natural to guess that \( R_n \) is equal to \( 2^{n-1} \) for all \( n \geq 1 \). To prove this, we have to argue that the number of regions doubles if we increase \( n \) by 1. If you try to do this, however, then you will fail! The reason is that \( R_n \) is not equal to \( 2^{n-1} \) for all \( n \geq 1 \); our guess was correct only for \( 1 \leq n \leq 5 \).

We will prove below that \( R_n \) grows only polynomially in \( n \). This will imply that \( R_n \) cannot be equal to \( 2^{n-1} \) for all \( n \), because the latter function grows exponentially.

### 4.6.1 A Polynomial Upper Bound on \( R_n \)

Let \( n \) be a (large) integer, consider a placement of \( n \) points on a circle, and connect each of the \( \binom{n}{2} \) pairs of points by a straight-line segment. Recall that we assume that no three segments pass through one point. We define the following graph:

- Each of the \( n \) points on the circle is a vertex.
- Each intersection point between two segments is a vertex.
- These vertices divide the segments into subsegments and the circle into arcs in a natural way. Each such subsegment and arc is an edge of the graph.

The figure below illustrates this for the case when \( n = 5 \). The graph on the right has 10 = 5 + 5 vertices: Each of the 5 points on the circle leads to one vertex and each of the 5 intersection points leads to one vertex. These 10 vertices divide the \( \binom{5}{2} = 10 \) segments into 20 straight-line edges and the circle into 5 circular edges. Therefore, the graph has 20 + 5 = 25 edges.
4.6. Counting Regions when Cutting a Circle

Note that, strictly speaking, this process does not define a proper graph, because any two consecutive vertices on the circle are connected by two edges (one straight-line edge and one circular edge), whereas in a proper graph, there can be only one edge between any pair of vertices. For simplicity, however, we will refer to the resulting structure as a graph.

Define $V_n$ and $E_n$ to be the number of vertices and edges of the graph, respectively. We claim that

$$V_n \leq n + \binom{n}{2}.$$  \hfill (4.6)

This claim follows from the following observations:

- There are exactly $n$ vertices on the circle.

- The $n$ points on the circle are connected by $\binom{n}{2}$ segments, and any two such segments intersect at most once. Therefore, the number of vertices inside the circle is at most the number of pairs of segments. The latter quantity is equal to

$$\binom{n}{2}.$$

We next claim that

$$E_n \leq n + \binom{V_n}{2}. \quad (4.7)$$

This claim follows from the following observations:

- There are exactly $n$ edges on the circle.

- Any straight-line edge joins two vertices. Therefore, the number of straight-line edges is at most the number of pairs of vertices, which is $\binom{V_n}{2}$. 
The final claim is that

\[ R_n \leq E_n. \]  

To prove this claim, we do the following. For each region \( r \), choose a point \( p_r \) inside \( r \), such that the \( y \)-coordinate of \( p_r \) is not equal to the \( y \)-coordinate of any vertex. Define \( f(r) \) to be the edge that is reached when walking from \( p_r \) horizontally to the right.

This defines a one-to-one function \( f \) from the set of regions to the set of edges. Therefore, the number of regions, which is \( R_n \), is at most the number of edges, which is \( E_n \).

By combining (4.6), (4.7), and (4.8), we get

\[
R_n \leq E_n \\
\leq n + \binom{V_n}{2} \\
\leq n + \left( n + \frac{\binom{n}{2}}{2} \right).
\]

In order to estimate the last quantity, we are going to use asymptotic notation; see Section 2.3. First observe that

\[
\binom{n}{2} = \frac{n(n-1)}{2} = O(n^2).
\]

This implies that

\[
\binom{\binom{n}{2}}{2} = \frac{O(n^2)}{2} = O(n^4),
\]
which implies that

\[ n + \left( \binom{n}{2} \right) = n + O(n^4) = O(n^4), \]

which implies that

\[ \left( n + \left( \binom{n}{2} \right) \right) = \left( O(n^4) \right) = O(n^8), \]

which implies that

\[ R_n \leq n + \left( n + \binom{n}{2} \right) = n + O(n^8) = O(n^8). \]

Thus, we have proved our claim that \( R_n \) grows polynomially in \( n \) and, therefore, for large values of \( n \), \( R_n \) is not equal to \( 2^{n-1} \). (Using results that we will see in Section 7.4.1, it can be shown that, in fact, \( R_n = O(n^4) \).)

We remark that there is a shorter way to prove that \( R_n \) is not equal to \( 2^{n-1} \) for all \( n \geq 1 \): You can verify by hand that \( R_6 = 31 \). Still, this single example does not rule out the possibility that \( R_n \) grows exponentially. The analysis that we gave above does rule this out.

We have proved above that \( R_n = O(n^8) \). We also mentioned that this upper bound can be improved to \( O(n^4) \). In the following subsections, we will prove that the latter upper bound cannot be improved. That is, we will prove that \( R_n = \Theta(n^4) \). In fact, we will determine an exact formula, in terms of \( n \), for the value of \( R_n \).

### 4.6.2 A Recurrence Relation for \( R_n \)

Let \( n \geq 2 \) be an integer. Consider a placement of \( n \) points on a circle. We denote these points by \( p_1, p_2, \ldots, p_n \) and assume that they are numbered in counterclockwise order. As before, we connect each of the \( \binom{n}{2} \) pairs of points by a straight-line segment. We assume that no three segments pass through one point. We are going to derive a recurrence relation for the number \( R_n \) of regions in the following way:

- Remove all segments that have \( p_n \) as an endpoint. At this moment, the number of regions is, by definition, equal to \( R_{n-1} \).
Add the \( n - 1 \) line segments \( p_1p_n, p_2p_n, \ldots, p_{n-1}p_n \) one by one. For each segment \( p_kp_n \) added, determine the increase \( I_k \) in the number of regions.

- Take the sum of \( R_{n-1} \) and all increases \( I_k \), i.e.,

\[
R_{n-1} + \sum_{k=1}^{n-1} I_k.
\]

This sum is equal to \( R_n \), because in the entire process, we have counted each of the regions for \( n \) points exactly once.

- Thus, together with the base case \( R_1 = 1 \), we obtain a recurrence for the values \( R_n \).

We start by illustrating this process for the case when \( n = 6 \). The figure below shows the situation after we have removed all segments that have \( p_6 \) as an endpoint. The number of regions is equal to \( R_5 = 16 \).

We are going to add, one by one, the five segments that have \( p_6 \) as an endpoint. When we add \( p_1p_6 \), one region gets cut into two. Thus, the number of regions increases by one. Using the notation introduced above, we have \( I_1 = 1 \).

When we add \( p_2p_6 \), four regions get cut into two. Thus, the number of regions increases by four, and we have \( I_2 = 4 \).
4.6. Counting Regions when Cutting a Circle

When we add $p_3p_6$, five regions get cut into two. Thus, the number of regions increases by five, and we have $I_3 = 5$.

When we add $p_4p_6$, four regions get cut into two. Thus, the number of regions increases by four, and we have $I_4 = 4$.

Finally, when we add $p_5p_6$, one region gets cut into two. Thus, the number of regions increases by one, and we have $I_5 = 1$.

After having added the five segments with endpoint $p_6$, we have accounted for all regions determined by the six points. In other words, the number of regions we have at the end is equal to $R_6$. Since the number of regions at
the end is also equal to the sum of (i) the number of regions we started with, which is $R_5$, and (ii) the total increase, we have

$$R_6 = R_5 + I_1 + I_2 + I_3 + I_4 + I_5 = 31.$$  

Let us look at this more carefully. We have seen that $I_3 = 5$. That is, when adding the segment $p_3p_6$, the number of regions increases by 5. Where does this number 5 come from? The segment $p_3p_6$ intersects 4 segments, namely $p_1p_4$, $p_1p_5$, $p_2p_4$, and $p_2p_5$. The increase in the number of regions is one more than the number of intersections. Thus, when adding a segment, if we determine the number $X$ of intersections between this new segment and existing segments, then the increase in the number of regions is equal to $1 + X$.

When we add $p_3p_6$, we have $X = 4$. Where does this number 4 come from? We make the following observations:

- Any segment that intersects $p_3p_6$ has one endpoint above $p_3p_6$ and one endpoint below $p_3p_6$.

- Any pair $(a, b)$ of points, with $a$ above $p_3p_6$ and $b$ below $p_3p_6$, defines a segment $ab$ that intersects $p_3p_6$.

- Thus, the value of $X$ is equal to the number of pairs $(a, b)$ of points in \{p_1, p_2, p_4, p_5\}, where $a$ is above $p_3p_6$ and $b$ is below $p_3p_6$. Since there are 2 choices for $a$ (viz., $p_1$ and $p_2$) and 2 choices for $b$ (viz., $p_4$ and $p_5$), it follows from the Product Rule that $X = 2 \cdot 2 = 4$.

Now that we have seen the basic approach, we are going to derive the recurrence for $R_n$ for an arbitrary integer $n \geq 2$. After having removed all segments that have $p_n$ as an endpoint, we have $R_{n-1}$ regions. For each integer $k$ with $1 \leq k \leq n - 1$, we add the segment $p_kp_n$. What is the number of existing segments that are intersected by this new segment?
4.6. Counting Regions when Cutting a Circle

We observe that for $i < j$,

- $p_i p_j$ intersects $p_k p_n$ if and only if $1 \leq i \leq k - 1$ and $k + 1 \leq j \leq n - 1$.

Since there are $k - 1$ choices for $i$ and $n - k - 1$ choices for $j$, the Product Rule implies that the number of intersections due to $p_k p_n$ is equal to $(k - 1)(n - k - 1)$. Thus, the segment $p_k p_n$ goes through $1 + (k - 1)(n - k - 1)$ regions, and each of them is cut into two. It follows that, when adding $p_k p_n$, the increase $I_k$ in the number of regions is equal to

$$I_k = 1 + (k - 1)(n - k - 1).$$

We conclude that

$$R_n = R_{n-1} + \sum_{k=1}^{n-1} I_k = R_{n-1} + \sum_{k=1}^{n-1} (1 + (k - 1)(n - k - 1)).$$

In the summation on the right-hand side

- the term 1 occurs exactly $n - 1$ times, and

- the term $(k - 1)(n - k - 1)$ is non-zero only if $2 \leq k \leq n - 2$.

It follows that, for $n \geq 2$,

$$R_n = R_{n-1} + (n - 1) + \sum_{k=2}^{n-2} (k - 1)(n - k - 1). \quad (4.9)$$

Thus, together with the base case

$$R_1 = 1,$$  \quad (4.10)

we have determined the recurrence we were looking for.
4.6.3 Simplifying the Recurrence Relation

In this subsection, we will use a combinatorial proof (see Section 3.7) to show that the summation on the right-hand side of (4.9) satisfies

$$\sum_{k=2}^{n-2} (k-1)(n-k-1) = \binom{n-1}{3},$$  \hspace{1cm} (4.11)

Consider the set $S = \{1, 2, \ldots, n-1\}$. We know that the number of 3-element subsets of $S$ is equal to $\binom{n-1}{3}$. As we will see below, the summation on the left-hand side of (4.11) counts exactly the same subsets.

We divide the 3-element subsets of $S$ into groups based on their middle element. Observe that the middle element can be any of the values 2, 3, \ldots, $n-2$. Thus, for any $k$ with $2 \leq k \leq n-2$, the $k$-th group $G_k$ consists of all 3-element subsets of $S$ whose middle element is equal to $k$. Since the groups are pairwise disjoint, we have

$$\binom{n-1}{3} = \sum_{k=2}^{n-2} |G_k|.$$

What is the size of the $k$-th group $G_k$? Any 3-element subset in $G_k$ consists of

- one element from $\{1, 2, \ldots, k-1\}$,
- the element $k$, and
- one element from $\{k+1, k+2, \ldots, n-1\}$.

It then follows from the Product Rule that

$$|G_k| = (k-1) \cdot 1 \cdot (n-k-1) = (k-1)(n-k-1).$$

Thus, we have proved the identity in (4.11), and the recurrence relation in (4.9) and (4.10) becomes

$$R_1 = 1,$$

$$R_n = R_{n-1} + (n-1) + \binom{n-1}{3}, \text{ if } n \geq 2. \hspace{1cm} (4.12)$$
4.6.4 Solving the Recurrence Relation

Now that we have a recurrence that looks reasonable, we are going to apply the unfolding technique of Section 4.5 to solve it: By repeatedly applying the recurrence in (4.12), we get

\[
R_n = (n - 1) + \binom{n - 1}{3} + R_{n-1}
\]

\[
= (n - 1) + (n - 2) + \binom{n - 1}{3} + \binom{n - 2}{3} + R_{n-2}
\]

\[
= (n - 1) + (n - 2) + (n - 3) + \binom{n - 1}{3} + \binom{n - 2}{3} + \binom{n - 3}{3} + R_{n-3}.
\]

By continuing, we get

\[
R_n = (n - 1) + (n - 2) + (n - 3) + \cdots + 3 + 2 + 1 + \binom{n - 1}{3} + \binom{n - 2}{3} + \binom{n - 3}{3} + \cdots + \binom{3}{3} + \binom{2}{3} + \binom{1}{3} + R_1.
\]

Since \( \binom{2}{3} = \binom{1}{3} = 0 \) and \( R_1 = 1 \), we get

\[
R_n = (n - 1) + (n - 2) + (n - 3) + \cdots + 3 + 2 + 1 + \binom{n - 1}{3} + \binom{n - 2}{3} + \binom{n - 3}{3} + \cdots + \binom{3}{3} + \binom{2}{3} + \binom{1}{3} + 1.
\]

Since, by Theorem 2.2.10, the first summation is equal to

\[
1 + 2 + 3 + \cdots + (n - 1) = n(n - 1)/2 = \binom{n}{2},
\]

we get

\[
R_n = 1 + \binom{n}{2} + \sum_{k=3}^{n-1} \binom{k}{3}.
\]

The final step is to simplify the summation on the right-hand side. We will use a combinatorial proof to show that

\[
\sum_{k=3}^{n-1} \binom{k}{3} = \binom{n}{4}.
\]

(4.13)
Consider all 4-element subsets of the set $S = \{1, 2, \ldots, n\}$. We know that there are $\binom{n}{4}$ many such subsets. We divide these subsets into groups based on their largest element. For any $k$ with $3 \leq k \leq n - 1$, the $k$-th group $G_k$ consists of all 4-element subsets of $S$ whose largest element is equal to $k + 1$. It should be clear that

$$\binom{n}{4} = \sum_{k=3}^{n-1} |G_k|.$$

To determine the size of the group $G_k$, we observe that any 4-element subset in $G_k$ consists of

- three elements from $\{1, 2, \ldots, k\}$ and
- the element $k + 1$.

It then follows from the Product Rule that

$$|G_k| = \binom{k}{3} \cdot 1 = \binom{k}{3},$$

completing the proof of (4.13).

After (finally!) having solved and simplified our recurrence, we conclude that for any integer $n \geq 1$,

$$R_n = 1 + \binom{n}{2} + \binom{n}{4}.$$

In Exercise 4.43 you will see a shorter way to determine the exact value of $R_n$. We went for the long derivation, because it allowed us to illustrate, along the way, several techniques from previous sections.

### 4.7 Exercises

**4.1** The function $f : \mathbb{N} \to \mathbb{Z}$ is recursively defined as follows:

\[
\begin{align*}
    f(0) &= 7, \\
    f(n) &= f(n-1) + 6n - 3 \text{ if } n \geq 1.
\end{align*}
\]

Prove that $f(n) = 3n^2 + 7$ for all integers $n \geq 0$. 

4.2 The function \( f : \mathbb{N} \rightarrow \mathbb{Z} \) is recursively defined as follows:

\[
\begin{align*}
  f(0) &= -18, \\
  f(n) &= 9(n-2)(n-3) + f(n-1) \quad \text{if } n \geq 1.
\end{align*}
\]

Prove that

\[ f(n) = 3(n-1)(n-2)(n-3) \]

for all \( n \geq 0 \).

4.3 The function \( f : \mathbb{N} \rightarrow \mathbb{Z} \) is recursively defined as follows:

\[
\begin{align*}
  f(0) &= 3, \\
  f(n) &= 2 \cdot f(n-1) - (f(n-1))^2 \quad \text{if } n \geq 1.
\end{align*}
\]

Prove that \( f(n) = 1 - 2^{2^n} \) for all integers \( n \geq 1 \). (\( 2^{2^n} \) means 2 to the power of \( 2^n \).)

4.4 The function \( f : \mathbb{N} \rightarrow \mathbb{Z} \) is defined by

\[ f(n) = 2n(n-6) \]

for each integer \( n \geq 0 \). Derive a recursive form of this function.

4.5 The function \( f : \mathbb{N}^3 \rightarrow \mathbb{N} \) is defined as follows:

\[
\begin{align*}
  f(k, n, 0) &= k + n \quad \text{if } k \geq 0 \text{ and } n \geq 0, \\
  f(k, 0, 1) &= 0 \quad \text{if } k \geq 0, \\
  f(k, 0, 2) &= 1 \quad \text{if } k \geq 0, \\
  f(k, 0, i) &= k \quad \text{if } k \geq 0 \text{ and } i \geq 3, \\
  f(k, n, i) &= f(k, f(k, n-1, i), i-1) \quad \text{if } k \geq 0, i \geq 1, \text{ and } n \geq 1.
\end{align*}
\]

Determine \( f(2, 3, 2) \).

4.6 The functions \( f : \mathbb{N} \rightarrow \mathbb{N} \) and \( g : \mathbb{N}^2 \rightarrow \mathbb{N} \) are recursively defined as follows:

\[
\begin{align*}
  f(0) &= 1, \\
  f(n) &= g(n, f(n-1)) \quad \text{if } n \geq 1, \\
  g(m, 0) &= 0 \quad \text{if } m \geq 0, \\
  g(m, n) &= m + g(m, n-1) \quad \text{if } m \geq 0 \text{ and } n \geq 1.
\end{align*}
\]

Solve these recurrences for \( f \), i.e., express \( f(n) \) in terms of \( n \).
4.7 The functions \( f : \mathbb{N} \to \mathbb{N} \) and \( g : \mathbb{N}^2 \to \mathbb{N} \) are recursively defined as follows:

\[
\begin{align*}
  f(0) &= 1, \\
  f(1) &= 2, \\
  f(n) &= g(f(n-2), f(n-1)) \quad \text{if } n \geq 2, \\
  g(m,0) &= 2m \quad \text{if } m \geq 0, \\
  g(m,n) &= g(m,n-1) + 1 \quad \text{if } m \geq 0 \text{ and } n \geq 1.
\end{align*}
\]

Solve these recurrences for \( f \), i.e., express \( f(n) \) in terms of \( n \).

4.8 In Section 4.2, we have defined the Fibonacci numbers \( f_0, f_1, f_2, \ldots \) Prove that for each integer \( n \geq 1 \),

\[
\sum_{i=1}^{n} f_{2i} = f_{2n+1} - 1
\]

and

\[
f_1^2 + f_2^2 + f_3^2 + \cdots + f_n^2 = f_nf_{n+1}.
\]

4.9 In Section 4.2, we have defined the Fibonacci numbers \( f_0, f_1, f_2, \ldots \) Prove that for each integer \( n \geq 0 \),

- \( f_{3n} \) is even,
- \( f_{3n+1} \) is odd,
- \( f_{3n+2} \) is odd,
- \( f_{4n} \) is a multiple of 3.

4.10 Let \( \varphi = \frac{1+\sqrt{5}}{2} \) and \( \psi = \frac{1-\sqrt{5}}{2} \), and let \( n \geq 0 \) be an integer. We have seen in Theorem 4.2.1 that

\[
\frac{\varphi^n - \psi^n}{\sqrt{5}} \tag{4.14}
\]

is equal to the \( n \)-th Fibonacci number \( f_n \). Since the Fibonacci numbers are obviously integers, the number in (4.14) is an integer as well.

Prove that the number in (4.14) is a rational number using only Newton’s Binomial Theorem.
4.11 Let \( n \geq 1 \) be an integer and consider a \( 2 \times n \) board \( B_n \) consisting of \( 2n \) square cells. The top part of the figure below shows \( B_{13} \).

![Figure showing \( B_{13} \)](image)

A *brick* is a horizontal or vertical board consisting of 2 square cells; see the bottom part of the figure above. A *tiling* of the board \( B_n \) is a placement of bricks on the board such that

- the bricks exactly cover \( B_n \) and
- no two bricks overlap.

The figure below shows a tiling of \( B_{13} \).

![Figure showing a tiling of \( B_{13} \)](image)

For \( n \geq 1 \), let \( T_n \) be the number of different tilings of the board \( B_n \). Determine the value of \( T_n \), i.e., express \( T_n \) in terms of numbers that we have seen in this chapter.

4.12 Let \( n \) be a positive integer and consider a \( 5 \times n \) board \( B_n \) consisting of \( 5n \) cells, each one having sides of length one. The top part of the figure below shows \( B_{12} \).

![Figure showing \( B_{12} \)](image)
A brick is a horizontal or vertical board consisting of $2 \times 3 = 6$ cells; see the bottom part of the figure above. A tiling of the board $B_n$ is a placement of bricks on the board such that

- the bricks exactly cover $B_n$ and
- no two bricks overlap.

The figure below shows a tiling of $B_{12}$.

Let $T_n$ be the number of different tilings of the board $B_n$.

- Let $n \geq 6$ be a multiple of 6. Determine the value of $T_n$.
- Let $n$ be a positive integer that is not a multiple of 6. Prove that $T_n = 0$.

4.13 Let $n$ be a positive integer and consider a $1 \times n$ board $B_n$ consisting of $n$ cells, each one having sides of length one. The top part of the figure below shows $B_9$.

You have an unlimited supply of bricks, which are of the following types (see the bottom part of the figure above):

- There are red ($R$) and blue ($B$) bricks, both of which are $1 \times 1$ cells.
- There are white ($W$), yellow ($Y$), and green ($G$) bricks, all of which are $1 \times 2$ cells.

A tiling of the board $B_n$ is a placement of bricks on the board such that
4.7. Exercises

- the bricks exactly cover \( B_n \) and
- no two bricks overlap.

In a tiling, a color can be used more than once and some colors may not be used at all. The figure below shows a tiling of \( B_9 \), in which each color is used and the color red is used twice.

\[
\begin{array}{ccccccc}
B & W & R & G & R & Y \\
\end{array}
\]

Let \( T_n \) be the number of different tilings of the board \( B_n \).

- Determine \( T_1 \) and \( T_2 \).
- Let \( n \geq 3 \) be an integer. Prove that
  \[
  T_n = 2 \cdot T_{n-1} + 3 \cdot T_{n-2}.
  \]
- Prove that for any integer \( n \),
  \[
  2(-1)^{n-1} + 3(-1)^{n-2} = (-1)^n.
  \]
- Prove that for any integer \( n \geq 1 \),
  \[
  T_n = \frac{3^{n+1} + (-1)^n}{4}.
  \]

4.14 Consider strings of \( n \) characters, each character being \( a, b, c, \) or \( d \), that contain an even number of as. (Recall that 0 is even.) Let \( E_n \) be the number of such strings. Prove that for any integer \( n \geq 1 \),

\[
E_{n+1} = 2 \cdot E_n + 4^n.
\]

4.15 Let \( A_n \) be the number of bitstrings of length \( n \) that contain 000. Prove that for \( n \geq 4 \),

\[
A_n = A_{n-1} + A_{n-2} + A_{n-3} + 2^{n-3}.
\]

4.16 Let \( n \geq 1 \) be an integer and define \( A_n \) to be the number of bitstrings of length \( n \) that do not contain 101.

- Determine \( A_1, A_2, A_3, \) and \( A_4 \).
• Prove that for each integer \( n \geq 4 \),
\[
A_n = 3 + A_1 + A_2 + A_3 + \cdots + A_{n-4} + A_{n-3} + A_{n-1}
\]
\[
= 3 + \sum_{k=1}^{n-3} A_k + A_{n-1}.
\]

*Hint*: Divide the strings into groups depending on the number of leading 1s.

4.17 Let \( n \geq 1 \) be an integer and consider \( n \) people \( P_1, P_2, \ldots, P_n \). Let \( A_n \) be the number of ways these \( n \) people can be divided into groups, such that each group consists of either one or two people.

• Determine \( A_1, A_2, \) and \( A_3 \).

• Prove that for each integer \( n \geq 3 \),
\[
A_n = A_{n-1} + (n - 1) \cdot A_{n-2}.
\]

4.18 A *maximal run of ones* in a bitstring is a maximal consecutive substring of ones. For example, the bitstring 110011110100111 has four maximal runs of ones: 11, 1111, 1, and 111. These maximal runs have lengths 2, 4, 1, and 3, respectively.

Let \( n \geq 1 \) be an integer and let \( B_n \) be the number of bitstrings of length \( n \) that do not contain any maximal run of ones of odd length; in other words, every maximal run of ones in these bitstrings has an even length.

• Determine \( B_1, B_2, \) and \( B_3 \).

• Determine the value of \( B_n \), i.e., express \( B_n \) in terms of numbers that we have seen in this chapter.

4.19 Let \( n \geq 1 \) be an integer and let \( S_n \) be the number of ways in which \( n \) can be written as a sum of 1s and 2s; the order in which the 1s and 2s occur in the sum matters. For example, \( S_3 = 3 \), because
\[
3 = 1 + 1 + 1 = 1 + 2 = 2 + 1.
\]

• Determine \( S_1, S_2, \) and \( S_4 \).

• Determine the value of \( S_n \), i.e., express \( S_n \) in terms of numbers that we have seen in this chapter.
4.20 Let \( n \geq 1 \) be an integer and consider the set \( S_n = \{1, 2, \ldots, n\} \). A non-neighbor subset of \( S_n \) is any subset \( T \) of \( S \) having the following property: If \( k \) is any element of \( T \), then \( k + 1 \) is not an element of \( T \). (Observe that the empty set is a non-neighbor subset of \( S_n \).)

For example, if \( n = 3 \), then \( \{1, 3\} \) is a non-neighbor subset, whereas \( \{2, 3\} \) is not a non-neighbor subset.

Let \( N_n \) denote the number of non-neighbor subsets of the set \( S_n \).

- Determine \( N_1 \), \( N_2 \), and \( N_3 \).
- Determine the value of \( N_n \), i.e., express \( N_n \) in terms of numbers that we have seen in class.

4.21 Let \( S \) be the set of ordered pairs of integers that is recursively defined in the following way:

- \((0, 0) \in S\).
- If \((a, b) \in S\) then \((a + 2, b + 3) \in S\).
- If \((a, b) \in S\) then \((a + 3, b + 2) \in S\).

Prove that for every element \((a, b)\) in \( S \), \( a + b \) is divisible by 5.

4.22 Let \( S \) be the set of integers that is recursively defined in the following way:

- 4 is an element of \( S \).
- If \( x \) and \( y \) are elements of \( S \), then \( x + y^2 \) is an element of \( S \).

Prove that every element of \( S \) is divisible by 4.

4.23 Let \( S \) be the set of ordered triples of integers that is recursively defined in the following way:

- \((66, 55, 1331) \in S\).
- If \((a, b, c) \in S\) then \((a + 7, b + 5, 14a - 10b + c + 24) \in S\).

Prove that for every element \((a, b, c)\) in \( S \),

\[
a^2 - b^2 = c.
\]
4.24 Let $S$ be the set of integers that is recursively defined in the following way:

- 1 is an element of $S$.
- If $x$ is an element of $S$, then $x + 2\sqrt{x} + 1$ is also an element of $S$.

Give a simple description of the set $S$ and prove that your answer is correct.

4.25 The set $S$ of binary strings is recursively defined in the following way:

- The string 00 is an element of the set $S$.
- The string 01 is an element of the set $S$.
- The string 10 is an element of the set $S$.
- If the string $s$ is an element of the set $S$, then the string 0$s$ (i.e., the string obtained by adding the bit 0 at the front of $s$) is also an element of the set $S$.
- If the string $s$ is an element of the set $S$, then the string 10$s$ (i.e., the string obtained by adding the bits 10 at the front of $s$) is also an element of the set $S$.

Let $s$ be an arbitrary string in the set $S$. Prove that $s$ does not contain the substring 11.

4.26 A binary tree is

- either one single node
- or a node whose left subtree is a binary tree and whose right subtree is a binary tree.
4.7. Exercises  

Prove that any binary tree with \( n \) leaves has exactly \( 2n - 1 \) nodes.

4.27 In Section 4.4, we have seen the recursive algorithm GOSSIP(\( n \)), which computes a sequence of phone calls for the persons \( P_1, P_2, \ldots, P_n \). The base case for this algorithm was when \( n = 4 \). Assume we change the base case to \( n = 2 \): In this new base case, there are only two people \( P_1 \) and \( P_2 \), and only one phone call is needed. The rest of the algorithm remains unchanged.

Prove that the modified algorithm GOSSIP(\( n \)) results in a sequence of \( 2n - 3 \) phone calls for any integer \( n \geq 2 \). (Thus, for \( n \geq 4 \), it makes one more phone call than the algorithm in Section 4.4.)

4.28 In Section 4.4, we have seen the recursive algorithm GOSSIP(\( n \)), which computes a sequence of phone calls for the persons \( P_1, P_2, \ldots, P_n \), for any integer \( n \geq 4 \).

Give an iterative (i.e., non-recursive) version of this algorithm in pseudocode. Your algorithm must produce exactly the same sequence of phone calls as algorithm GOSSIP(\( n \)).

4.29 The following recursive algorithm FIB takes as input an integer \( n \geq 0 \) and returns the \( n \)-th Fibonacci number \( f_n \):

```
Algorithm FIB(n):
    if \( n = 0 \) or \( n = 1 \)
    then \( f = n \)
    else \( f = FIB(n - 1) + FIB(n - 2) \)
    endif;
    return \( f \)
```

Let \( a_n \) be the number of additions made by algorithm FIB(\( n \)), i.e., the total number of times the \( + \)-function in the else-case is called. Prove that for all \( n \geq 0 \),

\[
a_n = f_{n+1} - 1.
\]

4.30 Consider the following recursive algorithm BEER(\( n \)), which takes as input an integer \( n \geq 1 \):
Algorithm Beer(n):

    if \( n = 1 \)
    then eat some peanuts
    else choose an arbitrary integer \( m \) with \( 1 \leq m \leq n - 1 \);
        Beer(m);
        drink one pint of beer;
        Beer(n - m)
    endif

- Explain why, for any integer \( n \geq 1 \), algorithm Beer(n) terminates.

- Let \( B(n) \) be the number of pints of beer you drink when running algorithm Beer(n). Determine the value of \( B(n) \).

4.31 Consider the following recursive algorithm Silly, which takes as input an integer \( n \geq 1 \) which is a power of 2:

Algorithm Silly(n):

    if \( n = 1 \)
    then drink one pint of beer
    else if \( n = 2 \)
    then fart once
    else fart once;
        Silly(n/2);
        fart once
    endif
    endif

For \( n \) a power of 2, let \( F(n) \) be the number of times you fart when running algorithm Silly(n). Determine the value of \( F(n) \).

4.32 Jennifer loves to drink India Pale Ale (IPA). After a week of hard work, Jennifer goes to the pub and runs the following recursive algorithm, which takes as input an integer \( n \geq 1 \), which is a power of 4:
Algorithm $\text{JENNIFERDRINKSIPA}(n)$:

if $n = 1$
then place one order of chicken wings
else for $k = 1$ to 4
  do $\text{JENNIFERDRINKSIPA}(n/4)$;
     drink $n$ pints of IPA
endfor
endif

For $n$ a power of 4, let

- $P(n)$ be the number of pints of IPA that Jennifer drinks when running algorithm $\text{JENNIFERDRINKSIPA}(n)$,
- $C(n)$ be the number of orders of chicken wings that Jennifer places when running algorithm $\text{JENNIFERDRINKSIPA}(n)$.

Determine the values of $P(n)$ and $C(n)$.

4.33 Let $n \geq 2$ be an integer and consider a sequence $s_1, s_2, \ldots, s_n$ of $n$ pairwise distinct numbers. The following algorithm computes the smallest and largest elements in this sequence:

Algorithm $\text{MINMAX}(s_1, s_2, \ldots, s_n)$:

$min = s_1$;
$max = s_1$;
for $i = 2$ to $n$
do if $s_i < min$ (1)
  then $min = s_i$
endif;
if $s_i > max$ (2)
  then $max = s_i$
endif
endwhile;
return $(min, max)$

This algorithm makes comparisons between input elements in lines (1) and (2). Determine the total number of comparisons as a function of $n$. 
4.34 Let \( n \geq 2 \) be a power of 2 and consider a sequence \( S \) of \( n \) pairwise distinct numbers. The following algorithm computes the smallest and largest elements in this sequence:

**Algorithm** \texttt{FastMinMax}(\( S, n \)):

\[
\text{if } n = 2 \\
\text{then let } x \text{ and } y \text{ be the two elements in } S; \\
\text{if } x < y \quad (1) \\
\quad \text{then } \text{min} = x; \\
\quad \text{max} = y \\
\quad \text{else } \text{min} = y; \\
\quad \text{max} = x \\
\text{endif} \\
\text{else} \\
\text{divide } S \text{ into two subsequences } S_1 \text{ and } S_2, \text{ both of size } n/2; \\
\quad (\text{min}_1, \text{max}_1) = \text{FastMinMax}(S_1, n/2); \\
\quad (\text{min}_2, \text{max}_2) = \text{FastMinMax}(S_2, n/2); \\
\text{if } \text{min}_1 < \text{min}_2 \quad (2) \\
\text{then } \text{min} = \text{min}_1 \\
\text{else } \text{min} = \text{min}_2 \\
\text{endif;} \\
\text{if } \text{max}_1 < \text{max}_2 \quad (3) \\
\text{then } \text{max} = \text{max}_2 \\
\text{else } \text{max} = \text{max}_1 \\
\text{endif} \\
\text{endif;} \\
\text{return } (\text{min}, \text{max})
\]

This algorithm makes comparisons between input elements in lines (1), (2), and (3). Let \( C(n) \) be the total number of comparisons made by algorithm \texttt{FastMinMax} on an input sequence of length \( n \).

- Derive a recurrence for \( C(n) \).
- Use this recurrence to prove that \( C(n) = \frac{3}{2} n - 2 \) for each \( n \geq 2 \) that is a power of 2.

4.35 Let \( k \) be a positive integer and let \( n = 2^k \). You are given an \( n \times n \) board
A _tromino_ is an L-shaped object consisting of three $1 \times 1$ cells. Each tromino can appear in four different orientations; see the right part of the figure below.

A _tiling_ of the board $B_n$ is a placement of trominoes on the board such that

- the trominoes cover exactly all white cells (thus, the black cell is not covered by any tromino) and

- no two trominoes overlap.

Here is a tiling of the board given above:

Describe a recursive algorithm that

- takes as input a board $B_n$ having exactly one black cell (which can be anywhere on the board) and

- returns a tiling of this board.

*Hint:* Look at the following figure:
4.36 Let $S$ be a set of $n$ points in the plane. Each point $p$ of $S$ is given by its $x$- and $y$-coordinates $p_x$ and $p_y$, respectively. We assume that no two points of $S$ have the same $x$-coordinate and no two points of $S$ have the same $y$-coordinate.

A point $p$ of $S$ is called maximal in $S$ if there is no point in $S$ that is to the north-east of $p$, i.e.,

$$\{q \in S : q_x > p_x \text{ and } q_y > p_y\} = \emptyset.$$ 

The figure below shows an example, in which the $\bullet$-points are maximal and the $\times$-points are not maximal. Observe that, in general, there is more than one maximal element in $S$.

Describe a recursive algorithm $\text{MAXELEM}(S)$ that has the same structure as algorithm $\text{MERGESORT}$ in Section 4.5 and does the following:

**Input:** A set $S$ of $n$ points in the plane, in sorted order from left to right.

**Output:** All maximal elements of $S$, in sorted order from left to right.

The running time $T(n)$ of your algorithm must be $O(n \log n)$.

4.37 The Hadamard matrices $H_0, H_1, H_2, \ldots$ are recursively defined as follows:

$$H_0 = (1)$$

and for $k \geq 1$,

$$H_k = \begin{pmatrix} H_{k-1} & H_{k-1} \\ H_{k-1} & -H_{k-1} \end{pmatrix}.$$ 

Thus, $H_0$ is a $1 \times 1$ matrix whose only entry is 1,

$$H_1 = \begin{pmatrix} 1 & 1 \\ 1 & -1 \end{pmatrix},$$
and

\[
H_2 = \begin{pmatrix}
1 & 1 & 1 & 1 \\
1 & -1 & 1 & -1 \\
1 & 1 & -1 & -1 \\
1 & -1 & -1 & 1 \\
\end{pmatrix}.
\]

Observe that \(H_k\) has \(2^k\) rows and \(2^k\) columns.

If \(x\) is a column vector of length \(2^k\), then \(H_kx\) is the column vector of length \(2^k\) obtained by multiplying the matrix \(H_k\) with the vector \(x\).

Describe a recursive algorithm \(\text{Mult}(k, x)\) that does the following:

**Input:** An integer \(k \geq 0\) and a column vector \(x\) of length \(n = 2^k\).

**Output:** The column vector \(H_kx\) (having length \(n\)).

The running time \(T(n)\) of your algorithm must be \(O(n \log n)\).

**Hint:** The input only consists of \(k\) and \(x\). The matrix \(H_k\), which has \(n^2\) entries, is not given as part of the input. Since you are aiming for an \(O(n \log n)\)-time algorithm, you cannot compute all entries of the matrix \(H_k\).

**4.38** Prove, for example by induction, that for \(n \geq 1\),

\[
1 + 2 + 3 + \cdots + n = \frac{n(n + 1)}{2},
\]

and

\[
1^2 + 2^2 + 3^2 + \cdots + n^2 = \frac{n(n + 1)(2n + 1)}{6}.
\]

**4.39** Assume you remember that

\[
1^2 + 2^2 + 3^2 + \cdots + n^2
\]

is equal to a polynomial of degree three, i.e.,

\[
1^2 + 2^2 + 3^2 + \cdots + n^2 = An^3 + Bn^2 + Cn + D,
\]

but you have forgotten the values of \(A, B, C,\) and \(D\). How can you determine these four values?

**4.40** In Section 4.6.3, we have shown that

\[
\sum_{k=2}^{n-2} (k-1)(n-k-1) = \binom{n-1}{3}.
\]

Use Exercise 4.38 to give an alternative proof.
4.41 In Section 4.6.4, we have used the fact that
\[ \sum_{k=1}^{n-1} k = \binom{n}{2}, \]
which follows from Theorem 2.2.10. Give an alternative proof that uses the approach that we used to prove the identity in (4.13).

