Random processes are introduced in the context of several simple physical systems, including random walks on a lattice, polymers, and diffusion controlled chemical reactions. The generation of random number sequences also is discussed.

7.1 Order to Disorder

In Chapter 6 we saw several examples of how under certain conditions, the behavior of a nonlinear deterministic system can appear to be random. In this chapter we will see some examples of how chance can generate statistically predictable outcomes. For example, we know that if we bet often on the outcome of a game for which the probability of winning is less than 50%, we will lose money eventually.

We first discuss an example that illustrates the tendency of systems of many particles to evolve to a well defined state. Imagine a closed box that is divided into two parts of equal volume (see Figure 7.1). The left half contains a gas of $N$ identical particles and the right half is initially empty. We then make a small hole in the partition between the two halves. What happens? We know that after some time, the average number of particles in each half of the box will become $N/2$, and we say that the system has reached equilibrium.

How can we simulate this process? One way is to give each particle an initial velocity and position and adopt a deterministic model of the motion of the particles. For example, we could assume that each particle moves in a straight line until it hits a wall of the box or another particle and undergoes an elastic collision. We will consider similar deterministic models in Chapter 9. Instead, we first simulate a probabilistic model based on a random process.

The basic assumptions of this model are that the motion of the particles is random and the particles do not interact with one another. Hence, the probability per unit time that a particle goes through the hole in the partition is the same for all $N$ particles regardless of the number of
particles in either half. We also assume that the size of the hole is such that only one particle can pass through at a time. We first model the motion of a particle passing through the hole by choosing one of the \( N \) particles at random and moving it to the other side. For visualization purposes we will use arrays to specify the position of each particle. We then randomly generate an integer \( i \) between 0 and \( N - 1 \) and change the arrays appropriately. A more efficient Monte Carlo algorithm is discussed in Problem 7.2b. The tool we need to simulate this random process is a random number generator.

It is counterintuitive that we can use a deterministic computer to generate sequences of random numbers. In Section 7.9 we discuss some of the methods for computing a set of numbers that appear statistically random, but are in fact generated by a deterministic algorithm. These algorithms are sometimes called pseudorandom number generators to distinguish their output from intrinsically random physical processes such as the time between clicks in a Geiger counter near a radioactive sample.

For the present we will be content to use the random number generator supplied with Java, although the random number generators included with various programming languages vary in quality. The method \( \text{Math.random()} \) produces a random number \( r \) that is uniformly distributed in the interval \( 0 \leq r < 1 \). To generate a random integer \( i \) between 0 and \( N - 1 \), we write:

```java
int i = (int)(N*Math.random());
```

The effect of the \((\text{int})\) cast is to eliminate the decimal digits from a floating point number. For example, \((\text{int})(5.7) = 5\).

The algorithm for simulating the evolution of the model can be summarized by the following steps:

1. Use a random number generator to choose a particle at random.
2. Move this particle to the other side of the box.
3. Give the particle a random position on the new side of the box. This step is for visualization purposes only.
4. Increase the “time” by unity.
Note that this definition of time is arbitrary. Class \texttt{Box} implements this algorithm and class \texttt{BoxApp} plots the evolution of the number of particles on the left half of the box.

Listing 7.1: Class \texttt{Box} for the simulation of the approach to equilibrium.

```java
package org.opensourcephysics.sip.ch07;
import java.awt.*;
import org.opensourcephysics.display.*;

public class Box implements Drawable {
    public double x[], y[];
    public int N, nleft, time;

    public void initialize() {
        x = new double[N]; // location of particles (for visualization purposes only)
        y = new double[N];
        nleft = N; // start with all particles on the left
        time = 0;
        for (int i = 0; i < N; i++) {
            x[i] = 0.5*Math.random(); // needed only for visualization
            y[i] = Math.random();
        }
    }

    public void step() {
        int i = (int)(Math.random()*N);
        if (x[i] < 0.5) {
            nleft--; // move to right
            x[i] = 0.5*(1+Math.random());
            y[i] = Math.random();
        } else {
            nleft++; // move to left
            x[i] = 0.5*Math.random();
            y[i] = Math.random();
        }
        time++;
    }

    public void draw(DrawingPanel panel, Graphics g) {
        if (x==null) {
            return;
        }
        int size = 2;
        int xMiddle = panel.xToPix(0.5); // position of partition in middle of box
        g.setColor(Color.black);
        g.drawLine(xMiddle, panel.yToPix(0), xMiddle, panel.yToPix(0.45));
        g.drawLine(xMiddle, panel.yToPix(0.55), xMiddle, panel.yToPix(1.0));
        g.setColor(Color.red);
        for (int i = 0; i < N; i++) {
            int xpix = panel.xToPix(x[i]);
            int ypix = panel.yToPix(y[i]);
```
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Listing 7.2: Target class for plotting the approach to equilibrium.

```java
package org.opensourcephysics.sip.ch07;
import org.opensourcephysics.controls.*;
import org.opensourcephysics.frames.*;

public class BoxApp extends AbstractSimulation {
    Box box = new Box();
    PlotFrame plotFrame = new PlotFrame("time", "number on left", "Box data");
    DisplayFrame displayFrame = new DisplayFrame("Partitioned box");

    public void initialize() {
        displayFrame.clearDrawables();
        displayFrame.addDrawable(box);
        box.N = control.getInt("Number of particles");
        box.initialize();
        plotFrame.clearData();
        displayFrame.setPreferredMinMax(0, 1, 0, 1);
    }

    public void doStep() {
        box.step();
        plotFrame.append(0, box.time, box.nleft);
    }

    public void reset() {
        // clicking reset should erase positions of particles
        control.setValue("Number of particles", 64);
        plotFrame.clearData();
        enableStepsPerDisplay(true);
        setStepsPerDisplay(10);
    }

    public static void main(String[] args) {
        SimulationControl.createApp(new BoxApp());
    }
}
```

How long does it take for the system to reach equilibrium? How does this time depend on
the number of particles? After the system reaches equilibrium, what is the magnitude of the
fluctuations? How do the fluctuations depend on the number of particles? Problems 7.2 and 7.3
explore such questions.

Exercise 7.1. Simple tests of operators and methods
a. It is frequently quicker to write a short program to test how the operators and methods of a computer language work than to look them up in a manual or online. Write a test class to determine the values of \(3/2\), \(3.0/2.0\), \((\text{int}) (3.0/2.0)\), \(2/3\), and \((\text{int}) (-3/2)\).

b. Determine the behavior of \(\text{Math.round}(\text{arg})\), \(\text{Math.ceil}(\text{arg})\), and \(\text{Math.rint}(\text{arg})\).

c. Write a program to test whether the same sequence of random numbers appears each time the program is run if we use the method \(\text{Math.random}()\) to generate the sequence.

d. Create an object from class \(\text{Random}\) and use the methods \(\text{setSeed}(\text{long seed})\) and \(\text{nextDouble}()\). Show that you obtain the same sequence of random numbers if the same seed is used. One reason to specify the seed rather than to choose it at random from the time (as is the default) is that it is convenient to use the same random number sequence when testing a program. Suppose that your program gives a strange result for a particular run. If you notice an error in the program and change the program, you would want to use the same random number sequence to test whether your changes corrected the error. Another reason for specifying the seed is that another user could obtain the same results if you tell them the seed that you used.

**Problem 7.2.** Approach to equilibrium

a. Use \(\text{BoxApp}\) and \(\text{Box}\) and describe the nature of the evolution of \(n\), the number of particles on the left side of the box. Choose the total number of particles \(N\) to be \(N = 8, 16, 64, 400, 800, \text{ and } 3600\). Does the system reach equilibrium? What is your qualitative criterion for equilibrium? Does \(n\), the number of particles on the left-hand side, change when the system is in equilibrium?

b. The algorithm we have used is needlessly cumbersome, because our only interest is the number of particles on each side. We used the positions only for visualization purposes. Because each particle has the same chance to go through the hole, the probability per unit time that a particle moves from left to right equals the number of particles on the left divided by the total number of particles, that is, \(p = n/N\). Modify the program so that the following algorithm is implemented.

(i) Generate a random number \(r\) from a uniformly distributed set of random numbers in the interval \(0 \leq r < 1\).

(ii) If \(r \leq p = n/N\), move a particle from left to right, that is \(n \rightarrow n-1\); otherwise, \(n \rightarrow n+1\).

c. Does the time dependence of \(n\) appear to be deterministic for sufficiently large \(N\)? What is the qualitative behavior of \(n(t)\)? Estimate the time for the system to reach equilibrium from the plots. How does this time depend on \(N\)?

**Problem 7.3.** Equilibrium fluctuations

a. As a rough measure of the equilibrium fluctuations, visually estimate the deviation of \(n(t)\) from \(N/2\) for \(N = 16, 64, 400, 800, \text{ and } 3600\)? Choose a time interval that is bigger than the time needed to reach equilibrium. How do your results for the deviation depend on \(N\)?

b. A better measure of the equilibrium fluctuations is the mean square fluctuations \(\Delta n^2\), which is defined as

\[
\Delta n^2 = \langle (n - \langle n \rangle)^2 \rangle = \langle n^2 \rangle - 2\langle n \langle n \rangle \rangle + \langle n \rangle^2 = \langle n^2 \rangle - 2\langle n \rangle + \langle n \rangle^2 = \langle n^2 \rangle - \langle n \rangle^2.
\]  

(7.1)
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The brackets, \( \langle \cdots \rangle \), denote an average taken after the system has reached equilibrium. The relative magnitude of the fluctuations is \( \Delta n / \langle n \rangle \). Modify your program so that averages are taken after equilibrium has been reached. Run for a time that is long enough to obtain meaningful results. Compute the mean square fluctuations \( \Delta n^2 \) for the same values of \( N \) considered in part (a). How do the relative fluctuations, \( \Delta n / \langle n \rangle \), depend on \( N \)? (You might find it helpful to see how averages are computed in Listings 7.3 and 7.4.)

From Problem 7.2 we see that \( n(t) \) decreases in time from its initial value to its equilibrium value in an almost deterministic manner if \( N \gg 1 \). It is instructive to derive the time dependence of \( n(t) \) to show explicitly how chance can generate deterministic behavior. If there are \( n(t) \) particles on the left side after \( t \) moves, then the change in \( \langle n \rangle (t) \) in the time interval \( \Delta t \) is given by

\[
\Delta \langle n \rangle = \left[ \frac{\langle n \rangle(t)}{N} - \frac{N - \langle n \rangle(t)}{N} \right] \Delta t. \tag{7.2}
\]

(We defined the time so that the time interval \( \Delta t = 1 \) in our simulations.) What is the meaning of the two terms in (7.2)? If we treat \( \langle n \rangle \) and \( t \) as continuous variables and take the limit \( \Delta t \to 0 \), we have

\[
\frac{\Delta \langle n \rangle}{\Delta t} \to \frac{d\langle n \rangle}{dt} = 1 - \frac{2\langle n \rangle(t)}{N}. \tag{7.3}
\]

The solution of the differential equation (7.3) is

\[
\langle n \rangle(t) = \frac{N}{2} [1 + e^{-2t/N}], \tag{7.4}
\]

where we have used the initial condition \( \langle n \rangle(t = 0) = N \). Note that \( \langle n \rangle(t) \) decays exponentially to its equilibrium value \( N/2 \). How does this form (7.4) compare to your simulation results for various values of \( N \)? We can define a relaxation time \( \tau \) as the time it takes the difference \( \langle n \rangle(t) - N/2 \) to decrease to \( 1/e \) of its initial value. How does \( \tau \) depend on \( N \)? Does this prediction for \( \tau \) agree with your results from Problem 7.2?

*Problem 7.4. A simple modification*

Modify your program so that each side of the box is chosen with equal probability. One particle is then moved from the side chosen to the other side. If the side chosen does not have a particle, then no particle is moved during this time interval. Do you expect that the system behaves in the same way as before? Do the simulation starting with all the particles on the left side of the box and choose \( N = 800 \). Do not keep track of the positions of the particles. Compare the behavior of \( n(t) \) with the behavior of \( n(t) \) found in Problem 7.3. How do the values of \( \langle n \rangle \) and \( \Delta n \) compare?

The probabilistic method discussed on page 206 for simulating the approach to equilibrium is an example of a Monte Carlo algorithm, that is, the random sampling of the most probable outcomes. An alternative method is to use exact enumeration and determine all the possibilities at each time interval. For example, suppose that \( N = 8 \) and \( n(t = 0) = 8 \). At \( t = 1 \), the only possibility is \( n = 7 \) and \( n' = 1 \). Hence, \( P(n = 7, t = 1) = 1 \) and all other probabilities are zero. At \( t = 2 \), one of the seven particles on the left can move to the right, or the one particle on the right can move to the left. Because the first possibility can occur in seven different ways, we have
the nonzero probabilities, \( P(n = 6, t = 2) = 7/8 \) and \( P(n = 8, t = 2) = 1/8 \). Hence at \( t = 2 \), the average number of particles on the left side of the box is

\[
\langle n(t = 2) \rangle = 6P(n = 6, t = 2) + 8P(n = 8, t = 2) = \frac{7}{8}[6 \times 7 + 8 \times 1] = 6.25.
\]

Is this exact result consistent with what you found in Problem 7.2? In this example \( N \) is small, and we could continue the enumeration of all the possibilities indefinitely. However for larger \( N \), the number of possibilities becomes very large after a few time intervals, and we need to use Monte Carlo methods to sample the most probable outcomes.

### 7.2 Random Walks

In Section 7.1 we considered the random motion of many particles in a box, but we did not care about their positions – all we needed to know was the number of particles on each side. Suppose that we want to characterize the motion of a dust particle in the atmosphere. We know that as a given dust particle collides with molecules in the atmosphere, it changes its direction frequently, and its motion appears to be random. A simple model for the trajectory of a dust particle in the atmosphere is based on the assumption that the particle moves in any direction with equal probability. Such a model is an example of a random walk.

The original statement of a random walk was formulated in the context of a drunken sailor. If a drunkard begins at a lamp post and takes \( N \) steps of equal length in random directions, how far will the drunkard be from the lamp post? We will find that the mean square displacement of a random walker, for example, a dust particle or a drunkard, grows linearly with time. This result and its relation to diffusion leads to many applications that might seem to be unrelated to the original drunken sailor problem.

We first consider an idealized example of a random walker that can move only along a line. Suppose that the walker begins at \( x = 0 \) and that each step is of equal length \( a \). At each time interval the walker has a probability \( p \) of a step to the right and a probability \( q = 1 - p \) of a step to the left. The direction of each step is independent of the preceding one. After \( N \) steps the displacement \( x \) of the walker from the origin is given by

\[
x_N = \sum_{i=1}^{N} s_i,
\]

where \( s_i = \pm a \). For \( p = 1/2 \) we can generate one walk of \( N \) steps by flipping a coin \( N \) times and increasing \( x \) by \( a \) each time the coin is heads and decreasing \( x \) by \( a \) each time the coin is tails.

We expect that if we average over a sufficient number of walks of \( N \) steps, then the average of \( x_N \), denoted by \( \langle x_N \rangle \), would be \( (p - q)Na \). We can derive this result by writing \( \langle x \rangle = \sum_{i=1}^{N} \langle s_i \rangle = N\langle s \rangle \), because the average of the sum is the sum of the averages, and the average of each step is the same. We have \( \langle s \rangle = p(a) + q(-a) = (p - q)a \), and the result follows. For simplicity, we will frequently drop the subscript \( N \) and write \( \langle x \rangle \).

Because \( \langle x \rangle = 0 \) for \( p = 1/2 \), we need a better measure of the extent of the walk. One measure
is the displacement squared:

\[ x_N^2 = \left[ \sum_{i=1}^{N} s_i \right]^2. \]  

(7.6)

For \( p \neq 1/2 \), it is convenient to consider the mean square net displacement \( \Delta x^2 \) defined as

\[ \Delta x^2 = \langle (x - \langle x \rangle)^2 \rangle = \langle x^2 \rangle - \langle x \rangle^2. \]  

(7.7a)

(We have written \( \Delta x^2 \) rather than \( \langle \Delta x^2 \rangle \) for simplicity.) To determine \( \Delta x^2 \), we write (7.7a) as the sum of two terms:

\[ \Delta x^2 = \langle \sum_{i=1}^{N} \Delta_i \sum_{j=1}^{N} \Delta_j \rangle = \langle \sum_{i=1}^{N} \Delta_i^2 \rangle + \langle \sum_{i \neq j=1}^{N} \Delta_i \Delta_j \rangle, \]  

(7.8)

where \( \Delta_i = s_i - \langle s \rangle \). The first sum on the right-hand side of (7.8) includes terms for which \( i = j \); the second sum is over \( i \) and \( j \) with \( i \neq j \). Because each step is independent, we have \( \langle \Delta_i \Delta_j \rangle = \langle \Delta_i \rangle \langle \Delta_j \rangle \). This term is zero because \( \langle \Delta_i \rangle = 0 \) for any \( i \). The first term in (7.8) equals

\[ N \langle \Delta_i^2 \rangle = N \langle s^2 \rangle - \langle s \rangle^2 = N [ a^2 -(p-q)^2 a^2 ] = N 4pq a^2, \]  

where we have used the fact that \( p + q = 1 \). Hence

\[ \Delta x^2 = 4pq Na^2. \]  

(analytical result)  

(7.9)

We can gain more insight into the nature of random walks by doing a Monte Carlo simulation, that is, by using a computer to “flip coins.” The implementation of the random walk algorithm is simple, for example,

```java
if (p < Math.random()) {
    x++;
} else {
    x--;
}
```

Clearly we have to sample many \( N \) step walks because in general, each walk will give a different outcome. We need to do a Monte Carlo simulation many times and average over the results to obtain meaningful averages. Each \( N \)-step walk is called a trial. How do we know how many trials to use? The simple answer is to average over more and more trials until the average results don’t change within the desired accuracy. The more sophisticated answer is to do an error analysis similar to what we do in measurements in the laboratory. Such an analysis is discussed in Section 12.4.

