# CMSC 420: Data Structures ${ }^{1}$ <br> Spring 2001 <br> Dave Mount 

## Lecture 1: Course Introduction and Background

## (Tuesday, Jan 30, 2001)

Algorithms and Data Structures: The study of data structures and the algorithms that manipulate them is among the most fundamental topics in computer science. Most of what computer systems spend their time doing is storing, accessing, and manipulating data in one form or another. Some examples from computer science include:

Networking: Suppose you need to multicast a message from one source node to many other machines on the network. Along what paths should the message be sent, and how can this set of paths be determined from the network's structure?
Information Retrieval: How does a search engine like Google store the contents of the Web so that the few pages that are most relevant to a given query can be extracted quickly?
Compilers: You need to store a set of variable names along with their associated types. Given an assignment between two variables we need to look them up in a symbol table, determine their types, and determine whether it is possible to cast from one type to the other (say because one is a subtype of the other).

Computer Graphics: You are designing a virtual reality system for an architectural building walk-through. Given the location of the viewer, what portions of the architectural design are visible to the viewer?

In many areas of computer science, much of the content deals with the questions of how to store, access, and manipulate the data of importance for that area. In this course we will deal with the first two tasks of storage and access at a very general level. (The last issue of manipulation is further subdivided into two areas, manipulation of numeric or floating point data, which is the subject of numerical analysis, and the manipulation of discrete data, which is the subject of discrete algorithm design.) An good understanding of data structures is fundamental to all of these areas.

What is a data structure? Whenever we deal with the representation of real world objects in a computer program we must first consider a number of issues:

Modeling: the manner in which objects in the real world are modeled as abstract mathematical entities and basic data types,

Operations: the operations that are used to store, access, and manipulate these entities and the formal meaning of these operations,
Representation: the manner in which these entities are represented concretely in a computer's memory, and

Algorithms: the algorithms that are used to perform these operations.

[^0]Note that the first two items above are essentially mathematical in nature, and deal with the "what" of a data structure, whereas the last two items involve the implementation issues and the "how" of the data structure. The first two essentially encapsulate the essence of an abstract data type (or ADT). In contrast the second two items, the concrete issues of implementation, will be the focus of this course.

For example, you are all familiar with the concept of a stack from basic programming classes. This a sequence of objects (of unspecified type). Objects can be inserted into the stack by pushing and removed from the stack by popping. The pop operation removes the last unremoved object that was pushed. Stacks may be implemented in many ways, for example using arrays or using linked lists. Which representation is the fastest? Which is the most space efficient? Which is the most flexible? What are the tradeoffs involved with the use of one representation over another? In the case of a stack, the answers are all rather mundane. However, as data structures grow in complexity and sophistication, the answers are far from obvious.

In this course we will explore a number of different data structures, study their implementations, and analyze their efficiency (both in time and space). One of our goals will be to provide you with the tools that you will need to design and implement your own data structures to solve your own specific problems in data storage and retrieval.

Course Overview: In this course we will consider many different abstract data types, and we will consider many different data structures for storing each type. Note that there will generally be many possible data structures for each abstract type, and there will not generally be a "best" one for all circumstances. It will be important for you as a designer of data structures to understand each structure well enough to know the circumstances where one data structure is to be preferred over another.

How important is the choice of a data structure? There are numerous examples from all areas of computer science where a relatively simple application of good data structure techniques resulted in massive savings in computation time and, hence, money.
Perhaps a more important aspect of this course is a sense of how to design new data structures. The data structures we will cover in this course have grown out of the standard applications of computer science. But new applications will demand the creation of new domains of objects (which we cannot foresee at this time) and this will demand the creation of new data structures. It will fall on the students of today to create these data structures of the future. We will see that there are a few important elements which are shared by all good data structures. We will also discuss how one can apply simple mathematics and common sense to quickly ascertain the weaknesses or strengths of one data structure relative to another.

Algorithmics: It is easy to see that the topics of algorithms and data structures cannot be separated since the two are inextricably intertwined. So before we begin talking about data structures, we must begin with a quick review of the basics of algorithms, and in particular, how to measure the relative efficiency of algorithms. The main issue in studying the efficiency of algorithms is the amount of resources they use, usually measured in either the space or time used. There are usually two ways of measuring these quantities. One is a mathematical analysis of the general algorithm being used, called an asymptotic analysis, which can capture gross aspects of efficiency for all possible inputs but not exact execution times. The second is an empirical analysis of an actual implementation to determine exact running times for a sample of specific inputs, but it cannot predict the performance of the algorithm on all inputs. In class we will deal mostly with the former, but the latter is important also.
There is another aspect of complexity, that we will not discuss at length (but needs to be considered) and that is the complexity of programming. Some of the data structures that we will discuss will be quite simple to implement and others much more complex. The issue
of which data structure to choose may be dependent on issues that have nothing to do with run-time issues, but instead on the software engineering issues of what data structures are most flexible, which are easiest to implement and maintain, etc. These are important issues, but we will not dwell on them excessively, since they are really outside of our scope.
For now let us concentrate on running time. (What we are saying can also be applied to space, but space is somewhat easier to deal with than time.) Given a program, its running time is not a fixed number, but rather a function. For each input (or instance of the data structure), there may be a different running time. Presumably as input size increases so does running time, so we often describe running time as a function of input/data structure size $n$, denoted $T(n)$. We want our notion of time to be largely machine-independent, so rather than measuring CPU seconds, it is more common to measure basic "steps" that the algorithm makes (e.g. the number of statements executed or the number of memory accesses). This will not exactly predict the true running time, since some compilers do a better job of optimization than others, but its will get us within a small constant factor of the true running time most of the time.
Even measuring running time as a function of input size is not really well defined, because, for example, it may be possible to sort a list that is already sorted, than it is to sort a list that is randomly permuted. For this reason, we usually talk about worst case running time. Over all possible inputs of size $n$, what is the maximum running time. It is often more reasonable to consider expected case running time where we average over all inputs of size $n$. We will usually do worst-case analysis, except where it is clear that the worst case is significantly different from the expected case.

Review of Asymptotics: There are particular bag of tricks that most algorithm analyzers use to study the running time of algorithms. For this class we will try to stick to the basics. The first element is the notion of asymptotic notation. Suppose that we have already performed an analysis of an algorithm and we have discovered through our worst-case analysis that

$$
T(n)=13 n^{3}+42 n^{2}+2 n \log n+3 \sqrt{n} .
$$

(This function was just made up as an illustration.) Unless we say otherwise, assume that logarithms are taken base 2 . When the value $n$ is small, we do not worry too much about this function since it will not be too large, but as $n$ increases in size, we will have to worry about the running time. Observe that as $n$ grows larger, the size of $n^{3}$ is much larger than $n^{2}$, which is much larger than $n \log n$ (note that $0<\log n<n$ whenever $n>1$ ) which is much larger than $\sqrt{n}$. Thus the $n^{3}$ term dominates for large $n$. Also note that the leading factor 13 is a constant. Such constant factors can be affected by the machine speed, or compiler, so we may ignore it (as long as it is relatively small). We could summarize this function succinctly by saying that the running time grows "roughly on the order of $n^{3 "}$, and this is written notationally as $T(n) \in O\left(n^{3}\right)$.
Informally, the statement $T(n) \in O\left(n^{3}\right)$ means, "when you ignore constant multiplicative factors, and consider the leading (i.e. fastest growing) term, you get $n^{3 "}$. This intuition can be made more formal, however. (It is not the most standard one, but is good enough for most uses and is the easiest one to apply.)

Definition: $T(n) \in O(f(n))$ if $\lim _{n \rightarrow \infty} T(n) / f(n)$ is either zero or a constant (but not $\infty$ ).
For example, we said that the function above $T(n) \in O\left(n^{3}\right)$. Using the definition we have

$$
\lim _{n \rightarrow \infty} \frac{T(n)}{f(n)}=\lim _{n \rightarrow \infty} \frac{13 n^{3}+42 n^{2}+2 n \log n+3 \sqrt{n}}{n^{3}}
$$

$$
\begin{aligned}
& =\lim _{n \rightarrow \infty}\left(13+\frac{42}{n}+\frac{2 \log n}{n^{2}}+\frac{3}{n^{2.5}}\right) \\
& =13
\end{aligned}
$$

Since this is a constant, we can assert that $T(n) \in O\left(n^{3}\right)$.
The $O$ notation is good for putting an upper bound on a function. Notice that if $T(n)$ is $O\left(n^{3}\right)$ it is also $O\left(n^{4}\right), O\left(n^{5}\right)$, etc. since the limit will just go to zero. We will try to avoid getting bogged down in this notation, but it is important to know the definitions. To get a feeling what various growth rates mean here is a summary.
$T(n) \in O(1):$ Great. This means your algorithm takes only constant time. You can't beat this.
$T(n) \in O(\log \log n):$ Super fast! For all intents this is as fast as a constant time.
$T(n) \in O(\log n):$ Very good. This is called logarithmic time. It is the running time of binary search and the height of a balanced binary tree. This is about the best that can be achieved for data structures based on binary trees. Note that $\log 1000 \approx 10$ and $\log 1,000,000 \approx 20$ (log's base 2).
$T(n) \in O\left((\log n)^{k}\right):($ where $k$ is a constant). This is called polylogarithmic time. Not bad, when simple logarithmic time is not achievable. We will often write this as $O\left(\log ^{k} n\right)$.
$T(n) \in O\left(n^{p}\right):($ where $0<p<1$ is a constant). An example is $O(\sqrt{n})$. This is slower than polylogarithmic (no matter how big $k$ is or how small $p$ ), but is still faster than linear time, which is acceptable for data structure use.
$T(n) \in O(n)$ : This is called linear time. It is about the best that one can hope for if your algorithm has to look at all the data. (In data structures the goal is usually to avoid this though.)
$T(n) \in O(n \log n):$ This one is famous, because this is the time needed to sort a list of numbers. It arises in a number of other problems as well.
$T(n) \in O\left(n^{2}\right): Q u a d r a t i c$ time. Okay if $n$ is in the thousands, but rough when $n$ gets into the millions.
$T(n) \in O\left(n^{k}\right):($ where $k$ is a constant). This is called polynomial time. Practical if $k$ is not too large.
$T(n) \in O\left(2^{n}\right), O\left(n^{n}\right), O(n!)$ : Exponential time. Algorithms taking this much time are only practical for the smallest values of $n$ (e.g. $n \leq 10$ or maybe $n \leq 20$ ).

## Lecture 2: Mathematical Prelimaries

(Thursday, Feb 1, 2001)
Read: Chapt 1 of Weiss and skim Chapt 2.
Mathematics: Although this course will not be a "theory" course, it is important to have a basic understanding of the mathematical tools that will be needed to reason about the data structures and algorithms we will be working with. A good understanding of mathematics helps greatly in the ability to design good data structures, since through mathematics it is possible to get a clearer understanding of the nature of the data structures, and a general feeling for their efficiency in time and space. Last time we gave a brief introduction to asymptotic (big"Oh" notation), and later this semester we will see how to apply that. Today we consider a few other preliminary notions: summations and proofs by induction.

Summations: Summations are important in the analysis of programs that operate iteratively. For example, in the following code fragment

$$
\text { for }(i=0 ; i<n ; i++)\{\ldots\}
$$

Where the loop body (the "...") takes $f(i)$ time to run the total running time is given by the summation

$$
T(n)=\sum_{i=0}^{n-1} f(i)
$$

Observe that nested loops naturally lead to nested sums. Solving summations breaks down into two basic steps. First simplify the summation as much as possible by removing constant terms (note that a constant here means anything that is independent of the loop variable, $i)$ and separating individual terms into separate summations. Then each of the remaining simplified sums can be solved. Some important sums to know are

$$
\begin{aligned}
\sum_{i=1}^{n} 1 & =n \quad \text { (The constant series) } \\
\sum_{i=1}^{n} i & =\frac{n(n+1)}{2} \quad \text { (The arithmetic series) } \\
\sum_{i=1}^{n} \frac{1}{i} & =\ln n+O(1) \quad \text { (The harmonic series) } \\
\sum_{i=0}^{n} c^{i} & =\frac{c^{n+1}-1}{c-1} \quad c \neq 1 \quad \text { (The geometric series) }
\end{aligned}
$$

Note that complex sums can often be broken down into simpler terms, which can then be solved. For example

$$
\begin{aligned}
T(n) & =\sum_{i=n}^{2 n-1}(3+4 i)=\sum_{i=0}^{2 n-1}(3+4 i)-\sum_{i=0}^{n-1}(3+4 i) \\
& =\left(3 \sum_{i=0}^{2 n-1} 1+4 \sum_{i=0}^{2 n-1} i\right)-\left(3 \sum_{i=0}^{n-1} 1+4 \sum_{i=0}^{n-1} i\right) \\
& =\left(3(2 n)+4 \frac{2 n(2 n-1)}{2}\right)-\left(3(n)+4 \frac{n(n-1)}{2}\right)=\left(n+6 n^{2}\right)
\end{aligned}
$$

The last summation is probably the most important one for data structures. For example, suppose you want to know how many nodes are in a complete 3 -ary tree of height $h$. (We have not given a formal definition of tree's yet, but consider the figure below.
The height of a tree is the maximum number of edges from the root to a leaf.) One way to break this computation down is to look at the tree level by level. At the top level (level 0) there is 1 node, at level 1 there are 3 nodes, at level 2,9 nodes, and in general at level $i$ there $3^{i}$ nodes. To find the total number of nodes we sum over all levels, 0 through $h$. Plugging into the above equation with $h=n$ we have:

$$
\sum_{i=0}^{h} 3^{i}=\frac{3^{h+1}-1}{2} \in O\left(3^{h}\right)
$$



Figure 1: Complete 3-ary tree of height 2.

Conversely, if someone told you that he had a 3 -ary tree with $n$ nodes, you could determine the height by inverting this. Since $n=\left(3^{(h+1)}-1\right) / 2$ then we have

$$
3^{(h+1)}=(2 n+1)
$$

implying that

$$
h=\left(\log _{3}(2 n+1)\right)-1 \in O(\log n) .
$$

Another important fact to keep in mind about summations is that they can be approximated using integrals.

$$
\sum_{i=a}^{b} f(i) \approx \int_{x=a}^{b} f(x) d x
$$

Given an obscure summation, it is often possible to find it in a book on integrals, and use the formula to approximate the sum.

Recurrences: A second mathematical construct that arises when studying recursive programs (as are many described in this class) is that of a recurrence. A recurrence is a mathematical formula that is defined recursively. For example, let's go back to our example of a 3 -ary tree of height $h$. There is another way to describe the number of nodes in a complete 3-ary tree. If $h=0$ then the tree consists of a single node. Otherwise that the tree consists of a root node and 3 copies of a 3 -ary tree of height $h-1$. This suggests the following recurrence which defines the number of nodes $N(h)$ in a 3-ary tree of height $h$ :

$$
\begin{aligned}
& N(0)=1 \\
& N(h)=3 N(h-1)+1 \quad \text { if } h \geq 1
\end{aligned}
$$

Although the definition appears circular, it is well grounded since we eventually reduce to $N(0)$.

$$
\begin{aligned}
& N(1)=3 N(0)+1=3 \cdot 1+1=4 \\
& N(2)=3 N(1)+1=3 \cdot 4+1=13 \\
& N(3)=3 N(2)+1=3 \cdot 13+1=40
\end{aligned}
$$

and so on.
There are two common methods for solving recurrences. One (which works well for simple regular recurrences) is to repeatedly expand the recurrence definition, eventually reducing it to a summation, and the other is to just guess an answer and use induction. Here is an example of the former technique.

$$
\begin{aligned}
N(h) & =3 N(h-1)+1 \\
& =3(3 N(h-2)+1)+1=9 N(h-2)+3+1
\end{aligned}
$$

$$
\begin{aligned}
& =9(3 N(h-3)+1)+3+1=27 N(h-3)+9+3+1 \\
& \vdots \\
& =3^{k} N(h-k)+\left(3^{k-1}+\ldots+9+3+1\right)
\end{aligned}
$$

When does this all end? We know that $N(0)=1$, so let's set $k=h$ implying that

$$
N(h)=3^{h} N(0)+\left(3^{h-1}+\ldots+3+1\right)=3^{h}+3^{h-1}+\ldots+3+1=\sum_{i=0}^{h} 3^{i}
$$

This is the same thing we saw before, just derived in a different way.
Proofs by Induction: The last mathematical technique of importance is that of proofs by induction. Induction proofs are critical to all aspects of computer science and data structures, not just efficiency proofs. In particular, virtually all correctness arguments are based on induction. From courses on discrete mathematics you have probably learned about the standard approach to induction. You have some theorem that you want to prove that is of the form, "For all integers $n \geq 1$, blah, blah, blah", where the statement of the theorem involves $n$ in some way. The idea is to prove the theorem for some basis set of $n$-values (e.g. $n=1$ in this case), and then show that if the theorem holds when you plug in a specific value $n-1$ into the theorem then it holds when you plug in $n$ itself. (You may be more familiar with going from $n$ to $n+1$ but obviously the two are equivalent.)
In data structures, and especially when dealing with trees, this type of induction is not particularly helpful. Instead a slight variant called strong induction seems to be more relevant. The idea is to assume that if the theorem holds for all values of $n$ that are strictly less than $n$ then it is true for $n$. As the semester goes on we will see examples of strong induction proofs.
Let's go back to our previous example problem. Suppose we want to prove the following theorem.

Theorem: Let $T$ be a complete 3-ary tree with $n \geq 1$ nodes. Let $H(n)$ denote the height of this tree. Then

$$
H(n)=\left(\log _{3}(2 n+1)\right)-1
$$

Basis Case: (Take the smallest legal value of $n, n=1$ in this case.) A tree with a single node has height 0 , so $H(1)=0$. Plugging $n=1$ into the formula gives $\left(\log _{3}(2 \cdot 1+1)\right)-1$ which is equal to $\left(\log _{3} 3\right)-1$ or 0 , as desired.
Induction Step: We want to prove the theorem for the specific value $n>1$. Note that we cannot apply standard induction here, because there is no complete 3 -ary tree with 2 nodes in it (the next larger one has 4 nodes).
We will assume the induction hypothesis, that for all smaller $n^{\prime}, 1 \leq n^{\prime}<n, H\left(n^{\prime}\right)$ is given by the formula above. (This is sometimes called strong induction, and it is good to learn since most induction proofs in data structures work this way.)
Let's consider a complete 3 -ary tree with $n>1$ nodes. Since $n>1$, it must consist of a root node plus 3 identical subtrees, each being a complete 3 -ary tree of $n^{\prime}<n$ nodes. How many nodes are in these subtrees? Since they are identical, if we exclude the root node, each subtree has one third of the remaining number nodes, so $n^{\prime}=(n-1) / 3$. Since $n^{\prime}<n$ we can apply the induction hypothesis. This tells us that

$$
H\left(n^{\prime}\right)=\left(\log _{3}\left(2 n^{\prime}+1\right)\right)-1=\left(\log _{3}(2(n-1) / 3+1)\right)-1
$$

$$
\begin{aligned}
& =\left(\log _{3}(2(n-1)+3) / 3\right)-1=\left(\log _{3}(2 n+1) / 3\right)-1 \\
& =\log _{3}(2 n+1)-\log _{3} 3-1=\log _{3}(2 n+1)-2
\end{aligned}
$$

Note that the height of the entire tree is one more than the heights of the subtrees so $H(n)=H\left(n^{\prime}\right)+1$. Thus we have:

$$
H(n)=\log _{3}(2 n+1)-2+1=\log _{3}(2 n+1)-1,
$$

as desired.
This may seem like an awfully long-winded way of proving such a simple fact. But induction is a very powerful technique for proving many more complex facts that arise in data structure analyses.
You need to be careful when attempting proofs by induction that involve $O(n)$ notation. Here is an example of a common error.

False! Theorem: For $n \geq 1$, let $T(n)$ be given by the following summation

$$
T(n)=\sum_{i=0}^{n} i
$$

then $T(n) \in O(n)$. (We know from the formula for the linear series above that $T(n)=$ $n(n+1) / 2 \in O\left(n^{2}\right)$. So this must be false. Can you spot the error in the following "proof"?)
Basis Case: For $n=1$ we have $T(1)=1$, and 1 is $O(1)$.
Induction Step: We want to prove the theorem for the specific value $n>1$. Suppose that for any $n^{\prime}<n, T\left(n^{\prime}\right) \in O\left(n^{\prime}\right)$. Now, for the case $n$, we have (by definition)

$$
T(n)=\sum_{i=0}^{n} i=\left(\sum_{i=0}^{n-1} i\right)+n=T(n-1)+n
$$

Now, since $n-1<n$, we can apply the induction hypothesis, giving $T(n-1) \in O(n-1)$. Plugging this back in we have

$$
T(n) \in O(n-1)+n
$$

But $(n-1)+n \leq 2 n-1 \in O(n)$, so we have $T(n) \in O(n)$.
What is the error? Recall asymptotic notation applies to arbitrarily large $n$ (for $n$ in the limit). However induction proofs by their very nature only apply to specific values of $n$. The proper way to prove this by induction would be to come up with a concrete expression, which does not involve $O$-notation. For example, try to prove that for all $n \geq 1, T(n) \leq 50 n$. If you attempt to prove this by induction (try it!) you will see that it fails.

## Lecture 3: A Quick Introduction to Java

(Tuesday, Feb 6, 2001)
Read: Chapt 1 of Weiss (from 1.4 to the end).
Disclaimer: I am still learning Java, and there is a lot that I do not know. If you spot any errors in this lecture, please bring them to my attention.

Overview: This lecture will present a very quick description of the Java language (at least the portions that will be needed for the class projects. The Java programming language has a number of nice features, which make it an excellent programming language for implementing data structures. Based on my rather limited experience with Java, I find it to be somewhat cleaner and easier than C++ for data structure implementation (because of its cleaner support for strings, exceptions and threads, its automatic memory management, and its more integrated view of classes). Nonetheless, there are instances where C ++ would be better choice than Java for implementation (because of its template capability, its better control over memory layout and memory management).

Java API: The Java language is very small and simple compared to C++. However, it comes with a very large and constantly growing library of utility classes. Fortunately, you only need to know about the parts of this library that you need, and the documentation is all online. The libraries are grouped into packages, which altogether are called the "Java 2 Platform API" (application programming interface). Some examples of these packages include
java.lang: Contains things like strings, that are essentially built in to the language.
java.io: Contains support for input and output, and
java.util: Contains some handy data structures such as lists and hash tables.
Simple Types: The basic elements of Java and C++ are quite similar. Both languages use the same basic lexical structure and have the same methods for writing comments. Java supports the following simple types, from which more complex types can be constructed.

Integer types: These include byte (8-bit), short (16-bit), int (32-bit), and long (64-bit). All are signed integers.
Floating point types: These include float (32-bit) and double (64-bit).
Characters: The char type is a 16 -bit unicode character (unlike $\mathrm{C}++$ where it behaves essentially like a byte.)
Booleans: The boolean type is the analogue to C++'s bool, either true or false.
Variables, Literals, and Casting: Variable declarations, literals (constants) and casting work much like they do in $\mathrm{C}++$.

```
int i = 10 * 4; // integer assignment
byte j = (byte) i; // example of a cast
long l = i; // cast not needed when "widening"
float pi1 = 3.14f; // floating assignment
double pi2 = 3.141593653;
char c = 'a'; // character
boolean b = (i < pi1); // boolean
```

All the familiar arithmetic operators are just as they are in $\mathrm{C}++$.
Arrays: In $\mathrm{C}++$ an array is just a pointer to its first element. There is no index range checking at run time. In Java, arrays have definite lengths and run-time index checking is done. They are indexed starting from 0 just as in $\mathrm{C}++$. When an array is declared, no storage is allocated. Instead, storage is allocated by using the new operator, as shown below.

```
int A[]; // A is an array, initially null
A = new int[4] = {12, 4, -6, 41}; // this allocates and defines array
char B[][] = new char[3] [8]; // B is a 2-dim, 3 x 8 array of char
B[0] [3] = 'c';
```

In Java you do not need to delete things created by new. It is up to the system to do garbage collection to get rid of objects that you no longer refer to. In Java a 2-dimension array is actually an array of arrays.

String Object: In C ++ a string is implemented as an array of characters, terminated by a null character. Java has much better integrated support for strings, in terms of a special String object. For example, each string knows its length. Strings can be concatenated using the "+" operator (which is the only operator overloading supported in Java).

```
String S = "abc"; // S is the string "abc"
String T = S + "def"; // T is the string "abcdef"
char c = T[3]; // c = 'd'
S = S + T; // now S = "abcabcdef"
int k = T.length(); // k = 6
```

Note that strings can be appended to by concatenation (as shown above with $S$ ). If you plan on making many changes to a string variable, there is a type StringBuffer which provides a more efficient method for "growable" strings.

Java supports automatic casting of simple types into strings, which comes in handy for doing output. For example,

```
int cost = 25;
String Q = "The value is " + cost + " dollars";
```

Would set Q to "The value is 25 dollars". If you define your own object, you can define a method toString(), which produces a String representation for your object. Then any attempt to cast your object to a string will invoke your method.
There are a number of other string functions that are supported. One thing to watch out for is string comparison. The standard == operator will not do what you might expect. The reason is that a String object (like all objects in Java) behaves much like a pointer or reference does in C++. Comparisons using == test whether the pointers are the same, not whether the contents are the same. The equals() or compareTo() methods to test the string contents.

Flow Control: All the standard flow-control statements, if-then-else, for, while, switch, continue, etc., are the same as in $\mathrm{C}++$.

Standard I/O: Input and output from the standard input/output streams are not quite as simple in Java as in C++. Output is done by the print() and println() methods, each of which prints a string to the standard output. The latter prints an end-of-line character. In Java there are no global variables or functions. Everything is referenced as a member of some object. These functions exist as part of the System. out object. Here is an example of its usage.

```
int age = 12;
System.out.println("He is " + age + " years old");
```

Input in Java is much messier. Since our time is limited, we'll refer you to our textbook (by Weiss) for an example of how to use the java.io.BufferedReader, java.io.InputStreamReader, StringTokenizer to input lines and read them into numeric, string, and character variables.

Classes and Objects: Virtually everything in Java is organized around the notion of hierarchy of objects and classes. As in C++ an object is an entity which stores data and provides methods for accessing and modifying the data. An object is an instance of a class, which defines the specifics of the object. As in $\mathrm{C}++$ a class consists of instance variables and methods. Although classes are superficially similar to classes in $\mathrm{C}++$, there are some important differences in how they operate.

File Structure: Java differs from C++ in that there is a definite relationship between files and classes. Normally each Java class resides within its own file, whose name is derived from the name of the class. For example, a class Student would be stored in file Student.java. Unlike C++, in which class declarations are stored in different files (.h and .cc), in Java everything is stored in the same file and all definitions are made within the class.

Packages: Another difference with $\mathrm{C}++$ is the notion of packages. A package is a collection of related classes (e.g. the various classes that make up one data structure). The classes that make up a package are stored in a common directory (with the same name as the package). In C++ you access predefined objects using an \#include command to load the definitions. In Java this is done by importing the desired package. For example, if you wanted to use the Stack data structure in the java.util library, you would load this using the following.

```
import java.util.Stack;
```

The Dot (.) Operator: Class members are accessed as in C++ using the dot operator. However in Java, objects all behave as if they were pointers (or perhaps more accurately as references that are allowed to be null).

```
class Person {
    String name;
    int age;
    public Person() {...} // constructor (details omitted)
}
Person p; // p has the value null
p = new Person(); // p points to a new object
Person q = p; // q and p now point to the same object
p = null; // q still points to the object
```

Note that unlike C++ the declaration "Person p" did not allocate an object as it would in $\mathrm{C}++$. It creates a null pointer to a Person. Objects are only created using the "new" operator. Also, unlike $\mathrm{C}++$ the statement " $\mathrm{q}=\mathrm{p}$ " does not copy the contents of $p$, but instead merely assigns the pointer $q$ to point to the same object as $p$. Note that the pointer and reference operators of $\mathrm{C}++(" * ", " \& ", "->")$ do not exist in Java (since everything is a pointer).
While were at it, also note that there is no semicolon (;) at the end of the class definition. This is one small difference between Java and C++.
Inheritance and Subclasses: As with C++, Java supports class inheritance (but only single not multiple inheritance). A class which is derived from another class is said to be a child or subclass, and it is said to extend the original parent class, or superclass.

```
class Person { ... }
class Student extends Person { ... }
```

The Student class inherits the structure of the Person class and can add new instance variables and new methods. If it redefines a method, then this new method overrides the
old method. Methods behave like virtual member functions in $\mathrm{C}++$, in the sense that each object "remembers" how it was created, even if it assigned to a superclass.

```
class Person {
    public void print() { System.out.println("Person"); }
}
class Student extends Person {
    public void print() { System.out.println("Student"); }
}
Person P = new Person();
Student S = new Student();
Person Q = S;
Q.print();
```

What does the last line print? It prints "Student", because $Q$ is assigned to an object that was created as a Student object. This called dynamic dispatch, because the system decides at run time which function is to be called.

Variable and Method Modifiers: Getting back to classes, instance variables and methods behave much as they do in $\mathrm{C}++$. The following variable modifiers and method modifiers determine various characteristics of variables and methods.
public: Any class can access this item.
protected: Methods in the same package or subclasses can access this item.
private: Only this class can access this item.
static: There is only one copy of this item, which is shared between all instances of the class.
final: This item cannot be overridden by subclasses. (This is allows the compiler to perform optimizations.)
abstract: (Only for methods.) This is like a pure virtual method in C++. This method is not defined here (no code is provided), but rather must be defined by a subclass.
Note that Java does not have const variables like $\mathrm{C}++$, nor does it have enumerations. A constant is defined as a static final variable. (Enumerations can be faked by defining many constants.)

```
class Person {
    static final int TALL = 0;
    static final int MEDIUM = 1;
    static final int SHORT = 2;
}
```

Class Modifiers: Unlike C++, classes can also be associated with modifiers.
(Default:) If no modifier is given then the class is friendly, which means that it can be instantiated only by classes in the same package.
public: This class can be instantiated or extended by anything in the same package or anything that imports this class.
final: This class can have no subclasses.
abstract: This class has some abstract methods.
"this" and "super": As in C++, this is a pointer to the current object. In addition, Java defines super to be a pointer to the object's superclass. The super pointer is nice because it provides a method for explicitly accessing overridden methods in the parent's class and invoking the parent's constructor.

```
class Person {
    int age;
    public Person(int a) { age = a; } // Person constructor
}
class Student extends Person {
    String major;
    public Student(int a, String m) { // Student constructor
        super(a);
        major = m;
    }
}
```

(By the way, this is a dangerous way to initialize a String object, since recall that assignment does not copy contents, it just copies the pointer.)

Functions, Parameters, and Wrappers: Unlike C++ all functions must be part of some class. For example, the main procedure must itself be declared as part of some class, as a public, static, void method. The command line arguments are passed in as an array of strings. Since arrays can determine their own length from the length () method, there is no need to indicate the number of arguments. For example, here is what a typical Hello-World program would look like. This is stored in a file HelloWorld.java

```
class HelloWorld {
    public static void main(String args[]) {
        System.out.println("Hello World");
        System.out.println("The number of arguments is " + args.length());
    }
}
```

In Java simple types are passed by value and objects are passed by reference. Thus, modifying an object which has been passed in through a parameter of the function has the effect of modifying the original object in the calling routine.
What if you want to pass an integer to a function and have its value modified. Java provides a simple trick to do this. You simply wrap the integer with a class, and now it is passed by reference. Java efines wrapper classes Integer (for int), Long, Character (for char), Boolean, etc.