4.42 In Section 4.6.4, we have shown that
\[ \sum_{k=3}^{n-1} \binom{k}{3} = \binom{n}{4}. \]
Use induction and Pascal’s Identity (see Theorem 3.7.2) to give an alternative proof.

4.43 Consider the numbers \( R_n \) that we defined in Section 4.6. The \( n \) points on the circle define \( \binom{n}{2} \) line segments, one segment for each pair of points. Let \( X \) be the total number of intersections among these \( \binom{n}{2} \) line segments.

- Prove that
  \[ R_n = 1 + \binom{n}{2} + X. \]
  
  *Hint:* Start with only the circle and the \( n \) points. Then add the \( \binom{n}{2} \) line segments one by one.

- Prove that
  \[ X = \binom{n}{4}. \]

4.44 For an integer \( n \geq 1 \), draw \( n \) straight lines, such that no two of them are parallel and no three of them intersect in one single point. These lines divide the plane into regions (some of which are bounded and some of which are unbounded). Denote the number of these regions by \( C_n \).

Derive a recurrence for the numbers \( C_n \) and use it to prove that for \( n \geq 1 \),
\[ C_n = 1 + n(n + 1)/2. \]

4.45 Let \( n \geq 1 \) be an integer. Consider \( 2n \) straight lines \( \ell_1, \ell'_1, \ldots, \ell_n, \ell'_n \) such that
4.7. Exercises

- for each $i$ with $1 \leq i \leq n$, $\ell_i$ and $\ell'_i$ are parallel,
- no two of the lines $\ell_1, \ldots, \ell_n$ are parallel,
- no two of the lines $\ell'_1, \ldots, \ell'_n$ are parallel,
- no three of the $2n$ lines intersect in one single point.

These lines divide the plane into regions (some of which are bounded and some of which are unbounded). Denote the number of these regions by $R_n$. From the figure below, you can see that $R_1 = 3$, $R_2 = 9$, and $R_3 = 19$.

- Derive a recurrence for the numbers $R_n$ and use it to prove that $R_n = 2n^2 + 1$ for $n \geq 1$.

4.46 Let $m \geq 1$ and $n \geq 1$ be integers. Consider $m$ horizontal lines and $n$ non-horizontal lines such that
- no two of the non-horizontal lines are parallel,
- no three of the $m + n$ lines intersect in one single point.

These lines divide the plane into regions (some of which are bounded and some of which are unbounded). Denote the number of these regions by $R_{m,n}$. From the figure below, you can see that $R_{4,3} = 23$. 
Derive a recurrence for the numbers $R_{m,n}$ and use it to prove that

$$R_{m,n} = 1 + m(n + 1) + \binom{n+1}{2}.$$  

4.47 A line is called *slanted* if it is neither horizontal nor vertical. Let $k \geq 1$, $m \geq 1$, and $n \geq 0$ be integers. Consider $k$ horizontal lines, $m$ vertical lines, and $n$ slanted lines, such that

- no two of the slanted lines are parallel,
- no three of the $k + m + n$ lines intersect in one single point.

These lines divide the plane into regions (some of which are bounded and some of which are unbounded). Denote the number of these regions by $R_{k,m,n}$. From the figure below, you can see that $R_{4,2,2} = 30$.

Prove that

$$R_{k,m,0} = (k + 1)(m + 1).$$

Derive a recurrence for the numbers $R_{k,m,n}$ and use it to prove that

$$R_{k,m,n} = (k + 1)(m + 1) + (k + m)n + \binom{n+1}{2}.$$
Chapter 5

Discrete Probability

We all have some intuitive understanding of the notions of “chance” and “probability”. When buying a lottery ticket, we know that there is a chance of winning the jackpot, but we also know that this chance is very small. Before leaving home in the morning, we check the weather forecast and see that, with probability 80%, we get 15 centimetres of snow in Ottawa. In this chapter, we will give a formal definition of this notion of “probability”. We start by presenting a surprising application of probability and random numbers.

5.1 Anonymous Broadcasting

Consider a group of \( n \) people \( P_1, P_2, \ldots, P_n \), for some integer \( n \geq 3 \). One person in this group, say \( P_k \), would like to broadcast, *anonymously*, a message to all other people in the group. That is, \( P_k \) wants to broadcast a message such that

- everybody in the group receives the message,
- nobody knows that the message was broadcast by \( P_k \).

In other words, when \( P_i \) (with \( i \neq k \)) receives the message, he only knows that it was broadcast by one of \( P_1, \ldots, P_{i-1}, P_{i+1}, \ldots, P_n \); \( P_i \) cannot determine who broadcast the message.

At first sight, it seems to be impossible to do this. In 1988, however, David Chaum published, in the Journal of Cryptology, a surprisingly simple
protocol that does achieve this. Chaum referred to the problem as the *Dining Cryptographers Problem.*

We will present and analyze the protocol for the case when $n = 3$. Thus, there are three people $P_1$, $P_2$, and $P_3$. We assume that exactly one of them broadcasts a message and refer to this person as the *broadcaster*. We also assume that the message is a binary string. The broadcaster will announce the message one bit at a time.

The three people $P_1$, $P_2$, and $P_3$ sit at a table, in clockwise order of their indices. Let $b$ be the current bit that the broadcaster wants to announce. The protocol for broadcasting this bit is as follows:

**Step 1:** Each person $P_i$ generates a *random* bit $b_i$, for example by flipping a coin. Thus, with 50% probability, $b_i = 0$ and with 50% probability, $b_i = 1$.

**Step 2:** Each person $P_i$ looks at his own bit $b_i$ and also shows it to his clockwise neighbor.

At the end of this second step,

- $P_1$ knows $b_1$ and $b_3$, but not $b_2$,
- $P_2$ knows $b_1$ and $b_2$, but not $b_3$,
- $P_3$ knows $b_2$ and $b_3$, but not $b_1$.

**Step 3:** Each person $P_i$ computes the sum $s_i$ (modulo 2) of the bits that he knows. Thus,

- $P_1$ computes $s_1 = (b_1 + b_3) \mod 2$,
- $P_2$ computes $s_2 = (b_1 + b_2) \mod 2$,
- $P_3$ computes $s_3 = (b_2 + b_3) \mod 2$. 
5.1. Anonymous Broadcasting

**Step 4:** Each person $P_i$ does the following:

- If $P_i$ is not the broadcaster, he sets $t_i = s_i$.
- If $P_i$ is the broadcaster, he sets $t_i = (s_i + b) \mod 2$. Recall that $b$ is the current bit that the broadcaster wants to announce. (Thus, if $b = 1$, then $P_i$ “secretly” flips the bit $s_i$ and stores the result in $t_i$.)

**Step 5:** Each person $P_i$ shows his bit $t_i$ to the other two people.

**Step 6:** Each person $P_i$ computes the sum (modulo 2) of the three bits $t_1$, $t_2$, and $t_3$, i.e., the value $(t_1 + t_2 + t_3) \mod 2$.

This concludes the description of the protocol for broadcasting one bit $b$. Observe that for any bit $x$, we have $(x + x) \mod 2 = 0$. Therefore, the bit computed in the last step is equal to

$$t_1 + t_2 + t_3 = s_1 + s_2 + s_3 + b$$

$$= (b_1 + b_3) + (b_1 + b_2) + (b_2 + b_3) + b$$

$$= (b_1 + b_1) + (b_2 + b_2) + (b_3 + b_3) + b$$

$$= b,$$

where all arithmetic is done modulo 2. In other words, the bit computed in the last step is equal to the bit that the broadcaster wants to announce. This shows that each person in the group receives this bit.

It remains to show that a non-broadcaster cannot determine who broadcast the bit. In the analysis below, we assume that

- $b = 1$, i.e., the broadcaster announces the bit 1,
- $P_2$ is not the broadcaster.

We have to show that $P_2$ cannot determine whether $P_1$ or $P_3$ is the broadcaster. Note that $P_2$ knows the values

$$b_1, b_2, t_1, t_2, t_3,$$

but does not know the bit $b_3$. We consider the cases when $b_1 = b_2$ and $b_1 \neq b_2$ separately.

**Case 1:** $b_1 = b_2$. This case has two subcases depending on the value of $b_3$.

**Case 1.1:** $b_3 = b_1$; thus, all three $b$-bits are equal.
• If $P_1$ is the broadcaster, then (all arithmetic is done modulo 2)

$$t_1 = s_1 + 1 = b_1 + b_3 + 1 = 1$$

and

$$t_3 = s_3 = b_2 + b_3 = 0.$$ 

• If $P_3$ is the broadcaster, then

$$t_1 = s_1 = b_1 + b_3 = 0$$

and

$$t_3 = s_3 + 1 = b_2 + b_3 + 1 = 1.$$ 

Thus, the broadcaster is the person whose $t$-bit is equal to 1.

**Case 1.2:** $b_3 \neq b_1$ and, thus, $b_3 \neq b_2$.

• If $P_1$ is the broadcaster, then

$$t_1 = s_1 + 1 = b_1 + b_3 + 1 = 0$$

and

$$t_3 = s_3 = b_2 + b_3 = 1.$$ 

• If $P_3$ is the broadcaster, then

$$t_1 = s_1 = b_1 + b_3 = 1$$

and

$$t_3 = s_3 + 1 = b_2 + b_3 + 1 = 0.$$ 

Thus, the broadcaster is the person whose $t$-bit is equal to 0.

Since $P_2$ knows $b_1$ and $b_2$, he knows when Case 1 occurs. Since $P_2$ does not know $b_3$, however, he cannot determine whether Case 1.1 or 1.2 occurs. As a result, $P_2$ cannot determine whether $P_1$ or $P_3$ is the broadcaster.

**Case 2:** $b_1 \neq b_2$. This case has two subcases depending on the value of $b_3$.

**Case 2.1:** $b_3 = b_1$ and, thus, $b_3 \neq b_2$. 

**Case 2.2:** $b_3 \neq b_1$ and, thus, $b_3 \neq b_2$. 

**Case 2.3:** $b_3 = b_2$ and, thus, $b_3 \neq b_1$. 

**Case 2.4:** $b_3 \neq b_2$ and, thus, $b_3 \neq b_1$. 

Thus, the broadcaster is the person whose $t$-bit is equal to 1.
5.1. Anonymous Broadcasting

- If $P_1$ is the broadcaster, then
  
  \[ t_1 = s_1 + 1 = b_1 + b_3 + 1 = 1 \]

  and
  
  \[ t_3 = s_3 = b_2 + b_3 = 1. \]

- If $P_3$ is the broadcaster, then
  
  \[ t_1 = s_1 = b_1 + b_3 = 0 \]

  and
  
  \[ t_3 = s_3 + 1 = b_2 + b_3 + 1 = 0. \]

Thus, $t_1$ is always equal to $t_3$, no matter whether $P_1$ or $P_3$ is the broadcaster.

**Case 2.2:** $b_3 \neq b_1$ and, thus, $b_3 = b_2$.

- If $P_1$ is the broadcaster, then
  
  \[ t_1 = s_1 + 1 = b_1 + b_3 + 1 = 0 \]

  and
  
  \[ t_3 = s_3 = b_2 + b_3 = 0. \]

- If $P_3$ is the broadcaster, then
  
  \[ t_1 = s_1 = b_1 + b_3 = 1 \]

  and
  
  \[ t_3 = s_3 + 1 = b_2 + b_3 + 1 = 1. \]

Thus, $t_1$ is always equal to $t_3$, no matter whether $P_1$ or $P_3$ is the broadcaster.

Since $P_2$ knows $b_1$ and $b_2$, he knows when Case 2 occurs. Since $P_2$ does not know $b_3$, however, he cannot determine whether Case 2.1 or 2.2 occurs. As in Case 1, $P_2$ cannot determine whether $P_1$ or $P_3$ is the broadcaster.

We conclude from Cases 1 and 2 that the broadcasting of the bit $b = 1$ is indeed anonymous. Now consider the case when the bit $b$ to be announced is equal to 0. It follows from the protocol that in this case, there is no “secret bit flipping” done in Step 4 and all three people use the same rules to compute the $s$-values and the $t$-values. In this case, $t_1 = s_1$, $t_2 = s_2$, and
$t_3 = s_3$, and $P_2$ can determine the bit $b_3$. He cannot, however, determine whether $P_1$ or $P_3$ is the broadcaster.

To conclude this section, we remark that for each bit to be announced, the entire protocol must be followed. That is, in each round of the protocol, one bit is broadcast and each person $P_i$ must flip a coin to determine the bit $b_i$ that is used in this round. We also remark that the protocol works only if exactly one person is the broadcaster.

## 5.2 Probability Spaces

In this section, we give a formal definition of the notion of “probability” in terms of sets and functions.

**Definition 5.2.1** A *sample space* $S$ is a non-empty countable set. Each element of $S$ is called an *outcome* and each subset of $S$ is called an *event*.

In daily life, we express probabilities in terms of percentages. For example, the weather forecast may tell us that, with 80% probability, we will be getting a snow storm today. In probability theory, probabilities are expressed in terms of numbers in the interval $[0, 1]$. A probability of 80% becomes a probability of 0.8.

**Definition 5.2.2** Let $S$ be a sample space. A *probability function* on $S$ is a function $\text{Pr} : S \rightarrow \mathbb{R}$ such that

- for all $\omega \in S$, $0 \leq \text{Pr}(\omega) \leq 1$, and
- $\sum_{\omega \in S} \text{Pr}(\omega) = 1$.

For any outcome $\omega$ in the sample space $S$, we will refer to $\text{Pr}(\omega)$ as the probability that the outcome is equal to $\omega$.

**Definition 5.2.3** A *probability space* is a pair $(S, \text{Pr})$, where $S$ is a sample space and $\text{Pr}$ is a probability function on $S$.

A probability function $\text{Pr}$ maps each element of the sample space $S$ (i.e., each outcome) to a real number in the interval $[0, 1]$. It turns out to be useful
to extend this function so that it maps any event to a real number in $[0, 1]$. If $A$ is an event (i.e., $A \subseteq S$), then we define

$$\Pr(A) = \sum_{\omega \in A} \Pr(\omega). \quad (5.1)$$

We will refer to $\Pr(A)$ as the probability that the event $A$ occurs. Note that since $S \subseteq S$, the entire sample space $S$ is an event and

$$\Pr(S) = \sum_{\omega \in S} \Pr(\omega) = 1,$$

where the last equality follows from the second condition in Definition 5.2.2.

### 5.2.1 Examples

**Flipping a coin:** Assume we flip a coin. Since there are two possible outcomes (the coin either comes up heads ($H$) or tails ($T$)), the sample space is the set $S = \{H, T\}$. If the coin is fair, i.e., the probabilities of $H$ and $T$ are equal, then the probability function $\Pr : S \rightarrow \mathbb{R}$ is given by

$$\Pr(H) = 1/2, \quad \Pr(T) = 1/2.$$ 

Observe that this function $\Pr$ satisfies the two conditions in Definition 5.2.2. Since this sample space has two elements, there are four events, one event for each subset. These events are

$$\emptyset, \{H\}, \{T\}, \{H, T\},$$

and it follows from (5.1) that

$$\Pr(\emptyset) = 0, \quad \Pr(\{H\}) = \Pr(H) = 1/2, \quad \Pr(\{T\}) = \Pr(T) = 1/2, \quad \Pr(\{H, T\}) = \Pr(H) + \Pr(T) = 1/2 + 1/2 = 1.$$

**Flipping a coin twice:** If we flip a fair coin twice, then there are four possible outcomes, and the sample space becomes $S = \{HH, HT, TH, TT\}$. For example, $HT$ indicates that the first flip resulted in heads, whereas the
second flip resulted in tails. In this case, the probability function \( \Pr : S \to \mathbb{R} \) is given by

\[
\Pr(\text{HH}) = \Pr(\text{HT}) = \Pr(\text{TH}) = \Pr(\text{TT}) = 1/4.
\]

Observe again that this function \( \Pr \) satisfies the two conditions in Definition 5.2.2. Since the sample space consists of 4 elements, the number of events is equal to \( 2^4 = 16 \). For example, \( A = \{HT, TH\} \) is an event and it follows from (5.1) that

\[
\Pr(A) = \Pr(HT) + \Pr(TH) = 1/4 + 1/4 = 1/2.
\]

In words, when flipping a fair coin twice, the probability that we see one heads and one tails (without specifying the order) is equal to 1/2.

**Rolling a die twice:** If we roll a fair die, then there are six possible outcomes (1, 2, 3, 4, 5, and 6), each one occurring with probability 1/6. If we roll this die twice, we obtain the sample space

\[
S = \{(i, j) : 1 \leq i \leq 6, 1 \leq j \leq 6\},
\]

where \( i \) is the result of the first roll and \( j \) is the result of the second roll. Note that \(|S| = 6 \times 6 = 36\). Since the die is fair, each outcome has the same probability. Therefore, in order to satisfy the two conditions in Definition 5.2.2 we must have

\[
\Pr(i, j) = 1/36
\]

for each outcome \((i, j)\) in \( S \).

If we are interested in the sum of the results of the two rolls, then we define the event

\[
A_k = \text{“the sum of the results of the two rolls is equal to } k\text{”},
\]

which, using the notation of sets, is the same as

\[
A_k = \{(i, j) \in S : i + j = k\}.
\]

In the matrix below, the leftmost column indicates the result of the first roll, the top row indicates the result of the second roll, and each entry is the sum of the results of the two corresponding rolls.
As can be seen from this matrix, the event $A_k$ is non-empty only if $k \in \{2, 3, \ldots, 12\}$. For any other $k$, the event $A_k$ is empty, which means that it can never occur.

It follows from (5.1) that

$$\Pr(A_k) = \sum_{(i,j) \in A_k} \Pr(i,j) = \sum_{(i,j) \in A_k} 1/36 = |A_k|/36.$$ 

For example, if $k = 4$, you can see in the matrix that the event $A_4$ has three elements, i.e., there are 3 possible outcomes of two rolls that result in a sum of 4. These outcomes are (3, 1), (2, 2), and (1, 3). It follows that

$$\Pr(A_4) = 3/36 = 1/12.$$ 

In a similar way, we see that

$$\begin{align*}
\Pr(A_2) &= 1/36, \\
\Pr(A_3) &= 2/36 = 1/18, \\
\Pr(A_4) &= 3/36 = 1/12, \\
\Pr(A_5) &= 4/36 = 1/9, \\
\Pr(A_6) &= 5/36, \\
\Pr(A_7) &= 6/36 = 1/6, \\
\Pr(A_8) &= 5/36, \\
\Pr(A_9) &= 4/36 = 1/9, \\
\Pr(A_{10}) &= 3/36 = 1/12, \\
\Pr(A_{11}) &= 2/36 = 1/18, \\
\Pr(A_{12}) &= 1/36.
\end{align*}$$

A sample space is not necessarily uniquely defined. In the last example, where we are interested in the sum of the results of two rolls of a die, we could also have taken the sample space to be the set

$$S' = \{2, 3, 4, 5, 6, 7, 8, 9, 10, 11, 12\}.$$
The probability function $P'_r$ corresponding to this sample space $S'$ is given by

$$P'_r(k) = P(A_k),$$

because $P'_r(k)$ is the probability that we get the outcome $k$ in the sample space $S'$, which is the same as the probability that event $A_k$ occurs in the sample space $S$. You should verify that this function $P'_r$ satisfies the two conditions in Definition 5.2.2 and, thus, is a valid probability function on $S'$.

## 5.3 Basic Rules of Probability

In this section, we prove some basic properties of probability functions. As we will see, all these properties follow from Definition 5.2.2. Throughout this section, $(S, P)$ is a probability space.

Recall that an event is a subset of the sample space $S$. In particular, the empty set $\emptyset$ is an event. Intuitively, $P(\emptyset)$ must be zero, because it is the probability that there is no outcome, which can never occur. The following lemma states that this is indeed the case.

**Lemma 5.3.1** $P(\emptyset) = 0$.

**Proof.** By (5.1), we have

$$P(\emptyset) = \sum_{\omega \in \emptyset} P(\omega).$$

Since there are zero terms in this summation, its value is equal to zero. $\blacksquare$

We say that two events $A$ and $B$ are *disjoint*, if $A \cap B = \emptyset$. A sequence $A_1, A_2, \ldots, A_n$ of events is *pairwise disjoint*, if any pair in this sequence is disjoint. The following lemma is similar to the Sum Rule of Section 3.4.

**Lemma 5.3.2** If $A_1, A_2, \ldots, A_n$ is a sequence of pairwise disjoint events, then

$$P(A_1 \cup A_2 \cup \cdots \cup A_n) = \sum_{i=1}^{n} P(A_i).$$
Proof. Define \( A = A_1 \cup A_2 \cup \cdots \cup A_n \). Using (5.1), we get

\[
\Pr(A) = \sum_{\omega \in A} \Pr(\omega) = \sum_{i=1}^{n} \sum_{\omega \in A_i} \Pr(\omega) = \sum_{i=1}^{n} \Pr(A_i).
\]

To give an example, assume we roll a fair die twice. What is the probability that the sum of the two results is even? If you look at the matrix in Section 5.2, then you see that there are 18 entries, out of 36, that are even. Therefore, the probability of having an even sum is equal to \( 18/36 = 1/2 \). Below we will give a different way to determine this probability.

The sample space is the set

\[
S = \{(i,j) : 1 \leq i \leq 6, 1 \leq j \leq 6\},
\]

where \( i \) is the result of the first roll and \( j \) is the result of the second roll. Each element of \( S \) has the same probability \( 1/36 \) of being an outcome of rolling the die twice.

The event we are interested in is

\[
A = \{(i,j) \in S : i + j \text{ is even}\}.
\]

Observe that \( i + j \) is even if and only if both \( i \) and \( j \) are even or both \( i \) and \( j \) are odd. Therefore, we split the event \( A \) into two disjoint events

\[
A_1 = \{(i,j) \in S : \text{both } i \text{ and } j \text{ are even}\}
\]

and

\[
A_2 = \{(i,j) \in S : \text{both } i \text{ and } j \text{ are odd}\}.
\]

By Lemma 5.3.2, we have

\[
\Pr(A) = \Pr(A_1) + \Pr(A_2).
\]
The set $A_1$ has $3 \cdot 3 = 9$ elements, because there are 3 choices for $i$ and 3 choices for $j$. Similarly, the set $A_2$ has 9 elements. It follows that

$$\Pr(A) = \Pr(A_1) + \Pr(A_2) = 9/36 + 9/36 = 1/2.$$ 

In the next lemma, we relate the probability that an event occurs to the probability that the event does not occur. If $A$ is an event, then $\overline{A}$ denotes its complement, i.e., $\overline{A} = S \setminus A$. Intuitively, the sum of $\Pr(A)$ and $\Pr(\overline{A})$ must be equal to one, because the event $A$ either occurs or does not occur. The following lemma states that this is indeed the case. Observe that this is similar to the Complement Rule of Section 3.3.

Lemma 5.3.3 For any event $A$,

$$\Pr(A) = 1 - \Pr(\overline{A}).$$

Proof. Since $A$ and $\overline{A}$ are disjoint and $S = A \cup \overline{A}$, it follows from Lemma 5.3.2 that

$$\Pr(S) = \Pr(A \cup \overline{A}) = \Pr(A) + \Pr(\overline{A}).$$

We have seen in Section 5.2 that $\Pr(S) = 1$. 

Consider again the sample space that we saw after Lemma 5.3.2. We showed that, when rolling a fair die twice, we get an even sum with probability 1/2. It follows from Lemma 5.3.3 that we get an odd sum with probability $1 - 1/2 = 1/2$.

The next lemma is similar to the Principle of Inclusion and Exclusion that we have seen in Section 3.5.

Lemma 5.3.4 If $A$ and $B$ are events, then

$$\Pr(A \cup B) = \Pr(A) + \Pr(B) - \Pr(A \cap B).$$

Proof. Since $B \setminus A$ and $A \cap B$ are disjoint and $B = (B \setminus A) \cup (A \cap B)$, it follows from Lemma 5.3.2 that

$$\Pr(B) = \Pr(B \setminus A) + \Pr(A \cap B).$$

Next observe that $A$ and $B \setminus A$ are disjoint. Since $A \cup B = A \cup (B \setminus A)$, we again apply Lemma 5.3.2 and obtain

$$\Pr(A \cup B) = \Pr(A) + \Pr(B \setminus A).$$
By combining these two equations, we obtain the claim in the lemma.

To give an example, assume we choose a number $x$ in the sample space $S = \{1, 2, \ldots, 1000\}$, such that each element has the same probability $1/1000$ of being chosen. What is the probability that $x$ is divisible by 2 or 3? Define the events

$$A = \{i \in S : i \text{ is divisible by 2}\}$$

and

$$B = \{i \in S : i \text{ is divisible by 3}\}.$$

Then we want to determine $\Pr(A \cup B)$, which, by Lemma 5.3.4, is equal to

$$\Pr(A \cup B) = \Pr(A) + \Pr(B) - \Pr(A \cap B).$$

Since there are $\lfloor 1000/2 \rfloor = 500$ even numbers in $S$, we have

$$\Pr(A) = 500/1000.$$

Since there are $\lfloor 1000/3 \rfloor = 333$ elements in $S$ that are divisible by 3, we have

$$\Pr(B) = 333/1000.$$

Observe that $i$ belongs to $A \cap B$ if and only if $i$ is divisible by 6, i.e.,

$$A \cap B = \{i \in S : i \text{ is divisible by 6}\}.$$

Since there are $\lfloor 1000/6 \rfloor = 166$ elements in $S$ that are divisible by 6, we have

$$\Pr(A \cap B) = 166/1000.$$

We conclude that

$$\Pr(A \cup B) = \Pr(A) + \Pr(B) - \Pr(A \cap B)$$

$$= 500/1000 + 333/1000 - 166/1000$$

$$= 667/1000.$$

**Lemma 5.3.5 (Union Bound)** For any integer $n \geq 1$, if $A_1, A_2, \ldots, A_n$ is a sequence of events, then

$$\Pr(A_1 \cup A_2 \cup \cdots \cup A_n) \leq \sum_{i=1}^{n} \Pr(A_i).$$
Proof. The proof is by induction on \( n \). If \( n = 1 \), we have equality and, thus, the claim obviously holds. Let \( n \geq 2 \) and assume the claim is true for \( n - 1 \), i.e., assume that

\[
\Pr(A_1 \cup A_2 \cup \cdots \cup A_{n-1}) \leq \sum_{i=1}^{n-1} \Pr(A_i).
\]

Define \( B = A_1 \cup A_2 \cup \cdots \cup A_{n-1} \). Then it follows from Lemma 5.3.4 that

\[
\Pr(B \cup A_n) = \Pr(B) + \Pr(A_n) - \Pr(B \cap A_n) \leq \Pr(B) + \Pr(A_n),
\]

because \( \Pr(B \cap A_n) \geq 0 \) (this follows from the first condition in Definition 5.2.2). Since we assumed that

\[
\Pr(B) \leq \sum_{i=1}^{n-1} \Pr(A_i),
\]

it follows that

\[
\Pr(A_1 \cup A_2 \cup \cdots \cup A_n) = \Pr(B \cup A_n)
\leq \Pr(B) + \Pr(A_n)
\leq \sum_{i=1}^{n-1} \Pr(A_i) + \Pr(A_n)
= \sum_{i=1}^{n} \Pr(A_i).
\]

Lemma 5.3.6 If \( A \) and \( B \) are events with \( A \subseteq B \), then

\[
\Pr(A) \leq \Pr(B).
\]

Proof. Using (5.1) and the fact that \( \Pr(\omega) \geq 0 \) for each \( \omega \) in \( S \), we get

\[
\Pr(A) = \sum_{\omega \in A} \Pr(\omega)
\leq \sum_{\omega \in B} \Pr(\omega)
= \Pr(B).
\]
5.4 Uniform Probability Spaces

In this section, we consider finite sample spaces $S$ in which each outcome has the same probability. Since, by Definition 5.2.2, all probabilities add up to one, the probability of each outcome must be equal to $\frac{1}{|S|}$.

**Definition 5.4.1** A *uniform probability space* is a pair $(S, \Pr)$, where $S$ is a finite sample space and the probability function $\Pr$ satisfies

$$\Pr(\omega) = \frac{1}{|S|} \text{ for each outcome } \omega \text{ in } S.$$ 

The probability spaces that we have seen in Section 5.2.1 are all uniform, except the space $(S', \Pr')$ that we saw at the end of that section.

To give another example, when playing Lotto 6/49, you choose a 6-element subset of the set $A = \{1, 2, \ldots, 49\}$. Twice a week, the Ontario Lottery and Gaming Corporation (OLG) draws the six “winning numbers” uniformly at random from $A$. If your numbers are equal to those drawn by OLG, then you can withdraw from this course and spend the rest of your life on the beach. Most people find it silly to choose the subset $\{1, 2, 3, 4, 5, 6\}$. They think that it is better to choose, for example, the subset $\{2, 5, 16, 36, 41, 43\}$. Is this true?

For this example, the sample space is the set $S$ consisting of all 6-element subsets of $A$. Since $S$ has size $\binom{49}{6}$ and the subset drawn by OLG is uniform, each outcome (i.e., each 6-element subset of $S$) has a probability of

$$\frac{1}{\binom{49}{6}} = \frac{1}{13,983,816} \approx 0.000000072.$$ 

In particular, both $\{1, 2, 3, 4, 5, 6\}$ and $\{2, 5, 16, 36, 41, 43\}$ have the same probability of being the winning numbers. (Still, the latter subset was drawn by OLG on February 8, 2014.)

The lemma below states that in a uniform probability space $(S, \Pr)$, the probability of an event $A$ is just the ratio of the size of $A$ and the size of $S$.

**Lemma 5.4.2** If $(S, \Pr)$ is a uniform probability space and $A$ is an event, then

$$\Pr(A) = \frac{|A|}{|S|}.$$
Proof. By applying (5.1), we get

\[
\Pr(A) = \sum_{\omega \in A} \Pr(\omega) = \sum_{\omega \in A} \frac{1}{|S|} = \frac{1}{|S|} \sum_{\omega \in A} 1 = \frac{|A|}{|S|}.
\]

5.4.1 The Probability of Getting a Full House

In a standard deck of 52 cards, each card has a suit and a rank. There are four suits (spades ♠, hearts ♥, clubs ♣, and diamonds ♦), and 13 ranks (ace, 2, 3, 4, 5, 6, 7, 8, 9, 10, jack, queen, and king).

A hand of five cards is called a full house, if three of the cards are of the same rank and the other two cards are also of the same (but necessarily different) rank. For example, a hand consisting of three sevens and two queens is a full house.

Assume we get a uniformly random hand of five cards. What is the probability that this hand is a full house? To answer this question, we start by observing that the sample space is the set \( S \) consisting of all hands of five cards. Note that

\[
|S| = \binom{52}{5} = 2,598,960.
\]

We assume that the 5-card hand is chosen uniformly at random, so that each hand has a probability of \( 1/|S| \) of being chosen.

Since we are interested in the probability of a random hand being a full house, we define the event \( A \) to be the set of all elements in \( S \) that are full houses. By Lemma 5.4.2 we have

\[
\Pr(A) = \frac{|A|}{|S|}.
\]

Thus, to determine \( \Pr(A) \), it remains to determine the size of the set \( A \), i.e., the total number of full houses. For this, we will use the Product Rule of Section 3.1.

- The procedure is “choose a full house”.

- First task: Choose the rank of the three cards in the full house. There are 13 ways to do this.
Second task: Choose the suits of these three cards. There are $\binom{4}{3}$ ways to do this.

Third task: Choose the rank of the other two cards in the full house. There are 12 ways to do this.

Fourth task: Choose the suits of these two cards. There are $\binom{4}{2}$ ways to do this.

Thus, the number of full houses is equal to

$$|A| = 13 \cdot \binom{4}{3} \cdot 12 \cdot \binom{4}{2} = 3,744.$$  

We conclude that the probability of getting a full house is equal to

$$\Pr(A) = \frac{|A|}{|S|} = \frac{3,744}{2,598,960} \approx 0.00144.$$  

5.5 The Birthday Paradox

Let $n \geq 2$ be an integer and consider a group of $n$ people. In this section, we will determine the probability $p_n$ that at least two of them have the same birthday. We will ignore leap years, so that there are 365 days in one year.

Below, we will show that $p_2 = 1/365$. If $n \geq 366$, then it follows from the Pigeonhole Principle (see Section 3.10) that there must be at least two people with the same birthday and, therefore, $p_n = 1$. Intuitively, if $n$ increases from 2 to 365, the value of $p_n$ should increase as well. What is the value of $n$ such that $p_n$ is larger than $1/2$ for the first time? That is, what is the value of $n$ for which $p_{n-1} \leq 1/2 < p_n$? In this section, we will see that this question can be answered using simple counting techniques that we have seen in Chapter 3.

We denote the people by $P_1, P_2, \ldots, P_n$, we denote the number of days in one year by $d$, and we number the days in one year as $1, 2, \ldots, d$. The sample space is the set

$$S_n = \{(b_1, b_2, \ldots, b_n) : b_i \in \{1, 2, \ldots, d\} \text{ for each } 1 \leq i \leq n\},$$  

where $b_i$ denotes the birthday of $P_i$. Note that

$$|S_n| = d^n.$$
We consider the uniform probability space; thus, for each element \((b_1, b_2, \ldots, b_n)\) in \(S_n\), we have
\[
\Pr (b_1, b_2, \ldots, b_n) = \frac{1}{|S_n|} = \frac{1}{d^n}.
\]
The event we are interested in is
\[
A_n = \text{"at least two of the numbers in } b_1, b_2, \ldots, b_n \text{ are equal"}.
\]
Using the notation of sets, this is the same as
\[
A_n = \{(b_1, b_2, \ldots, b_n) \in S_n : b_1, b_2, \ldots, b_n \text{ contains duplicates}\}.
\]
The probability \(p_n\) that we introduced above is equal to
\[
p_n = \Pr (A_n).
\]
As mentioned above, the Pigeonhole Principle implies that \(p_n = 1\) for \(n > d\). Therefore, we assume from now on that \(n \leq d\).

Let us start by determining \(p_2\). Since we consider the uniform probability space, we have, by Lemma 5.4.2,
\[
p_2 = \Pr (A_2) = \frac{|A_2|}{|S_2|}.
\]
We know already that \(|S_2| = d^2\). The event \(A_2\) is equal to
\[
A_2 = \{(1, 1), (2, 2), \ldots, (d, d)\}.
\]
Thus, \(|A_2| = d\) and we obtain
\[
p_2 = \frac{|A_2|}{|S_2|} = \frac{d}{d^2} = \frac{1}{d}.
\]
To determine \(p_n\) for larger values of \(n\), it is easier to determine the probability of the complement \(\overline{A}_n\). The latter probability, together with Lemma 5.3.3, will give us the value of \(p_n\). Note that
\[
\overline{A}_n = \{(b_1, b_2, \ldots, b_n) \in S_n : b_1, b_2, \ldots, b_n \text{ are pairwise distinct}\}.
\]
In other words, \(\overline{A}_n\) is the set of all ordered sequences consisting of \(n\) pairwise distinct elements of \(\{1, 2, \ldots, d\}\). In Section 3.6 see (3.1), we have seen that
\[
|\overline{A}_n| = \frac{d!}{(d-n)!}.
\]
We conclude that, for any \( n \) with \( 2 \leq n \leq d \),

\[
\begin{align*}
p_n &= \Pr(A_n) \\
     &= 1 - \Pr(\bar{A}_n) \\
     &= 1 - \frac{|A_n|}{|S_n|} \\
     &= 1 - \frac{d!}{(d - n)!d^n}.
\end{align*}
\]

By taking \( d = 365 \), we get \( p_{22} = 0.476 \) and \( p_{23} = 0.507 \). Thus, in a random group of 23 people\(^1\), the probability that at least two of them have the same birthday is more than 50%. Most people are very surprised when they see this for the first time, because our intuition says that a much larger group is needed to have a probability of more than 50%. The values \( p_n \) approach 1 pretty fast. For example, \( p_{40} = 0.891 \) and \( p_{100} = 0.9999997 \).

### 5.5.1 Throwing Balls into Boxes

When we derived the expression for \( p_n \), we did not use the fact that the value of \( d \) is equal to 365. In other words, the expression is valid for any value of \( d \). For general values of \( d \), we can interpret the birthday problem in the following way: Consider \( d \) boxes \( B_1, B_2, \ldots, B_d \), where \( d \) is a large integer. Assume that we throw \( n \) balls into these boxes so that each ball lands in a uniformly random box. Then \( p_n \) is the probability that there is at least one box that contains more than one ball. Since it is not easy to see how the expression

\[
p_n = 1 - \frac{d!}{(d - n)!d^n}
\]

depends on \( n \), we will approximate it by a simpler expression. For this, we will use the inequality

\[
1 - x \leq e^{-x}, \tag{5.2}
\]

which is valid for any real number \( x \). If \( x \) is close to zero, then the inequality is very accurate. The easiest way to prove this inequality is by showing that the minimum of the function \( f(x) = x + e^{-x} \) is equal to \( f(0) = 1 \), using techniques from calculus.

\(^1\)two soccer teams plus the referee
If we define \( q_n = 1 - p_n \), then we have

\[
q_n = \frac{d!}{(d - n)!d^n}.
\]

Using (5.2), we get

\[
q_n = \frac{d \cdot d - 1 \cdot d - 2 \cdot d - 3 \cdots d - (n - 1)}{d}
= \frac{d - 1 \cdot d - 2 \cdot d - 3 \cdots d - (n - 1)}{d}
= \frac{(1 - 1/d) \cdot (1 - 2/d) \cdot (1 - 3/d) \cdots (1 - (n - 1)/d)}{d}
\leq e^{-1/d} \cdot e^{-2/d} \cdot e^{-3/d} \cdots e^{-(n-1)/d}
= e^{-(1+2+3+\cdots+(n-1))/d}.
\]

Using the equality

\[
1 + 2 + 3 + \cdots + (n - 1) = n(n - 1)/2,
\]

see Theorem 2.2.10 we thus get

\[
q_n \leq e^{-n(n-1)/(2d)},
\]

and therefore,

\[
p_n = 1 - q_n \geq 1 - e^{-n(n-1)/(2d)}.
\]

If \( n \) is large, then \( n(n - 1)/(2d) \) is very close to \( n^2/(2d) \) and, thus,

\[
p_n \gtrsim 1 - e^{-n^2/(2d)}.
\]

If we take \( n = \sqrt{2d} \), then we get

\[
p_n \gtrsim 1 - e^{-1} \approx 0.632.
\]

Thus, for large values of \( d \), if we throw \( \sqrt{2d} \) balls into \( d \) boxes, then with probability (approximately) at least \( 1 - 1/e \), there is a box that contains more than one ball.
5.6 The Big Box Problem

Keith chooses two distinct elements $x$ and $y$, with $x < y$, from the set $A = \{0, 1, 2, \ldots, 100\}$. He takes two identical boxes, and puts $x$ dollars in one box and $y$ dollars in the other box. Then Keith closes the two boxes, shuffles them, and puts them on a table. At this moment, we neither know $x$ nor $y$; the only information we have is that these are distinct elements from the set $A$.