The more difficult parts of a program to simulate random walks are associated with bookkeeping. The walker takes a total of \( N \) steps in each trial and the net displacement \( x \) is computed after every step. Our convention will be to use a variable name ending in \texttt{Accumulator} (or \texttt{Accum}) to denote a variable that accumulates the value of some variable. In Listing 7.3 we provide two classes to be used to simulate random walks.
Listing 7.3: Listing of Walker class.

```java
package org.opensourcephysics.sip.ch07;

public class Walker {
    int xAccum[], xSquaredAccum[]; // accumulated data on displacement of walkers, index is time
    int N; // maximum number of steps
    double p; // probability of step to the right
    int position; // position of walker

    public void initialize() {
        xAccum = new int[N+1];
        xSquaredAccum = new int[N+1];
    }

    public void step() {
        position = 0;
        for(int t = 0; t < N; t++) {
            if(Math.random() < p) {
                position++;
            } else {
                position--;
            }
            xAccum[t+1] += position; // determine displacement of walker after each step
            xSquaredAccum[t+1] += position*position;
        }
    }
}
```

Listing 7.4: Target class for random walk simulation.

```java
package org.opensourcephysics.sip.ch07;
import org.opensourcephysics.controls.s;
import org.opensourcephysics.frames.s;

public class WalkerApp extends AbstractSimulation {
    Walker walker = new Walker();
    PlotFrame plotFrame = new PlotFrame("time", ",x,,<x^2>", "Averages");
    HistogramFrame distribution = new HistogramFrame("x", "H(x)", "Histogram");
    int trials; // number of trials

    public WalkerApp() {
        plotFrame.setXYColumnNames(0, "t", "<x>");
        plotFrame.setXYColumnNames(1, "t", "<x^2>");
    }

    public void initialize() {
        walker.p = control.getDouble("Probability p of step to right");
        walker.N = control.getInt("Number of steps N");
        walker.initialize();
        trials = 0;
    }
}
Problem 7.5. Random walks in one dimension

a. In class Walker the steps are of unit length so that $a = 1$. Use Walker and WalkerApp to estimate the number of trials needed to obtain $\Delta x^2$ for $N = 20$ and $p = 1/2$ with an accuracy of approximately 5%. Compare your result for $\Delta x^2$ to the exact answer in (7.9). Approximately how many trials do you need to obtain the same relative accuracy for $N = 100$?

b. Is $\langle x \rangle$ exactly zero in your simulations? Explain the difference between the analytical result and the results of your simulations. Note that we have used the same notation $\langle \ldots \rangle$ to denote the exact average calculated analytically and the approximate average computed by averaging over many trials. The distinction between the two averages should be clear from the context.

c. How do your results for $\langle x \rangle$ and $\Delta x^2$ change for $p \neq q$? Choose $p = 0.7$ and determine the $N$ dependence of $\langle x \rangle$ and $\Delta x^2$.

d. Determine $\Delta x^2$ for $N = 1$ to $N = 5$ by enumerating all the possible walks. For simplicity, choose $p = 1/2$ so that $\langle x \rangle = 0$. For $N = 1$, there are two possible walks: one step to the right and one step to the left. In both cases $x^2 = 1$ and hence $\langle x_1^2 \rangle = 1$. For $N = 2$ there are four possible walks with the same probability: (i) two steps to the right, (ii) two steps to the left, (iii) first step to the right and second step to the left, and (iv) first step to the left and second
step to the right. The value of \( x_2^2 \) for these walks is 4, 4, 0, and 0 respectively, and hence \( \langle x_2^2 \rangle = (4 + 4 + 0 + 0)/4 = 2 \). Write a program that enumerates all the possible walks of a given number of steps and compute the various averages of interest exactly.

The class WalkerApp displays the distribution of values of the displacement \( x \) after \( N \) steps. One way of determining the number of times that the variable \( x \) has a certain value would be to define a one-dimensional array, \( \text{probability} \), and let

\[
\text{probability}[x] += 1;
\]

In this case because \( x \) takes only integer values, the array index of \( \text{probability} \) is the same as \( x \) itself. However, the above statement does not work in Java because \( x \) can be negative as well as positive. What we need is a way of mapping the value \( x \) to a bin or index number. The HistogramFrame class, which is part of the Open Source Physics display package, does this mapping automatically using the Java Hashtable class. In simple data structures data is accessed by an index that indicates the location of the data in the data structure. Hashtable data is accessed by a key, which in our case is the value of \( x \). A hashing function converts the key to an index. The append method of the HistogramFrame class takes a value, finds the index using a hashing function, and then increments the data associated with that key. The HistogramFrame class also draws itself.

The HistogramFrame class is very useful for taking a quick look at the distribution of values in a data set. You do not need to know how to group the data into bins or the range of values of the data. The default bin width is unity, but the bin width can be set using the setBinWidth method. See WalkerApp for an example of the use of the HistogramFrame class. Frequently, we wish to use the histogram data to compute other quantities. You can collect the data using the Data Table menu item in HistogramFrame and copy the data to a file. Another option is to include additional code in your program to analyze the data. The following statements assume that a HistogramFrame object called histogram has been created and data entered into it.