Generic Data Structures without Templates: Java does not have many of the things that C++ has, templates for example. So, how do you define a generic data structure (e.g. a Stack) that can contain objects of various types (int, float, Person, etc.)? In Java every class automatically extends a most generic class called Object. We set up our generic stack to hold objects of type Object. Since this is a superclass of all classes, we can store any objects from any class in our stack. Of course, to keep our sanity, we should store object of only one type in any given stack. (Since it will be very difficulty to figure out what comes out whenever we pop something off the stack.)
Consider the ArrayVector class defined below. Suppose we want to store integers in this vector. An int is not an object, but we can use the wrapper class Integer to convert an integer into an object.

```
ArrayVector A = new ArrayVector();
Integer x = (Integer) A.elemAtRank(0); // accesses 999
```

A.insertAtRank(0, Integer(999)); // insert 999 at position 0
A.insertAtRank(1, Integer(888)); // insert 888 at position 1

Note that the ArrayVector contains Objects. Thus, when we access something from the data structure, we need to cast it back to the appropriate type.

```
// This is an example of a simple Java class, which implements a vector
// object using an array implementation. This would be stored in a file
// named "ArrayVector.java".
public class ArrayVector {
    private Object[] a; // Array storing the elements of the vector
    private int capacity = 16; // Length of array a
    private int size = 0; // Number of elements stored in the vector
    // Constructor
    public ArrayVector() { a = new Object[capacity]; }
    // Accessor methods
    public Object elemAtRank(int r) { return a[r]; }
    public int size() { return size; }
    public boolean isEmpty() { return size() == 0; }
    // Modifier methods
    public Object replaceAtRank(int r, Object e) {
        Object temp = a[r];
        a[r] = e;
        return temp;
    }
    public Object removeAtRank(int r) {
        Object temp = a[r];
        for (int i=r; i < size-1; i++) // Shift elements down
            a[i] = a[i+1];
        size--;
        return temp;
    }
    public void insertAtRank(int r, Object e) {
        if (size == capacity) { // An overflow
            capacity *= 2;
            Object[] b = new Object[capacity];
            for (int i=0; i < size; i++)
                b[i] = a[i];
            a = b;
        }
        for (int i=size-1; i>=r; i--) // Shift elements up
            a[i+1] = a[i];
        a[r] = e;
        size++;
    }
}
```


## Lecture 4: Basic Data Structures

(Thursday, Feb 8, 2001)

Read: Chapt. 3 in Weiss.
Basic Data Structures: Before we go into our coverage of complex data structures, it is good to remember that in many applications, simple data structures are sufficient. This is true, for example, if the number of data objects is small enough that efficiency is not so much an issue, and hence a complex data structure is not called for. In many instances where you need a data structure for the purposes of prototyping an application, these simple data structures are quick and easy to implement.

Abstract Data Types: An important element to good data structure design is to distinguish between the functional definition of a data structure and its implementation. By an abstract data structure (ADT) we mean a set of objects and a set of operations defined on these objects. For example, a stack ADT is a structure which supports operations such as push and pop (whose definition you are no doubt familiar with). A stack may be implemented in a number of ways, for example using an array or using a linked list. An important aspect of object oriented languages like C++ and Java is the capability to present the user of a data structure with an abstract definition of its function without revealing the methods with which it operates. To a large extent, this course will be concerned with the various options for implementing simple abstract data types and the tradeoffs between these options.

Linear Lists: A linear list or simply list is perhaps the most basic abstract data types. A list is simply an ordered sequence of elements $\left\langle a_{1}, a_{2}, \ldots, a_{n}\right\rangle$. We will not specify the actual type of these elements here, since it is not relevant to our presentation. (In $\mathrm{C}++$ this might be done through the use of templates. In Java this can be handled by defining the elements to be of type Object. Since the Object class is a superclass of all other classes, we can store any class type in our list.)
The size or length of such a list is $n$. There is no agreed upon specification of the list ADT, but typical operations include:
get(i): Returns element $a_{i}$.
set (i,x): Sets the $i$ th element to $x$.
length(): Returns the length of the list.
insert (i, x ): Insert element $x$ just prior to element $a_{i}$ (causing the index of all subsequent items to be increased by one).
delete(i): Delete the $i$ th element (causing the indices of all subsequent elements to be decreased by 1).

I am sure that you can imagine many other useful operations, for example search the list for an item, split or concatenate lists, return a sublist, make the list empty. There are often a number of programming language related operations, such as returning an iterator object for the list. Lists of various types are among the most primitive abstract data types, and because these are taught in virtually all basic programming classes we will not cover them in detail here.

There are a number of possible implementations of lists. The most basic question is whether to use sequential allocation (meaning storing the elements sequentially in an array) or linked allocation (meaning storing the elements in a linked list. With linked allocation there are many other options to be considered. Is the list singly linked, doubly linked, circularly linked? Another question is whether we have an internal list, in which the nodes that constitute the linked list contain the actual data items $a_{i}$, or we have an external list, in which each linked list item contains a pointer to the associated data item.

| $a_{1}$ |
| :--- |
| $a_{2}$ |
| $a_{3}$ |
| $a_{4}$ |

Sequential allocation


Internal linked list


External linked list

Figure 2: Common types of list allocation.

Multilists and Sparse Matrices: Although lists are very basic structures, they can be combined in various nontrivial ways. One such example is the notion of a multilist, which can be thought of as two sets of linked lists that are interleaved with each other. One application of multilists is in representing sparse matrices. Suppose that you want to represent a very large $m \times n$ matrix. Such a matrix can store $O(m n)$ entries. But in many applications the number of nonzero entries is much smaller, say on the order of $O(m+n)$. (For example, if $n=m=10,000$ this might means that only around $0.01 \%$ of the entries are being used.) A common approach for representing such a sparse matrix is by creating $m$ linked lists, one for each row and $n$ linked lists, one for each column. Each linked list stores the nonzero entries of the matrix. Each entry contains the row and column indices $[i][j]$ as well as the value stored in this entry.


Figure 3: Sparse matrix representation using multilists.
Stacks, Queues, and Deques: There are two very special types of lists: stacks and queues and their generalization, called the deque.

Stack: Supports insertions (called pushes) and deletions (pops) from only one end (called the top). Stacks are used often in processing tree-structured objects, in compilers (in processing nested structures), and is used in systems to implement recursion.
Queue: Supports insertions (called enqueues) at one end (called the tail or rear) and deletions (called dequeues) from the other end (called the head or front). Queues are used in operating systems and networking to store a list of items that are waiting for some resource.
Deque: This is a play on words. It is written like "d-e-que" for a "double-ended queue", but it is pronounced like deck, because it behaves like a deck of cards, where you can deal off the top or the bottom. A deque supports insertions and deletions from either end. Clearly, given a deque, you immediately have an implementation of a stack or queue by simple inheritance.

Both stacks and queues can be implemented efficiently as arrays or as linked lists. There are a number of interesting issues involving these data structures. However, since you have probably seen these already, we will skip the details here.

Graphs: Intuitively, a graph is a collection of vertices or nodes, connected by a collection of edges. Graphs are extremely important because they are a very flexible mathematical model for many application problems. Basically, any time you have a set of objects, and there is some "connection" or "relationship" or "interaction" between pairs of objects, a graph is a good way to model this. Examples of graphs in application include communication and transportation networks, VLSI and other sorts of logic circuits, surface meshes used for shape description in computer-aided design and geographic information systems, precedence constraints in scheduling systems. The list of application is almost too long to even consider enumerating it.

Definition: A directed graph (or digraph) $G=(V, E)$ consists of a finite set $V$, called the vertices or nodes, and $E$, a set of ordered pairs, called the edges of $G$. (Another way of saying this is that $E$ is a binary relation on $V$.)

Observe that self-loops are allowed by this definition. Some definitions of graphs disallow this. Multiple edges are not permitted (although the edges $(v, w)$ and ( $w, v$ ) are distinct).


Digraph


Graph

Figure 4: Digraph and graph example.

Definition: An undirected graph (or graph) $G=(V, E)$ consists of a finite set $V$ of vertices, and a set $E$ of unordered pairs of distinct vertices, called the edges. (Note that self-loops are not allowed).

Note that directed graphs and undirected graphs are different (but similar) objects mathematically. We say that vertex $v$ is adjacent to vertex $u$ if there is an edge $(u, v)$. In a directed graph, given the edge $e=(u, v)$, we say that $u$ is the origin of $e$ and $v$ is the destination of $e$. In undirected graphs $u$ and $v$ are the endpoints of the edge. The edge $e$ is incident (meaning that it touches) both $u$ and $v$.
In a digraph, the number of edges coming out of a vertex is called the out-degree of that vertex, and the number of edges coming in is called the in-degree. In an undirected graph we just talk about the degree of a vertex as the number of incident edges. By the degree of a graph, we usually mean the maximum degree of its vertices.
When discussing the size of a graph, we typically consider both the number of vertices and the number of edges. The number of vertices is typically written as $n$ or $V$, and the number of edges is written as $m$ or $E$. Here are some basic combinatorial facts about graphs and digraphs. We will leave the proofs to you. Given a graph with $V$ vertices and $E$ edges then:

## In a graph:

- $0 \leq E \leq\binom{ n}{2}=n(n-1) / 2 \in O\left(n^{2}\right)$.
- $\sum_{v \in V} \operatorname{deg}(v)=2 E$.


## In a digraph:

- $0 \leq E \leq n^{2}$.
- $\sum_{v \in V}$ in- $\operatorname{deg}(v)=\sum_{v \in V}$ out-deg $(v)=E$.

Notice that generally the number of edges in a graph may be as large as quadratic in the number of vertices. However, the large graphs that arise in practice typically have much fewer edges. A graph is said to be sparse if $E \in O(V)$, and dense, otherwise. When giving the running times of algorithms, we will usually express it as a function of both $V$ and $E$, so that the performance on sparse and dense graphs will be apparent.

Adjacency Matrix: An $n \times n$ matrix defined for $1 \leq v, w \leq n$.

$$
A[v, w]= \begin{cases}1 & \text { if }(v, w) \in E \\ 0 & \text { otherwise }\end{cases}
$$

If the digraph has weights we can store the weights in the matrix. For example if $(v, w) \in$ $E$ then $A[v, w]=W(v, w)$ (the weight on edge $(v, w)$ ). If $(v, w) \notin E$ then generally $W(v, w)$ need not be defined, but often we set it to some "special" value, e.g. $A(v, w)=$ -1 , or $\infty$. (By $\infty$ we mean (in practice) some number which is larger than any allowable weight. In practice, this might be some machine dependent constant like MAXINT.)
Adjacency List: An array $\operatorname{Adj}[1 \ldots n]$ of pointers where for $1 \leq v \leq n, \operatorname{Adj}[v]$ points to a linked list containing the vertices which are adjacent to $v$ (i.e. the vertices that can be reached from $v$ by a single edge). If the edges have weights then these weights may also be stored in the linked list elements.


Figure 5: Adjacency matrix and adjacency list for digraphs.
We can represent undirected graphs using exactly the same representation, but we will store each edge twice. In particular, we representing the undirected edge $\{v, w\}$ by the two oppositely directed edges $(v, w)$ and $(w, v)$. Notice that even though we represent undirected graphs in the same way that we represent digraphs, it is important to remember that these two classes of objects are mathematically distinct from one another.
This can cause some complications. For example, suppose you write an algorithm that operates by marking edges of a graph. You need to be careful when you mark edge $(v, w)$ in the representation that you also mark $(w, v)$, since they are both the same edge in reality. When dealing with adjacency lists, it may not be convenient to walk down the entire linked list, so it is common to include cross links between corresponding edges.
An adjacency matrix requires $O\left(V^{2}\right)$ storage and an adjacency list requires $O(V+E)$ storage. The $V$ arises because there is one entry for each vertex in $A d j$. Since each list has out- $\operatorname{deg}(v)$ entries, when this is summed over all vertices, the total number of adjacency list records is $O(E)$. For sparse graphs the adjacency list representation is more space efficient.



Adjacency matrix


Adjacency list (with crosslinks)

Figure 6: Adjacency matrix and adjacency list for graphs.

## Lecture 5: Trees

(Tuesday, Feb 13, 2001)
Read: Chapt. 4 in Weiss.
Trees: Trees and their variants are among the most common data structures. In its most general form, a free tree is a connected, undirected graph that has no cycles. Since we will want to use our trees for applications in searching, it will be more meaningful to assign some sense of order and direction to our trees.

Formally a tree (actually a rooted tree) is defined recursively as follows. It consists of one or more items called nodes. (Our textbook allows for the possibility of an empty tree with no nodes.) It consists of a distinguished node called the root, and a set of zero or more nonempty subsets of nodes, denoted $T_{1}, T_{2}, \ldots, T_{k}$, where each is itself a tree. These are called the subtrees of the root.
The root of each subtree $T_{1}, \ldots, T_{k}$ is said to be a child of $r$, and $r$ is the parent of each root. The children of $r$ are said to be siblings of one another. Trees are typically drawn like graphs, where there is an edge (sometimes directed) from a node to each of its children. See the figure below.


Figure 7: Trees.
If there is an order among the $T_{i}$ 's, then we say that the tree is an ordered tree. The degree of a node in a tree is the number of children it has. A leaf is a node of degree 0 . A path between two nodes is a sequence of nodes $u_{1}, u_{2}, \ldots, u_{k}$ such that $u_{i}$ is a parent of $u_{i+1}$. The length of a path is the number of edges on the path (in this case $k-1$ ). There is a path of length 0 from every node to itself.

The depth of a node in the tree is the length of the unique path from the root to that node. The root is at depth 0 . The height of a node is the length of the longest path from the node to a leaf. Thus all leaves are at height 0 . If there is a path from $u$ to $v$ we say that $v$ is a descendants of $u$. We say it is a proper descendent if $u \neq v$. Similarly, $u$ is an ancestor of $v$.

Implementation of Trees: One difficulty with representing general trees is that since there is no bound on the number of children a node can have, there is no obvious bound on the size of a given node (assuming each node must store pointers to all its children). The more common representation of general trees is to store two pointers with each node: the firstChild and the nextSibling. The figure below illustrates how the above tree would be represented using this technique.


Figure 8: A binary representation of general trees.
Trees arise in many applications in which hierarchies exist. Examples include the Unix file system, corporate managerial structures, and anything that can be described in "outline form" (like the chapters, sections, and subsections of a user's manual). One special case of trees will be very important for our purposes, and that is the notion of a binary tree.

Binary Trees: Our text defines a binary tree as a tree in which each node has no more than two children. However, this definition is subtly flawed. A binary tree is defined recursively as follows. A binary tree can be empty. Otherwise, a binary tree consists of a root node and two disjoint binary trees, called the left and right subtrees. The difference in the two definitions is important. There is a distinction between a tree with a single left child, and one with a single right child (whereas in our normal definition of tree we would not make any distinction between the two).
The typical Java representation of a tree as a data structure is given below. The element field contains the data for the node and is of some abstract type, which in Java might be Object. (In C++ the element type could be specified using a template.) When we need to be concrete we will often assume that element fields are just integers. The left field is a pointer to the left child (or null if this tree is empty) and the right field is analogous for the right child.

```
class BinaryTreeNode {
    Object element; // data item
    BinaryTreeNode left; // left child
    BinaryTreeNode right; // right child
}
```

Binary trees come up in many applications. One that we will see a lot of this semester is for representing ordered sets of objects, a binary search tree. Another one that is used often in compiler design is expression trees which are used as an intermediate representation for expressions when a compiler is parsing a statement of some programming language. For
example, in the figure below right, we show an expression tree for the expression $((a+b) *$ $c) /(d-e))$.


Expression Tree


Binary tree and extended binary tree

Figure 9: Expression tree.

Traversals: There are three natural ways of visiting or traversing every node of a tree, preorder, postorder, and (for binary trees) inorder. Let $T$ be a tree whose root is $r$ and whose subtrees are $T_{1}, T_{2}, \ldots, T_{m}$ for $m \geq 0$.

Preorder: Visit the root $r$, then recursively do a preorder traversal of $T_{1}, T_{2}, \ldots, T_{k}$. For example: $\langle/, *,+, a, b, c,-, d, e\rangle$ for the expression tree shown above.
Postorder: Recursively do a postorder traversal of $T_{1}, T_{2}, \ldots, T_{k}$ and then visit $r$. Example: $\langle a, b,+, c, *, d, e,-, /\rangle$. (Note that this is not the same as reversing the preorder traversal.)
Inorder: (for binary trees) Do an inorder traversal of $T_{L}$, visit $r$, do an inorder traversal of $T_{R}$. Example: $\langle a,+, b, *, c, /, d,-, e\rangle$.

Note that theses traversal correspond to the familiar prefix, postfix, and infix notations for arithmetic expressions.

Preorder arises in game-tree applications in AI, where one is searching a tree of possible strategies by depth-first search. Postorder arises naturally in code generation in compilers. Inorder arises naturally in binary search trees which we will see more of.
These traversals are most easily coded using recursion. If recursion is not desired (either for greater efficiency or for fear of using excessive system stack space) it is possible to use your own stack to implement the traversal. Either way the algorithm is quite efficient in that its running time is proportional to the size of the tree. That is, if the tree has $n$ nodes then the running time of these traversal algorithms are all $O(n)$.

Extended Binary Trees: Binary trees are often used in search applications, where the tree is searched by starting at the root and then proceeding either to the left or right child, depending on some condition. In some instances the search stops at a node in the tree. However, in some cases it attempts to cross a null link to a nonexistent child. The search is said to "fall out" of the tree. The problem with falling out of a tree is that you have no information of where this happened, and often this knowledge is useful information. Given any binary tree, an extended binary tree is one which is formed by replacing each missing child (a null pointer) with a special leaf node, called an external node. The remaining nodes are called internal nodes. An example is shown in the figure above right, where the external nodes are shown as squares and internal nodes are shown as circles. Note that if the original tree is empty, the extended tree consists of a single external node. Also observe that each internal node has exactly two children and each external node has no children.

Let $n$ denote the number of internal nodes in an extended binary tree. Can we predict how many external nodes there will be? It is a bit surprising but the answer is yes, and in fact the number of extended nodes is $n+1$. The proof is by induction. This sort of induction is so common on binary trees, that it is worth going through this simple proof to see how such proofs work in general.

Claim: An extended binary tree with $n$ internal nodes has $n+1$ external nodes.
Proof: By (strong) induction on the size of the tree. Let $X(n)$ denote the number of external nodes in a binary tree of $n$ nodes. We want to show that for all $n \geq 0, X(n)=n+1$.
The basis case is for a binary tree with 0 nodes. In this case the extended tree consists of a single external node, and hence $X(0)=1$.
Now let us consider the case of $n \geq 1$. By the induction hypothesis, for all $0 \leq n^{\prime}<n$, we have $X\left(n^{\prime}\right)=n^{\prime}+1$. We want to show that it is true for $n$. Since $n \geq 1$ the tree contains a root node. Among the remaining $n-1$ nodes, some number $k$ are in the left subtree, and the other $(n-1)-k$ are in the right subtree. Note that $k$ and $(n-1)-k$ are both less than $n$ and so we may apply the induction hypothesis. Thus, there are $X(k)=k+1$ external nodes in the left subtree and $X((n-1)-k)=n-k$ external nodes in the right subtree. Summing these we have

$$
X(n)=X(k)+X((n-1)-k)=(k+1)+(n-k)=n+1
$$

which is what we wanted to show.
Remember this general proof structure. When asked to prove any theorem on binary trees by induction, the same general structure applies.

Complete Binary Trees: We have discussed linked allocation strategies for general trees and binary trees. Is it possible to allocate trees using sequential (that is, array) allocation? In general it is not possible because of the somewhat unpredictable structure of trees. However, there is a very important case where sequential allocation is possible.

Complete Binary Tree: is a binary tree in which every level of the tree is completely filled, except possibly the bottom level, which is filled from left to right.

It is easy to verify that a complete binary tree of height $h$ has between $2^{h}$ and $2^{h+1}-1$ nodes, implying that a tree with $n$ nodes has height $O(\log n)$. (We leave these as exercises involving geometric series.) An example is provided in the figure below.
The extreme regularity of complete binary trees allows them to be stored in arrays, so no additional space is wasted on pointers. Consider an indexing of nodes of a complete tree from 1 to $n$ in increasing level order (so that the root is numbered 1 and the last leaf is numbered $n)$. Observe that there is a simple mathematical relationship between the index of a node and the indices of its children and parents.
In particular:
$\operatorname{leftChild}(i):$ if $(2 i \leq n)$ then $2 i$, else null.
$\operatorname{rightChild}(i):$ if $(2 i+1 \leq n)$ then $2 i+1$, else null.
parent $(i)$ : if $(i \geq 2)$ then $\lfloor i / 2\rfloor$, else null.
Observe that the last leaf in the tree is at position $n$, so adding a new leaf simply means inserting a value at position $n+1$ in the list and updating $n$.


Figure 10: A complete binary tree.

Threaded Binary Trees: Note that binary tree traversals are defined recursively. Therefore a straightforward implementation would require extra space to store the stack for the recursion. The stack will save the contents of all the ancestors of the current node, and hence the additional space required is proportional to the height of the tree. (Either you do it explicitly or the system handles it for you.) When trees are balanced (meaning that a tree with $n$ nodes has $O(\log n)$ height) this is not a big issue, because $\log n$ is so smaller compared to $n$. However, with arbitrary trees, the height of the tree can be as high as $n-1$. Thus the required stack space can be considerable.
This raises the question of whether there is some way to traverse a tree without using additional storage. There are two tricks for doing this. The first one involves altering the links in the tree as we do the traversal. When we descend from parent to child, we reverse the parent-child link to point from parent to the grandparent. These reversed links provide a way to back up the tree when the recursion bottoms out. On backing out of the tree, we "unreverse" the links, thus restoring the original tree structure. We will leave the details as an exercise.

The second method involves a clever idea of using the space occupied for the null pointers to store information to aid in the traversal. In particular, each left-child pointer that would normally be null is set to the inorder predecessor of this node. Similarly each right-child pointer that would normally be null is set to the inorder successor of this node. The resulting links are called threads. This is illustrated in the figure below (where threads are shown as broken curves).

Each such pointer needs to have a special "mark bit" to indicate whether it used as a parentchild link or as a thread. So the additional cost is only two bits per node. Now, suppose that you are currently visiting a node $u$. How do we get to the inorder successor of $u$ ? If the right child pointer is a thread, then we just follow it. Otherwise, we go the right child, and then traverse left-child links until reaching the bottom of the tree (namely a threaded link).

```
BinaryTreeNode nextInOrder() {
    BinaryTreeNode q = right;
    if (rightIsThread) return q;
    while (!q.leftIsThread) {
        q = q.left;
    }
    return q;
}
```

```
// inorder successor of "this"
// go to right child
// if thread, then done
// else q is right child
// go to left child
// ...until hitting thread
```



Figure 11: A Threaded Tree.

For example, in the figure below, suppose that we start with node $A$ in each case. In the left case, we immediately hit a thread, and return $B$ as the result. In the right case, the right pointer is not a thread, so we follow it to $B$, and then follow left links to $C$ and $D$. But $D$ 's left child pointer is a thread, so we return $D$ as the final successor. Note that the entire process starts at the first node in an inorder traversal, that is, the leftmost node in the tree.


Figure 12: Inorder successor in a threaded tree.

## Lecture 6: Minimum Spanning Trees and Prim's Algorithm

(Thursday, Feb 15, 2001)
Read: Section 9.5 in Weiss.
Minimum Spanning Trees: A common problem in communications networks and circuit design is that of connecting together a set of nodes (communication sites or circuit components) by a network of minimal total length (where length is the sum of the lengths of connecting wires). We assume that the network is undirected. To minimize the length of the connecting network, it never pays to have any cycles (since we could break any cycle without destroying connectivity and decrease the total length). Since the resulting connection graph is connected, undirected, and acyclic, it is a free tree.

The computational problem is called the minimum spanning tree problem (MST for short). More formally, given a connected, undirected graph $G=(V, E)$, a spanning tree is an acyclic subset of edges $T \subseteq E$ that connects all the vertices together. Assuming that each edge $(u, v)$ of $G$ has a numeric weight or cost, $w(u, v)$, (may be zero or negative) we define the cost of a
spanning tree $T$ to be the sum of edges in the spanning tree

$$
w(T)=\sum_{(u, v) \in T} w(u, v)
$$

A minimum spanning tree (MST) is a spanning tree of minimum weight. Note that the minimum spanning tree may not be unique, but it is true that if all the edge weights are distinct, then the MST will be distinct (this is a rather subtle fact, which we will not prove). The figure below shows three spanning trees for the same graph, where the shaded rectangles indicate the edges in the spanning tree. The one on the left is not a minimum spanning tree, and the other two are. (An interesting observation is that not only do the edges sum to the same value, but in fact the same set of edge weights appear in the two MST's. Is this a coincidence? We'll see later.)


Figure 13: Spanning trees (the middle and right are minimum spanning trees.

Steiner Minimum Trees: Minimum spanning trees are actually mentioned in the U.S. legal code. The reason is that AT\&T was a government supported monopoly at one time, and was responsible for handling all telephone connections. If a company wanted to connect a collection of installations by an private internal phone system, AT\&T was required (by law) to connect them in the minimum cost manner, which is clearly a spanning tree ... or is it?
Some companies discovered that they could actually reduce their connection costs by opening a new bogus installation. Such an installation served no purpose other than to act as an intermediate point for connections. An example is shown in the figure below. On the left, consider four installations that that lie at the corners of a $1 \times 1$ square. Assume that all edge lengths are just Euclidean distances. It is easy to see that the cost of any MST for this configuration is 3 (as shown on the left). However, if you introduce a new installation at the center, whose distance to each of the other four points is $1 / \sqrt{2}$. It is now possible to connect these five points with a total cost of of $4 / \sqrt{2}=2 \sqrt{2} \approx 2.83$. This is better than the MST.


Figure 14: Steiner Minimum tree.
In general, the problem of determining the lowest cost interconnection tree between a given set of nodes, assuming that you are allowed additional nodes (called Steiner points) is called
the Steiner minimum tree (or SMT for short). An interesting fact is that although there is a simple greedy algorithm for MST's (as we will see below), the SMT problem is much harder, and in fact is NP-hard. (By the way, the US Legal code is rather ambiguous on the point as to whether the phone company was required to use MST's or SMT's in making connections.)

Generic approach: We will present a greedy algorithm (called Prim's algorithm) for computing a minimum spanning tree. A greedy algorithm is one that builds a solution by repeated selecting the cheapest (or generally locally optimal choice) among all options at each stage. An important characteristic of greedy algorithms is that once they make a choice, they never "unmake" this choice. Before presenting these algorithms, let us review some basic facts about free trees. They are all quite easy to prove.

## Lemma:

- A free tree with $n$ vertices has exactly $n-1$ edges.
- There exists a unique path between any two vertices of a free tree.
- Adding any edge to a free tree creates a unique cycle. Breaking any edge on this cycle restores a free tree.

Let $G=(V, E)$ be an undirected, connected graph whose edges have numeric edge weights (which may be positive, negative or zero). The intuition behind Prim's algorithms is simple, we maintain a subtree $A$ of the edges in the MST. Initially this set is empty, and we will add edges one at a time, until $A$ equals the MST. We say that a subset $A \subseteq E$ is viable if $A$ is a subset of edges in some MST (recall that it is not unique). We say that an edge $(u, v) \in E-A$ is safe if $A \cup\{(u, v)\}$ is viable. In other words, the choice $(u, v)$ is a safe choice to add so that $A$ can still be extended to form an MST. Note that if $A$ is viable it cannot contain a cycle. Prim's algorithm operates by repeatedly adding a safe edge to the current spanning tree.

When is an edge safe? We consider the theoretical issues behind determining whether an edge is safe or not. Let $S$ be a subset of the vertices $S \subseteq V$. A cut $(S, V-S)$ is just a partition of the vertices into two disjoint subsets. An edge $(u, v)$ crosses the cut if one endpoint is in $S$ and the other is in $V-S$. Given a subset of edges $A$, we say that a cut respects $A$ if no edge in $A$ crosses the cut. It is not hard to see why respecting cuts are important to this problem. If we have computed a partial MST, and we wish to know which edges can be added that do not induce a cycle in the current MST, any edge that crosses a respecting cut is a possible candidate.
An edge of $E$ is a light edge crossing a cut, if among all edges crossing the cut, it has the minimum weight (the light edge may not be unique if there are duplicate edge weights). Intuition says that since all the edges that cross a respecting cut do not induce a cycle, then the lightest edge crossing a cut is a natural choice. The main theorem which drives both algorithms is the following. It essentially says that we can always augment $A$ by adding the minimum weight edge that crosses a cut which respects $A$. (It is stated in complete generality, so that it can be applied to both algorithms.)

MST Lemma: Let $G=(V, E)$ be a connected, undirected graph with real-valued weights on the edges. Let $A$ be a viable subset of $E$ (i.e. a subset of some MST), let $(S, V-S)$ be any cut that respects $A$, and let $(u, v)$ be a light edge crossing this cut. Then the edge $(u, v)$ is safe for $A$.
Proof: It will simplify the proof to assume that all the edge weights are distinct. Let $T$ be any MST for $G$. If $T$ contains $(u, v)$ then we are done. Suppose that no MST contains $(u, v)$. We will derive a contradiction.


A

$T+(u, v)$

$T^{\prime}=T-(x, y)+(u, v)$
$w\left(T^{\prime}\right)=w(T)-8+4$

Figure 15: MST Lemma.

Add the edge $(u, v)$ to $T$, thus creating a cycle. Since $u$ and $v$ are on opposite sides of the cut, and since any cycle must cross the cut an even number of times, there must be at least one other edge $(x, y)$ in $T$ that crosses the cut.
The edge $(x, y)$ is not in $A$ (because the cut respects $A$ ). By removing $(x, y)$ we restore a spanning tree, call it $T^{\prime}$. We have

$$
w\left(T^{\prime}\right)=w(T)-w(x, y)+w(u, v)
$$

Since $(u, v)$ is lightest edge crossing the cut, we have $w(u, v)<w(x, y)$. Thus $w\left(T^{\prime}\right)<$ $w(T)$. This contradicts the assumption that $T$ was an MST.

Prim's Algorithm: There are two well-known greedy algorithms for computing MST's: Prim's algorithm and Kruskal's algorithm. We will discuss Prim's algorithm here. Prim's algorithm runs in $O((V+E) \log V)$ time. Prim's algorithm builds the tree up by adding leaves one at a time to the current tree. We start with a root vertex $r$ (it can be any vertex). At any time, the subset of edges $A$ forms a single tree. We look to add a single vertex as a leaf to the tree. The process is illustrated in the following figure.


Figure 16: Prim's Algorithm.
Observe that if we consider the set of vertices $S$ currently part of the tree, and its complement $(V-S)$, we have a cut of the graph and the current set of tree edges $A$ respects this cut. Which edge should we add next? The MST Lemma from the previous lecture tells us that it is safe to add the light edge. In the figure, this is the edge of weight 4 going to vertex $u$. Then
$u$ is added to the vertices of $S$, and the cut changes. Note that some edges that crossed the cut before are no longer crossing it, and others that were not crossing the cut are.
It is easy to see, that the key questions in the efficient implementation of Prim's algorithm is how to update the cut efficiently, and how to determine the light edge quickly. To do this, we will make use of a priority queue data structure. Recall that this is the data structure used in HeapSort. This is a data structure that stores a set of items, where each item is associated with a key value. The priority queue supports three operations.
ref = insert (u, key): Insert vertex $u$ with the key value key in the priority queue. The operation returns a pointer ref to allow future access to this item.
$\mathrm{u}=$ extractMin(): Extract the vertex with the minimum key value in $Q$.
decreaseKey (ref, newKey): Decrease the value of the entry with the given reference ref. The new key value is newKey.