We will refer to the box containing $x$ dollars as the small box and to the box containing $y$ dollars as the big box. Our goal is to find the big box. We are allowed to do the following:

1. We can choose one of the two boxes, open it, and count how much money is inside it.

2. Now we have to make our final decision: Either we keep the box we just opened or we take the other box.

For example, assume that the box we pick in the first step contains $33$ dollars. Then we know that the other box contains either less than $33$ or more than $33$. It seems that the only reasonable thing to do is to flip a fair coin when making our final decision. If we do that, then we find the big box with probability 0.5.

In the rest of this section, we will show the surprising result that we can find the big box with probability at least 0.505.

The idea is as follows. Assume that we know a number $z$ such that $x < z < y$. (Keep in mind that we do not know $x$ and we do not know $y$. Thus, we assume that we know a number $z$ that is between the two unknown numbers $x$ and $y$.)

- If the box we choose in the first step contains more than $z$ dollars, then we know that this is the big box and, therefore, we keep it.

- If the box we choose in the first step contains less than $z$ dollars, then we know that this is the small box and, therefore, we take the other box.
Thus, if we know this number $z$ with $x < z < y$, then we are guaranteed to find the big box.

Of course, it is not realistic to assume that we know this magic number $z$. The trick is to choose a random $z$ and hope that it is between $x$ and $y$. If $z$ is between $x$ and $y$, then we find the big box with probability $1/2$; otherwise, we find the big box with probability $1/2$. As we will see later, the overall probability of finding the big box will be at least 0.505.

In order to avoid the case when $z = x$ or $z = y$, we will choose $z$ from the set

$$B = \{1/2, 3/2, 5/2, \ldots, 100 - 1/2\}.$$  

Note that $|B| = 100$. Our algorithm that attempts to find the big box does the following:

\textbf{Algorithm FindBigBox:}

\textbf{Step 1:} Choose one of the two boxes uniformly at random, open it, and count the amount of money inside it; let this amount be $a$.

\textbf{Step 2:} Choose $z$ uniformly at random from the set $B$.

\textbf{Step 3:} Do the following:

- If $a > z$, then keep the box chosen in Step 1.
- Otherwise (i.e., if $a < z$), take the other box.

\subsection{The Probability of Finding the Big Box}

We are now going to determine the probability that this algorithm finds the big box. First, we have to ask ourselves what the sample space is. There are two places in the algorithm where a random element is chosen:

- In Step 1, we choose the element $a$, which is a random element from the set $\{x, y\}$. We know that this value $a$ is equal to one of $x$ and $y$. However, at the end of Step 1, we do not know whether $a = x$ or $a = y$.

- In Step 2, we choose a random element from the set $B$.

Based on this, the sample space $S$ is the Cartesian product

$$S = \{x, y\} \times B = \{(a, z) : a \in \{x, y\}, z \in B\}$$
and Steps 1 and 2 can be replaced by

- choose a uniformly random element \((a, z)\) in \(S\).

Note that \(|S| = 200\).

We say that algorithm \textsc{FindBigBox} is \textit{successful} if it finds the big box. Thus, we want to determine \(\Pr(W)\), where \(W\) is the event

\[ W = \text{“algorithm \textsc{FindBigBox} is successful”}. \]

We are going to write this event as a subset of the sample space \(S\). For this, we have to determine all elements \((a, z)\) in \(S\) for which algorithm \textsc{FindBigBox} is successful.

First consider the case when \(a = x\). In this case, the box we choose in Step 1 is the small box. There are two possibilities for \(z\):

- If \(x = a > z\), then the algorithm keeps the small box and, thus, is not successful.
- If \(x = a < z\), then the algorithm takes the other box (which is the big box) and, thus, is successful.

Thus, the event \(W\) contains the set

\[ W_x = \{(x, z) : z \in \{x + 1/2, x + 3/2, \ldots, 100 - 1/2\}\}. \]

You can verify that

\[ |W_x| = 100 - x. \]

The second case to consider is when \(a = y\). In this case, the box we choose in Step 1 is the big box. Again, there are two possibilities for \(z\):

- If \(y = a > z\), then the algorithm keeps the big box and, thus, is successful.
- If \(y = a < z\), then the algorithm takes the other box (which is the small box) and, thus, is not successful.

Thus, the event \(W\) contains the set

\[ W_y = \{(y, z) : z \in \{1/2, 3/2, \ldots, y - 1/2\}\}. \]
You can verify that

\[ |W_y| = y. \]

We conclude that

\[ W = W_x \cup W_y. \]

Since the events \( W_x \) and \( W_y \) are disjoint, we can use Lemma 5.3.2 to determine the probability that algorithm FINDBIGBOX is successful:

\[
\Pr(W) = \Pr(W_x \cup W_y) \\
= \Pr(W_x) + \Pr(W_y) \\
= \frac{|W_x|}{|S|} + \frac{|W_y|}{|S|} \\
= \frac{100 - x}{200} + \frac{y}{200} \\
= \frac{1}{2} + \frac{y - x}{200}.
\]

Since \( x \) and \( y \) are distinct integers and \( x < y \), we have \( y - x \geq 1 \), and we get

\[
\Pr(W) \geq \frac{1}{2} + \frac{1}{200} = 0.505.
\]

## 5.7 The Monty Hall Problem

The Monty Hall Problem is a well-known puzzle in probability theory. It is named after the host, Monty Hall, of the American television game show Let’s Make a Deal. The problem became famous in 1990, when (part of) a reader’s letter was published in Marilyn vos Savant’s column Ask Marilyn in the magazine Parade:

Suppose you’re on a game show, and you’re given the choice of three doors: Behind one door is a car; behind the others, goats. You pick a door, say No. 1, and the host, who knows what’s behind the doors, opens another door, say No. 3, which has a goat. He then says to you, “Do you want to pick door No. 2?” Is it to your advantage to switch your choice?

Note that the host can always open a door that has a goat behind it. After the host has opened No. 3, we know that the car is either behind No. 1
or No. 2, and it seems that both these doors have the same probability (i.e., 50%) of having the car behind them. We will prove below, however, that this is not true: It is indeed to our advantage to switch our choice.

We assume that the car is equally likely to be behind any of the three doors. Moreover, the host knows what is behind each door.

- We initially choose one of the three doors uniformly at random; this door remains closed.
- The host opens one of the other two doors that has a goat behind it.
- Our final choice is to switch to the other door that is still closed.

Let $A$ be the event that we win the car and let $B$ be the event that the initial door has a goat behind it. Then it is not difficult to see that event $A$ occurs if and only if event $B$ occurs. Therefore, the probability that we win the car is equal to

$$\Pr(A) = \Pr(B) = \frac{2}{3}.$$ 

5.8 Conditional Probability

Anil Maheshwari has two children. We are told that one of them is a boy. What is the probability that the other child is also a boy? Most people will say that this probability is $\frac{1}{2}$. We will show below that this is not the correct answer.

Since Anil has two children, the sample space is

$$S = \{(b,b), (b,g), (g,b), (g,g)\},$$

where, for example, $(b,g)$ indicates that the youngest child is a boy and the oldest child is a girl. We assume a uniform probability function, so that each outcome has a probability of $1/4$.

We are given the additional information that one of the two children is a boy, or, to be more precise, that at least one of the two children is a boy. This means that the actual sample space is not $S$, but

$$\{(b,b), (b,g), (g,b)\}.$$

When asking for the probability that the other child is also a boy, we are really asking for the probability that both children are boys. Since there is
only one possibility (out of three) for both children to be boys, it follows that this probability is equal to $1/3$.

This is an example of a \textit{conditional probability}: We are asking for the probability of an event (both children are boys) given that another event (at least one of the two children is a boy) occurs.

**Definition 5.8.1** Let $(S, \Pr)$ be a probability space and let $A$ and $B$ be two events with $\Pr(B) > 0$. The \textit{conditional probability} $\Pr(A \mid B)$, pronounced as “the probability of $A$ given $B$”, is defined to be

$$\Pr(A \mid B) = \frac{\Pr(A \cap B)}{\Pr(B)}.$$ 

Let us try to understand where this definition comes from. Initially, the sample space is equal to $S$. When we are given the additional information that event $B$ occurs, the sample space “shrinks” to $B$, and event $A$ occurs if and only if event $A \cap B$ occurs.

![Venn Diagram](image)

You may think that $\Pr(A \mid B)$ should therefore be defined to be $\Pr(A \cap B)$. However, since the sum of all probabilities must be equal to 1, we have to normalize, i.e., divide by $\Pr(B)$. Equivalently, if $A = B$, we get $\Pr(A \mid A)$, which is the probability that event $A$ occurs, given that event $A$ occurs. This probability should be equal to 1. Indeed, using the definition, we do get

$$\Pr(A \mid A) = \frac{\Pr(A \cap A)}{\Pr(A)} = \frac{\Pr(A)}{\Pr(A)} = 1.$$ 

In Exercise 5.13 you are asked to give a formal proof that our definition gives a valid probability function on the sample space $S$. 
5.8.1 Anil’s Children

Returning to Anil’s two children, we saw that the sample space is
\[ S = \{ (b,b), (b,g), (g,b), (g,g) \} \]
and we assumed a uniform probability function. The events we considered are
\[ A = \text{“both children are boys”} \]
and
\[ B = \text{“at least one of the two children is a boy”}, \]
and we wanted to know \( \Pr(A \mid B) \). Writing \( A \) and \( B \) as subsets of the sample space \( S \), we get
\[ A = \{(b,b)\} \]
\[ B = \{(b,b), (b,g), (g,b)\}. \]
It follows that
\[ \Pr(A \mid B) = \frac{\Pr(A \cap B)}{\Pr(B)} = \frac{\Pr(A)}{\Pr(B)} = \frac{1/4}{3/4} = \frac{1}{3}, \]
which is the same answer as we got before.

5.8.2 Rolling a Die

Assume we roll a fair die and consider the events
\[ A = \text{“the result is 3”} \]
and
\[ B = \text{“the result is odd”}. \]
Then
\[ \Pr(A \mid B) = \frac{\Pr(A \cap B)}{\Pr(B)} = \frac{\Pr(A)}{\Pr(B)} = \frac{1/6}{3/6} = \frac{1}{3} \]
and
\[ \Pr(B \mid A) = \frac{\Pr(B \cap A)}{\Pr(A)} = \frac{\Pr(A)}{\Pr(A)} = 1. \]
This shows that, in general, \( \Pr(A \mid B) \) is not equal to \( \Pr(B \mid A) \).
Consider the event

\[ C = \text{“the result is a prime number”}, \]

i.e., the result is 2, 3, or 5. We have

\[
\Pr(C \mid B) = \frac{\Pr(C \cap B)}{\Pr(B)} = \frac{2/6}{3/6} = \frac{2}{3}
\]

and

\[
\Pr(C \mid A) = \frac{\Pr(C \cap A)}{\Pr(A)} = \frac{\Pr(A)}{\Pr(A)} = 1.
\]

Recall that \( \overline{A} \) denotes the complement of \( A \). Thus, this is the event

\[ \overline{A} = \text{“the result is not 3”}. \]

We have

\[
\Pr(C \mid \overline{A}) = \frac{\Pr(C \cap \overline{A})}{\Pr(\overline{A})} = \frac{2/6}{5/6} = \frac{2}{5}.
\]

This shows that, in general, \( \Pr(C \mid A) + \Pr(C \mid \overline{A}) \) is not equal to 1. It should be an easy exercise to verify that, for the above events \( A \) and \( C \), \( \Pr(A \mid C) + \Pr(\overline{A} \mid C) \) is equal to 1. Intuitively, this should be true for any two events \( A \) and \( C \) with \( \Pr(C) > 0 \): When we are given that event \( C \) occurs, then either \( A \) occurs or \( A \) does not occur (in which case \( \overline{A} \) occurs). The following lemma states that this intuition is indeed correct.

**Lemma 5.8.2** Let \( A \) and \( B \) be events with \( \Pr(B) > 0 \). Then

\[
\Pr(A \mid B) + \Pr(\overline{A} \mid B) = 1.
\]

**Proof.** By definition, we have

\[
\Pr(A \mid B) + \Pr(\overline{A} \mid B) = \frac{\Pr(A \cap B)}{\Pr(B)} + \frac{\Pr(\overline{A} \cap B)}{\Pr(B)} = \frac{\Pr(A \cap B) + \Pr(\overline{A} \cap B)}{\Pr(B)}.
\]

Since the events \( A \cap B \) and \( \overline{A} \cap B \) are disjoint, we have, by Lemma 5.3.2

\[
\Pr(A \cap B) + \Pr(\overline{A} \cap B) = \Pr((A \cap B) \cup (\overline{A} \cap B)).
\]
By drawing a Venn diagram, you will see that

\[(A \cap B) \cup (\overline{A} \cap B) = B,\]

implying that

\[\Pr(A \cap B) + \Pr(\overline{A} \cap B) = \Pr(B).\]

We conclude that

\[\Pr(A \mid B) + \Pr(\overline{A} \mid B) = \frac{\Pr(B)}{\Pr(B)} = 1.\]

5.9 The Law of Total Probability

Both Mick and Keith have a random birthday. What is the probability that they have the same birthday? We have seen in Section 5.5 that this probability is equal to \(1/365\). A common way to determine this probability is as follows: Consider Mick’s birthday, which can be any of the 365 days of the year. By symmetry, it does not really matter what Mick’s birthday is, so we just assume that it is July 26. Then Mick and Keith have the same birthday if and only Keith’s birthday is also on July 26. Therefore, since Keith has a random birthday, the probability that Mick and Keith have the same birthday is equal to \(1/365\). The following theorem explains this reasoning.

**Theorem 5.9.1 (Law of Total Probability)** Let \((S, \Pr)\) be a probability space and let \(A\) be an event. Assume that \(B_1, B_2, \ldots, B_n\) is a sequence of events such that

1. \(\Pr(B_i) > 0\) for all \(i\) with \(1 \leq i \leq n\),
2. the events \(B_1, B_2, \ldots, B_n\) are pairwise disjoint, and
3. \(\bigcup_{i=1}^{n} B_i = S\).

Then

\[\Pr(A) = \sum_{i=1}^{n} \Pr(A \mid B_i) \cdot \Pr(B_i).\]
Proof. The assumptions imply that

\[ A = A \cap S = A \cap \left( \bigcup_{i=1}^{n} B_i \right) = \bigcup_{i=1}^{n} (A \cap B_i). \]

Since the events \( A \cap B_1, A \cap B_2, \ldots, A \cap B_n \) are pairwise disjoint, it follows from Lemma 5.3.2 that

\[ \Pr(A) = \Pr\left( \bigcup_{i=1}^{n} (A \cap B_i) \right) = \sum_{i=1}^{n} \Pr(A \cap B_i). \]

The theorem follows by observing that, from Definition 5.8.1

\[ \Pr(A \cap B_i) = \Pr(A \mid B_i) \cdot \Pr(B_i). \]

Let us consider the three conditions in this theorem. The first condition is that \( \Pr(B_i) > 0 \), i.e., there is a positive probability that event \( B_i \) occurs. The second and third conditions, i.e.,

- the events \( B_1, B_2, \ldots, B_n \) are pairwise disjoint, and
- \( \bigcup_{i=1}^{n} B_i = S \),

are equivalent to

- exactly one of the events \( B_1, B_2, \ldots, B_n \) is guaranteed to occur.

In the example in the beginning of this section, we wanted to know \( \Pr(A) \), where \( A \) is the event

\[ A = \text{“Mick and Keith have the same birthday”}. \]

In order to apply Theorem 5.9.1, we define a sequence \( B_1, B_2, \ldots \) of events that satisfy the conditions in this theorem and for which \( \Pr(A \mid B_i) \) is easy
to determine. For this example, we define the event $B_i$, for each $i$ with $1 \leq i \leq 365$, to be

$$B_i = \text{“Mick’s birthday is on the } i\text{-th day of the year”}.$$ 

It is clear that $\Pr(B_i) = 1/365 > 0$ and exactly one of the events $B_1, B_2, \ldots, B_{365}$ is guaranteed to occur. It follows that

$$\Pr(A) = \sum_{i=1}^{365} \Pr(A \mid B_i) \cdot \Pr(B_i).$$

To determine $\Pr(A \mid B_i)$, we assume that the event $B_i$ occurs, i.e., we fix Mick’s birthday to be the $i$-th day of the year. Given this event $B_i$, event $A$ occurs if and only if Keith’s birthday is also on the $i$-th day. Thus, we have $\Pr(A \mid B_i) = 1/365$ and it follows that

$$\Pr(A) = \sum_{i=1}^{365} (1/365) \cdot \Pr(B_i)$$

$$= (1/365) \sum_{i=1}^{365} \Pr(B_i)$$

$$= (1/365) \cdot 1$$

$$= 1/365,$$

which is the same answer as we got in the beginning of this section.

### 5.9.1 Flipping a Coin and Rolling Dice

Consider the following experiment:

- We flip a fair coin.
  - If the coin comes up heads, then we roll a fair die. Let $R$ denote the result of this die.
  - If the coin comes up tails, then we roll two fair dice. Let $R$ denote the sum of the results of these dice.
What is the probability that the value of \( R \) is equal to 2? That is, if we define the event \( A \) to be

\[
A = \text{“the value of } R \text{ is equal to } 2\text{”},
\]

then we want to know \( \Pr(A) \). Since the value of \( R \) depends on whether the coin comes up heads or tails, we define the event

\[
B = \text{“the coin comes up heads”}.
\]

Since (i) both \( B \) and its complement \( \overline{B} \) occur with a positive probability and (ii) exactly one of \( B \) and \( \overline{B} \) is guaranteed to occur, we can apply Theorem 5.9.1 and get

\[
\Pr(A) = \Pr(A \mid B) \cdot \Pr(B) + \Pr(A \mid \overline{B}) \cdot \Pr(\overline{B}).
\]

We determine the four terms on the right-hand side:

- It should be clear that

\[
\Pr(B) = \Pr(\overline{B}) = \frac{1}{2}.
\]

- To determine \( \Pr(A \mid B) \), we assume that the event \( B \) occurs, i.e., the coin comes up heads. Because of this assumption, we roll one die, and the event \( A \) occurs if and only if the result of this roll is 2. It follows that

\[
\Pr(A \mid B) = \frac{1}{6}.
\]

- To determine \( \Pr(A \mid \overline{B}) \), we assume that the event \( \overline{B} \) occurs, i.e., the coin comes up tails. Because of this assumption, we roll two dice, and the event \( A \) occurs if and only if both rolls result in 1. Since there are 36 possible outcomes when rolling two dice, it follows that

\[
\Pr(A \mid \overline{B}) = \frac{1}{36}.
\]

We conclude that

\[
\Pr(A) = \Pr(A \mid B) \cdot \Pr(B) + \Pr(A \mid \overline{B}) \cdot \Pr(\overline{B})
\]

\[
= \frac{1}{6} \cdot \frac{1}{2} + \frac{1}{36} \cdot \frac{1}{2}
\]

\[
= \frac{7}{72}.
\]
5.10 Independent Events

Consider two events $A$ and $B$ in a sample space $S$. We would like to define the notion of these two events being “independent”. Intuitively, this should express that (i) whether or not event $A$ occurs does not depend on whether or not event $B$ occurs and (ii) whether or not event $B$ occurs does not depend on whether or not event $A$ occurs. In other words, if we assume that $\Pr(A) > 0$ and $\Pr(B) > 0$, then (i) $\Pr(A)$ should be equal to the conditional probability $\Pr(A \mid B)$ and (ii) $\Pr(B)$ should be equal to conditional probability $\Pr(B \mid A)$. As we will show below, the following definition exactly captures this.

**Definition 5.10.1** Let $(S, \Pr)$ be a probability space and let $A$ and $B$ be two events. We say that $A$ and $B$ are independent if

$$\Pr(A \cap B) = \Pr(A) \cdot \Pr(B).$$

Note that in this definition, it is not assumed that $\Pr(A) > 0$ and $\Pr(B) > 0$. If $\Pr(B) > 0$, then

$$\Pr(A \mid B) = \frac{\Pr(A \cap B)}{\Pr(B)},$$

and $A$ and $B$ are independent if and only if

$$\Pr(A \mid B) = \Pr(A).$$

Similarly, if $\Pr(A) > 0$, then $A$ and $B$ are independent if and only if

$$\Pr(B \mid A) = \Pr(B).$$

5.10.1 Rolling Two Dice

Assume we roll a red die and a blue die; thus, the sample space is

$$S = \{(i, j) : 1 \leq i \leq 6, 1 \leq j \leq 6\},$$

where $i$ is the result of the red die and $j$ is the result of the blue die. We assume a uniform probability function. Thus, each outcome has a probability of $1/36$.

Let $D_1$ denote the result of the red die and let $D_2$ denote the result of the blue die. Consider the events

$$A = \text{"}D_1 + D_2 = 7\text{"}$$
and

\[ B = "D_1 = 4". \]

Are these events independent?

- Since

\[ A = \{(1, 6), (2, 5), (3, 4), (4, 3), (5, 2), (6, 1)\}, \]

we have \( \Pr(A) = 6/36 = 1/6. \)

- Since

\[ B = \{(4, 1), (4, 2), (4, 3), (4, 4), (4, 5), (4, 6)\}, \]

we have \( \Pr(B) = 6/36 = 1/6. \)

- Since

\[ A \cap B = \{(4, 3)\}, \]

we have \( \Pr(A \cap B) = 1/36. \)

- It follows that \( \Pr(A \cap B) = \Pr(A) \cdot \Pr(B) \) and we conclude that \( A \) and \( B \) are independent.

As an exercise, you should verify that the events

\[ A' = "D_1 + D_2 = 11" \]

and

\[ B' = "D_1 = 5" \]

are not independent.

### 5.10.2 A Basic Property of Independent Events

Consider two events \( A \) and \( B \). If these events are independent, then whether or not \( A \) occurs does not depend on whether or not \( B \) occurs. Since whether or not \( B \) occurs is the same as whether the complement \( \overline{B} \) does not or does occur, it should not be a surprise that the events \( A \) and \( \overline{B} \) are independent as well. The following lemma states that this is indeed the case.

**Lemma 5.10.2** If \( A \) and \( B \) are independent events, then \( A \) and \( \overline{B} \) are also independent.
5.10. Independent Events

Proof. To prove that $A$ and $\overline{B}$ are independent, we have to show that

$$\Pr(A \cap \overline{B}) = \Pr(A) \cdot \Pr(\overline{B}).$$

Using Lemma 5.3.3, this is equivalent to showing that

$$\Pr(A \cap \overline{B}) = \Pr(A) \cdot (1 - \Pr(B)). \quad (5.3)$$

Since the events $A \cap B$ and $A \cap \overline{B}$ are disjoint and

$$A = (A \cap B) \cup (A \cap \overline{B}),$$

it follows from Lemma 5.3.2 that

$$\Pr(A) = \Pr(A \cap B) + \Pr(A \cap \overline{B}).$$

Since $A$ and $B$ are independent, we have

$$\Pr(A \cap B) = \Pr(A) \cdot \Pr(B).$$

It follows that

$$\Pr(A) = \Pr(A) \cdot \Pr(B) + \Pr(A \cap \overline{B}),$$

which is equivalent to (5.3). \qed

5.10.3 Pairwise and Mutually Independent Events

We have defined the notion of two events being independent. The following definition generalizes this in two ways to sequences of events:

Definition 5.10.3 Let $(S, \Pr)$ be a probability space, let $n \geq 2$, and let $A_1, A_2, \ldots, A_n$ be a sequence of events.

1. We say that this sequence is pairwise independent if for any two distinct indices $i$ and $j$, the events $A_i$ and $A_j$ are independent, i.e.,

$$\Pr(A_i \cap A_j) = \Pr(A_i) \cdot \Pr(A_j).$$

2. We say that this sequence is mutually independent if for all $k$ with $2 \leq k \leq n$ and all $i_1 < i_2 < \ldots < i_k$,

$$\Pr(A_{i_1} \cap A_{i_2} \cap \cdots \cap A_{i_k}) = \Pr(A_{i_1}) \cdot \Pr(A_{i_2}) \cdots \Pr(A_{i_k}).$$
Thus, in order to show that the sequence $A_1, A_2, \ldots, A_n$ is pairwise independent, we have to verify $\binom{n}{2}$ equalities. On the other hand, to show that this sequence is mutually independent, we have to verify $\sum_{k=2}^{n} \binom{n}{k} = 2^n - 1 - n$ equalities.

For example, if we want to prove that the sequence $A, B, C$ of three events is mutually independent, then we have to show that

$$\Pr(A \cap B) = \Pr(A) \cdot \Pr(B),$$
$$\Pr(A \cap C) = \Pr(A) \cdot \Pr(C),$$
$$\Pr(B \cap C) = \Pr(B) \cdot \Pr(C),$$

and

$$\Pr(A \cap B \cap C) = \Pr(A) \cdot \Pr(B) \cdot \Pr(C).$$

To give an example, consider flipping three coins and assume that the result is a uniformly random element from the sample space

$$S = \{HHH, HHT, HTH, THH, HTT, THT, TTH, TTT\},$$

where, e.g., $HHT$ indicates that the first two coins come up heads and the third coin comes up tails. For $i = 1, 2, 3$, let $f_i$ denote the result of the $i$-th flip, and define the events

$$A = \text{“} f_1 = f_2 \text{”},$$
$$B = \text{“} f_2 = f_3 \text{”},$$

and

$$C = \text{“} f_1 = f_3 \text{”}.$$

If we write these events as subsets of the sample space, then we get

$$A = \{HHH, HHT, TTH, TTT\},$$
$$B = \{HHH, THH, HTT, TTT\},$$

and

$$C = \{HHH, HTH, THT, TTT\}.$$

It follows that

$$\Pr(A) = \frac{|A|/|S|} = \frac{4}{8} = \frac{1}{2},$$
$$\Pr(B) = \frac{|B|/|S|} = \frac{4}{8} = \frac{1}{2},$$
$$\Pr(C) = \frac{|C|/|S|} = \frac{4}{8} = \frac{1}{2},$$
$$\Pr(A \cap B) = \frac{|A \cap B|/|S|} = \frac{2}{8} = \frac{1}{4},$$
$$\Pr(A \cap C) = \frac{|A \cap C|/|S|} = \frac{2}{8} = \frac{1}{4},$$
$$\Pr(B \cap C) = \frac{|B \cap C|/|S|} = \frac{2}{8} = \frac{1}{4}.$$
Thus, the sequence $A, B, C$ is pairwise independent. Since

$$\Pr(A \cap B \cap C) = |A \cap B \cap C|/|S| = 2/8 = 1/4,$$

we have

$$\Pr(A \cap B \cap C) \neq \Pr(A) \cdot \Pr(B) \cdot \Pr(C)$$

and, therefore, the sequence $A, B, C$ is not mutually independent. Of course, this is not surprising: If both events $A$ and $B$ occur, then event $C$ also occurs.

## 5.11 Describing Events by Logical Propositions

We have defined an event to be a subset of a sample space. In several examples, however, we have described events in plain English or as logical propositions.

- Since the intersection ($\cap$) of sets corresponds to the conjunction ($\land$) of propositions, we often write $A \land B$ for the event “both $A$ and $B$ occur”.

- Similarly, since the union ($\cup$) of sets corresponds to the disjunction ($\lor$) of propositions, we often write $A \lor B$ for the event “$A$ or $B$ occurs”.

### 5.11.1 Flipping a Coin and Rolling a Die

If we flip a coin and roll a die, the sample space is

$$S = \{H1, H2, H3, H4, H5, H6, T1, T2, T3, T4, T5, T6\}.$$

The events

$$A = \text{“the coin comes up heads”}$$

and

$$B = \text{“the result of the die is even”}$$

correspond to the subsets

$$A = \{H1, H2, H3, H4, H5, H6\}$$

and

$$B = \{H2, H4, H6, T2, T4, T6\}.$$
of the sample space $S$, respectively. The event that both $A$ and $B$ occur is written as $A \cap B$ and corresponds to the subset

$$A \cap B = \{H2, H4, H6\}$$

of $S$. The event that $A$ or $B$ occurs is written as $A \cup B$ and corresponds to the subset

$$A \cup B = \{H1, H2, H3, H4, H5, H6, T2, T4, T6\}$$

of $S$.

Assume that both the coin and the die are fair and the result of rolling the die is independent of the result of flipping the coin. The probability that both $A$ and $B$ occur, i.e., $\Pr(A \cap B)$, is equal to $|A \cap B|/|S| = 3/12 = 1/4$. We can also use independence to determine this probability:

$$\Pr(A \cap B) = \Pr(A) \cdot \Pr(B) = (1/2)(3/6) = 1/4.$$  

Observe that when we determine $\Pr(A)$, we do not consider the entire sample space $S$. Instead, we consider the coin’s sample space, which is $\{H, T\}$. Similarly, when we determine $\Pr(B)$, we consider the die’s sample space, which is $\{1, 2, 3, 4, 5, 6\}$.

The probability that $A$ or $B$ occurs, i.e., $\Pr(A \cup B)$, is equal to $|A \cup B|/|S| = 9/12 = 3/4$.

### 5.11.2 Flipping Coins

Let $n \geq 1$ be an integer and assume we flip $n$ fair coins. For each $i$ with $1 \leq i \leq n$, define the event

$$A_i = \text{“the } i\text{-th coin comes up heads”}.$$ 

We assume that the coin flips are independent of each other, by which we mean that the sequence $A_1, A_2, \ldots, A_n$ of events is mutually independent. Define the event

$$A = A_1 \cap A_2 \cap \cdots \cap A_n.$$ 

What is $\Pr(A)$, i.e., the probability that all $n$ coins come up heads? Since there are $2^n$ many possible outcomes for $n$ coin flips and only one of them
5.11. Describing Events by Logical Propositions

satisfies event $A$, this probability is equal to $1/2^n$. Alternatively, we can use independence to determine $\Pr(A)$:

$$\Pr(A) = \Pr(A_1 \land A_2 \land \cdots \land A_n) = \Pr(A_1) \cdot \Pr(A_2) \cdots \Pr(A_n).$$

Since each coin is fair, we have $\Pr(A_i) = 1/2$ and, thus, we get

$$\Pr(A) = \left(\frac{1}{2}\right)\left(\frac{1}{2}\right)\cdots\left(\frac{1}{2}\right) = \frac{1}{2^n}.$$  

5.11.3 The Probability of a Circuit Failing

Consider a circuit $C$ that consists of $n$ components $C_1, C_2, \ldots, C_n$. Let $p$ be a real number with $0 < p < 1$ and assume that any component fails with probability $p$, independently of the other components. For each $i$ with $1 \leq i \leq n$, define the event

$$A_i = \text{"component } C_i \text{ fails"}.$$ 

Let $A$ be the event

$$A = \text{"the entire circuit fails"}.$$ 

- Assume that the entire circuit fails when at least one component fails. What is $\Pr(A)$, i.e., the probability that the circuit fails? By our assumption, we have

$$A = A_1 \lor A_2 \lor \cdots \lor A_n$$

and, thus, using De Morgan,

$$\overline{A} = \overline{A_1} \land \overline{A_2} \land \cdots \land \overline{A_n}.$$ 

Using independence and Lemmas 5.3.3 and 5.10.2, we get

$$\Pr(A) = 1 - \Pr(\overline{A})$$

$$= 1 - \Pr(\overline{A_1} \land \overline{A_2} \land \cdots \land \overline{A_n})$$

$$= 1 - \Pr(\overline{A_1}) \cdot \Pr(\overline{A_2}) \cdots \Pr(\overline{A_n})$$

$$= 1 - (1 - p)(1 - p)\cdots(1 - p)$$

$$n \text{ times}$$

$$= 1 - (1 - p)^n.$$
Since $0 < p < 1$, we have $\lim_{n \to \infty} \Pr(A) = 1$. We conclude that for large values of $n$, it is very likely that the circuit fails.

- Now assume that the entire circuit fails when all components fail. Again, we want to know the probability $\Pr(A)$ that the circuit fails. In this case, we have

$$A = A_1 \land A_2 \land \cdots \land A_n,$$

and we get

$$\Pr(A) = \Pr(A_1 \land A_2 \land \cdots \land A_n)$$

$$= \Pr(A_1) \cdot \Pr(A_2) \cdots \Pr(A_n)$$

$$= p \cdot p \cdots p$$

$$= p^n \text{ n times}$$

Since $0 < p < 1$, we have $\lim_{n \to \infty} \Pr(A) = 0$. Thus, for large values of $n$, it is very likely that the circuit does not fail.

5.12 Choosing a Random Element in a Linked List

Consider a linked list $L$. Each node $u$ in $L$ stores a pointer to its successor node $\text{succ}(u)$. If $u$ is the last node in $L$, then $u$ does not have a successor and $\text{succ}(u) = \text{nil}$. We are also given a pointer to the first node $\text{head}(L)$ of $L$.

\[ \text{head}(L) \rightarrow u \rightarrow \text{succ}(u) \rightarrow \text{nil} \]

Our task is to choose, uniformly at random, a node in $L$. Thus, if this list has $n$ nodes, then each node must have a probability of $1/n$ of being chosen.

We assume that we are given a function $\text{RANDOM}$: For any integer $i \geq 1$, a call to $\text{RANDOM}(i)$ returns a uniformly random element from the set $\{1, 2, \ldots, i\}$; the value returned is independent of all other calls to this function.

To make the problem interesting, we assume that we do not know the value of $n$, i.e., at the start, we do not know the number of nodes in the
list $L$. Also, we are allowed to only make one pass over this list. We will prove below that the following algorithm solves the problem:

**Algorithm** \textsc{ChooseRandomNode}(\(L\)):

\[
\begin{align*}
&u = \text{head}(L); \\
&i = 1; \\
&\textbf{while } u \neq \text{nil} \\
&\quad \textbf{do } r = \text{Random}(i); \\
&\quad \quad \textbf{if } r = 1 \\
&\quad \quad \quad \textbf{then } x = u \\
&\quad \quad \quad \textbf{endif}; \\
&\quad u = \text{succ}(u); \\
&\quad i = i + 1 \\
&\textbf{endwhile}; \\
&\textbf{return } x
\end{align*}
\]

In one iteration of the while-loop, the call to \textsc{Random}(\(i\)) returns a uniformly random element \(r\) from the set \(\{1, 2, \ldots, i\}\). If \(r = 1\), which happens with probability \(1/i\), the value of \(x\) is set to the currently visited node. If \(r \neq 1\), which happens with probability \(1 − 1/i\), the value of \(x\) does not change during this iteration of the while-loop. Thus,

- in the first iteration, \(x\) is set to the first node of \(L\) with probability 1,
- in the second iteration, \(x\) is set to the second node of \(L\) with probability \(1/2\), whereas the value of \(x\) does not change with probability \(1/2\),
- in the third iteration, \(x\) is set to the third node of \(L\) with probability \(1/3\), whereas the value of \(x\) does not change with probability \(2/3\),
- in the last iteration, \(x\) is set to the last node of \(L\) with probability \(1/|L|\), whereas the value of \(x\) does not change with probability \((|L|−1)/|L|\).

We now prove that the output \(x\) of algorithm \textsc{ChooseRandomNode}(\(L\)) is a uniformly random node of the list \(L\). Let \(n\) denote the number of nodes in \(L\) and let \(v\) be an arbitrary node in \(L\). We will prove that, after the algorithm has terminated, \(x = v\) with probability \(1/n\).

Let \(k\) be the integer such that \(v\) is the \(k\)-th node in \(L\); thus, \(1 \leq k \leq n\). We observe that, after the algorithm has terminated, \(x = v\) if and only if
• during the \( k \)-th iteration, the value of \( x \) is set to \( v \), and

• for all \( i = k + 1, k + 2, \ldots, n \), during the \( i \)-th iteration, the value of \( x \) does not change.

Define the event

\[ A = \text{“after the algorithm has terminated, } x = v \text{”}. \]

For each \( i \) with \( 1 \leq i \leq n \), define the event

\[ A_i = \text{“the value of } x \text{ changes during the } i \text{-th iteration”}. \]

Then

\[ A = A_k \land A_{k+1} \land A_{k+2} \land A_{k+3} \land \cdots \land A_n. \]

Recall that we assume that the output of the function \textsc{Random} is independent of all other calls to this function. This implies that the events \( A_1, A_2, \ldots, A_n \) are mutually independent. It follows that

\[
\Pr(A) = \Pr(A_k \land A_{k+1} \land A_{k+2} \land A_{k+3} \land \cdots \land A_n) \\
= \Pr(A_k) \cdot \Pr(A_{k+1}) \cdot \Pr(A_{k+2}) \cdot \Pr(A_{k+3}) \cdots \Pr(A_n) \\
= \frac{1}{k} \cdot \left(1 - \frac{1}{k+1}\right) \cdot \left(1 - \frac{1}{k+2}\right) \cdots \left(1 - \frac{1}{n}\right) \\
= \frac{1}{k} \cdot \frac{k+1}{k+2} \cdot \frac{k+2}{k+3} \cdots \frac{n-1}{n} \\
= \frac{1}{n}.
\]

5.13 Long Runs in Random Bitstrings

Let \( n \) be a large integer and assume we flip a fair coin \( n \) times, where all flips are mutually independent. If we write 0 for heads and 1 for tails, then we obtain a random bitstring

\[ S = s_1s_2 \ldots s_n. \]

A \textit{run of length } \( k \) \textit{is a consecutive subsequence of } \( S \), all of whose bits are the same. For example, the bitstring

\[ 0011100101000011000 \]
contains, among others, the following consecutive subsequences in bold,

\[00111100101000011000,\]

which are runs of lengths 4, 2, and 1, respectively.

Would you be surprised to see a “long” run in the random bitstring \(S\), say a run of length about \(\log n\)? Most people will answer this question with “yes”. We will prove below, however, that the correct answer is “no”: The probability that this happens is about \(1 - 1/n^2\); thus, it converges to 1 when \(n\) goes to infinity. In other words, you should be surprised if a random bitstring does not contain a run of length about \(\log n\).

We choose a positive integer \(k\) and define the event \[A = \text{“}S \text{ contains a run of length at least } k\text{”}.\]

We are going to prove a lower bound on \(\Pr(A)\) in terms of \(n\) and \(k\). At the end, we will show that by taking \(k\) to be slightly less than \(\log n\), we have \(\Pr(A) \geq 1 - 1/n^2\).