```java
// creates array entries of data from histogram.
java.util.Map.Entry[] entries = histogram.entries();
for (int i = 0, length = entries.length; i < length; i++) {
    Integer binNumber = (Integer) entries[i].getKey(); // gets bin number
    Double occurrences = (Double) entries[i].getValue(); // gets number of occurrences for bin
    // gets value of left edge of bin
    double value = histogram.getLeftMostBinPosition(binNumber.intValue()); // sets value to middle of bin
    double number = occurrences.doubleValue(); // convert from Double class to double data type
    // use value and number in your analysis
}
```

Problem 7.6. Nature of the probability distribution

a. Compute \( P_N(x) \), the probability that the displacement of the walker from the origin is \( x \) after \( N \) steps. What is the difference between the histogram, that is, the number of occurrences, and the probability? Consider \( N = 10 \) and \( N = 40 \) and at least 1000 trials. Does the qualitative form of \( P_N(x) \) change as the number of trials increases? What is the approximate width of \( P_N(x) \) and the value of \( P_N(x) \) at its maximum for each value of \( N \)?
b. What is the approximate shape of the envelope of $P_N(x)$? Does the shape change as $N$ is increased?

c. Fit the envelope $P_N(x)$ for sufficiently large $N$ to the continuous function

$$C \frac{1}{\sqrt{2\pi \Delta x^2}} e^{-(x-\langle x \rangle)^2/2\Delta x^2}. \quad (7.10)$$

The form of (7.10) is the standard form of the Gaussian distribution with $C = 1$. The easiest way to do this fit is to plot your results for $P_N(x)$ and the form (7.10) on the same graph using your results for $\langle x \rangle$ and $\Delta x^2$ as input parameters. Visually choose the constant $C$ to obtain a reasonable fit. What are the possible values of $x$ for a given value of $N$? What is the minimum difference between these values? How does this difference compare to your value for $C$?

Problem 7.7. More random walks in one dimension

a. Suppose that the probability of a step to the right is $p = 0.7$. Compute $\langle x \rangle$ and $\Delta x^2$ for $N = 4$, 8, 16, and 32. What is the interpretation of $\langle x \rangle$ in this case? What is the qualitative dependence of $\Delta x^2$ on $N$?

b. An interesting property of random walks is the mean number $\langle D_N \rangle$ of distinct lattice sites visited during the course of an $N$ step walk. Do a Monte Carlo simulation of $\langle D_N \rangle$ and determine its $N$ dependence.

We can equally well consider either a large number of successive walks as in Problem 7.7 or a large number of noninteracting walkers moving at the same time as in Problem 7.8.
Problem 7.8. A random walk in two dimensions

a. Consider a collection of walkers initially at the origin of a square lattice (see Figure 7.2). At each unit of time, each of the walkers moves at random with equal probability in one of the four possible directions. Create a drawable class, Walker2D, which contains the positions of \( M \) walkers moving in two dimensions and draws their location, and modify WalkerApp. Unlike WalkerApp, this new class need not specify the maximum number of steps. Instead the number of walkers should be specified.

b. Run your application with the number of walkers \( M \geq 1000 \) and allow the walkers to take at least 500 steps. If each walker represents a bee, what is the qualitative nature of the shape of the swarm of bees? Describe the qualitative nature of the surface of the swarm as a function of the number of steps, \( N \). Is the surface jagged or smooth?

c. Compute the quantities \( \langle x \rangle \), \( \langle y \rangle \), \( \Delta x^2 \), and \( \Delta y^2 \) as a function of \( N \). The average is over the \( M \) walkers. Also compute the mean square displacement \( R^2 \) given by

\[
R^2 = \langle x^2 \rangle - \langle x \rangle^2 + \langle y^2 \rangle - \langle y \rangle^2 = \Delta x^2 + \Delta y^2. \tag{7.11}
\]

What is the dependence of each quantity on \( N \)? (As before, we will frequently write \( R^2 \) instead of \( R^2_N \).)

d. Estimate \( R^2 \) for \( N = 8, 16, 32, \) and 64 by averaging over a large number of walkers for each value of \( N \). Assume that \( R = \sqrt{R^2} \) has the asymptotic \( N \) dependence:

\[
R \sim N^\nu, \quad (N \gg 1) \tag{7.12}
\]

and estimate the exponent \( \nu \) from a log-log plot of \( R^2 \) versus \( N \). We will see in Chapter 14 that the exponent \( 1/\nu \) is related to how a random walk fills space. If \( \nu \approx 1/2 \), estimate the magnitude of the self-diffusion coefficient \( D \) from the relation \( R^2 = 4DN \).

e. Do a Monte Carlo simulation of \( R^2 \) on a triangular lattice (see Figure 9.5) and estimate \( \nu \). Can you conclude that the exponent \( \nu \) is independent of the symmetry of the lattice? Does \( D \) depend on the symmetry of the lattice? If so, give a qualitative explanation for this dependence.

f. Enumerate all the random walks on a square lattice for \( N = 4 \) and obtain exact results for \( \langle x \rangle \), \( \langle y \rangle \), and \( R^2 \). Assume that all four directions are equally probable. Verify your program by comparing the Monte Carlo and exact enumeration results.

Problem 7.9. The fall of a rain drop

Consider a random walk that starts at a site a distance \( y = h \) above a horizontal line (see Figure 7.3). If the probability of a step down is greater than the probability of a step up, we expect that the walker will eventually reach a site on the horizontal line. This walk is a simple model of the fall of a rain drop in the presence of a random swirling breeze. Do a Monte Carlo simulation to determine the mean time \( \tau \) for the walker to reach any site on the line \( y = 0 \) and find the functional dependence of \( \tau \) on \( h \). Is it possible to define a velocity in the vertical direction? Because the walker does not always move vertically, it suffers a net displacement \( x \) in the horizontal direction. How does \( \Delta x^2 \) depend on \( h \) and \( \tau \)? Reasonable values for the step probabilities are 0.1, 0.6, 0.15, 0.15, corresponding to up, down, right, and left, respectively.
Figure 7.3: Examples of the random path of a raindrop to the ground. The step probabilities are given in Problem 7.9. The walker starts at $x = 0, y = h$.

### 7.3 Modified Random Walks

So far we have considered random walks on one- and two-dimensional lattices where the walker has no “memory” of the previous step. What happens if the walkers remember the nature of their previous steps? What happens if there are multiple random walkers, with the condition that no double occupancy is allowed? We explore these and other variations of the simple random walk in this section. All these variations have applications to physical systems, but the applications are more difficult to understand than the models themselves.

The fall of a raindrop considered in Problem 7.9 is an example of a restricted random walk, that is, a walk in the presence of a boundary. In the following problem, we discuss in a more general context the effects of various types of restrictions or boundaries on random walks. Other examples of a restricted random walk are given in Problems 7.17 and 7.23.

**Problem 7.10.** Restricted random walks

a. Consider a one-dimensional lattice with trap sites at $x = 0$ and $x = L$ ($L > 0$). A walker begins at site $x_0$ ($0 < x_0 < L$) and takes unit steps to the left and right with equal probability. When the walker arrives at a trap site, it can no longer move. Do a Monte Carlo simulation and verify that the mean number of steps $\tau$ for the particle to be trapped (the mean first passage time) is given by

$$\tau = (2D)^{-1}x_0(L - x_0), \quad (7.13)$$

where $D$ is the self-diffusion coefficient in the absence of the traps, and the average is over all possible walks.

b. Random walk models in the presence of traps have had an important role in condensed matter physics. For example, consider the following idealized model of energy transport in solids. The solid is represented as a lattice with two types of sites: hosts and traps. An incident photon is absorbed at a host site and excites the host molecule or atom. The excitation energy or exciton is transferred at random to one of the host’s nearest neighbors and the original excited molecule
returns to its ground state. In this way the exciton wanders through the lattice until it reaches a trap site at which a chemical reaction occurs. A simple version of this energy transport model is given by a one-dimensional lattice with traps placed on a periodic sublattice. Because the traps are placed at regular intervals, we can replace the random walk on an infinite lattice by a random walk on a circular ring. Consider a lattice of $N$ host or nontrapping sites and one trap site. If a walker has an equal probability of starting from any host site and an equal probability of a step to each nearest neighbor site, what is the $N$ dependence of the mean survival time $\tau$ (the mean number of steps taken before a trap site is reached)? Use the results of part (a) rather than doing a simulation.

c. Consider a one-dimensional lattice with reflecting sites at $x = -L$ and $x = L$. For example, if a walker reaches the reflecting site at $x = L$, it is reflected at the next step to $x = L - 1$. At $t = 0$, the walker starts at $x = 0$ and steps with equal probability to the left and right. Write a Monte Carlo program to determine $P_N(x)$, the probability that the walker is at site $x$ after $N$ steps. Compare the form of $P_N(x)$ with and without the presence of the reflecting sites. Can you distinguish the two probability distributions if $N$ is the order of $L$? At what value of $N$ can you first distinguish the two distributions?

Problem 7.11. A persistent random walk

a. In a persistent random walk, the transition or jump probability depends on the previous step. Consider a walk on a one-dimensional lattice, and suppose that step $N - 1$ has been made. Then step $N$ is made in the same direction with probability $\alpha$; a step in the opposite direction occurs with probability $1 - \alpha$. Write a program to do a Monte Carlo simulation of the persistent random walk in one dimension. Estimate $\langle x \rangle$, $\Delta x^2$, and $P_N(x)$. Note that it is necessary to specify both the initial position and an initial direction of the walker. What is the $\alpha = 1/2$ limit of the persistent random walk?

b. Consider $\alpha = 0.25$ and $\alpha = 0.75$ and determine $\Delta x^2$ for $N = 8, 64, 256, \text{ and } 512$. Assume that $\Delta x^2 \sim N^{2\nu}$ for large $N$, and estimate the value of $\nu$ from a log-log plot of $\Delta x^2$ versus $N$ for large $N$. Does $\nu$ depend on $\alpha$? If $\nu \approx 1/2$, determine the self-diffusion coefficient $D$ for $\alpha = 0.25$ and 0.75. In general, $D$ is given by

$$D = \frac{1}{2d} \lim_{N \to \infty} \frac{\Delta x^2}{N},$$

(7.14)

where $d$ is the dimension of space. That is, $D$ is given by the asymptotic behavior of the mean square displacement. (For the simple random walk considered in Section 7.2, $\Delta x^2 \propto N$ for all $N$.) Give a physical argument why $D(\alpha \neq 0.5)$ is greater (smaller) than $D(\alpha = 0.5)$.

c. You might have expected that the persistent random walk yields a nonzero value for $\langle x \rangle$. Verify that $\langle x \rangle = 0$, and explain why this result is exact. How does the persistent random walk differ from the biased random walk for which $p \neq q$?

d. A persistent random walk can be considered as an example of a multistate walk in which the state of the walk is defined by the last transition. The walker is in one of two states; at each step the probabilities of remaining in the same state or switching states are $\alpha$ and $1 - \alpha$, respectively.
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One of the earliest applications of a two state random walk was to the study of diffusion in a chromatographic column. Suppose that a molecule in a chromatographic column can be either in a mobile phase (constant velocity $v$) or in a trapped phase (zero velocity). Instead of each step changing the position by $\pm 1$, the position at each step changes by $+v$ or 0. A quantity of experimental interest is the probability $P_N(x)$ that a molecule has moved a distance $x$ in $N$ steps. Choose $v = 1$ and $\alpha = 0.75$ and determine the qualitative behavior of $P_N(x)$.

Problem 7.12. Synchronized random walks

a. Randomly place two walkers on a one-dimensional lattice of $L$ sites, so that both walkers are not at the same site. At each time step randomly choose whether the walkers move to the left or to the right. Both walkers move in the same direction. If a walker cannot move in the chosen direction because it is at a boundary, then this walker remains at the same site for this time step. A trial ends when both walkers are at the same site. Write a program to determine the mean time and the mean square fluctuations of the time for two walkers to reach the same site. This model is relevant to a method of doing cryptography using neural networks (see Rutter et al.).

b. Change your program so that you use biased random walkers for which $p \neq q$. How does this change affect your results?

Problem 7.13. Random walk on a continuum

One of the first continuum models of a random walk was proposed by Rayleigh in 1919. In this model the length $a$ of each step is a random variable and the direction of each step is uniformly random. In this case the variable of interest is $R$, the distance of the walker from the origin after $N$ steps. The model is known as the freely jointed chain in polymer physics (see Section 7.7) in which case $R$ is the end-to-end distance of the polymer. For simplicity, we first consider a walker in two dimensions with steps of equal (unit) length at a random angle.

a. Write a Monte Carlo program to compute $\langle R \rangle$ and determine its dependence on $N$.

b. Because $R$ is a continuous variable, we need to compute $p_N(R)\Delta R$, the probability that $R$ is between $R$ and $R + \Delta R$ after $N$ steps. The quantity $p_N(R)$ is the probability density. Because the area of the ring between $R$ and $R + \Delta R$ is $\pi(R + \Delta R)^2 - \pi R^2 = 2\pi R \Delta R + \pi(\Delta R)^2 \approx 2\pi R \Delta R$, we see that $p_N(R)\Delta R$ is proportional to $R\Delta R$. Verify that for sufficiently large $N$ that $p_N(R)\Delta R$ has the form

$$p_N(R)\Delta R \propto 2\pi R \Delta R e^{-(R-\langle R \rangle)^2/2\Delta R^2}, \quad (7.15)$$

where $\Delta R^2 = \langle R^2 \rangle - \langle R \rangle^2$.

Problem 7.14. Random walks with steps of variable length

a. Consider a random walk in one dimension with jumps of all lengths. The probability that the length of a single step is between $a$ and $a + \Delta a$ is $f(a)\Delta a$, where $f(a)$ is the probability density. If the form of $f(a)$ is given by $f(a) = Ce^{-a}$ for $a > 0$ with the normalization condition $\int_0^\infty f(a)da = 1$, the code needed to generate step lengths according to this probability density is given by (see Section 12.5)
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stepLength = -Math.log(1 - Math.random());

Modify Walker and WalkerApp to simulate walks of variable length with this probability density. Note that the bin width ∆a is one of the input parameters. Consider N ≥ 100 and visualize the motion of the walker. Generate many walks of N steps and determine p(x)Δx, the probability that the displacement is between x and x + Δx after N steps. Plot p(x) versus x and confirm that the form of p(x) is consistent with a Gaussian distribution.

b. Assume that the probability density f(a) is given by f(a) = C/a^2 for a ≥ 1. Determine the normalization constant C using the condition C ∫[1^∞] a^-2 da = 1. In this case, we will learn in Section 12.5 that the statement

stepLength = 1.0 / (1.0 - Math.random());

generates values of a according to this form of f(a). Do a Monte Carlo simulation as in part (a) and determine p(x)Δx. Is the form of p(x) a Gaussian? This type of random walk for which f(a) decreases as a power law, a^{-1-α}, is known as a Levy flight for α ≤ 2.

Problem 7.15. Exploring the central limit theorem

Consider a continuous random variable x with probability density f(x). That is, f(x)Δx is the probability that x has a value between x and x + Δx. The mth moment of f(x) is defined as

⟨x^m⟩ = ∫x^m f(x) dx. (7.16)

The mean value ⟨x⟩ is given by (7.16) with m = 1. The variance σ_x^2 of f(x) is defined as

σ_x^2 = ⟨x^2⟩ - ⟨x⟩^2. (7.17)

Consider the sum y_n corresponding to the average of n values of x:

y = y_n = 1/n (x_1 + x_2 + ... + x_n). (7.18)

Suppose that we make many measurements of y. We know that the values of y will not be identical, but will be distributed according to a probability density p(y), where p(y)Δy is the probability that the measured value of y is in the range y to y + Δy. The main quantities of interest are ⟨y⟩, p(y), and an estimate of the probable variability of y in a series of measurements.

a. Suppose that f(x) is uniform in the interval [−1,1]. Calculate ⟨x⟩, ⟨x^2⟩, and σ_x analytically.

b. Write a program to make a sufficient number of measurements of y and determine ⟨y⟩ and p(y)Δy. Use the HistogramFrame class to determine and plot p(y)Δy. Choose at least 10^4 measurements of y for n = 4, 16, 32, and 64. What is the qualitative form of p(y)? Does the qualitative form of p(y) change as the number of measurements of y is increased for a given value of n? Does the qualitative form of p(y) change as n is increased?
c. Each value of $y$ can be considered to be a measurement. How much does the value of $y$ vary on the average from one measurement to another? Make a rough estimate of this variability by comparing several measurements of $y$ for a given value of $n$. Increase $n$ by a factor of four and estimate the variability of $y$ again. Does the variability from one measurement to another decrease (on the average) as $n$ is increased?

d. The sample variance $\hat{\sigma}^2$ is given by

$$\hat{\sigma}^2 = \frac{\sum_{i=1}^{n}|y_i - \langle y \rangle|^2}{n-1}. \quad (7.19)$$

The reason for the factor of $n-1$ rather than $n$ in (7.19) is that to compute $\hat{\sigma}^2$, we need to use the $n$ values of $x$ to compute the mean, $\langle y \rangle$, and thus, loosely speaking, we have only $n-1$ independent values of $x$ remaining to calculate $\hat{\sigma}^2$. Show that if $n >> 1$, then $\hat{\sigma}^2 \approx \sigma_y^2$, where $\sigma_y^2$ is given by

$$\sigma_y^2 = \langle y^2 \rangle - \langle y \rangle^2. \quad (7.20)$$

e. The quantity $\hat{\sigma}$ is known as the standard deviation of the mean. That is, $\hat{\sigma}$ gives a measure of how much variation we expect to find if we make repeated measurements of $y$. How does the value of $\hat{\sigma}$ compare with your estimate of the variability in part (b)?

f. What is the qualitative shape of the probability density $p(y)$ that you obtained in part (b)? What is the order of magnitude of the width of the probability?

g. Verify from your results that $\hat{\sigma} \approx \sigma_y \approx \sigma_x/\sqrt{n-1} \approx \sigma_x/\sqrt{n}$.

h. To test the generality of your results, consider the exponential probability density

$$f(x) = \begin{cases} e^{-x}, & x \geq 0 \\ 0, & x < 0. \end{cases} \quad (7.21)$$

Calculate $\langle x \rangle$ and $\sigma_x$ analytically. Modify your Monte Carlo program and estimate $\langle y \rangle$, $\hat{\sigma}$, $\sigma_y$, and $p(y)$. How are $\hat{\sigma}$, $\sigma_y$, and $\sigma_x$ related for a given value of $n$? Plot $p(y)$ and discuss its qualitative form and its dependence on $n$ and on the number of measurements of $y$.

Problem 7.15 illustrates the central limit theorem, which states that the probability distribution of a sum of random variables, the random variable $y$, is a Gaussian centered at $\langle y \rangle$ with a standard deviation approximately given by $1/\sqrt{n}$ times the standard deviation of $f(x)$. The requirements are that $f(x)$ has finite first and second moments, that the measurements of $y$ are statistically independent, and that $n$ is large. What is the relation of the central limit theorem to the calculations of the probability distribution in the random walk models that we have considered?

**Problem 7.16.** Generation of the Gaussian distribution

Consider the sum

$$y = \sum_{i=1}^{12} r_i, \quad (7.22)$$
where \( r_i \) is a uniform random number in the unit interval. Make many measurements of \( y \) and show that the probability distribution of \( y \) approximates the Gaussian distribution with mean value 6 and variance 1. What is the relation of this result to the central limit theorem? Discuss how to use this result to generate a Gaussian distribution with arbitrary mean and variance. This way of generating a Gaussian distribution is particularly useful when a “quick and dirty” approximation is appropriate. A better method for generating a sequence of random numbers distributed according to the Gaussian distribution is discussed in Section 12.5.

Many of the problems we have considered have revealed the slow convergence of Monte Carlo simulations and the difficulty of obtaining quantitative results for asymptotic quantities. We conclude this section with a cautionary note and consider a “simple” problem for which straightforward Monte Carlo methods give misleading asymptotic results.

*Problem 7.17. Random walk on lattices containing random traps*

a. In Problem 7.10 we considered the mean survival time of a one-dimensional random walker in the presence of a periodic distribution of traps. Now suppose that the trap sites are distributed at random on a one-dimensional lattice with density \( \rho = N/L \). For example, if \( \rho = 0.01 \), the probability that site is a trap site is 1%. (A site is a trap site if \( r \leq \rho \), where as usual \( r \) is uniformly distributed in the interval \( 0 \leq r < 1 \).) If a walker is placed at random at any nontrapping site, determine its mean survival time \( \tau \), that is, the mean number of steps before a trap site is reached. Assume that the walker has an equal probability of moving to a nearest neighbor site at each step and use periodic boundary conditions, that is, the lattice sites are located on a ring. The major complication is that it is necessary to perform three averages: the distribution of traps, the origin of the walker, and the different walks for a given trap distribution and origin. Choose reasonable values for the number of trials associated with each average and do a Monte Carlo simulation to estimate the mean survival time \( \tau \). If \( \tau \) exhibits a power law dependence on \( \rho \), for example, \( \tau \approx \tau_0 \rho^{-z} \), estimate the exponent \( z \).

b. A seemingly straightforward extension of part (a) is to estimate the survival probability \( S_N \) after \( N \) steps. Choose \( \rho = 0.5 \) and do a Monte Carlo simulation of \( S_N \) for \( N \) as large as possible. (Published results are for \( N = 3 \times 10^4 \) on lattices large enough that a walker doesn’t reach the boundary, and about 54,000 trials.) Assume that the asymptotic form of \( S_N \) for large \( N \) is given by

\[
S_N \sim e^{-bN^\alpha},
\]

where the exponent \( \alpha \) is the quantity of interest and \( b \) is a constant that depends on \( \rho \). Are your results consistent with this form? Is it possible to make a meaningful estimate of the exponent \( \alpha \)?

c. It has been proven that the asymptotic \( N \) dependence of \( S_N \) has the form (7.23) with \( \alpha = 1/3 \). Are your Monte Carlo results consistent with this value of \( \alpha \)? The object of part (b) is to convince you that it is not possible to use simple Monte Carlo methods directly to obtain the correct asymptotic behavior of \( S_N \). The difficulty is that we are trying to estimate \( S_N \) in the asymptotic region where \( S_N \) is very small, and the small number of trials in this region prevent us from obtaining meaningful results.
d. One way to reduce the number of required averages is to determine exactly the probability that the walker is at site $i$ after $N$ steps for a given distribution of trap sites. The method is illustrated in Figure 7.4. The first line represents a given configuration of traps distributed randomly on a one-dimensional lattice. One walker is placed at each non-trap site; trap sites are assigned the value 0. Because each walker moves with probability $1/2$ to each neighbor, the number of walkers $w_i(N+1)$ on site $i$ at step $N+1$ is given by

$$w_i(N+1) = \frac{1}{2}[w_{i+1}(N) + w_{i-1}(N)].$$  \hfill (7.24)

(Compare the relation (7.24) to the relation that you found in Problem 7.5d.) The survival probability $S_N$ after $N$ steps for a given configuration of traps is given exactly by

$$S_N = \frac{1}{w_0} \sum_i w_i(N),$$ \hfill (7.25)

where $w_0$ is the initial number of walkers and the sum is over all sites in the lattice. Explain the relation (7.25), and write a program that computes $S_N$ using (7.24) and (7.25). Then obtain $\langle S_N \rangle$ by averaging over several configurations of traps. Choose $\rho = 0.5$ and determine $S_N$ for $N = 32, 64, 128, 512, \text{ and } 1024$. Choose periodic boundary conditions and as large a lattice as possible. How well can you estimate the exponent $\alpha$? For comparison, Havlin et al. consider a lattice of $L = 50000$ and values of $N$ up to $10^7$.

One reason that random walks are very useful in simulating many physical processes is that they are closely related to solutions of the diffusion equation. The one-dimensional diffusion equation can be written as

$$\frac{\partial P(x, t)}{\partial t} = D \frac{\partial^2 P(x, t)}{\partial x^2},$$ \hfill (7.26)

where $D$ is the self-diffusion coefficient and $P(x, t) \Delta x$ is the probability of a particle being in the interval between $x$ and $x + \Delta x$ at time $t$. In a typical application $P(x, t)$ might represent the concentration of ink molecules diffusing in a fluid. In three dimensions the second derivative
\[ \frac{\partial^2}{\partial x^2} \] is replaced by the Laplacian \( \nabla^2 \). In Appendix 7B we show that the solution to the diffusion equation with the boundary condition \( P(x = \pm \infty, t) = 0 \) yields
\[ \langle x(t) \rangle = 0 \] (7.27)
and
\[ \langle x^2(t) \rangle = 2Dt. \quad \text{(one dimension)} \] (7.28)
If we compare the form of of (7.28) with (7.9), we see that the random walk on a one-dimensional lattice and the diffusion equation give the same time dependence if we identify \( t \) with \( N \Delta t \) and \( D \) with \( a^2/\Delta t \).

The relation of discrete random walks to the diffusion equation is an example of how we can approach many problems in several ways. The traditional way to treat diffusion is to formulate the problem as a partial differential equation as in (7.26). The usual method for solving (7.26) numerically is known as the Crank-Nicholson method (see Press et al.). One difficulty with this approach is the treatment of complicated boundary conditions. An alternative is to formulate the problem as a random walk on a lattice for which it is straightforward to incorporate various boundary conditions. We will consider random walks in many contexts (see for example, Section 11.5 and Chapter 17).

### 7.4 The Poisson Distribution and Nuclear Decay

As we have seen, we often can change variable names and consider a seemingly different physical problem. Our goal in this section is to discuss the decay of unstable nuclei, but we first discuss a conceptually easier problem related to throwing darts at random. Related physical problems are the distribution of stars in the sky and the distribution of photons on a photographic plate.

Suppose we randomly throw \( N = 100 \) darts at a board that has been divided into \( M = 1000 \) equal size regions. The probability that a dart lands in a given region or cell in any one throw is \( p = 1/M \). If we count the number of darts in the different regions, we would find that most cells are empty, some cells have one dart, and other cells have more than one dart. What is the probability \( P(n) \) that a given cell has \( n \) darts?

**Problem 7.18.** Throwing darts

Write a program that simulates the throwing of \( N \) darts at random into \( M \) cells in a dart board. Throwing a dart at random at the board is equivalent to choosing an integer at random between 1 and \( M \). Determine \( H(n) \), the number of cells with \( n \) darts. Average \( H(n) \) over many trials, and then compute the probability distribution
\[ P(n) = \frac{H(n)}{M}. \] (7.29)
As an example, choose \( N = 50 \) and \( M = 500 \). Choose the number of trials to be sufficiently large so that you can determine the qualitative form of \( P(n) \). What is \( \langle n \rangle \)?

In this case, the probability \( p \) that a dart lands in a given cell, is much less then unity. The conditions \( N >> 1 \) and \( p << 1 \) with \( \langle n \rangle = Np \) fixed and the independence of the events (the
presence of a dart in a particular cell) satisfy the requirements for a Poisson distribution. The form of the Poisson distribution is

\[ P(n) = \frac{\langle n \rangle^n}{n!} e^{-\langle n \rangle}, \tag{7.30} \]

where \( n \) is the number of darts in a given cell and \( \langle n \rangle \) is the mean number, \( \langle n \rangle = \sum_{n=0}^{N} nP(n) \). Because \( N >> 1 \), we can take the upper limit of this sum to be \( \infty \) when it is convenient.

**Problem 7.19.** Darts and the Poisson distribution

a. Write a program to compute \( \sum_{n=0}^{\infty} P(n) \), \( \sum_{n=0}^{\infty} nP(n) \), and \( \sum_{n=0}^{\infty} n^2 P(n) \) using the form (7.30) for \( P(n) \) and reasonable values of \( p \) and \( N \). Verify that \( P(n) \) in (7.30) is normalized. What is the value of \( \sigma_n^2 = \langle n^2 \rangle - \langle n \rangle^2 \) for the Poisson distribution?

b. Modify the program that you developed in Problem 7.18 to compute \( \langle n \rangle \) as well as \( P(n) \). Choose \( N = 50 \) and \( M = 1000 \). How does your computed values of \( P(n) \) compare to the Poisson distribution in (7.30) using your measured value of \( \langle n \rangle \) as input? If time permits, use larger values of \( N \) and \( M \).

c. Choose \( N = 50 \) and \( M = 100 \) and redo part (b). Are your results consistent with a Poisson distribution? What happens if \( M = N = 50 \)?

Now that we are more familiar with the Poisson distribution, we consider the decay of radioactive nuclei. We know that a collection of radioactive nuclei will decay, and we cannot know a priori which nucleus will decay next. If all nuclei of a particular type are identical, why do they not all decay at the same time? The answer is based on the uncertainty inherent in the quantum description of matter at the microscopic level. In the following, we will see that a simple model of the decay process leads to exponential decay. This approach complements the continuum approach discussed in Section 3.9.

Because each nucleus is identical, we assume that during any time interval \( \Delta t \), each nucleus has the same probability per unit time \( p \) of decaying. The basic algorithm is simple – choose an unstable nucleus and generate a random number \( r \) uniformly distributed in the unit interval \( 0 \leq r < 1 \). If \( r \leq p \), the unstable nucleus decays; otherwise, it does not. Each unstable nucleus is tested once during each time interval. Note that for a system of unstable nuclei, there are many events that can happen during each time interval, for example, 0, 1, 2, . . . , \( n \) nuclei can decay. Once a nucleus decays, it is no longer in the group of unstable nuclei that is tested at each time interval. Class Nuclei in Listing 7.5 implements the nuclear decay algorithm.

**Listing 7.5: The Nuclei class.**