A priority queue can be implemented using the same heap data structure used in heapsort. The above operations can be performed in $O(\log n)$ time, where $n$ is the number of items in the heap. Some care needs to be taken in how references are handled. The various heap operations cause items in the heap to be moved around. But the references cannot be moved, because they are our only way of accessing heap nodes from the outside. This can be done by having the heap store pointers to the reference objects.

What do we store in the priority queue? At first you might think that we should store the edges that cross the cut, since this is what we are removing with each step of the algorithm. The problem is that when a vertex is moved from one side of the cut to the other, this results in a complicated sequence of updates.
There is a much more elegant solution, and this is what makes Prim's algorithm so nice. For each vertex in $u \in V-S$ (not part of the current spanning tree) we associate $u$ with a key value key [u], which is the weight of the lightest edge going from $u$ to any vertex in $S$. We also store in pred [u] the end vertex of this edge in $S$. If there is not edge from $u$ to a vertex in $V-S$, then we set its key value to $+\infty$. We will also need to know which vertices are in $S$ and which are not. We do this by coloring the vertices in $S$ black.
Here is Prim's algorithm. The root vertex $r$ can be any vertex in $V$.
Prim's Algorithm

```
Prim(Graph G, Vertex r) {
    for each (u in V) { // initialization
        key[u] = +infinity;
        color[u] = white;
        pred[u] = null;
    }
    key[r] = 0; // start at root
    Q = new PriQueue(V); // put vertices in Q
    while (Q.nonEmpty()) { // until all vertices in MST
        u = Q.extractMin(); // vertex with lightest edge
        for each (v in Adj[u]) {
            if ((color[v] == white) && (w(u,v) < key[v])) {
                key[v] = w(u,v); // new lighter edge out of v
                Q.decreaseKey(v, key[v]);
                pred[v] = u;
            }
        }
        color[u] = black;
```

```
    }
    [The pred pointers define the MST as an inverted tree rooted at r]
}
```

The following figure illustrates Prim's algorithm. The arrows on edges indicate the predecessor pointers, and the numeric label in each vertex is the key value.


Figure 17: Prim's Algorithm.
To analyze Prim's algorithm, we account for the time spent on each vertex as it is extracted from the priority queue. It takes $O(\log V)$ to extract this vertex from the queue. For each incident edge, we spend potentially $O(\log V)$ time decreasing the key of the neighboring vertex. Thus the time is $O(\log V+\operatorname{deg}(u) \log V)$ time. The other steps of the update are constant time. So the overall running time is

$$
\begin{aligned}
T(V, E) & =\sum_{u \in V}(\log V+\operatorname{deg}(u) \log V)=\sum_{u \in V}(1+\operatorname{deg}(u)) \log V \\
& =\log V \sum_{u \in V}(1+\operatorname{deg}(u))=(\log V)(V+2 E)=\Theta((V+E) \log V)
\end{aligned}
$$

Since $G$ is connected, $V$ is asymptotically no greater than $E$, so this is $\Theta(E \log V)$.

## Lecture 7: Binary Search Trees

(Tuesday, Feb 20, 2001)
Read: Chapt 4 of Weiss, through 4.3.6.
Searching: Searching is among the most fundamental problems in data structure design. Traditionally, we assume that we are given a set of elements $\left\{E_{1}, E_{2}, \ldots, E_{n}\right\}$, where each $E_{i}$ is associated with a distinct key value, $x_{i}$ over some totally ordered domain. Given an arbitrary search key $x$, we wish to determine whether there is a element containing this key.

Assuming that each object contains a key field is sometimes a bit awkward (since it may involve accessing private data). A somewhat more attractive way to handle this in modern languages like C++ and Java is to assume that the element $E$ has a comparison method defined for it. (In Java we would say that the object is a subclass of type Comparable.) This means that the class provides a method int compareTo(), whose argument is of the same type as this class. Since we are not constrained to using one programming language, we will allow ourselves to overload the comparison operators (which Java does not permit) to make our pseudo-code easier to read.

| Relationship | Expressed in Java | What we'll write |
| :--- | :--- | :--- |
| $r_{1}<r_{2}$ | r1.compareTo (r2) < 0 | r1 < r2 |
| $r_{1}=r_{2}$ | r1.compareTo (r2) == 0 | r1 == r2 |
| $r_{1}>r_{2}$ | r1.compareTo (r2) >0 | r1 > r2 |

Henceforth, we will not mention keys, but instead assume that our base Element type has such a comparison method. But we will abuse our own convention from time to time by refering to $x$ sometimes as "key" (when talking about ordering and comparisons) and sometimes as an "element" (when talking about insertion).

The Dictionary ADT: Perhaps the most basic example of a data structure based on the structure is the dictionary. A dictionary is an ADT which supports the operations of insertion, deletion, and finding. There are a number of additional operations that one may like to have supported, but these seem to be the core operations. Throughout, let us assume that $x$ is of type Element.
insert (Element x): Insert $x$ into the dictionary. Recall that we assume that keys uniquely identify records. Thus, if an element with the same key as $x$ already exists in the structure, there are a number of possible responses. These include ignoring the operation, returning an error status (or printing an error message), or throwing an exception. Throwing the exception is probably the cleanest option, but not all languages support exceptions.
delete (Element x): Delete the element matching $x$ 's key from the dictionary. If this key does not appear in the dictionary, then again, some special error handling is needed.
Element find(Element x): Determine whether there is an element matching $x$ 's key in the dictionary? This returns a pointer to the associated object, or null, if the object does not appear in the dictionary.

Other operations that might like to see in a dictionary include printing (or generally iterating through) the entries, range queries (find or count all objects in a range of values), returning or extracting the minimum or maximum element, and computing set operations such as union and intersection. There are three common methods for storing dictionaries: sorted arrays, hash tables, and binary search trees. We discuss two of these below. Hash tables will be presented later this semester.

Sequential Allocation: The most naive idea is to simply store the keys in a linear array and run sequentially through the list to search for an element. Although this is simple, it is not efficient unless the set is small. Given a simple unsorted list insertion is very fast $O(1)$ time (by appending to the end of the list). However searching takes $O(n)$ time in the worst case.
In some applications, there is some sort of locality of reference, which causes the same small set of keys to be requested more frequently than others over some subsequence of access requests. In this case there are heuristics for moving these frequently accessed items to the front of the list, so that searches are more efficient. One such example is called move to front. Each time a key is accessed, it is moved from its current location to the front of the list. Another example
is called transpose, which causes an item to be moved up one position in the list whenever it is accessed. It can be proven formally that these two heuristics do a good job in placing the most frequently accessed items near the front of the list.
Instead, if the keys are sorted by key value then the expected search time can be reduced to to $O(\log n)$ through binary search. Given that you want to search for a key $x$ in a sorted array, we access the middle element of the array. If $x$ is less than this element then recursively search the left sublist, if $x$ is greater then recursively search the right sublist. You stop when you either find the element or the sublist becomes empty. It is a well known fact that the number of probes needed by binary search is $O(\log n)$. The reason is quite simple, each probe eliminates roughly one half of the remaining items from further consideration. The number of times you can "halve" a list of size $n$ is $\lg n$ (where $\lg$ means $\log$ base 2).
Although binary search is fast, it is hard to update the list dynamically, since the list must be maintained in sorted order. This would require $O(n)$ time for insertion and deletion. To fix this we use binary trees, which we describe next.

Binary Search Trees: In order to provide the type of rapid access that binary search offers, but at the same time allows efficient insertion and deletion of keys, the simplest generalization is called a binary search tree. The idea is to store the records in the nodes of a binary tree, such that an inorder traversal visits the nodes in increasing key order. In particular, if $x$ is the key stored in the root node, then the left subtree contains all keys that are less than $x$, and the right subtree stores all keys that are greater than $x$. (Recall that we assume that keys are distinct, so no other keys are equal to $x$.)
Defining such an object in C++ or Java typically involves two class definitions. One class for the tree itself, BinarySearchTree (which is publically accessible) and another one for the individual nodes of this tree, BinaryNode (which is a protected friend of the BinarySearchTree). The main data object in the tree is a pointer to root node. The methods associated with the binary search tree class are then transfered to operations on the nodes of the tree. (For example, finding a node in the tree is might be transformed to a find operation starting with the root node of the tree.) See our text for an example.

A node in the binary search tree would typically consist of three members: an associated Element, which we will call data, and a left and right pointer to the two children. Of course other information (parent pointers, level links, threads) can be added to this.

Search in Binary Search Trees: The search for a key $x$ proceeds as follows. The search key is compared to the current node's key, say $y$. If it matches, then we are done. Otherwise if $x<y$ we recursively search the left subtree and if $x>y$ then we recursively search the right subtree. A natural way to handle this would be to make the search procedure a recursive member function of the BinaryNode class. The one technical hassle in implementing this in $\mathrm{C}++$ or Java is that we cannot call a member function for a null pointer, which happens naturally if the search fails and we fall out the bottom of the tree. There are a number of ways to deal with this. Our book proposes making the search function a member of the BinarySearchTree class, and passing in the pointer to the current BinaryNode as an argument. (But this is not the only way to handle this.)

By the way, the initial call is made from the find() method associated with the BinarySearchTree class, which invokes find(x, root), where root is the root of the tree.

Recursive Binary Tree Search

```
Element find(Element x, BinaryNode p) {
        if (p == null) return null; // didn't find it
    else if (x < p.data) // x is smaller?
```

```
        return find(x, p.left); // search left
    else if (x > p.data) // x is larger?
        return find(x, p.right); // search right
    else return p.data; // found it!
}
```

It should be pretty easy to see how this works, so we will leave it to you to verify its correctness. We will often express such simple algorithms in recursive form. Such simple procedures can often be implemented iteratively, thus saving some of the overhead induced by recursion. Here is a nonrecursive version of the same algorithm.

Nonrecursive Binary Tree Search

```
Element find(Element x) {
    BinaryTreeNode p = root;
    while (p != null) {
        if (x < p.data) p = p.left;
        else if (x > p.data) p = p.right;
        else return p.data;
    }
    return null;
}
```

Given the simplicity of the nonrecursive version, why would anyone ever use the recursive version? The answer is the no one would. However, we will see many more complicated algorithms that operate on trees, and these algorithms are almost always much easier to present and understand in recursive form (once you get use to recursive thinking!).
What is the worst-case running time of the search algorithm? Both of them essentially take constant time for each node they visit, and then descend to one of the descendents. In the worst-case the search will visit the deepest node in the tree. Hence the running time is $O(h)$, where $h$ is the height of the tree. If the tree contains $n$ elements, $h$ could vary anywhere from $\lg n$ (for a balanced tree) up to $n-1$ for a degenerate (path) tree.

Insertion: To insert a new element in a binary search tree, we essentially try to locate the key in the tree. At the point that we "fall out" of the tree, we insert a new leaf node containing the desired element. It turns out that this is always the right place to put the new node.
As in the previous case, it is probably easier to write the code in its nonrecursive form, but let's try to do the recursive version, since the general form will be useful for more complex tree operations. The one technical difficulty here is that when a new node is created, we need to "reach up" and alter one of the pointer fields in the parent's node. To do this the procedure returns a pointer to the subtree with the newly added element. Note that "most" of the time, this is returning the same value that is passed in, and hence there is no real change taking place.
The initial call from the BinarySearchTree object is root $=$ insert( $x$, root). We assume that there is a constructor for the BinaryNode of the form BinaryNode (data, left, right). An example is shown below.

```
BinaryNode insert(Element x, BinaryNode p) {
    if (p == NULL)
        p = new BinaryNode(x, null, null);
    else if (x < p.data) p.left = insert(x, p.left);
```

```
    else if (x > p.data) p.right = insert(x, p.right);
```

    else ...Error: attempt to insert duplicate!...
    return p
    \}


Figure 18: Binary tree insertion.

## Lecture 8: More on Binary Search Trees

(Thursday, Feb 22, 2001)
Read: Chapt 4 of Weiss, through 4.4.
Deletion: We continue our discussion of binary search trees, and consider how to delete an element from the tree. Deletion is a more involed than insertion. There are a couple of cases to consider. If the node is a leaf, it can just be deleted with no problem. If it has no left child (or equivalently no right child) then we can just replace the node with it's left (or right) child.


Figure 19: Binary tree deletion: Case of one child.
If both children are present then things are trickier. Removal of the node would create a "hole" in the tree, and we need to fill this hole. The idea is to fill the hole with the element either immediately preceding or immediately following the deleted element. Because the tree is ordered according to an inorder traversal, the candidate elements are the inorder predecessor or successor to this node. We will delete, say, the inorder successor, and make this the replacement. You might ask, what if the inorder successor has two children? It turns out this cannot happen. (You should convince yourself of this. Note that because this node has two children, the inorder successor will be the leftmost node in its right subtree.)


Figure 20: Binary tree deletion: Case of two children.

Before giving the code we first give a utility function, findMin(), which returns a pointer to the element with the minimum key value in a tree. This will be used to find the inorder successor. This is found by just following left links as far as possible. We assume that the argument $p$ is non-null (which will be the case as we shall see later). (Also note that variables are being passed by value, so modifying $p$ will have no effect on the actual argument.) As with the insertion method, the initial call is made to the root of the tree, delete(x, root). Also, as with insertion, rather than using a parent link to reach up and modify the parent's link, we return the new result, and the parent stores this value in the appropriate link.

Binary Tree Deletion

```
BinaryNode findMin(BinaryNode p) { // assume (p != null)
    while (p.left != null) p = p.left;
    return p;
}
BinaryNode delete(Element x, BinaryNode p) {
    if (p == null) // fell out of tree?
            ... Error: deletion of nonexistent element!...
    else {
        if (x < p.data) // x in left subtree
                p.left = delete(x, p.left);
            else if (x > p.data) // x in right subtree
            p.right = delete(x, p.right);
                    // x here, either child empty?
            else if (p.left == null || p.right == null) {
                BinaryNode repl; // get replacement
                if (p.left == null) repl = p.right;
                if (p.right == null) repl = p.left;
                return repl;
            }
            else { // both children present
                p.data = findMin(p.right).data; // copy replacement
                p.right = delete(p.data, p.right); // now delete the replacement
            }
    }
```

```
    return p;
```

\}

In typical Java fashion, we did not worry about deleting nodes. Where would this be added to the above code? Notice that we did not have a special case for handling leaves. Can you see why this is correct as it is?

Analysis of Binary Search Trees: It is not hard to see that all of the procedures find(), insert(), and delete() run in time that is proportional to the height of the tree being considered. (The delete() procedure is the only one for which this is not obvious. Note that each recursive call moves us lower in the tree. Also note that findMin(), is only called once, since the node that results from this operation does not have a left child, and hence will itself not generate another call to findMin.)
The question is, given a binary search tree $T$ containing $n$ keys, what is the height of the tree? It is not hard to see that in the worst case, if we insert keys in either strictly increasing or strictly decreasing order, then the resulting tree will be completely degenerate, and have height $n-1$. On the other hand, following the analogy from binary search, if the first key to be inserted is the median element of the set of keys, then it will nicely break the set of keys into two sets of sizes roughly $n / 2$ each. This will result in a nicely balanced tree, whose height will be $O(\log n)$. Thus, if you had a million nodes, the worst case height is a million and the best is around 20, which is quite a difference. An interesting question is what is the expected height of the tree, assuming that keys are inserted in random order.
Our textbook gives a careful analysis of the average case, but it is based on solving a complex recurrence. Instead, we will present a simple probabilistic analysis to show that the number of nodes along the leftmost chain of the tree is expected to be $\ln n$. Since the leftmost path is a representative path in the tree, it is not surprising the that average path length is still $O(\log n)$ (and a detailed analysis shows that the expected path length is actually $2 \ln n$ ).
The proof is based on the following observation. Consider a fixed set of key values and suppose that these keys are inserted in random order. One way to visualize the tree is in terms of the graphical structure shown in the following figure. As each new node is added, it locates the current interval containing it, and it splits a given interval into two subintervals. As we insert keys, when is it that we add a new node onto the left chain? This happens when the current key is smaller than all the other keys that we have seen up to now. For example, for the insertion order $\langle 9,3,10,6,1,8,5\rangle$ the minimum changes three times, when 9,3 , and 1 are inserted. The corresponding tree has three nodes along the leftmost chain ( 9,3 , and 1 ). For the insertion order $\langle 8,9,5,10,3,6,1\rangle$ the minimum changes four times, when $8,5,3$, and 1 are inserted. The leftmost chain has four nodes $(8,5,3$, and 1$)$.
So this suggests the following probabilistic question: As you scan a randomly permuted list of distinct numbers from left to right, how often do you expect the minimum value to change? Let $p_{i}$ denote the probability that the minimum changes when the $i$ th key is being considered. What is the probability that the $i$ th element is a new minimum? For this to happen, this key must be the smallest among all the first $i$ keys. We assert that the probability that this happens is $1 / i$. If you consider the first $i$ keys, exactly one of them is the minimum among these keys. Since we are considering random permutations, this minimum is equally likely to be in any of the $i$ positions. Thus, there is a $1 / i$ chance that it is in the last ( $i$ th) position.
So with probability $p_{i}=1 / i$ the minimum changes when we consider the $i$ th key. Let $x_{i}$ be a random indicator variable which is 1 if the minimum changes with the $i$ th key and 0 otherwise. The total number of minimum changes is the random variable $X=x_{1}+x_{2}+\ldots+x_{n}$. We have $x_{i}=1$ with probability $p_{i}$, and $x_{i}=0$ with probability $1-p_{i}$. Thus the expected number

Insertion order: 93106185


Left chain has 3 nodes


Minimum changes 3 times

Insertion order: 89510361


Left chain has 4 nodes


Minimum changes 4 times

Figure 21: Length of the leftmost chain.
of changes is

$$
E(X)=\sum_{i=1}^{n}\left(1 \cdot p_{i}+0 \cdot\left(1-p_{i}\right)\right)=\sum_{i=1}^{n} p_{i}=\sum_{i=1}^{n} \frac{1}{i}=\ln n+O(1)
$$

The last line follows from the fact that the sum $(1 / i)$ is the Harmonic series, which we saw in an earlier lecture. This completes the proof that the expected length of the leftmost chain is roughly $\ln n$.
Interestingly this analysis breaks down if we are doing deletions. It can be shown that if we alternate random insertions and random deletions (keeping the size of the tree steady around $n$ ), then the height of the tree will settle down at $O(\sqrt{n})$, which is worse that $O(\log n)$. The reason has to do with the fact that the replacement element was chosen in a skew manner (always taking the minimum from the right subtree). Over the course of many deletions, this can result in a tree that is left-heavy. This can be fixed by alternating taking the replacement from the right subtree with the left subtree resulting in a tree with expected height $O(\log n)$.

## Lecture 9: AVL Trees

(Tuesday, Feb 27, 2001)
Read: Chapt 4 of Weiss, through 4.4.
Balanced Binary Trees: The (unbalanced) binary tree described earlier is fairly easy to implement, but suffers from the fact that if nodes are inserted in increasing or decreasing order (which is not at all an uncommon phenomenon) then the height of the tree can be quite bad. Although worst-case insertion orders are relatively rare (as indicated by the expected case analysis), even inserting keys in increasing order produces a worst-case result. This raises the question of whether we can design a binary search tree which is guaranteed to have $O(\log n)$ height, irrespective of the order of insertions and deletions. The simplest example is that of the AVL tree.

AVL Trees: AVL tree's are height balanced trees. The idea is at each node we need to keep track of balance information, which indicates the differences in height between the left and right subtrees. In a perfectly balanced (complete) binary tree, the two children of any internal node have equal heights. However, maintaining a complete binary tree is tricky, because even a single insertion can cause a large disruption to the tree's structure. But we do not have to do much to remedy this situation. Rather than requiring that both children have exactly the same height, we only require that their heights differ by at most one. The resulting search
trees are called $A V L$ trees (named after the inventors, Adelson-Velskii and Landis). These trees maintain the following invariant:

AVL balance condition: For every node in the tree, the heights of its left subtree and right subtree differ by at most 1 . (The height of a null subtree is defined to be -1 by convention.)

In order to maintain the balance condition we can add a new field, balance to each node, which stores the difference in the height of the right subtree and the height of the left subtree. This number will always be either $-1,0$, or +1 in an AVL tree. (Thus it can be stored using only 2 bits from each node.) Rather than using the balance field in the code below, we will assume a simpler method in which we store the height of each subtree.
Before discussing how we maintain this balance condition we should consider the question of whether this condition is strong enough to guarantee that the height of an AVL tree with $n$ nodes will be $O(\log n)$. To prove this, let's let $N(h)$ denote the minimum number of nodes that can be in an AVL tree of height $h$. We can generate a recurrence for $N(h)$. Clearly $N(0)=1$. In general $N(h)$ will be 1 (for the root) plus $N\left(h_{L}\right)$ and $N\left(h_{R}\right)$ where $h_{L}$ and $h_{R}$ are the heights of the two subtrees. Since the overall tree has height $h$, one of the subtrees must have height $h-1$, suppose $h_{L}$. To make the other subtree as small as possible we minimize its height. It's height can be no smaller than $h-2$ without violating the AVL condition. Thus we have the recurrence

$$
\begin{aligned}
& N(0)=1 \\
& N(h)=N(h-1)+N(h-2)+1 .
\end{aligned}
$$

This recurrence is not well defined since $N(1)$ cannot be computed from these rules, so we add the additional case $N(1)=2$. This recurrence looks very similar to the Fibonacci recurrence $(F(h)=F(h-1)+F(h-2))$. In fact, it can be argued (by a little approximating, a little cheating, and a little constructive induction) that

$$
N(h) \approx\left(\frac{1+\sqrt{5}}{2}\right)^{h}
$$

The quantity $(1+\sqrt{5}) / 2 \approx 1.618$ is the famous Golden ratio. Thus, by inverting this we find that the height of the worst case AVL tree with $n$ nodes is roughly $\log _{\phi} n$, where $\phi$ is the Golden ratio. This is $O(\log n)$ (because log's of different bases differ only by a constant factor).
All that remains is to show how to perform insertions and deletions in AVL trees, and how to restore the AVL balance condition after each insertion or deletion. We will do this next time.

Insertion: The insertion routine for AVL trees starts exactly the same as the insertion routine for binary search trees, but after the insertion of the node in a subtree, we must ask whether the subtree has become unbalanced. If so, we perform a rebalancing step.
Rebalancing itself is a purely local operation (that is, you only need constant time and actions on nearby nodes), but requires a little careful thought. The basic operation we perform is called a rotation. The type of rotation depends on the nature of the imbalance. Let us assume that the source of imbalance is that the left subtree of the left child is too deep. (The right subtree of the right child is handled symmetrically.) See the figure below. The operation performed in this case is a right single rotation. Notice that after this rotation has been performed, the balance factors change. The heights of the subtrees of $b$ and $d$ are now both even with each other.


Figure 22: Single rotation.

On the other hand, suppose that the heavy grandchild is the right subtree of the left subtree. (Again the the left subtree of the right child is symmetrical.) In this case note that a single rotation will not fix the imbalance. However, two rotations do suffice. See the figure below. In particular, do a left rotation on the left-right grandchild (the right child of the left child), and then a right rotation on the left child, you will restore balance. This operation is called a double rotation, and is shown in the figure below. (We list multiple balanced factors to indicate the possible values. We leave it as an exercise to determine how to update the balance factors.)


Figure 23: Left-right double rotation.
Before presenting the code for insertion, we present the utilities routines used. The rotations are called rotateL() and rotateR() for the left and right single rotations, respectively, and rotateLR() and rotateRL() for the left-right and right-left double rotations, respectively. We also use a utility function height ( t ), which returns t . height if $t$ is nonnull and -1 otherwise (the height of a null tree).

AVL Tree Utilities

```
int height(AvlNode t) { return t == null ? -1 : t.height}
AvlNode rotateR(AvlNode t) // right single rotation
{
    AvlNode s = t.left; // (t.bal = -2; s.bal = 0 or 1)
    t.left = s.right; // swap inner child
```

```
    s.right = t; // bring s above t
    t.height = ...exercise... // update subtree heights
    s.height = ...exercise...
    return s; // s replaces t
}
AvlNode rotateLR(AvlNode t)
// left-right double rotation
{
    t.left = rotateL(t.left);
    return rotateR(t);
}
```

The recursive insertion code is given below. We use the same structure given in Weiss. We use the usual trick of passing pointer information up the tree by returning a pointer to the resulting subtree after insertion.

AVL Tree Insertion

```
AvlNode insert(Element x, AvlNode t) { // recursive insertion method
    if (t == null) { // bottom or tree: create new node
            t = new AvlNode(x); // create node and initialize
    }
    else if (x < t.element) {
        t.left = insert(x, t.left); // insert recursively on left
            // check height condition
        if (height(t.left) - height(t.right) == 2) {
            // rotate on the left side
            // left-left insertion
            // left-right insertion
            t = rotateLR(t);
        }
                                    // update height
        t.height = max(height(t.left), height(t.right)) + 1;
    }
    else if (x > t.element) {
        ...symmetric with left insertion...
    }
    else {
        ...Error: duplicate insertion ignored...
    }
    return t;
}
```

As mentioned earlier, and interesting feature of this algorithm (which is not at all obvious) is that after the first rotation, the height of the affected subtree is the same as it was before the insertion, and hence no further rotations are required.

Deletion: Deletion is similar to insertion in that we start by applying the deletion algorithm for unbalanced binary trees. Recall that this breaks into three cases, leaf, single child, and two children. In the two children case we need to find a replacement key. Once the deletion is finished, we walk back up the tree (by returning from the recursive calls), updating the balance factors (or heights) as we go. Whenever we comes to an unbalanced node, we apply an appropriate rotation to remedy the situation.

The deletion code is messier than the insertion code. (We will not present it.) But the idea is the same. Suppose that we have deleted a key from the left subtree and that as a result this subtree's height has decreased by one, and this has caused some ancestor to violate the height balance condition. There are two cases. First, if balance factor for the right child is either 0 or +1 (that is, it is not left heavy), we can perform a single rotation as shown in the figure below. We list multiple balance factors, because there are a number of possibilities. (We leave it as an exercise to figure out which balance factors apply to which cases.)


Figure 24: Single rotation for deletion.
On the other hand, if the right child has a balance factor of -1 (it is left heavy) then we need to perform a double rotation, as shown in the following figure. A more complete example of a deletion is shown in the figure below. We delete element 1 . The causes node 2 to be unbalanced. We perform a single left rotation at 2 . However, now the root is unbalanced. We perform a right-left double rotation to the root.


Figure 25: Double rotation for deletion.

Lazy Deletion: The deletion code for AVL trees is almost twice as long as the insertion code. It is not all that more complicated conceptually, but the number of cases is significantly larger. Our text suggests an interesting alternative for avoiding the pain of coding up the deletion algorithm, called lazy deletion. The idea is to not go through the entire deletion process. For each node we maintain a boolean value indicating whether this element is alive or dead. When a key is deleted, we simply declare it to be dead, but leave it in the tree. If an attempt is made to insert the same key value again, we make the element alive again. Of course, after a


Figure 26: AVL Deletion example.
long sequence of deletions and insertions, it is possible that the tree has many dead nodes. To fix this we periodically perform a garbage collection phase, which traverses the tree, selecting only the live elements, and then building a new AVL tree with these elements.

## Lecture 10: Splay Trees

## (Thursday, March 1, 2001)

Read: Chapt 4 of Weiss, through 4.5.
Splay Trees and Amortization: Recall that we have discussed binary trees, which have the nice property that if keys are inserted and deleted randomly, then the expected times for insert, delete, and member are $O(\log n)$. Because worst case scenarios can lead $O(n)$ behavior per operation, we were lead to the idea of the height balanced tree, or AVL tree, which guarantees $O(\log n)$ time per operation because it maintains a balanced tree at all times. The basic operations that AVL trees use to maintain balance are called rotations (either single or double). The disadvantages of AVL trees are that we need to maintain balance information in the nodes, and the routines for updating the AVL tree are somewhat more complicated that one might generally like.
Today we will introduce a new data structure, called a splay tree. Like the AVL tree, a splay tree is a binary tree, and we will use rotation operations in modifying the tree. Unlike an AVL tree no balance information needs to be stored. Because a splay tree has no balance information, it is possible to create unbalanced splay trees. Splay trees have an interesting self-adjusting nature to them. In particular, whenever the tree becomes unbalanced, accesses to unbalanced portions of the tree will naturally tend to balance themselves out. This is really quite clever, when you consider the fact that the tree has no idea whether it is balanced or not! Thus, like an unbalanced binary tree, it is possible that a single access operation could take as long as $O(n)$ time (and not the $O(\log n)$ that we would like to see). However, the nice property that splay trees have is the following:

Splay Tree Amortized Performance Bound: Starting with an empty tree, the total time needed to perform any sequence of $m$ insertion/deletion/find operations on a splay tree is $O(m \log n)$, where $n$ is the maximum number of nodes in the tree.

Thus, although any one operation may be quite costly, over any sequence of operations there must be a large number of efficient operations to balance out the few costly ones. In other words, over the sequence of $m$ operations, the average cost of an operation is $O(\log n)$.
This idea of arguing about a series of operations may seem a little odd at first. Note that this is not the same as the average case analysis done for unbalanced binary trees. In that case, the average was over the possible insertion orders, which an adversary could choose to make arbitrarily bad. In this case, an adversary could pick the worst sequence imaginable, and it would still be the case that the time to execute the entire sequence is $O(m \log n)$. Thus, unlike the case of the unbalanced binary tree, where the adversary can force bad behavior time after time, in splay trees, the adversary can force bad behavior only once in a while. The rest of the time, the operations execute quite efficiently. Observe that in many computational problems, the user is only interested in executing an algorithm once. However, with data structures, operations are typically performed over and over again, and we are more interested in the overall running time of the algorithm than we are in the time of a single operation.

This type of analysis based on sequences of operations is called an amortized analysis. In the business world, amortization refers to the process of paying off a large payment over time in small installments. Here we are paying for the total running time of the data structure's
algorithms over a sequence of operations in small installments (of $O(\log n)$ each) even though each individual operation may cost much more. Amortized analyses are extremely important in data structure theory, because it is often the case that if one is willing to give up the requirement that every access be efficient, it is often possible to design data structures that are simpler than ones that must perform well for every operation.
Splay trees are potentially even better than standard search trees in one sense. They tend to bring recently accessed data to up near the root, so over time, we may need to search less than $O(\log n)$ time to find frequently accessed elements.