For each \(i\) with \(1 \leq i \leq n - k + 1\), we define the event \[A_i = \text{“}the subsequence of length } k\text{ starting at position } i\text{ is a run”}.\]

Since \[A = A_1 \lor A_2 \lor \cdots \lor A_{n-k+1},\]
we have \[\Pr(A) = \Pr(A_1 \lor A_2 \lor \cdots \lor A_{n-k+1}).\]

Since the probability on the right-hand side is difficult to analyze (it requires the Principle of Inclusion and Exclusion), we consider the complement of the event \(A\), i.e., the event \[\overline{A} = \text{“}each run in } S\text{ has length less than } k\text{”}.\]

Note that \[\overline{A} = \overline{A}_1 \land \overline{A}_2 \land \cdots \land \overline{A}_{n-k+1},\]
where \(\overline{A}_i\) is the complement of \(A_i\), i.e, the event \[\overline{A}_i = \text{“}the subsequence of length } k\text{ starting at position } i\text{ is not a run”}.\]
It follows that
\[
\Pr (\overline{A}) = \Pr (\overline{A}_1 \land \overline{A}_2 \land \cdots \land \overline{A}_{n-k+1}).
\] (5.4)

We determine \(\Pr (\overline{A}_i)\), by first determining \(\Pr (A_i)\). The event \(A_i\) occurs if and only if
\[
s_i = s_{i+1} = \cdots = s_{i+k-1} = 0
\]
or
\[
s_i = s_{i+1} = \cdots = s_{i+k-1} = 1.
\]
Since the coin flips are mutually independent, it follows that
\[
\Pr (A_i) = \frac{1}{2^k} + \frac{1}{2^k} = \frac{1}{2^{k-1}}
\]
and, therefore,
\[
\Pr (\overline{A}_i) = 1 - \Pr (A_i) = 1 - \frac{1}{2^{k-1}}.
\]

Is the probability on the right-hand side of (5.4) equal to the product of the individual probabilities? If the events \(\overline{A}_1, \overline{A}_2, \ldots, \overline{A}_{n-k+1}\) are mutually independent, then the answer is “yes”. However, it should be clear that, for example, the events \(\overline{A}_1\) and \(\overline{A}_2\) are not independent: If we are told that event \(A_1\) occurs, then the first \(k\) bits in the sequence \(S\) are equal; let us say they are all equal to 0. In this case, the probability that event \(A_2\) occurs is equal to the probability that the \((k+1)\)-st bit in \(S\) is 0, which is equal to \(1/2\) and not \(1/2^{k-1}\) (assuming that \(k \geq 3\)). So it seems that we are stuck.

Fortunately, there is a way out:

Let us assume that the integer \(k\) is chosen such that \(n/k\) is an integer. We divide the sequence \(S = s_1 s_2 \ldots s_n\) into \(n/k\) blocks, each having length \(k\). Thus,

- the first block is the subsequence \(s_1 s_2 \ldots s_k\),
- the second block is the subsequence \(s_{k+1} s_{k+2} \ldots s_{2k}\),
- the third block is the subsequence \(s_{2k+1} s_{2k+2} \ldots s_{3k}\),
- the \((n/k)\)-th block is the subsequence \(s_{n-k+1} s_{n-k+2} \ldots s_n\).

For each \(i\) with \(1 \leq i \leq n/k\), we define the event
\[
B_i = \text{“the } i\text{-th block is a run”}.
\]
Thus, the complement of $B_i$ is the event
$$
B_i = \text{“the } i\text{-th block is not a run”}.
$$

Since $B_i = A_{(i-1)k+1}$ and $B_i = \overline{A}_{(i-1)k+1}$, we have
$$
\Pr(B_i) = 1 - 1/2^{k-1}.
$$

Observe that

- the events $\overline{B}_1, \overline{B}_2, \ldots, \overline{B}_{n/k}$ are mutually independent, because the blocks do not overlap, and
- if the event $\overline{A}$ occurs, then the event $\overline{B}_1 \land \overline{B}_2 \land \cdots \land \overline{B}_{n/k}$ also occurs.

Using Lemma 5.3.6, it follows that
$$
\Pr(\overline{A}) \leq \Pr(\overline{B}_1 \land \overline{B}_2 \land \cdots \land \overline{B}_{n/k})
= \Pr(\overline{B}_1) \cdot \Pr(\overline{B}_2) \cdots \Pr(\overline{B}_{n/k})
= (1 - 1/2^{k-1}) \cdot (1 - 1/2^{k-1}) \cdots (1 - 1/2^{k-1})
= (1 - 1/2^{k-1})^{n/k}.
$$

Using the inequality $1 - x \leq e^{-x}$, see (5.2), we get
$$
1 - 1/2^{k-1} \leq e^{-1/2^{k-1}} = e^{-2^{k-1}}
$$

and, thus,
$$
\Pr(\overline{A}) \leq \left(e^{-2^{k-1}}\right)^{n/k} = e^{-2n/(k2^k)}. \tag{5.5}
$$

Note that until now, $k$ was arbitrary. We choose $k$ to be
$$
k = \log n - 2 \log \log n.
$$

Using basic properties of logarithms, see Section 2.4, we will show below that, for this choice of $k$, the right-hand side in (5.5) is a “nice” function of $n$.

In Section 2.4 we have seen that
$$
2^{\log n} = n
$$

and
$$
2^{2\log \log n} = \log^2 n.
$$
It follows that
\[ 2^k = 2^{\log n - 2 \log \log n} = \frac{2^{\log n}}{2^{2 \log \log n}} = \frac{n}{\log^2 n}. \]
Thus,
\[
\frac{2n}{k2^k} = \frac{2 \log^2 n}{k} \geq \frac{2 \log^2 n}{\log n - 2 \log \log n} \geq \frac{2 \log^2 n}{\log n} = 2 \log n = 2 \frac{\ln n}{\ln 2} \geq 2 \ln n,
\]
implying that
\[
\Pr(\overline{A}) \leq e^{-2n/(k2^k)} \leq e^{-2\ln n} = 1/n^2.
\]
We conclude that, for the value of \(k\) chosen above,
\[
\Pr(A) = 1 - \Pr(\overline{A}) \geq 1 - 1/n^2.
\]
Thus, with probability at least \(1 - 1/n^2\), a random bitstring of length \(n\) contains a run of length at least \(\log n - 2 \log \log n\).

We remark that we have been cheating, because we assumed that both \(k\) and \(n/k\) are integers. Assume that \(n\) is of the form \(2^{2^m}\), for some positive integer \(m\). Then both \(\log n\) and \(\log \log n\) are integers and, thus, \(k\) is an integer as well. In a correct derivation, we divide the sequence \(S\) into \([n/k]\) blocks of size \(k\) and, if \(n/k\) is not an integer, one block of length less than \(k\). We then get
\[
\Pr(\overline{A}) \leq (1 - 1/2^{k-1})^{[n/k]} \leq (e^{-2/2^k})^{[n/k]} = e^{-2 \lfloor n/k \rfloor / 2^k}.\]
In Section 5.2, we defined a sample space to be any non-empty countable set. All sample spaces that we have seen so far are finite. In some cases, infinite (but countable) sample spaces arise in a natural way. To give an example, assume we flip a fair coin repeatedly and independently until it comes up heads for the first time. The sample space $S$ is the set of sequences of all coin flips that can occur. If we denote by $T^n H$ the sequence consisting of $n$ tails followed by one heads, then

$$S = \{H, TH, TTH, TTTH, TTTTH, \ldots\} = \{T^n H : n \geq 0\},$$

This upper bound is larger than the upper bound we had before by only a small additive factor of $O((\log^2 n)/n^3)$. 

5.14 Infinite Probability Spaces

In Section 5.2, we defined a sample space to be any non-empty countable set. All sample spaces that we have seen so far are finite. In some cases, infinite (but countable) sample spaces arise in a natural way. To give an example, assume we flip a fair coin repeatedly and independently until it comes up heads for the first time. The sample space $S$ is the set of sequences of all coin flips that can occur. If we denote by $T^n H$ the sequence consisting of $n$ tails followed by one heads, then

$$S = \{H, TH, TTH, TTTH, TTTTH, \ldots\} = \{T^n H : n \geq 0\},$$
which is indeed an infinite set.

Since the coin is fair and the coin flips are mutually independent, the outcome $T^nH$ has a probability of $(1/2)^{n+1}$, i.e.,

$$\Pr(T^nH) = (1/2)^{n+1}.$$ 

Recall that according to Definition 5.2.2 in order for this to be a valid probability function, the sum of all probabilities must be equal to 1, i.e., the infinite series

$$\sum_{n=0}^{\infty} \Pr(T^nH) = \sum_{n=0}^{\infty} (1/2)^{n+1}$$

must be equal to 1. Since you may have forgotten about infinite series, we recall the definition in the following subsection.

### 5.14.1 Infinite Series

**Definition 5.14.1** Let $a_0, a_1, a_2, \ldots$ be an infinite sequence of real numbers. If

$$\lim_{N \to \infty} \sum_{n=0}^{N} a_n = \lim_{N \to \infty} (a_0 + a_1 + a_2 + \cdots + a_N)$$

exists, then we say that the infinite series $\sum_{n=0}^{\infty} a_n$ converges. In this case, the value of this infinite series is equal to

$$\sum_{n=0}^{\infty} a_n = \lim_{N \to \infty} \sum_{n=0}^{N} a_n.$$ 

For example, let $x$ be a real number with $x \neq 1$, and define $a_n = x^n$ for $n \geq 0$. We claim that

$$\sum_{n=0}^{N} a_n = \sum_{n=0}^{N} x^n = 1 + x + x^2 + \cdots + x^N = \frac{1 - x^{N+1}}{1 - x},$$

which can be proved either by induction on $N$ or by verifying that

$$(1 - x)(1 + x + x^2 + \cdots + x^N) = 1 - x^{N+1}.$$
If \(-1 < x < 1\), then \(\lim_{N \to \infty} x^{N+1} = 0\) and it follows that
\[
\sum_{n=0}^{\infty} x^n = \lim_{N \to \infty} \sum_{n=0}^{N} x^n = \lim_{N \to \infty} \frac{1 - x^{N+1}}{1 - x} = \frac{1}{1 - x}.
\]

We have proved the following result:

**Lemma 5.14.2** If \(x\) is a real number with \(-1 < x < 1\), then
\[
\sum_{n=0}^{\infty} x^n = \frac{1}{1 - x}.
\]

Now we can return to the coin flipping example that we saw in the beginning of Section 5.14. If we take \(x = 1/2\) in Lemma 5.14.2, then we get
\[
\sum_{n=0}^{\infty} \Pr(T^n H) = \sum_{n=0}^{\infty} (1/2)^n+1 = (1/2) \sum_{n=0}^{\infty} (1/2)^n = (1/2) \cdot \frac{1}{1 - 1/2} = 1.
\]

Thus, we indeed have a valid probability function on the infinite sample space \(S = \{T^n H : n \geq 0\}\).

We have seen in Lemma 5.14.2 that the infinite series \(\sum_{n=0}^{\infty} x^n\) converges if \(-1 < x < 1\). It is not difficult to see that for all other values of \(x\), the limit
\[
\lim_{N \to \infty} \sum_{n=0}^{N} x^n
\]
does not exist. As a result, if \(x \geq 1\) or \(x \leq -1\), the infinite series \(\sum_{n=0}^{\infty} x^n\) does not converge. Another example of an infinite series that does not converge is
\[
\sum_{n=1}^{\infty} 1/n = 1 + 1/2 + 1/3 + 1/4 + \cdots
\]
In Section 6.8.3 we will prove that

\[ \sum_{n=1}^{N} \frac{1}{n} = 1 + \frac{1}{2} + \frac{1}{3} + \frac{1}{4} + \cdots + \frac{1}{N} \]

is about \( \ln N \). It follows that

\[ \lim_{N \to \infty} \sum_{n=0}^{N} \frac{1}{n} \]

is about

\[ \lim_{N \to \infty} \ln N, \]

which clearly does not exist.

### 5.14.2 Who Flips the First Heads

Consider a game in which two players \( P_1 \) and \( P_2 \) take turns flipping, independently, a fair coin. Thus, first \( P_1 \) flips the coin, then \( P_2 \) flips the coin, then \( P_1 \) flips the coin, then \( P_2 \) flips the coin, etc. The player who flips heads first is the winner of the game.

Who is more likely to win this game? Our intuition says that \( P_1 \) has an advantage, because he is the player who starts: If the first flip is heads, then the game is over and \( P_1 \) wins. We will prove below that this intuition is correct: \( P_1 \) has a probability of \( \frac{2}{3} \) of winning the game and, thus, the winning probability of \( P_2 \) is only \( \frac{1}{3} \).

The sample space \( S \) is the set of sequences of all coin flips that can occur. Since the game is over as soon as a heads is flipped, we have

\[ S = \{ T^n H : n \geq 0 \}. \]

Since \( P_1 \) starts, the event

\[ A = " P_1 \text{ wins the game}" \]

corresponds to the subset

\[ A = \{ T^n H : n \geq 0 \text{ and } n \text{ is even} \}, \]

which we rewrite as

\[ A = \{ T^{2m} H : m \geq 0 \}. \]
The probability that $P_1$ wins the game is equal to

$$\Pr(A) = \sum_{m=0}^{\infty} \Pr(T^{2m}H)$$

$$= \sum_{m=0}^{\infty} \frac{1}{2}^{2m+1}$$

$$= \frac{1}{2} \sum_{m=0}^{\infty} \frac{1}{2}^{2m}$$

$$= \frac{1}{2} \sum_{m=0}^{\infty} \frac{1}{4}^m.$$ 

By taking $x = 1/4$ in Lemma 5.14.2, we get

$$\Pr(A) = \frac{1}{2} \cdot \frac{1}{1 - 1/4} = \frac{2}{3}.$$ 

Let $B$ be the event

$$B = \text{“}P_2\text{ wins the game”}.$$ 

Since either $P_1$ or $P_2$ wins the game, we have

$$\Pr(B) = 1 - \Pr(A) = 1 - \frac{2}{3} = \frac{1}{3}.$$ 

Let us verify, using an infinite series, that $\Pr(B)$ is indeed equal to $1/3$. The event $B$ corresponds to the subset

$$B = \{T^nH : n \geq 0 \text{ and } n \text{ is odd}\},$$

which we rewrite as

$$B = \{T^{2m+1}H : m \geq 0\}.$$
The probability that $P_2$ wins the game is thus equal to

$$\Pr(B) = \sum_{m=0}^{\infty} \Pr(T^{2m+1}H)$$

$$= \sum_{m=0}^{\infty} (1/2)^{2m+2}$$

$$= (1/4) \sum_{m=0}^{\infty} (1/2)^{2m}$$

$$= (1/4) \sum_{m=0}^{\infty} (1/4)^m$$

$$= (1/4) \cdot \frac{1}{1 - 1/4} = 1/3.$$

### 5.14.3 Who Flips the Second Heads

Let us change the game from the previous subsection: Again, the two players $P_1$ and $P_2$ take turns flipping, independently, a fair coin, where $P_1$ starts. The game ends as soon as a second heads comes up. The player who flips the second heads wins the game.

Before you continue reading: Who do you think has a higher probability of winning this game?

In this game, a sequence of coin flips can occur if and only if (i) the sequence contains exactly two heads and (ii) the last element in the sequence is heads. Thus, the sample space $S$ is given by

$$S = \{T^mHT^nH : m \geq 0, n \geq 0\}.$$  

The event

$$A = \text{“}P_1 \text{ wins the game”}$$

corresponds to the subset

$$A = \{T^mHT^nH : m \geq 0, n \geq 0, m + n \text{ is odd}\}.$$  

Below, we will determine $\Pr(A)$, i.e., the probability that $P_1$ wins the game.

We split the event $A$ into two events

$$A_1 = \text{“}P_1 \text{ flips both the first and the second heads”}$$
and

\[ A_2 = \text{“} P_2 \text{ flips the first heads and } P_1 \text{ flips the second heads”}. \]

If we write these two events as subsets of the sample space \( S \), we get

\[
A_1 = \{T^m H T^n H : m \geq 0, n \geq 0, m \text{ is even and } n \text{ is odd}\}
\]
\[
= \{T^{2k} H T^{2\ell+1} H : k \geq 0, \ell \geq 0\}
\]

and

\[
A_2 = \{T^m H T^n H : m \geq 0, n \geq 0, m \text{ is odd and } n \text{ is even}\}
\]
\[
= \{T^{2k+1} H T^{2\ell} H : k \geq 0, \ell \geq 0\}. 
\]

Observe that \( A_1 \cap A_2 = \emptyset \) and \( A = A_1 \cup A_2 \), implying that

\[
\Pr(A) = \Pr(A_1) + \Pr(A_2).
\]

We determine the two probabilities on the right-hand side.

We have

\[
\Pr(A_1) = \sum_{k=0}^{\infty} \sum_{\ell=0}^{\infty} \Pr(T^{2k} H T^{2\ell+1} H)
\]
\[
= \sum_{k=0}^{\infty} \sum_{\ell=0}^{\infty} \left(\frac{1}{2}\right)^{2k+2\ell+3}
\]
\[
= \left(\frac{1}{2^3}\right) \sum_{k=0}^{\infty} \left(\frac{1}{2}\right)^{2k} \sum_{\ell=0}^{\infty} \left(\frac{1}{2}\right)^{2\ell}
\]
\[
= \left(\frac{1}{8}\right) \sum_{k=0}^{\infty} \left(\frac{1}{4}\right)^k \sum_{\ell=0}^{\infty} \left(\frac{1}{4}\right)^\ell
\]
\[
= \left(\frac{1}{8}\right) \sum_{k=0}^{\infty} \left(\frac{1}{4}\right)^k \cdot \frac{1}{1 - \frac{1}{4}}
\]
\[
= \left(\frac{1}{6}\right) \sum_{k=0}^{\infty} \left(\frac{1}{4}\right)^k
\]
\[
= \left(\frac{1}{6}\right) \cdot \frac{1}{1 - \frac{1}{4}}
\]
\[
= 2/9
\]
\[ \Pr(A_2) = \sum_{k=0}^{\infty} \sum_{\ell=0}^{\infty} \Pr(T^{2k+1}HT^{2\ell}H) = \sum_{k=0}^{\infty} \sum_{\ell=0}^{\infty} \left(\frac{1}{2}\right)^{2k+2\ell+3} = \frac{2}{9}. \]

Thus, the probability that \( P_1 \) wins the game is equal to
\[ \Pr(A) = \Pr(A_1) + \Pr(A_2) = \frac{2}{9} + \frac{2}{9} = \frac{4}{9}. \]

The probability that \( P_2 \) wins the game is equal to
\[ 1 - \Pr(A) = \frac{5}{9}. \]

Thus, \( P_2 \) has a slightly larger probability of winning the game.

### 5.15 Exercises

5.1 Consider a coin that has 0 on one side and 1 on the other side. We flip this coin once and roll a die twice, and are interested in the product of the three numbers.

- What is the sample space?
- How many possible events are there?
- If both the coin and the die are fair, how would you define the probability function \( \Pr \) for this sample space?

5.2 Consider the sample space \( S = \{a, b, c, d\} \) and a probability function \( \Pr : S \to \mathbb{R} \) on \( S \). Define the events \( A = \{a\}, B = \{a, b\}, C = \{a, b, c\}, \) and \( D = \{b, d\} \). You are given that \( \Pr(A) = 1/10, \Pr(B) = 1/2, \) and \( \Pr(C) = 7/10. \) Determine \( \Pr(D) \).
5.3 Let $n$ be a positive integer. We flip a fair coin $2n$ times and consider the possible outcomes, which are strings of length $2n$ with each character being $H$ (= heads) or $T$ (= tails). Thus, we take the sample space $S$ to be the set of all such strings. Since our coin is fair, each string of $S$ should have the same probability. Thus, we define $\Pr(s) = 1/|S|$ for each string $s$ in $S$. In other words, we have a uniform probability space.

You are asked to determine the probability that in the sequence of $2n$ flips, the coin comes up heads exactly $n$ times:

- What is the event $A$ that describes this?
- Determine $\Pr(A)$.

5.4 A cup contains two pennies (P), one nickel (N), and one dime (D). You choose one coin uniformly at random, and then you choose a second coin from the remaining coins, again uniformly at random.

- Let $S$ be the sample space consisting of all ordered pairs of letters P, N, and D that represent the possible outcomes. Write out all elements of $S$.
- Determine the probability for each element in this sample space.

5.5 You flip a fair coin, independently, six times.

- What is the sample space?
- Define the events

$$A = \text{"the coin comes up heads at least four times"},$$
$$B = \text{"the number of heads is equal to the number of tails"},$$
$$C = \text{"there are at least four consecutive heads"}.$$

Determine $\Pr(A)$, $\Pr(B)$, $\Pr(C)$, $\Pr(A \mid B)$, and $\Pr(C \mid A)$.

5.6 Let $k \geq 2$ be an integer and consider the sample space $S$ consisting of all sequences of $k$ characters, where each character is one of the digits $0, 1, 2, \ldots, 9$.

If we choose a sequence $s$ uniformly at random from the sample space $S$, what is the probability that none of the digits in $s$ is equal to 5?
5.7 The Fibonacci numbers are defined as follows: \( f_0 = 0 \), \( f_1 = 1 \), and 
\( f_n = f_{n-1} + f_{n-2} \) for \( n \geq 2 \).

Let \( n \) be a large integer. A Fibonacci die is a die that has \( f_n \) faces. Such a die is fair: If we roll it, each face is on top with the same probability \( 1/f_n \).

There are three different types of Fibonacci dice:

- \( D_1 \): \( f_{n-2} \) of its faces show the number 1 and the other \( f_{n-1} \) faces show the number 4.

- \( D_2 \): Each face shows the number 3.

- \( D_3 \): \( f_{n-2} \) of its faces show the number 5 and the other \( f_{n-1} \) faces show the number 2.

Assume we roll each of \( D_1 \), \( D_2 \), and \( D_3 \) once, independently of each other. Let \( R_1 \), \( R_2 \), and \( R_3 \) be the numbers on the top face of \( D_1 \), \( D_2 \), and \( D_3 \), respectively. Determine 
\[
\Pr(R_1 > R_2)
\]
and 
\[
\Pr(R_2 > R_3),
\]
and show that 
\[
\Pr(R_3 > R_1) = \frac{f_{n-2} f_{n+1}}{f_n^2}.
\]

5.8 A group of ten people sits down, uniformly at random, around a table. Lindsay and Simon are part of this group. Determine the probability that Lindsay and Simon sit next to each other.

5.9 In Section 5.4.1, we have seen the different cards that are part of a standard deck of cards.

- You choose 2 cards uniformly at random from the 13 spades in a deck of 52 cards. Determine the probability that you choose an ace and a king.

- You choose 2 cards uniformly at random from a deck of 52 cards. Determine the probability that you choose an ace and a king.
5.15. Exercises

- You choose 2 cards uniformly at random from a deck of 52 cards. Determine the probability that you choose an ace and a king of the same suit.

5.10 You are given three events $A$, $B$, and $C$ in a sample space $S$. Is the following true or false?

$$\Pr(A \cap B \cap C) = \Pr(A \cup B \cup C) - \Pr(B) - \Pr(C) + \Pr(B \cap C).$$

5.11 Let $S$ be a set consisting of 6 positive integers and 8 negative integers. Choose a 4-element subset of $S$ uniformly at random, and multiply the elements in this subset. Denote the product by $x$. Determine the probability that $x > 0$.

5.12 Prove the inequality in [5.2], i.e., prove that

$$1 - x \leq e^{-x}$$

for all real numbers $x$.

5.13 Let $(S, \Pr)$ be a probability space and let $B$ be an event with $\Pr(B) > 0$. Define the function $\Pr': S \rightarrow \mathbb{R}$ by

$$\Pr'(\omega) = \begin{cases} \frac{\Pr(\omega)}{\Pr(B)} & \text{if } \omega \in B, \\ 0 & \text{if } \omega \notin B. \end{cases}$$

- Prove that $\Pr'$ is a probability function on $S$ according to Definition [5.2.2]

- Prove that for any event $A$,

$$\Pr'(A) = \frac{\Pr(A \cap B)}{\Pr(B)}.$$

5.14 Consider two events $A$ and $B$ in a sample space $S$.

- Assume that $\Pr(A) = 1/2$ and $\Pr(B \mid \overline{A}) = 3/5$. Determine $\Pr(A \cup B)$.

- Assume that $\Pr(A \cup B) = 5/6$ and $\Pr(\overline{A} \mid B) = 1/3$. Determine $\Pr(B)$. 
5.15 Give an example of a sample space $S$ and six events $A$, $B$, $C$, $D$, $E$, and $F$ such that

- $\Pr(A \mid B) = \Pr(A)$,
- $\Pr(C \mid D) < \Pr(C)$,
- $\Pr(E \mid F) > \Pr(E)$.

*Hint:* The sequence of six events may contain duplicates. Try to make the sample space $S$ as small as you can.

5.16 You roll two dice $D_1$ and $D_2$. Let $A$ be the event “$D_1 + D_2 = 7$” and let $B$ be the event “$D_1 = 4$”. Determine the conditional probabilities $\Pr(A \mid B)$ and $\Pr(B \mid A)$.

5.17 You flip a fair coin three times. Define the four events (recall that zero is even)

- $A =$ “the coin comes up heads an odd number of times”,
- $B =$ “the coin comes up heads an even number of times”,
- $C =$ “the coin comes up tails an odd number of times”,
- $D =$ “the coin comes up tails an even number of times”.

- Determine $\Pr(A)$, $\Pr(B)$, $\Pr(C)$, $\Pr(D)$, $\Pr(A \mid C)$, and $\Pr(A \mid D)$.
- Are there any two events in the sequence $A$, $B$, $C$, and $D$ that are independent?

5.18 A hand of 13 cards is chosen uniformly at random from a standard deck of 52 cards. Define the following events:

- $A =$ “the hand has at least one ace”,
- $B =$ “the hand has at least two aces”,
- $C =$ “the hand has the ace of spades”.

Determine the conditional probabilities $\Pr(A \mid B)$, $\Pr(B \mid A)$, and $\Pr(B \mid C)$.

5.19 We take a uniformly random permutation of a standard deck of 52 cards, so that each permutation has a probability of $1/52!$. Define the following events:
• $A = \text{“the top card is an ace”}$,
• $B = \text{“the bottom card is the ace of spades”}$,
• $C = \text{“the bottom card is the queen of spades”}$.

Determine the conditional probabilities $\Pr(A \mid B)$ and $\Pr(A \mid C)$.

5.20 Consider two dice, each one having one face showing the letter $a$, two faces showing the letter $b$, and the remaining three faces showing the letter $c$. You roll each die once, independently of the other die.

• What is the sample space?
• Define the events

$A = \text{“at least one of the two dice shows the letter } b \text{ on its top face”}$

and

$B = \text{“both dice show the same letter on their top faces”}$.

Determine $\Pr(A)$, $\Pr(B)$, and $\Pr(A \mid B)$.

5.21 You flip a fair coin, independently, three times. Define the events $A = \text{“the first flip results in heads”}$, and $B = \text{“the coin comes up heads exactly once”}$.

Determine $\Pr(A \mid B)$ and $\Pr(B \mid A)$.

5.22 You roll a fair die twice. Define the events $A = \text{“the sum of the results is even”}$ and $B = \text{“the sum of the results is at least 10”}$.

Determine $\Pr(A \mid B)$.

5.23 According to Statistics Canada, a random person in Canada has
• a probability of 4/5 to live to at least 70 years old and
• a probability of 1/2 to live to at least 80 years old.

John (a random person in Canada) has just celebrated his 70-th birthday. What is the probability that John will celebrate his 80-th birthday?

5.24 Let $A$ and $B$ be events with $\Pr(A) \neq 0$ and $\Pr(B) \neq 0$. Use the definition of conditional probability to prove Bayes’ Theorem:

$$\Pr(A \mid B) = \frac{\Pr(B \mid A) \cdot \Pr(A)}{\Pr(B)}.$$ 

5.25 Medical doctors have developed a test for detecting disease $X$.

• The test is 98% effective on people who have $X$: If a person has $X$, then with probability 0.98, the test says that the person indeed has $X$.

• The test gives a false reading for 3% of the population without the disease: If a person does not have $X$, then with probability 0.03, the test says that the person does have $X$.

• It is known that 0.1% of the population has $X$.

Assume we choose a person uniformly at random from the population and test this person for disease $X$.

• Determine the probability that the test says that the person has $X$.

• Assume the test says that the person has $X$. Use Exercise 5.24 to determine the probability that the person indeed has $X$.

5.26 Consider three events $A$, $B$, and $C$ in a sample space $S$, and assume that $\Pr(B \cap C) \neq 0$ and $\Pr(C) \neq 0$. Prove that

$$\Pr(A \cap B \cap C) = \Pr(A \mid B \cap C) \cdot \Pr(B \mid C) \cdot \Pr(C).$$

5.27 You have a fair die and do the following experiment:

• Roll the die once; let $x$ be the outcome.

• Roll the die $x$ times (independently); let $y$ be the smallest outcome of these $x$ rolls.
• Roll the die \( y \) times (independently); let \( z \) be the largest outcome of these \( y \) rolls.

Use Exercise 5.26 to determine

\[
\Pr(x = 1 \text{ and } y = 2 \text{ and } z = 3).
\]

5.28 In this exercise, we will use the product notation. In case you are not familiar with this notation:

- For \( k \leq m \), \( \prod_{i=k}^{m} x_i \) denotes the product
  \[
x_k \cdot x_{k+1} \cdot x_{k+2} \cdots x_m.
\]

- If \( k > m \), then \( \prod_{i=k}^{m} x_i \) is an “empty” product, which we define to be equal to 1.

Let \( n \geq 1 \) be an integer, and for each \( i = 1, 2, \ldots, n \), let \( p_i \) be a real number such that \( 0 < p_i < 1 \). In this exercise, you will prove that

\[
\sum_{i=1}^{n} p_i \prod_{j=i+1}^{n} (1 - p_j) = 1 - \prod_{i=1}^{n} (1 - p_i).
\] (5.6)

For example,

- for \( n = 1 \), (5.6) becomes
  \[
p_1 = 1 - (1 - p_1),
\]

- for \( n = 2 \), (5.6) becomes
  \[
p_1(1 - p_2) + p_2 = 1 - (1 - p_1)(1 - p_2),
\]

- for \( n = 3 \), (5.6) becomes
  \[
p_1(1 - p_2)(1 - p_3) + p_2(1 - p_3) + p_3 = 1 - (1 - p_1)(1 - p_2)(1 - p_3).
\]

Assume we do an experiment consisting of \( n \) tasks \( T_1, T_2, \ldots, T_n \). Each task is either a success or a failure, independently of the other tasks. For each \( i = 1, 2, \ldots, n \), let \( p_i \) be the probability that \( T_i \) is a success. Define the event

\[A = \text{“at least one task is a success”}.
\]
• Prove (5.6) by determining $\Pr(A)$ in two different ways.

5.29 A hand of 5 cards is chosen uniformly at random from a standard deck of 52 cards. Define the event

$$A = \text{“the hand has at least one ace”}.$$

• Explain what is wrong with the following argument:

We are going to determine $\Pr(A)$. Event $A$ states that the hand has at least one ace. By symmetry, we may assume that $A$ is the event that the hand has the ace of spades. Since there are $\binom{52}{5}$ hands of five cards and exactly $\binom{51}{4}$ of them contain the ace of spades, it follows that

$$\Pr(A) = \frac{\binom{51}{4}}{\binom{52}{5}} = \frac{5}{52}.$$
We are going to determine Pr(A) using the Law of Total Probability (Theorem 5.9.1). For each $x \in \{\heartsuit, \diamondsuit, \clubsuit, \spadesuit\}$, define the event

$$B_x = \text{“the hand has the ace of suit } x\text{”}.$$  

We observe that

$$\Pr(B_x) = \frac{\binom{4}{4}}{\binom{52}{5}} = \frac{5}{52}.$$  

We next observe that

$$\Pr(A \mid B_x) = 1,$$

because if event $B_x$ occurs, then event $A$ also occurs. Thus, using the Law of Total Probability, we get

$$\Pr(A) = \sum_x \Pr(A \mid B_x) \cdot \Pr(B_x)$$

$$= \sum_x 1 \cdot \Pr(B_x)$$

$$= \sum_x \frac{5}{52}$$

$$= \frac{4 \cdot 5}{52}$$

$$= \frac{5}{13}.$$  

- Determine the value of Pr(A).

5.30 You are doing two projects $P$ and $Q$. The probability that project $P$ is successful is equal to $2/3$ and the probability that project $Q$ is successful is equal to $4/5$. Whether or not these two projects are successful are independent of each other. What is the probability that both $P$ and $Q$ are not successful?

5.31 Consider two independent events $A$ and $B$ in a sample space $S$. Assume that $A$ and $B$ are disjoint, i.e., $A \cap B = \emptyset$. What can you say about Pr(A) and Pr(B)?
5.32 You flip three fair coins independently of each other. Let $A$ be the event “at least two flips in a row are heads” and let $B$ be the event “the number of heads is even”. (Note that zero is even.) Are $A$ and $B$ independent?

5.33 You flip three fair coins independently of each other. Define the events

\[ A = \text{“there is at most one tails”} \]

and

\[ B = \text{“not all flips are identical”}. \]

Are $A$ and $B$ independent?

5.34 Let $n \geq 2$ be an integer and consider two fixed integers $a$ and $b$ with $1 \leq a < b \leq n$.

- Use the Product Rule to determine the number of permutations of \{1, 2, \ldots, n\} in which $a$ is to the left of $b$.

- Consider a uniformly random permutation of the set \{1, 2, \ldots, n\}, and define the event

\[ A = \text{“in this permutation, $a$ is to the left of $b$”}. \]

Use your answer to the first part of this exercise to determine $\Pr(A)$.

5.35 Let $n \geq 3$ be an integer, consider a uniformly random permutation of the set \{1, 2, \ldots, n\}, and define the events

\[ A = \text{“in this permutation, 2 is to the left of 3”} \]

and

\[ B = \text{“in this permutation, 1 is to the left of 2 and 1 is to the left of 3”}. \]

Are these two events independent?

Hint: Use the Product Rule to determine the number of permutations that define $A$, $B$, and $A \cap B$. 
5.36 You flip two fair coins independently of each other. Define the following events:

\[ A = \text{“the number of heads is odd”}, \]
\[ B = \text{“the first coin comes up heads”}, \]
and
\[ C = \text{“the second coin comes up heads”}. \]

• Are the events \( A \) and \( B \) independent?
• Are the events \( A \) and \( C \) independent?
• Are the events \( B \) and \( C \) independent?
• Are the events \( A, B, \) and \( C \) pairwise independent?
• Are the events \( A, B, \) and \( C \) mutually independent?

5.37 You are given a tetrahedron, which is a die with four faces. Each of these faces has one of the bitstrings 110, 101, 011, and 000 written on it. Different faces have different bitstrings.

We roll the tetrahedron so that each face is at the bottom with equal probability 1/4. For \( k = 1, 2, 3 \), define the event

\[ A_k = \text{“the bitstring written on the bottom face has 0 at position } k\text{”}. \]

For example, if the bitstring at the bottom face is 101, then \( A_1 \) is false, \( A_2 \) is true, and \( A_3 \) is false.

• Are the events \( A_1 \) and \( A_2 \) independent?
• Are the events \( A_1 \) and \( A_3 \) independent?
• Are the events \( A_2 \) and \( A_3 \) independent?
• Are the events \( A_1, A_2, A_3 \) pairwise independent?
• Are the events \( A_1, A_2, A_3 \) mutually independent?

5.38 In a group of 100 children, 34 are boys and 66 are girls. You are given the following information about the girls:
• Each girl has green eyes or is blond or is left-handed.
• 20 of the girls have green eyes.
• 40 of the girls are blond.
• 50 of the girls are left-handed.
• 10 of the girls have green eyes and are blond.
• 14 of the girls have green eyes and are left-handed.
• 4 of the girls have green eyes, are blond, and are left-handed.

We choose one of these 100 children uniformly at random. Define the events
\( G = \) “the kid chosen is a girl with green eyes”,
\( B = \) “the kid chosen is a blond girl”,
and
\( L = \) “the kid chosen is a left-handed girl”.

• Are the events \( G \) and \( B \) independent?
• Are the events \( G \) and \( L \) independent?
• Are the events \( B \) and \( L \) independent?
• Verify whether or not the following equation holds:
\[
\Pr(G \land B \land L) = \Pr(G) \cdot \Pr(B) \cdot \Pr(L).
\]

5.39 Annie, Boris, and Charlie write an exam that consists of only one question: What is 26 times 26? Calculators are not allowed during the exam. Both Annie and Boris are pretty clever and each of them gives the correct answer with probability \( \frac{9}{10} \). Charlie has trouble with two-digit numbers and gives the correct answer with probability \( \frac{6}{10} \).

• Assume that the three students do not cheat, i.e., each student answers the question independently of the other two students. Determine the probability that at least two of them give the correct answer.
5.15. Exercises

- Assume that Annie and Boris do not cheat, but Charlie copies Annie’s answer. Determine the probability that at least two of them give the correct answer.

*Hint:* The answer to the second part is smaller than the answer to the first part.

5.40 Prove that for any real number \( x \neq 1 \) and any integer \( N \geq 0 \),

\[
\sum_{n=0}^{N} x^n = \frac{1 - x^{N+1}}{1 - x}.
\]

5.41 Use the following argumentation to convince yourself that

\[
\sum_{n=0}^{\infty} \frac{1}{2^n} = 2.
\]

Take the interval \( I = [0, 2) \) of length 2 on the real line and, for each \( n \geq 0 \), an interval \( I_n \) of length \( 1/2^n \). It is possible to place all intervals \( I_n \) with \( n \geq 0 \) in \( I \) such that

- no two intervals \( I_n \) and \( I_m \), with \( m \neq n \), overlap and
- all intervals \( I_n \) with \( n \geq 0 \) completely cover the interval \( I \).

5.42 Two players \( P \) and \( Q \) take turns rolling two fair and independent dice. The first player who gets a sum of seven wins the game. Determine the probability that player \( P \) wins the game.

5.43 By flipping a fair coin repeatedly and independently, we obtain a sequence of \( H \)'s and \( T \)'s. We stop flipping the coin as soon as the sequence contains either \( HH \) or \( TH \).

Two players play a game, in which Player 1 wins if the last two symbols in the sequence are \( HH \). Otherwise, the last two symbols in the sequence are \( TH \), in which case Player 2 wins. Define the events

\[ A = \text{“Player 1 wins”} \]

and

\[ B = \text{“Player 2 wins.”} \]

Determine \( \Pr(A) \) and \( \Pr(B) \).
5.44 We flip a fair coin repeatedly and independently, and stop as soon as we see one of the two sequences HTT and HHT. Let $A$ be the event that the process stops because HTT is seen.