```java
package org.opensourcephysics.sip.ch07;

public class Nuclei {
    int n[]; // accumulated data on number of unstable nuclei, index is time
    int tmax; // maximum time to record data
    int n0; // initial number of unstable nuclei
    double p; // decay probability

    public void initialize () {
        n = new int[tmax+1];
```
Problem 7.20. Monte Carlo simulation of nuclear decay

a. Write a target class that extends AbstractSimulation, does many trials, and plots the average number of unstable nuclei as a function of time. Assume that the time interval $\Delta t$ is one second. Choose the initial number of unstable nuclei $n_0 = 10000$, $p = 0.01$, and $t_{\text{max}} = 100$, and average over 100 trials. Is your result for $n(t)$, the mean number of unstable nuclei at time $t$, consistent with the expected behavior, $n(t) = n(0) e^{-\lambda t}$ found in Section 3.9? What is the value of $\lambda$ for this value of $p$?

b. There are a very large number of unstable nuclei in a typical radioactive source. We also know that over any reasonable time interval, only a relatively small number decay. Because $N >> 1$ and $p << 1$, we expect that $P(n)$, the probability that $n$ nuclei decay during a specified time interval, is a Poisson distribution. Modify your target class so that it outputs the probability that $n$ unstable nuclei decay during the first time interval. Choose $n_0 = 1000$, $p = 0.001$, $t_{\text{max}} = 1$, and average over 1000 trials. What is the mean number $\langle n \rangle$ of nuclei that decay during this interval? What is the associated variance? Plot $P(n)$ versus $n$ and compare your results to the Poisson distribution (7.30) with your measured value of $\langle n \rangle$ as input. Then consider $p = 0.02$ and determine if $P(n)$ is a Poisson distribution.

c. Modify your target class so that it outputs the probability that $n$ unstable nuclei decay during the first two time intervals. Choose $n_0 = 10000$, $p = 0.001$, and $t_{\text{max}} = 2$. Average over 1000 trials. Compare the probability you obtain with your results from part (b). How do your results change as the time interval becomes larger?

d. Increase $p$ for fixed $n_0 = 10000$ and determine $P(n)$ for a fixed time interval. Estimate the values of $p$ and $n$ for which the Poisson distribution is no longer applicable.

e. Modify your program so that it flashes a small circle on the screen or makes a sound (like that of a Geiger counter) when a nucleus decays. You can have the computer make a beep by using the method Toolkit.getDefaultToolkit().beep() in java.awt. Choose the location of the small circle at random. Do a single run and describe the qualitative differences between the visual or audio patterns for the cases in parts (a)–(d)? Choose $n_0 \geq 5000$. Such a visualization
might be somewhat misleading on a serial computer because only one nuclei can be considered at a time. In contrast, for a real system, the nuclei can decay simultaneously.

7.5 Problems in Probability

Why have we bothered to simulate many random processes that can be solved by analytical methods? The main reason is that it is simpler to introduce new methods in a familiar context. Another reason is that if we change the nature of many random processes slightly, it often is the case that it would be difficult or impossible to obtain the answers by familiar methods. Still another reason is that writing a program and doing a simulation can aid your intuitive understanding of the system, especially if the questions involve the subtle concept of probability. Probability is an elusive concept in part because it cannot be measured at one time. To reinforce the importance of thinking about how to solve a problem on a computer, we suggest some problems in probability in the following. Does thinking about how to write a program to simulate these problems help you to find a pencil and paper solution?

Problem 7.21. Three boxes: stick or switch?
Suppose that there are three identical boxes, each with a lid. When you leave the room, a friend places a $10 bill in one of the boxes and closes the lid of each box. The friend knows the location of the $10 bill, but you do not. You then reenter the room and guess which one of the boxes has the $10 bill. As soon as you do, your friend opens the lid of a box that is empty. If you have chosen an empty box, your friend will open the lid of the other empty box. If you have chosen the right box, your friend will open the lid of one of the two empty boxes. You now have the opportunity to stay with your original choice or switch to the other unopened box. Suppose that you play this contest many times and that each time you guess correctly, you keep the money. To maximize your winnings, should you maintain your initial choice or should you switch? Which strategy is better? This contest is known as the Monty Hall problem. Write a program to simulate this game and output the probability of winning for switching and for not switching. It is likely that before you finish your program, the correct strategy will become clear. To make your program more useful, consider an arbitrary number of boxes.

Problem 7.22. Conditional probability
Suppose that many people in a community are tested at random for HIV. The accuracy of the test is 87% and the incidence of the disease in the general population, independent of any test, is 1%. If a person tests positive for HIV, what is the probability that this person really has HIV? Write a program to compute the probability. The answer can be found by using Bayes’ theorem (cf. Bernardo and Smith). The answer is much less than 87%.

Problem 7.23. The roll of the dice
Suppose that two gamblers each begin with $100 in capital and on each throw of a coin, one gambler must win $1 and the other must lose $1. How long can they play on the average until the capital of the loser is exhausted? How long can they play if they each begin with $1000? Neither gambler is allowed to go into debt. The eventual outcome is known as the gambler’s ruin.


<table>
<thead>
<tr>
<th>Team</th>
<th>Won</th>
<th>Lost</th>
<th>Percentage</th>
</tr>
</thead>
<tbody>
<tr>
<td>St. Louis Cardinals</td>
<td>105</td>
<td>57</td>
<td>0.648</td>
</tr>
<tr>
<td>Houston Astros</td>
<td>92</td>
<td>70</td>
<td>0.568</td>
</tr>
<tr>
<td>Chicago Cubs</td>
<td>89</td>
<td>73</td>
<td>0.469</td>
</tr>
<tr>
<td>Pittsburgh Pirates</td>
<td>72</td>
<td>89</td>
<td>0.447</td>
</tr>
<tr>
<td>Cincinnati Reds</td>
<td>76</td>
<td>86</td>
<td>0.426</td>
</tr>
<tr>
<td>Milwaukee Brewers</td>
<td>67</td>
<td>94</td>
<td>0.416</td>
</tr>
</tbody>
</table>


Problem 7.24. The boys of summer

Luck plays a large role in the outcome of any baseball season. The National League Central Division standings for 2004 are given in Table 7.1. Suppose that the teams remain unchanged and their probability of winning a particular game is given by their 2004 winning percentage. Do a simulation to determine the probability that the Cardinals would lead the division for another season. For simplicity, assume that the teams play only each other.

Much of the present day motivation for the development of probability comes from science rather than from gambling. The next problem has much to do with statistical physics even though this application is not apparent.

Problem 7.25. Money exchange

Consider a line that has been subdivided into bins. There can be an indefinite number of coins in each bin. For simplicity, we initially assign one coin to each bin. The money exchange proceeds as follows. Select two bins at random. If there is at least one coin in the first bin, move one coin to the second bin. If the first bin is empty, then do nothing. After many coin exchanges, how is the occupancy of the bins distributed? Are the coins uniformly distributed as in the initial state or are many bins empty? Write a Monte Carlo program to simulate this money exchange and show the state of the bins visually. Consider a system with at least 256 bins. Plot the histogram \( H(n) \) versus \( n \), where \( H(n) \) is the number of bins with \( n \) coins. Do your results change qualitatively if you consider bigger systems or begin with more coins in each bin?

Problem 7.26. Distribution of cooking times

An industrious physics major finds a job at a local fast food restaurant to help her pay her way through college. Her task is to cook 20 hamburgers on a grill at any one time. When a hamburger is cooked, she is supposed to replace it with an uncooked hamburger. However, our physics major does not pay attention to whether the hamburger is cooked or not. Her method is to choose a hamburger at random and replace it by a uncooked one. She does not check if the hamburger that she removes from the grill is ready. What is the distribution of cooking times of the hamburgers that she removes? For simplicity, assume that she replaces a hamburger at regular intervals of thirty seconds and that there is an indefinite supply of uncooked hamburgers. Does the qualitative nature of the distribution change if she cooks 40 hamburgers at any one time?
7.6 Method of Least Squares

In Problem 7.20 we did a simulation of \( N(t) \), the number of unstable nuclei at time \( t \). Given the finite accuracy of our data, how do we know if our simulation results are consistent with the exponential relation between \( N \) and \( t \)? The approach that we have been using is to plot the computed values of \( \log N(t) \) as a function of \( t \) and to rely on our eye to help us draw the curve that best fits the data points. Such a visual approach works best when the curve is a straight line, that is, when the relation is linear. The advantages of this approach are that it is straightforward and allows us to see what we are doing. For example, if a data point is far from the straight line, or if there is a gap in the data, we will notice it easily. If the analytical relation is not linear, it is likely that we will notice that the data points do not fit a simple straight line, but instead show curvature. If we blindly let a computer fit the data to a straight line, we might not notice that the fit is not very good unless we already have had experience fitting data. Finally, the visceral experience of fitting the data manually gives us some feeling for the nature of the data that might otherwise be missed. It is a good idea to plot some data in this way even though a computer can do it much faster.

Although the visual approach is simple, it does not yield precise results, and we also need to use more systematic fitting methods. The most common method for finding the best straight line fit to a series of measured points is called linear regression or least squares. Suppose we have \( n \) pairs of measurements \( (x_1, y_1), (x_2, y_2), \ldots, (x_n, y_n) \) and that the errors are entirely in the values of \( y \). For simplicity, we assume that the uncertainties in \( \{y_i\} \) all have the same magnitude. Our goal is to obtain the best fit to the linear function

\[
y = mx + b. \tag{7.31}
\]

The problem is to calculate the values of the parameters \( m \) and \( b \) for the best straight line through the \( n \) data points. The difference,

\[
d_i = y_i - mx_i - b, \tag{7.32}
\]

is a measure of the discrepancy in \( y_i \). It is reasonable to assume that the best pair of values of \( m \) and \( b \) are those that minimize the quantity

\[
\chi^2 = \sum_{i=1}^{n} (y_i - mx_i - b)^2. \tag{7.33}
\]

Why should we minimize the sum of the squared differences between the experimental values, \( y_i \), and the analytical values, \( mx_i + b \), and not some other function of the differences? The justification is based on the assumption that if we did many simulations or measurements, then the values of \( d_i \) would be distributed according to the Gaussian distribution (see Problems 7.5 and 7.15). Based on this assumption, it can be shown that the values of \( m \) and \( b \) that minimize \( \chi \) yield a set of values of \( mx_i + b \) that are the most probable set of measurements that we would find based on the available information. This link to probability is the reason we have discussed least squares fits in this chapter, even though we will not explicitly show that the difference \( d_i \) is distributed according to a Gaussian distribution.
CHAPTER 7. RANDOM PROCESSES

To minimize $\chi$, we take the partial derivative of $S$ with respect to $b$ and $m$:

$$\frac{\partial \chi}{\partial m} = -2 \sum_{i=1}^{n} x_i (y_i - mx_i - b) = 0,$$

(7.34a)

$$\frac{\partial \chi}{\partial b} = -2 \sum_{i=1}^{n} (y_i - mx_i - b) = 0.$$

(7.34b)

From (7.34) we obtain two simultaneous equations:

$$m \sum_{i=1}^{n} x_i^2 + b \sum_{i=1}^{n} x_i = \sum_{i=1}^{n} x_i y_i$$

(7.35a)

$$m \sum_{i=1}^{n} x_i + bn = \sum_{i=1}^{n} y_i.$$

(7.35b)

It is convenient to define the quantities

$$\langle c \rangle = \frac{1}{n} \sum_{i=1}^{n} x_i$$

(7.36a)

$$\langle y \rangle = \frac{1}{n} \sum_{i=1}^{n} y_i$$

(7.36b)

$$\langle xy \rangle = \frac{1}{n} \sum_{i=1}^{n} x_i y_i,$$

(7.36c)

and rewrite (7.35) as

$$m \langle x^2 \rangle + b \langle x \rangle = \langle xy \rangle,$$

(7.37a)

$$m \langle x \rangle + b = \langle y \rangle.$$

(7.37b)

The solution of (7.37) can be expressed as

$$m = \frac{\langle xy \rangle - \langle x \rangle \langle y \rangle}{\sigma_x^2}$$

(7.38a)

$$b = \langle y \rangle - m \langle x \rangle,$$

(7.38b)

where

$$\sigma_x^2 = \langle x^2 \rangle - \langle x \rangle^2.$$  

(7.38c)

Equation (7.38) determines the slope $m$ and the intercept $b$ of the best straight line through the $n$ data points. (Note that the average in the equations for the coefficients $m$ and $b$ are over the data points.)

As an example, consider the data shown in Table 7.2 for a one-dimensional random walk. To make the example more interesting, suppose that the walker takes steps of length 1 or 2 with equal
Table 7.2: Computed values of the mean square displacement $\Delta x^2$ as a function of the total number of steps $N$. The mean square displacement was averaged over 1000 trials. The one-dimensional random walker takes steps of length 1 or 2 with equal probability, and the direction of the step is random with $p = 1/2$.

<table>
<thead>
<tr>
<th>$N$</th>
<th>$\Delta x^2$</th>
</tr>
</thead>
<tbody>
<tr>
<td>8</td>
<td>19.43</td>
</tr>
<tr>
<td>16</td>
<td>37.65</td>
</tr>
<tr>
<td>32</td>
<td>76.98</td>
</tr>
<tr>
<td>64</td>
<td>160.38</td>
</tr>
</tbody>
</table>

probability. The direction of the step is random and $p = 1/2$. As in Section 7.2, we assume that the mean square displacement $\Delta x^2$ obeys the general relation

$$\Delta x^2 = aN^{2\nu},$$

(7.39)

with an unknown exponent $\nu$ and $a$ is a constant. Note that the fitting problem in (7.39) is nonlinear, that is, $\langle x^2 \rangle - \langle x \rangle^2$ depends on $N^\nu$ rather than $N$. Often a problem that looks nonlinear can be turned into a linear problem by a change of variables. In this case we convert the nonlinear relation (7.39) to a linear relation by taking the logarithm of both sides:

$$\ln(\Delta x^2) = \ln a + 2\nu \ln N.$$  

(7.40)

The values of $y = \ln(\Delta x^2)$ and $x = \ln N$ in Table 7.2 and the least squares fit are shown in Figure 7.5. We use (7.38) and find that $m = 1.02$ and $b = 0.83$. Hence, we conclude from our limited data and the relation $2\nu = m$ that $\nu \approx 0.51$, which is consistent with the expected result $\nu = 1/2$.

The least squares fitting procedure also allows us to estimate the uncertainty or the most probable error in $m$ and $b$ by analyzing the measurements themselves. The result of this analysis is that the most probable error in $m$ and $b$, $\sigma_m$ and $\sigma_b$ respectively, is given by

$$\sigma_m = \frac{1}{\sqrt{n}} \Delta \frac{1}{\sigma_x},$$

(7.41a)

$$\sigma_b = \frac{1}{\sqrt{n}} \left( \frac{\langle x^2 \rangle}{\sigma_x} \right)^{1/2} \Delta,$$

(7.41b)

where

$$\Delta^2 = \frac{1}{n-2} \sum_{i=1}^{n} d_i^2,$$

(7.42)

and $d_i$ is given by (7.32).

Because there are $n$ data points, we might have guessed that $n$ rather than $n - 2$ would be present in the denominator of (7.42). The reason for the factor of $n - 2$ is related to the fact that to determine $\Delta$, we first need to calculate two quantities $m$ and $b$, leaving only $n - 2$ independent degrees of freedom. To see that the $n - 2$ factor is reasonable, consider the special case of $n = 2$. 
Figure 7.5: Plot of $\Delta x^2$ versus $\ln N$ for the data listed in Table 7.2. The straight line $y = 1.02x + 0.83$ through the points is found by minimizing the sum (7.33).

In this case we can find a line that passes exactly through the two data points, but we cannot deduce anything about the reliability of the set of measurements because the fit is exact. If we use (7.42), we see that both the numerator and denominator would be zero, and hence $\Delta$ would be undetermined. If a factor of $n$ rather than $n - 2$ appeared in (7.42), we would conclude that $\Delta = 0/2 = 0$, an absurd conclusion. Usually $n >> 1$, and the difference between $n$ and $n - 2$ is negligible.

For our example, $\Delta = 0.03$, $\sigma_b = 0.07$, and $\sigma_m = 0.02$. The uncertainties $\delta m$ and $\delta \nu$ are related by $2\delta \nu = \delta m$. Because $\delta m = \sigma_m$, we conclude that our best estimate for $\nu$ is $\nu = 0.51 \pm 0.01$.

If the values of $y_i$ have different uncertainties $\sigma_i$, then the data points are weighted by the quantity $w_i = 1/\sigma_i^2$. In this case it is reasonable to minimize the quantity

$$
\chi^2 = \sum_{i=1}^{n} w_i(y_i - mx_i - b)^2. \tag{7.43}
$$

The resulting expressions in (7.38) for $m$ and $b$ are unchanged if we generalize the definition of the averages to be

$$
\langle f \rangle = \frac{1}{\langle w \rangle} \sum_{i=1}^{n} w_i f_i, \tag{7.44}
$$

where

$$
\langle w \rangle = \frac{1}{n} \sum_{i=1}^{n} w_i. \tag{7.45}
$$
Problem 7.27. Example of least squares fit

a. Write a program to find the least squares fit for a set of data. As a check on your program, compute the most probable values of \( m \) and \( b \) for the data shown in Table 7.2.

b. Modify the random walk program so that steps of length 1 and 2 are taken with equal probability. Use at least 10000 trials and do a least squares fit to \( \Delta x^2 \) as done in the text. Is your most probable estimate for \( \nu \) closer to \( \nu = 1/2 \)?

For simple random walk problems the relation \( \Delta x^2 = aN\nu \) holds for all \( N \). However, in many random walk problems a power law relation between \( \Delta x^2 \) and \( N \) holds only asymptotically for large \( N \), and hence we should use only the larger values of \( N \) to estimate the slope. Also, because we are finding the best fit for the logarithm of the independent variable \( N \), we need to give equal weight to all intervals of \( \ln N \). In the above example, we used \( N = 8, 16, 32, \) and 64, so that the values of \( \ln N \) are equally spaced.

7.7 Applications to Polymers

Random walk models play an important role in polymer physics (cf. de Gennes). A polymer consists of \( N \) repeat units (monomers) with \( N \gg 1 \) \((N \sim 10^3 - 10^5)\). For example, polyethylene can be represented as \( \cdots -\text{CH}_2-\text{CH}_2-\text{CH}_2-\cdots \). The detailed structure of the polymer is important for many practical applications. For example, if we wish to improve the fabrication of rubber, a good understanding of the local motions of the monomers in the rubber chain is essential. However, if we are interested in the global properties of the polymer, the details of the chain structure can be ignored.

Let us consider a familiar example of a polymer chain in a good solvent: a noodle in warm water. A short time after we place a noodle in warm water, the noodle becomes flexible, and it neither collapses into a little ball or becomes fully stretched. Instead, it adopts a random structure as shown schematically in Figure 7.6. If we do not add too many noodles, we can say that the noodles behave as a dilute solution of polymer chains in a good solvent. The dilute nature of the solution implies that we can ignore entanglement effects of the noodles and consider each noodle individually. The presence of a good solvent implies that the polymers can move freely and adopt many different configurations.

A fundamental geometrical property that characterizes a polymer in a good solvent is the mean square end-to-end distance \( \langle R^2 \rangle \), where \( N \) is the number of monomers. (For simplicity, we will frequently write \( R^2 \) in the following.) For a dilute solution of polymer chains in a good solvent, it is known that the asymptotic dependence of \( R^2 \) is given by (7.12) with \( \nu \approx 0.5874 \) in three dimensions. If we were to ignore the interactions of the monomers, the simple random walk model would yield \( \nu = 1/2 \), independent of the dimension and symmetry of the lattice. Because this result for \( \nu \) does not agree with experiment, we know that we are overlooking an important physical feature of polymers.

We now discuss a random walk that incorporates the global features of dilute linear polymers in solution. We already have introduced a model of a polymer chain consisting of straight line segments of the same size joined together at random angles (see Problem 7.13). A further idealization
Figure 7.6: (a) Schematic illustration of a linear polymer in a good solvent. (b) Example of the corresponding self-avoiding walk on a square lattice.

is to place the polymer chain on a lattice (see Figure 7.6). A more realistic model of linear polymers accounts for its most important physical feature, that is, two monomers cannot occupy the same spatial position. This constraint is known as the excluded volume condition, which is ignored in a simple random walk. A well known lattice model for a linear polymer chain that incorporates this constraint is known as the self-avoiding walk (SAW). This model consists of the set of all \( N \) step walks starting from the origin subject to the global constraint that no lattice site can be visited more than once in each walk; this constraint accounts for the excluded volume condition.