Splaying: As we mentioned earlier, the key idea behind splay trees is that of self-organization. Imagine that over a series of insertions and deletions, our tree has become rather unbalanced. If it is possible to repeatedly access the unbalanced portion of the tree, then we are doomed to poor performance. However, if we can perform an operation that takes unbalanced regions of the tree, and makes them more balanced then that operation is of interest to us. As we said before, since splay trees contain no balance information, we cannot selectively apply this operation at positions of known imbalance. Rather we will perform it everywhere along the access path to every node. This basic operation is called splaying. The word splaying means "spreading", and the splay operation has a tendency to "mix up" trees and make them more random. As we know, random binary trees tend towards $O(\log n)$ height, and this is why splay trees seem to work as they do.
Basically all operations in splay trees begin with a call to a function splay ( $\mathrm{x}, \mathrm{t}$ ) which will reorganize the subtree rooted at node $t$, bringing the node with key value $x$ to the root of the tree, and generally reorganizing the tree along the way. If $x$ is not in the tree, either the node immediately preceding or following $x$ will be brought to the root.
Here is how splay ( $\mathrm{x}, \mathrm{t}$ ) works. We perform the normal binary search descent to find the node $v$ with key value $x$. If $x$ is not in the tree, then let $v$ be the last node visited before we fall out of the tree. If $v$ is the root then we are done. If $v$ is a child of the root, then we perform a single rotation (just as we did in AVL trees) at the root to bring $v$ up to the root's position. Otherwise, if $v$ is at least two levels deep in the tree, we perform one of four possible double rotations from $v$ 's grandparent. In each case the double rotations will have the effect of pulling $v$ up two levels in the tree. We then go up to the new grandparent and repeat the operation. Eventually $v$ will be carried to the top of the tree. In general there are many ways to rotate a node to the root of a tree, but the choice of rotations used in splay trees is very important to their efficiency.
The rotations are selected as follows. Recall that $v$ is the node containing $x$ (or its immediate predecessor or successor). Let $p$ denote $x$ 's parent and let $g$ denote $x$ 's grandparent. There are four possible cases for rotations. If $x$ is the left child of a right child, or the right child of a left child, then we call this a zig-zag case. In this case we perform a double rotation to bring $x$ up to the top. See the figure below. Note that this can be accomplished by performing one single rotation at $p$ and a second at $g$.
Otherwise, if $x$ is the left child of a left child or the right child of a right child we call this a zig-zig case. In this case we perform a new type of double rotation, by first rotating at $g$ and then rotating at $p$. The result is shown in the figure below.

A complete example of splay $(3, t)$ is shown in the next figure.
Splay Tree Operations: Let us suppose that we have implemented the splay operation. How can we use this operation to help us perform the basic dictionary operations of insert, delete, and find?
To find key $x$, we simply call splay ( $\mathrm{x}, \mathrm{t}$ ). If $x$ is in the tree it will be transported to the root. (This is nice, because in many situations there are a small number of nodes that are repeatedly


Figure 27: Zig-zag case.


Figure 28: Zig-zig case.
being accessed. This operation brings the object to the root so the subsequent accesses will be even faster. Note that the other data structures we have seen, repeated find's do nothing to alter the tree's structure.)
Insertion of $x$ operates by first calling $\operatorname{splay}(\mathrm{x}, \mathrm{t})$. If $x$ is already in the tree, it will be transported to the root, and we can take appropriate action (e.g. error message). Otherwise, the root will consist of some key $w$ that is either the key immediately before $x$ or immediately after $x$ in the set of keys. Let us suppose that it is the former case $(w<x)$. Let $R$ be the right subtree of the root. We know that all the nodes in $R$ are greater than $x$ so we make a new root node with $x$ as data value, and make $R$ its right subtree. The remaining nodes are hung off as the left subtree of this node. See the figure below.
Finally to delete a node $x$, we execute $\operatorname{splay}(\mathrm{x}, \mathrm{t})$ to bring the deleted node to the root. If it is not the root we can take appropriate error action. Let $L$ and $R$ be the left and right subtrees of the resulting tree. If $L$ is empty, then $x$ is the smallest key in the tree. We can delete $x$ by setting the root to the right subtree $R$, and deleting the node containing $x$. Otherwise, we perform splay ( $\mathrm{x}, \mathrm{L}$ ) to form a tree $L^{\prime}$. Since all the nodes in this subtree are already less than $x$, this will bring the predecessor $w$ of $x$ (i.e. it will bring the largest key in $L$ to the root of $L$, and since all keys in $L$ are less than $x$, this will be the immediate predecessor of $x$ ). Since this is the largest value in the subtree, it will have no right child. We then make $R$ the right subtree of the root of $L^{\prime}$. We discard the node containing $x$. See the figure below.

Why Splay Trees Work: As mentioned earlier, if you start with an empty tree, and perform any sequence of $m$ splay tree operations (insert, delete, find, say), then the total running time will be $O(m \log n)$, where $n$ is the maximum number of elements in the tree at any time. Thus the average time per operation is $O(\log n)$, as you would like. Proving this involves a complex potential argument, which we will skip. However it is possible to give some intuition as to why splay trees work.


Figure 29: Example showing the result of Splay(3).


Figure 30: Insertion of $x$.


Figure 31: Deletion of $x$.

It is possible to arrange things so that the tree has $O(n)$ height, and hence a splay applied to the lowest leaf of the tree will take $O(n)$ time. However, the process of splaying on this node will result in a tree in which the average node depth decreases by about half. (See the example on page 135 of Weiss.) To analyze a splay tree's performance, we model two quantities:
Real cost: the time that the splay takes (this is proportional to the depth of the node being splayed), and
Increase in balance: the extent to which the splay operation improves the balance of the tree.

The analysis shows that if the real cost is high, then the tree becomes much more balanced (and hence subsequent real costs will be lower). On the other hand, if the tree becomes less balanced, then the real cost will be low. So you win either way.
Consider, for example, the zig-zig case. Label the subtrees as shown in the zig-zig figure above. Notice that the element we were searching for is in the subtree rooted at $x$, and hence is either in $A$ or $B$. Also note that after performing the rotation, depths of the in subtrees $A$ and $B$ have decreased (by two levels for $A$ and one level for $B$ ). Consider the total number of elements in all the subtrees $A, B, C$ and $D$.
$A$ and $B$ are light: If less than half of the elements lie in $A$ and $B$, then we are happy, because as part of the initial search, when we passed through nodes $g$ and $p$, we eliminated over half of the elements from consideration. If this were to happen at every node, then the total real search time would be $O(\log n)$.
$A$ and $B$ are heavy: On the other hand, if over half of the elements were in $A$ and $B$, then the real cost may be much higher. However, in this case over half of these elements have decreased their depth by at least one level. If this were to happen at every level of the tree, the average depth would decrease considerably.

Thus, we are happy in either case: either because the real cost is low or the tree becomes much more balanced.

## Lecture 11: Skip Lists

(Tuesday, March 6, 2001)
Read: Section 10.4.2 in Weiss, and Samet's notes Section 5.1.
Recap: So far we have seen three different method for storing dictionaries. Unbalanced binary trees are simple, and worked well on average, but an adversary could force very bad running time. AVL trees guarantee good performance, but are somewhat harder to implement. Splay trees provided an interesting alternative to AVL trees, because they are simple and self-organizing.
Today we are going to continue our investigation of different data structures for storing dictionaries. The data structure we will consider today is called a skip list. This is among the most practical of the data structures we have seen, since it is quite simple and (based on my own experience) seems to be the fastest of all these methods. A skip list is an interesting generalization of a linked list. As such, it has much of the simplicity of linked lists, but provides optimal $O(\log n)$ performance. Another interesting feature of skip lists is that they are a randomized data structure. In other words, we use a random number generator in creating these trees. We will show that skip lists are efficient in the expected case. However, unlike unbalanced binary search trees, the expectation has nothing to do with the distribution of the keys. It depends only on the random number generator. Hence an adversary cannot pick a sequence of operations for our tree that will always be bad. And in fact, the probability that a skip list might perform badly is very small.

Perfect Skip Lists: Skip lists began with the idea, "how can we make sorted linked lists better?" It is easy to do operations like insertion and deletion into linked lists, but it is hard to locate items efficiently because we have to walk through the list one item at a time. If we could "skip" over lots of items at a time, then we could solve this problem. One way to think of skip lists is as a hierarchy of sorted linked lists, stacked one on top of the other.

To make this more concrete, imagine a linked list, sorted by key value. Let us assume that we have a special sentinel node at the end, called nil, which behaves as if it has a key of $\infty$. Take every other entry of this linked list (say the even numbered entries) and lift them up to a new linked list with $1 / 2$ as many entries. Now take every other entry of this linked list and lift it up to another linked with $1 / 4$ as many entries as the original list. The head and nil nodes are always lifted. We could repeat this process $\lceil\lg n\rceil$ times, until there are is only one element in the topmost list. To search in such a list you would use the pointers at high levels to "skip" over lots of elements, and then descend to lower levels only as needed. An example of such a "perfect" skip list is shown below.


Figure 32: Perfect skip list.
To search for a key $x$ we would start at the highest level. We scan linearly along the list at the current level $i$ searching for first item that is greater than to $x$ (recalling that the nil key value is $\infty$ ). Let $p$ point to the node just before this step. If $p$ 's data value is equal to $x$ then we stop. Otherwise, we descend to the next lower level $i-1$ and repeat the search. At level 0 we have all the keys stored, so if we do not find it at this level we quit. For example, the figure shows the search for $x=19$ in dotted lines.
The search time in the worst case may have to go through all $\lceil\lg n\rceil$ levels (if the key is not in the list). We claim this search visits at most two noded per level of this perfect skip list. This is true because you know that at the previous (higher) level you lie between two consecutive nodes $p$ and $q$, where $p$ 's data value is less than $x$ and $q$ 's data value is greater than $x$. Between any two consecutive nodes at the one level there is exactly one new node at the next lower level. Thus as we descend a level, the search will visit the current node and at most one additional node. Thus the there are at most two nodes visited per level and $O(\log n)$ levels, for a total of $O(\log n)$ time.

Randomized Skip Lists: The problem with the data structure mentioned above is that it is exactly balanced (somewhat like a perfectly balanced binary tree). The insertion of any node would result in a complete restructuring of the list if we insisted on this much structure. Skip lists (like all good balanced data structures) allow a certain amount of imbalance to be present. In fact, skip lists achieve this extra "slop factor" through randomization.
Let's take a look at the probabilistic structure of a skip list at any point in time. (By the way, this is not exactly how the structure is actually built, but it serves to give the intuition behind its structure.) In a skip list we do not demand that exactly every other node at level $i$ be promoted to level $i+1$, instead think of each node at level $i$ tossing a coin. If the coin comes up heads (i.e. with probability $1 / 2$ ) this node promotes itself to the next higher level linked list, and otherwise it stays where it is. Randomization being what it is, it follows that the
expected number of nodes at level 1 is $n / 2$, the expected number at level 2 is $n / 4$, and so on. Furthermore, since nodes appear randomly at each of the levels, we would expect the nodes at a given level to be well distributed throughout (not all bunching up at one end). Thus a randomized skip list behaves much like an idealized skip list in the expected case. The search procedure is exactly the same as it was in the idealized case. See the figure below.


Figure 33: Randomized skip list.
The interesting thing about skip lists is that it is possible to insert and delete nodes into a list, so that this probabilistic structure will hold at any time. For insertion of key $x$ we first do a search on key $x$ to find its immediate predecessors in the skip list (at each level of the structure). If $x$ is not in the list, we create a new node $x$ and insert it at the lowest level of the skip list. We then toss a coin (or equivalently generate a random integer). If the result is tails (if the random number is even) we stop. Otherwise we insert $x$ at the next higher level of the structure. We repeat this process until the coin comes up tails, or we hit the maximum level in the structure. Since this is just repeated link list insertion the code is very simple. To do deletion, we simply delete the node from every level it appears in.
Note that at any point in time, the linked list will have the desired probabilistic structure we mentioned earlier. The reason is that (1) because the adversary cannot see our random number generator, he has no way to selectively delete nodes at a particular level, and (2) each node tosses it's coins independently of the other nodes, so the levels of nodes in the skip list are independent of one another. (This seems to be one important difference between skip lists and trees, since it is hard to do anything independently in a tree, without affecting your children.)

Analysis: The analysis of skip lists is an example of a probabilistic analysis. We want to argue that in the expected case, the search time is $O(\log n)$. Clearly this is the time that dominates in insertion and deletion. First observe that the expected number of levels in a skip list is $O(\log n)$. The reason is that at level 0 we have $n$ keys, at level 1 we expect that $n / 2$ keys survive, at level 2 we expect $n / 4$ keys survive, and so forth. By the same argument we used in the ideal case, after $O(\log n)$ levels, there will be no keys left.
The argument to prove the expected bound on the search time is rather interesting. What we do is look at the reversal of the search path. (This is a common technique is probabilistic algorithms, and is sometimes called backwards analysis.) Observe that the forward search path drops down a level whenever the next link would take us "beyond" the node we are searching
for. When we reverse the search path, observe that it will always take a step up if it can (i.e. if the node it is visiting appears at the next higher level), otherwise it will take a step to the left.


Header
Sentinel
Figure 34: Reversal of search path to $x=11$.
Now, when we arrive at level $i$ of any node in the skip list, we argue that the probability that there is a level above us is just $1 / 2$. The reason is that when we inserted the node, this is the probability that it promoted itself to the next higher level. Therefore with probability $1 / 2$ we step to the next higher level. With the remaining probability $1-(1 / 2)=1 / 2$ we stay at the same level. The expected number of steps needed to walk through $j$ levels of the skip list is given by the following recurrence.

$$
C(j)=1+\frac{1}{2} C(j-1)+\frac{1}{2} C(j) .
$$

The 1 counts the current step. With probability $1 / 2$ we step to the next higher level and so have one fewer level to pass through, and with probability $1 / 2$ we stay at the same level. This can be rewritten as

$$
C(j)=2+C(j-1) .
$$

By expansion it is easy to verify that $C(j)=2 j$. Since $j$ is at most the number of levels in the tree, we have that the expected search time is at most $O(\log n)$.

Implementation Notes: One of the appeals of skip lists is their ease of implementation. Most of procedures that operate on skip lists make use of the same simple code that is used for linked-list operations. One additional element is that you need to keep track of the level that you are currently on. The way that this is done most efficiently, is to have nodes of variable size, where the size of a node is determined randomly as it is created. We take advantage of the fact that Java (and C ++ ) allow us to dynamically allocated arrays of variable size.
The skip list object consists of a header node (header) and the constructor creates a sentinel node, whose key value is set to some special infinite value (which presumably depends on the key type). We assume that the constructor is given the maximum allowable number of levels. (A somewhat smarter implementation would adaptively determine the proper number of levels.)

Skip List Classes

```
class SkipListNode {
    Element data; // key data
    SkipListNode forward[]; // array of forward pointers
                            // Constructor given data and level
    SkipListNode(Element d, int level) {
        data = d
        forward = new SkipListNode[level+1];
    }
}
```

```
class SkipList {
    int maxLevel; // maximum level
    SkipListNode header; // header node
    SkipList(int maxLev) { // constructor given max level number
maxLevel = maxLev; // allocate header node
            header = new SkipListNode(null, maxLevel);
                            // append the "nil" node to the header
        SkipListNode sentinel = new SkipListNode(INFINITY, maxLevel);
        for (int i = 0; i <= maxLevel; i++)
            header.forward[i] = sentinel;
    }
}
```

The code for finding an element in the skip list is given below. Observe that the search is very much the same as a standard linked list search, except for the loop that moves us down one level at a time.

Find an object in the skip list

```
Element find(Element key) {
    SkipListNode current = header; // start at header
    // start search at max level
    for (int i = maxLevel; i >= 0; i--) {
        SkipListNode next = current.forward[i];
        while (next.data < key) { // search forward on level i
            current = next;
            next = current.forward[i];
        }
    }
    current = current.forward[0]; // this must be it
    if (current.data == key) return current.data;
    else return null;
}
```

It is worth noting that you do not need to store the of a node as part of the forward pointer field. The skip list routines maintain their own knowledge of the level as they move around the data structure. Also notice that (if correctly implemented) you will never attempt to index beyond the limit of a node.
We will leave the other elements of the skip list implementation as an exercise. One of the important elements needed for skip list insertion is the generation of the random level number for a node. Here is the code for generating a random level.

```
int generateRandomLevel() {
    int newLevel = 0;
    while (newLevel < maxLevel && Math.random() < 0.5) newLevel++;
    return newLevel;
}
```


## Lecture 12: B-trees

(Thursday, March 8, 2001)

Read: Section 4.7 of Weiss and 5.2 in Samet.
B-trees: Although it was realized quite early it was possible to use binary trees for rapid searching, insertion and deletion in main memory, these data structures were really not appropriate for data stored on disks. When accessing data on a disk, an entire block (or page) is input at once. So it makes sense to design the tree so that each node of the tree essentially occupies one entire block. Although it is possible to segment binary trees in a clever way so that many neighboring nodes are squeezed into one disk block it is difficult to perform updates on the tree to preserve this proximity.
An alternative strategy is that of the $B$-tree, which solves the problem by using multiway trees rather than binary trees. Standard binary search tree have two pointers, left and right, and a single key value $k$ that is used to split the keys in the left subtree from the right subtree. In a $j$-ary multiway search tree node, we have $j$ child pointers, $p_{1}, p_{2}, \ldots, p_{j}$, and $j-1$ key values, $k_{1}<k_{2}<\ldots<k_{j-1}$. We use these key values to split the elements among the subtrees. We assume that the data values $k$ located in subtree pointed to be $p_{i}$ satisfy $k_{i} \leq k<k_{i+1}$, for $1 \leq i \leq j$. For convenience of notatoin, we imagine that there are two sentinel keys, $k_{0}=-\infty$ and $k_{j}=+\infty$. See the figure below.


Figure 35: Binary and 4-ary search tree nodes.
B-trees are multiway search trees, in which we achieve balance by constraining the "width" of each node. As with all other balanced structures, we allow for a certain degree of flexibility in our constraints, so that updates can be performed efficiently. A B-tree is described in terms of a parameter $m$, which controls the maximum degree of a node in the tree, and which typically is set so that each node fits into one disk block.
Our definition differs from the one given in Weiss. Weiss's definition is more reasonable for actual implementation on disks, but ours is designed to make the similarities with other search trees a little more apparent. For $m \geq 3, B$-tree of order $m$ has the following properties:

- The root is either a leaf or has between 2 and $m$ children.
- Each nonleaf node except the root has between $\lceil m / 2\rceil$ and $m$ (nonnull) children. A node with $j$ children contains $j-1$ key values.
- All leaves are at the same level of the tree, and each leaf contains between $\lceil m / 2\rceil-1$ to $m-1$ keys.

So for example, when $m=4$ each internal node has from 2 to 4 children and from 1 to 3 keys. When $m=7$ each internal node has from 4 to 7 children and from 3 to 6 keys (except the root which may have as few as 2 children and 1 key). Why is the requirement on the number of children not enforced on the root of the tree? We will see why later.
The definition in the book differs in that all data is stored in the leaves, and the leaves are allowed to have a different structure from nonleaf nodes. This makes sense for storing large
data records on disks. In the figure below we show a B-tree of order 3. In such a tree each node has either 2 or 3 children, and hence is also called a 2-3 tree. Note that even though we draw nodes of different sizes, all nodes have the same storage capacity.


Figure 36: B-tree of order 3, also known as a 2-3 tree.

Height analysis: Consider a B-tree (of order $m$ ) of height $h$. Since every node has degree at most $m$, then the tree has at most $m^{h}$ leaves. Since each node contains at least one key, this means that $n \geq m^{h}$. We assume that $m$ is a constant, so the height of the tree is given asymptotically by

$$
h \leq \log _{m} n=\frac{\lg n}{\lg m} \in O(\log n)
$$

This is very small as $m$ gets large. Even in the worst case (excluding the root) the tree essentially is splitting $m / 2$ ways, implying that

$$
h \leq \log _{(m / 2)} n=\frac{\lg n}{\lg (m / 2)} \in O(\log n)
$$

Thus, the height of the tree is $O(\log n)$ in either case. For example, if $m=256$ we can store 100,000 records with a height of only three. This is important, since disk accesses are typically many orders of magnitude slower than main memory accesses, it is important to minimize the number of accesses.

Node structure: Unlike skip-lists, where nodes are typically allocated with different sizes, here every node is allocated with the maximum possible size, but not all nodes are fully utilized. A typical B-tree node might have the following Java class structure. In this case, we are storing integers as the elements. We place twice as many elements in each leaf node as in each internal node, since we do not need the storage for child pointers.

```
const int M = ... // order of the tree
class BTreeNode {
    int nChildren; // number of children
    BTreeNode child[M]; // children
    int key[M-1]; // keys
};
```

Search: Searching a B-tree for a key $x$ is a straightforward generalization of binary tree searching. When you arrive at an internal node with keys $k_{1}<k_{2}<\ldots<k_{j-1}$ search (either linearly or by binary search) for $x$ in this list. If you find $x$ in the list, then you are done. Otherwise, determine the index $i$ such that $k_{i}<x<k_{i+1}$. (Recall that $k_{0}=-\infty$ and $k_{j}=+\infty$.) Then recursively search the subtree $p_{i}$. When you arrive at a leaf, search all the keys in this node. If it is not here, then it is not in the B-tree.

Insertion: To insert a key into a B-tree of order $m$, we perform a search to find the appropriate leaf into which to insert the node. If the leaf is not at full capacity (it has fewer than $m-1$ keys) then we simply insert it and are done. (Note that when $m$ is large this is the typical case.) Note that this will involve sliding keys around within the leaf node to make room for the new entry. Since $m$ is assumed to be a constant, this constant time overhead is ignored.
Otherwise the node overflows and we have to take action to restore balance. There are two methods for restoring order in a B-tree.

Key rotation: In spite of this name, this is not the same as rotation in AVL trees. This method is quick and involves little data movement. Unfortunately, it is not always applicable. Let's let $p$ be the node that overflows. We look at our immediate left and right siblings in the B-tree. (Note that one or the other may not exist). Suppose that one of them, the right sibling say, has fewer than the maximum of $m-1$ keys. Let $q$ denote this sibling, and let $k$ denote the key in the parent node that splits the elements in $q$ from those in $p$.
We take the key $k$ from the parent and make it the minimum key in $q$. We take the maximum key from $p$ and use it to replace $k$ in the parent node. Finally we transfer the rightmost subtree from $p$ to become the leftmost subtree of $q$ (and readjust all the pointers accordingly). At this point the tree is balanced and no further rebalancing is needed. An example of this is shown in the following figure. We have shown all the child links (even at the leaf level where they would not normally be used) to make clearer how the update works at any level of the tree. Key rotation to the right child is analogous.


Figure 37: Key rotation.
Node split: If both siblings are full, then key rotation is not possible. Recall that a node can have at most $m-1$ keys and at least $\lceil m / 2\rceil-1$ keys. Suppose after inserting the new key we have an overfull node with $m$ keys, $k_{1}<k_{2}<\ldots<k_{m}$. We split this node into three groups: one with the smallest $\lceil(m-1) / 2\rceil$ keys, a single central element, and one with the largest $\lfloor(m-1) / 2\rfloor$ keys. We copy the smallest group to a new B-tree node, and leave the largest in the existing node. We take the extra key and insert it into the parent.
At this point the parent may overflow, and so we repeat this process recursively. When the root overflows, we split the root into two nodes, and create a new root with two children. (This is why the root is exempt from the normal number of children requirement.) See the figure below for an example. This shows an interesting feature of B-trees, namely that they seem to grow from the root up, rather than from the leaves down.
To see that the new nodes are of proper size it suffices to show that

$$
\left\lceil\frac{m}{2}\right\rceil-1 \leq\left\lfloor\frac{m-1}{2}\right\rfloor \leq\left\lceil\frac{m-1}{2}\right\rceil \leq m
$$



Figure 38: Node split.

The upper bound $(\leq m)$ holds for any integer $m$. To see the lower bound, we observe that if $m$ is odd then

$$
\left\lceil\frac{m}{2}\right\rceil-1=\frac{m+1}{2}-1=\frac{m-1}{2}=\left\lfloor\frac{m-1}{2}\right\rfloor
$$

If $m$ is even then

$$
\left\lceil\frac{m}{2}\right\rceil-1=\frac{m}{2}-1=\frac{m-2}{2}=\left\lfloor\frac{m-1}{2}\right\rfloor .
$$

So we actually have equality in each case.
Deletion: As in binary tree deletion we begin by finding the node to be deleted. We need to find a suitable replacement for this node. This is done by finding the largest key in the left child (or equivalently the smallest key in the right child), and moving this key up to fill the hole. This key will always be at the leaf level. This creates a hole at the leaf node. If this leaf node still has sufficient capacity (at least $\lceil m / 2\rceil-1$ keys) then we are done. Otherwise we need restore the size constraints.
Suppose that after deletion the current node has $\lceil m / 2\rceil-2$ keys. As with insertion we first check whether a key rotation is possible. If one of the two siblings has at least one key more than the minimum, then rotate the extra into this node, and we are done.

Otherwise we know that at least one of our neighbors (say the left) has exactly $\lceil m / 2\rceil-1$ keys. In this case we need to perform a node merge. This is the inverse of a node split. We combine the contents of the current node, together with the contents of its sibling, together with the key from the parent that separates these two. The resulting node has

$$
\left(\left\lceil\frac{m}{2}\right\rceil-2\right)+\left(\left\lceil\frac{m}{2}\right\rceil-1\right)+1
$$

total keys. We will leave it as an exercise to prove that this is never greater than $m-1$.

Note that the removal of a key from the parent's node may cause it to underflow. Thus, the process needs to be repeated recursively. In the worst case this continues up to the root. If the root looses its only key then we simply remove the root, and make its root's only child the new root of the B-tree.
An example of deletion of 16 where $m=3$ (meaning each node has from 1 to 2 keys) is shown in the figure below. (As before, we have shown all the child links, even at the leaf level, where they are irrelevant in order to make clearer how the operation works at all levels.)


Figure 39: Deletion of key 16.

## Lecture 13: Review for Midterm

(Tue, March 13, 2001)
The midterm exam will be Thurs, Mar 15. The exam is closed-book and closed-notes, but you will be allowed one sheet of notes (front and back).
Read: You are responsible only for the material covered in class, but it is a good idea to consult Weiss's book and Samet's notes for additional insights on the concepts discussed in class. We have covered Section 1.2, Chapters 2, 3, 4, and Section 10.4.2 in Weiss and portions of Chapter 1 and Chapter 5 (Sections 5.1 and 5.2) in Samet's notes.

Overview: So far this semester we have covered some of the basic elements of data structures and search trees. Here is an overview of the major topics that we have seen. We have also discussed Java, but you are not responsible for this on the exam.

Basics of Analysis: Asymptotic notation, induction proofs, summations, and recurrences.
Basic Data Structures: Stacks, queues, deques, graphs (undirected and directed) and their representations. In particular we discussed both sequential (array) and linked allocations for these objects. For graphs and digraphs we discussed the adjacency matrix and the adjacency list representations and Prim's algorithm for computing minimum spanning trees.

Trees: We discussed multiway and binary trees and their representations. For complete trees we discussed sequential allocation (as used in the binary heap) and discussed traversals and threaded binary trees.

Search Trees: Next we discussed the dictionary abstract data type and search trees (as well as skip lists). We presented a number of data structures, each of which provided some special features.

Unbalanced binary search trees: Expected height $O(\log n)$ assuming random insertions. Expected height over a series of random insertions and deletions is $O(\sqrt{n})$, but this is due to the skew in the choice of replacement node (always selecting the minimum from the right subtree). If we randomly select min from right, and max from left, the trees will be $O(\log n)$ height on average (assuming random insertions and deletions).
AVL trees: Height balanced binary search trees. These guarantee $O(\log n)$ time for insertions, deletions, and finds. Use single and double rotations to restore balance. Because deletion is somewhat more complex, we discussed lazy deletion as an alternative approach (in which deleted nodes are just marked rather than removed).
Splay trees: Self organizing binary search trees, which are based on the splay operation, which brings a given key to the root. Splay trees guarantee $O(m \log n)$ time for a series of $m$ insertions, deletions, and finds. The tree stores no balance information or makes any structural assumptions. Any one operation could be slow, but any long series of operations will be fast. It is a good for data sets where the frequency of access of elements is nonuniform (a small fraction of elements are accessed very often), since frequently accessed elements tend to rise to the top of the tree.
Skip lists: A simple randomized data structure with $O(\log n)$ expected time for insert, delete, and find (with high probability). Unlike balanced binary trees, skip lists perform this well on average, no matter what the input looks like (in other words, there are no bad inputs, just unlucky coin tosses). They also tend to be very efficient in practice because the underlying manipulations are essentially the same as simple linked list operations.
B-trees: Widely considered the best data structure for disk access, since node size can be adjusted to the size of a disk page. Guarantees $O(\log n)$ time for dictionary operations (but is rather messy to code). The basic restructuring operations are node splits, node merges, and key-rotations.

## Lecture 14: Leftist and Skew Heaps

(Thursday, March 29, 2001)
Reading: Sections 6.6 and 6.7 in Weiss.
Leftist Heaps: The standard binary heap data structure is an simple and efficient data structure for the basic priority queue operations insert ( x ) and $\mathrm{x}=$ extractMin(). It is often the case in data structure design that the user of the data structure wants to add additional capabilities to the abstract data structure. When this happens, it may be necessary to redesign components of the data structure to achieve efficient performance.
For example, consider an application in which in addition to insert and extractMin, we want to be able to merge the contents of two different queues into one queue. As an application, suppose that a set of jobs in a computer system are separate queues waiting for the use of two resources. If one of the resources fails, we need to merge these two queues into a single queue.
We introduce a new operation $Q=$ merge (Q1, Q2), which takes two existing priority queues $Q_{1}$ and $Q_{2}$, and merges them into a new priority queue, $Q$. (Duplicate keys are allowed.) This
operation is destructive, which means that the priority queues $Q_{1}$ and $Q_{2}$ are destroyed in order to form $Q$. (Destructiveness is quite common for operations that map two data structures into a single combined data structure, since we can simply reuse the same nodes without having to create duplicate copies.)
We would like to be able to implement merge() in $O(\log n)$ time, where $n$ is the total number of keys in priority queues $Q_{1}$ and $Q_{2}$. Unfortunately, it does not seem to be possible to do this with the standard binary heap data structure (because of its highly rigid structure and the fact that it is stored in an array, without the use of pointers).
We introduce a new data structure called a leftist heap, which is fairly simple, and can provide the operations insert(x), extractMin(), and merge(Q1, Q2). This data structure has many of similar features to binary heaps. It is a binary tree which is partially ordered, which means that the key value in each parent node is less than or equal to the key values in its children's nodes. However, unlike a binary heap, we will not require that the tree is complete, or even balanced. In fact, it is entirely possible that the tree may be quite unbalanced.

Leftist Heap Property: Define the null path length, $\operatorname{npl}(v)$, of any node $v$ to be the length of the shortest path to a descendent has either 0 or 1 child. The value of $\operatorname{npl}(v)$ can be defined recursively as follows.

$$
\operatorname{npl}(v)= \begin{cases}-1 & \text { if } v=\text { null } \\ 1+\min (\operatorname{npl}(v . l e f t), \operatorname{npl}(v . \text { right })) & \text { otherwise }\end{cases}
$$

We will assume that each node has an extra field, v.npl that contains the node's npl value. The leftist property states that for every node $v$ in the tree, the npl of its left child is at least as large as the npl of its right child. We say that the keys of a tree are partially ordered if each node's key is greater than or equal to its parent's key.

Leftist heap: Is a binary tree whose keys are partially ordered (parent is less than or equal to child) and which satisfies the leftist property ( $\mathrm{npl}($ left $) \geq \mathrm{npl}($ right $)$ ). An example is shown in the figure below, where the npl values are shown next to each node.