- Prove that the event $A$ is given by the set

\[ \{ T^m(HT)^n HTT : m \geq 0, n \geq 0 \}. \]

In other words, event $A$ holds if and only if the sequence of coin flips is equal to $T^m(HT)^n HTT$ for some $m \geq 0$ and $n \geq 0$.

- Prove that $Pr(A) = 1/3$.

5.45 For $i \in \{1, 2\}$, consider the game $G_i$, in which two players $P_1$ and $P_2$ take turns flipping, independently, a fair coin, where $P_i$ starts. The game ends as soon as heads comes up. The player who flips heads first is the winner of the game $G_i$. For $j \in \{1, 2\}$, define the event

\[ B_{ij} = "P_j \text{ wins the game } G_i". \]

In Section 5.14.2, we have seen that

\[ Pr(B_{11}) = Pr(B_{22}) = 2/3 \quad (5.7) \]

and

\[ Pr(B_{12}) = Pr(B_{21}) = 1/3. \quad (5.8) \]

Consider the game $G$, in which $P_1$ and $P_2$ take turns flipping, independently, a fair coin, where $P_1$ starts. The game ends as soon as a second heads comes up. The player who flips the second heads wins the game. Define the event

\[ A = "P_1 \text{ wins the game } G". \]

In Section 5.14.3, we used an infinite series to show that

\[ Pr(A) = 4/9. \quad (5.9) \]

Use the Law of Total Probability (Theorem 5.9.1) to give an alternative proof of (5.9). You are allowed to use (5.7) and (5.8).

5.46 Consider two players $P_1$ and $P_2$: 
5.15. Exercises

- $P_1$ has one fair coin.
- $P_2$ has two coins. One of them is fair, whereas the other one is 2-headed (Her Majesty is on both sides of this coin).

The two players $P_1$ and $P_2$ play a game in which they alternate making turns: $P_1$ starts, after which it is $P_2$’s turn, after which it is $P_1$’s turn, after which it is $P_2$’s turn, etc.

- When it is $P_1$’s turn, she flips her coin once.
- When it is $P_2$’s turn, he does the following:
  - $P_2$ chooses one of his two coins uniformly at random. Then he flips the chosen coin once.
  - If the first flip did not result in heads, then $P_2$ repeats this process one more time: $P_2$ again chooses one of his two coins uniformly at random and flips the chosen coin once.

The player who flips heads first is the winner of the game.

- Determine the probability that $P_2$ wins this game, assuming that all random choices and coin flips made are mutually independent.

5.47 Jennifer loves to drink India Pale Ale (IPA), whereas Connor Hillen prefers Black IPA. Jennifer and Connor decide to go to their favorite pub Chez Lindsay et Simon. The beer menu shows that this pub has ten beers on tap:

- Phillips Cabin Fever Imperial Black IPA,
- Big Rig Black IPA,
- Leo’s Early Breakfast IPA,
- Goose Island IPA,
- Caboose IPA,
- and five other beers, neither of which is an IPA.
Each of the first five beers is an IPA, whereas each of the first two beers is a Black IPA.

Jennifer and Connor play a game, in which they alternate ordering beer: Connor starts, after which it is Jennifer’s turn, after which it is Connor’s turn, after which it is Jennifer’s turn, etc.

- When it is Connor’s turn, he orders two beers; each of these is chosen uniformly at random from the ten beers (thus, these two beers may be equal).
- When it is Jennifer’s turn, she orders one of the ten beers, uniformly at random.

The game ends as soon as (i) Connor has ordered at least one Black IPA, in which case he pays the bill, or (ii) Jennifer has ordered at least one IPA, in which case she pays the bill.

- Determine the probability that Connor pays the bill, assuming that all random choices made are mutually independent.

5.48 You would like to generate a uniformly random bit, i.e., with probability 1/2, this bit is 0, and with probability 1/2, it is 1. You find a coin in your pocket, but you are not sure if it is a fair coin: It comes up heads (H) with probability \( p \) and tails (T) with probability \( 1 - p \), for some real number \( p \) that is unknown to you. In particular, you do not know if \( p = 1/2 \). In this exercise, you will show that this coin can be used to generate a uniformly random bit.

Consider the following recursive algorithm `GetRandomBit`, which does not take any input:

```
Algorithm GetRandomBit:
    // all coin flips made are mutually independent
    flip the coin twice;
    if the result is HT
        then return 0
    else if the result is TH
        then return 1
    else GetRandomBit
    endif
endif
```
• The sample space \( S \) is the set of all sequences of coin flips that can occur when running algorithm \text{GetRandomBit}. Determine this sample space \( S \).

• Prove that algorithm \text{GetRandomBit} returns a uniformly random bit.

5.49 You would like to generate a \textit{biased} random bit: With probability \( 2/3 \), this bit is 0, and with probability \( 1/3 \), it is 1. You find a \textit{fair} coin in your pocket: This coin comes up heads (\( H \)) with probability \( 1/2 \) and tails (\( T \)) with probability \( 1/2 \). In this exercise, you will show that this coin can be used to generate a biased random bit.

Consider the following recursive algorithm \text{GetBiasedBit}, which does not take any input:

```
Algorithm \text{GetBiasedBit}:

// all coin flips made are mutually independent
flip the coin;
if the result is \( H \)
then return 0
else \( b = \) \text{GetBiasedBit};
return \( 1 - b \)
endif
```

• The sample space \( S \) is the set of all sequences of coin flips that can occur when running algorithm \text{GetBiasedBit}. Determine this sample space \( S \).

• Prove that algorithm \text{GetBiasedBit} returns 0 with probability \( 2/3 \).

5.50 Let \( n \geq 2 \) be an integer and consider a uniformly random permutation \((a_1, a_2, \ldots, a_n)\) of the set \( \{1, 2, \ldots, n\} \). For each \( k \) with \( 1 \leq k \leq n \), define the event

\[ A_k = \text{“}a_k\text{ is the largest element among } a_1, a_2, \ldots, a_k\text{”}. \]

• Let \( k \) and \( \ell \) be two integers with \( 1 \leq k < \ell \leq n \). Prove that the events \( A_k \) and \( A_\ell \) are independent.

\textbf{Hint:} Use the Product Rule to determine the number of permutations that define \( A_k, A_\ell, \text{ and } A_k \cap A_\ell \).
• Prove that the sequence $A_1, A_2, \ldots, A_n$ of events is mutually independent.

5.51 Let $n \geq 2$ be an integer. We generate a random bitstring $S = s_1 s_2 \cdots s_n$, by setting, for each $i = 1, 2, \ldots, n$, $s_i = 1$ with probability $1/i$ and, thus, $s_i = 0$ with probability $1 - 1/i$. All random choices made when setting these bits are mutually independent.

For each $i$ with $1 \leq i \leq n$, define the events

$$B_i = \text{“} s_i = 1 \text{”}$$

and

$$R_i = \text{“} \text{the rightmost 1 in the bitstring } S \text{ is at position } i \text{”}.$$

• Determine $\Pr(R_i)$.

The following algorithm TryToFindRightmostOne$(S, n, m)$ takes as input the binary string $S = s_1 s_2 \cdots s_n$ of length $n$ and an integer $m$ with $1 \leq m \leq n$. As the name suggests, this algorithm tries to find the position of the rightmost 1 in the string $S$.

**Algorithm** TryToFindRightmostOne$(S, n, m)$:

```plaintext
for i = 1 to m
    do if $s_i = 1$
        then $k = i$
        endif
endfor;
// $k$ is the position of the rightmost 1 in the substring $s_1 s_2 \cdots s_m$
// the next while-loop finds the position of the leftmost 1 in the substring $s_{m+1} s_{m+2} \cdots s_n$, if this position exists
$\ell = m + 1$;
while $\ell \leq n$ and $s_\ell = 0$
    do $\ell = \ell + 1$
endwhile;
// if $\ell \leq n$, then $\ell$ is the position of the leftmost 1 in the substring $s_{m+1} s_{m+2} \cdots s_n$
if $\ell \leq n$
    then return $\ell$
else return $k$
endif
```
Define the event

\[ E_m = \text{“there is exactly one 1 in the substring } s_{m+1}s_{m+2}\ldots s_n \text{”}. \]

• Prove that

\[ \Pr(E_m) = \frac{m}{n} \left( \frac{1}{m} + \frac{1}{m+1} + \cdots + \frac{1}{n-1} \right). \]

Define the event

\[ A = \text{“TRYToFINDRIGHTMOSTONE}(S, n, m) \text{ returns the position of the rightmost 1 in the string } S \text{”}. \]

• Prove that

\[ \Pr(A) = \frac{m}{n} \left( 1 + \frac{1}{m} + \frac{1}{m+1} + \cdots + \frac{1}{n-1} \right). \]

5.52 You realize that it is time to buy a new pair of shoes. You look up all \( n \) shoe stores in Ottawa and visit them in random order. While shopping, you create a bitstring \( s_1s_2\ldots s_n \) of length \( n \): For each \( i \) with \( 1 \leq i \leq n \), you set \( s_i \) to 1 if and only if the \( i \)-th store has the best pair of shoes, among the first \( i \) stores that you have visited.

• Use Exercise 5.50 to prove that this bitstring satisfies the condition in Exercise 5.51.

After you have visited the first \( m \) shoe stores, you are bored of shopping. You keep on visiting shoe stores, but as soon as you visit a store that has a pair of shoes that you like more than the previously best pair you have found, you buy the former pair of shoes.

• Use Exercise 5.51 to determine the probability that you buy the best pair of shoes that is available in Ottawa.
Chapter 6

Random Variables and Expectation

6.1 Random Variables

We have already seen random variables in Chapter 5, even though we did not use that term there. For example, in Section 5.2.1 we rolled a die twice and were interested in the sum of the results of these two rolls. In other words, we did an “experiment” (rolling a die twice) and asked for a function of the outcome (the sum of the results of the two rolls).

Definition 6.1.1 Let $S$ be a sample space. A random variable on the sample space $S$ is a function $X : S \rightarrow \mathbb{R}$.

In the example given above, the sample space is

$$S = \{(i, j) : 1 \leq i \leq 6, 1 \leq j \leq 6\}$$

and the random variable is the function $X : S \rightarrow \mathbb{R}$ defined by

$$X(i, j) = i + j$$

for all $(i, j)$ in $S$.

Note that the term “random variable” is misleading: A random variable is not a variable, but a function.
6.1.1 Flipping Three Coins

Assume we flip three coins. The sample space is

\[ S = \{HHH, HHT, HTH, HTT, THH, THT, TTH, TTT\}, \]

where, e.g., \( TTH \) indicates that the first two coins come up tails and the third coin comes up heads.

Let \( X : S \to \mathbb{R} \) be the random variable that maps any outcome (i.e., any element of \( S \)) to the number of heads in the outcome. Thus,

\[
\begin{align*}
X(HHH) &= 3, \\
X(HHT) &= 2, \\
X(HTH) &= 2, \\
X(HTT) &= 1, \\
X(THH) &= 2, \\
X(THT) &= 1, \\
X(TTH) &= 1, \\
X(TTT) &= 0.
\end{align*}
\]

If we define the random variable \( Y \) to be the function \( Y : S \to \mathbb{R} \) that

- maps an outcome to 1 if all three coins come up heads or all three coins come up tails, and
- maps an outcome to 0 in all other cases,

then we have

\[
\begin{align*}
Y(HHH) &= 1, \\
Y(HHT) &= 0, \\
Y(HTH) &= 0, \\
Y(HTT) &= 0, \\
Y(THH) &= 0, \\
Y(THT) &= 0, \\
Y(TTH) &= 0, \\
Y(TTT) &= 1.
\end{align*}
\]

Since a random variable is a function \( X : S \to \mathbb{R} \), it maps any outcome \( \omega \) to a real number \( X(\omega) \). Usually, we just write \( X \) instead of \( X(\omega) \). Thus, for any outcome in the sample space \( S \), we denote the value of the random
variable, for this outcome, by \( X \). In the example above, we flip three coins and write

\[ X = \text{the number of heads} \]

and

\[ Y = \begin{cases} 
1 & \text{if all three coins come up heads or all three coins come up tails,} \\
0 & \text{otherwise.}
\end{cases} \]

### 6.1.2 Random Variables and Events

Random variables give rise to events in a natural way. In the three-coin example, “\( X = 0 \)” corresponds to the event \( \{TTT\} \), whereas “\( X = 2 \)” corresponds to the event \( \{HHT, HTH, THH\} \). The table below gives some values of the random variables \( X \) and \( Y \), together with the corresponding events.

<table>
<thead>
<tr>
<th>value</th>
<th>event</th>
</tr>
</thead>
<tbody>
<tr>
<td>( X = 0 )</td>
<td>( {TTT} )</td>
</tr>
<tr>
<td>( X = 1 )</td>
<td>( {HHT, HTH, TTH} )</td>
</tr>
<tr>
<td>( X = 2 )</td>
<td>( {HHT, HTH, THH} )</td>
</tr>
<tr>
<td>( X = 3 )</td>
<td>( {HHH} )</td>
</tr>
<tr>
<td>( X = 4 )</td>
<td>( \emptyset )</td>
</tr>
<tr>
<td>( Y = 0 )</td>
<td>( {HHT, HTH, HTH, THH, THT, TTH} )</td>
</tr>
<tr>
<td>( Y = 1 )</td>
<td>( {HHH, TTT} )</td>
</tr>
<tr>
<td>( Y = 2 )</td>
<td>( \emptyset )</td>
</tr>
</tbody>
</table>

Thus, the event “\( X = x \)” corresponds to the set of all outcomes that are mapped, by the function \( X \), to the value \( x \):

**Definition 6.1.2** Let \( S \) be a sample space and let \( X : S \to \mathbb{R} \) be a random variable. For any real number \( x \), we define “\( X = x \)” to be the event \( \{\omega \in S : X(\omega) = x\} \).

Let us return to the example in which we flip three coins. Assume that the coins are fair and the three flips are mutually independent. Consider again the corresponding random variables \( X \) and \( Y \). It should be clear how we determine, for example, the probability that \( X \) is equal to 0, which we
will write as $\Pr(X = 0)$. Using our interpretation of “$X = 0$” as being the event $\{TTT\}$, we get

$$\Pr(X = 0) = \Pr(TTT) = \frac{1}{8}.$$  

Similarly, we get

$$\Pr(X = 1) = \Pr(\{HTT, THT, TTH\}) = \frac{3}{8},$$

$$\Pr(X = 2) = \Pr(\{HHT, HTH, THH\}) = \frac{3}{8},$$

$$\Pr(X = 3) = \Pr(\{HHH\}) = \frac{1}{8},$$

$$\Pr(X = 4) = \Pr(\emptyset) = 0,$$

$$\Pr(Y = 0) = \Pr(\{HHT, HTH, HTT, THH, THT, TTH\}) = \frac{6}{8} = \frac{3}{4},$$

$$\Pr(Y = 1) = \Pr(\{HHH, TTT\}) = \frac{2}{8} = \frac{1}{4},$$

$$\Pr(Y = 2) = \Pr(\emptyset) = 0.$$

Consider an arbitrary probability space $(S, \Pr)$ and let $X : S \to \mathbb{R}$ be a random variable. Using (5.1) and Definition 6.1.2, the probability of the event “$X = x$”, i.e., the probability that $X$ is equal to $x$, is equal to

$$\Pr(X = x) = \Pr(\{\omega \in S : X(\omega) = x\}) = \sum_{\omega : X(\omega) = x} \Pr(\omega).$$

We have interpreted “$X = x$” as being an event. We extend this to more general statements involving $X$. For example, “$X \geq x$” denotes the event

$$\{\omega \in S : X(\omega) \geq x\}.$$
For our three-coin example, the random variable $X$ can take each of the values 0, 1, 2, and 3 with a positive probability. As a result, “$X \geq 2$” denotes the event “$X = 2$ or $X = 3$”, and we have
\[
\Pr(X \geq 2) = \Pr(X = 2 \lor X = 3) = \Pr(X = 2) + \Pr(X = 3) = \frac{3}{8} + \frac{1}{8} = \frac{1}{2}.
\]

### 6.2 Independent Random Variables

In Section 5.10, we have defined the notion of two events being independent. The following definition extends this to random variables.

**Definition 6.2.1** Let $(S, \Pr)$ be a probability space and let $X$ and $Y$ be two random variables on $S$. We say that $X$ and $Y$ are independent if for all real numbers $x$ and $y$, the events “$X = x$” and “$Y = y$” are independent, i.e.,
\[
\Pr(X = x \land Y = y) = \Pr(X = x) \cdot \Pr(Y = y).
\]

Assume we flip three fair coins independently and, as in Section 6.1.1, consider the random variables
\[
X = \text{the number of heads}
\]
and
\[
Y = \begin{cases} 
1 & \text{if all three coins come up heads or all three coins come up tails,} \\
0 & \text{otherwise.}
\end{cases}
\]

Are these two random variables independent? Observe the following: If $Y = 1$, then $X = 0$ or $X = 3$. In other words, if we are given some information about the random variable $Y$ (in this case, $Y = 1$), then the random variable $X$ cannot take, for example, the value 2. Based on this, we take $x = 2$ and $y = 1$ in Definition 6.2.1 Since the event “$X = 2 \land Y = 1$” is equal to $\emptyset$, we have
\[
\Pr(X = 2 \land Y = 1) = \Pr(\emptyset) = 0.
\]
On the other hand, we have seen in Section 6.1.2 that $\Pr(X = 2) = \frac{3}{8}$ and $\Pr(Y = 1) = \frac{1}{4}$. It follows that

$$\Pr(X = 2 \land Y = 1) \neq \Pr(X = 2) \cdot \Pr(Y = 1)$$

and, therefore, the random variables $X$ and $Y$ are not independent.

Now consider the random variable

$$Z = \begin{cases} 1 & \text{if the first coin comes up heads,} \\ 0 & \text{if the first coin comes up tails.} \end{cases}$$

We claim that the random variables $Y$ and $Z$ are independent. To verify this, we have to show that for all real numbers $y$ and $z$,

$$\Pr(Y = y \land Z = z) = \Pr(Y = y) \cdot \Pr(Z = z). \tag{6.1}$$

Recall from Section 6.1.2 that $\Pr(Y = 1) = \frac{1}{4}$ and $\Pr(Y = 0) = \frac{3}{4}$. Since the coin flips are independent, we have $\Pr(Z = 1) = \frac{1}{2}$ and $\Pr(Z = 0) = \frac{1}{2}$. Furthermore,

$$\begin{align*}
\Pr(Y = 1 \land Z = 1) &= \Pr(HHH) \\
&= \frac{1}{8}, \\
\Pr(Y = 1 \land Z = 0) &= \Pr(TTT) \\
&= \frac{1}{8}, \\
\Pr(Y = 0 \land Z = 1) &= \Pr(HHT, HTH, HTT) \\
&= \frac{3}{8}, \\
\Pr(Y = 0 \land Z = 0) &= \Pr(THH, THT, TTH) \\
&= \frac{3}{8}.
\end{align*}$$

It follows that

$$\begin{align*}
\Pr(Y = 1 \land Z = 1) &= \Pr(Y = 1) \cdot \Pr(Z = 1), \\
\Pr(Y = 1 \land Z = 0) &= \Pr(Y = 1) \cdot \Pr(Z = 0), \\
\Pr(Y = 0 \land Z = 1) &= \Pr(Y = 0) \cdot \Pr(Z = 1), \\
\text{and} \\
\Pr(Y = 0 \land Z = 0) &= \Pr(Y = 0) \cdot \Pr(Z = 0).
\end{align*}$$
6.3. Distribution Functions

Thus, (6.1) holds if \((y, z) \in \{(1, 1), (1, 0), (0, 1), (0, 0)\}\). For any other pair \((y, z)\), such as \((y, z) = (3, 5)\) or \((y, z) = (1, 2)\), at least one of the events “\(Y = y\)” and “\(Z = z\)” is the empty set, i.e., cannot occur. Therefore, for such pairs, we have

\[
\Pr(Y = y \land Z = z) = 0 = \Pr(Y = y) \cdot \Pr(Z = z).
\]

Thus, we have indeed verified that (6.1) holds for all real numbers \(y\) and \(z\). As a result, we have shown that the random variables \(Y\) and \(Z\) are independent.

Are the random variables \(X\) and \(Z\) independent? If \(X = 0\), then all three coins come up tails and, therefore, \(Z = 0\). Thus,

\[
\Pr(X = 0 \land Z = 1) = \Pr(\emptyset) = 0,
\]

whereas

\[
\Pr(X = 0) \cdot \Pr(Z = 1) = 1/8 \cdot 1/2 = 0.
\]

As a result, the random variables \(X\) and \(Z\) are not independent.

We have defined the notion of two random variables being independent. As in Definition 5.10.3, there are two ways to generalize this to sequences of random variables:

**Definition 6.2.2** Let \((S, \Pr)\) be a probability space, let \(n \geq 2\), and let \(X_1, X_2, \ldots, X_n\) be a sequence of random variables on \(S\).

1. We say that this sequence is **pairwise independent** if for all real numbers \(x_1, x_2, \ldots, x_n\), the sequence “\(X_1 = x_1\)”, “\(X_2 = x_2\)”, \ldots, “\(X_n = x_n\)” of events is pairwise independent.

2. We say that this sequence is **mutually independent** if for all real numbers \(x_1, x_2, \ldots, x_n\), the sequence “\(X_1 = x_1\)”, “\(X_2 = x_2\)”, \ldots, “\(X_n = x_n\)” of events is mutually independent.

### 6.3 Distribution Functions

Consider a random variable \(X\) on a sample space \(S\). In Section 6.1.2, we have defined \(\Pr(X = x)\), i.e., the probability of the event “\(X = x\)”, to be

\[
\Pr(X = x) = \Pr(\{\omega \in S : X(\omega) = x\}).
\]
This defines a function that maps any real number \( x \) to the real number \( \Pr(X = x) \). This function is called the distribution function of the random variable \( X \):

**Definition 6.3.1** Let \((S, \Pr)\) be a probability space and let \( X : S \rightarrow \mathbb{R} \) be a random variable. The *distribution function* of \( X \) is the function \( D : \mathbb{R} \rightarrow \mathbb{R} \) defined by

\[
D(x) = \Pr(X = x)
\]

for all \( x \in \mathbb{R} \).

For example, consider a fair red die and a fair blue die, and assume we roll them independently. The sample space is

\[
S = \{(i, j) : 1 \leq i \leq 6, 1 \leq j \leq 6\},
\]

where \( i \) is the result of the red die and \( j \) is the result of the blue die. Each outcome \((i, j)\) in \( S \) has the same probability of \( 1/36 \).

Let \( X \) be the random variable whose value is equal to the sum of the results of the two dies. The matrix below gives all possible values of \( X \). The leftmost column gives the result of the red die, the top row gives the result of the blue die, and each other entry is the corresponding value of \( X \).
As can be seen from this matrix, the random variable $X$ can take any value in $\{2, 3, 4, \ldots, 12\}$. The distribution function $D$ of $X$ is given by

\[
\begin{align*}
D(2) &= \Pr(X = 2) = 1/36, \\
D(3) &= \Pr(X = 3) = 2/36, \\
D(4) &= \Pr(X = 4) = 3/36, \\
D(5) &= \Pr(X = 5) = 4/36, \\
D(6) &= \Pr(X = 6) = 5/36, \\
D(7) &= \Pr(X = 7) = 6/36, \\
D(8) &= \Pr(X = 8) = 5/36, \\
D(9) &= \Pr(X = 9) = 4/36, \\
D(10) &= \Pr(X = 10) = 3/36, \\
D(11) &= \Pr(X = 11) = 2/36, \\
D(12) &= \Pr(X = 12) = 1/36,
\end{align*}
\]

whereas for all $x \not\in \{2, 3, 4, \ldots, 12\}$,

\[D(x) = \Pr(X = x) = 0.\]

In Sections 6.6 and 6.7 we will see other examples of distribution functions.

6.4 Expected Values

Consider the probability space $(S, \Pr)$ with sample space $S = \{1, 2, 3\}$ and probability function $\Pr$ defined by $\Pr(1) = 4/5$, $\Pr(2) = 1/10$, and $\Pr(3) = 1/10$. Assume we choose an element in $S$ according to this probability function. Let $X$ be the random variable whose value is equal to the element in $S$ that is chosen. Thus, as a function $X : S \to \mathbb{R}$, we have $X(1) = 1$, $X(2) = 2$, and $X(3) = 3$.

The “expected value” of $X$ is the value of $X$ that we observe “on average”. How should we define this? Since $X$ has a much higher probability to take the value 1 than the other two values 2 and 3, the value 1 should get a larger “weight” in the expected value of $X$. Based on this, it is natural to define the expected value of $X$ to be

\[
1 \cdot \Pr(1) + 2 \cdot \Pr(2) + 3 \cdot \Pr(3) = 1 \cdot \frac{4}{5} + 2 \cdot \frac{1}{10} + 3 \cdot \frac{1}{10} = \frac{13}{10}.
\]
Definition 6.4.1 Let \((S, \Pr)\) be a probability space and let \(X : S \to \mathbb{R}\) be a random variable. The expected value of \(X\) is defined to be

\[
E(X) = \sum_{\omega \in S} X(\omega) \cdot \Pr(\omega),
\]

provided this summation converges absolutely.\(^1\)

6.4.1 Some Examples

Flipping a coin: Assume we flip a fair coin, in which case the sample space is \(S = \{H, T\}\) and \(\Pr(H) = \Pr(T) = 1/2\). Define the random variable \(X\) to have value

\[
X = \begin{cases} 
1 & \text{if the coin comes up heads,} \\
0 & \text{if the coin comes up tails.}
\end{cases}
\]

Thus, as a function \(X : S \to \mathbb{R}\), we have \(X(H) = 1\) and \(X(T) = 0\). The expected value \(E(X)\) of \(X\) is equal to

\[
E(X) = X(H) \cdot \Pr(H) + X(T) \cdot \Pr(T) \\
= 1 \cdot \frac{1}{2} + 0 \cdot \frac{1}{2} \\
= \frac{1}{2}.
\]

This example shows that the term “expected value” is a bit misleading: \(E(X)\) is not the value that we expect to observe, because the value of \(X\) is never equal to its expected value.

Rolling a die: Assume we roll a fair die. Define the random variable \(X\) to be the value of the result. Then, \(X\) takes each of the values in \(\{1, 2, 3, 4, 5, 6\}\) with equal probability \(1/6\), and we get

\[
E(X) = 1 \cdot \frac{1}{6} + 2 \cdot \frac{1}{6} + 3 \cdot \frac{1}{6} + 4 \cdot \frac{1}{6} + 5 \cdot \frac{1}{6} + 6 \cdot \frac{1}{6} \\
= \frac{7}{2}.
\]

\(^1\)The series \(\sum_{n=0}^{\infty} a_n\) converges absolutely if the series \(\sum_{n=0}^{\infty} |a_n|\) converges. If a series converges absolutely, then we can change the order of summation without changing the value of the series.
Now define the random variable $Y$ to be equal to one divided by the result of the die. In other words, $Y = 1/X$. This random variable takes each of the values in \{1, 1/2, 1/3, 1/4, 1/5, 1/6\} with equal probability $1/6$, and we get

\[
\mathbb{E}(Y) = \frac{1}{6} \cdot \frac{1}{1} + \frac{1}{2} \cdot \frac{1}{6} + \frac{1}{3} \cdot \frac{1}{6} + \frac{1}{4} \cdot \frac{1}{6} + \frac{1}{5} \cdot \frac{1}{6} + \frac{1}{6} \cdot \frac{1}{6} = \frac{49}{120}.
\]

Note that $\mathbb{E}(Y) \neq 1/\mathbb{E}(X)$. Thus, this example shows that, in general, $\mathbb{E}(1/X) \neq 1/\mathbb{E}(X)$.

**Rolling two dice:** Consider a fair red die and a fair blue die, and assume we roll them independently. The sample space is

\[S = \{(i,j) : 1 \leq i \leq 6, 1 \leq j \leq 6\},\]

where $i$ is the result of the red die and $j$ is the result of the blue die. Each outcome $(i,j)$ in $S$ has the same probability of $1/36$.

Let $X$ be the random variable whose value is equal to the sum of the results of the two rolls. As a function $X : S \to \mathbb{R}$, we have $X(i,j) = i + j$. The matrix below gives all possible values of $X$. The leftmost column indicates the result of the red die, the top row indicates the result of the blue die, and each other entry is the corresponding value of $X$.

<table>
<thead>
<tr>
<th></th>
<th>1</th>
<th>2</th>
<th>3</th>
<th>4</th>
<th>5</th>
<th>6</th>
</tr>
</thead>
<tbody>
<tr>
<td>1</td>
<td>2</td>
<td>3</td>
<td>4</td>
<td>5</td>
<td>6</td>
<td>7</td>
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<tr>
<td>5</td>
<td>6</td>
<td>7</td>
<td>8</td>
<td>9</td>
<td>10</td>
<td>11</td>
</tr>
<tr>
<td>6</td>
<td>7</td>
<td>8</td>
<td>9</td>
<td>10</td>
<td>11</td>
<td>12</td>
</tr>
</tbody>
</table>
Chapter 6. Random Variables and Expectation

The expected value \( \mathbb{E}(X) \) of \( X \) is equal to

\[
\mathbb{E}(X) = \sum_{(i,j) \in S} X(i, j) \cdot \Pr(i, j)
\]

\[
= \sum_{(i,j) \in S} (i + j) \cdot \frac{1}{36}
\]

\[
= \frac{1}{36} \sum_{(i,j) \in S} (i + j)
\]

\[
= \frac{1}{36} \cdot \text{the sum of all 36 entries in the matrix}
\]

\[
= \frac{1}{36} \cdot 252
\]

\[
= 7.
\]

6.4.2 Comparing the Expected Values of Comparable Random Variables

Consider a probability space \((S, \Pr)\), and let \(X\) and \(Y\) be two random variables on \(S\). Recall that \(X\) and \(Y\) are functions that map elements of \(S\) to real numbers. We will write \(X \leq Y\), if for each element \(\omega \in S\), we have \(X(\omega) \leq Y(\omega)\). In other words, the value of \(X\) is at most the value of \(Y\), no matter which outcome \(\omega\) is chosen. The following lemma should not be surprising:

**Lemma 6.4.2** Let \((S, \Pr)\) be a probability space and let \(X\) and \(Y\) be two random variables on \(S\). If \(X \leq Y\), then \(\mathbb{E}(X) \leq \mathbb{E}(Y)\).

**Proof.** Using Definition 6.4.1 and the assumption that \(X \leq Y\), we obtain

\[
\mathbb{E}(X) = \sum_{\omega \in S} X(\omega) \cdot \Pr(\omega)
\]

\[
\leq \sum_{\omega \in S} Y(\omega) \cdot \Pr(\omega)
\]

\[
= \mathbb{E}(Y).
\]
6.4.3 An Alternative Formula for the Expected Value

In the last example, we used Definition 6.4.1 to compute the expected value $\mathbb{E}(X)$ of the random variable $X$ that was defined to be the sum of the results when rolling two dice. This was a painful way to compute $\mathbb{E}(X)$, because we added all 36 entries in the matrix. There is a slightly easier way to determine $\mathbb{E}(X)$: By looking at the matrix, we see, for example, that there is 1 way for the event “$X = 2$” to occur, whereas there are 2 ways for the event “$X = 3$” to occur. The table below lists the number of ways that each possible event can occur.

<table>
<thead>
<tr>
<th>event</th>
<th>number of ways</th>
</tr>
</thead>
<tbody>
<tr>
<td>$X = 2$</td>
<td>1</td>
</tr>
<tr>
<td>$X = 3$</td>
<td>2</td>
</tr>
<tr>
<td>$X = 4$</td>
<td>3</td>
</tr>
<tr>
<td>$X = 5$</td>
<td>4</td>
</tr>
<tr>
<td>$X = 6$</td>
<td>5</td>
</tr>
<tr>
<td>$X = 7$</td>
<td>6</td>
</tr>
<tr>
<td>$X = 8$</td>
<td>5</td>
</tr>
<tr>
<td>$X = 9$</td>
<td>4</td>
</tr>
<tr>
<td>$X = 10$</td>
<td>3</td>
</tr>
<tr>
<td>$X = 11$</td>
<td>2</td>
</tr>
<tr>
<td>$X = 12$</td>
<td>1</td>
</tr>
</tbody>
</table>

Based on this, we get

$$
\mathbb{E}(X) = 2 \cdot \frac{1}{36} + 3 \cdot \frac{2}{36} + 4 \cdot \frac{3}{36} + 5 \cdot \frac{4}{36} + 6 \cdot \frac{5}{36} + 7 \cdot \frac{6}{36} + 8 \cdot \frac{5}{36} + 9 \cdot \frac{4}{36} + 10 \cdot \frac{3}{36} + 11 \cdot \frac{2}{36} + 12 \cdot \frac{1}{36}
$$

$$
= 7.
$$

Even though this is still quite painful, less computation is needed. What we have done is the following: In the definition of $\mathbb{E}(X)$, i.e.,

$$
\mathbb{E}(X) = \sum_{(i,j) \in S} X(i, j) \cdot \Pr(i, j),
$$

we rearranged the terms in the summation. That is, instead of taking the sum over all elements $(i, j)$ in $S$, we
• grouped together all outcomes \((i, j)\) for which \(X(i, j) = i + j\) has the same value, say, \(k\),

• multiplied this common value \(k\) by the probability that \(X\) is equal to \(k\),

• and took the sum of the resulting products over all possible values of \(k\).

This resulted in

\[
E(X) = \sum_{k=2}^{12} k \cdot \Pr(X = k).
\]

The following lemma states that this can be done for any random variable.

**Lemma 6.4.3** Let \((S, \Pr)\) be a probability space and let \(X : S \rightarrow \mathbb{R}\) be a random variable. The expected value of \(X\) is equal to

\[
E(X) = \sum_x x \cdot \Pr(X = x).
\]

**Proof.** Recall that the event \("X = x"\) corresponds to the subset

\[
A_x = \{\omega \in S : X(\omega) = x\}
\]

of the sample space \(S\). We have

\[
E(X) = \sum_{\omega \in S} X(\omega) \cdot \Pr(\omega)
\]

\[
= \sum_x \sum_{\omega : X(\omega) = x} X(\omega) \cdot \Pr(\omega)
\]

\[
= \sum_x \sum_{\omega : X(\omega) = x} x \cdot \Pr(\omega)
\]

\[
= \sum_x \sum_{\omega \in A_x} x \cdot \Pr(\omega)
\]

\[
= \sum_x x \sum_{\omega \in A_x} \Pr(\omega)
\]

\[
= \sum_x x \cdot \Pr(A_x)
\]

\[
= \sum_x x \cdot \Pr(X = x).
\]
6.5. Linearity of Expectation

When determining the expected value of a random variable $X$, it is usually easier to use Lemma 6.4.3 than Definition 6.4.1. To use Lemma 6.4.3 you have to do the following:

- Determine all values $x$ that $X$ can take.
- For each such value $x$, determine $\Pr(X = x)$.
- Compute the sum of all products $x \cdot \Pr(X = x)$.

6.5 Linearity of Expectation

We now come to one of the most useful tools for determining expected values:

**Theorem 6.5.1** Let $(S, \Pr)$ be a probability space. For any two random variables $X$ and $Y$ on $S$, and for any two real numbers $a$ and $b$,

$$E(aX + bY) = a \cdot E(X) + b \cdot E(Y).$$

**Proof.** Recall that both $X$ and $Y$ are functions from $S$ to $\mathbb{R}$. Define the random variable $Z$ to be $Z = aX + bY$. That is, as a function $Z : S \to \mathbb{R}$, $Z$ is defined by

$$Z(\omega) = a \cdot X(\omega) + b \cdot Y(\omega)$$

for all $\omega$ in $S$. Using Definition 6.4.1, we get

$$E(Z) = \sum_{\omega \in S} Z(\omega) \cdot \Pr(\omega)$$

$$= \sum_{\omega \in S} (a \cdot X(\omega) + b \cdot Y(\omega)) \cdot \Pr(\omega)$$

$$= a \sum_{\omega \in S} X(\omega) \cdot \Pr(\omega) + b \sum_{\omega \in S} Y(\omega) \cdot \Pr(\omega)$$

$$= a \cdot E(X) + b \cdot E(Y).$$

Let us return to the example in which we roll two independent fair dice, one being red and the other being blue. Define the random variable $X$ to be
the sum of the results of the two rolls. We have seen two ways to compute the expected value $E(X)$ of $X$. We now present a third way, which is the easiest one: We define two random variables

\[ Y = \text{the result of the red die} \]

and

\[ Z = \text{the result of the blue die}. \]

We have already seen that

\[ E(Y) = 1 \cdot \frac{1}{6} + 2 \cdot \frac{1}{6} + 3 \cdot \frac{1}{6} + 4 \cdot \frac{1}{6} + 5 \cdot \frac{1}{6} + 6 \cdot \frac{1}{6} = \frac{7}{2}. \]

Of course, by the same computation, we have

\[ E(Z) = \frac{7}{2}. \]

Observe that

\[ X = Y + Z. \]

Then, by the Linearity of Expectation (i.e., Theorem 6.5.1), we have

\[
E(X) = E(Y + Z) = E(Y) + E(Z) = \frac{7}{2} + \frac{7}{2} = \frac{7}{2}.
\]

We have stated the Linearity of Expectation for two random variables. The proof of Theorem 6.5.1 can easily be generalized to any finite sequence of random variables:

**Theorem 6.5.2** Let $(S, \Pr)$ be a probability space, let $n \geq 2$ be an integer, let $X_1, X_2, \ldots, X_n$ be a sequence of random variables on $S$, and let $a_1, a_2, \ldots, a_n$ be a sequence of real numbers. Then,

\[
E \left( \sum_{i=1}^{n} a_iX_i \right) = \sum_{i=1}^{n} a_i \cdot E(X_i).
\]
The following theorem states that the Linearity of Expectation also holds for infinite sequences of random variables:

**Theorem 6.5.3** Let \( (S, \Pr) \) be a probability space and let \( X_1, X_2, \ldots \) be an infinite sequence of random variables on \( S \) such that the infinite series

\[
\sum_{i=1}^{\infty} \mathbb{E}(|X_i|)
\]

converges. Then,

\[
\mathbb{E}\left( \sum_{i=1}^{\infty} X_i \right) = \sum_{i=1}^{\infty} \mathbb{E}(X_i).
\]

**Proof.** Define the random variable \( X \) to be

\[
X = \sum_{i=1}^{\infty} X_i.
\]

The derivation below uses Definition 6.4.1 and the assumption that the infinite series \( \sum_{i=1}^{\infty} \mathbb{E}(|X_i|) \) converges, which allows us to change the order of summation without changing the value of the series:

\[
\sum_{i=1}^{\infty} \mathbb{E}(X_i) = \sum_{i=1}^{\infty} \sum_{\omega \in S} X_i(\omega) \cdot \Pr(\omega)
\]

\[
= \sum_{\omega \in S} \sum_{i=1}^{\infty} X_i(\omega) \cdot \Pr(\omega)
\]

\[
= \sum_{\omega \in S} \Pr(\omega) \sum_{i=1}^{\infty} X_i(\omega)
\]

\[
= \sum_{\omega \in S} \Pr(\omega) \cdot X(\omega)
\]

\[
= \mathbb{E}(X)
\]

\[
= \mathbb{E}\left( \sum_{i=1}^{\infty} X_i \right).
\]
6.6 The Geometric Distribution

Let \( p \) be a real number with \( 0 < p < 1 \) and consider an experiment that is \textit{successful} with probability \( p \) and \textit{fails} with probability \( 1 - p \). We repeat this experiment independently until it is successful for the first time. What is the expected number of times that we perform the experiment?