Self-avoiding walks have many applications, such as the physics of magnetic materials and the study of phase transitions, and they are of interest as purely mathematical objects. Many of the obvious questions have resisted rigorous analysis, and exact enumeration and Monte Carlo simulation have played an important role in our current understanding. The result for \( \nu \) in two dimensions for the self-avoiding walk is known to be exactly \( \nu = 3/4 \). The proportionality constant in (7.12) depends on the structure of the monomers and on the solvent. In contrast, the exponent \( \nu \) is independent of these details and depends only on the spatial dimension.

We consider Monte Carlo simulations of the self-avoiding walk in two dimensions in Problems 7.28–7.30. Another algorithm for the self-avoiding walk is considered in Project 7.41.

**Problem 7.28.** The two-dimensional self-avoiding walk
Consider the self-avoiding walk on the square lattice. Choose an arbitrary site as the origin and assume that the first step is “up.” The walks generated by the three other possible initial directions only differ by a rotation of the whole lattice and do not have to be considered explicitly. The second step can be in three rather than four possible directions because of the constraint that the walk cannot return to the origin. To obtain unbiased results, we generate a random number to choose
one of the three directions. Successive steps are generated in the same way. Unfortunately, the walk will very likely not continue indefinitely. To obtain unbiased results, we must choose at random one of the three steps, even though one or more of these steps might lead to a self-intersection. If the next step does lead to a self-intersection, the walk must be terminated to keep the statistics unbiased. An example of a three step walk is shown in Figure 7.7a. The next step leads to a self-intersection and violates the constraint. In this case we must start a new walk at the origin.

a. Write a program that implements this algorithm and record the fraction \( f(N) \) of successful attempts at constructing polymer chains with \( N \) total monomers. Represent the lattice as an array so that you can record the sites that already have been visited. What is the qualitative dependence of \( f(N) \) on \( N \)? What is the maximum value of \( N \) that you can reasonably consider?

b. Determine the mean square end-to-end distance \( \langle R_N^2 \rangle \) for values of \( N \) that you can reasonably consider with this sampling method.

The disadvantage of the straightforward sampling method in Problem 7.28 is that it becomes very inefficient for long chains, that is, the fraction of successful attempts decreases exponentially. To overcome this attrition, several “enrichment” techniques have been developed. We first discuss a relatively simple algorithm proposed by Rosenbluth and Rosenbluth in which each walk of \( N \) steps is associated with a weighting function \( w(N) \). Because the first step to the north is always possible, we have \( w(1) = 1 \). In order that all allowed configurations of a given \( N \) are counted equally, the weights \( w(N) \) for \( N > 1 \) are determined according to the following possibilities:

1. All three possible steps violate the self-intersection constraint (see Figure 7.7b). The walk is terminated with a weight \( w(N) = 0 \), and a new walk is generated at the origin.
2. All three steps are possible and \( w(N) = w(N - 1) \).
3. Only \( m \) steps are possible with \( 1 \leq m < 3 \) (see Figure 7.7c). In this case \( w(N) = (m/3)w(N - 1) \), and one of the \( m \) possible steps is chosen at random.
The desired unbiased value of $\langle R^2 \rangle$ is obtained by weighting $R^2_i$, the value of $R^2$ obtained in the $i$th trial, by the value of $w_i(N)$, the weight found for this trial. Hence we write

$$\langle R^2 \rangle = \sum_i w_i(N) R^2_i \sum_i w_i(N),$$

(7.46)

where the sum is over all trials.

**Problem 7.29.** Rosenbluth and Rosenbluth enrichment method

Incorporate the Rosenbluth method into your Monte Carlo program and compute $R^2$ for $N = 4, 8, 16, \text{ and } 32$. Estimate the exponent $\nu$ from a log-log plot of $R^2$ versus $N$. Can you distinguish your estimate for $\nu$ from its random walk value $\nu = 1/2$?

The Rosenbluth and Rosenbluth procedure is not particularly efficient because many walks still terminate, and thus we do not obtain many walkers for large $N$. Grassberger improved this algorithm by increasing the population of walkers with high weights and reducing the population of walkers with low weights. The idea is that if $w(N)$ for a given trial is above a certain threshold, we add a new walker and give the new and old walker half of the original weight. If $w(N)$ is below a certain threshold, then we eliminate half of the walkers with weights below this threshold (for example, every second walker) and double the weights of the remaining half. It is a good idea to adjust the thresholds as the simulation runs in order to maintain a relatively constant number of walkers.

More recently Prellberg and Krawczyk further improved the Rosenbluth and Rosenbluth enrichment method so that there is no need to provide a threshold value. After each step the average weight of the walkers, $\langle w(N) \rangle$ is computed for a given trial and the ratio $r = w(N)/\langle w(N) \rangle$ is used to determine whether to add walkers (enrichment) or eliminate walkers (pruning). If $r > 1$, then $c = \min(r, m)$ copies of the walker are made each with weight $w(N)/c$. If $r < 1$, then remove this walker with probability $1 - r$. This algorithm leads to an approximately constant number of walkers and is related to the Wang-landau method which we will discuss in Problem 16.30.

Another enrichment algorithms is the “reptation” method (see Wall and Mandel). For simplicity, consider a model polymer chain in which all bond angles are $\pm 90^\circ$. As an example of this model, the five independent $N = 5$ polymer chains are shown in Figure 7.8. (Other chains differ only by a rotation or a reflection.) The reptation method can be stated as follows:

1. Choose a chain at random and remove the tail link.
2. Attempt to add a link to the head of the chain. There is a maximum of two directions in which the new head link can be added.
3. If the attempt violates the self-intersection constraint, return to the original chain and inter-change the head and tail. Include the chain in the statistical sample.

The above steps are repeated many times to obtain an estimate of $R^2$.

As an example of the reptation method, consider chain $a$ of Figure 7.8. A new link can be added in two directions (see Figure 7.9a), so that on the average we find, $a \rightarrow \frac{1}{2}c + \frac{1}{2}d$. In contrast, a link can be added to chain $b$ in only one direction, and we obtain $b \rightarrow \frac{1}{2}c + \frac{1}{2}b$, where the tail and
head of chain b have been interchanged (see Figure 7.9b). Confirm that $c \rightarrow \frac{1}{2}e + \frac{1}{2}a$, $d \rightarrow \frac{1}{2}c + \frac{1}{2}d$, and $e \rightarrow \frac{1}{2}a + \frac{1}{2}b$, and that all five chains are equally probable. That is, the transformations in the reptation method preserve the proper statistical weights of the chains without attrition. There is just one problem: unless we begin with a double ended “cul-de-sac” configuration such as shown in Figure 7.10, we will never obtain such a configuration using the above transformation. Hence, the reptation method introduces a small statistical bias, and the calculated mean end-to-end distance will be slightly larger than if all configurations were considered. However, the probability of such trapped configurations is very small, and the bias can be neglected for most purposes.

*Problem 7.30. The reptation method

a. Adopt the $\pm 90^\circ$ bond angle restriction and calculate by hand the exact value of $\langle R^2 \rangle$ for $N = 5$. Then write a Monte Carlo program that implements the reptation method. Generate one walk of $N = 5$ and use the reptation method to generate a statistical sample of chains. As a check on your program, compute $\langle R^2 \rangle$ for $N = 5$ and compare your result with the exact result. Then extend your Monte Carlo computations of $\langle R^2 \rangle$ to larger $N$.

b. Modify the reptation model so that the bond angle also can be $180^\circ$. This modification leads to a maximum of three directions for a new bond. Compare your results with those from part (a).

In principle, the dynamics of a polymer chain undergoing collisions with solvent molecules can be simulated by using a molecular dynamics method. However, in practice only relatively small chains can be simulated in this way. An alternative approach is to use a Monte Carlo model that simplifies the effect of the random collisions of the solvent molecules with the atoms of the chain. Most of these models (cf. Verdier and Stockmayer) consider the chain to be composed of beads.
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Figure 7.9: The possible transformations of chains a and b. One of the two possible transformations of chain b violates the self-intersection restriction and the head and tail are interchanged.

Figure 7.10: Example of a double cul-de-sac configuration for the self-avoiding walk that cannot be obtained by the reptation method.

connected by bonds and restrict the positions of the beads to the sites of a lattice. For simplicity, we assume that the bond angles can be either $\pm 90^\circ$ or $180^\circ$. The idea is to begin with an allowed configuration of $N$ beads ($N - 1$ bonds). A possible starting configuration can be generated by taking successive steps in the positive $y$ direction and positive $x$ directions. The dynamics of the Verdier-Stockmayer algorithm is summarized by the following steps.

1. Select at random a bead (occupied site) on the polymer chain. If the bead is not an end site, then the bead can move to a nearest neighbor site of another bead if this site is empty and if the new angle between adjacent bonds is either $\pm 90^\circ$ or $180^\circ$. For example, bead 4 in Figure 7.11 can move to position $4'$ while bead 3 cannot move if selected. That is, a selected bead can move to a diagonally opposite unoccupied site only if the two bonds to which it is attached are mutually perpendicular.

2. If the selected bead is an end site, move it to one of two (maximum) possible unoccupied sites so that the bond to which it is connected changes its orientation by $\pm 90^\circ$ (see Figure 7.11).
3. If the selected bead cannot move, retain the previous configuration.

The physical quantities of interest include \( \langle R^2 \rangle \) and the mean square displacement of the center of mass of the chain \( \langle r^2 \rangle = \langle x^2 \rangle - \langle x \rangle^2 + \langle y^2 \rangle - \langle y \rangle^2 \), where \( x \) and \( y \) are the coordinates of the center of mass. The unit of time is the number of Monte Carlo steps per bead; in one Monte Carlo step per bead each bead has one chance on the average to move to a different site.

Another efficient method for simulating the dynamics of a polymer chain is the bond fluctuation model (see Carmesin and Kremer).

**Problem 7.31.** The dynamics of polymers in a dilute solution

a. Consider a two-dimensional lattice and compute \( \langle R^2 \rangle \) and \( \langle r^2 \rangle \) for various values of \( N \). How do these quantities depend on \( N \)? (The first published results for three dimensions were limited to 32 Monte Carlo steps per bead for \( N = 8, 16, \) and 32 and only 8 Monte Carlo steps per bead for \( N = 64 \).) Also compute the probability density \( P(R) \) that the end-to-end distance is \( R \). How does this probability compare to a Gaussian distribution?

b. Two configurations are strongly correlated if they differ by only the position of one bead. Hence, it would be a waste of computer time to measure the end-to-end distance and the position of the center of mass after every single move. Ideally, we wish to compute these quantities for configurations that are approximately statistically independent. Because we do not know \textit{a priori} the mean number of Monte Carlo steps per bead needed to obtain configurations that are statistically independent, it is a good idea to estimate this time in our preliminary calculations. The correlation time, \( \tau \), is the time needed to obtain statistically independent configurations and can be obtained by computing the equilibrium averaged time-autocorrelation function for a chain of fixed \( N \):

\[
C(t) = \frac{\langle R^2(t' + t)R^2(t') \rangle - \langle R^2 \rangle^2}{\langle R^4 \rangle - \langle R^2 \rangle^2}.
\] (7.47)

\( C(t) \) is defined so that \( C(t = 0) = 1 \) and \( C(t) = 0 \) if the configurations are not correlated. Because the configurations will become uncorrelated if the time \( t \) between the configurations
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is sufficiently long, we expect that \( C(t) \to 0 \) for \( t >> 1 \). We expect that \( C(t) \sim e^{-t/\tau} \), that is, \( C(t) \) decays exponentially with a decay or correlation time \( \tau \). Estimate \( \tau \) from a plot of \( \ln C(t) \) versus \( t \). Another way of estimating \( \tau \) is from the integral \( \int_0^\infty dt \, C(t) \), where \( C(t) \) is normalized so that \( C(0) = 1 \). (Because we determine \( C(t) \) at discrete values of \( t \), this integral is actually a sum.) How do your two estimates of \( \tau \) compare? A more detailed discussion of the estimation of correlation times can be found in Section 16.7.

Another type of random walk that is less constrained than the self-avoiding random walk is the “true” self-avoiding walk. This walk describes the path of a random walker that avoids visiting a lattice site with a probability that is a function of the number of times the site has been visited already. This constraint leads to a reduced excluded volume interaction in comparison to the usual self-avoiding walk.

Problem 7.32. The true self-avoiding walk in one dimension

In one dimension the true self-avoiding walk corresponds to a walker that can jump to one of its two nearest neighbors with a probability that depends on the number of times these neighbors already have been visited. Suppose that the walker is at site \( i \) at step \( t \). The probability that the walker will jump to site \( i + 1 \) at time \( t + 1 \) is given by

\[
p_{i+1} = \frac{e^{-gn_{i+1}}}{e^{-gn_{i+1}} + e^{-gn_{i-1}}},
\]

(7.48)

where \( n_{i\pm1} \) is the number of times that the walker has already visited site \( i \pm 1 \). The probability of a jump to site \( i - 1 \) is \( p_{i-1} = 1 - p_{i+1} \). The parameter \( g \) \((g > 0)\) is a measure of the “desire” of the path to avoid itself. The first few steps of a typical true self-avoiding walk walk are shown in Figure 7.12. The main quantity of interest is the exponent \( \nu \). We know that \( g = 0 \) corresponds

Figure 7.12: Example of the evolution of the true self-avoiding walk with \( g = 1 \) (see (7.48)). The shaded site represents the location of the walker at time \( t \). The number of visits to each site are given within each site and the probability of a step to a nearest neighbor site is given below it. Note the use of periodic boundary conditions.
to the usual random walk with $\nu = 1/2$ and that the limit $g \to \infty$ corresponds to the self-avoiding walk. What is the value of $\nu$ for a self-avoiding walk in one dimension? Is the value of $\nu$ for any finite value of $g$ different than these two limiting cases?

Write a program to do a Monte Carlo simulation of the true self-avoiding walk in one dimension. Use an array to record the number of visits to every site. At each step calculate the probability $p$ of a jump to the right. Generate a random number $r$ and compare it to $p$. If $r \leq p$, move the walker to the right; otherwise move the walker to the left. Compute $\langle \Delta x^2 \rangle$, where $x$ is the distance of the walker from the origin, as a function of the number of steps $N$. Make a log-log plot of $\langle \Delta x^2 \rangle$ versus $N$ and estimate $\nu$. Can you distinguish $\nu$ from its random walk and self-avoiding walk values? Reasonable choices of parameters are $g = 0.1$ and $N \sim 10^3$. Averages over $10^3$ trials yield qualitative results. For comparison, published results are for $N \sim 10^4$ and for $10^3$ trials; extended results for $g = 2$ are given for $N = 2 \times 10^5$ and $10^4$ trials (see Bernasconi and Pietronero).

### 7.8 Diffusion Controlled Chemical Reactions

Imagine a system containing particles of a single species $A$. The particles diffuse, and when two particles “collide,” a reaction occurs such that the two combine to form an inert species which is no longer involved in the reaction. We can represent this chemical reaction as

$$A + A \to 0. \quad (7.49)$$

If we ignore the spatial distribution of the particles, we can describe the kinetics by a simple rate equation:

$$\frac{dA(t)}{dt} = -kA^2(t), \quad (7.50)$$

where $A$ is the concentration of $A$ particles at time $t$ and $k$ is the rate constant. (In the chemical kinetics literature it is traditional to use the term concentration rather than the number density.) For simplicity, we assume that all reactants are entered into the system at $t = 0$ and that no reactants are added later (the system is closed). It is easy to show that the solution of the first-order differential equation (7.50) is

$$A(t) = \frac{A(0)}{1 + ktA(0)}. \quad (7.51)$$

Hence, $A(t) \sim t^{-1}$ in the limit of long times.

Another interesting case is the bimolecular reaction

$$A + B \to 0. \quad (7.52)$$

If we neglect spatial fluctuations in the concentration as before (this neglect yields what is known as a mean-field approximation), we can write the corresponding rate equation as

$$\frac{dA(t)}{dt} = \frac{dB(t)}{dt} = -kA(t)B(t). \quad (7.53)$$
We also have that
\[ A(t) - B(t) = \text{constant}, \]  
\tag{7.54} 
because each reaction leaves the difference between the concentration of \( A \) and \( B \) particles unchanged. For the special case of equal initial concentrations, the solution of (7.53) with (7.54) is the same as (7.51). What is the solution for the case \( A(0) \neq B(0) \)?

This derivation of the time dependence of \( A \) for the kinetics of the one and two species annihilation process is straightforward, but is based on the assumption that the particles are distributed uniformly. In the following two problems, we simulate the kinetics of these processes and test this assumption.

**Problem 7.33.** Diffusion controlled chemical reactions in one dimension

a. Assume that \( N \) particles do a random walk on a one-dimensional lattice of length \( L \) with periodic boundary conditions. Every particle moves once in one unit of time. Use the array \( \text{site}[j] \) to record the label of the particle, if any, at site \( j \). Because we are interested in the long time behavior of the system when the concentration \( A = N/L \) of particles is small, it is efficient to also maintain an array of particle positions, \( x[i] \) such that \( \text{site}[x[i]] = i \). For example, if particle 5 is located at site 12, then \( x[5] = 12 \) and \( \text{site}[12] = 5 \). We also need an array, \( \text{newSite} \), to maintain the new positions of the walkers as they are moved one at a time. After each walker is moved, we check to see if two walkers have landed on the same position \( k \). If they have, we set \( \text{newSite}[k] = -1 \), and the value of \( x[i] \) for these two walkers to \(-1\). The value \(-1\) indicates that no particle exists at the site. After all the walkers have moved, we let \( \text{site} = \text{newSite} \) for all sites, and remove all the reacting particles in \( x \) that have values equal to \(-1\). This operation can be accomplished by replacing any reacting particle in \( x \) by the last particle in the array. Begin with all sites occupied, \( A(t = 0) = 1 \).

b. Make a log-log plot of the quantity \( A(t)^{-1} - A(0)^{-1} \) versus the time \( t \). The times should be separated by exponential intervals so that your data is equally spaced on a logarithmic plot. For example, you might include data with times equal to \( 2^p \), with \( p = 1, 2, 3, \ldots \). Does your log-log plot yield a straight line for long times? If so, calculate its slope. Is the mean-field approximation for \( A(t) \) valid in one dimension? You can obtain crude results for small lattices of order \( L = 100 \) and times of order \( t = 10^2 \). To obtain results to within 10%, you will need lattices of order \( L = 10^4 \) and times of order \( t = 2^{13} \).

c. More insight into the origin of the time dependence of \( A(t) \) can be gained from the behavior of the quantity \( P(r, t) \), the probability that the nearest neighbor distance is \( r \) at time \( t \). The nearest neighbor distance of a given particle is defined as the minimum distance between it and all other particles. The distribution of these distances changes dramatically as the reaction proceeds, and this change can give information about the reaction mechanism. Place the particles at random on a one-dimensional lattice and verify that the most probable nearest neighbor distance is \( r = 1 \) (one lattice constant) for all concentrations. (This result is true in any dimension.) Then verify that the distribution of nearest neighbor distances on a one-dimensional lattice is given by
\[ P(r, t = 0) = 2A e^{-2A(r-1)}, \]  
\tag{random distribution}  