Figure 40: Leftist heap structure (with npl values shown).
Note that any tree that does not satisfy leftist property can always be made to do so by swapping left and right subtrees at any nodes that violate the property. Observe that this does not affect the partial ordering property. Also observe that satisfying the leftist heap property does not imply that the tree is balanced. Indeed, a degenerate binary tree in which
is formed from a chain of nodes each attached as the left child of its parent does satisfy this property.
The key to the efficiency of leftist heap operations is that there exists a short $(O(\log n)$ length $)$ path in every leftist heap (namely the rightmost path). We prove the following lemma, which implies that the rightmost path in the tree cannot be of length greater than $O(\log n)$.

Lemma: A leftist heap with $r \geq 1$ nodes on its rightmost path has has at least $2^{r}-1$ nodes.
Proof: The proof is by induction on the size of the rightmost path. Before beginning the proof, we begin with two observations, which are easy to see: (1) the shortest path in any leftist heap is the rightmost path in the heap, and (2) any subtree of a leftist heap is a leftist heap. For the basis case, if there is only one node on the rightmost path, then the tree has at least one node. Since $1 \geq 2^{1}-1$, the basis case is satisfied.
For the induction step, let us suppose that the lemma is true for any leftist heap with strictly fewer than $r$ nodes on its rightmost path, and we will prove it for a binary tree with exactly $r$ nodes on its rightmost path. Remove the root of the tree, resulting in two subtrees. The right subtree has exactly $r-1$ nodes on its rightmost path (since we have eliminated only the root), and the left subtree must have a at least $r-1$ nodes on its rightmost path (since otherwise the rightmost path in the original tree would not be the shortest, violating (1)). Thus, by applying the induction hypothesis, it follows that the right and left subtrees have at least $2^{r-1}-1$ nodes each, and summing them, together with the root node we get a total of at least

$$
2\left(2^{r-1}-1\right)+1=2^{r}-1
$$

nodes in the entire tree.
Leftist Heap Operations: The basic operation upon which leftist heaps are based is the merge operation. Observe, for example, that both the operations insert() and extractMin() can be implemented by using the operation merge(). (Note the similarity with splay trees, in which all operations were centered around the splay operation.)
The basic node of the leftist heap is a LeftHeapNode. Each such node contains an data field of type Element (upon which comparisons can be made) a left and right child, and an npl value. The constructor is given each of these values. The leftist heap class consists of a single root, which points to the root node of the heap. Later we will describe the main procedure merge (LeftHeapNode h1, LeftHeapNode h2), which merges the two (sub)heaps rooted at $h_{1}$ and $h_{2}$. For now, assuming that we have this operation we define the main heap operations. Recall that merge is a destructive operation.

Leftist Operations

```
    void insert(Element x) {
        root = merge(root, new LeftHeapNode(x, null, null, 0))
    }
    Element extractMin() {
        if (root == null) return null // empty heap
        Element minItem = root.data // minItem is root's element
        root = merge(root.left, root.right)
return minItem
    }
    void merge(LeftistHeap Q1, LeftistHeap Q2) {
        root = merge(Q1.root, Q2.root)
    }
```

Leftist Heap Merge: All that remains, is to show how merge(h1, h2) is implemented. The formal description of the procedure is recursive. However it is somewhat easier to understand in its nonrecursive form. Let $h_{1}$ and $h_{2}$ be the two leftist heaps to be merged. Consider the rightmost paths of both heaps. The keys along each of these paths form an increasing sequence. We could merge these paths into a single sorted path (as in merging two lists of sorted keys in the merge-sort algorithm). However the resulting tree might not satisfy the leftist property. Thus we update the npl values, and swap left and right children at each node along this path where the leftist property is violated. A recursive implementation of the algorithm is given below. It is essentially the same as the one given in Weiss, with some minor coding changes.

Merge two leftist heaps

```
LeftHeapNode merge(LeftHeapNode h1, LeftHeapNode h2) {
    if (h1 == null) return h2 // if one is empty, return the other
    if (h2 == null) return h1
    if (h1.data > h2.data) // swap so h1 has the smaller root
        swap(h1, h2)
    if (h1.left == null) // h1 must be a leaf in this case
        h1.left = h2
    else { // merge h2 on right and swap if needed
        h1.right = merge(h1.right, h2)
        if (h1.left.npl < h1.right.npl) {
                swap(h1.left, h1.right) // swap children to make leftist
        }
        h1.npl = h1.right.npl + 1 // update npl value
    }
    return h1
}
```

For the analysis, observe that because the recursive algorithm spends $O(1)$ time for each node on the rightmost path of either $h_{1}$ or $h_{2}$, the total running time is $O(\log n)$, where $n$ is the total number of nodes in both heaps.

This recursive procedure is a bit hard to understand. A somewhat easier 2-step interpretation is as follows. First, merge the two right paths, and update the npl values as you do so. Second, walk back up along this path and swap children whenever the leftist condition is violated. The figure below illustrates this way of thinking about the algorithm. The recursive algorithm just combines these two steps, rather than making the separate.

Skew Heaps: Recall that a splay tree is a self-adjusting balanced binary tree. Unlike the AVL tree, which maintains balance information at each node and performs rotations to maintain balance, the splay tree performs a series of rotations to continuously "mix-up" the tree's structure, and hence achieve a sort of balance from an amortized perspective.

A skew heap is a self-organizing heap, and applies this same idea as in splay trees, but to the leftist heaps instead. As with splay trees we store no balance information with each node (no npl values). We perform the merge operation as described above, but we swap the left and right children from every node along the rightmost path. (Note that in the process the old rightmost path gets flipped over to become the leftmost path.) An example is given in Weiss, in Section 6.7.

It can be shown that if the heap contains at most $n$ elements, then a sequence of $m$ heap operations (insert, extractMin, merge) starting from empty heaps takes $O(m \log n)$ time. Thus


Figure 41: Example of merging two leftist heaps (2-step interpretation.
the average time per operation is $O(\log n)$. Because no balance information is needed, the code is quite simple.

## Lecture 15: Disjoint Set Union-Find

(Tuesday, April 3, 2001)
Reading: Weiss, Chapt 8 through Sect 8.6.
Equivalence relations: An equivalence relation over some set $S$ is a relation that satisfies the following properties for all elements of $a, b, c \in S$.

Reflexive: $a \equiv a$.
Symmetric: $a \equiv b$ then $b \equiv a$
Transitive: $a \equiv b$ and $b \equiv c$ then $a \equiv c$.
Equivalence relations arise in numerous applications. An example includes any sort of "grouping" operation, where every object belongs to some group (perhaps in a group by itself) and no object belongs to more than one group. More formally these groups are called equivalent classes and the subdivision of the set into such classes is called a partition. For example, suppose we are maintaining a bidirectional communication network. The ability to communicate is an equivalence relation, since if machine $a$ can communicate with machine $b$, and machine $b$ can communicate with machine $c$, then machine $a$ can communicate with machine $c$ (e.g. by sending messages through $b$ ). Now suppose that a new link is created between two groups, which previously were unable to communicate. This has the effect of merging two equivalence classes into one class.
We discuss a data structure that can be used for maintaining equivalence partitions with two operations: (1) union, merging to groups together, and (2) find, determining which group an element belongs to. This data structure should not be thought of as a general purpose data structure for storing sets. In particular, it cannot perform many important set operations, such as splitting two sets, or computing set operations such as intersection and complementation. And its structure is tailored to handle just these two operations. However, there are many applications for which this structure is useful. As we shall see, the data structure is simple and amazingly efficient.

Union-Find ADT: We assume that we have an underlying finite set of elements $S$. We want to maintain a partition of the set. In addition to the constructor, the (abstract) data structure supports the following operations.

Set $\mathrm{s}=\mathrm{find}($ Element x$):$ Return an set identifier of the set $s$ that contains the element $x$. A set identifier is simply a special value (of unspecified type) with the property that find $(x)==\operatorname{find}(y)$ if and only if $x$ and $y$ are in the same set.
Set $\mathrm{r}=\mathrm{union}($ Set s , Set t ): Merge two sets named $s$ and $t$ into a single set $r$ containing their union. We assume that $s, t$ and $r$ are given as set identifiers. This operations destroys the sets $s$ and $t$.

Note that there are no key values used here. The arguments to the find and union operations are pointers to objects stored in the data structure. The constructor for the data structure is given the elements in the set $S$ and produces a structure in which every element $x \in S$ is in a singleton set $\{x\}$ containing just $x$.

Inverted Tree Implementation: We will derive our implementation of a data structure for the union-find ADT by starting with a simple structure based on a forest of inverted trees. You think of an inverted tree as a multiway tree in which we only store parent links (no child or sibling links). The root's parent pointer is null. There is no limit on how many children a node can have. The sets are represented by storing the elements of each set in separate tree.
In this implementation a set identifier is simply a pointer to the root of an inverted tree. Thus the type Set is just an alias for the type Element (which is perhaps not very good practice, but is very convenient). We will assume that each element $x$ of the set has a pointer x.parent, which points to its parent in this tree. To perform the operation find( x ), we walk along parent links and return the root of the tree. This root element is set identifier for the set containing $x$. Notice that this satisfies the above requirement, since two elements are in the same set if and only if they have the same root. We call this find1(). Later we will propose an improvement

```
Set find1(Element x) {
    while (x.parent != null) x = x.parent; // follow chain to root
    return x; // return the root
}
```

For example, suppose that $S=\{a, b, c, \ldots, m\}$ and the current partition is:

$$
\{a, f, g, h, k, l\},\{b\},\{c, d, e, m\},\{i, j\} .
$$

This might be stored in an inverted tree as shown in the following figure. The operation find ( $k$ ) would return a pointer to node $g$.


Figure 42: Union-find Tree Example.
Note that there is no particular order to how the individual trees are structured, as long as they contain the proper elements. (Again recall that unlike existing data structures that we have discussed in which there are key values, here the arguments are pointers to the nodes of the inverted tree. Thus, there is never a question of "what" is in the data structure, the issue is "which" tree are you in.) Initially each element is in its own set. To do this we just set all parent pointers to null.
A union is straightforward to implement. To perform the union of two sets, e.g. to take the union of the set containing $\{b\}$ with the set $\{i, j\}$ we just link the root of one tree into the root of the other tree. But here is where it pays to be smart. If we link the root of $\{i, j\}$ into $\{b\}$, the height of the resulting tree is 2 , whereas if we do it the other way the height of the tree is only 1. (Recall that height of a tree is the maximum number of edges from any leaf to the root.) It is best to keep the tree's height small, because in doing so we make the find's run faster.

In order to perform union's intelligently, we maintain an extra piece of information, which records the height of each tree. For historic reasons we call this height the rank of the tree. We assume that the rank is stored as a field in each element. The intelligent rule for doing unions is to link the root of the set of smaller rank as a child of the root of the set of larger rank.

Observe will only need the rank information for each tree root, and each tree root has no need for a parent pointer. So in a really tricky implementation, the ranks and parent pointers can share the same storage field, provided that you have some way of distinguishing them. For example, in our text, parent pointers (actually indices) are stored with positive integers and ranks are stored as negative integers. We will not worry about being so clever.


Figure 43: Union-find with ranks.
The code for union operation is shown in the following figure. When updating the ranks there are two cases. If the trees have the same ranks, then the rank of the result is one larger than this value. If not, then the rank of the result is the same as the rank of the higher ranked tree. (You should convince yourself of this.)
_Union operation

```
Set union(Set s, Set t) {
    if (s.rank > t.rank) { // s has strictly higher rank
        t.parent = s // make s the root (rank does not change)
        return s
    }
    else { // t has equal or higher rank
        if (s.rank == t.rank) t.rank++
        s.parent = t // make t the root
        return t
    }
}
```

Analysis of Running Time: Consider a sequence of $m$ union-find operations, performed on a domain with $n$ total elements. Observe that the running time of the initialization is proportional to $n$, the number of elements in the set, but this is done only once. Each union takes only constant time, $O(1)$.

In the worst case, find takes time proportional to the height of the tree. The key to the efficiency of this procedure is the following observation, which implies that the tree height is never greater than $\lg m$. (Recall that $\lg m$ denotes the logarithm base 2 of $m$.)

Lemma: Using the union-find procedures described above any tree with height $h$ has at least $2^{h}$ elements.

Proof: Given a union-find tree $T$, let $h$ denote the height of $T$, let $n$ denote the number of elements in $T$, and let $u$ denote the number of unions needed to build $T$. We prove the lemma by strong induction on the number of unions performed to build the tree. For the basis (no unions) we have a tree with 1 element of height 0 . Since $2^{0}=1$, the basis case is established.
For the induction step, assume that the hypothesis is true for all trees built with strictly fewer than $u$ union operations, and we want to prove the lemma for a union-find tree built with exactly $u$ union operations. Let $T$ be such a tree. Let $T^{\prime}$ and $T^{\prime \prime}$ be the two subtrees that were joined as part of the final union. Let $h^{\prime}$ and $h^{\prime \prime}$ be the heights of $T^{\prime}$ and $T^{\prime \prime}$, respectively, and define $n^{\prime}$ and $n^{\prime \prime}$ similarly. Each of $T^{\prime}$ and $T^{\prime \prime}$ were built with strictly fewer than $u$ union operations. By the induction hypothesis we have

$$
n^{\prime} \geq 2^{h^{\prime}} \quad \text { and } \quad n^{\prime \prime} \geq 2^{h^{\prime \prime}}
$$

Let us assume that $T^{\prime}$ was made the child of $T^{\prime \prime}$ (the other case is symmetrical). This implies that $h^{\prime} \leq h^{\prime \prime}$, and $h=\max \left(h^{\prime}+1, h^{\prime \prime}\right)$. There are two cases. First, if $h^{\prime}=h^{\prime \prime}$ then $h=h^{\prime}+1$ and we have

$$
n=n^{\prime}+n^{\prime \prime} \geq 2^{h^{\prime}}+2^{h^{\prime \prime}}=2^{h^{\prime}+1}=2^{h}
$$

Second, if $h^{\prime}<h^{\prime \prime}$ then $h=h^{\prime \prime}$ and we have

$$
n=n^{\prime}+n^{\prime \prime} \geq n^{\prime \prime} \geq 2^{h^{\prime \prime}}=2^{h}
$$

In either case we get the desired result.
Since the unions's take $O(1)$ time each, we have the following.
Theorem: After initialization, any sequence of $m$ union's and find's can be performed in time $O(m \log m)$.

Path Compression: It is possible to apply a very simple heuristic improvement to this data structure which provides a significant improvement in the running time. Here is the intuition. If the user of your data structure repeatedly performs find's on a leaf at a very low level in the tree then each such find takes as much as $O(\log n)$ time. Can we improve on this?
Once we know the result of the find, we can go back and "short-cut" each pointer along the path to point directly to the root of the tree. This only increases the time needed to perform the find by a constant factor, but any subsequent find on this node (or any of its ancestors) will take only $O(1)$ time. The operation of short-cutting the pointers so they all point directly to the root is called path-compression and an example is shown below. Notice that only the pointers along the path to the root are altered. We present a slick recursive version below as well. Trace it to see how it works.

```
Set find2(Element x) {
    if (x.parent == null) return x // return root
    else return x.parent = find2(x.parent) // find root and update parent
}
```

The running time of find2 is still proportional to the depth of node being found, but observe that each time you spend a lot of time in a find, you flatten the search path. Thus the work you do provides a benefit for later find operations. (This is the sort of thing that we observed in earlier amortized analyses.)


Figure 44: Find using path compression.

Does the savings really amount to anything? The answer is yes. It was actually believed at one time that if path compression is used, then (after initialization) the running time of the algorithm for a sequence of $m$ union and finds was $O(m)$, and hence the amortized cost of each operation is $O(1)$. However, this turned out to be false, but the amortized time is much less than $O(\log m)$.
Analyzing this algorithm is quite tricky. (Our text gives the full analysis if you are interested.) In order to create a bad situation you need to do lots of unions to build up a tree with some real depth. But as soon as you start doing finds on this tree, it very quickly becomes very flat again. In the worst case we need to an immense number of unions to get high costs for the finds.

To give the analysis (which we won't prove) we introduce two new functions, $A(m, n)$ and $\alpha(n)$. The function $A(m, n)$ is called Ackerman's function. It is famous for being just about the fastest growing function imaginable.

$$
\begin{aligned}
A(1, j) & =2^{j} \quad \text { for } j \geq 1 \\
A(i, 1) & =A(i-1,2) \quad \text { for } i \geq 2 \\
A(i, j) & =A(i-1, A(i, j-1)) \quad \text { for } i, j \geq 2
\end{aligned}
$$

In spite of its innocuous appearance, this function is a real monster. To get a feeling for how fast this function grows, observe that

$$
A(2, j)=2^{2^{2}}
$$

where the tower of powers of 2 is $j$ high. (Try expanding the recurrence to prove this.) Hence,

$$
A(4,1)=A(3,2)=A(2, A(3,1))=A(2, A(2,2))=A\left(2,2^{2^{2}}\right)=A(2,16) \approx 10^{80}
$$

which is already greater than the estimated number of atoms in the observable universe.
Since Ackerman's function grows so fast, its inverse, called $\alpha$ grows incredibly slowly. Define

$$
\alpha(m, n)=\min \{i \geq 1 \mid A(i,\lfloor m / n\rfloor)>\lg n\}
$$

This definition is somewhat hard to interpret, but the important bottom line is that assuming $\lfloor m / n\rfloor \geq 1$, we have $\alpha(m, n) \leq 4$ as long as $m$ is less than the number of atoms in the universe
(which is certainly true for most input sets!) The following result (which we present without proof) shows that an sequence of union-find operations take amortized time $O(\alpha(m, n)$ ), and so each operation (for all practical purposes) takes amortized constant time.

Theorem: After initialization, any sequence of $m$ union's and find's (using path compression) on an initial set of $n$ elements can be performed in time $O(m \alpha(m, n))$ time. Thus the amortized cost of each union-find operation is $O(\alpha(m, n))$.

## Lecture 16: Geometric Preliminaries

(Thursday, April 5, 2001)
Today's material is not discussed in any of our readings.
Geometric Information: A large number of modern data structures and algorithms problems involve geometric data. The reason is that rapidly growing fields such as computer graphics, robotics, computer vision, computer-aided design, visualization, human-computer interaction, virtual reality, and others deal primarily with geometric data. Geometric applications give rise to new data structures problems, which we will be studying for a number of lectures.
Before discussing geometric data structures, we need to provide some background on what geometric data is, and how we compute with it. With nongeometric data we stated that we are storing records, and each record is associated with an identifying key value. These key values served a number of functions for us. They served a function of identification the the objects of the data structure. Because of the underlying ordering relationship, they also provided a means for searching for objects in the data structure, by giving us a way to direct the search by subdividing the space of keys into subsets that are greater than or less than the key.
In geometric data structures we will need to generalize the notion of a key. A geometric point object in the plane can be identified by its $(x, y)$ coordinates, which can be thought of as a type of key value. However, if we are storing more complex objects, such as rectangles, line segments, spheres, and polygons, the notion of a identifying key is not as appropriate. As with one-dimensional data, we also have associated application data. For example, in a graphics system, the geometric description of an object is augmented with information about the object's color, texture, and its surface reflectivity properties. Since our interest will be in storing and retrieving objects based on their geometric properties, we will not discuss these associated data values.

Primitive Objects: Before we start discussing geometric data structures, we will digress to discuss a bit about geometric objects, their representation, and manipulation. Here is a list of common geometric objects and possible representations. The list is far from complete. Let $R^{d}$ denote $d$-dimensional space with real coordinates.

Scalar: This is a single (1-dimensional) real number. It is represented as float or double.
Point: Points are locations in space. A typical representation is as a d-tuple of scalars, e.g. $P=\left(p_{0}, p_{1}, \ldots, p_{d-1}\right) \in R^{d}$. Is it better to represent a point as an array of scalars, or as an object with data members labeled $x, y$, and $z$ ? The array representation is more general and often more convenient, since it is easier to generalized to higher dimensions and coordinates can be parameterized. You can define "names" for the coordinates. If the dimension might vary then the array representation is necessary.
For example, in Java we might represent a point in 3-space using something like the following. (Note that "static final" is Java's way of saying "const".)

```
class Point {
    public static final int DIM = 3;
    protected float coord[DIM];
}
```

Vector: Vectors are used to denote direction and magnitude in space. Vectors and points are represented in essentially the same way, as a $d$-tuple of scalars, $\vec{v}=\left(v_{0}, v_{1}, \ldots, v_{d-1}\right)$. It is often convenient to think of vectors as free vectors, meaning that they are not tied down to a particular origin, but instead are free to roam around space. The reason for distinguishing vectors from points is that they often serve significantly different functions. For example, velocities are frequently described as vectors, but locations are usually described as points. This provides the reader of your program a bit more insight into your intent.
Line Segment: A line segment can be represented by giving its two endpoints ( $p_{1}, p_{2}$ ). In some applications it is important to distinguish between the line segments $\overrightarrow{p_{1} p_{2}}$ and $\overrightarrow{p_{2} p_{1}}$. In this case they would be called directed line segments.
Ray: Directed lines in 3- and higher dimensions are not usually represented by equations but as rays. A ray can be represented by storing an origin point $P$ and a nonzero directional vector $\vec{v}$.


Figure 45: Basic geometric objects.

Line: A line in the plane can be represented by a line equation

$$
y=a x+b \quad \text { or } \quad a x+b y=c
$$

The former definition is the familiar slope-intercept representation, and the latter takes an extra coefficient but is more general since it can easily represent a vertical line. The representation consists of just storing the pair $(a, b)$ or $(a, b, c)$.
Another way to represent a line is by giving two points through which the line passes, or by giving a point and a directional vector. These latter two methods have the advantage that they generalize to higher dimensions.
Hyperplanes and halfspaces: In general, in dimension $d$, a linear equation of the form

$$
a_{0} p_{0}+a_{1} p_{1}+\ldots+a_{d-1} p_{d-1}=c
$$

defines a $d-1$ dimensional hyperplane. It can be represented by the $d$-tuple ( $a_{0}, a_{1}, \ldots, a_{d-1}$ ) together with $c$. Note that the vector $\left(a_{0}, a_{1}, \ldots, a_{d-1}\right)$ is orthogonal to the hyperplane. The set of points that lie on one side or the other of a hyperplane is called a halfspace. The formula is the same as above, but the " $=$ " is replaced with an inequality such as " $<$ " or " $\geq$ ".


Figure 46: Planes and Halfspaces.

Orthogonal Hyperplane: In many data structures it is common to use hyperplanes that are orthogonal to one of the coordinate axes, called orthogonal hyperplanes. In this case it is much easier to store such a hyperplane by storing (1) an integer cutDim from the set $\{0,1, \ldots, d-1\}$, which indicates which axis the plane is perpendicular to and (2) a scalar cutVal, which indicates where the plane cuts this axis.
Simple Polygon: Solid objects can be represented as polygons (in dimension 2) and polyhedra (in higher dimensions). A polygon is a cycle of line segments joined end-to-end. It is said to be simple if its boundary does not self-intersect. It is convex if it is simple and no internal angle is greater than $\pi$.


Figure 47: Polygons, rectangles and spheres.
The standard representation of a simple polygon is just a list of points (stored either in an array or a circularly linked list). In general representing solids in higher dimensions involves more complex structures, which we will discuss later.
Orthogonal Rectangle: Rectangles can be represented as polygons, but rectangles whose sides are parallel to the coordinate axes are common. A couple of representations are often used. One is to store the two points of opposite corners (e.g. lower left and upper right).
Circle/Sphere: A $d$-dimensional sphere, can be represented by point $P$ indicating the center of the sphere, and a positive scalar representing its radius. A points $X$ lies within the sphere if

$$
\left(x_{0}-p_{0}\right)^{2}+\left(x_{1}-p_{1}\right)^{2}+\ldots+\left(x_{d-1}-p_{d-1}\right)^{2} \leq r^{2} .
$$

Topological Notions: When discussing geometric objects such as circles and polygons, it is often important to distinguish between what is inside, what is outside and what is on the boundary. For example, given a triangle $T$ in the plane we can talk about the points that are in the interior $(\operatorname{int}(T))$ of the triangle, the points that lie on the boundary (bnd $(T))$ of the triangle, and the points that lie on the exterior $(\operatorname{ext}(T))$ of the triangle. A set is closed if it includes its boundary, and open if it does not. Sometimes it is convenient to define a set as being
semi-open, meaning that some parts of the boundary are included and some are not. Making these notions precise is tricky, so we will just leave this on an intuitive level.


Figure 48: Topological notation.
Operations on Primitive Objects: When dealing with simple numeric objects in 1-dimensional data structures, the set of possible operations needed to be performed on each primitive object was quite simple, e.g. compare one key with another, add or subtract keys, print keys, etc. With geometric objects, there are many more operations that one can imagine.

Basic point/vector operators: Let $\alpha$ be any scalar, let $P$ and $Q$ be points, and $\vec{u}, \vec{v}, \vec{w}$ be vectors. We think of vectors as being free to roam about the space whereas points are tied down to a particular location. We can do all the standard operations on vectors you learned in linear algebra (adding, subtracting, etc.) The difference of two points $P-Q$ results in the vector directed from $Q$ to $P$. The sum of a point $P$ and vector $\vec{u}$ is the point lying on the head of the vector when its tail is placed on $P$.

vector addition

point subtraction

point-vector addition

Figure 49: Point-vector operators.
As a programming note, observe that $\mathrm{C}++$ has a very nice mechanism for handling these operations on Points and Vectors using operator overloading. Java does not allow overloading of operators.
Affine combinations: Given two points $P$ and $Q$, the point $(1-\alpha) P+\alpha Q$ is point on the line joining $P$ and $Q$. We can think of this as a weighted average, so that as $\alpha$ approaches 0 the point is closer to $P$ and as $\alpha$ approaches 1 the point is closer to $Q$. This is called an affine combination of $P$ and $Q$. If $0 \leq \alpha \leq 1$, then the point lies on the line segment $\overline{P Q}$.
Length and distance: The length of a vector $v$ is defined to be

$$
\|\vec{v}\|=\sqrt{v_{0}^{2}+v_{1}^{2}+\ldots+v_{d-1}^{2}} .
$$

The distance between two points $P$ and $Q$ is the length of the vector between them, that is

$$
\operatorname{dist}(P, Q)=\|P-Q\| .
$$

Computing distances between other geometric objects is also important. For example, what is the distance between two triangles? When discussing complex objects, distance usually means the closest distance between objects.


Figure 50: Affine combinations.

Orientation and Membership: There are a number of geometric operations that deal with the relationship between geometric objects. Some are easy to solve. (E.g., does a point $P$ lie within a rectangle $R$ ?) Some are harder. (E.g., given points $A, B, C$, and $D$, does $D$ lie within the unique circle defined by the other three points?) We will discuss these as the need arises.
Intersections: The other sort of question that is important is whether two geometric objects intersect each other. Again, some of these are easy to answer, and others can be quite complex.

Example: Circle/Rectangle Intersection: As an example of a typical problem involving geometric primitives, let us consider the question of whether a circle in the plane intersects a rectangle. Let us represent the circle by its center point $C$ and radius $r$ and represent the rectangle $R$ by its lower left and upper right corner points, $R_{l o}$ and $R_{h i}$ (for low and high).
Problems involving circles are often more easily recast as problems involving distances. So, instead, let us consider the problem of determining the distance $d$ from $C$ to its closest point on the rectangle $R$. If $d>r$ then the circle does not intersect the rectangle, and otherwise it does.
In order to determine the distance from $C$ to the rectangle, we first observe that if $C$ lies inside $R$, then the distance is 0 . Otherwise we need to determine which point of the boundary of $R$ is closest to $C$. Observe that we can subdivide the exterior of the rectangle into 8 regions, as shown in the figure below. If $C$ lies in one of the four corner regions, then $C$ is closest to the corresponding vertex of $R$ and otherwise $C$ is closest to one of the four sides of $R$. Thus, all we need to do is to classify which region $C$ lies is, and compute the corresponding distance. This would usually lead a lengthy collection of if-then-else statements, involving 9 different cases.


Figure 51: Distance from a point to a rectangle, and squared distance contributions for each region.
We will take a different approach. Sometimes things are actually to code if you consider the problem in its general $d$-dimensional form. Rather than computing the distance, let us first
concentrate on computing the squared distance instead. The distance is the sum of the squares of the distance along the $x$-axis and distance along the $y$-axis. Consider just the $x$-coordinates, if $C$ lies to are to the left of the rectangle then the contribution is $\left(R_{l o, x}-c_{x}\right)^{2}$, if $C$ lies to the right then the contribution is $\left(c_{x}-R_{h i, x}\right)^{2}$, and otherwise there is no $x$-contribution to the distance. A similar analysis applies for $y$. This suggests the following code, which works in all dimensions.

Distance from Point $C$ to Rectangle $R$

```
float distance(Point C, Rectangle R) {
    sumSq = 0 // sum of squares
    for (int i = 0; i < Point.DIM; i++) {
        if (C[i] < R.lo[i]) // left of rectangle
            sumSq += square(R.lo[i] - C[i])
        else if (C[i] > R.hi[i]) // right of rectangle
            sumSq += square(C[i] - R.hi[i])
    }
    return sqrt(sumSq)
}
```

Finally, once the distance has been computed, we test whether it is less than the radius $r$. If so the circle intersects the rectangle and otherwise it does not.

## Lecture 17: Quadtrees and kd-trees

(Tuesday, April 10, 2001)
Reading: Today's material is discussed in Chapt. 2 of Samet's book on spatial data structures.
Geometric Data Structures: Geometric data structures are based on many of the concepts presented in typical one-dimensional (i.e. single key) data structures. There are a few reasons why generalizing 1-dimensional data structures to multi-dimensional data structures is not always straightforward. The first is that most 1-dimensional data structures are built in one way or another around the notion of sorting the keys. However, what does it mean to sort 2-dimensional data? (Obviously your could sort lexicographically, first by $x$-coordinate and then by $y$-coordinate, but this is a rather artificial ordering and does not convey any of the 2-dimensional structure.) The idea of sorting and binary search is typically replaced with a more general notion hierarchical subdivision, namely that of partitioning space into local regions using a tree structure. Let $S$ denote a set of geometric objects (e.g. points, line segments, lines, sphere, rectangles, polygons). In addition to the standard operations of inserting and deleting elements from $S$, the following are example of common queries.

Find: Is a given object in the set?
Nearest Neighbor: Find the closest object (or the closest $k$ objects) to a given point.
Range Query: Given a region (e.g. a rectangle, triangle, or circle) list (or count) all the objects that lie entirely/partially inside of this region.
Ray Shooting: Given a ray, what is the first object (if any) that is hit by this ray.
As with 1-dimensional data structures, the goal is to store our objects in some sort of intelligent way so that queries of the above form can be answered efficiently. Notice with 1-dimensional data it was important for us to consider the dynamic case (i.e. objects are being inserted and deleted). The reason was that the static case can always be solved by simply sorting the objects. Note that in higher dimensions sorting is not an option, so solving these problems is nontrivial, even for static sets.

Point quadtree: We first consider a very simple way of generalizing unbalanced binary trees in the 1-dimensional case to 4 -ary trees in the two dimensional case. The resulting data structure is called a point quadtree. As with standard binary trees, the points are inserted one by one. The insertion of each point results in a subdivision of a rectangular region into four smaller rectangles, by passing horizontal and vertical splitting lines through the new point. Consider the insertion of the following points:

$$
(35,40),(50,10),(60,75),(80,65),(85,15),(5,45),(25,35),(90,5)
$$

The result is shown in the following figure. Each node has four children, labeled NW, NE, SW, and SE (corresponding to the compass directions). The tree is shown on the left and the associated spatial subdivision is shown on the right. Each node in the tree is associated with a rectangular cell. Note that some rectangles are special in that they extend to infinity. Since semi-infinite rectangles sometimes bother people, it is not uncommon to assume that everything is contained within one large bounding rectangle, which may be provided by the user when the tree is first constructed. Once you get used to the mapping between the tree and the spatial subdivision, it is common just to draw the subdivision and omit the tree.