We model this problem in the following way: Assume we have a coin that comes up heads with probability \( p \) and, thus, comes up tails with probability \( 1 - p \). We flip this coin repeatedly and independently until it comes up heads for the first time. (We have seen this process in Section 5.14 for the case when \( p = 1/2 \).) Define the random variable \( X \) to be the number of times that we flip the coin; this includes the last coin flip, which resulted in heads. We want to determine the expected value \( \mathbb{E}(X) \) of \( X \).

The sample space is given by

\[
S = \{T^{k-1}H : k \geq 1\},
\]

where \( T^{k-1}H \) denotes the sequence consisting of \( k - 1 \) tails followed by one heads. Since the coin flips are independent, the outcome \( T^{k-1}H \) has a probability of \( (1 - p)^{k-1}p = p(1 - p)^{k-1} \), i.e.,

\[
\Pr(T^{k-1}H) = p(1 - p)^{k-1}.
\]

Let us first verify that all probabilities add up to 1: Using Lemma 5.14.2, we have

\[
\sum_{k=1}^{\infty} \Pr(T^{k-1}H) = \sum_{k=1}^{\infty} p(1 - p)^{k-1} = p \sum_{k=1}^{\infty} (1 - p)^{k-1} = p \sum_{\ell=0}^{\infty} (1 - p)^{\ell} = p \frac{1}{1 - (1 - p)} = 1.
\]
6.6.1 Determining the Expected Value

We are going to use Lemma 6.4.3 to determine the expected value $E(X)$. We first observe that $X$ can take any value in $\{1, 2, 3, \ldots\}$. For any integer $k \geq 1$, $X = k$ if and only if the coin flips give the sequence $T^{k-1}H$. It follows that

$$\Pr(X = k) = \Pr(T^{k-1}H) = p(1 - p)^{k-1}. \quad (6.2)$$

By Lemma 6.4.3, we have

$$E(X) = \sum_{k=1}^{\infty} k \cdot Pr(X = k) = \sum_{k=1}^{\infty} kp(1 - p)^{k-1} = p \sum_{k=1}^{\infty} k(1 - p)^{k-1}.$$ 

How do we determine the infinite series on the right-hand side?

According to Lemma 5.14.2, we have

$$\sum_{k=0}^{\infty} x^k = \frac{1}{1-x},$$

for any real number $x$ with $-1 < x < 1$. Both sides of this equation are functions of $x$ and these two functions are equal to each other. If we differentiate both sides, we get two derivatives that are also equal to each other:

$$\sum_{k=0}^{\infty} kx^{k-1} = \frac{1}{(1-x)^2}.$$ 

Since for $k = 0$, the term $kx^{k-1}$ is equal to 0, we have

$$\sum_{k=1}^{\infty} kx^{k-1} = \frac{1}{(1-x)^2}.$$
If we take \( x = 1 - p \), we get

\[
E(X) = p \sum_{k=1}^{\infty} k(1-p)^{k-1} \\
= p \cdot \frac{1}{(1-(1-p))^2} \\
= \frac{p}{p^2} \\
= \frac{1}{p}.
\]

In Section 6.3, we have defined the distribution function of a random variable. The distribution function of the coin-flipping random variable \( X \) is given by (6.2). This function is called a geometric distribution:

**Definition 6.6.1** Let \( p \) be a real number with \( 0 < p < 1 \). A random variable \( X \) has a geometric distribution with parameter \( p \), if its distribution function satisfies

\[
\Pr(X = k) = p(1-p)^{k-1}
\]

for any integer \( k \geq 1 \).

Our calculation that led to the value of \( E(X) \) proves the following theorem:

**Theorem 6.6.2** Let \( p \) be a real number with \( 0 < p < 1 \) and let \( X \) be a random variable that has a geometric distribution with parameter \( p \). Then

\[
E(X) = 1/p.
\]

For example, if we flip a fair coin (in which case \( p = 1/2 \)) repeatedly and independently until it comes up heads for the first time, then the expected number of coin flips is equal to 2.

### 6.7 The Binomial Distribution

As in Section 6.6, we consider an experiment that is successful with probability \( p \) and fails with probability \( 1 - p \). For an integer \( n \geq 1 \), we repeat the
6.7. The Binomial Distribution

experiment, independently, \(n\) times. What is the expected number of times that the experiment is successful?

We again model this problem using a coin that comes up heads with probability \(p\) and, thus, comes up tails with probability \(1-p\). We flip the coin, independently, \(n\) times and define the random variable \(X\) to be the number of times the coin comes up heads. We want to determine the expected value \(E(X)\) of \(X\).

Since our coin comes up heads with probability \(p\), it is reasonable to guess that \(E(X) = np\). For example, if \(p = 1/2\), then, on average, half of the coin flips should come up heads. We will prove below that \(E(X) = np\).

6.7.1 Determining the Expected Value

Since the random variable \(X\) can take any value in \(\{0,1,2,\ldots,n\}\), we have, by Lemma 6.4.3,

\[
E(X) = \sum_{k=0}^{n} k \cdot \Pr(X = k).
\]

Thus, we have to determine \(\Pr(X = k)\), i.e., the probability that in a sequence of \(n\) independent coin flips, the coin comes up heads exactly \(k\) times.

To give an example, assume that \(n = 4\) and \(k = 2\). The table below gives all \(\binom{4}{2} = 6\) sequences of 4 coin flips that contain exactly 2 \(H\)’s, together with their probabilities:

<table>
<thead>
<tr>
<th>sequence</th>
<th>probability</th>
</tr>
</thead>
<tbody>
<tr>
<td>HHTT</td>
<td>(p \cdot p \cdot (1-p) \cdot (1-p) = p^2(1-p)^2)</td>
</tr>
<tr>
<td>HTHT</td>
<td>(p \cdot (1-p) \cdot p \cdot (1-p) = p^2(1-p)^2)</td>
</tr>
<tr>
<td>HTTH</td>
<td>(p \cdot (1-p) \cdot (1-p) \cdot p = p^2(1-p)^2)</td>
</tr>
<tr>
<td>THHT</td>
<td>((1-p) \cdot p \cdot p \cdot (1-p) = p^2(1-p)^2)</td>
</tr>
<tr>
<td>THTH</td>
<td>((1-p) \cdot p \cdot (1-p) \cdot p = p^2(1-p)^2)</td>
</tr>
<tr>
<td>TTHH</td>
<td>((1-p) \cdot (1-p) \cdot p \cdot p = p^2(1-p)^2)</td>
</tr>
</tbody>
</table>

As can be seen from this table, each of the \(\binom{4}{2}\) sequences has the same probability \(p^2(1-p)^2\). It follows that, if \(n = 4\),

\[
\Pr(X = 2) = \binom{4}{2} p^2(1-p)^2.
\]
We now consider the general case. Let \( n \geq 1 \) and \( k \) be integers with \( 0 \leq k \leq n \). Then, \( X = k \) if and only if there are exactly \( k \) \( H \)'s in the sequence of \( n \) coin flips. The number of such sequences is equal to \( \binom{n}{k} \), and each one of them has probability \( p^k(1-p)^{n-k} \). Therefore, we have

\[
Pr(X = k) = \binom{n}{k} p^k(1-p)^{n-k}.
\]

(6.3)

As a sanity check, let us use Newton’s Binomial Theorem (i.e., Theorem 3.6.5) to verify that all probabilities add up to 1:

\[
\sum_{k=0}^{n} Pr(X = k) = \sum_{k=0}^{n} \binom{n}{k} p^k(1-p)^{n-k}
= ((1-p) + p)^n
= 1.
\]

We are now ready to compute the expected value of the random variable \( X \):

\[
E(X) = \sum_{k=0}^{n} k \cdot Pr(X = k)
= \sum_{k=0}^{n} k \binom{n}{k} p^k(1-p)^{n-k}
= \sum_{k=1}^{n} k \binom{n}{k} p^k(1-p)^{n-k}.
\]

Since

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

we get

\[
E(X) = \sum_{k=1}^{n} \binom{n-1}{k-1} p^k(1-p)^{n-k}.
\]
By changing the summation variable from \( k \) to \( \ell + 1 \), we get
\[
E(X) = \sum_{\ell=0}^{n-1} n \binom{n-1}{\ell} p^{\ell+1}(1-p)^{n-1-\ell} = pn \sum_{\ell=0}^{n-1} \binom{n-1}{\ell} p^{\ell}(1-p)^{n-1-\ell}.
\]

By Newton’s Binomial Theorem (i.e., Theorem 3.6.5), the summation is equal to
\[
((1-p) + p)^{n-1} = 1.
\]
Therefore, we get
\[
E(X) = pn \cdot 1 = pn.
\]

We have done the following: Our intuition told us that \( E(X) = pn \). Then, we went through a painful calculation to show that our intuition was correct. There must be an easier way to show that \( E(X) = pn \). We will show below that there is indeed a much easier way.

### 6.7.2 Using the Linearity of Expectation

We define a sequence \( X_1, X_2, \ldots, X_n \) of random variables as follows: For each \( i \) with \( 1 \leq i \leq n \),
\[
X_i = \begin{cases} 
1 & \text{if the } i\text{-th coin flip results in heads,} \\
0 & \text{if the } i\text{-th coin flip results in tails.} 
\end{cases}
\]

Observe that
\[
X = X_1 + X_2 + \cdots + X_n,
\]
because
- \( X \) counts the number of heads in the sequence of \( n \) coin flips, and
- the summation on the right-hand side is equal to the number of 1’s in the sequence \( X_1, X_2, \ldots, X_n \), which, by definition, is equal to the number of heads in the sequence of \( n \) coin flips.
Using the Linearity of Expectation (see Theorem 6.5.2), we get
\[
\mathbb{E}(X) = \mathbb{E}\left(\sum_{i=1}^{n} X_i\right)
= \sum_{i=1}^{n} \mathbb{E}(X_i).
\]
Thus, we have to determine the expected value of \(X_i\). Since \(X_i\) is either 1 or 0, we have, using Lemma 6.4.3,
\[
\mathbb{E}(X_i) = 1 \cdot \Pr(X_i = 1) + 0 \cdot \Pr(X_i = 0)
= \Pr(X_i = 1)
= \Pr(\text{the } i\text{-th coin flip results in heads})
= p.
\]
We conclude that
\[
\mathbb{E}(X) = \sum_{i=1}^{n} \mathbb{E}(X_i)
= \sum_{i=1}^{n} p
= pn.
\]
I hope you agree that this is much easier than what we did before.

The distribution function of the random variable \(X\) is given by (6.3). This function is called a binomial distribution:

**Definition 6.7.1** Let \(n \geq 1\) be an integer and let \(p\) be a real number with \(0 < p < 1\). A random variable \(X\) has a **binomial distribution with parameters** \(n\) and \(p\), if its distribution function satisfies
\[
\Pr(X = k) = \binom{n}{k} p^k (1 - p)^{n-k}
\]
for any integer \(k\) with \(0 \leq k \leq n\).

Our calculation that led to the value of \(\mathbb{E}(X)\) proves the following theorem:
Theorem 6.7.2 Let $n \geq 1$ be an integer, let $p$ be a real number with $0 < p < 1$, and let $X$ be a random variable that has a binomial distribution with parameters $n$ and $p$. Then

$$\mathbb{E}(X) = pn.$$ 

6.8 Indicator Random Variables

In Section 6.7, we considered the random variable $X$ whose value is equal to the number of heads in a sequence of $n$ independent coin flips. In Section 6.7.2, we defined a sequence $X_1, X_2, \ldots, X_n$ of random variables, where $X_i = 1$ if the $i$-th coin flip results in heads and $X_i = 0$ otherwise. This random variable $X_i$ indicates whether or not the $i$-th flip in the sequence is heads. Because of this, we call $X_i$ an indicator random variable.

Definition 6.8.1 A random variable $X$ is an indicator random variable, if it can only take values in $\{0, 1\}$.

As we have already seen in Section 6.7.2, the expected value of an indicator random variable is easy to determine:

Lemma 6.8.2 If $X$ is an indicator random variable, then

$$\mathbb{E}(X) = \Pr(X = 1).$$

Proof. Since $X$ is either 0 or 1, we have, using Lemma 6.4.3

$$\mathbb{E}(X) = 0 \cdot \Pr(X = 0) + 1 \cdot \Pr(X = 1) = \Pr(X = 1).$$

In the following subsections, we will see some examples of how indicator random variables can be used to compute the expected value of non-trivial random variables.
6.8.1 Runs in Random Bitstrings

Let \( n \) be a large integer. We generate a random bitstring \( S = s_1s_2 \ldots s_n \)
by flipping a fair coin, independently, \( n \) times. Let \( k \geq 1 \) be an integer. Recall from Section 5.13 that a run of length \( k \) is a consecutive subsequence of \( S \), all of whose bits are equal. Define the random variable \( X \) to be the number of runs of length \( k \).

For example, if \( S \) is the bitstring

\[
\begin{array}{cccccccccccccccc}
0 & 0 & 1 & 1 & 1 & 1 & 0 & 0 & 1 & 1 & 0 & 0 & 0 & 0 & 1 \\
1 & 2 & 3 & 4 & 5 & 6 & 7 & 8 & 9 & 10 & 11 & 12 & 13 & 14 & 15 & 16
\end{array}
\]

and \( k = 3 \), then \( X = 6 \), because \( S \) contains 6 runs of length 3, starting at positions 3, 4, 5, 8, 13, and 14.

We want to determine the expected value \( \mathbb{E}(X) \) of \( X \).

A run of length \( k \) can start at any of the positions \( 1, 2, \ldots, n-k+1 \). Our approach will be to define an indicator random variable that tells us whether or not the subsequence of length \( k \) that starts at any such position is a run. Thus, for any \( i \) with \( 1 \leq i \leq n-k+1 \), we define the indicator random variable

\[
X_i = \begin{cases} 
1 & \text{if the subsequence } s_is_{i+1} \ldots s_{i+k-1} \text{ is a run,} \\
0 & \text{otherwise.}
\end{cases}
\]

Using Lemma 6.8.2 we get

\[
\mathbb{E}(X_i) = \Pr(X_i = 1).
\]

Since \( X_i = 1 \) if and only if all bits in the subsequence \( s_is_{i+1} \ldots s_{i+k-1} \) are 0 or all bits in this subsequence are 1, we have

\[
\mathbb{E}(X_i) = \Pr(X_i = 1) = (1/2)^k + (1/2)^k = 1/2^{k-1}.
\]

Since

\[
X = \sum_{i=1}^{n-k+1} X_i,
\]

\[
\mathbb{E}(X) = \sum_{i=1}^{n-k+1} \mathbb{E}(X_i) = \sum_{i=1}^{n-k+1} \frac{1}{2^{k-1}} = \frac{n-k+1}{2^{k-1}}.
\]
the Linearity of Expectation (see Theorem 6.5.2) implies that

\[ \mathbb{E}(X) = \mathbb{E}\left( \sum_{i=1}^{n-k+1} X_i \right) \]
\[ = \sum_{i=1}^{n-k+1} \mathbb{E}(X_i) \]
\[ = \sum_{i=1}^{n-k+1} \frac{1}{2^{k-1}} \]
\[ = \frac{n - k + 1}{2^{k-1}}. \]

Observe that the random variables \( X_1, X_2, \ldots, X_{n+k-1} \) are not independent. (Can you see why?) Nevertheless, our derivation is correct, because the Linearity of Expectation is valid for any sequence of random variables.

For example, if we take \( k = 1 + \log n \), then \( 2^{k-1} = 2^{\log n} = n \), so that

\[ \mathbb{E}(X) = \frac{n - \log n}{n} = 1 - \frac{\log n}{n}. \]

Thus, for large values of \( n \), the expected number of runs of length \( 1 + \log n \) is very close to 1. This is in line with Section 5.13, because we proved there that it is very likely that the sequence contains a run of length about \( \log n \).

If we take \( k = 1 + \frac{1}{2} \log n \), then

\[ 2^{k-1} = 2^{(\log n)/2} = 2^{\log \sqrt{n}} = \sqrt{n} \]

and

\[ \mathbb{E}(X) = \frac{n - \frac{1}{2} \log n}{\sqrt{n}} = \sqrt{n} - \frac{\log n}{2\sqrt{n}}. \]

Thus, for large values of \( n \), the expected number of runs of length \( 1 + \frac{1}{2} \log n \) is very close to \( \sqrt{n} \).

### 6.8.2 Largest Elements in Prefixes of Random Permutations

Consider a sequence \( s_1, s_2, \ldots, s_n \) of \( n \) numbers. The following algorithm computes the largest element in this sequence:
We would like to know the number of times that line (*) is executed, i.e., the number of times that the value of the variable \( max \) changes. For example, if the input sequence is

\[
3, 2, 5, 4, 6, 1,
\]

then the value of \( max \) changes 3 times, namely when we encounter 3, 5, and 6. On the other hand, for the sequence

\[
6, 5, 4, 3, 2, 1,
\]

the value of \( max \) changes only once, whereas for

\[
1, 2, 3, 4, 5, 6,
\]

it changes 6 times.

Assume that the input sequence \( s_1, s_2, \ldots, s_n \) is a random permutation of the integers 1, 2, \ldots, \( n \); thus, each permutation has probability \( 1/n! \) of being the input. We define the random variable \( X \) to be the number of times that line (*) is executed when running algorithm \( \text{FindMax}(s_1, s_2, \ldots, s_n) \). We are interested in the expected value \( \mathbb{E}(X) \) of \( X \).

The algorithm makes \( n \) iterations and in each one, line (*) is either executed or not executed. We define, for each iteration, an indicator random variable that tells us whether or not line (*) is executed during that iteration. That is, for any \( i \) with \( 1 \leq i \leq n \), we define

\[
X_i = \begin{cases} 
1 & \text{if line (*) is executed in the } i\text{-th iteration,} \\
0 & \text{otherwise.}
\end{cases}
\]

Since

\[
X = \sum_{i=1}^n X_i,
\]
it follows from the Linearity of Expectation (see Theorem \ref{thm:linearity}) that

\[
\mathbb{E}(X) = \mathbb{E}\left(\sum_{i=1}^{n} X_i\right) = \sum_{i=1}^{n} \mathbb{E}(X_i) = \sum_{i=1}^{n} \Pr(X_i = 1).
\]

How do we determine \(\Pr(X_i = 1)\)? Observe that \(X_i = 1\) if and only if the maximum of the subsequence \(s_1, s_2, \ldots, s_i\) is at the last position in this subsequence. Since the entire sequence \(s_1, s_2, \ldots, s_n\) is a random permutation of \(1, 2, \ldots, n\), the elements in the subsequence \(s_1, s_2, \ldots, s_i\) are in a random order as well. The largest element in this subsequence is in any of the \(i\) positions with equal probability \(1/i\). In particular, the probability that the largest element is at the last position in this subsequence is equal to \(1/i\). It follows that

\[
\Pr(X_i = 1) = 1/i.
\]

This can be proved in a more formal way as follows: By the Product Rule, the number of permutations \(s_1, s_2, \ldots, s_n\) of \(1, 2, \ldots, n\) for which \(s_i\) is the largest element among \(s_1, s_2, \ldots, s_i\) is equal to

\[
\binom{n}{i} (i-1)! (n-i)! = n!/i.
\]

(Do you see why?) Therefore,

\[
\Pr(X_i = 1) = \frac{n!/i}{n!} = 1/i.
\]

Thus,

\[
\mathbb{E}(X) = \sum_{i=1}^{n} \Pr(X_i = 1) = \sum_{i=1}^{n} 1/i = 1 + \frac{1}{2} + \frac{1}{3} + \cdots + \frac{1}{n}.
\]
The number on the right-hand side is called the harmonic number and denoted by $H_n$. In the following subsection, we will show that $H_n$ is approximately equal to $\ln n$. Thus, the expected number of times that line (*) of algorithm FINDMAX is executed, when given as input a random permutation of 1, 2, \ldots, $n$, is about $\ln n$.

### 6.8.3 Estimating the Harmonic Number

Consider a positive real-valued decreasing function $f$. Thus, if $x \leq x'$, then $f(x) \geq f(x')$. We would like to estimate the summation

$$\sum_{i=1}^{n} f(i).$$

For example, if we take $f(x) = 1/x$, then the summation is the harmonic number $H_n$ of the previous subsection.

For each $i$ with $2 \leq i \leq n$, draw the rectangle with bottom-left corner at the point $(i-1, 0)$ and top-right corner at the point $(i, f(i))$, as in the figure below.

The area of the $i$-th rectangle is equal to $f(i)$ and, thus,

$$\sum_{i=1}^{n} f(i)$$

is equal to the sum of

- $f(1)$ and
the total area of the $n - 1$ rectangles.

Since $f$ is decreasing, the rectangles are below the graph $y = f(x)$. It follows that the total area of the $n - 1$ rectangles is less than or equal to the area between $f$ and the $x$-axis, between $x = 1$ and $x = n$. We conclude that

$$\sum_{i=1}^{n} f(i) \leq f(1) + \int_{1}^{n} f(x) \, dx. \quad (6.4)$$

To obtain a lower bound on the summation, we modify the figure as indicated below: For each $i$ with $1 \leq i \leq n$, we draw the rectangle with bottom-left corner at the point $(i, 0)$ and top-right corner at the point $(i + 1, f(i))$.

In this case, the graph $y = f(x)$ is below the top sides of the rectangles and, therefore,

$$\sum_{i=1}^{n} f(i) \geq \int_{1}^{n+1} f(x) \, dx. \quad (6.5)$$

If we apply (6.4) and (6.5) to the function $f(x) = 1/x$, then we get

$$H_n = \sum_{i=1}^{n} \frac{1}{i} \leq 1 + \int_{1}^{n} \frac{dx}{x} = 1 + \ln n$$
and

\[ H_n = \sum_{i=1}^{n} \frac{1}{i} \geq \int_{1}^{n+1} \frac{dx}{x} = \ln(n+1). \]

We have proved the following result:

**Lemma 6.8.3** The harmonic number \( H_n = \sum_{i=1}^{n} \frac{1}{i} \) satisfies

\[ \ln(n+1) \leq H_n \leq 1 + \ln n. \]

### 6.9 The Insertion-Sort Algorithm

**InsertionSort** is a simple sorting algorithm that takes as input an array \( A[1...n] \) of numbers. The algorithm uses a for-loop in which a variable \( i \) runs from 2 to \( n \). At the start of the \( i \)-th iteration,

- the subarray \( A[1...i-1] \) is sorted, whereas
- the algorithm has not yet seen any of the elements in the subarray \( A[i...n] \).

In the \( i \)-th iteration, the algorithm takes the element \( A[i] \) and repeatedly swaps it with its left neighbor until the subarray \( A[1...i] \) is sorted. The pseudocode of this algorithm is given below.

**Algorithm** **InsertionSort**(*A[1...n]*):

```plaintext
for i = 2 to n
    do j = i;
        while j > 1 and A[j] < A[j - 1]
            do swap A[j] and A[j - 1];
                j = j - 1
        endwhile
endfor
```
We are interested in the total number of swaps that are made by this algorithm. The worst-case happens when the input array is sorted in reverse order, in which case the total number of swaps is equal to

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

Thus, in the worst case, each of the \(\binom{n}{2}\) pairs of input elements is swapped.

Assume that the input array \(A[1\ldots n]\) contains a uniformly random permutation of the integers \(1, 2, \ldots, n\). Thus, each permutation has probability \(1/n!\) of being the input. We define the random variable \(X\) to be the total number of swaps made when running algorithm \(\text{INSERTIONSORT}(A[1\ldots n])\). We will determine the expected value \(\mathbb{E}(X)\) of \(X\).

Since we want to count the number of pairs of input elements that are swapped, we will use, for each pair of input elements, an indicator random variable that indicates whether or not this pair gets swapped by the algorithm. That is, for each \(a\) and \(b\) with \(1 \leq a < b \leq n\), we define

\[
X_{ab} = \begin{cases} 
1 & \text{if } a \text{ and } b \text{ get swapped by the algorithm,} \\
0 & \text{otherwise.}
\end{cases}
\]

We observe that, since \(a < b\), these two elements get swapped if and only if in the input array, \(b\) is to the left of \(a\). Since the input array contains a uniformly random permutation, the events “\(b\) is to the left of \(a\)” and “\(a\) is to the left of \(b\)” are symmetric. Therefore, we have

\[
\mathbb{E}(X_{ab}) = \Pr(X_{ab} = 1) = 1/2.
\]

A formal proof of this is obtained by showing that there are \(n!/2\) permutations of \(1, 2, \ldots, n\) in which \(b\) appears to the left of \(a\) and, thus, \(n!/2\) permutations in which \(a\) appears to the left of \(b\). (See also Exercise [5.34](#).)

Since each pair of input elements is swapped at most once, we have

\[
X = \sum_{a=1}^{n-1} \sum_{b=a+1}^{n} X_{ab}.
\]
It follows from the Linearity of Expectation (see Theorem 6.5.2) that

$$\mathbb{E}(X) = \mathbb{E}\left(\sum_{a=1}^{n-1} \sum_{b=a+1}^{n} X_{ab}\right)$$

$$= \sum_{a=1}^{n-1} \sum_{b=a+1}^{n} \mathbb{E}(X_{ab})$$

$$= \sum_{a=1}^{n-1} \sum_{b=a+1}^{n} \frac{1}{2}$$

$$= \frac{1}{2} \binom{n}{2}. $$

Thus, the expected number of swaps on a random input array is one half times the worst-case number of swaps.

6.10 The Quick-Sort Algorithm

We have already seen algorithm QUICKSORT in Section 1.3. This algorithm takes as input an array $A[1\ldots n]$ of numbers, which we assume for simplicity to be distinct. A generic call QUICKSORT($A, i, j$) takes two indices $i$ and $j$ and sorts the subarray $A[i\ldots j]$. Thus, the call QUICKSORT($A, 1, n$) sorts the entire array.
The Quick-Sort Algorithm

Algorithm QuickSort(A, i, j):

if $i < j$

then $p$ = uniformly random element in $A[i \ldots j]$;

compare $p$ with all other elements in $A[i \ldots j]$;

rearrange $A[i \ldots j]$ such that it has the following form (this rearranging defines the value of $k$):

\[
\begin{array}{ccc}
<i & p & >p \\
\hline
i & k & j
\end{array}
\]

QuickSort(A, i, $k - 1$);

QuickSort(A, $k + 1$, j)

endif

The element $p$ is called the pivot. We have seen in Section 1.3 that the worst-case running time of algorithm QuickSort(A, 1, n) is $\Theta(n^2)$. In this section, we will prove that the expected running time is only $O(n \log n)$.

We assume for simplicity that the input array is a permutation of the integers 1, 2, \ldots, n. We do not make any other assumption about the input. In particular, we do not assume that the input is a random permutation. The only place where randomization is used is when the pivot is chosen: It is chosen uniformly at random in the subarray on which QuickSort is called.

The quantity that we will analyze is the total number of comparisons (between pairs of input elements) that are made during the entire execution of algorithm QuickSort(A, 1, n). In such a comparison, the algorithm takes two distinct input elements, say $a$ and $b$, and decides whether $a < b$ or $a > b$. Observe from the pseudocode that the only comparisons being made are between the pivot and all other elements in the subarray that is the input to the current call to QuickSort. Since the operation “compare $a$ to $b$” is the same as the operation “compare $b$ to $a$” (even though the outcomes are opposite), we will assume below that in such a comparison, $a < b$.

We define the random variable $X$ to be the total number of comparisons that are made by algorithm QuickSort(A, 1, n). We will prove that $E(X) = O(n \log n)$.

For each $a$ and $b$ with $1 \leq a < b \leq n$, we define the indicator random
variable

\[ X_{ab} = \begin{cases} 
1 & \text{if } a \text{ and } b \text{ are compared to each other when} \\
& \text{running } \text{QUICKSORT}(A, 1, n), \\
0 & \text{otherwise.} 
\end{cases} \]

Since each pair of input elements is compared at most once, we have

\[ X = \sum_{a=1}^{n-1} \sum_{b=a+1}^{n} X_{ab}. \]

It follows from the Linearity of Expectation (see Theorem \ref{thm:linearity}) that

\[ \mathbb{E}(X) = \mathbb{E} \left( \sum_{a=1}^{n-1} \sum_{b=a+1}^{n} X_{ab} \right) \]

\[ = \sum_{a=1}^{n-1} \sum_{b=a+1}^{n} \mathbb{E}(X_{ab}) \]

\[ = \sum_{a=1}^{n-1} \sum_{b=a+1}^{n} \Pr(X_{ab} = 1). \]

We consider two input elements \( a \) and \( b \) with \( 1 \leq a < b \leq n \). We are going to determine \( \Pr(X_{ab} = 1) \), i.e., the probability that the elements \( a \) and \( b \) are compared to each other when running algorithm \( \text{QUICKSORT}(A, 1, n) \). Define the set

\[ S_{ab} = \{a, a + 1, \ldots, b\}. \]

At the start of algorithm \( \text{QUICKSORT}(A, 1, n) \), all elements of the set \( S_{ab} \) are part of the input. Consider the first pivot \( p \) that is chosen. We observe the following:

- Assume that \( p \notin S_{ab} \).

  - If \( p < a \), then after the algorithm has rearranged the input array, all elements of the set \( S_{ab} \) are to the right of \( p \) and, thus, all these elements are part of the input for the recursive call \( \text{QUICKSORT}(A, k+1, n) \). During the rearranging, \( a \) and \( b \) have not been compared to each other. However, they may be compared to each other during later recursive calls.
– If \( p > b \), then after the algorithm has rearranged the input array, all elements of the set \( S_{ab} \) are to the left of \( p \) and, thus, all these elements are part of the input for the recursive call \( \text{QUICKSORT}(A, 1, k-1) \). During the rearranging, \( a \) and \( b \) have not been compared to each other. However, they \textit{may} be compared to each other during later recursive calls.

\[
\bullet \text{ Assume that } p \in S_{ab}.
\]

– If \( p \neq a \) and \( p \neq b \), then after the algorithm has rearranged the input array, \( a \) is to the left of \( p \) and \( b \) is to the right of \( p \). During the rearranging, \( a \) and \( b \) have not been compared to each other. Also, since \( a \) and \( b \) have been “separated”, they will not be compared to each other during later recursive calls. Thus, we have \( X_{ab} = 0 \).

– If \( p = a \) or \( p = b \), then during the rearranging, \( a \) and \( b \) have been compared to each other. Thus, we have \( X_{ab} = 1 \). (Note that in later recursive calls, \( a \) and \( b \) will not be compared to each other again.)

We conclude that whether or not \( a \) and \( b \) are compared to each other is completely determined by the element of the set \( S_{ab} \) that is the first element in this set to be chosen as a pivot. If this element is equal to \( a \) or \( b \), then \( X_{ab} = 1 \). On the other hand, if this element belongs to \( S_{ab} \setminus \{a, b\} \), then \( X_{ab} = 0 \). Since

\[
\bullet \text{ in any recursive call, the pivot is chosen uniformly at random from the subarray that is the input for this call, and}
\]

\[
\bullet \text{ at the start of the first recursive call in which the pivot belongs to the set } S_{ab}, \text{ all elements of this set are part of the input for this call,}
\]

each of the \( b - a + 1 \) elements of \( S_{ab} \) has the same probability of being the first element of \( S_{ab} \) that is chosen as a pivot. It follows that

\[
\Pr(X_{ab} = 1) = \frac{2}{b - a + 1}.
\]
We conclude that

\[ \mathbb{E}(X) = \sum_{a=1}^{n-1} \sum_{b=a+1}^{n} \Pr(X_{ab} = 1) = \sum_{a=1}^{n-1} \sum_{b=a+1}^{n} \frac{2}{b-a+1} \leq 2 \sum_{a=1}^{n-1} \left( \frac{1}{2} + \frac{1}{3} + \cdots + \frac{1}{n-a+1} \right) \]

where \( H_n \) is the harmonic number that we have seen in Sections 6.8.2 and 6.8.3. Using Lemma 6.8.3, it follows that

\[ \mathbb{E}(X) \leq 2n \ln n. \]

### 6.11 Skip Lists

Consider a set \( S \) of \( n \) numbers. We would like to store these numbers in a data structure that supports the following operations:

- **SEARCH\((x)\):** This operation returns the largest element in the set \( S \) that is less than or equal to \( x \).
- **INSERT\((x)\):** This operation inserts the number \( x \) into the set \( S \).
- **DELETE\((x)\):** This operation deletes the number \( x \) from the set \( S \).

A standard data structure for this problem is a balanced binary search tree (such as a red-black tree or an AVL-tree), which allows each of these three operations to be performed in \( O(\log n) \) time. Searching in a binary search
tree is straightforward, but keeping the tree balanced after an insertion or deletion is cumbersome.

In this section, we introduce skip lists as an alternative data structure. A skip list is constructed using the outcomes of coin flips, which result in a structure that is balanced in the expected sense. Because of this, the insertion and deletion algorithms become straightforward: We, as a programmer, do not have to take care of rebalancing operations, because the coin flips take care of this.

To define a skip list for the set $S$, we first construct a sequence $S_0, S_1, S_2, \ldots$ of subsets of $S$:

- Let $S_0 = S$.
- For $i = 0, 1, 2, \ldots$, assume that the set $S_i$ has already been constructed. If $S_i$ is non-empty, we do the following:
  - Initialize an empty set $S_{i+1}$.
  - For each element $y$ in the set $S_i$, flip a fair and independent coin. If the coin comes up heads, element $y$ is added to the set $S_{i+1}$.

The process terminates as soon as the next set $S_{i+1}$ is empty.

Let $h$ be the number of non-empty sets that are constructed by this process, and consider the sequence $S_0, S_1, \ldots, S_h$ of sets. Observe that $h$ is a random variable and each of the sets $S_1, S_2, \ldots, S_h$ is a random subset of $S$.

The skip list for $S$ consists of the following:

- For each $i$ with $0 \leq i \leq h$, we store the sorted sequence of elements of the set $S_i$ in a linked list $L_i$.
  - Each node $u$ of $L_i$ stores one element of $S_i$, which is denoted by $key(u)$.
  - Each node $u$ of $L_i$ stores a pointer to its successor node in $L_i$, which is denoted by $right(u)$. If $u$ is the rightmost node in $L_i$, then $right(u) = \text{nil}$.
  - We add a dummy node at the beginning of $L_i$. The key of this node is $\text{nil}$ and its successor is the node of $L_i$ whose key is the smallest element in $S_i$. 

- For each \( i \) with \( 1 \leq i \leq h \) and each node \( u \) of \( L_i \), \( u \) stores a pointer to the node \( u' \) in \( L_{i-1} \) for which \( key(u') = key(u) \). The node \( u' \) is denoted by \( \text{down}(u) \).

- There is a pointer to the dummy node in the list \( L_h \). We will refer to this node as the root of the skip list.

The value of \( h \) is called the height of the skip list. An example of a skip list of height 3 for the set \( S = \{1, 2, 3, 4, 6, 7, 9\} \) is shown in the figure below.

6.11.1 Algorithm Search

The algorithm that searches for a number \( x \) keeps track of the current node \( u \) and the index \( i \) of the list \( L_i \) that contains \( u \). Initially, \( u \) is the root of the skip list and \( i = h \). At any moment, if \( i \geq 1 \), the algorithm tests if the key of \( \text{right}(u) \) is less than \( x \). If this is the case, then \( u \) moves one node to the right in the list \( L_i \); otherwise, \( u \) moves to the node \( \text{down}(u) \) in the list \( L_{i-1} \). Once \( i = 0 \), node \( u \) moves to the right in the list \( L_0 \) and stops at the last node whose key is at most equal to \( x \). The pseudocode of this algorithm \( \text{Search}(x) \) is given below.
The dashed edges in the figure below show the path that is followed when running algorithm \texttt{Search(7)}. Note that if we replace \textquotedblleft key(right(u)) < x\textquotedblright in the first while-loop by \textquotedblleft key(right(u)) \leq x\textquotedblright, we obtain a different path that ends in the same node: This path moves from the root to the node in $L_3$ whose key is 7, and then it moves down to the list $L_0$. As we will see later, using the condition \textquotedblleft key(right(u)) < x\textquotedblright simplifies the algorithm for deleting an element from the skip list.

6.11.2 **Algorithms** \texttt{INSERT} and \texttt{DELETE}

Algorithm \texttt{INSERT(x)} takes as input a number $x$ and inserts it into the skip list. This algorithm works as follows:
Run algorithm \texttt{SEARCH}(x) and consider the node \( u \) that is returned. We assume that \( \text{key}(u) \neq x \) and, thus, \( x \) is not in the skip list yet. Observe that the new element \( x \) belongs between the nodes \( u \) and \( \text{right}(u) \).

Flip a fair and independent coin repeatedly until it comes up tails for the first time. Let \( k \) be the number of flips.

Add the new element \( x \) to the lists \( L_0, L_1, \ldots, L_{k-1} \). Note that if \( k \geq h+2 \), we have to add new lists \( L_{h+1}, \ldots, L_{k-1} \) to the skip list (each one containing a dummy node and a node storing \( x \)), set \( h = k - 1 \), and update the pointer to the root of the new skip list.

When adding \( x \) to a list \( L_i \), we have to know its predecessor in this list.

- To find these predecessors, we modify algorithm \texttt{SEARCH}(x) as follows: Each time the current node \( u \) moves down, we push \( u \) onto an initially empty stack. In this way, the predecessors that we need are stored, in the correct order, on the stack.