Is the form (7.55) properly normalized? Start with \( A(t = 0) = 0.1 \) and find \( P(r, t) \) for \( t = 10, 100, \) and \( 1000 \). Average over all particles. How does \( P(r, t) \) change as the reaction proceeds? Does it retain the same form as the concentration decreases?
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d.\(^*\) Compute the quantity \(D(t)\), the number of distinct sites visited by an individual walker. How does the time dependence of \(D(t)\) compare to the computed time dependence of \(A(t)^{-1} - 1\)?

e.\(^*\) Write a program to simulate the reaction \(A + B = 0\). For simplicity, assume that multiple occupancy of the same site is not allowed, for example, an \(A\) particle cannot jump to a site already occupied by an \(A\) particle. The easiest procedure is to allow a walker to choose one of its nearest neighbor sites at random, but to not move the walker if the chosen site is already occupied by a particle of the same type. If the site is occupied by a walker of another type, then the pair of reacting particles is annihilated. Keep separate arrays for the \(A\) and \(B\) particles, with the value of the array denoting the label of the particle as before. One way to distinguish \(A\) and \(B\) walkers is to make the array element \(\text{site}(k)\) positive if the site is occupied by an \(A\) particle and negative if the site is occupied by a \(B\) particle. Start with equal concentrations of \(A\) and \(B\) particles and occupy the sites at random. Some of the interesting questions are similar to those that we posed in parts (b)–(d). Color code the particles and observe what happens to the relative positions of the particles.

\(^*\)Problem 7.34. Reaction diffusion in two dimensions

a. Do a similar simulation as in Problem 7.33 on a two-dimensional lattice for the reaction \(A + A \rightarrow 0\). In this case it is convenient to have one array for each dimension, for example, \(\text{siteX}\) and \(\text{siteY}\), or to store the lattice as a one-dimensional array (see Section 13.2). Set \(A(t = 0) = 1\), and choose \(L = 50\). Show the walkers after each Monte Carlo step per walker and describe their distribution as they diffuse. Are the particles uniformly distributed throughout the lattice for all times? Calculate \(A(t)\) and compare your results for \(A(t)^{-1} - A(0)^{-1}\) to the \(t\)-dependence of \(D(t)\), the number of distinct lattice sites that are visited in time \(t\). (In two dimensions, \(D(t) \sim t / \log t\).) How well do the slopes compare? Do a similar simulation with \(A(t = 0) = 0.01\). What slope do you obtain in this case? What can you conclude about the initial density dependence? Is the mean-field approximation valid in this case?

b. Begin with \(A\) and \(B\) type random walkers initially segregated on the left and right halves (in the \(x\) direction) of a square lattice. The process \(A + B \rightarrow C\) exhibits a reaction front where the production of particles of type \(C\) is nonzero. Some of the quantities of interest are the time dependence of the mean position \(\langle x \rangle(t)\) and the width \(w(t)\) of the reaction front. The rules of this process are the same as in part (a) except that a particle of type \(C\) is added to a site when a reaction occurs. A particular site can be occupied by one particle of type \(A\) or type \(B\) as well as any number of particles of type \(C\). If \(n(x,t)\) is the number of particles of type \(C\) at a distance \(x\) from the initial boundary of the reactants, then \(x(t)\) and \(w(t)\) can be written as

\[
\langle x \rangle(t) = \frac{\sum_x x n(x,t)}{\sum_x n(x,t)} \quad (7.56)
\]

\[
w(t)^2 = \frac{\sum_x [x - \langle x \rangle(t)]^2 n(x,t)}{\sum_x n(x,t)} \quad (7.57)
\]

Choose lattice sizes of order 100 \(\times\) 100, and average over at least 10 trials. The fluctuations in \(x(t)\) and \(w(t)\) can be reduced by averaging \(n(x,t)\) over the order of 100 time units centered about \(t\). More details can be found in Jiang and Ebner.
7.9 Random Number Sequences

So far we have used the random number generator supplied with Java to generate the desired random numbers. In principle, we could have generated these numbers from a random physical process, such as the decay of radioactive nuclei or the thermal noise from a semiconductor device. In practice, random number sequences are generated from a physical process only for specific purposes such as a lottery. Although we could store the outcome of a random physical process so that the random number sequence would be both truly random and reproducible, such a method would usually be inconvenient and inefficient in part because we often require very long sequences. In practice, we use a digital computer, a deterministic machine, to generate sequences of pseudorandom numbers. Although these sequences cannot be truly random, such a distinction is unimportant if the sequence satisfies all our criteria for randomness. It is common to refer to random number generators even though we really mean pseudorandom number generators.

Most random number generators yield a sequence in which each number is used to find the succeeding one according to a well defined algorithm. The most important features of a desirable random number generator are that its sequence satisfies the known statistical tests for randomness, which we will explore in the following problems. We also want the generator to be efficient and machine independent, and the sequence to be reproducible.

The most widely used random number generator is based on the linear congruential method. One advantage of the linear congruential method is that it is very fast. For a given seed $x_0$, each number in the sequence is determined by the one-dimensional map

$$x_n = (ax_{n-1} + c) \mod m,$$

where $a$, $c$, and $m$ as well as $x_n$ are integers. The notation $y = z \mod m$ means that $m$ is subtracted from $z$ until $0 \leq y < m$. The map (7.58) is characterized by three parameters, the multiplier $a$, the increment $c$, and the modulus $m$. Because $m$ is the largest integer generated by (7.58), the maximum possible period is $m$.

In general, the period depends on all three parameters. For example, if $a = 3$, $c = 4$, $m = 32$, and $x_0 = 1$, the sequence generated by (7.58) is 1, 7, 25, 15, 17, 23, 9, 31, 1, 7, 25, ..., and the period is 8 rather than the maximum possible value of $m = 32$. If we choose $a$, $c$, and $m$ carefully such that the maximum period is obtained, then all possible integers between 0 and $m - 1$ would occur in the sequence. Because we usually wish to have random numbers $r$ in the unit interval $0 \leq r < 1$ rather than random integers, random number generators usually return the ratio $x_n/m$ which is always less than unity. Several rules have been developed (see Knuth) to obtain the longest period. Some of the properties of the linear congruential method are explored in Problem 7.35.

Another popular random number generator is the generalized feedback shift register method which uses bit manipulation (see Sections 15.1 and 15.6). Every integer is represented as a series of 1s and 0s called bits. These bits can be shuffled by using the bitwise exclusive or (xor) operator defined by $a \oplus b = 1$ if the bits $a \neq b$; $a \oplus b = 0$ if $a = b$. The $n$th member of the sequence is given by

$$x_n = x_{n-p} \oplus x_{n-q},$$

where $p > q$, and $p$, $q$, and $x_n$ are integers. The first $p$ random integers must be supplied by another random number generator. As an example of how the operator $\oplus$ works, suppose that $n = 6$, $p = 5$, $q = 3$, $x_3 = 11$, and $x_1 = 6$. Then $x_6 = x_1 \oplus x_3 = 0110 \oplus 1011 = 1101 = 2^3 + 2^2 + 2^0 = 8 + 4 + 1 = 13.$
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Not all values of \( p \) and \( q \) lead to good results. Some common pairs are \((p, q) = (31, 3), (250, 103), \) and \((521, 168)\).

In Java and C the exclusive or operation on the integers \( m \) and \( n \) is written as \( m \oplus n \). The algorithm for producing the random numbers after \( p \) integers have been produced is shown in the following. Initially the index, \( k \), can be set to 0.

1. If \( k < q \), set \( j = k + q \), else set \( j = k - p + q \).
2. Set \( x_k = x_k \oplus x_j \); \( x_k \) is the desired random number for this iteration. If a random number between 0 and 1 is desired, divide \( x_k \) by the maximum possible integer that the computer can hold.
3. Increment \( k \) to \((k + 1) \mod p\).

Because the exclusive or operator and bit manipulation is very fast, this random number generator is very fast. However, the period may not be long enough for some applications and the correlations between numbers might not be as good as needed. The shuffling algorithm discussed in Problem 7.36 should be used to improve this generator.

These two examples of random number generators illustrate their general nature. That is, numbers in the sequence are used to find the succeeding ones according to a well defined algorithm. The sequence is determined by the seed, the first number of the sequence, or the first \( p \) members of the sequence for the generalized feedback shift register and related generators. Usually, the maximum possible period is related to the size of the computer word, for example, 32, or 64 bits. The choice of the constants and the proper initialization of the sequence is very important and thus these algorithms must be implemented with care.

There is no necessary and sufficient test for the randomness of a finite sequence of numbers; the most that can be said about any finite sequence of numbers is that it is apparently random. Because no single statistical test is a reliable indicator, we need to consider several tests. Some of the best known tests are discussed in Problem 7.35. Many of these tests can be stated in terms of random walks.

Problem 7.35. Statistical tests of randomness

a. Period. An obvious requirement for a random number generator is that its period be much greater than the number of random numbers needed in a specific calculation. One way to visualize the period of the random number generator is to use it to generate a plot of the displacement \( x \) of a random walker as a function of the number of steps \( N \). When the period of the random number is reached, the plot will begin to repeat itself. Generate such a plot using (7.58) for \( a = 899, c = 0, \) and \( m = 32768, \) and for \( a = 16807, c = 0, m = 2^{31} - 1 \) with \( x_0 = 12 \). What are the periods of the corresponding random number generators? Obtain similar plots using different values for the parameters \( a, c, \) and \( m \). Why is the seed value \( x_0 = 0 \) forbidden for the choice \( c = 0? \) Do some combinations of \( a, c, \) and \( m \) give longer periods than others?

b. Uniformity. A random number sequence should contain numbers distributed in the unit interval with equal probability. The simplest test of uniformity is to divide this interval into \( M \) equal size subintervals or bins. For example, consider the first \( N = 10^4 \) numbers generated by (7.58)
with $a = 106$, $c = 1283$, and $m = 6075$ (see Press et al.). Place each number into one of $M = 100$ bins. Is the number of entries in each bin approximately equal? What happens if you increase $N$?

c. **Chi-square test.** Is the distribution of numbers in the bins of part (b) consistent with the laws of statistics? The most common test of this consistency is the *chi-square* or $\chi^2$ test. Let $y_i$ be the observed number in bin $i$ and $E_i$ be the expected value. The chi-square statistic is

$$\chi^2 = \sum_{i=1}^{M} \frac{(y_i - E_i)^2}{E_i}. \quad (7.60)$$

For the example in part (b) with $N = 10^4$ and $M = 100$, we have $E_i = 100$. The magnitude of the number $\chi^2$ is a measure of the agreement between the observed and expected distributions; $\chi^2$ should not be too big or too small. In general, the individual terms in the sum (7.60) are expected to be order one, and because there are $M$ terms in the sum, we expect $\chi^2 \leq M$. As an example, we did five independent runs of a random number generator with $N = 10^4$ and $M = 100$, and found $\chi^2 \approx 92, 124, 85, 91,$ and $99$. These values of $\chi^2$ are consistent with this expectation. Although we usually want $\chi^2$ to be as small as possible, we would be suspicious if $\chi^2 \approx 0$, because such a small value suggests that $N$ is a multiple of the period of the generator and that each value in the sequence appears an equal number of times.

d. **Filling sites.** Although a random number sequence might be distributed in the unit interval with equal probability, the consecutive numbers might be correlated in some way. One test of this correlation is to fill a square lattice of $L^2$ sites at random. Consider an array $n(x,y)$ that is initially empty, where $1 \leq x, y \leq L$. A site is selected randomly by choosing its two coordinates $x_i$ and $y_i$ from two consecutive numbers in the sequence. If the site is empty, it is filled and $n(x_i,y_i) = 1$; otherwise it is not changed. This procedure is repeated $t$ times, where $t$ is the number of Monte Carlo steps per site. That is, the time is increased by $1/L^2$ each time a pair of random numbers is generated. Because this process is analogous to the decay of radioactive nuclei, we expect that the fraction of empty lattice sites should decay as $e^{-t}$. Determine the fraction of unfilled sites using the random number generator that you have been using for $L = 10, 15, \text{and } 20$. Are your results consistent with the expected fraction? Repeat the same test using (7.58) with $a = 65549, c = 0, \text{and } m = 231$. The existence of triplet correlations can be determined by a similar test on a simple cubic lattice by choosing the three coordinates $x_i, y_i, \text{and } z_i$ from three consecutive random numbers.

e. **Parking lot test.** Fill sites as in part (d) and draw the sites that have been filled. Do the filled sites look random, or are there stripes of filled sites? Try $a = 65549, c = 0, \text{and } m = 231$.

f. **Hidden correlations.** Another way of checking for correlations is to plot $x_{i+k}$ versus $x_i$. If there are any obvious patterns in the plot, then there is something wrong with the generator. Use the generator (7.58) with $a = 16807, c = 0, \text{and } m = 2^{31} - 1$. Can you detect any structure in the plotted points for $k = 1$ to $k = 5$? Test the random number generator that you have been using. Do you see any evidence of lattice structure, for example, equidistant parallel lines? Is the logistic map $x_{n+1} = 4x_n(1-x_n)$ a suitable random number generator?
g. **Short-term correlations.** Another measure of short term correlations is the autocorrelation function

\[ C(k) = \frac{\langle x_{i+k}x_i \rangle - \langle x_i \rangle^2}{\langle x_i^2 \rangle - \langle x_i \rangle^2}, \tag{7.61} \]

where \( x_i \) is the \( i \)th term in the sequence. We have used the fact that \( \langle x_{i+k} \rangle = \langle x_i \rangle \), that is, the choice of the origin of the sequence is irrelevant. The quantity \( \langle x_{i+k}x_i \rangle \) is found for a particular choice of \( k \) by forming all the possible products of \( x_{i+k}x_i \) and dividing by the number of products. If \( x_{i+k} \) and \( x_i \) are not correlated, then \( \langle x_{i+k}x_i \rangle = \langle x_{i+k} \rangle \langle x_i \rangle \) and \( C(k) = 0 \). Is \( C(k) \) identically zero for any finite sequence? Compute \( C(k) \) for \( a = 106, c = 1283, \) and \( m = 6075 \).

h. **Random walk.** A test based on the properties of random walks has been proposed by Vattulainen et al. Assume that a walker begins at the origin of the \( x-y \) plane and walks for \( N \) steps. Average over \( M \) walkers and count the number of walks that end in each quadrant \( q_i \). Use the \( \chi^2 \) test (7.60) with \( y_i \rightarrow q_i, M = 4, \) and \( E_i = M/4 \). If \( \chi^2 > 7.815 \) (a 5% probability if the random number generator is perfect), we say that the run fails. The random number generator fails if two out of three independent runs fail. The probability of a perfect generator failing two out of three runs is approximately \( 3 \times 0.95 \times (0.05)^2 \approx 0.007 \). Test several random number generators.

**Problem 7.36.** Improving random number generators

One way to reduce sequential correlation and to lengthen the period is to mix or shuffle the random numbers produced by a random number generator. A standard procedure is to begin with a list of \( N \) random numbers (between 0 and 1) using a given generator \( \text{rng} \). The number \( N \) is arbitrary, but should be less than the period of \( \text{rng} \). Also generate one more random number, \( r_{\text{extra}} \). Then for each desired random number use the following procedure.

(i) Calculate the integer \( k \) given by \( \text{int}(N \times r_{\text{extra}}) \). Use the \( k \)th random number, \( r_k \), from your list as the desired random number.

(ii) Set \( r_{\text{extra}} \) equal to the random number, \( r_k \), chosen in step (i).

(iii) Generate a new random number, \( r \), from \( \text{rng} \) and use it to replace the number chosen in step (i), that is, \( r_k = r \).

Consider a random number generator with a relatively short period and strong sequential correlation and show that this shuffling scheme improves the quality of the random number sequence.

At least some of the statistical tests given in Problem 7.35 should be done whenever serious calculations are contemplated. However, even if a random number generator passes all these tests, there still can be problems in rare cases. Typically, these problems arise when a small number of events have a large weight. In these cases a very small bias in the random number generator might lead to systematic errors, and two generators, which appear equally good as determined by various statistical tests, might give statistically different results in a specific application (see Project 16.34). For this reason, it is important that the particular random number generator used be reported along with the actual results. Confidence in the results also can be increased by repeating the calculation with another random number generator.

Because all random number generators are based on a deterministic algorithm, it always is possible to construct a test generator for which a particular algorithm will fail. The success of a
random number generator in passing various statistical tests is necessary, but it is not a sufficient condition for its use in all applications. In Project 16.34 we discuss an application of Monte Carlo methods to the Ising model for which some popular random number generators give incorrect results.

7.10 Variational Methods

Many problems in physics can be formulated in terms of a variational principle. In the following, we consider examples of variational principles in geometrical optics and classical mechanics. We then discuss how Monte Carlo methods can be applied to these problems. A more sophisticated application of Monte Carlo methods to a variational problem in quantum mechanics is discussed in Chapter 17.

Our everyday experience of light leads naturally to the concept of light rays. This description of light propagation, called geometrical or ray optics, is applicable when the wavelength of light is small compared to the linear dimensions of any obstacles or openings. The path of a light ray can be formulated in terms of Fermat’s principle of least time: A ray of light follows the path between two points (consistent with any constraints) that requires the least amount of time. Fermat’s principle can be adopted as the basis of geometrical optics. For example, Fermat’s principle implies that light travels from a point $A$ to a point $B$ in a straight line in a homogeneous medium. Because the speed of light is constant along any path within the medium, the path of shortest time is the path of shortest distance, that is, a straight line from $A$ to $B$. What happens if we impose the constraint that the light must strike a mirror before reaching $B$?

The speed of light in a medium can be expressed in terms of $c$, the speed of light in a vacuum, and the index of refraction $n$ of the medium:

$$v = \frac{c}{n}.$$  \hspace{1cm} (7.62)

Suppose that a light ray in a medium with index of refraction $n_1$ passes through a second medium with index of refraction $n_2$. The two media are separated by a plane surface. We now show how we can use Fermat’s principle and a simple Monte Carlo method to find the path of the light. The analytical solution to this problem using Fermat’s principle is found in many texts (cf. Feynman et al.).

Our strategy, as implemented in class Fermat, is to begin with a straight path and to make changes in the path at random. These changes are accepted only if they reduce the travel time of the light. Some of the features of Fermat and FermatApp include:

1. Light propagates from left to right through $N$ regions. The index of refraction $n[i]$ is uniform in each region $[i]$. The index $i$ increases from left to right. We have chosen units such that the speed of light in vacuum equals unity.

2. Because the light propagates in a straight line in each medium, the path of the light is given by the coordinates $y[i]$ at each boundary.

3. The coordinates of the light source and the detector are at $(0,y[0])$ and $(N,y[N])$ respectively, where $y[0]$ and $y[N]$ are fixed.
4. The path is the connection of the set of points at the boundary of each region.

5. The path of the light is found by choosing the boundary $i$ at random and generating a trial value of $y[i]$ that differs from its previous value by a random number between $-dy$ to $dy$. If the trial value of $y[i]$ yields a shorter travel time, this value becomes the new value for $y[i]$.

6. The path is redrawn whenever it is changed.

Listing 7.6: Fermat class.