Figure 52: Point quadtree.
We will not discuss algorithms for the point quad-tree in detail. Instead we will defer this discussion to kd-trees, and simply note that for each operation on a kd-tree, there is a similar algorithm for quadtrees.

Point kd-tree: Point quadtrees can be generalized to higher dimensions, but observe that for each splitting point the number of children of a node increases exponentially with dimension. In dimension $d$, each node has $2^{d}$ children. Thus, in 3 -space the resulting trees have 8 children each (called point octrees). The problem is that it becomes hard to generalize algorithms from one dimension to another, and if you are in a very high dimensional space (e.g. 20 dimensional space) every node has $2^{20}$ (or roughly a million) children. Obviously just storing all the null pointers in a single node would be a ridiculous exercise. (You can probably imagine better ways of doing this, but why bother?)
An alternative is to retain the familiar binary tree structure, by altering the subdivision method. As in the case of a quadtree the cell associated with each node is associated with a rectangle (assuming the planar case) or a hyper-rectangle in $d$-dimensional space. When a new point is inserted into some node (equivalently into some cell) we split the cell by a horizontal or vertical splitting line, which passes through this point. In higher dimensions, we split the cell by a $(d-1)$ dimensional hyperplane that is orthogonal to one of the coordinate axes. In any dimension, the split can be specified by giving the splitting axes ( $x, y$ or whatever), also
called the cutting dimension. The associated splitting value is the associated coordinate of the data point.
The resulting data structure is called a point $k d$-tree. Actually, this is a bit of a misnomer. The data structure was named by its inventor Jon Bentley to be a 2-d tree in the plane, a 3-d tree in 3 -space, and a $k$ - $d$ tree in dimension $k$. However, over time the name "kd-tree" became commonly used irrespective of dimension. Thus it is common to say a "kd-tree in dimension 3 " rather than a " 3 -d tree".
How is the cutting dimension chosen? There are a number of ways, but the most common is just to alternate among the possible axes at each new level of the tree. For example, at the root node we cut orthogonal to the $x$-axis, for its children we cut orthogonal to $y$, for the grandchildren we cut along $x$, and so on. In general, we can just cycle through the various axes. In this way the cutting dimension need not be stored explicitly in the node, but can be determined implicitly as the tree is being traversed. An example is shown below (given the same points as in the previous example). Again we show both the tree and the spatial subdivision.


Figure 53: Point kd-tree decomposition.
The tree representation is much the same as it is for quad-trees except now every node has only two children. The left child holds all the points that are less than the splitting point along the cutting dimension, and the right subtree holds those that are larger. What should we do if the point coordinate is the same? A good solution is just to pick one subtree (e.g. the right) and store it there.
As with unbalanced binary search trees, it is possible to prove that if keys are inserted in random order, then the expected height of the tree is $O(\log n)$, where $n$ is the number of points in the tree.

Insertion into kd-trees: Insertion operates as it would for a regular binary tree. The code is given below. The initial call is root $=\operatorname{insert}(\mathrm{x}$, root, 0 ). We assume that the dimensions are indexed $0,1, \ldots, D I M-1$. The parameter $c d$ keeps track of the current cutting dimension. We advance the cutting dimension each time we descend to the next level of the tree, by adding one to its value modulo DIM. (For example, in the plane this value alternates between 0 and 1 as we recurse from level to level.) We assume that the + operator has been overloaded for the cutting dimension to operate modulo $D I M$. Thus, when we say $c d+1$ in the pseudocode, this a shorthand for (cd+1) \% DIM. Notice that if a point has the same coordinate value as the cutting dimension it is always inserted in the right subtree.

```
    if (t == null) // fell out of tree
        t = new KDNode(x)
    else if (x == t.data) // duplicate data point?
        ...error duplicate point...
    else if (x[cd] < t.data[cd]) // left of splitting line
    t.left = insert(x, t.left, cd+1)
    else // on or right of splitting line
    t.right = insert(x, t.right, cd+1)
    return t
}
```

FindMin and FindMax: We will discuss deletion from kd-trees next time. As a prelude to this discussion, we will discuss an interesting query problem, which is needed as part of the deletion operation. The problem is, given a node $t$ in a kd-tree, and given the index of a coordinate axis $i$ we wish to find the point in the subtree rooted at $t$ whose $i$-th coordinate is minimum. This is a good example of a typical type of query problem in kd-trees, because it shows the fundamental feature of all such query algorithms, namely to recurse only on subtrees that might lead to a possible solution.
For the sake of concreteness, assume that the coordinate is $i=0$ (meaning the $x$-axis in the plane). The procedure operates recursively. When we arrive at a node $t$ if the cutting dimension for $t$ is the $x$-axis, then we know that the minimum coordinate cannot be in $t$ 's right subtree. If the left subtree is nonnull, we recursively search the left subtree and otherwise we return $t$ 's data point. On the other hand, if $t$ 's cutting dimension is any other axis (the $y$-axis say) then we cannot know whether the minimum coordinate will lie in its right subtree or its left subtree, or might be $t$ 's data point itself. So, we recurse on both the left and right subtrees, and return the minimum of these and $t$ 's data point. We assume we have access to a function minimum ( $\mathrm{p} 1, \mathrm{p} 2, \mathrm{p} 3, \mathrm{i}$ ) which returns the point reference $p_{1}, p_{2}$, or $p_{3}$ which is nonnull and minimum on coordinate $i$. The code and a figure showing which nodes are visited are presented below.

```
    // find minimum along axis i
Point findMin(KDNode t, int i, int cd) {
    if (t == null) return null // fell out of tree
    if (cd == i) { // i == splitting axis
        if (t.left == null) return t.data // no left child?
        else return findMin(t.left, i, cd+1)
    }
    else { // else min of both children and t
        return minimum(t.data,
            findMin(t.left, i, cd+1),
            findMin(t.right, i, cd+1),
            i)
    }
}
```


## Lecture 18: More on kd-trees

(Thursday, April 12, 2001)
Reading: Today's material is discussed in Chapt. 2 of Samet's book on spatial data structures.


Figure 54: Example of findMin() for x-axis.

Deletion from a kd-tree: As with insertion we can model the deletion code after the deletion code for unbalanced binary search trees. However, there are some interesting twists here. Recall that in the 1-dimensional case we needed to consider a number of different cases. If the node is a leaf we just delete the node. Otherwise, its deletion would result in a "hole" in the tree. We need to find an appropriate replacement element. In the 1-dimensional case, we were able to simplify this if the node has a single child (by making this child the new child of our parent). However, this would move the child from an even level to an odd level, or vice versa. This would violate the entire structure of the kd-tree. Instead, in both the 1-child and 2 -child cases we will need to find a replacement node, from whichever subtree we can.
Let us assume first that the right subtree is non-empty. Recall that in the 1-dimensional case, the replacement key was taken to be the smallest key from the right child, and after using the replacement to fill the hole, we recursively deleted the replacement. How do we generalize this to the multi-dimensional case? What does it mean to find the "smallest" element in a such a set? The right thing to do is to find the point whose coordinate along the current cutting dimension is minimum. Thus, if the cutting dimension is the $x$-axis, say, then the replacement key is the point with the smallest $x$-coordinate in the right subtree. We use the findMin() function (given last time) to do this for us.

What do we do if the right subtree is empty? At first, it might seem that the right thing to do is to select the maximum node from the left subtree. However, there is a subtle trap here. Recall that we maintain the invariant that points whose coordinates are equal to the cutting dimension are stored in the right subtree. If we select the replacement point to be the point with the maximum coordinate from the left subtree, and if there are other points with the same coordinate value in this subtree, then we will have violated our invariant. There is a clever trick for getting around this though. For the replacement element we will select the minimum (not maximum) point from the left subtree, and we move the left subtree over and becomes the new right subtree. The left child pointer is set to null. The code and an example are given below. As before recall that cd+1 means that we increment the cutting dimension modulo the dimension of the space.
kd-tree Deletion

```
KDNode delete(Point x, KDNode t, int cd) {
    if (t == null) // fell out of tree
            ...error deletion of nonexistant point...
    else if (x == t.data) { // found it
            if (t.right != null) { // take replacement from right
                t.data = findMin(t.right, cd, cd+1)
                t.right = delete(t.data, t.right, cd+1)
```

```
        }
        else if (t.left != null) { // take replacement from left
            t.data = findMin(t.left, cd, cd+1)
            t.right = delete(t.data, t.left, cd+1)
            t.left = null
        }
        else t = null // delete this leaf
        else if (x[cd] < t.data[cd]) // search left subtree
        t.left = delete(x, t.left, cd+1)
        else // search right subtree
        t.right = delete(x, t.right, cd+1)
        return t
```

        \}
    $\}$


Figure 55: Deletion from a kd-tree.
Recall that in the 1-dimensional case, in the 2-child case the replacement node was guaranteed to have either zero or one child. However this is not necessarily the case here. Thus we may do many 2-child deletions. As with insertion the running time is proportional to the height of the tree.

Nearest Neighbor Queries: Next we consider how to perform retrieval queries on a kd-tree. Nearest neighbor queries are among the most important queries. We are given a set of points $S$ stored in a kd-tree, and a query point $q$, and we want to return the point of $S$ that is closest to $q$. We assume that distances are measured using Euclidean distances, but generalizations to other sorts of distance functions is often possible.

An intuitively appealing approach to nearest neighbor queries would be to find the leaf node of the kd-tree that contains $q$ and then search this and the neighboring cells of the kd-tree. The problem is that the nearest neighbor may actually be very far away, in the sense of the tree's structure. For example, in the figure shown below, the cell containing the query point belongs to node $(70,30)$ but the nearest neighbor is far away in the tree at $(20,50)$. neighbor
We will need a more systematic approach to finding nearest neighbors. Nearest neighbor queries illustrate three important elements of range and nearest neighbor processing.

Partial results: Store the intermediate results of the query and update these results as the query proceeds.
Traversal order: Visit the subtree first that is more likely to be relevant to the final results.
Pruning: Do not visit any subtree that be judged to be irrelevant to the final results.
Throughout the search we will maintain an object close which contains the partial results of the query. For nearest neighbor search this has two fields, close.dist, the distance to the closest point seen so far and close.point, the data point that is the closest so far. Initially close.point = null and close.dist = INFINITY. This structure will be passed in as an argument to the recursive search procedure, and will be returned from the procedure. (Alternatively, it could be passed in as a reference parameter or a pointer to the object and updated.)
We will also pass in two other arguments. One is the cutting dimension cd, which as always, cycles among the various dimensions. The other is a rectangle object $r$ that represents the current node's cell. This argument is important, because it is used to determine whether we should visit a node's descendents. For example, if the (minimum) distance from the query point to this cell is greater than the distance from the query point to the closest point seen so far, there is no need to search the subtree rooted at this node. This is how we prune the search space.
When we arrive at a node $t$, we first determine whether the associated cell $r$ is close enough to provide the nearest neighbor. If not, we return immediately. Next we test whether the associated data point $t$. data is closer to the query $q$ than the current closest. If so, we update the current closest. Finally we search $t$ 's subtrees. Rather than just visiting the left subtree first (as we usually do), instead we will visit the subtree whose cell is closer to $q$. To do this, we test which side of $t$ 's splitting line $q$ lies. We recursively search this subtree first and then visit the farther subtree second. Each call updates the nearest neighbor result. It is important to prioritize the visit order, because the sooner we find the nearest neighbor the more effective our pruning step will be. So it makes sense to search first where we think we have the greater chance of finding the nearest neighbor.
In support of this, we need to implement a rectangle class. We assume that this class supports two utilities for manipulating rectangles. Given a rectangle $r$ representing the current cell and given the current cutting dimension cd and the cutting value t.data[cd], one method r.trimLeft() returns the rectangular cell for the left child and the other r.trimRight() returns the rectangular cell for the right child. We leave their implementation as an exercise. We also assume there are functions distance ( $\mathrm{q}, \mathrm{p}$ ) and distance ( $\mathrm{q}, \mathrm{r}$ ) that compute the distance from $q$ to a point $p$ and from $q$ to (the closest part of) a rectangle $r$. The complete code and a detailed example are given below.
kd-tree Nearest Neighbor Searching

```
NNResult nearNeigh(Point q, KDNode t, int cd, Rectangle r, NNResult close) {
    if (t == null) return close // fell out of tree
    if (distance(q, r) >= close.dist) // this cell is too far away
```

```
        return close;
Scalar dist = distance(q, t.data) // distance to t's data
if (dist < close.dist) // is this point closer?
    close = NNResult(t.data, dist) // yes, update
if (q[cd] < t.data[cd]) { // q is closer to left child
    close = nearNeigh(q, t.left, cd+1, r.trimLeft(cd, t.data), close)
    close = nearNeigh(q, t.right, cd+1, r.trimRight(cd, t.data), close)
}
else { // q is closer to right child
    close = nearNeigh(q, t.right, cd+1, r.trimRight(cd, t.data), close)
    close = nearNeigh(q, t.left, cd+1, r.trimLeft(cd, t.data), close)
}
return close
}
```

This structure is typical of complex range queries in geometric data structures. One common simplification (or cheat) that is often used to avoid having to maintain the current rectangular cell $r$, is to replace the line that computes the distance from $q$ to $r$ with a simpler computation. Just prior to visiting the farther child, we compute the distance from $q$ to $t$ 's splitting line. If this distance is greater than the current closest distance we omit the second recursive call to the farther child. However, using the cell produces more accurate pruning, which usually leads to better run times.
The running time of nearest neighbor searching can be quite bad in the worst case. In particular, it is possible to contrive examples where the search visits every node in the tree, and hence the running time is $O(n)$, for a tree of size $n$. However, the running time can be shown to be much closer to $O\left(2^{d}+\log n\right)$, where $d$ is the dimension of the space. Generally, you expect to visit some set of nodes that are in the neighborhood of the query point (giving rise to the $2^{d}$ term) and require $O(\log n)$ time to descend the tree to find these nodes.

## Lecture 19: Range Searching in kd-trees

## (Tuesday, April 17, 2001)

Reading: Today's material is discussed in Chapt. 2 of Samet's book on spatial data structures.
Range Queries: So far we have discussed insert, deletion, and nearest neighbor queries in kd-trees. Today we consider an important class of queries, called orthogonal range queries, or just range queries for short. Given a set of points $S$, the query consists of an axis-aligned rectangle $r$, and the problem is to report the points lying within this rectangle. Another variant is to count the number of points lying within the rectangle. These variants are called range reporting queries and range counting queries, respectively. (Range queries can also be asked with nonorthogonal objects, such as triangles and circles.) We will consider how to answer range queries using the kd-tree data structure. For example, in the figure below the query would report the points $\{7,8,10,11\}$ or the count of 4 points.

Range Searching in kd-trees: Let us consider how to answer range counting queries. We assume that the range $Q$ is a class which supports three basic utilities.
Q.contains (Point p): True if $p$ is contained within $Q$.
Q.contains (Rectangle C) : True if rectangle $C$ is entirely contained within $Q$.
Q.isDisjointFrom(Rectangle C) : True if rectangle $C$ is entirely disjoint from $Q$.

$60,10 \mathrm{y}$


Figure 56: Nearest neighbor search.


Figure 57: Orthogonal range query.

Let $t$ denote the current node in the tree, and let $C$ denote the current rectangular cell associated with $t$. As usual let $c d$ denote the current cutting dimension. The search proceeds as follows. First, if we fall out of the tree then there is nothing to count. If the current node's cell is disjoint from the query range, then again there is nothing to count. If the query range contains the current cell, then we can return all the points in the subtree rooted at $t$ ( t . size). Otherwise, we determine whether this point is in the range and count it if so, and then we recurse on the two children and update the count accordingly. Recall that we have overloaded the addition operation for cutting dimensions, so $c d+1$ is taken modulo the dimension. kd-tree Range Counting Query

```
int range(Range Q, KDNode t, int cd, Rectangle C) {
    if (t == null) return 0 // fell out of tree
    if (Q.isDisjointFrom(C)) return 0 // cell doesn't overlap query
    if (Q.contains(C)) return t.size // query contains cell
    int count = 0
    if (Q.contains(t.data)) count++ // count this data point
    // recurse on children
    count += range(Q, t.left, cd+1, C.trimLeft(cd, t.data))
    count += range(Q, t.right, cd+1, C.trimRight(cd, t.data))
    return count
}
```

The figure below shows an example of a range search. Next to each node we store the size of the associated subtree. We say that a node is visited if a call to range() is made on this node. We say that a node is processed if both of its children are visited. Observe that for a node to be processed, its cell must overlap the range without being contained within the range. In the example, the shaded nodes are those that are not processed. For example the subtree rooted at $b$ is entirely disjoint from the query and the subtrees rooted at $n$ and $r$ are entirely contained in the query. The nodes with squares surrounding them those whose points have been added individually to the count (by the second to last line of the procedure). There are 5 such nodes, and including the 3 points from the subtree rooted at $r$ and the 2 points from the subtree rooted at $n$, the final count returned is 10 .

Analysis of query time: How many nodes does this method visit altogether? We claim that the total number of nodes is $O(\sqrt{n})$ assuming a balanced kd-tree (which is a reasonable assumption in the average case).

Theorem: Given a balanced kd-tree with $n$ points range queries can be answered in $O(\sqrt{n})$ time.

Recall from the discussion above that a node is processed (both children visited) if and only if the cell overlaps the range without being contained within the range. We say that such a cell


Figure 58: Range search in kd-trees.
is stabbed by the query. To bound the total number of nodes that are processed in the search, it suffices to count the total number of nodes whose cells are stabbed by the query rectangle. Rather than prove the above theorem directly, we will prove a simpler result, which illustrates the essential ideas behind the proof. Rather than using a 4 -sided rectangle, we consider an orthogonal range having a only one side, that is, an orthogonal halfplane. In this case, the query stabs a cell if the vertical or horizontal line that defines the halfplane intersects the cell.

Lemma: Given a balanced kd-tree with $n$ points, any vertical or horizontal line stabs $O(\sqrt{n})$ cells of the tree.
Proof: Let us consider the case of a vertical line $x=x_{0}$. The horizontal case is symmetrical. Consider a processed node which has a cutting dimension along $x$. The vertical line $x=x_{0}$ either stabs the left child or the right child but not both. If it fails to stab one of the children, then it cannot stab any of the cells belonging to the descendents of this child either. If the cutting dimension is along the $y$-axis (or generally any other axis in higher dimensions), then the line $x=x_{0}$ stabs both children's cells.
Since we alternate splitting on left and right, this means that after descending two levels in the tree, we may stab at most two of the possible four grandchildren of each node. In general each time we descend two more levels we double the number of nodes being stabbed. Thus, we stab the root node, at most 2 nodes at level 2 of the tree, at most 4 nodes at level 4,8 nodes at level 6 , and generally at most $2^{i}$ nodes at level $2 i$.
Because we have an exponentially increasing number, the total sum is dominated by its last term. Thus is suffices to count the number of nodes stabbed at the lowest level of the tree. If we assume that the kd-tree is balanced, then the tree has roughly $\lg n$ levels. Thus the number of leaf nodes processed at the bottommost level is $2^{(\lg n) / 2}=2^{\lg \sqrt{n}}=\sqrt{n}$. This completes the proof.

We have shown that any vertical or horizontal line can stab only $O(\sqrt{n})$ cells of the tree. Thus, if we were to extend the four sides of $Q$ into lines, the total number of cells stabbed by all these lines is at most $O(4 \sqrt{n})=O(\sqrt{n})$. Thus the total number of cells stabbed by the query range is $O(\sqrt{n})$, and hence the total query time is $O(\sqrt{n})$. Again, this assumes that the kd-tree is balanced (having $O(\log n)$ depth). If the points were inserted in random order, this will be true on average.

## Lecture 20: Range Trees

(Thursday, April 19, 2001)

Reading: Samet's book, Section 2.5.
Range Queries: Last time we saw how to use kd-trees to solve orthogonal range queries. Recall that we are given a set of points in the plane, and we want to count or report all the points that lie within a given axis-aligned rectangle. We argued that if the tree is balanced, then the running time is close to $O(\sqrt{n})$ to count the points in the range. If there are $k$ points in the range, then we can list them in additional $O(k)$ time for a total of $O(k+\sqrt{n})$ time to answer range reporting queries.
Although $\sqrt{n}$ is much smaller than $n$, it is still larger than $\log n$. Is this the best we can do? Today we will present a data structure called a range tree which can answer orthogonal counting range queries in $O\left(\log ^{2} n\right)$. (Recall that $\log ^{2} n$ means $\left.(\log n)^{2}\right)$. If there are $k$ points in the range it can also report these points in $O\left(k+\log ^{2} n\right)$ time. It uses $O(n \log n)$ space, which is somewhat larger than the $O(n)$ space used by kd-trees. (There are actually two versions of range trees. We will present the simpler version. There is a significantly more complex version that can answer queries in $O(k+\log n)$ time, thus shaving off a log factor in the running time.) The data structure can be generalized to higher dimensions. In dimension $d$ it answers range queries in $O\left(\log ^{d} n\right)$ time.

Range Trees (Basics): The range tree data structure works by reducing an orthogonal range query in 2-dimensions to a collection of $O(\log n)$ range queries in 1-dimension, then it solves each of these in $O(\log n)$ time, for a combined time of $O\left(\log ^{2} n\right)$. (In higher dimensions, it reduces a range query in dimension $d$ to $O(\log n)$ range queries in dimension $d-1$.) It works by a technique called a multi-level tree, which involves cascading multiple data structures together in a clever way. Throughout we assume that a range is given by a pair of points [lo, hi], and we wish to report all points $p$ such that

$$
l o_{x} \leq p_{x} \leq h i_{x} \quad \text { and } \quad l o_{y} \leq p_{y} \leq h i_{y}
$$

1-dimensional Range Tree: Before discussing 2-dimensional range trees, let us first consider what a 1-dimensional range tree would look like. Given a set of points $S$, we want to preprocess these points so that given a 1 -dimensional interval $Q=[l o, h i]$ along the $x$-axis, we can count all the points that lie in this interval. There are a number of simple solutions to this, but we will consider a particular method that generalizes to higher dimensions.
Let us begin by storing all the points of our data set in the external nodes (leaves) of a balanced binary tree sorted by $x$-coordinates (e.g., an AVL tree). The data values in the internal nodes will just be used for searching purposes. They may or may not correspond to actual data values stored in the leaves. We assume that if an internal node contains a value $x_{0}$ then the leaves in the left subtree are strictly less than $x_{0}$, and the leaves in the right subtree are greater than or equal to $x_{0}$. Each node $t$ in this tree is implicitly associated with a subset $S(t) \subseteq S$ of elements of $S$ that are in the leaves descended from $t$. (For example $S($ root $)=S$.) We assume that for each node $t$, we store the number of leaves that are descended from $t$, denoted $t$.size. Thus t. size is equal to the number of elements in $S(t)$.
Let us introduce a few definitions before continuing. Given the interval $Q=[l o, h i]$, we say that a node $t$ is relevant to the query if $S(t) \subseteq Q$. That is, all the descendents of $t$ lie within the interval. If $t$ is relevant then clearly all of the nodes descended from $t$ are also relevant. A relevant node $t$ is canonical if $t$ is relevant, but its parent it not. The canonical nodes are the roots of the maximal subtrees that are contained within $Q$. For each canonical node $t$, the subset $S(t)$ is called a canonical subset. Because of the hierarchical structure of the tree, it is
easy to see that the canonical subsets are disjoint from each other, and they cover the interval $Q$. In other words, the subset of points of $S$ lying within the interval $Q$ is equal to the disjoint union of the canonical subsets. Thus, solving a range counting query reduces to finding the canonical nodes for the query range, and returning the sum of their sizes.

We claim that the canonical subsets corresponding to any range can be identified in $O(\log n)$ time from this structure. Intuitively, given any interval [lo, hi], we search the tree to find the leftmost leaf $u$ whose key is greater than or equal to $l o$ and the rightmost leaf $v$ whose key is less than or equal to $h i$. Clearly all the leaves between $u$ and $v$ (including $u$ and $v$ ) constitute the points that lie within the range. Since these two paths are of length at most $O(\log n)$, there are at most $O(2 \log n)$ such trees possible, which is $O(\log n)$. To form the canonical subsets, we take the subsets of all the maximal subtrees lying between $u$ and $v$. This is illustrated in the following figure.


Figure 59: Canonical sets for interval queries.
There are a few different ways to map this intuition into an algorithm. Our approach will be modeled after the approach used for range searching in kd-trees. We will maintain for each node a cell $C$, which in this 1-dimensional case is just an interval $\left[C_{l o}, C_{h i}\right]$. As with kd-trees, the cell for node $t$ contains all the points in $S(T)$.

The arguments to the procedure are the current node, the range $Q$, and the current cell. Let $C_{0}=[-\infty,+\infty]$ be the initial cell for the root node. The initial call is range1D (root, Q , $\mathrm{C} 0)$. Let $t . x$ denote the key associated with $t$. If $C=\left[x_{0}, x_{1}\right]$ denotes the current interval for node $t$, then when we recurse on the left subtree we trim this to the interval $\left[x_{0}, t . x\right]$ and when we recurse on the right subtree we trim the interval to $\left[t \cdot x, x_{1}\right]$. We assume that given two ranges $A$ and $B$, we have utility functions A.contains (B) which determined whether interval $A$ contains interval $B$, and there is a similar function A.contains(x) that determines whether point $x$ is contained within $A$.
Since the data are only stored in the leaf nodes, when we encounter such a node we consider whether it lies in the range and count it if so. Otherwise, observe that if $t . x \leq Q$.lo then all the points in the left subtree are less than the interval, and hence it does not need to be visited. Similarly if $t . x>Q . h i$ then the right subtree does not need to be visited. Otherwise, we need to visit these subtrees recursively.

1-Dimensional Range Counting Query

```
int range1Dx(Node t, Range Q, Interval C=[x0,x1]) {
    if (t.isExternal) // hit the leaf level?
```

```
            return (Q.contains(t.point) ? 1 : 0) // count if point in range
    if (Q.contains(C)) // Q contains entire cell?
            return t.size // return entire subtree size
    int count = 0
    if (t.x > Q.lo) // overlap left subtree?
        count += range1Dx(t.left, Q, [x0, t.x]) // count left subtree
    if (t.x <= Q.hi) // overlap right subtree?
        count += range1Dx(t.right, Q, [t.x, x1])// count right subtree
    return count
}
```



Figure 60: Range search cases.

The external nodes counted in the second line and the internal nodes for which we return $t . s i z e$ are the canonical nodes. The above procedure answers range counting queries. To extend this to range reporting queries, we simply replace the step that counts the points in the subtree with a procedure that traverses the subtree and prints the data in the leaves. Each tree can be traversed in time proportional to the number of leaves in each subtree. Combining the observations of this section we have the following results.

Lemma: Given a 1-dimensional range tree and any query range $Q$, in $O(\log n)$ time we can compute a set of $O(\log n)$ canonical nodes $t$, such that the answer to the query is the disjoint union of the associated canonical subsets $S(t)$.
Theorem: 1-dimensional range counting queries can be answered in $O(\log n)$ time and range reporting queries can be answered in $O(k+\log n)$ time, where $k$ is the number of values reported.

Range Trees: Now let us consider how to answer range queries in 2-dimensional space. We first create 1-dimensional tree $T$ as described in the previous section sorted by the $x$-coordinate. For each internal node $t$ of $T$, recall that $S(t)$ denotes the points associated with the leaves descended from $t$. For each node $t$ of this tree we build a 1-dimensional range tree for the points of $S(t)$, but sorted on $y$-coordinates. This called the auxiliary tree associated with $t$. Thus, there are $O(n)$ auxiliary trees, one for each internal node of $T$. An example of such a structure is shown in the following figure.
Notice that there is duplication here, because a given point in a leaf occurs in the point sets associated with each of its ancestors. We claim that the total sizes of all the auxiliary trees is $O(n \log n)$. To see why, observe that each point in a leaf of $T$ has $O(\log n)$ ancestors, and so each point appears in $O(\log n)$ auxiliary trees. The total number of nodes in each auxiliary tree is proportional to the number of leaves in this tree. (Recall that the number of internal nodes in an extended tree one less than the number of leaves.) Since each of the $n$ points appears as a leaf in at most $O(\log n)$ auxiliary trees, the sum of the number of leaves in all the auxiliary trees is at most $O(n \log n)$. Since $T$ has $O(n)$ nodes, the overall total is $O(n \log n)$.

Now, when a 2-dimensional range is presented we do the following. First, we invoke a variant of the 1-dimensional range search algorithm to identify the $O(\log n)$ canonical nodes. For each such node $t$, we know that all the points of the set lie within the $x$ portion of the range, but not necessarily in the $y$ part of the range. So we search the 1-dimensional auxiliary range and return a count of the resulting points. The algorithm below is almost identical the previous one, except that we make explicit reference to the $x$-coordinates in the search, and rather than adding t.size to the count, we invoke a 1-dimensional version of the above procedure using the $y$-coordinate instead. Let $\mathbf{Q} . \mathrm{x}$ denote the $x$-part of $Q$ 's range, consisting of the interval [Q.lo.x,Q.hi.x]. The call Q.contains(t.point) is applied on both coordinates, but the call Q.x.contains (C) only checks the $x$-part of $Q$ 's range. The algorithm is given below. The procedure range1Dy() is the same procedure described above, except that it searches on $y$ rather than $x$.

2-Dimensional Range Counting Query

```
int range2D(Node t, Range2D Q, Range1D C=[x0,x1]) {
    if (t.isExternal) // hit the leaf level?
        return (Q.contains(t.point) ? 1 : 0) // count if point in range
        if (Q.x.contains(C)) { // Q's x-range contains C
        [y0,y1] = [-infinity, +infinity] // initial y-cell
        return range1Dy(t.aux.root, Q, [y0, y1])// search auxiliary tree
    }
    int count = 0
    if (t.x > Q.lo.x) // overlap left subtree?
        count += range2D(t.left, Q, [x0, t.x]) // count left subtree
    if (t.x <= Q.hi.x) // overlap right subtree?
        count += range2D(t.right, Q, [t.x, x1]) // count right subtree
    return count
}
```



Figure 61: Range tree.

Analysis: It takes $O(\log n)$ time to identify the canonical nodes along the $x$-coordinates. For each of these $O(\log n)$ nodes we make a call to a 1-dimensional range tree which contains no more
than $n$ points. As we argued above, this takes $O(\log n)$ time for each. Thus the total running time is $O\left(\log ^{2} n\right)$. As above, we can replace the counting code with code in range1Dy() with code that traverses the tree and reports the points. This results in a total time of $O\left(k+\log ^{2} n\right)$, assuming $k$ points are reported.

Thus, each node of the 2-dimensional range tree has a pointer to a auxiliary 1-dimensional range tree. We can extend this to any number of dimensions. At the highest level the $d$ dimensional range tree consists of a 1-dimensional tree based on the first coordinate. Each of these trees has an auxiliary tree which is a $(d-1)$-dimensional range tree, based on the remaining coordinates. A straightforward generalization of the arguments presented here show that the resulting data structure requires $O\left(n \log ^{d} n\right)$ space and can answer queries in $O\left(\log ^{d} n\right)$ time.