- An easier way that avoids using a stack is to flip the coin and, thus, determine \( k \), before running algorithm \texttt{SEARCH}(x). We then modify algorithm \texttt{SEARCH}(x): If \( i < k \) and the current node \( u \) moves down, we add the new element \( x \) to \( L_i \) between the nodes \( u \) and \( \text{right}(u) \).

The figure below shows the skip list that results when inserting the number 5 into the skip list of the previous figure. In this case, \( k = 3 \) and the new element is added to the lists \( L_0, L_1, \) and \( L_2 \). The dashed edges show the pointers that are changed during this insertion.
Algorithm \textsc{Delete}(x) takes as input a number \( x \) and deletes it from the skip list. This algorithm does the following:

- Run a modified version of algorithm \textsc{Search}(x): Each time the current node \( u \) moves down, test if \( \text{key(right}(u)) = x \). If this is the case, delete the node \( \text{right}(u) \) by setting \( \text{right}(u) = \text{right(right}(u)) \). Finally, delete the node in \( L_0 \) whose key is equal to \( x \).

- At this moment, it may happen that some of the lists \( L_h, L_{h-1}, \ldots \) only consist of dummy nodes. If this is the case, delete these lists, and update the height \( h \) and the root of the new skip list.

Implementation details of skip lists and algorithms \textsc{Search}, \textsc{Insert}, and \textsc{Delete} can be found in Pat Morin’s free textbook \textit{Open Data Structures}, which is available at \url{http://opendatastructures.org/}

### 6.11.3 Analysis of Skip Lists

In this subsection, we will prove that the expected size of a skip list is \( O(n) \) and the expected running time of algorithm \textsc{Search} is \( O(\log n) \). This will imply that the expected running times of algorithms \textsc{Insert} and \textsc{Delete} are \( O(\log n) \) as well. Throughout this subsection, we assume for simplicity that \( n \) is a power of \( 2 \), so that \( \log n \) is an integer.

Consider again the lists \( L_0, L_1, \ldots, L_h \) in the skip list. For the purpose of analysis, we define for each integer \( i > h \), \( L_i \) to be an empty list.

For each element \( x \) in the list \( L_0 \), we define the random variable \( h(x) \) to be the largest value of \( i \) such that \( x \) is contained in the list \( L_i \). Thus, \( x \) occurs in the lists \( L_0, L_1, \ldots, L_{h(x)} \), but not in the list \( L_{h(x) + 1} \).

**Lemma 6.11.1** For any element \( x \) in the list \( L_0 \), we have

\[
\mathbb{E}(h(x)) = 1.
\]

**Proof.** The value of \( h(x) \) is determined by the following process: flip a fair coin repeatedly and independently until it comes up tails for the first time. The value of \( h(x) \) is then equal to the number of flips minus one. For example, if we flip the coin three times (i.e., obtain the sequence \( HHT \)), then \( x \) is contained in the lists \( L_0, L_1, \) and \( L_2 \), but not in \( L_3 \); thus, \( h(x) = 2 \). By Theorem 6.6.2, the expected number of coin flips is equal to two. As a result, the expected value of \( h(x) \) is equal to one. \( \blacksquare \)
Lemma 6.11.2 For any element $x$ in the list $L_0$ and for any $i \geq 0$,
$$\Pr(x \in L_i) = \frac{1}{2^i}.$$  

Proof. The claim follows from the fact that $x$ is contained in the list $L_i$ if and only if the first $i$ coin flips for $x$ all result in heads. \[\Box\]

Lemma 6.11.3 Ignoring the dummy element, we have, for each $i \geq 0$,
$$\mathbb{E}(|L_i|) = \frac{n}{2^i}.$$  

Proof. We know from Lemma 6.11.2 that each element $x$ is contained in $L_i$ with probability $1/2^i$, independently of the other elements in $L_i$. Therefore, the size of $L_i$ is a random variable that has a binomial distribution with parameters $n$ and $p = 1/2^i$. The claim then follows from Theorem 6.7.2. \[\Box\]

Lemma 6.11.4 Let $X$ be the random variable that is equal to the total number of nodes in all lists $L_0, L_1, L_2, \ldots$, ignoring the dummy nodes. Then
$$\mathbb{E}(X) = 2n.$$  

Proof. We will give two proofs. In the first proof, we observe that
$$X = \sum_{i=0}^{h} |L_i|$$
and, thus,
$$\mathbb{E}(X) = \mathbb{E} \left( \sum_{i=0}^{h} |L_i| \right).$$

Observe that the number of terms in the summation on the right-hand side is equal to $h + 1$, which is a random variable. In general, the Linearity of Expectation does not apply to summations consisting of a random number of terms; see Exercise 6.37 for an example. Therefore, we proceed as follows. Recall that, for the purpose of analysis, we have defined for each integer $i > h$, $L_i$ to be an empty list. It follows that
$$X = \sum_{i=0}^{\infty} |L_i|.$$
Using the Linearity of Expectation (i.e., Theorem 6.5.3) and Lemmas 6.11.3 and 5.14.2 we get

\[
E(X) = E\left( \sum_{i=0}^{\infty} |L_i| \right) = \sum_{i=0}^{\infty} E(|L_i|) = \sum_{i=0}^{\infty} n/2^i = n \sum_{i=0}^{\infty} (1/2)^i = 2n.
\]

In the second proof, we use the fact that each element \( x \) occurs in exactly \( h(x) + 1 \) lists, namely \( L_0, L_1, \ldots, L_{h(x)} \). Thus, we have

\[
X = \sum_{x} (h(x) + 1).
\]

Using the Linearity of Expectation (i.e., Theorem 6.5.2) and Lemma 6.11.1 we get

\[
E(X) = E\left( \sum_{x} (h(x) + 1) \right) = \sum_{x} E(h(x) + 1) = \sum_{x} (E(h(x)) + 1) = \sum_{x} 2 = 2n.
\]

Lemma 6.11.5 The expected height \( h \) of the skip list satisfies

\[
E(h) \leq \log n + 1.
\]
Proof. Since
\[ h = \max_x h(x), \]
we have
\[ \mathbb{E}(h) = \mathbb{E}\left(\max_x h(x)\right). \]
It is tempting, but wrong, to think that this is equal to
\[ \max_x \mathbb{E}(h(x)), \]
which is equal to 1 by Lemma 6.11.1 (In Exercise 6.36 you will find a simple example showing that, in general, the expected value of a maximum is not equal to the maximum of the expected values.)

To prove a correct upper bound on \( \mathbb{E}(h) \), we introduce for each integer \( i \geq 1 \), an indicator random variable
\[ X_i = \begin{cases} 1 & \text{if the list } L_i \text{ is non-empty,} \\ 0 & \text{otherwise.} \end{cases} \]
We observe that
\[ h = \sum_{i=1}^{\infty} X_i. \]
Since \( X_i \) is either 0 or 1, it is obvious that
\[ \mathbb{E}(X_i) \leq 1. \quad (6.6) \]
We next claim that
\[ X_i \leq |L_i|. \quad (6.7) \]
To justify this, if the list \( L_i \) is empty, then (6.7) becomes 0 \( \leq 0 \), which is a true statement. On the other hand, if the list \( L_i \) is non-empty, then (6.7) becomes 1 \( \leq |L_i| \), which is again a true statement. Combining (6.7) with Lemmas 6.4.2 and 6.11.3, we obtain
\[ \mathbb{E}(X_i) \leq \mathbb{E}(|L_i|) = n/2^i. \quad (6.8) \]
Using the Linearity of Expectation (i.e., Theorem 6.5.3), we get

\[ \mathbb{E}(h) = \mathbb{E} \left( \sum_{i=1}^{\infty} X_i \right) \]

\[ = \sum_{i=1}^{\infty} \mathbb{E}(X_i) \]

\[ = \sum_{i=1}^{\log n} \mathbb{E}(X_i) + \sum_{i=\log n+1}^{\infty} \mathbb{E}(X_i). \]

If we apply (6.6) to the first summation and (6.8) to the second summation, we get

\[ \mathbb{E}(h) \leq \log n \sum_{i=1}^{\log n} 1 + \sum_{i=\log n+1}^{\infty} \frac{n}{2^i} \]

\[ = \log n + \sum_{j=0}^{\infty} \frac{n}{2^{\log n+1+j}} \]

\[ = \log n + \sum_{j=0}^{\infty} \frac{n}{n \cdot 2^{1+j}} \]

\[ = \log n + \sum_{j=0}^{\infty} \frac{1}{2^{1+j}} \]

\[ = \log n + \frac{1}{2} \sum_{j=0}^{\infty} \frac{1}{2^j} \]

\[ = \log n + \frac{1}{2} \cdot 2 \]

\[ = \log n + 1. \]

\[ \blacktriangleleft \]

**Lemma 6.11.6** Let \( Y \) be the random variable that is equal to the total number of nodes in all lists \( L_0, L_1, L_2, \ldots \), including the dummy nodes. Then

\[ \mathbb{E}(Y) \leq 2n + \log n + 2. \]
Proof. The total number of dummy nodes is equal to \( h + 1 \). Using the notation of Lemma 6.11.4, we have

\[
Y = X + h + 1.
\]

Thus, using the Linearity of Expectation (i.e., Theorem 6.5.2) and Lemmas 6.11.4 and 6.11.5, we get

\[
E(Y) = E(X + h + 1) = E(X) + E(h) + 1 \leq 2n + (\log n + 1) + 1 = 2n + \log n + 2.
\]

Consider any number \( x \). As we have seen in Section 6.11.1, algorithm \textsc{Search}(x) starts at the root of the skip list and follows a path to the rightmost node \( u \) in the bottom list \( L_0 \) for which \( \text{key}(u) \leq x \). We will refer to this path as the search path of the algorithm. In the figure below, you see the same skip list as we have seen before. The dashed edges form the search path of algorithm \textsc{Search}(7).

\[\text{Lemma 6.11.7} \quad \text{For any number } x, \text{ let } N \text{ be the random variable that is equal to the number of nodes on the search path of algorithm } \textsc{Search}(x). \text{ Then } E(N) \leq 2 \log n + 5.\]

Proof. Consider the node \( u \) that is returned by algorithm \textsc{Search}(x), let \( v \) be the second last node on the search path, and let \( P \) be the part of this search path from the root to \( v \). In the example above, \( u \) is the node in \( L_0 \)}
whose key is 7, \( v \) is the node in \( L_0 \) whose key is 6, and \( P \) is the part of the dashed path from the root to \( v \).

Let \( M \) be the random variable that is equal to the number of nodes on \( P \). Then, \( N = M + 1 \) and
\[
\mathbb{E}(N) = \mathbb{E}(M + 1) = \mathbb{E}(M) + 1.
\]

Thus, it suffices to prove that
\[
\mathbb{E}(M) \leq 2 \log n + 4.
\]

Define the following path \( P' \) in the skip list:

- \( P' \) starts at node \( v \).
- At any node on \( P' \), the path \( P' \) goes up one level if this is possible, and goes left one node otherwise.

You should convince yourself that this path \( P' \) is the reverse of \( P \) and, therefore, \( M \) is the number of nodes on \( P' \). You should also convince yourself that this may not be true, if we take for \( P \) the path from the root to \( u \).

For each \( i \geq 0 \), define the random variable \( M_i \) to be equal to the number of nodes in the list \( L_i \) at which the path \( P' \) goes left. Then, \( M \) is the sum of

- \( h \): these are the nodes on \( P' \) at which \( P' \) goes up one level,
- \( 1 \): this accounts for the last node on \( P' \), which is the root, and
- \( \sum_{i=0}^{h} M_i \).

Thus,
\[
\mathbb{E}(M) = \mathbb{E} \left( h + 1 + \sum_{i=0}^{h} M_i \right)
= \mathbb{E}(h) + 1 + \mathbb{E} \left( \sum_{i=0}^{h} M_i \right).
\]

As in the proof of Lemma \[6.11.4\] the number of terms in the latter summation is equal to \( h + 1 \), which is a random variable. Therefore, we cannot apply
the Linearity of Expectation to this sum. As in the proof of Lemma 6.11.4, we proceed as follows. We first observe that

\[ M = h + 1 + \sum_{i=0}^{\infty} M_i. \]

As the figure below indicates, the random variable \( M_i \) can be interpreted as being the number of tails obtained when flipping a fair coin until it comes up heads for the first time. Since (i) the list \( L_i \) may be empty (in which case \( M_i = 0 \)) or (ii) the portion of the path \( P' \) in \( L_i \) may terminate because it reaches the dummy node, \( M_i \) is in fact less than or equal to the number of tails.

\[
\begin{align*}
L_{i+1} : & \quad \square \\
L_i : & \quad H \rightarrow T \rightarrow T \rightarrow T \rightarrow T \\
L_{i-1} : & \quad \square
\end{align*}
\]

Therefore, by Lemma 6.4.2 and Theorem 6.6.2,

\[ \mathbb{E}(M_i) \leq 1. \quad (6.9) \]

Also, since \( M_i \) is less than or equal to the size \( |L_i| \) of the list \( L_i \) (ignoring the dummy node), we have, using Lemmas 6.4.2 and 6.11.3,

\[ \mathbb{E}(M_i) \leq \mathbb{E}(|L_i|) = n/2^i. \quad (6.10) \]

Using the Linearity of Expectation (i.e., Theorem 6.5.3), we get

\[
\mathbb{E}(M) = \mathbb{E} \left( h + 1 + \sum_{i=0}^{\infty} M_i \right) \\
= \mathbb{E}(h) + 1 + \sum_{i=0}^{\infty} \mathbb{E}(M_i) \\
= \mathbb{E}(h) + 1 + \sum_{i=0}^{\log n} \mathbb{E}(M_i) + \sum_{i=\log n+1}^{\infty} \mathbb{E}(M_i). \]

6.12. Exercises

We know from Lemma [6.11.5] that $E(h) \leq \log n + 1$. If we apply (6.9) to the first summation and (6.10) to the second summation, we get

$$E(M) \leq (\log n + 1) + 1 + \sum_{i=0}^{\log n} 1 + \sum_{i=\log n+1}^{\infty} n/2^i$$

$$= 2\log n + 3 + \sum_{i=\log n+1}^{\infty} n/2^i.$$  

We have seen the infinite series in the proof of Lemma [6.11.5] and showed that it is equal to 1. Thus, we conclude that

$$E(M) \leq 2\log n + 4.$$  

6.12 Exercises

6.1 Consider a fair coin that has 0 on one side and 1 on the other side. We flip this coin once and roll a fair die twice. Define the following random variables:

- $X$ = the result of the coin,
- $Y$ = the sum of the two dice,
- $Z = X \cdot Y$.

- Determine the distribution functions of $X$, $Y$, and $Z$.
- Are $X$ and $Y$ independent random variables?
- Are $X$ and $Z$ independent random variables?
- Are $Y$ and $Z$ independent random variables?
- Are $X$, $Y$ and $Z$ mutually independent random variables?

6.2 Lindsay and Simon want to play a game in which the expected amount of money that each of them wins is equal to zero. After having chosen a number $x$, the game is played as follows: Lindsay rolls a fair die, independently, three times.
If none of the three rolls results in 6, then Lindsay pays one dollar to Simon.

If exactly one of the rolls results in 6, then Simon pays one dollar to Lindsay.

If exactly two rolls result in 6, then Simon pays two dollars to Lindsay.

If all three rolls result in 6, then Simon pays \( x \) dollars to Lindsay.

Determine the value of \( x \).

6.3 We flip a fair coin twice (independently). For each heads, you win 3 dollars, whereas for each tails, you lose 2 dollars. Define the random variable \( X \) to be the amount of money that you win.

- Use the definition of expected value to compute \( \mathbb{E}(X) \).
- Use the Linearity of Expectation to compute \( \mathbb{E}(X) \).

We flip a fair coin 27 times (independently). For each heads, you win 3 dollars, whereas for each tails, you lose 2 dollars. Define the random variable \( Y \) to be the amount of money that you win.

- Compute the expected value \( \mathbb{E}(Y) \) of \( Y \).

6.4 Let \( r \) and \( b \) be positive integers and define \( \alpha = \frac{r}{r+b} \). A bowl contains \( r \) red balls and \( b \) blue balls; thus, \( \alpha \) is the fraction of the balls that are red. Consider the following experiment:

- Choose one ball uniformly at random.
  - If the chosen ball is red, then put it back, together with an additional red ball.
  - If the chosen ball is blue, then put it back, together with an additional blue ball.

Define the random variable \( X \) to be the fraction of the balls that are red, after this experiment. Prove that \( \mathbb{E}(X) = \alpha \).

6.5 The Ontario Lottery and Gaming Corporation (OLG) offers the following lottery game:
・ OLG chooses a winning number \( w \) in the set \( S = \{0, 1, 2, \ldots, 999\} \).

・ If John wants to play, he pays $1 and chooses a number \( x \) in \( S \).
  
  - If \( x = w \), then John receives $700 from OLG. In this case, John wins $699.
  
  - Otherwise, \( x \neq w \) and John does not receive anything. In this case, John loses $1.

Assume that

・ John plays this game once per day for one year (i.e., for 365 days),

・ each day, OLG chooses a new winning number,

・ each day, John chooses \( x \) uniformly at random from the set \( S \), independently from previous choices.

Define the random variable \( X \) to be the total amount of dollars that John wins during one year. Determine the expected value \( \mathbb{E}(X) \). (Hint: Use the Linearity of Expectation.)

6.6 Assume we roll two fair and independent dice. Let \((i, j)\) be the outcome, where \( i \) is the result of the first die and \( j \) is the result of the second die. Define the random variables

\[ X = |i - j| \]

and

\[ Y = \max(i, j). \]

Are \( X \) and \( Y \) independent random variables?

6.7 Consider two random variables \( X \) and \( Y \). If \( X \) and \( Y \) are independent, then it can be shown that

\[ \mathbb{E}(XY) = \mathbb{E}(X) \cdot \mathbb{E}(Y). \]

In this exercise, you will show that the converse of this statement is, in general, not true.

Let \( X \) be the random variable that takes each of the values \(-1, 0, \) and \( 1 \) with probability \( 1/3 \). Let \( Y \) be the random variable with value \( Y = X^2 \).
• Prove that $X$ and $Y$ are not independent.

• Prove that $\mathbb{E}(XY) = \mathbb{E}(X) \cdot \mathbb{E}(Y)$.

6.8 Jennifer loves to drink India Pale Ale (IPA), whereas Lindsay Bangs prefers wheat beer. Jennifer and Lindsay decide to go to their favorite pub *Chez Connor et Simon*. The beer menu shows that this pub has ten beers on tap:

• Five of these beers are of the IPA style.

• Three of these beers are of the wheat beer style.

• Two of these beers are of the pilsner style.

Jennifer and Lindsay order a uniformly random subset of seven beers (thus, there are no duplicates). Define the following random variables:

$$J = \text{the number of IPAs in this order},$$
$$L = \text{the number of wheat beers in this order}.$$

• Determine the expected value $\mathbb{E}(L)$ of the random variable $L$.

• Are $J$ and $L$ independent random variables?

6.9 We flip a fair coin independently and stop as soon as we get a $T$ followed by an $H$. Define the random variable $X$ to be the number of coin flips. For example, if the coin flips are $HHHTTTTH$, then $X = 8$.

• Determine the expected value $\mathbb{E}(X)$ of $X$. (*Hint:* Use the Linearity of Expectation.)

6.10 In Section 6.6, we have shown that for $-1 < x < 1$,

$$\sum_{k=1}^{\infty} k x^k = \frac{x}{(1-x)^2}.$$

In this question, you will prove this formula in a different way.
Consider the following infinite matrix:

\[
\begin{pmatrix}
x & 0 & 0 & 0 & 0 & \ldots \\
x^2 & x^2 & 0 & 0 & 0 & \ldots \\
x^3 & x^3 & x^3 & 0 & 0 & \ldots \\
x^4 & x^4 & x^4 & x^4 & 0 & \ldots \\
x^5 & x^5 & x^5 & x^5 & x^5 & 0 & \ldots \\
\vdots & \vdots & \vdots & \vdots & \vdots & \ddots \\
\end{pmatrix}
\]

We are going to add all elements in this matrix in two different ways. A \textit{row-sum} is the sum of all elements in one row, whereas a \textit{column-sum} is the sum of all elements in one column.

Note that the sum of all row-sums is equal to

\[x + 2x^2 + 3x^3 + 4x^4 + 5x^5 + \cdots = \sum_{k=1}^{\infty} kx^k.\]

- Using only the formula \(\sum_{k=0}^{\infty} x^k = \frac{1}{1-x}\) and algebraic manipulation, prove that the sum of all column-sums is equal to

\[\frac{x}{(1-x)^2}.\]

6.11 Let \(X\) be a random variable that takes values in \(\{0, 1, 2, 3, \ldots\}\). By Lemma 6.4.3, we have

\[\mathbb{E}(X) = \sum_{k=1}^{\infty} k \cdot \Pr(X = k).\]

As in Exercise 6.10, define an infinite matrix and use it to prove that

\[\mathbb{E}(X) = \sum_{k=1}^{\infty} \Pr(X \geq k).\]

6.12 Let \(0 < p < 1\) and consider a coin that comes up heads with probability \(p\) and tails with probability \(1 - p\). We flip the coin independently until it comes up heads for the first time. Define the random variable \(X\) to be the number of times that we flip the coin. In Section 6.6, we have shown that \(\mathbb{E}(X) = 1/p\). Below, you will prove this in a different way.
Let \( k \geq 1 \) be an integer. Determine \( \Pr(X \geq k) \).

Using only the formula \( \sum_{k=0}^{\infty} x^k = \frac{1}{1-x} \), the expression for \( \mathbb{E}(X) \) from Exercise 6.11 and your answer for \( \Pr(X \geq k) \), prove that \( \mathbb{E}(X) = 1/p \).

6.13 Consider an experiment that is successful with probability 0.8. We repeat this experiment (independently) until it is successful for the first time. The first 5 times we do the experiment, we have to pay $10 per experiment. After this, we have to pay $5 per experiment. Define the random variable \( X \) to be the total amount of money that we have to pay during all experiments. Determine the expected value \( \mathbb{E}(X) \).

Hint: Recall that \( \sum_{k=1}^{\infty} kx^{k-1} = 1/(1-x)^2 \).

6.14 When Lindsay and Simon have a child, this child is a boy with probability 1/2 and a girl with probability 1/2, independently of the gender of previous children. Lindsay and Simon stop having children as soon as they have a girl. Define the random variables

\[
B = \text{the number of boys that Lindsay and Simon have}
\]

and

\[
G = \text{the number of girls that Lindsay and Simon have}.
\]

Determine the expected values \( \mathbb{E}(B) \) and \( \mathbb{E}(G) \).

6.15 Let \( p \) be a real number with \( 0 < p < 1 \). When Lindsay and Simon have a child, this child is a boy with probability \( p \) and a girl with probability \( 1 - p \), independently of the gender of previous children. Lindsay and Simon stop having children as soon as they have a child that has the same gender as their first child. Define the random variable \( X \) to be the number of children that Lindsay and Simon have. Determine the expected value \( \mathbb{E}(X) \).

Hint: Recall that \( \sum_{k=1}^{\infty} kx^{k-1} = 1/(1-x)^2 \).

6.16 Let \( X_1, X_2, \ldots, X_n \) be a sequence of mutually independent random variables. For each \( i \) with \( 1 \leq i \leq n \), assume that

- the variable \( X_i \) is either equal to 0 or equal to \( n + 1 \), and

- \( \mathbb{E}(X_i) = 1 \).
Determine
\[ \Pr(X_1 + X_2 + \cdots + X_n \leq n). \]

**6.17** The Ottawa Senators and the Toronto Maple Leafs play a best-of-seven series: These two hockey teams play against each other until one of them has won four games. Assume that

- in any game, the Sens have a probability of 3/4 of defeating the Leafs,
- the results of the games are mutually independent.

Determine the probability that seven games are played in this series.

**6.18** Let \( n \geq 1 \) be an integer, let \( p \) be a real number with \( 0 < p < 1 \), and let \( X \) be a random variable that has a binomial distribution with parameters \( n \) and \( p \). In Section 6.7.1 we have seen that the expected value \( \mathbb{E}(X) \) of \( X \) satisfies

\[ \mathbb{E}(X) = \sum_{k=1}^{n} k \binom{n}{k} p^k (1 - p)^{n-k}. \]

(6.11)

Recall Newton’s Binomial Theorem, i.e., Theorem 3.6.5:

\[ (x + y)^n = \sum_{k=0}^{n} \binom{n}{k} x^{n-k} y^k. \]

- Use (6.11) to prove that \( \mathbb{E}(X) = pn \), by taking the derivative, with respect to \( y \), in Newton’s Binomial Theorem.

**6.19** A **maximal run of ones** in a bitstring is a maximal consecutive substring of ones. For example, the bitstring 100111110100111 has four maximal runs of ones: 1, 11111, 1, and 111.

Let \( n \geq 1 \) be an integer and consider a random bitstring of length \( n \). Define the random variable \( X \) to be the number of maximal runs of ones in this bitstring. Determine the expected value \( \mathbb{E}(X) \) of \( X \). (Hint: Use indicator random variables.)

**6.20** Let \( n \geq 1 \) be an integer and consider a uniformly random permutation \( a_1, a_2, \ldots, a_n \) of the set \( \{1, 2, \ldots, n\} \). Define the random variable \( X \) to be the number of indices \( i \) for which \( 1 \leq i < n \) and \( a_i < a_{i+1} \).

Determine the expected value \( \mathbb{E}(X) \) of \( X \). (Hint: Use indicator random variables.)
6.21 Let \( n \geq 1 \) be an integer and consider a permutation \( a_1, a_2, \ldots, a_n \) of the set \( \{1, 2, \ldots, n\} \). We partition this permutation into *increasing subsequences*. For example, for \( n = 10 \), the permutation

\[
3, 5, 8, 1, 2, 4, 10, 7, 6, 9
\]
is partitioned into four increasing subsequences: (i) 3, 5, 8, (ii) 1, 2, 4, 10, (iii) 7, and (iv) 6, 9.

Let \( a_1, a_2, \ldots, a_n \) be a uniformly random permutation of the set \( \{1, 2, \ldots, n\} \). Define the random variable \( X \) to be the number of increasing subsequences in the partition of this permutation. For the example above, we have \( X = 4 \).

In this exercise, you will determine the expected value \( \mathbb{E}(X) \) of \( X \) in two different ways.

- For each \( i \) with \( 1 \leq i \leq n \), let
  \[
  X_i = \begin{cases} 
  1 & \text{if an increasing subsequence starts at position } i, \\
  0 & \text{otherwise.}
  \end{cases}
  \]

  For the example above, we have \( X_1 = 1, X_2 = 0, X_3 = 0, \) and \( X_8 = 1 \).
  - Determine \( \mathbb{E}(X_1) \).
  - Let \( i \) be an integer with \( 2 \leq i \leq n \). Use the Product Rule to determine the number of permutations of \( \{1, 2, \ldots, n\} \) for which \( X_i = 1 \).
  - Use these indicator random variables to determine \( \mathbb{E}(X) \).

- For each \( i \) with \( 1 \leq i \leq n \), let
  \[
  Y_i = \begin{cases} 
  1 & \text{if the value } i \text{ is the leftmost element of an increasing subsequence,} \\
  0 & \text{otherwise.}
  \end{cases}
  \]

  For the example above, we have \( Y_1 = 1, Y_3 = 1, Y_5 = 0, \) and \( Y_7 = 1 \).
  - Determine \( \mathbb{E}(Y_1) \).
  - Let \( i \) be an integer with \( 2 \leq i \leq n \). Use the Product Rule to determine the number of permutations of \( \{1, 2, \ldots, n\} \) for which \( Y_i = 1 \).
  - Use these indicator random variables to determine \( \mathbb{E}(X) \).
6.22 Lindsay Bangs and Simon Pratt visit their favorite pub that has 10 different beers on tap. Both Lindsay and Simon order, independently of each other, a uniformly random subset of 5 beers.

- One of the beers available is Leo’s Early Breakfast IPA. Determine the probability that this is one of the beers that Lindsay orders.

- Let $X$ be the random variable whose value is the number of beers that are ordered by both Lindsay and Simon. Determine the expected value $\mathbb{E}(X)$ of $X$. (Hint: Use indicator random variables.)

6.23 Lindsay and Simon have discovered a new pub that has $n$ different beers $B_1, B_2, \ldots, B_n$ on tap, where $n \geq 1$ is an integer. They want to try all different beers in this pub and agree on the following approach: During a period of $n$ days, they visit the pub every day. On each day, they drink one of the beers. Lindsay drinks the beers in order, i.e., on the $i$-th day, she drinks beer $B_i$. Simon takes a uniformly random permutation $a_1, a_2, \ldots, a_n$ of the set $\{1, 2, \ldots, n\}$ and drinks beer $B_{a_i}$ on the $i$-th day.

Let $X$ be the random variable whose value is the number of days during which Lindsay and Simon drink the same beer. Determine the expected value $\mathbb{E}(X)$ of $X$. (Hint: Use indicator random variables.)

6.24 Consider the following recursive algorithm TWOTAILS, which takes as input a positive integer $n$:

```
Algorithm TWOTAILS(n):
    // all coin flips made are mutually independent
    flip a fair coin twice;
    if the coin came up tails exactly twice
        then return $2^n$
    else TWOTAILS(n + 1)
    endif
```

- You run algorithm TWOTAILS(1), i.e., with $n = 1$. Define the random variable $X$ to be the value of the output of this algorithm. Let $k \geq 1$ be an integer. Determine $\Pr (X = 2^k)$. 
• Is the expected value $\mathbb{E}(X)$ of the random variable $X$ finite or infinite?

6.25 Let $A[1 \ldots n]$ be an array of $n$ numbers. Consider the following two algorithms, which take as input the array $A$ and a number $x$. If $x$ is not present in $A$, then these algorithms return the message “not present”. Otherwise, they return an index $i$ such that $A[i] = x$. The first algorithm runs linear search from left to right, whereas the second algorithm runs linear search from right to left.

**Algorithm** LinearSearchLeftToRight($A, x$):

```plaintext
i := 1;
while $i \leq n$ and $A[i] \neq x$ do $i := i + 1$ endwhile;
if $i = n + 1$ then return “not present” else return $i$ endif
```

**Algorithm** LinearSearchRightToLeft($A, x$):

```plaintext
i := n;
while $i \geq 1$ and $A[i] \neq x$ do $i := i - 1$ endwhile;
if $i = 0$ then return “not present” else return $i$ endif
```

Consider the following algorithm, which again take as input the array $A$ and a number $x$. If $x$ is not present in $A$, then it returns the message “not present”. Otherwise, it returns an index $i$ such that $A[i] = x$.

**Algorithm** RandomLinearSearch($A, x$):

```plaintext
flip a fair coin;
if the coin comes up heads then LinearSearchLeftToRight($A, x$)
else LinearSearchRightToLeft($A, x$)
endif
```

Assume that the number $x$ occurs exactly once in the array $A$ and let $k$ be the index such that $A[k] = x$. Let $X$ be the random variable whose value is the number of times the test “$A[i] \neq x$” is made in algorithm RandomLinearSearch($A, x$). (In words, $X$ is the number of comparisons made by algorithm RandomLinearSearch($A, x$).) Determine the expected value $\mathbb{E}(X)$ of $X$. 

6.26 Consider the set $V = \{1, 2, \ldots, n\}$ and let $p$ be a real number with $0 < p < 1$. We construct a graph $G = (V, E)$ with vertex set $V$, whose edge set $E$ is determined by the following random process: Each unordered pair $\{i, j\}$ of vertices, where $i \neq j$, occurs as an edge in $E$ with probability $p$, independently of the other unordered pairs.

A triangle in $G$ is an unordered triple $\{i, j, k\}$ of distinct vertices, such that $\{i, j\}$, $\{j, k\}$, and $\{k, i\}$ are edges in $G$.

Define the random variable $X$ to be the total number of triangles in the graph $G$. Determine the expected value $\mathbb{E}(X)$. (Hint: Use indicator random variables.)

6.27 In Section 6.9, we have seen the following algorithm \textsc{InsertionSort}, which sorts any input array $A[1 \ldots n]$:

```plaintext
Algorithm \textsc{InsertionSort}(A[1 \ldots n]):

for $i = 2$ to $n$
do $j = i$;
   while $j > 1$ and $A[j] < A[j - 1]$
   $j = j - 1$
endwhile
endfor
```

Consider an input array $A[1 \ldots n]$, where each element $A[i]$ is chosen independently and uniformly at random from the set $\{1, 2, \ldots, m\}$.

- Let $i$ and $j$ be two indices with $1 \leq i < j \leq n$, and consider the values $A[i]$ and $A[j]$ (just before the algorithm starts). Prove that

$$\Pr(A[i] > A[j]) = \frac{1}{2} - \frac{1}{2m}.$$

- Let $X$ be the random variable that is equal to the number of times the swap-operation is performed when running \textsc{InsertionSort}(A[1 \ldots n]). Determine the expected value $\mathbb{E}(X)$ of $X$.

6.28 Assume we have $n$ balls and $m$ boxes. We throw the balls independently and uniformly at random in the boxes. Thus, for $1 \leq k \leq n$ and $1 \leq i \leq m$,

$$\Pr(\text{the } k\text{-th ball falls in the } i\text{-th box}) = 1/m.$$
Define the following three random variables:

\[ X = \text{the number of boxes that do not contain any ball}, \]
\[ Y = \text{the number of boxes that contain at least one ball}, \]
\[ Z = \text{the number of boxes that contain exactly one ball}. \]

- Determine the expected values \( E(X), E(Y), \) and \( E(Z) \).

- Assuming that \( m = n \), determine the limits
  
  1. \( \lim_{n \to \infty} E(X)/n \),
  2. \( \lim_{n \to \infty} E(Y)/n \),
  3. \( \lim_{n \to \infty} E(Z)/n \).

(Use the fact that \( \lim_{n \to \infty} (1 - 1/n)n = 1/e \).)

**6.29** Let \( 0 < p < 1 \) and consider a coin that comes up heads with probability \( p \) and tails with probability \( 1 - p \). For each integer \( n \), let \( b_n \) be the outcome when flipping this coin; thus, \( b_n \in \{H, T\} \). The values \( b_n \) partition the set of integers into intervals, where each interval is a maximal consecutive sequence of zero or more \( T \)'s followed by one \( H \):

\[
\begin{array}{cccccccccccc}
\cdots & H & T & T & T & H & T & H & T & H & H & \cdots \\
\cdots & -2 & -1 & 0 & 1 & 2 & 3 & 4 & 5 & 6 & 7 & \cdots
\end{array}
\]

- Consider the interval that contains the integer 0, and let \( X \) be its length. (In the example above, \( X = 4 \).) Determine the expected value \( E(X) \) of \( X \).

(\textit{Hint: Use the Linearity of Expectation. The answer is not } 1/p, \textit{which is the expected number of coin flips until the first } H.)

**6.30** Your friend Mick takes a permutation of 1, 2, \ldots, \( n \), stores it in boxes \( B_1, B_2, \ldots, B_n \) (so that each box stores exactly one number), and then closes all boxes. You have no idea what the permutation is.

Mick opens the boxes \( B_1, B_2, \ldots, B_n \), one after another. For each \( i \) with \( 1 \leq i \leq n \), just before opening box \( B_i \), you have to guess which number is stored in it.
6.12. Exercises

- Assume that, when you guess the number in box $B_i$, you do not remember the numbers stored in $B_1, B_2, \ldots, B_{i-1}$. Then, the only reasonable thing you can do is to take a random element in $\{1, 2, \ldots, n\}$ and guess that this random element is stored in $B_i$.

Assume that you do this for each $i$ with $1 \leq i \leq n$. Let $X$ be the random variable whose value is equal to the number of times that your guess is correct. Compute the expected value $E(X)$ of $X$.

- Now assume that your memory is perfect, so that, when you guess the number in box $B_i$, you know the numbers stored in $B_1, B_2, \ldots, B_{i-1}$.

How would you make the $n$ guesses such that the following is true: If $Y$ is the random variable whose value is equal to the number of times that your guess is correct, then the expected value $E(Y)$ of $Y$ satisfies $E(Y) = \Omega(\log n)$.

6.31 Recall that a permutation of the set $S = \{1, 2, \ldots, n\}$ is a bijection $f : S \to S$ (i.e., $f$ is one-to-one and onto).

- Consider a fixed element $k$ in the set $S$. How many permutations $f : S \to S$ are there for which $f(k) = k$?

- We choose a permutation $f : S \to S$ uniformly at random. Define the random variable $X = |\{k \in S : f(k) = k\}|$.

Determine the expected value $E(X)$.

6.32 Let $d$ be the number of days in one year, and consider a group $P_1, P_2, \ldots, P_n$ of $n$ people. Assume that each person has a random and independent birthday, i.e., any of the $d$ days has a probability of $1/d$ of being the birthday of $P_i$, which is independent of the birthdays of the other people.

Define the random variable $X$ to be the number of pairs $\{P_i, P_j\}$ of people that have the same birthday. Prove that

$$E(X) = \frac{1}{d} \binom{n}{2}.$$ 

6.33 Michiel’s Craft Beer Company (MCBC) sells $n$ different brands of India Pale Ale (IPA). When you place an order, MCBC sends you one bottle of
IPA, chosen uniformly at random from the \( n \) different brands, independently of previous orders.

Simon Pratt wants to try all different brands of IPA. He repeatedly places orders at MCBC (one bottle per order) until he has received at least one bottle of each brand.

Define the random variable \( X \) to be the total number of orders that Simon places. Determine the expected value \( E(X) \). (Hint: Use the Linearity of Expectation. If Simon has received exactly \( i \) different brands of IPA, how many orders does he expect to place until he receives a new brand?)

6.34 MCBC still sells \( n \) different brands of IPA. As in Exercise 6.33, when you place an order, MCBC sends you one bottle of IPA, chosen uniformly at random from the \( n \) different brands, independently of previous orders.

Simon Pratt places \( m \) orders at MCBC. Define the random variable \( X \) to be the total number of distinct brands that Simon receives. Determine the expected value \( E(X) \). (Hint: Use indicator random variables.)

6.35 You are given an array \( A[0 \ldots n-1] \) of \( n \) numbers. Let \( D \) be the number of distinct numbers that occur in this array. For each \( i \) with \( 0 \leq i \leq n-1 \), let \( N_i \) be the number of elements in the array that are equal to \( A[i] \).

- Show that \( D = \sum_{i=0}^{n-1} 1/N_i \).

Consider the following algorithm:

**Algorithm** \( \text{ESTIMATE}(A[1\ldots n]) \):

**Step 1:** Choose an integer \( k \) in \( \{0, 1, 2, \ldots, n-1\} \) uniformly at random, and let \( a = A[k] \).
**Step 2:** Traverse the array and compute the number \( N_k \) of times that \( a \) occurs.
**Step 3:** Return the value \( X = n/N_k \).