```java
package org.opensourcephysics.sip.ch07;

public class Fermat {
    double y[]; // y coordinate of light ray, index is x coordinate
    double v[]; // light speed of ray for medium starting at index value
    int N; // number of media
    double dn; // change in index of refraction from one region to the next
    double dy = 0.1; // maximum change in y position
    int steps;

    public void initialize() {
        y = new double[N+1];
        v = new double[N];
        double indexOfRefraction = 1.0;
        for(int i = 0; i =N; i++) {
            y[i] = i; // initial path is a straight line
        }
        for(int i = 0; i <N; i++) {
            v[i] = 1.0/indexOfRefraction;
            indexOfRefraction += dn;
        }
        steps = 0;
    }

    public void step() {
        int i = 1+(int) (Math.random()*(N-1));
        double yTrial = y[i]+2.0*dy*(Math.random() - 0.5);
        double previousTime = Math.sqrt(Math.pow(y[i-1]-y[i], 2)+1)/v[i-1]; // left medium
        previousTime += Math.sqrt(Math.pow(y[i+1]-y[i], 2)+1)/v[i]; // right medium
        double trialTime = Math.sqrt(Math.pow(y[i-1]-yTrial, 2)+1)/v[i-1]; // left medium
        trialTime += Math.sqrt(Math.pow(y[i+1]-yTrial, 2)+1)/v[i]; // right medium
        if(trialTime<previousTime) {
            y[i] = yTrial;
        }
        steps++;
    }
}
```

Listing 7.7: Target class for Fermat’s principle.

```java
package org.opensourcephysics.sip.ch07;
```
Problem 7.37. The law of refraction

a. Use Fermat and FermatApp to determine the angle of incidence $\theta_1$ and the angle of refraction $\theta_2$ between two media with different indices of refraction. The angles $\theta_1$ and $\theta_2$ are measured from the normal to the boundary. Set $N = 2$ and let the first medium be air ($n_1 \approx 1$) and the second medium be glass ($n_2 \approx 1.5$). Describe the path of the light after a number of trial paths are attempted. Add statements to the program to determine $\theta_1$ and $\theta_2$, the vertical position of the intersection of the light at the boundary between the two media, and the total time for the
Figure 7.13: Near the horizon, the apparent (exaggerated) position of the sun is higher than the true position of the sun. Note that the light rays from the true sun are curved due to refraction.

light to go from $(0, y_0)$ to $(2, y_2)$.

b. Modify the program so that the first medium represents glass ($n_1 \approx 1.5$) and the second medium represents water ($n_2 \approx 1.33$). Verify that your results are consistent with $n_2 \sin \theta_2 = n_1 \sin \theta_1$.

Problem 7.38. Inhomogeneous media

a. The earth’s atmosphere is thin at the top and dense near the earth’s surface. We can model this inhomogeneous medium by dividing the atmosphere into equal width segments each of which is homogeneous. To simulation this atmosphere run your program with $N = 10$ and $\Delta n = 0.1$, and find the path of least time. Use your results to explain why when we see the sun set, the sun already is below the horizon (see Figure 7.13).

b.* Modify your program to find the appropriate distribution $n(y)$ for a fiber optic cable, which we take to be a flat, long ribbon. In this case the $i$th region corresponds to a cross sectional slab through the cable. Although a real cable is three-dimensional, we consider a two-dimensional cable for simplicity. We want the cable to have the property that if a ray of light starts from one side of the cable and ends at the other, the slope $dy/dx$ of the path should be near zero at the edges so that light does not escape from the cable.

Fermat’s principle is an example of an extremum principle. An extremum means that a small change $\epsilon$ in an independent variable leads to a change in a function (more precisely, a function of functions) that is proportional to $\epsilon^2$ or a higher power of $\epsilon$. An important extremum principle in classical mechanics is based on the action $S$:

$$S = \int_{t_0}^{t_{\text{final}}} L \, dt,$$  \hspace{1cm} (7.63)

where $t_0$ and $t_{\text{final}}$ are the initial and final times, respectively. The Lagrangian $L$ in (7.63) is the kinetic energy minus the potential energy. The extremum principle for the action is known as the principle of least action or Hamilton’s action principle. The path where (7.63) is stationary (either
a minimum or a saddle point) satisfies Newton’s second law (for conservative forces). One reason for the importance of the principle of least action is that quantum mechanics can be formulated in terms of an integral over the action (see Section 17.10).

To use (7.63) to find the motion of a single particle in one dimension, we fix the position at the chosen initial and final times, \(x(t_0)\) and \(x(t_{\text{final}})\), and then choose the velocities and positions for the intermediate times \(t_0 < t < t_{\text{final}}\) to minimize the action. One way to implement this procedure numerically is to convert the integral in (7.63) to a sum:

\[
S \approx \sum_{i=1}^{N-1} L(t_i) \Delta t, \tag{7.64}
\]

where \(t_i = t_0 + i\Delta t\). (The approximation used to obtain (7.64) is known as the rectangular approximation and is discussed in Chapter 12.) For a single particle in one dimension moving in an external potential \(u(x)\), we can write

\[
L_i \approx \frac{m}{2(\Delta t)^2}(x_{i+1} - x_i)^2 - u(x_i), \tag{7.65}
\]

where \(m\) is the mass of the particle and \(u(x_i)\) is the potential energy of the particle at \(x_i\). The velocity has been approximated as the difference in position divided by the change in time \(\Delta t\).

**Problem 7.39.** Principle of least action

a. Write a program to minimize the action \(S\) given in (7.63) for the motion of a single particle in one dimension. Use the approximate form of the Lagrangian given in (7.65). One way to write the program is to modify class `Fermat` so that the vertical coordinate for the light ray becomes the position of the particle, and the horizontal region number \(i\) becomes the discrete time interval of duration \(\Delta t\).

b. Verify your program for the case of free fall for which the potential energy is \(u(y) = mgy\). Choose \(y(t = 0) = 2\) m and \(y(t = 10\) s\) = 8 m, and begin with \(N = 20\). Allow the maximum change in the position to be 5 m.

c. Consider the harmonic potential \(u(x) = \frac{1}{2}kx^2\). What shape do you expect the path \(x(t)\) to be? Increase \(N\) to approximately 50 and estimate the path by minimizing the action.

It is possible to extend the principle of least action to more dimensions or particles, but it is necessary to begin with a path close to the optimum one to obtain a good approximation to the optimum path in a reasonable time.

In Problems 7.37–7.39 a simple Monte Carlo algorithm that always accepts paths that reduce the time or action is sufficient. However, for more complicated index of refraction distributions or potentials, it is possible that such a simple algorithm will find only a local minimum and the global minimum will be missed. The problem of finding the global minimum is very general and is shared by all optimization algorithms if the system has many relative minima. Optimization is a very active area of research in many fields of science and engineering. Ideas from physics, biology, and computer science have led to many improved algorithms. We will discuss some of these algorithms in Chapter 16. In most of these algorithms paths that are worse than the current path...
are sometimes accepted in an attempt to climb out of a local minimum. Other algorithms involve ways of sampling over a wider range of possible paths. Another approach is to convert the Monte Carlo algorithm into a deterministic algorithm. We have already mentioned that an analytical variational calculation leads to Newton’s second law. Passerone and Parrinello discuss an algorithm for looking for extrema in the action by maintaining the discrete structure in Eq. (7.65), and then finding the extremum by taking the derivative with respect to each coordinate, \( x_i \) and setting the resulting equations equal to zero. This procedure leads to a set of deterministic equations that need to be solved numerically. The performance can be improved by enforcing energy conservation and using some other tricks.

### 7.11 Projects

Almost all of the problems in this chapter can be done using more efficient programs, greater number of trials, and larger systems. More applications of random walks and random number sequences are discussed in subsequent chapters. Many more ideas for projects can be gained from the references.

**Project 7.40.** Competition between diffusion and fragmentation

As we have discussed, random walks are useful for understanding diffusion in contexts more general than the movement of a particle. Consider a particle in solution whose mass can grow either by the absorption of particles or shrink by the loss of small particles, including fragmentation. We can model this process as a random walk by replacing the position of the particle by its mass. One difference between this case and the random walks we have studied so far is that the random variable, the mass, must be positive. The model of Ferkinghoff-Berg et al. can be summarized as follows:

(i) Begin with \( N \) objects with some distribution of lengths. Let the integer \( L_i \) represent the length of the \( i \)th object.

(ii) All the objects change their length by \( \pm 1 \). This step is analogous to a random walk. If the length of an object becomes equal to 0, it is removed from the system. An easy way to eliminate the \( i \)th object is to set its length equal to the length of the last object and reduce \( N \) by unity.

(iii) Choose one object at random with a probability that is proportional to the length of the object. Fragment this object into two objects, where the fraction of the mass going to each object is random.

(iv) Repeat steps (ii) and (iii).

a. Write a program to implement this algorithm in one dimension. One way to implement step (iii) is given in the following code, where `totalMass` is the sum of the lengths of all the objects.