Theorem: $d$-dimensional range counting queries can be answered in $O\left(\log ^{d} n\right)$ time, and range reporting queries can be answered in $O\left(k+\log ^{d} n\right)$ time, where $k$ is the number of values reported.

## Lecture 21: Interval Trees

(Tuesday, April 24, 2001)
Reading: The presentation is not taken from any of our readings. It is derived from a description in the book Computational Geometry: Algorithms and Applications, by M. de Berg, M. van Kreveld, M. Overmars, and O. Schwarzkopf, Chapt 10. A copy is on reserve in the computer science library in A.V. Williams.

Segment Data: So far we have considered geometric data structures for storing points. However, there are many others types of geometric data that we may want to store in a data structure. Today we consider how to store orthogonal (horizontal and vertical) line segments in the plane. We assume that a line segment is represented by giving its pair of endpoints. The segments are allowed to intersect one another.

As a basic motivating query, we consider the following window query. Given a set of orthogonal line segments $S$, which have been preprocessed, and given an orthogonal query rectangle $W$, count or report all the line segments of $S$ that intersect $W$. We will assume that $W$ is closed and solid rectangle, so that even if a line segment lies entirely inside of $W$ or intersects only the boundary of $W$, it is still reported. For example, given the window below, the query would report the segments that are shown with solid lines, and segments with broken lines would not be reported.


Figure 62: Window Query.

Window Queries for Orthogonal Segments: We will present a data structure, called the interval tree, which (combined with a range tree) can answer window counting queries for orthogonal
line segments in $O\left(\log ^{2} n\right)$ time, where $n$ is the number line segments. It can report these segments in $O\left(k+\log ^{2} n\right)$ time, where and $k$ is the total number of segments reported. The interval tree uses $O(n \log n)$ storage and can be built in $O(n \log n)$ time.
We will consider the case of range reporting queries. (There are some subtleties in making this work for counting queries.) We will derive our solution in steps, starting with easier subproblems and working up to the final solution. To begin with, observe that the set of segments that intersect the window can be partitioned into three types: those that have no endpoint in $W$, those that have one endpoint in $W$, and those that have two endpoints in $W$.
We already have a way to report segments of the second and third types. In particular, we may build a range tree just for the $2 n$ endpoints of the segments. We assume that each endpoint has a cross-link indicating the line segment with which it is associated. Now, by applying a range reporting query to $W$ we can report all these endpoints, and follow the cross-links to report the associated segments. Note that segments that have both endpoints in the window will be reported twice, which is somewhat unpleasant. We could fix this either by sorting the segments in some manner and removing duplicates, or by marking each segment as it is reported and ignoring segments that have already been marked. (If we use marking, after the query is finished we will need to go back an "unmark" all the reported segments in preparation for the next query.)
All that remains is how to report the segments that have no endpoint inside the rectangular window. We will do this by building two separate data structures, one for horizontal and one for vertical segments. A horizontal segment that intersects the window but neither of its endpoints intersects the window must pass entirely through the window. Observe that such a segment intersects any vertical line passing from the top of the window to the bottom. In particular, we could simply ask to report all horizontal segments that intersect the left side of $W$. This is called a vertical segment stabbing query. In summary, it suffices to solve the following subproblems (and remove duplicates):

Endpoint inside: Report all the segments of $S$ that have at least one endpoint inside $W$. (This can be done using a range query.)
Horizontal through segments: Report all the horizontal segments of $S$ that intersect the left side of $W$. (This reduces to a vertical segment stabbing query.)

Vertical through segments: Report all the vertical segments of $S$ that intersect the bottom side of $W$. (This reduces to a horizontal segment stabbing query.)

We will present a solution to the problem of vertical segment stabbing queries. Before dealing with this, we will first consider a somewhat simpler problem, and then modify this simple solution to deal with the general problem.

Vertical Line Stabbing Queries: Let us consider how to answer the following query, which is interesting in its own right. Suppose that we are given a collection of horizontal line segments $S$ in the plane and are given an (infinite) vertical query line $\ell_{q}: x=x_{q}$. We want to report all the line segments of $S$ that intersect $\ell_{q}$. Notice that for the purposes of this query, the $y$-coordinates are really irrelevant, and may be ignored. We can think of each horizontal line segment as being a closed interval along the $x$-axis. We show an example in the figure below on the left.
As is true for all our data structures, we want some balanced way to decompose the set of intervals into subsets. Since it is difficult to define some notion of order on intervals, we instead will order the endpoints. Sort the interval endpoints along the $x$-axis. Let $\left\langle x_{1}, x_{2}, \ldots, x_{2 n}\right\rangle$ be the resulting sorted sequence. Let $x_{m e d}$ be the median of these $2 n$ endpoints. Split the intervals into three groups, $L$, those that lie strictly to the left of $x_{\text {med }}, R$ those that lie strictly


Figure 63: Line Stabbing Query.
to the right of $x_{\text {med }}$, and $M$ those that contain the point $x_{\text {med }}$. We can then define a binary tree by putting the intervals of $L$ in the left subtree and recursing, putting the intervals of $R$ in the right subtree and recursing. Note that if $x_{q}<x_{\text {med }}$ we can eliminate the right subtree and if $x_{q}>x_{\text {med }}$ we can eliminate the left subtree. See the figure right.
But how do we handle the intervals of $M$ that contain $x_{\text {med }}$ ? We want to know which of these intervals intersects the vertical line $\ell_{q}$. At first it may seem that we have made no progress, since it appears that we are back to the same problem that we started with. However, we have gained the information that all these intervals intersect the vertical line $x=x_{\text {med }}$. How can we use this to our advantage?
Let us suppose for now that $x_{q} \leq x_{\text {med }}$. How can we store the intervals of $M$ to make it easier to report those that intersect $\ell_{q}$. The simple trick is to sort these lines in increasing order of their left endpoint. Let $M_{L}$ denote the resulting sorted list. Observe that if some interval in $M_{L}$ does not intersect $\ell_{q}$, then its left endpoint must be to the right of $x_{q}$, and hence none of the subsequent intervals intersects $\ell_{q}$. Thus, to report all the segments of $M_{L}$ that intersect $\ell_{q}$, we simply traverse the sorted list and list elements until we find one that does not intersect $\ell_{q}$, that is, whose left endpoint lies to the right of $x_{q}$. As soon as this happens we terminate. If $k^{\prime}$ denotes the total number of segments of $M$ that intersect $\ell_{q}$, then clearly this can be done in $O\left(k^{\prime}+1\right)$ time.
On the other hand, what do we do if $x_{q}>x_{\text {med }}$ ? This case is symmetrical. We simply sort all the segments of $M$ in a sequence, $M_{R}$, which is sorted from right to left based on the right endpoint of each segment. Thus each element of $M$ is stored twice, but this will not affect the size of the final data structure by more than a constant factor. The resulting data structure is called an interval tree.

Interval Trees: The general structure of the interval tree was derived above. Each node of the interval tree has a left child, right child, and itself contains the median $x$-value used to split the set, $x_{\text {med }}$, and the two sorted sets $M_{L}$ and $M_{R}$ (represented either as arrays or as linked lists) of intervals that overlap $x_{\text {med }}$. We assume that there is a constructor that builds a node given these three entities. The following high-level pseudocode describes the basic recursive step in the construction of the interval tree. The initial call is root $=\operatorname{IntTree}(S)$, where $S$ is the initial set of intervals. Unlike most of the data structures we have seen so far, this one is not built by the successive insertion of intervals (although it would be possible to do so). Rather we assume that a set of intervals $S$ is given as part of the constructor, and the entire structure
is built all at once. We assume that each interval in $S$ is represented as a pair $\left(x_{\mathrm{lo}}, x_{\mathrm{hi}}\right)$. An example is shown in the following figure.

```
IntTreeNode IntTree(IntervalSet S) {
    if (|S| == 0) return null // no more
    xMed = median endpoint of intervals in S // median endpoint
    L = {[xlo, xhi] in S | xhi < xMed} // left of median
    R}={[xlo, xhi] in S | xlo > xMed} // right of median
    M = {[xlo, xhi] in S | xlo <= xMed <= xhi} // contains median
    ML = sort M in increasing order of xlo // sort M
    MR = sort M in decreasing order of xhi
    t = new IntTreeNode(xMed, ML, MR) // this node
    t.left = IntTree(L)
// left subtree
    t.right = IntTree(R)
    return t
}
```



Figure 64: Interval Tree.

We assert that the height of the tree is $O(\log n)$. To see this observe that there are $2 n$ endpoints. Each time through the recursion we split this into two subsets $L$ and $R$ of sizes at most half the original size (minus the elements of $M$ ). Thus after at most $\lg (2 n)$ levels we will reduce the set sizes to 1 , after which the recursion bottoms out. Thus the height of the tree is $O(\log n)$.

Implementing this constructor efficiently is a bit subtle. We need to compute the median of the set of all endpoints, and we also need to sort intervals by left endpoint and right endpoint. The fastest way to do this is to presort all these values and store them in three separate lists. Then as the sets $L, R$, and $M$ are computed, we simply copy items from these sorted lists to the appropriate sorted lists, maintaining their order as we go. If we do so, it can be shown that this procedure builds the entire tree in $O(n \log n)$ time.

The algorithm for answering a stabbing query was derived above. We summarize this algorithm below. Let $x_{q}$ denote the $x$-coordinate of the query line.

Line Stabbing Queries for an Interval Tree

```
stab(IntTreeNode t, Scalar xq) {
    if (t == null) return // fell out of tree
    if (xq < t.xMed) { // left of median?
```

```
        for (i = 0; i < t.ML.length; i++) { // traverse ML
        if (t.ML[i].lo <= xq) print(t.ML[i])// ..report if in range
        else break // ..else done
        }
        stab(t.left, xq) // recurse on left
    }
    else { // right of median
        for (i = 0; i < t.MR.length; i++) { // traverse MR
        if (t.MR[i].hi >= xq) print(t.MR[i])// ..report if in range
        else break // ..else done
        }
        stab(t.right, xq) // recurse on right
    }
}
```

This procedure actually has one small source of inefficiency, which was intentionally included to make code look more symmetric. Can you spot it? Suppose that $x_{q}=t \cdot x_{\text {med }}$ ? In this case we will recursively search the right subtree. However this subtree contains only intervals that are strictly to the right of $x_{\text {med }}$ and so is a waste of effort. However it does not affect the asymptotic running time.
As mentioned earlier, the time spent processing each node is $O\left(1+k^{\prime}\right)$ where $k^{\prime}$ is the total number of points that were recorded at this node. Summing over all nodes, the total reporting time is $O(k+v)$, where $k$ is the total number of intervals reported, and $v$ is the total number of nodes visited. Since at each node we recurse on only one child or the other, the total number of nodes visited $v$ is $O(\log n)$, the height of the tree. Thus the total reporting time is $O(k+\log n)$.

Vertical Segment Stabbing Queries: Now let us return to the question that brought us here. Given a set of horizontal line segments in the plane, we want to know how many of these segments intersect a vertical line segment. Our approach will be exactly the same as in the interval tree, except for how the elements of $M$ (those that intersect the splitting line $x=x_{\text {med }}$ ) are handled.
Going back to our interval tree solution, let us consider the set $M$ of horizontal line segments that intersect the splitting line $x=x_{\text {med }}$ and as before let us consider the case where the query segment $q$ with endpoints $\left(x_{q}, y_{\text {lo }}\right)$ and $\left(x_{q}, y_{\text {hi }}\right)$ lies to the left of the splitting line. The simple trick of sorting the segments of $M$ by their left endpoints is not sufficient here, because we need to consider the $y$-coordinates as well. Observe that a segment of $M$ stabs the query segment $q$ if and only if the left endpoint of a segment lies in the following semi-infinite rectangular region.

$$
\left\{(x, y) \mid x \leq x_{q} \text { and } y_{\mathrm{lo}} \leq y \leq y_{\mathrm{hi}}\right\}
$$

This is illustrated in the figure below. Observe that this is just an orthogonal range query. (It is easy to generalize the procedure given last time to handle semi-infinite rectangles.) The case where $q$ lies to the right of $x_{\text {med }}$ is symmetrical.
So the solution is that rather than storing $M_{L}$ as a list sorted by the left endpoint, instead we store the left endpoints in a 2-dimensional range tree (with cross-links to the associated segments). Similarly, we create a range tree for the right endpoints and represent $M_{R}$ using this structure.

The segment stabbing queries are answered exactly as above for line stabbing queries, except that part that searches $M_{L}$ and $M_{R}$ (the for-loops) are replaced by searches to the appropriate range tree, using the semi-infinite range given above.
We will not discuss construction time for the tree. (It can be done in $O(n \log n)$ time, but this involves some thought as to how to build all the range trees efficiently). The space needed


Figure 65: The segments that stab $q$ lie within the shaded semi-infinite rectangle.
is $O(n \log n)$, dominated primarily from the $O(n \log n)$ space needed for the range trees. The query time is $O\left(k+\log ^{3} n\right)$, since we need to answer $O(\log n)$ range queries and each takes $O\left(\log ^{2} n\right)$ time plus the time for reporting. If we use the spiffy version of range trees (which we mentioned but never discussed) that can answer queries in $O(k+\log n)$ time, then we can reduce the total time to $O\left(k+\log ^{2} n\right)$.

## Lecture 22: Memory Management

(Thursday, April 26, 2001)
Reading: Chapter 3 in Samet's notes.
Memory Management: One of the major systems issues that arises when dealing with data structures is how storage is allocated and deallocated as objects are created and destroyed. Although memory management is really a operating systems issue, we will discuss this topic over the next couple of lectures because there are a number interesting data structures issues that arise. In addition, sometimes for the sake of efficiency, it is desirable to design a special-purpose memory management system for your data structure application, rather than using the system's memory manager.

We will not discuss the issue of how the runtime system maintains memory in great detail. Basically what you need to know is that there are two principal components of dynamic memory allocation. The first is the stack. When procedures are called, arguments and local variables are pushed onto the stack, and when a procedure returns these variables are popped. The stacks grows and shrinks in a very predictable way. Variables allocated through new are stored in a different section of memory called the heap. (In spite of the similarity of names, this heap has nothing to do with the binary heap data structure, which is used for priority queues.) As elements are allocated and deallocated, the heap storage becomes fragmented into pieces. How to maintain the heap efficiently is the main issue that we will consider.
There are two basic methods of doing memory management. This has to do with whether storage deallocation is done explicitly or implicitly. Explicit deallocation is what $\mathrm{C}++$ uses. Objects are deleted by invoking delete, which returns the storage back to the system. Objects that are inaccessible but not deleted become unusable waste. In contrast Java uses implicit deallocation. It is up to the system to determine which objects are no longer accessible and reclaim their storage. This process is called garbage collection. In both cases there are a number of choices that can be made, and these choices can have significant impacts on the performance of the memory allocation system. Unfortunately, there is not one system that is best for all circumstances. We begin by discussing explicit deallocation systems.

Explicit Allocation/Deallocation: There is one case in which explicit deallocation is very easy to handle. This is when all the objects being allocated are of the same size. A large contiguous portion of memory is used for the heap, and we partition this into blocks of size $b$, where $b$ is the size of each object. For each unallocated block, we use one word of the block to act as a next pointer, and simply link these blocks together in linked list, called the available space list. For each new request, we extract a block from the available space list and for each delete we return the block to the list.
If the records are of different sizes then things become much trickier. We will partition memory into blocks of varying sizes. Each block (allocated or available) contains information indicating how large it is. Available blocks are linked together to form an available space list. The main questions are: (1) what is the best way to allocate blocks for each new request, and (2) what is the fastest way to deallocate blocks for each delete request.

The biggest problem in such systems is the fact that after a series of allocation and deallocation requests the memory space tends to become fragmented into small blocks of memory. This is called external fragmentation, and is inherent to all dynamic memory allocators. Fragmentation increases the memory allocation system's running time by increasing the size of the available space list, and when a request comes for a large block, it may not be possible to satisfy this request, even though there is enough total memory available in these small blocks. Observe that it is not usually feasible to compact memory by moving fragments around. This is because there may be pointers stored in local variables that point into the heap. Moving blocks around would require finding these pointers and updating them, which is a very expensive proposition. We will consider it later in the context of garbage collection. A good memory manager is one that does a good job of controlling external fragmentation.

Overview: When allocating a block we must search through the list of available blocks of memory for one that is large enough to satisfy the request. The first question is, assuming that there does exist a block that is large enough to satisfy the request, what is the best block to select? There are two common but conflicting strategies:

First fit: Search the blocks sequentially until finding the first block that is big enough to satisfy the request.

Best fit: Search all available blocks and use the smallest one that is large enough to fulfill the request.

Both methods work well in some instances and poorly in others. Some writeups say that first fit is prefered because (1) it is fast (it need only search until it finds a block), (2) if best fit does not exactly fill a request, it tends to produce small "slivers" of available space, which tend to aggravate fragmentation.

One method which is commonly used to reduce fragmentation is when a request is just barely filled by an available block that is slightly larger than the request, we allocate the entire block (more than the request) to avoid the creation of the sliver. This keeps the list of available blocks free of large numbers of tiny fragments which increase the search time. The additional waste of space that results because we allocate a larger block of memory than the user requested is called internal fragmentation (since the waste is inside the allocated block).

When deallocating a block, it is important that if there is available storage we should merge the newly deallocated block with any neighboring available blocks to create large blocks of free space. This process is called merging. Merging is trickier than one might first imagine. For example, we want to know whether the preceding or following block is available. How would we do this? We could walk along the available space list and see whether we find it, but this would be very slow. We might store a special bit pattern at the end and start of each block to
indicate whether it is available or not, but what if the block is allocated that the data contents happen to match this bit pattern by accident? Let us consider the implementation of this scheme in greater detail.

Implementation: The main issues are how to we store the available blocks so we can search them quickly (skipping over used blocks), and how do we determine whether we can merge with neighboring blocks or not. We want the operations to be efficient, but we do not want to use up excessive pointer space (especially in used blocks). Here is a sketch of a solution (one of many possible).


Allocated Block


Figure 66: Block structure for dynamic storage allocation.

Allocated blocks: For each block of used memory we record the following information. It can all fit in the first word of the block.
size: An integer that indicates the size of the block of storage. This includes both the size that the user has requested and the additional space used for storing these extra fields.
inUse: A single bit that is set to 1 (true) to indicate that this block is in use.
prevInUse: A single bit that is set to 1 (true) if the previous block is in use and 0 otherwise. (Later we will see why this is useful.)
Available blocks: For an available block we store more information, which is okay because the user is not using this space. These blocks are stored in a doubly-linked circular list, called avail.
size: An integer that indicates the size of the block of storage (including these extra fields).
inUse: A bit that is set to 0 (false) to indicate that this block is not in use.
prevInUse: A bit that is set to (true) if the previous block is in use (which should always be true, since we should never have two consecutive unused blocks).
pred: A pointer to the predecessor block on the available space list. Note that the available space list is not sorted. Thus the predecessor may be anywhere in the heap.
next: A pointer to the next block on the available space list.
size2: Contains the same value as size. Unlike the previous fields, this size indicator is stored in the last word of the block. Since is not within a fixed offset of the head of the block, for block $p$ we access this field as $*(\mathrm{p}+\mathrm{p}$. size -1$)$.

Note that available blocks require more space for all this extra information. The system will never allow the creation of a fragment that is so small that it cannot contain all this information.

Block Allocation: To allocate a block we search through the linked list of available blocks until finding one of sufficient size. If the request is about the same size (or perhaps slightly smaller) as the block, we remove the block from the list of available blocks (performing the necessary relinkings) and return a pointer to it. We also may need to update the prevInUse bit of the next block since this block is no longer available. Otherwise we split the block into two smaller blocks, return one to the user, and leave the other on the available space list.
Here is pseudocode for the allocation routine. Note that we make use of pointer arithmetic here. The argument $b$ is the desired size of the allocation. Because we reserve one word of storage for our own use we increment this value on entry to the procedure. We keep a constant TOO_SMALL, which indicates the smallest allowable fragment size. If the allocation would result in a fragment of size less than this value, we return the entire block. The procedure returns a generic pointer to the newly allocated block. The utility function avail.unlink(p) simply unlinks block $p$ from the doubly-linked available space list. An example is provided in the figure below. Shaded blocks are available.

```
(void*) alloc(int b) { // allocate block of size b
    b += 1 // extra space for system overhead
    p = search available space list for block of size at least b
    if (p == null) { ...error: insufficient memory...}
    if (p.size - b < TOO_SMALL) { // allocate whole block
            avail.unlink(p) // unlink p from avail space list
            q = p
    }
    else { // split the block
        p.size -= b // decrease size by b
            *(p+p.size-1) = p.size // store size in p's size2 field
            q = p + p.size // offset to start of new block
            q.size = b // size of new block
            q.prevInUse = 0 // previous block is unused
    }
    q.inUse = true // new block is used
    (q+q.size).prevInUse = true // adjust prevInUse for following block
    return q
}
```

Block Deallocation: To deallocate a block, we check whether the next block or the preceding blocks are available. For the next block we can find its first word and check its inUse field. For the preceding block we use our own prevInUse field. (This is why this field is present). If the previous block is not in use, then we use the size value stored in the last word to find the block's header. If either of these blocks is available, we merge the two blocks and update the header values appropriately. If both the preceding and next blocks are available, then this result in one of these blocks being deleting from the available space list (since we do not want to have two consecutive available blocks). If both the preceding and next blocks are in-use, we simply link this block into the list of available blocks (e.g. at the head of the list). We will leave the details of deallocation as an exercise.

Analysis: There is not much theoretical analysis of this method of dynamic storage allocation. Because the system has no knowledge of the future sequence of allocation and deallocation requests, it is possible to contrive situations in which either first fit or best fit (or virtually any other method you can imagine) will perform poorly. Empirical studies based on simulations have shown that this method achieves utilizations of around $2 / 3$ of the total available storage


Figure 67: An example of block allocation and deallocation.
before failing to satisfy a request. Even higher utilizations can be achieved if the blocks are small on average and block sizes are similar (since this limits fragmentation). A rule of thumb is to allocate a heap that is at least 10 times larger than the largest block to be allocated.

## Lecture 23: More on Memory Management

(Tuesday, May 1, 2001)
Reading: Chapter 3 in Samet's notes.
Buddy System: The dynamic storage allocation method described last time suffers from the problem that long sequences of allocations and deallocations of objects of various sizes tends to result in a highly fragmented space. The buddy system is an alternative allocation system which limits the possible sizes of blocks and their positions, and so tends to produce a more well-structured allocation. Because it limits block sizes, internal fragmentation (the waste caused when an allocation request is mapped to a larger block size) becomes an issue.

The buddy system works by starting with a block of memory whose size is a power of 2 and then hierarchically subdivides each block into blocks of equal sizes (like a 1-dimensional version of a quadtree.) To make this intuition more formal, we introduce the two elements of the buddy system. The first element that the sizes of all blocks (allocated and unallocated) are powers of 2. When a request comes for an allocation, the request (including the overhead space needed for storing block size information) is artificially increased to the next larger power of 2. Note that the allocated size is never more than twice the size of the request. The second element is that blocks of size $2^{k}$ must start at addresses that are multiples of $2^{k}$. (We assume that addressing starts at 0 , but it is easy to update this scheme to start at any arbitrary address by simply shifting addresses by some offset.)


Figure 68: Buddy system block structure.

Note that the above requirements limits the ways in which blocks may be merged. For example the figure below illustrates a buddy system allocation of blocks, where the blocks of size $2^{k}$ are shown at the same level. Available blocks are shown in white and allocated blocks are shaded. The two available blocks at addresses 5 and 6 (the two white blocks between 4 and 8) cannot be merged because the result would be a block of length 2 , starting at address 5 , which is not a multiple of 2 . For each size group, there is a separate available space list.


Figure 69: Buddy system example.
For every block there is exactly one other block with which this block can be merged with. This is called its buddy. In general, if a block $b$ is of size $2^{k}$, and is located at address $x$, then its buddy is the block of size $2^{k}$ located at address

$$
b u d d y_{k}(x)= \begin{cases}x+2^{k} & \text { if } 2^{k+1} \text { divides } x \\ x-2^{k} & \text { otherwise }\end{cases}
$$

Although this function looks fancy, it is very easy to compute in a language which is capable of bit manipulation. Basically the buddy's address is formed by complementing bit $k$ in the binary representation of $x$, where the lowest order bit is bit 0 . In languages like $C++$ and Java this can be implemented efficiently by shifting a single bit to the left by $k$ positions and exclusive or-ing with $x$, that is, $(1 \ll k)^{\wedge} \mathrm{x}$. For example, for $k=3$ the blocks of length 8 at addresses 208 and 216 are buddies. If we look at their binary representations we see that they differ in bit position 3 , and because they must be multiples of 8 , bits $0-2$ are zero.

$$
\begin{aligned}
\text { bit position: } & 876543210 \\
208 & =011010000_{2} \\
216 & =011011000_{2}
\end{aligned}
$$

As we mentioned earlier, one principle advantage of the buddy system is that we can exploit the regular sizes of blocks to search efficiently for available blocks. We maintain an array of
linked lists, one for the available block list for each size, thus avail $[\mathrm{k}]$ is the header pointer to the available block list for blocks of size $k$.

Here is how the basic operations work. We assume that each block has the same structure as described in the dynamic storage allocation example from last time. The prevInUse bit and the size field at the end of each available block are not needed given the extra structure provided in the buddy system. Each block stores its size (actually the $\log k$ of its size is sufficient) and a bit indicating whether it is allocated or not. Also each available block has links to the previous and next entries on the available space list. There is not just one available space, but rather there are multiple lists, one for each level of the hierarchy (something like a skip list). This makes it possible to search quickly for a block of a given size.

Buddy System Allocation: We will give a high level description of allocation and dealloction. To allocate a block of size $b$, let $k=\lceil\lg (b+1)\rceil$. (Recall that we need to include one extra word for the block size. However, to simplify our figures we will ignore this extra word.) We will allocate a block of size $2^{k}$. In general there may not be a block of exactly this size available, so find the smallest $j \geq k$ such that there is an available block of size $2^{j}$. If $j>k$, repeatedly split this block until we create a block of size $2^{k}$. In the process we will create one or more new blocks, which are added to the appropriate available space lists.

For example, in the figure we request a block of length 2. There are no such blocks, so we remove the the block of length 16 at address 32 from its available space list, split it into subblocks of sizes $2,2,4$ and 8 , add the last three to the appropriate available space lists, and return the first block.


Figure 70: Buddy allocation.

Deallocation: To deallocate a block, we first mark this block as being available. We then check to see whether its buddy is available. This can be done in constant time. If so, we remove the buddy from its available space list, and merge them together into a single free block of twice the size. This process is repeated until we find that the buddy is allocated.

The figure shows the deletion of the block of size 1 at address 21. It is merged with its buddy at address 20 , forming a block of size 2 at 20 . This is then merged with its buddy at 22 , forming a block of size 4 at 20. Finally it is merged with its buddy at 16 , forming a block of size 8 at 16 .


Figure 71: Buddy deallocation.

## Lecture 24: Garbage Collection

(Thursday, May 3, 2001)
Reading: Chapter 3 in Samet's notes.
Garbage Collection: In contrast to the explicit deallocation methods discussed in the previous lectures, in some memory management systems such as Java, there is no explicit deallocation of memory. In such systems, when memory is exhausted it must perform garbage collection to reclaim storage and sometimes to reorganize memory for better future performance. We will consider some of the issues involved in the implementation of such systems.

Any garbage collection system must do two basic things. First, it must detect which blocks of memory are unreachable, and hence are "garbage". Second, it must reclaim the space used by these objects and make it available for future allocation requests. Garbage detection is typically performed by defining a set of roots, e.g., local variables that point to objects in the heap, and then finding everything that is reachable from these roots. An object is reachable (or live) if there is some path of pointers or references from the roots by which the executing program can access the object. The roots are always accessible to the program. Objects that are not reachable are considered garbage, because they can no longer affect the future course of program execution.

Reference counts: How do we know when a block of storage is able to be deleted? One simple way to do this is to maintain a reference count for each block. This is a counter associated with the block. It is set to one when the block is first allocated. Whenever the pointer to
this block is assigned to another variable, we increase the reference count. (For example, the compiler can overload the assignment operation to achieve this.) When a variable containing a pointer to the block is modified, deallocated or goes out of scope, we decrease the reference count. If the reference count ever equals 0 , then we know that no references to this object remain, and the object can be deallocated.

Reference counts have two significant shortcomings. First, there is a considerable overhead in maintaining reference counts, since each assignment needs to modify the reference counts. Second, there are situations where the reference count method may fail to recognize unreachable objects. For example, if there is a circular list of objects, for example, $X$ points to $Y$ and $Y$ points to $X$, then the reference counts of these objects may never go to zero, even though the entire list is unreachable.

Mark-and-sweep: A more complete alternative to reference counts involves waiting until space runs out, and then scavenging memory for unreachable cells. Then these unreachable regions of memory are returned to available storage. These available blocks can be stored in an available space list using the same method described in the previous lectures for dynamic storage allocation. This method works by finding all immediately accessible pointers, and then traces them down and marks these blocks as being accessible. Then we sweep through memory adding all unmarked blocks to the available space list. Hence this method is called mark-and-sweep. An example is illustrated in the figure below.


Figure 72: Mark-and-sweep garbage collection.
How do we implement the marking phase? One way is by performing simple depth-first traversal of the "directed graph" defined by all the pointers in the program. We start from all the root pointers $t$ (those that are immediately accessible from the program) and invoke the following procedure for each. Let us assume that for each pointer $t$ we know its type and hence we know the number of pointers this object contains, denoted t.nChild and these pointers are denoted $t . c h i l d[i]$. (Note that although we use the term "child" as if pointers form a tree, this is not the case, since there may be cycles.) We assume that each block has a bit value t.isMarked which indicates whether the object has been marked.

```
mark(Pointer t) {
    if (t == null || t.isMarked) return // null or already visited
```

```
    t.isMarked = true // mark t visited
for (i = 0; i < t.nChild; i++)
    mark(t.child[i])
}
```

The recursive calls persist in visiting everything that is reachable, and only backing off when we come to a null pointer or something that has already been marked. Note that we do not need to store the $t . n C h i l d$ field in each object. It is function of $t$ 's type, which presumably the run-time system is aware of (especially in languages like Java that support dynamic casting). However we definitely need to allocate an extra bit in each object to store the mark.