- Determine the expected value \( E(X) \) of the random variable \( X \). (Hint: Use the definition of expected value, i.e., Definition 6.4.1)

6.36 One of Jennifer and Thomas is chosen uniformly at random. The person who is chosen wins $100. Define the random variables \( J \) and \( T \) as follows:

\[ J = \text{the amount that Jennifer wins} \]
and

\[ T = \text{the amount that Thomas wins.} \]

Prove that

\[ \mathbb{E}(\max(J, T)) \neq \max(\mathbb{E}(J), \mathbb{E}(T)). \]

6.37 Consider the sample space

\[ S = \{(123), (132), (213), (231), (312), (321), (111), (222), (333)\}. \]

We choose an element \( u \) from \( S \) uniformly at random. For each \( i = 1, 2, 3 \), let \( X_i \) be the random variable whose value is the \( i \)-th number in \( u \). (For example, if \( u = (312) \), then \( X_1 = 3 \), \( X_2 = 1 \), and \( X_3 = 2 \).) Let \( N \) be the random variable whose value is equal to that of \( X_2 \).

- Verify that \( \Pr(X_i = k) = 1/3 \) for any \( i \) and \( k \) with \( 1 \leq i \leq 3 \) and \( 1 \leq k \leq 3 \).
- Verify that \( X_1, X_2 \) and \( X_3 \) are pairwise independent.
- Verify that \( X_1, X_2 \) and \( X_3 \) are not mutually independent.
- Verify that \( \sum_{i=1}^{N} \mathbb{E}(X_i) = 4 \).
- Verify that \( \mathbb{E}\left(\sum_{i=1}^{N} X_i\right) \neq \sum_{i=1}^{N} \mathbb{E}(X_i) \).

6.38 Let \( n \geq 2 \) be a power of two and consider a full binary tree with \( n \) leaves. Let \( a_1, a_2, \ldots, a_n \) be a random permutation of the numbers \( 1, 2, \ldots, n \). Store this permutation at the leaves of the tree, in the order \( a_1, a_2, \ldots, a_n \), from left to right. For example, if \( n = 8 \) and the permutation is \( 2, 8, 1, 4, 6, 3, 5, 7 \), then we obtain the following tree:

```
       O
      /|
     / |\n    O  O
   /   |
  O  O  O
 /   /|
O  O  O  O
/   /   |
O  O  O  O
/   /   |
O  O  O  O
       \|
       O
     /|
    / |\n   O  O
  /   |
 O  O
```

2 8 1 4 6 3 5 7
Perform the following process on the tree:

- Visit the levels of the tree from bottom to top.
- At each level, take all pairs of consecutive nodes that have the same parent. For each such pair, compare the numbers stored at the two nodes, and store the smaller of these two numbers at the common parent.

For our example tree, we obtain the following tree:

```
       1
      / \   
     3   5
    /     \
   1   2  4
  / \   /   \
 2  8  6  3  5
```

It is clear that at the end of this process, the root stores the number 1. Define the random variable $X$ to be the number that is not equal to 1 and that is stored at a child of the root; think of $X$ being the “loser of the final game”. For our example tree, $X = 3$.

In this exercise, you will determine the expected value $\mathbb{E}(X)$ of the random variable $X$.

- Prove that $2 \leq X \leq 1 + n/2$.
- Prove that the following is true for each $k$ with $1 \leq k \leq n/2$: $X \geq k + 1$ if and only if
  - all numbers $1, 2, \ldots, k$ are stored in the left subtree of the root
  - or all numbers $1, 2, \ldots, k$ are stored in the right subtree of the root.
- Prove that for each $k$ with $1 \leq k \leq n/2$,
  \[
  \Pr(X \geq k + 1) = 2 \cdot \frac{(n/2)^k k!(n-k)!}{n!} = 2 \cdot \frac{(n/2)^k}{(n \choose k)}.
  \]
• According to Exercise 6.11, we have

\[ \mathbb{E}(X) = \sum_{k=1}^{\infty} \Pr(X \geq k). \]

Prove that

\[ \mathbb{E}(X) = \Pr(X \geq 1) + \sum_{k=1}^{n/2} \Pr(X \geq k + 1). \]

• Use Exercise 3.32 to prove that

\[ \mathbb{E}(X) = 3 - \frac{4}{n + 2}. \]
Chapter 7

The Probabilistic Method

The Probabilistic Method is a very powerful and surprising tool that uses probability theory to prove results in discrete mathematics. In this chapter, we will illustrate this method using several examples.

7.1 Large Bipartite Subgraphs

Recall that a graph is a pair \( G = (V, E) \), where \( V \) is a finite set whose elements are called vertices and \( E \) is a set whose elements are unordered pairs of distinct vertices. The elements of \( E \) are called edges. Assume we partition the vertex set \( V \) of \( G \) into two subsets \( A \) and \( B \) (thus, \( A \cap B = \emptyset \) and \( A \cup B = V \)). We say that an edge of \( E \) is between \( A \) and \( B \), if one vertex of this edge is in \( A \) and the other vertex is in \( B \).

\[
\begin{array}{c}
\text{a} \\
\text{b} \\
\text{c} \\
\text{d} \\
\text{e} \\
\end{array}
\]

For example, in the graph above, let \( A = \{a, d\} \) and \( B = \{b, c, e\} \). Then four of the eight edges are between \( A \) and \( B \), namely \( \{a, b\} \), \( \{a, e\} \), \( \{d, c\} \), and \( \{d, e\} \). Thus, the vertex set of this graph can be partitioned into two
subsets $A$ and $B$, such that at least half of $G$’s edges are between $A$ and $B$. The following theorem states that this is true for any graph.

**Theorem 7.1.1** Let $G = (V, E)$ be a graph with $m$ edges. The vertex set $V$ of $G$ can be partitioned into two subsets $A$ and $B$ such that the number of edges between $A$ and $B$ is at least $m/2$.

**Proof.** Consider the following random process: Initialize $A = \emptyset$ and $B = \emptyset$. For each vertex $u$ of $G$, flip a fair and independent coin. If the coin comes up heads, add $u$ to $A$; otherwise, add $u$ to $B$.

Define the random variable $X$ to be the number of edges of $G$ that are between $A$ and $B$. We will determine the expected value $\mathbb{E}(X)$ of $X$.

Number the edges of $G$ arbitrarily as $e_1, e_2, \ldots, e_m$. For each $i$ with $1 \leq i \leq m$, define the indicator random variable

$$X_i = \begin{cases} 
1 & \text{if } e \text{ is an edge between } A \text{ and } B, \\
0 & \text{otherwise.}
\end{cases}$$

Then

$$X = \sum_{i=1}^{m} X_i$$

and

$$\mathbb{E}(X) = \mathbb{E}\left( \sum_{i=1}^{m} X_i \right) = \sum_{i=1}^{m} \mathbb{E}(X_i) = \sum_{i=1}^{m} \Pr(X_i = 1).$$

To determine $\Pr(X_i = 1)$, let $e_i$ have vertices $a$ and $b$. The following table shows the four possibilities for $a$ and $b$; each one of them occurs with probability $1/4$.

<table>
<thead>
<tr>
<th>$a \in A$, $b \in A$</th>
<th>$X_i = 0$</th>
</tr>
</thead>
<tbody>
<tr>
<td>$a \in A$, $b \in B$</td>
<td>$X_i = 1$</td>
</tr>
<tr>
<td>$a \in B$, $b \in A$</td>
<td>$X_i = 1$</td>
</tr>
<tr>
<td>$a \in B$, $b \in B$</td>
<td>$X_i = 0$</td>
</tr>
</tbody>
</table>
Since $X_i = 1$ in two out of the four cases, we have

$$\Pr (X_i = 1) = 2/4 = 1/2,$$

and it follows that

$$\mathbb{E}(X) = \sum_{i=1}^{m} 1/2 = m/2.$$

Assume the claim in the theorem does not hold. Then, no matter how we partition the vertex set $V$ into $A$ and $B$, the number of edges between $A$ and $B$ will be less than $m/2$. In particular, the random variable $X$ will always be less than $m/2$. But then, $\mathbb{E}(X) < m/2$ as well, contradicting that $\mathbb{E}(X) = m/2.$

7.2 Ramsey Theory

We return to a problem that we have seen in Section 1.1. Consider a complete graph with $n$ vertices, in which each vertex represents a person. Any pair of distinct vertices is connected by an edge. Such an edge is solid if the two persons representing the vertices of this edge are friends. If these persons are strangers, the edge is dashed. Consider a subset $S$ of $k$ vertices. We say that $S$ is a solid $k$-clique, if any two distinct vertices in $S$ are connected by a solid edge. Thus, a solid $k$-clique represents a group of $k$ mutual friends. If any two distinct vertices of $S$ are connected by a dashed edge, then we say that $S$ is a dashed $k$-clique; this represents a group of $k$ mutual strangers.

In Section 1.1, we stated, without proof, Theorem 1.1.3. We repeat the statement of this theorem and use the Probabilistic Method to prove it.

**Theorem 7.2.1** Let $k \geq 3$ and $n \geq 3$ be integers with $n \leq \lfloor 2^{k/2} \rfloor$. There exists a complete graph with $n$ vertices, in which each edge is either solid or dashed, such that this graph does not contain a solid $k$-clique and does not contain a dashed $k$-clique.

**Proof.** We denote the complete graph with $n$ vertices by $K_n$. Consider the following random process: For each edge $e$ of $K_n$, flip a fair and independent coin. If the coin comes up heads, make $e$ a solid edge; otherwise, make $e$ a dashed edge.
Define the event

\[ A = \text{“there is a solid } k\text{-clique or there is a dashed } k\text{-clique”}. \]

We will prove below that \( \Pr(A) < 1 \). This will imply that \( \Pr(\overline{A}) > 0 \), i.e., the event

\[ \overline{A} = \text{“there is no solid } k\text{-clique and there is no dashed } k\text{-clique”} \]

has a positive probability. This, in turn, will imply that the statement in the theorem holds: If the statement would not hold, then \( \Pr(\overline{A}) \) would be zero.

Thus, it remains to prove that \( \Pr(A) < 1 \). The vertex set of \( K_n \) has exactly \( \binom{n}{k} \) many subsets of size \( k \). We denote these subsets by \( V_i, i = 1, 2, \ldots, \binom{n}{k} \). For each \( i \) with \( 1 \leq i \leq \binom{n}{k} \), define the event

\[ A_i = \text{“} V_i \text{ is a solid } k\text{-clique or a dashed } k\text{-clique”}. \]

Since the event \( A_i \) occurs if and only if the edges joining the \( \binom{k}{2} \) pairs of vertices of \( V_i \) are either all solid or all dashed, we have

\[ \Pr(A_i) = \frac{2}{2^{\binom{k}{2}}}; \]

note that the denominator is equal to 2 to the power \( \binom{k}{2} \).

Since \( A \) occurs if and only if \( A_1 \lor A_2 \lor \cdots \lor A_{\binom{n}{k}} \) occurs, the Union Bound (i.e., Lemma 5.3.5) implies that

\[
\Pr(A) = \Pr\left( A_1 \lor A_2 \lor \cdots \lor A_{\binom{n}{k}} \right) \\
\leq \sum_{i=1}^{\binom{n}{k}} \Pr(A_i) \\
= \sum_{i=1}^{\binom{n}{k}} \frac{2}{2^{\binom{k}{2}}} \\
= \frac{2^{\binom{n}{k}}}{2^{\binom{k}{2}}};
\]
If we can show that the quantity in the last line is less than one, then the proof is complete. We have
\[
\frac{2\binom{n}{k}}{2\binom{k}{2}} = \frac{n(n-1)(n-2) \cdots (n-k+1)}{k!} \cdot \frac{2}{2^{(k^2-k)/2}}
\leq \frac{n^k}{k!} \cdot \frac{2^{1+k/2}}{2^{k^2/2}}.
\]
Since \(n \leq \lfloor 2^{k/2} \rfloor \leq 2^{k/2}\), we get
\[
\frac{2\binom{n}{k}}{2\binom{k}{2}} \leq \frac{(2^{k/2})^k}{k!} \cdot \frac{2^{1+k/2}}{2^{k^2/2}} = \frac{2^{1+k/2}}{k!}.
\]
By Exercise 2.8, we have \(k! > 2^{1+k/2}\) for \(k \geq 3\). Thus, we conclude that
\[
\frac{2\binom{n}{k}}{2\binom{k}{2}} < 1.
\]

Take, for example, \(k = 20\) and \(n = 1024\). Theorem 7.2.1 states that there exists a group of 1024 people that does not contain a subgroup of 20 mutual friends and does not contain a subgroup of 20 mutual strangers. In fact, the proof shows more: Consider a group of 1024 people such that any two are friends with probability 1/2, and strangers with probability 1/2. The above proof shows that \(\Pr(A)\), i.e., the probability that there is a subgroup of 20 mutual friends or there is a subgroup of 20 mutual strangers, satisfies
\[
\Pr(A) \leq \frac{2^{1+k/2}}{k!} = \frac{2^{11}}{20!}.
\]
Therefore, with probability at least
\[
1 - \frac{2^{11}}{20!} = 0.99999999999999989,
\]
(there are 15 nines) this group does not contain a subgroup of 20 mutual friends and does not contain a subgroup of 20 mutual strangers.
7.3 Sperner’s Theorem

In Section 1.2, we considered the following problem. Let $S$ be a set of size $n$ and consider a sequence $S_1, S_2, \ldots, S_m$ of $m$ subsets of $S$, such that for all $i$ and $j$ with $i \neq j$,

$$S_i \not\subseteq S_j \text{ and } S_j \not\subseteq S_i.$$  \hspace{1cm} (7.1)

What is the largest possible value of $m$ for which such a sequence exists?

The sequence consisting of all subsets of $S$ having size $\left\lfloor \frac{n}{2} \right\rfloor$ satisfies (7.1). This sequence has length $m = \binom{n}{\left\lfloor \frac{n}{2} \right\rfloor}$. In Section 1.2, we stated, without proof, that this is the largest possible value of $m$; see Theorem 1.2.1. After stating this theorem again, we will prove it using the Probabilistic Method.

**Theorem 7.3.1 (Sperner)** Let $n \geq 1$ be an integer and let $S$ be a set with $n$ elements. Let $S_1, S_2, \ldots, S_m$ be a sequence of $m$ subsets of $S$, such that for all $i$ and $j$ with $i \neq j$,

$$S_i \not\subseteq S_j \text{ and } S_j \not\subseteq S_i.$$  Then

$$m \leq \binom{n}{\left\lfloor \frac{n}{2} \right\rfloor}.$$

**Proof.** We assume that none of the subsets in the sequence $S_1, S_2, \ldots, S_m$ is empty, because otherwise, $m$ must be equal to 1, in which case the theorem clearly holds.

We assume that $S = \{1, 2, \ldots, n\}$. We choose a uniformly random permutation $a_1, a_2, \ldots, a_n$ of the elements of $S$; thus, each permutation has probability $1/n!$ of being chosen. Define the following sequence of subsets $A_1, A_2, \ldots, A_n$ of $S$: For $j = 1, 2, \ldots, n$,

$$A_j = \{a_1, a_2, \ldots, a_j\}.$$  

For example, if $n = 4$ and the permutation is 3, 1, 4, 2, then

$$A_1 = \{3\},$$  
$$A_2 = \{1, 3\},$$  
$$A_3 = \{1, 3, 4\},$$  
$$A_4 = \{1, 2, 3, 4\}.$$
Observe that the subsets $A_1, A_2, \ldots, A_n$ are random subsets of $S$, because the permutation was randomly chosen.

Consider a subset $S_i$ in the statement of the theorem. We say that $S_i$ occurs in the sequence $A_1, A_2, \ldots, A_n$ if there is an index $j$ such that $S_i = A_j$.

Define the random variable $X$ to be the number of subsets in the sequence $S_1, S_2, \ldots, S_m$ that occur in $A_1, A_2, \ldots, A_n$. Since the subsets $A_1, A_2, \ldots, A_n$ are properly nested, i.e.,

$$A_1 \subset A_2 \subset \cdots \subset A_n,$$

the assumption in the theorem implies that $X$ is either 0 or 1. It follows that the expected value of $X$ satisfies

\[ \mathbb{E}(X) \leq 1. \]

We now derive an exact expression for the value of $\mathbb{E}(X)$. For each $i$ with $1 \leq i \leq m$, define the indicator random variable

\[ X_i = \begin{cases} 1 & \text{if } S_i \text{ occurs in the sequence } A_1, A_2, \ldots, A_n, \\ 0 & \text{otherwise.} \end{cases} \]

Let $k$ denote the size of the subset $S_i$, i.e., $k = |S_i|$. Then

\[ X_i = 1 \text{ if and only if } S_i = A_k. \]

Since $A_k = \{a_1, a_2, \ldots, a_k\}$, $X_i = 1$ if and only if the first $k$ values in the permutation form a permutation of the subset $S_i$:

\[
\begin{array}{cccc}
  a_1, a_2, \ldots, a_k & a_{k+1}, a_{k+2}, \ldots, a_n \\
  \text{permutation of } S_i
\end{array}
\]

The Product Rule of Section 3.1 shows that there are $k!(n-k)!$ many permutations of $S$ that have this property. Therefore, since we chose a random permutation of $S$, we have

\[
\mathbb{E}(X_i) = \Pr(X_i = 1) = \frac{k!(n-k)!}{n!} = \frac{1}{\binom{n}{k}} = \frac{1}{\binom{|S_i|}{k}}.
\]
Thus, since
\[ X = \sum_{i=1}^{m} X_i, \]
we get
\[
\mathbb{E}(X) = \mathbb{E}\left( \sum_{i=1}^{m} X_i \right) \\
= \sum_{i=1}^{m} \mathbb{E}(X_i) \\
= \sum_{i=1}^{m} \frac{1}{|S_i|}. 
\]
If we combine this with our upper bound \( \mathbb{E}(X) \leq 1 \), we get
\[
\sum_{i=1}^{m} \frac{1}{|S_i|} \leq 1.
\]
For a fixed value of \( n \), the binomial coefficient \( \binom{n}{k} \) is maximized when \( k = \lfloor n/2 \rfloor \); i.e., the largest value in the \( n \)-th row of Pascal’s Triangle (see Section 3.8) is in the middle. Thus,
\[
\binom{n}{|S_i|} \leq \binom{n}{\lfloor n/2 \rfloor},
\]
implying that
\[
1 \geq \sum_{i=1}^{m} \frac{1}{\binom{n}{|S_i|}} \\
\geq \sum_{i=1}^{m} \frac{1}{\binom{n}{\lfloor n/2 \rfloor}} \\
= \frac{m}{\binom{n}{\lfloor n/2 \rfloor}}.
\]
We conclude that
\[ m \leq \binom{n}{\lfloor n/2 \rfloor}. \]
7.4 Planar Graphs and the Crossing Lemma

Consider a graph $G = (V, E)$. Any one-to-one function $f : V \rightarrow \mathbb{R}^2$ gives an embedding of $G$:

1. Each vertex $a$ of $V$ is drawn as the point $f(a)$ in the plane.
2. Each edge $\{a, b\}$ of $E$ is drawn as the straight-line segment $f(a)f(b)$ between the points $f(a)$ and $f(b)$.

Besides the function $f$ being one-to-one, we assume that it satisfies the following three properties:

1. For any two edges $\{a, b\}$ and $\{a', b'\}$ of $E$, the intersection of the line segments $f(a)f(b)$ and $f(a')f(b')$ is empty or consists of exactly one point.
2. For any edge $\{a, b\}$ in $E$ and any vertex $c$ in $V$, the point $f(c)$ is not in the interior of the line segment $f(a)f(b)$.
3. For any three edges $\{a, b\}$, $\{a', b'\}$, and $\{a'', b''\}$ of $E$, the line segments $f(a)f(b)$, $f(a')f(b')$, and $f(a'')f(b'')$ do not have a point in common that is in the interior of any of these line segments.

For simplicity, we do not distinguish any more between a graph and its embedding. That is, a vertex $a$ refers to both an element of $V$ and the point in the plane that represents $a$. Similarly, an edge refers to both an element of $E$ and the line segment that represents it.

7.4.1 Planar Graphs

**Definition 7.4.1** An embedding of a graph $G = (V, E)$ is called plane, if no two edges of $E$ intersect, except possibly at their endpoints. A graph $G$ is called planar if there is a plane embedding of $G$.

Consider a plane embedding of a planar graph $G$. Again for simplicity, we denote this embedding by $G$. This embedding consists of vertices, edges, and faces (one of them being the unbounded face). For example, in the following plane embedding, there are 11 vertices, 18 edges, and 9 faces.
In the rest of this section, we will use the following notation:

- $G$ denotes a plane embedding of a planar graph.
- $v$ denotes the number of vertices of $G$.
- $e$ denotes the number of edges of $G$.
- $f$ denotes the number of faces in the embedding of $G$.

How many edges can $G$ have? Since $G$ has $v$ vertices, we obviously have $e \leq \binom{v}{2} = \Theta(v^2)$, an upper bound which holds for any graph with $v$ vertices. Since our graph $G$ is planar, we expect a much smaller upper bound on $e$: If $G$ has $\Theta(v^2)$ edges, then it seems to be impossible to draw $G$ without edge crossings. Below, we will prove that $e$ is, in fact, at most linear in $v$. The proof will use Euler’s Theorem for planar graphs:

**Theorem 7.4.2 (Euler)** Consider any plane embedding of a planar graph $G$. Let $v$, $e$, and $f$ be the number of vertices, edges, and faces of this embedding, respectively. Moreover, let $c$ be the number of connected components of $G$. Then

$$v - e + f = c + 1. \tag{7.2}$$

**Proof.** The idea of the proof is as follows. We start by removing all edges from $G$ (but keep all vertices), and show that (7.2) holds. Then we add back the edges of $G$, one by one, and show that (7.2) remains valid throughout this process.

After having removed all edges, we have $e = 0$ and the embedding consists of a collection of $v$ points. Since $f = 1$ and $c = v$, the relation $v - e + f = c + 1$ holds.

Assume the relation $v - e + f = c + 1$ holds and consider what happens when we add an edge $ab$. There are two possible cases.
Case 1: Before adding the edge $ab$, the vertices $a$ and $b$ belong to the same connected component.

When adding the edge $ab$,

- the number $v$ of vertices does not change,
- the number $e$ of edges increases by one,
- the number $f$ of faces increases by one (because the edge $ab$ splits one face into two),
- the number $c$ of connected components does not change.

It follows that the relation $v - e + f = c + 1$ still holds after $ab$ has been added.

Case 2: Before adding the edge $ab$, the vertices $a$ and $b$ belong to different connected components.

When adding the edge $ab$,

- the number $v$ of vertices does not change,
- the number $e$ of edges increases by one,
- the number $f$ of faces does not change,
- the number $c$ of connected components decreases by one.
It again follows that the relation \( v - e + f = c + 1 \) still holds after \( ab \) has been added.

Usually, Euler’s Theorem is stated for connected planar graphs, i.e., planar graphs for which \( c = 1 \):

**Theorem 7.4.3 (Euler)** Consider any plane embedding of a connected planar graph \( G \). If \( v \), \( e \), and \( f \) denote the number of vertices, edges, and faces of this embedding, respectively, then

\[
v - e + f = 2.
\]

We now show how to use Euler’s Theorem to prove an upper bound on the number of edges and faces of any connected planar graph:

**Theorem 7.4.4** Let \( G \) be any plane embedding of a connected planar graph with \( v \geq 3 \) vertices. Then

1. \( G \) has at most \( 3v - 6 \) edges and

2. this embedding has at most \( 2v - 4 \) faces.

**Proof.** As before, let \( e \) and \( f \) denote the number of edges and faces of \( G \), respectively. If \( v = 3 \), then \( e \leq 3 \) and \( f \leq 2 \). Hence, in this case, we have \( e \leq 3v - 6 \) and \( f \leq 2v - 4 \).

Assume that \( v \geq 4 \). We number the faces of \( G \) arbitrarily from 1 to \( f \). For each \( i \) with \( 1 \leq i \leq f \), let \( m_i \) denote the number of edges on the \( i \)-th face of \( G \). Since each edge lies on the boundary of at most two faces, the summation \( \sum_{i=1}^{f} m_i \) counts each edge at most twice. Thus,

\[
\sum_{i=1}^{f} m_i \leq 2e.
\]

On the other hand, since \( G \) is connected and \( v \geq 4 \), each face has at least three edges on its boundary, i.e., \( m_i \geq 3 \). It follows that

\[
\sum_{i=1}^{f} m_i \geq 3f.
\]
Combining these two inequalities implies that $3f \leq 2e$, which we rewrite as

\[ f \leq 2e/3. \]

Using Euler’s formula (with $c = 1$, because $G$ is connected), we obtain

\[ e = v + f - 2 \leq v + 2e/3 - 2, \]

which is equivalent to

\[ e \leq 3v - 6. \]

We also obtain

\[ f \leq 2e/3 \leq 2(3v - 6)/3 = 2v - 4. \]

This completes the proof.

### 7.4.2 The Crossing Number of a Graph

Consider an embedding of a graph $G = (V, E)$. We say that two distinct edges of $E$ cross, if their interiors have a point in common. In this case, we call this common point a crossing.

The example below shows an embedding of the complete graph $K_6$ on six vertices, which are denoted by black dots. In this embedding, there are three crossings, which are numbered 1, 2, and 3.

![Diagram of K_6 with crossings labelled 1, 2, and 3.]

**Definition 7.4.5** The crossing number $cr(G)$ of a graph $G$ is defined to be the minimum number of crossings in any embedding of $G$.

Thus, a graph $G$ is planar if and only if $cr(G) = 0$. The example above shows that $cr(K_6) \leq 3$.

In the rest of this section, we consider the following problem: Given a graph $G$ with $v$ vertices and $e$ edges, can we prove good bounds, in terms of $v$ and $e$, on the crossing number $cr(G)$ of $G$?
A simple lower bound on the crossing number

Let $G$ be any graph with $v \geq 3$ vertices and $e$ edges. Consider an embedding of $G$ having $cr(G)$ crossings; hence, this embedding is the “best” one.

We “make” $G$ planar, by defining all crossings to be vertices. That is, let $H$ be the graph whose vertex set is the union of the vertex set of $G$ and the set of all crossings in the embedding. Edges of $G$ are cut by the crossings into smaller edges, which are edges in the graph $H$.

The figure below shows the planar version of the embedding of $K_6$ that we saw before. This new graph has 9 vertices and 21 edges.

![Planar embedding of $K_6$](image)

We make the following observations:

- The graph $H$ is planar, because it is embedded without any crossings.
- The graph $H$ has $v + cr(G)$ vertices.
- How many edges does $H$ have? Any crossing in $G$ is the intersection of exactly two edges of $G$; these two edges contribute four edges to $H$. Hence, for any crossing in $G$, the number of edges in $H$ increases by two. It follows that $H$ has $e + 2 \cdot cr(G)$ edges.

Since $H$ is planar, we know from Theorem 7.4.4 that the number of its edges is bounded from above by three times the number of its vertices minus six, i.e.,

$$e + 2 \cdot cr(G) \leq 3(v + cr(G)) - 6.$$

By rewriting this inequality, we obtain the following result:

**Theorem 7.4.6** For any graph $G$ with $v \geq 3$ vertices and $e$ edges, we have

$$cr(G) \geq e - 3v + 6.$$
For example, consider the complete graph $K_n$ on $n$ vertices, where $n \geq 3$. Since this graph has $\binom{n}{2}$ edges, we obtain
\[
\text{cr}(K_n) \geq \left(\frac{n}{2}\right) - 3n + 6 = \frac{1}{2}n^2 - \frac{7}{2}n + 6. \tag{7.3}
\]
For $n = 6$, we get $\text{cr}(K_6) \geq 3$. Since we have seen before that $\text{cr}(K_6) \leq 3$, it follows that $\text{cr}(K_6) = 3$.

Since $K_n$ has $\binom{n}{2}$ edges and any two of them cross at most once, we have the following obvious upper bound on the crossing number of $K_n$:
\[
\text{cr}(K_n) \leq \left(\frac{\binom{n}{2}}{2}\right) = O(n^4). \tag{7.4}
\]
(Of course (7.4) holds for any graph with $n$ vertices.)

To conclude this subsection, (7.3) gives an $n^2$-lower bound, whereas (7.4) gives an $n^4$-upper bound on the crossing number of $K_n$. In the next section, we will determine the true asymptotic behavior of $\text{cr}(K_n)$.

**A better lower bound on the crossing number**

As before, let $G$ be any graph with $v \geq 3$ vertices and $e$ edges. Again we consider an embedding of $G$ having $\text{cr}(G)$ crossings. In the rest of this subsection, we will use the Probabilistic Method to prove a lower bound on the crossing number of $G$.

We choose a real number $p$ such that $0 < p \leq 1$. Consider a coin that comes up heads with probability $p$ and comes up tails with probability $1 - p$. Let $G_p$ be the random subgraph of $G$, that is obtained as follows.

- For each vertex $a$ of $G$, flip the coin (independently of the other coin flips) and add $a$ as a vertex to $G_p$ if and only if the coin comes up heads.

- Each edge $ab$ of $G$ appears as an edge in $G_p$ if and only if both $a$ and $b$ are vertices of $G_p$.

Recall that we fixed the embedding of $G$. As a result, this random process gives us an embedding of $G_p$ (which may not be the best one in terms of the number of crossings.)
We denote the number of vertices, edges, and crossings in the embedding of $G_p$ by $v_p$, $e_p$, and $x_p$, respectively. Observe that these three quantities are random variables.

It follows from Theorem 7.4.6 that
\[
\text{cr}(G_p) - e_p + 3v_p \geq 6,
\]
provided that $v_p \geq 3$. This implies that
\[
\text{cr}(G_p) - e_p + 3v_p \geq 0,
\]
for any value of $v_p$ that results from our random choices.

Since $\text{cr}(G_p) \leq x_p$, it follows that
\[
x_p - e_p + 3v_p \geq 0.
\]
The left-hand side is a random variable, which is always non-negative, no matter what graph $G_p$ results from our random choices. Therefore, its expected value is also non-negative, i.e.,
\[
\mathbb{E}(x_p - e_p + 3v_p) \geq 0.
\]
Using the Linearity of Expectation (i.e., Theorem 6.5.2), we get
\[
\mathbb{E}(x_p) - \mathbb{E}(e_p) + 3 \cdot \mathbb{E}(v_p) \geq 0. \quad (7.5)
\]
We are now going to compute these three expected values separately.

The random variable $v_p$ is equal to the number of successes in $v$ independent trials, each one having success probability $p$. In other words, $v_p$ has a binomial distribution with parameters $v$ and $p$, and, therefore, by Theorem 6.7.2,
\[
\mathbb{E}(v_p) = pv.
\]
To compute $\mathbb{E}(e_p)$, we number the edges of $G$ arbitrarily from 1 to $e$. For each $i$ with $1 \leq i \leq e$, define $X_i$ to be the indicator random variable with value
\[
X_i = \begin{cases} 
1 & \text{if the } i\text{-th edge is an edge in } G_p, \\
0 & \text{otherwise.}
\end{cases}
\]
Since an edge of $G$ is in $G_p$ if and only if both its vertices are in $G_p$, it follows that
\[
\mathbb{E}(X_i) = \Pr(X_i = 1) = p^2.
\]
Then, since \( e_p = \sum_{i=1}^{e} X_i \), we get
\[
E(e_p) = E\left( \sum_{i=1}^{e} X_i \right) = \sum_{i=1}^{e} E(X_i) = \sum_{i=1}^{e} p^2 = p^2 e.
\]

Finally, we compute the expected value of the random variable \( x_p \). Number the crossings in the embedding of \( G \) arbitrarily from 1 to \( cr(G) \). For each \( i \) with \( 1 \leq i \leq cr(G) \), define \( Y_i \) to be the indicator random variable with value
\[
Y_i = \begin{cases} 
1 & \text{if the } i \text{-th crossing is a crossing in } G_p, \\
0 & \text{otherwise}.
\end{cases}
\]

Let \( ab \) and \( cd \) be the edges of \( G \) that cross in the \( i \)-th crossing. This crossing appears as a crossing in \( G_p \) if and only if both \( ab \) and \( cd \) are edges in \( G_p \). Since the points \( a, b, c, \) and \( d \) are pairwise distinct, it follows that the \( i \)-th crossing of \( G \) appears as a crossing in \( G_p \) with probability \( p^4 \). Thus,
\[
E(Y_i) = Pr(Y_i = 1) = p^4.
\]

Since \( x_p = \sum_{i=1}^{cr(G)} Y_i \), it follows that
\[
E(x_p) = E\left( \sum_{i=1}^{cr(G)} Y_i \right) = \sum_{i=1}^{cr(G)} E(Y_i) = \sum_{i=1}^{cr(G)} p^4 = p^4 \cdot cr(G).
\]

Substituting the three expected values into (7.5), we get
\[
p^4 \cdot cr(G) - p^2 e + 3 \cdot pv \geq 0,
\]
which we rewrite as
\[
cr(G) \geq \frac{p^2 e - 3pv}{p^4}. \tag{7.6}
\]

Observe that this inequality holds for any real number \( p \) with \( 0 < p \leq 1 \).

If we assume that \( e \geq 4v \), and take \( p = 4v/e \) (so that \( 0 < p \leq 1 \)), then we obtain a new lower bound on the crossing number:

\(^1\)By our definition of embedding, see Section 7.4.1 there are exactly two edges that determine the \( i \)-th crossing.
Theorem 7.4.7 (Crossing Lemma) Let $G$ be any graph with $v$ vertices and $e$ edges. If $e \geq 4v$, then

$$cr(G) \geq \frac{1}{64} \frac{e^3}{v^2}.$$ 

Applying this lower bound to the complete graph $K_n$ gives $cr(K_n) = \Omega(n^4)$. This lower bound is much better than the quadratic lower bound in (7.3) and it matches the upper bound in (7.4). Hence, we have shown that $cr(K_n) = \Theta(n^4)$.

Remark 7.4.8 Let $n$ be a very large integer and consider the complete graph $K_n$ with $v = n$ vertices and $e = \binom{n}{2}$ edges. Let us see what happens if we repeat the proof for this graph. We choose a random subgraph $G_p$ of $K_n$, where $p = 4v/e = 8/(n - 1)$. The expected number of vertices in $G_p$ is approximately equal to $8$. Thus, the random graph $G_p$ is, expected, extremely small. Then we apply the weak lower bound of Theorem 7.4.6 to this, again expected, extremely small graph. The result is a proof that in any embedding of the huge graph $K_n$, there are $\Omega(n^4)$ crossings!

7.5 Exercises

7.1 Prove that Theorem 7.4.4 also holds if $G$ is not connected.

7.2 Let $K_5$ be the complete graph on 5 vertices. In this graph, each pair of vertices is connected by an edge. Prove that $K_5$ is not planar.

7.3 Let $G$ be any embedding of a connected planar graph with $v \geq 4$ vertices. Assume that this embedding has no triangles, i.e., there are no three vertices $a$, $b$, and $c$, such that $ab$, $bc$, and $ac$ are edges of $G$.

- Prove that $G$ has at most $2v - 4$ edges.
- Let $K_{3,3}$ be the complete bipartite graph on 6 vertices. The vertex set of this graph consists of two sets $A$ and $B$, both of size three, and each vertex of $A$ is connected by an edge to each vertex of $B$. Prove that $K_{3,3}$ is not planar.

7.4 Consider the numbers $R_n$ that were defined in Section 4.6. In Section 4.6.1 we proved that $R_n = O(n^8)$. Prove that $R_n = O(n^4)$. 


7.5 Let $n$ be a sufficiently large positive integer and consider the complete graph $K_n$. This graph has vertex set $V = \{1, 2, \ldots, n\}$, and each pair of distinct vertices is connected by an undirected edge. (Thus, $K_n$ has $\binom{n}{2}$ edges.)

Let $\vec{K}_n$ be the directed graph obtained by making each edge $\{i, j\}$ of $K_n$ a directed edge; thus, in $\vec{K}_n$, this edge either occurs as the directed edge $(i, j)$ from $i$ to $j$ or as the directed edge $(j, i)$ from $j$ to $i$.

We say that three pairwise distinct vertices $i$, $j$, and $k$ define a directed triangle in $\vec{K}_n$, if

- $(i, j)$, $(j, k)$, and $(k, i)$ are edges in $\vec{K}_n$ or

- $(i, k)$, $(k, j)$, and $(j, i)$ are edges in $\vec{K}_n$.

Prove that there exists a way to direct the edges of $K_n$, such that the number of directed triangles in $\vec{K}_n$ is at least $\frac{1}{4} \binom{n}{3}$.

7.6 Let $G = (V, E)$ be a graph with vertex set $V$ and edge set $E$. A subset $I$ of $V$ is called an independent set if for any two distinct vertices $u$ and $v$ in $I$, $(u, v)$ is not an edge in $E$. For example, in the following graph, $I = \{a, e, i\}$ is an independent set.

Let $n$ and $m$ denote the number of vertices and edges in $G$, respectively, and assume that $m \geq n/2$. This exercise will lead you through a proof of the fact that $G$ contains an independent set of size at least $\frac{n^2}{4m}$.

Consider the following algorithm, in which all random choices made are mutually independent:
Algorithm **INDEPSET**\((G)\):

**Step 1:** Set \( H = G \).

**Step 2:** Let \( d = \frac{2m}{n} \). For each vertex \( v \) of \( H \), with probability \( 1 - \frac{1}{d} \), delete the vertex \( v \), together with its incident edges, from \( H \).

**Step 3:** As long as the graph \( H \) contains edges, do the following: Pick an arbitrary edge \((u, v)\) in \( H \), and remove the vertex \( u \), together with its incident edges, from \( H \).

**Step 4:** Let \( I \) be the vertex set of the graph \( H \). Return \( I \).

- Argue that the set \( I \) that is returned by this algorithm is an independent set in \( G \).
- Let \( X \) and \( Y \) be the random variables whose values are the number of vertices and edges in the graph \( H \) after Step 2, respectively. Prove that
  \[
  \mathbb{E}(X) = \frac{n^2}{2m}
  \]
  and
  \[
  \mathbb{E}(Y) = \frac{n^2}{4m}.
  \]
- Let \( Z \) be the random variable whose value is the size of the independent set \( I \) that is returned by the algorithm. Argue that
  \[
  Z \geq X - Y.
  \]
- Prove that
  \[
  \mathbb{E}(Z) \geq \frac{n^2}{4m}.
  \]
- Argue that this implies that the graph \( G \) contains an independent set of size at least \( \frac{n^2}{4m} \).