```c
int i = 0; // label of object
int sum = length[0]; // length of first object, all lengths are integers
// choose object to fragment so that choice is proportional to length
```
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```java
int x = (int)(Math.random() * totalMass);
while (sum < x) {
    i++;
    sum += length[i];
}
// if object big enough to fragment, choose random fraction for each part
if (length[i] > 1) {
    int partA = 1 + (int)(Math.random() * (length[i] - 1));
    int partB = length[i] - partA;
    length[i] = partA;
    length[numberOfObjects] = partB;  // new object
    numberOfObjects++;
}
```

The main quantity of interest is the distribution of lengths \( P(L) \). Explore a variety of initial length distributions with a total mass of 5000 for which the distribution is peaked at about 20 mass units. Is the long time behavior of \( P(L) \) similar in shape for any initial distribution? Compute the total mass (sum of the lengths) and output this value periodically. Although the total mass will fluctuate, it should remain approximately constant. Why?

b. Collect data for three different initial distributions with the same number of objects \( N \), and scale \( P(L) \) and \( L \) so that the three distributions roughly fall on the same curve. For example, you can scale \( P(L) \) so that the maximum of the three distributions has the same value. Then multiply each value of \( L \) by a factor so that the distributions overlap.

c. The analytical results suggest that the universal behavior can be obtained by scaling \( L \) by the total mass raised to the 1/3 power. Is this prediction consistent with your results? Test this hypothesis by adjusting the initial distributions so that they all have the same total mass. Your results for the long time behavior of \( P(L) \) should fall on a universal curve. Why is this universality interesting? How can this result be used to analyze different systems? Would you need to do a new simulation for each value of \( L \)?

d. What happens if step (iii) is done more or less often than each random change of length. Does the scaling change?

**Project 7.41. Application of the pivot algorithm to self-avoiding walks**

The algorithms that we have discussed for generating self-avoiding random walks are all based on making local deformations of the walk (polymer chain) for a given value of \( N \), the number of bonds. As discussed in Problem 7.31, the time \( \tau \) between statistically independent configurations is nonzero. The problem is that \( \tau \) increases with \( N \) as some power, for example, \( \tau \sim N^3 \). This power law dependence of \( \tau \) on \( N \) is called critical slowing down and implies that it becomes increasingly more time consuming to generate long walks. We now discuss an example of a global algorithm that reduces the dependence of \( \tau \) on \( N \). Another example of a global algorithm that reduces critical slowing down is discussed in Project 16.32.

a. Consider the walk shown in Figure 7.14a. Select a site at random and one of the four possible directions. The shorter portion of the walk is rotated (pivoted) to this new direction by treating the walk as a rigid structure. The new walk is accepted only if the new walk is self-avoiding;
Figure 7.14: Examples of the first several changes generated by the pivot algorithm for a self-avoiding walk of \( N = 10 \) bonds (11 sites). The open circle denotes the pivot point. This figure is adopted from the article by MacDonald et al.

otherwise the old walk is retained. (The shorter portion of the walk is chosen to save computer time.) Some typical moves are shown in Figure 7.14. Note that if an end point is chosen, the previous walk is retained. Write a program to implement this algorithm and compute the dependence of the mean square end-to-end distance \( R^2 \) on \( N \). Consider values of \( N \) in the range \( 10 \leq N \leq 80 \). A discussion of the results and the implementation of the algorithm can be found in MacDonald et al. and Madras and Sokal, respectively.

b. Compute the correlation time \( \tau \) for different values of \( N \) using the approach discussed in Problem 7.31b.

Project 7.42. Pattern formation

In Problem 7.34 we saw that simple patterns can develop as a result of random behavior. The phenomenon of pattern formation is of much interest in a variety of contexts ranging from the large scale structure of the universe to the roll patterns seen in convection (for example, smoke rings). In the following, we explore the patterns that can develop in a simple reaction diffusion model based on the reactions, \( A + 2B \rightarrow 3B \), and \( B \rightarrow C \), where \( C \) is inert. Such a reaction is called autocatalytic.

In Problem 7.34 we considered chemical reactions in a closed system where the reactions can proceed to equilibrium. In contrast, open systems allow a continuous supply of fresh reactants and a removal of products. These two processes allow steady states to be realized and oscillatory conditions to be maintained indefinitely. In this problem we assume that \( A \) is added at a constant rate and that both \( A \) and \( B \) are removed by the feed process. Pearson (see references) modeled
these processes by two coupled reaction diffusion equations:

\[
\begin{align*}
\frac{\partial A}{\partial t} &= D_A \nabla^2 A - AB^2 + f(1 - A) \quad \text{(7.66a)} \\
\frac{\partial B}{\partial t} &= D_B \nabla^2 B + AB^2 - (f + k)B. \quad \text{(7.66b)}
\end{align*}
\]

The $AB^2$ term represents the reaction $A + 2B \rightarrow 3B$. This term is negative in (7.66a) because the reactant $A$ decreases, and is positive in (7.66b) because the reactant $B$ increases. The term $+f$ represents the constant addition of $A$, and the terms $-fA$ and $-fB$ represent the removal process; the term $-kB$ represents the reaction $B \rightarrow C$. All the quantities in (7.66) are dimensionless. We assume that the diffusion coefficients are $D_A = 2 \times 10^{-5}$ and $D_B = 10^{-5}$, and the behavior of the system is determined by the values of the rate constant $k$ and the feed rate $f$.

These two coupled equations can be solved subject to appropriate initial and boundary conditions. However, without the diffusion terms, it is more straightforward to solve (7.66) numerically using a simple Euler algorithm. Choose a time step equal to unity, and let $A = 0.1$ and $B = 0.5$ at $t = 0$. Determine the steady state values for $0 < f \leq 0.3$ and $0 < k \leq 0.07$ in increments of $\Delta f = 0.02$ and $\Delta k = 0.005$. Record the steady state values of $A$ and $B$. Then repeat this exercise for the initial values $A = 0.5$ and $B = 0.1$. You should find that for some values of $f$ and $k$, only one steady state solution can be obtained for the two initial conditions, and for other initial values

Figure 7.15: Evolution of the pattern starting from the initial conditions suggested in Project e.
of \( A \) and \( B \) there are two steady state solutions. Try other initial conditions. If you obtain a new solution, change the initial \( A \) or \( B \) slightly to see if your new solution is stable. On an \( f \) versus \( k \) plot indicate where there are two solutions and where there are one. In this way you can determine the approximate phase diagram for this process.

b. There is a small region in \( f-k \) space where one of the steady state solutions becomes unstable and periodic solutions occur (the mechanism is known as a Hopf bifurcation). Try \( f = 0.009, \ k = 0.03 \), and set \( A = 0.1 \) and \( B = 0.5 \) at \( t = 0 \). Plot the values of \( A \) and \( B \) versus the time \( t \). Are they periodic? Try other values of \( f \) and \( k \) and estimate where the periodic solutions occur.

c. Numerical solutions of the full equation with diffusion (7.66) can be found by making a finite difference approximation to the spatial derivatives as in (3.16) and using a simple Euler algorithm for the time integration. Adopt periodic boundary conditions. Although it is straightforward to write a program to do the numerical integration, an exploration of the dynamics of this system requires much computer resources. However, we can find some preliminary results with a small system and a coarse grid. Consider a \( 0.5 \times 0.5 \) system with a spatial mesh of \( 128 \times 128 \) grid points on a square lattice. Choose \( f = 0.18, \ k = 0.057, \) and \( \Delta t = 0.1 \). Let the entire system be in the initial trivial state (\( A = 1, \ B = 0 \)) except for a \( 20 \times 20 \) grid located at the center of the system where the sites are \( A = 1/2, \ B = 1/4 \) with a \( \pm 1\% \) random noise. The effect of the noise is to break the square symmetry. Let the system evolve for approximately 80,000 time steps and look at the patterns that develop. Color code the grid according to the concentration of \( A \), with red representing \( A = 1 \) and blue representing \( A \approx 0.2 \) and with several intermediate colors. Very interesting patterns have been found by Pearson.

Appendix 7: Random Walks and the Diffusion Equation

To gain some insight into the relation between random walks and the diffusion equation, we first show that the latter implies that \( \langle x(t) \rangle \) is zero and \( \langle x^2(t) \rangle \) is proportional to \( t \). We rewrite the diffusion equation (7.26) here for convenience:

\[
\frac{\partial P(x, t)}{\partial t} = D \frac{\partial^2 P(x, t)}{\partial x^2}. \tag{7.67}
\]

To derive the \( t \) dependence of \( \langle x(t) \rangle \) and \( \langle x^2(t) \rangle \) from (7.67), we write the average of any function of \( x \) as

\[
\langle f(x, t) \rangle = \int_{-\infty}^{\infty} f(x) P(x, t) \, dx. \tag{7.68}
\]

The average displacement is given by

\[
\langle x(t) \rangle = \int_{-\infty}^{\infty} x P(x, t) \, dx. \tag{7.69}
\]

To do the integral on the right hand side of (7.69), we multiply both sides of (7.67) by \( x \) and formally integrate over \( x \):

\[
\int_{-\infty}^{\infty} x \frac{\partial P(x, t)}{\partial t} \, dx = D \int_{-\infty}^{\infty} x \frac{\partial^2 P(x, t)}{\partial x^2} \, dx. \tag{7.70}
\]
The left-hand side can be expressed as:

\[ \int_{-\infty}^{\infty} x \frac{\partial P(x,t)}{\partial t} \, dx = \frac{\partial}{\partial t} \int_{-\infty}^{\infty} xP(x,t) \, dx = \frac{d}{dt} \langle x \rangle. \tag{7.71} \]

The right-hand side of (7.70) can be written in the desired form by doing an integration by parts:

\[ D \int_{-\infty}^{\infty} x \frac{\partial^2 P(x,t)}{\partial x^2} \, dx = D \left[ x \frac{\partial P(x,t)}{\partial x} \right]_{x=-\infty}^{x=\infty} - D \int_{-\infty}^{\infty} \frac{\partial P(x,t)}{\partial x} \, dx. \tag{7.72} \]

The first term on the right-hand side of (7.72) is zero because \( P(x = \pm \infty, t) = 0 \) and all the spatial derivatives of \( P \) at \( x = \pm \infty \) are zero. The second term also is zero because it integrates to \( D[P(x = \infty, t) - P(x = -\infty, t)] \). Hence, we find that

\[ \frac{d}{dt} \langle x \rangle = 0, \tag{7.73} \]

or \( \langle x \rangle \) is a constant, independent of time. Because \( x = 0 \) at \( t = 0 \), we conclude that \( \langle x \rangle = 0 \) for all \( t \).

To calculate \( \langle x^2(t) \rangle \), we can use a similar procedure and perform two integrations by parts. The result is

\[ \frac{d}{dt} \langle x^2(t) \rangle = 2D, \tag{7.74} \]

or

\[ \langle x^2(t) \rangle = 2Dt. \tag{7.75} \]

We see that the random walk and the diffusion equation have the same time dependence. In \( d \)-dimensional space, \( 2D \) is replaced by \( 2dD \).

The solution of the diffusion equation shows that the time dependence of \( \langle x^2(t) \rangle \) is equivalent to the long time behavior of a simple random walk on a lattice. In the following, we show directly that the continuum limit of the one-dimensional random walk model is a diffusion equation.

If there is an equal probability of taking a step to the right or left, the random walk can be written in terms of the simple master equation

\[ P(i,N) = \frac{1}{2} [P(i+1,N-1) + P(i-1,N-1)], \tag{7.76} \]

where \( P(i,N) \) is the probability that the walker is at site \( i \) after \( N \) steps. To obtain a differential equation for the probability density \( P(x,t) \), we identify \( t = N\tau, \ x = ia \), and \( P(i,N) = aP(x,t) \), where \( \tau \) is the time between steps and \( a \) is the lattice spacing. This association allows us to rewrite (7.76) in the equivalent form

\[ P(x,t) = \frac{1}{2} [P(x+a,t-\tau) + P(x-a,t-\tau)]. \tag{7.77} \]

We rewrite (7.77) by subtracting \( P(x,t-\tau) \) from both sides of (7.77) and dividing by \( \tau \):

\[ \frac{1}{\tau} [P(x,t) - P(x,t-\tau)] = \frac{a^2}{2\tau} [P(x+a,t-\tau) - 2P(x,t-\tau) + P(x-a,t-\tau)] a^{-2}. \tag{7.78} \]
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If we expand \( P(x, t - \tau) \) and \( P(x \pm a, t - \tau) \) in a Taylor series and take the limit \( a \to 0 \) and \( \tau \to 0 \) with the ratio \( D \equiv a^2/2\tau \) finite, we obtain the diffusion equation

\[
\frac{\partial P(x, t)}{\partial t} = D \frac{\partial^2 P(x, t)}{\partial x^2}.
\]

(7.79a)

The generalization of (7.79a) to three dimensions is

\[
\frac{\partial P(x, y, z, t)}{\partial t} = D \nabla^2 P(x, y, z, t),
\]

(7.79b)

where \( \nabla^2 = \partial^2 / \partial x^2 + \partial^2 / \partial y^2 + \partial^2 / \partial z^2 \) is the Laplacian operator. Equation (7.79) is known as the diffusion equation and is frequently used to describe the dynamics of fluid molecules.

The direct numerical solution of the prototypical parabolic partial differential equation (7.79) is a nontrivial problem in numerical analysis (cf. Press et al. or Koonin and Meredith). An indirect method of solving (7.79) numerically is to use a Monte Carlo method, that is, replace the partial differential equation (7.79) by a corresponding random walk on a lattice with discrete time steps. Because the asymptotic behavior of the partial differential equation and the random walk model are equivalent, this approach uses the Monte Carlo technique as a method of numerical analysis. In contrast, if our goal is to understand a random walk lattice model directly, the Monte Carlo technique is a simulation method. The difference between simulation and numerical analysis is sometimes in the eyes of the beholder.

Problem 7.43. Biased random walk

Show that the form of the differential equation satisfied by \( P(x, t) \) corresponding to a random walk with a drift, that is, a walk for \( p \neq q \), is

\[
\frac{\partial P(x, t)}{\partial t} = D \nabla^2 P(x, y, z, t) - v \frac{\partial P(x, t)}{\partial x}.
\]

(7.80)

How is \( v \) related to \( p \) and \( q \)?

References and Suggestions for Further Reading


J. Bernasconi and L. Pietronero, “True self-avoiding walk in one dimension,” Phys. Rev. B 29, 5196–5198 (1984). The authors present results for the exponent \( \nu \) accurate to 1%.
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F. Reif, Fundamentals of Statistical and Thermal Physics, McGraw-Hill (1965). This popular text on statistical physics has a good discussion on random walks (Chapter 1) and diffusion (Chapter 12).