Marking using Link Redirection: There is a significant problem in the above marking algorithm. This procedure is necessarily recursive, implying that we need a stack to keep track of the recursive calls. However, we only invoke this procedure if we have run out of space. So we do not have enough space to allocate a stack. This poses the problem of how can we traverse space without the use of recursion or a stack. There is no method known that is very efficient and does not use a stack. However there is a rather cute idea which allows us to dispense with the stack, provided that we allocate a few extra bits of storage in each object.
The method is called link redirection or the Schorr-Deutsch-Waite method. Let us think of our objects as being nodes in a multiway tree. (Recall that we have cycles, but because of we never revisit marked nodes, so we can think of pointers to marked nodes as if they are null pointers.) Normally the stack would hold the parent of the current node. Instead, when we traverse the link to the $i$ th child, we redirect this link so that it points to the parent. When we return to a node and want to proceed to its next child, we fix the redirected child link and redirect the next child.
Pseudocode for this procedure is given below. As before, in the initial call the argument $t$ is a root pointer. In general $t$ is the current node. The variable $p$ points to $t$ 's parent. We have a new field $t$. currChild which is the index of the current child of $t$ that we are visiting. Whenever the search ascends to $t$ from one of its children, we increment $t$. currChild and visit this next child. When $t . c u r r C h i l d==t . n C h i l d ~ t h e n ~ w e ~ a r e ~ d o n e ~ w i t h ~ t ~ a n d ~ a s c e n d ~ t o ~ i t s ~$ parent. We include two utilities for pointer redirection. The call descend ( $p, t, t . c$ ) moves us from $t$ to its child $t . c$ and saves $p$ in the pointer field containing $t . c$. The call ascend ( p , $\mathrm{t}, \mathrm{p} . \mathrm{c}$ ) moves us from $t$ to its parent $p$ and restores the contents of the $p$ 's child p.c. Each are implemented by a performing a cyclic rotation of these three quantities. (Note that the arguments are reference arguments.) An example is shown in the figure below, in which we print the state after each descend or ascend.

$$
\operatorname{descend}(p, t, t . c):\left(\begin{array}{c}
p \\
t \\
t . c
\end{array}\right) \longleftarrow\left(\begin{array}{c}
t \\
t . c \\
p
\end{array}\right) \quad \operatorname{ascend}(p, t, p . c):\left(\begin{array}{c}
p \\
t \\
p . c
\end{array}\right) \longleftarrow\left(\begin{array}{c}
p . c \\
p \\
t
\end{array}\right)
$$

_Marking using Link Redirection

```
markByRedirect(Pointer t) {
    p = null // parent pointer
    while (true) {
        if (t != null && !t.isMarked) { // first time at t?
            t.isMarked = true
            if (t.nChild > 0) {
                t.currChild = 0
    // mark as visited
    // t has children
    // start with child 0
    // descend via child 0
        }
```

```
            }
                                else if (p != null) { // returning to t
    j = p.currChild // parent's current child
    ascend(p, t, p.child[j]) // ascend via child j
    j = ++t.currChild // next child
    if (j < t.nChild) { // more children left
                descend(p, t, t.child[j]) // descend via child j
    }
        }
        else return // no parent? we're done
        }
}
```


ascend


Figure 73: Marking using link redirection.
As we have described it, this method requires that we reserve enough spare bits in each object to be able to keep track of which child we are visiting in order to store the t.currChild value. If an object has $k$ children, then $\lceil\lg k\rceil$ bits would be needed (in addition to the mark bit). Since most objects have a fairly small number of pointers, this is a small number of bits. Since we only need to use these bits for the elements along the current search path, rather than storing them in the object, they could instead be packed together in the form of a stack. Because we only need a few bits per object, this stack would require much less storage than the one needed for the recursive version of mark.

Stop-and-Copy: The alternative strategy to mark-and-sweep is called stop-and-copy. This method achieves much lower memory allocation, but provides for very efficient allocation. Stop-andcopy divides memory into two large banks of equal size. One of these banks is active and contains all the allocated blocks, and the other is entirely unused, or dormant. Rather than maintaining an available space list, the allocated cells are packed contiguously, one after the other at the fron of the current bank. When the current bank is full, we stop and determine which blocks in the current bank are reachable or alive. For each such live block we find, we copy it from the active bank to the dormant bank (and mark it so it is not copied again). The copying process packs these live blocks one after next, so that memory is compacted in the process, and hence there is no fragmentation. Once all the reachable blocks have been copied, the roles of the two banks are swapped, and then control is returned to the program.

Because storage is always compacted, there is no need to maintain an available space list. Free space consists of one large contiguous chunk in the current bank. Another nice feature of this method is that it only accesses reachable blocks, that is, it does not need to touch the garbage. (In mark-and-sweep the sweeping process needs to run through all of memory.) However, one shortcoming of the method is that only half of the available memory is usable at any given time, and hence memory utilization cannot exceed one half.

The trickiest issue in stop-and-copy is how to deal with pointers. When a block is copied from one bank to the other, there may be pointers to this object, which would need to be redirected. In order to handle this, whenever a block is copied, we store a forwarding link in the first word of the old block, which points to the new location of the block. Then, whenever we detect a pointer in the current object that is pointing into the bank of memory that is about to become dormant, we redirect this link by accessing the forwarding link. An example is shown in the figure below. The forwarding links are shown as broken lines.


Figure 74: Stop-and-copy and pointer redirection.
Which is better, mark-and-sweep or stop-and-copy? There is no consensus as to which is best in all circumstances. The stop-and-copy method seems to be popular in systems where it is easy to determine which words are pointers and which are not (as in Java). In languages such as $\mathrm{C}++$ where a pointer can be cast to an integer and then back to a pointer, it may be very hard to determine what really is a pointer and what is not, and so mark-and-sweep is more popular here. Stop-and-copy suffers from lots of data movement. To save time needed for copying long-lived objects (say those that have survived 3 or more garbage collections), we may declare them to be immortal and assign them to a special area of memory that is never garbage collected (unless we are really in dire need of space).

## Lecture 25: Tries and Digital Search Trees

(Tuesday, May 8, 2001)
Reading: Section 5.3 in Samet's notes. (The material on suffix trees is not covered there.)
Strings and Digital Data: Earlier this semester we studied data structures for storing and retrieving data from an ordered domain through the use of binary search trees, and related data
structures such as skip lists. Since these data structures can store any type of sorted data, they can certainly be used for storing strings. However, this is not always the most efficient way to access and retrieve string data. One reason for this is that unlike floating point or integer values, which can be compared by performing a single machine-level operation (basically subtracting two numbers and comparing the sign bit of the result) strings are compared lexicographically, that is, character by character. Strings also possess additional structure that simple numeric keys do not have. It would be nice to have a data structure which takes better advantage of the structural properties of strings.
Character strings arise in a number of important applications. These include language dictionaries, computational linguistics, keyword retrieval systems for text databases and web search, and computational biology and genetics (where the strings may be strands of DNA encoded as sequences over the alphabet $\{C, G, T, A\}$ ).

Tries: As mentioned above, our goal in designing a search structure for strings is to avoid having to look at every character of the query string at every node of the data structure. The basic idea common to all string-based search structures is the notion of visiting the characters of the search string from left to right as we descend the search structure. The most straightforward implementation of this concept is a trie. The name is derived from the middle syllable of "retrieval", but is pronounced the same as "try". Each internal node of a trie is $k$-way rooted tree, which may be implemented as an array whose length is equal to the number of characters $k$ in the alphabet. It is common to assume that the alphabet includes a special character, indicated by '.' in our figures, which represents a special end of string character. (In languages such as C++ and Java this would normally just be the null character.) Thus each path starting at the root is associated with a sequence of characters. We store each string along the associated path. The last pointer of the sequence (associated with the end of string character) points to a leaf node which contains the complete data associated with this string. An example is shown in the figure below.


Figure 75: A trie containing the strings: est, este, ete, set, sets, stet, test and tete. Only the nonnull entries are shown.

The obvious disadvantage of this straightforward trie implementation is the amount of space it uses. Many of the entries in each of the arrays is empty. In most languages the distribution of characters is not uniform, and certain character sequences are highly unlikely. There are a number of ways to improve upon this implementation.

For example, rather than allocating nodes using the system storage allocation (e.g., new) we store nodes in a 2-dimensional array. The number of columns equals the size of the alphabet, and the number of rows equals the total number of nodes. The entry $T[i, j]$ contains the index of the $j$-th pointer in node $i$. If there are $m$ nodes in the trie, then $\lceil\lg m\rceil$ bits suffice to represent each index, which is much fewer than the number of bits needed for a pointer.
de la Brandais trees: Another idea for saving space is, rather than storing each node as an array whose size equals the alphabet size, we only store the nonnull entries in a linked list. Each entry of the list contains a character, a child link, and a link to the next entry in the linked list for the node. Note that this is essentially just a first-child, next-sibling representation of the trie structure. These are called de la Brandais trees. An example is shown in the figure below.


Figure 76: A de la Brandais tree containing the strings: est, este, ete, set, sets, stet, test and tete.
Although de la Brandais trees have the nice feature that they only use twice as much pointer space as there are characters in the strings, the search time at each level is potentially linear in the number of characters in the alphabet. A hybrid method would be to use regular trie nodes when the branching factor is high and then convert to de la Brandais trees when the branching factor is low.

Patricia Tries: In many applications of strings there can be long sequences of repeated substrings. As an extreme example, suppose that you are creating a trie with the words "demystificational" and "demystifications" but no other words that contain the prefix "demys". In order to stores these words in a trie, we would need to create 10 trie nodes for the common substring "tification", with each node having a branching factor of just one each. To avoid this, we would like the tree to inform us that after reading the common prefix "demys" we should skip over the next 10 characters, and check whether the 11 th is either ' $a$ ' or 's'. This is the idea behind patricia tries. (The word 'patricia' is an acronym for Practical Algorithm To Retrieve Information Coded In Alphanumeric.)
A patricia trie uses the same node structure as the standard trie, but in addition each node contains an index field, indicating which character is to be checked at this node. This index field value increases as we descend the tree, thus we still process strings from left to right. However, we may skip over any characters that provide no discriminating power between keys. An example is shown below. Note that once only one matching word remains, we can proceed immediately to a leaf node.
Observe that because we skip over characters in a patricia trie, it is generally necessary to verify the correctness of the final result. For example, if we had attempted to search for the word "survive" then we would match 's' at position 1, 'i' at position 5, and 'e' at position 7,


Figure 77: A patricia trie for the strings: essence, essential, estimate, estimation, sublease, sublime, subliminal.
and so we arrive at the leaf node for "sublime". This means that "sublime" is the only possible match, but it does not necessarily match this word. Hence the need for verification.

Suffix trees: In some applications of string pattern matching we want to perform a number of queries about a single long string. For example, this string might be a single DNA strand. We would like to store this string in a data structure so that we are able to perform queries on this string. For example, we might want to know how many occurrences there are of some given substring within this long string.
One interesting application of tries and patricia tries is for this purpose. Consider a string $s=" a_{1} a_{2} \ldots a_{n}$. . We assume that the $(n+1)$-st character is the unique string termination character. Such a string implicitly defines $n+1$ suffixes. The $i$-th suffix is the string " $a_{i} a_{i+1} \ldots a_{n}$.". For each position $i$ there is a minimum length substring starting at position $i$ which uniquely identifies this substring. For example, consider the string "yabbadabbadoo". The substring " $y$ " uniquely identifies the first position of the string. However the second position is not uniquely identified by "a" or "ab" or even "abbad", since all of these substrings occur at least twice in the string. However, "abbada" uniquely identifies the second position, because this substring occurs only once in the entire string. The substring identifier for position $i$ is the minimum length substring starting at position $i$ of the string which occurs uniquely in $s$. Note that because the end of string character is unique, every position has a unique substring identifier. An example is shown in the following figure.

A suffix tree for a string $s$ is a trie in which we store each of the $n+1$ substring identifiers for $s$. An example is shown in the following figure. (Following standard suffix tree conventions, we put labels on the edges rather than in the nodes, but this data structure is typically implemented as a trie or a patricia trie.)
As an example of answering a query, suppose that we want to know how many times the substring "abb" occurs within the string $s$. To do this we search for the string "abb" in the suffix tree. If we fall out of the tree, then it does not occur. Otherwise the search ends at some node $u$. The number of leaves descended from $u$ is equal to the number of times "abb" occurs within $s$. (In the example, this is 2.) By storing this information in each node of the tree, we can answer these queries in time proportional to the length of the substring query (irrespective of the length of $s$ ).
Since suffix trees are often used for storing very large texts upon which many queries are to be performed (e.g. they were used for storing the entire contents of the Oxford English Dictionary and have been used in computational biology) it is important that the space used by the data structure be $O(n)$ where $n$ is the number of characters in the string. Using the standard trie representation, this is not necessarily true. (You can try to generate your own

| Position | Substring identifier |
| :---: | :--- |
| 1 | y |
| 2 | abbada |
| 3 | bbada |
| 4 | bada |
| 5 | ada |
| 6 | da |
| 7 | abbado |
| 8 | bbado |
| 9 | bado |
| 10 | ado |
| 11 | do |
| 12 | oo |
| 13 | o. |
| 14 | . |

Figure 78: Substring identifiers for the string "yabbadabbadoo.".


Figure 79: A suffix tree for the string "yabbadabbadoo.".
counterexample where the data structure uses $O\left(n^{2}\right)$ space, even if the alphabet is limited to 2 symbols.) However, if a patricia trie is used the resulting suffix tree has $O(n)$ nodes. The reason is that the number of leaves in the suffix tree is equal to $n+1$. We showed earlier in the semester that if every internal node has two children (or generally at least two children) then the number of internal nodes is not greater than $n$. Hence the total number of nodes is $O(n)$.

## Lecture 26: Hashing

(Thursday, May 10, 2001)
Reading: Chapter 5 in Weiss and Chapter 6 in Samet's notes.
Hashing: We have seen various data structures (e.g., binary trees, AVL trees, splay trees, skip lists) that can perform the dictionary operations insert(), delete() and find(). We know that these data structures provide $O(\log n)$ time access. It is unreasonable to ask any sort of treebased structure to do better than this, since to store $n$ keys in a binary tree requires at least $\Omega(\log n)$ height. Thus one is inclined to think that it is impossible to do better. Remarkably, there is a better method, at least if one is willing to consider expected case rather than worst case performance.
Hashing is a method that performs all the dictionary operations in $O(1)$ (i.e. constant) expected time, under some assumptions about the hashing function being used. Hashing is considered so good, that in contexts where just these operations are being performed, hashing is the method of choice (e.g. symbol tables for compilers are almost always implemented using hashing). Tree-based data structures are generally prefered in the following situations:

- When storing data on secondary storage (e.g. using B-trees),
- When knowledge of the relative order of elements is important (e.g. if a find() fails, I may want to know the nearest key. Hashing cannot help us with this.)

The idea behind hashing is very simple. We have a table containing $m$ entries. We select a hash function $h(x)$, which is an easily computable function that maps a key $x$ to a "virtually random" index in the range [0..m-1]. We will then attempt to store the key in index $h(x)$ in the table. Of course, it may be that different keys are mapped to the same location. This is called a collision. We need to consider how collisions are to be handled, but observe that if the hashing function does a good job of scattering the keys around the table, then the chances of a collision occuring at any index of the table are about the same. As long as the table size is at least as large as the number of keys, then we would expect that the number of keys that are map to the same cell should be small.

Hashing is quite a versatile technique. One way to think about hashing is as a means of implementing a content-addressable array. We know that arrays can be addressed by an integer index. But it is often convenient to have a look-up table in which the elements are addressed by a key value which may be of any discrete type, strings for example or integers that are over such a large range of values that devising an array of this size would be impractical. Note that hashing is not usually used for continuous data, such as floating point values, because similar keys 3.14159 and 3.14158 may be mapped to entirely different locations.

There are two important issues that need to be addressed in the design of any hashing system. The first is how to select a hashing function and the second is how to resolve collisions.

Hash Functions: A good hashing function should have the following properties.

- It should be simple to compute (using simple arithmetic operations ideally).
- It should produce few collisions. In turn the following are good rules of thumb in the selection of a hashing function.
- It should be a function of every bit of the key.
- It should break up naturally occuring clusters of keys.

As an example of the last rule, observe that in writing programs it is not uncommon to use very similar variables names, temp1, temp2, and temp3. It is important such similar names be mapped to entirely different locations.
We will assume that our hash functions are being applied to integer keys. Of course, keys need not be integers generally. But noninteger data, such as strings, can be thought of as a sequence of integers assuming their ASCII or UNICODE encoding. Once our key has been converted into an integer, we can think of hash functions on integers. One very simple hashing function is the function

$$
h(x)=x \bmod m
$$

This certainly maps each key into the range $[0 . . m-1]$, and it is certainly fast. Unfortunately this function is not a good choice when it comes to collision properties, since it does not break up clusters.

A more robust strategy is to first multiply the key value by some large integer constant $a$ and then take the mod. For this to have good scattering properties either $m$ should be chosen to be a prime number, or $a$ should be prime relative to $m$ (i.e. share no common divisors other than 1).

$$
h(x)=(a \cdot x) \bmod m .
$$

An even better approach is to both add and multiply. Let $a$ and $b$ be two large integer constants. Ideally $a$ is prime relative to $m$.

$$
h(x)=(a x+b) \bmod m .
$$

By selecting $a$ and $b$ at random, it can be shown such a scheme produces a good performance in the sense that for any two keys the probability of them being mapped to the same address is roughly $1 / m$. In our examples, we will simplify things by taking the last digit as the hash value, although in practice this is a very bad choice.

Collision Resolution: Once a hash function has been selected, the next element that has to be solved is how to handle collisions. If two elements are hashed to the same address, then we need a way to resolve the situation.

Separate Chaining: The simplest approach is a method called separate chaining. The idea is that we think of each of the $m$ locations of the hash table as simple head pointers to $m$ linked lists. The link list table[i] holds all keys that hash to location $i$.
To insert a key $x$ we simply compute $h(x)$ and insert the new element into the linked list table $\mathrm{h}(\mathrm{x})]$. (We should first search the list to make sure it is not a duplicate.) To find a key we just search this linked list. To delete a key we delete it from this linked list. An example is shown below, where we just use the last digit as the hash function (a very bad choice normally).
The running time of this procedure will depend on the length of the linked list to which the key has been hashed. If $n$ denotes the number of keys stored in the table currently, then the ratio $\lambda=n / m$ indicates the load factor of the hash table. If we assume that keys are being hashed roughly randomly (so that clustering does not occur), then it follows that each linked list is expected to contain $\lambda$ elements. As mentioned before, we select $m$ to be with a constant factor of $n$, so this ratio is a constant.


Figure 80: Collision resolution by separate Chaining.

Thus, it follows from a straightforward analysis that the expected running time of a successful search is roughly

$$
S_{c h}=1+\frac{\lambda}{2}=O(1)
$$

since about half of an average list needs be searched. The running time of an unsuccessful search is roughly

$$
U_{c h}=1+\lambda=O(1)
$$

Thus as long as the load factor is a constant separate chaining provide expected $O(1)$ time for insertion and deletion.

The problem with separate chaining is that we require separate storage for pointers and the new nodes of the linked list. This also creates additional overhead for memory allocation. It would be nice to have a method that does not require the use of any pointers.

Open Addressing: To avoid the use of extra pointers we will simply store all the keys in the hash table, and use a special value (different from every key) called EMPTY, to determine which entries have keys and which do not. But we will need some way of finding out which entry to go to next when a collision occurs. Open addressing consists of a collection of different strategies for finding this location. In it most general form, an open addressing system involves a secondary search function, $f$, and if we find that location $h(x)$ is occupied, we next try locations

$$
(h(x)+f(1)) \bmod m,(h(x)+f(2)) \bmod m,(h(x)+f(3)) \bmod m, \ldots
$$

until finding an open location. This sequence is called a probe sequence, and ideally it should be capable of searching the entire list. How is this function $f$ chosen? There are a number of alternatives, which we consider below.

Linear Probing: The simplest idea is to simply search sequential locations until finding one that is open. Thus $f(i)=i$. Although this approach is very simple, as the table starts to get full its performance becomes very bad (much worse than chaining).
To see what is happening let's consider an example. Suppose that we insert the following 4 keys into the hash table (using the last digit rule as given earlier): 10, 50, 42, 92. Observe that the first 4 locations of the hash table are filled. Now, suppose we want to add the key 31. With chaining it would normally be the case that since no other key has been hashed to location 1, the insertion can be performed right away. But the bunching of lists implies that we have to search through 4 cells before finding an available slot.
This phenomenon is called secondary clustering. Primary clustering happens when the table contains many names with the same hash value (presumably implying a poor choice for the hashing function). Secondary clustering happens when keys with different hash values have

| insert(50) |  | insert(42) |  | insert(92) |  | insert(31) |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $\rightarrow 0$ | 10 | 0 | 10 | 0 | 10 | 0 | 10 |
| 1 | 50 | 1 | 50 | 1 | 50 | $\rightarrow 1$ | 50 |
| 2 |  | $\rightarrow 2$ | 42 | $\rightarrow 2$ | 42 | 2 | 42 |
| 3 |  | 3 |  | 3 | 92 | 3 | 92 |
| 4 |  | 4 |  | 4 |  | 4 | 31 |
| 5 |  | 5 |  | 5 |  | 5 |  |
| 6 |  | 6 |  | 6 |  | 6 |  |

Figure 81: Linear probing.
nearly the same probe sequence. Note that this does not occur in chaining because the lists for separate hash locations are kept separate from each other, but in open addressing they can interfere with each other.

As the table becomes denser, this affect becomes more and more pronounced, and it becomes harder and harder to find empty spaces in the table.

Recall that $\lambda=n / m$ is the load factor for the table. For open addressing we require that $\lambda \leq 1$, because the table cannot hold more than $m$ entries. It can be shown that the expected running times of a successful and unsuccessful searches using linear probing are

$$
\begin{aligned}
S_{l p} & =\frac{1}{2}\left(1+\frac{1}{1-\lambda}\right) \\
U_{l p} & =\frac{1}{2}\left(1+\left(\frac{1}{1-\lambda}\right)^{2}\right)
\end{aligned}
$$

This is quite hard to prove. Observe that as $\lambda$ approaches 1 (a full table) this grows to infinity. A rule of thumb is that as long as the table remains less than $75 \%$ full, linear probing performs fairly well.

Quadratic Probing: To avoid secondary clustering, one idea is to use a nonlinear probing function which scatters subsequent probes around more effectively. One such method is called quadratic probing, which works as follows. If the index hashed to $h(x)$ is full, then we consider next $h(x)+1, h(x)+4, h(x)+9, \ldots$ (again taking indices mod $m$ ). Thus the probing function is $f(i)=i^{2}$.

Here is the search algorithm (insertion is similar). Note that there is a clever trick to compute $i^{2}$ without using multiplication. It is based on the observation that $i^{2}=(i-1)^{2}+2 i-1$. Therefore, $f(i)=f(i-1)+2 i-1$. Since we already know $f(i-1)$ we only have to add in the $2 i-1$. Since multiplication by 2 can be performed by shifting this is more efficient. The argument $x$ is the key to find, $T$ is the table, and $m$ is the size of the table. We will just assume we are storing integers, but the extension to other types is straightforward.

Hashing with Quadratic Probing

```
int find(int x, int T[m]) {
    i = 0
    c = h(x) // first position
    while (T[c] != EMPTY) && (T[c] != x) { // found key or empty slot
        c += 2*(++i) - 1 // next position
```

```
        c=c % m // wrap around using mod
        }
        return c
}
```



Figure 82: Various open-addressing systems.
The above procedure is not quite complete, since it loops infinitely when the table is full. This is easy to fix, by adding a variable counting the number of entries being used in the table.
Experience shows that this succeeds in breaking up the secondary clusters that arise from linear probing, but there are some tricky questions to consider. With linear probing we were assured that as long as there is one free location in the array, we will eventually find it without repeating any probe locations. How do we know if we perform quadratic probing that this will be the case? It might be that we keep hitting the same index of the table over and over (because of the fact that we take the index $\bmod m$ ).
It turns out (fortunately) that quadratic probing does do a pretty good job of visiting different locations of the array without repeating. It can even be formally proved that if $m$ is prime, then the first $m / 2$ locations that quadratic probing hits will be distinct. See Weiss for a proof.

Double Hashing: As we saw, the problem with linear probing is that we may see clustering or piling up arise in the table. Quadratic probing was an attractive way to avoid this by scattering the successive probe sequences around. If you really want to scatter things around for probing, then why don't you just use another hashing function to do this?
The idea is use $h(x)$ to find the first location at which to start searching, and then let $f(i)=$ $i \cdot g(x)$ be the probing sequence where $g(x)$ is another hashing function. Some care needs to be taken in the choice of $g(x)$ (e.g. $g(x)=0$ would be a disaster). As long as the table size $m$ is prime and $(g(x) \bmod m \neq 0$ we are assured of visiting all cells before repeating. Note that
the second hash function does not tell us where to put the object, it gives us an increment to use in cycling around the table.

Performance: Performance analysis shows that as long as primary clustering is avoided, then open addressing using is an efficient alternative to chaining. The running times of successful and unsuccessful searches for open addressing using double hashing are

$$
\begin{aligned}
S_{d h} & =\frac{1}{\lambda} \ln \frac{1}{1-\lambda} \\
U_{d h} & =\frac{1}{1-\lambda}
\end{aligned}
$$

To give some sort of feeling for what these quantities mean, consider the following table.

| $\lambda$ | 0.50 | 0.75 | 0.90 | 0.95 | 0.99 |
| :--- | :--- | :--- | :--- | :--- | :--- |
| $U(\lambda)$ | 2.00 | 4.00 | 10.0 | 20.0 | 100. |
| $S(\lambda)$ | 1.39 | 1.89 | 2.56 | 3.15 | 4.65 |

## Lecture 27: Final Review

(Tuesday, May 15, 2001)
Final Overview: This semester we have covered a number of different data structures. The fundamental concept behind all data structures is that of arranging data in a way that permits a given set of queries or accesses to be performed efficiently. The aim has been at describing the general principles that lead to good data structures as opposed to giving a "cookbook" of standard techniques. In addition we have discussed the mathematics needed to prove properties and analyze the performance of these data structures (both theoretically and empirically).
Some of the important concepts that I see that you should take away from this course are the following:

Mathematical Models/Objects: Before you design any data structure or algorithm, first isolate the key mathematical elements of the task at hand. Identify what the mathematical objects are, and what operations you are performing on them. (E.g. Dictionary: insertion, deletion, and finding of numeric keys. Ray Tracing: insertion, and ray tracing queries for a set of spheres in 3 -space.)
Recursive Subdivision: A large number of data structures (particularly tree based structures) have been defined recursively, by splitting the underlying domain using a key value. The algorithms that deal with them are often most cleanly conceptualized in recursive form.
(Strong) Induction: In order to prove properties of recursively defined objects, like trees, the most natural mathematical technique is induction. In particular, strong induction is important when dealing with trees.
Balance: The fundamental property on which virtually all good data structures rely is a notion of information balancing. Balance is achieved in many ways (rebalancing by rotations, using hash functions to distribute information randomly, balanced merging in Union-Find trees).
Amortization: When the running time of each operation is less important than the running time of a string of operations, you can often arrive at simpler data structures and algorithms. Amortization is a technique to justify the efficiency of these algorithms, e.g. through the use of potential function which measures the imbalance present in a structure and how your access functions affect this balance.

Randomization: Another technique for improving the performance of data structures is: if you can't randomize the data, randomize the data structure. Randomized data structures are simpler and more efficient than deterministic data structures which have to worry about worst case scenarios.

Asymptotic analysis: Before fine-tuning your algorithm or data structure's performance, first be sure that your basic design is a good one. A half-hour spent solving a recurrence to analyze a data structure may tell you more than 2 days spent prototyping and profiling.
Empirical analysis: Asymptotics are not the only answer. The bottom line: the algorithms have to perform fast on my data for my application. Prototyping and measuring performance parameters can help uncover hidden inefficiencies in your data structure.

Topics: Since the midterm, these are the main topics that have been covered.
Priority Queues: We discussed leftist heaps and skew-heaps. Priority queues supported the operations of insert and deleteMin in $O(\log n)$ time. Leftist and skew heaps also supported the operation merge. An interesting feature of leftist heaps is that the entire tree does not need to balanced. Only the right paths need to be short, in particular the rightmost path of the tree is of length $O(\log n)$. Skew heaps are a self-adjusting version of leftist heaps.
Disjoint Set Union/Find: We discussed a data structure for maintaining partitions of sets. This data structure can support the operations of merging to sets together (Union) and determining which set contains a specific element (Find).
Geometric data structures: We discussed kd-trees and range trees for storing points. We also discussed interval trees for storing intervals (and a variant for storing horizontal line segments). In all cases the key concepts are those of (1) determining a way of subdividing space, and (2) in the case of query processing, determining which subtrees are relevant to the search and which are not. Although kd-trees are very general, they are not always the most efficient data structure for geometric queries. Range trees and interval trees are much more efficient for their specific tasks (range queries and stabbing queries, respectively) but they are of much more limited applicability.
Tries: We discussed a number of data structures for storing strings, including tries, de la Brandais tries, patricia tries, and suffix trees. The idea behind a trie is that of traversing the search structure as you traverse the string. Standard tries involved a lot of space. The de la Brandais trie and patricia trie were more space-efficient variants of this idea. The suffix trees is a data structure for storing all the identifying substrings of a single string, which is used for answering queries about this string.
Memory management and garbage collection: We discussed methods for allocating and deallocating blocks of storage from memory. We discussed the standard method for memory allocation and the buddy system. We also discussed two methods for garbage collection, mark-and-sweep and stop-and-copy.

Hashing: Hashing is considered the simplest and most efficient technique (from a practical perspective) for performing dictionary operations. The operations insert, delete, and find can be performed in $O(1)$ expected time. The main issues here are designing a good hash function which distributes the keys in a nearly random manner, and a good collision resolution method to handle the situation when keys hash to the same location. In open addressing hashing, clustering (primary and secondary) are significant issues. Hashing provides the fastest access for exact searches. For continuous types of queries (range searching and nearest neighbor queries), trees are preferable because they provide ordering information.

## Lecture X01: Red-Black Trees

## (Supplemental)

Reading: Section 5.2 in Samet's notes.
Red-Black Trees: Red-Black trees are balanced binary search trees. Like the other structures we have studied (AVL trees, B-trees) this structure provides $O(\log n)$ insertion, deletion, and search times. What makes this data structure interesting is that it bears a similarity to both AVL trees (because of the style of manipulation through rotations) and B-trees (through a type of isomorphism). Our presentation follows Samet's. Note that this is different from usual presentations, since Samet colors edges, and most other presentations color vertices.
A red-black tree is a binary search tree in which each edge of the tree is associated with a color, either red or black. Each node $t$ of the data structure contains an additional one-bit field, $t . c o l o r$, which indicates the color of the incoming edge from the parent. We will assume that the color of nonexistent edge coming into the root node is set to black. Define the black-depth of a node to be the number of black edges traversed from the root to that node. A red-black tree has the following properties, which assure balance.
(1) Consider the external nodes in the augmented binary tree (where each null pointer in the original tree is replaced by a black edge to a special node). The black-depths of all these external nodes are equal. (Note: The augmented tree is introduced for the purposes of definition only, we do not actually construct these external nodes in our representation of red-black trees.)
(2) No path from the root to a leaf has two consecutive red edges.
(3) (Optional) The two edges leaving a node to its children cannot both be red.

Property (3) can be omitted without altering the basic balance properties of red-black trees (but the rebalancing procedures will have to be modified to handle this case). We will assume that property (3) is enforced. An example is shown in the figure below. Red edge are indicated with dashed lines and black edges with solid lines. The small nodes are the external nodes (not really part of the tree). Note that all these nodes are of black-depth 3 (counting the edge coming into them).


Figure 83: Red-black tree.
Red-black trees have obvious connections with 2-3 trees (B-trees of order $m=3$ ), and 2-3-4 trees (B-trees of order $m=4$ ). When property (3) is enforced, the red edges are entirely isolated from one another. By merging two nodes connected by a red-edge together into a
single node, we get a 3 -node of a $2-3$ tree tree. This is shown in part (a) of the figure below. (Note that there is a symmetric case in which the red edge connects a right child. Thus, a given 2-3 tree can be converted into many different possible red-black trees.)

(a)

(b)

Figure 84: Red-black and 2-3 or 2-3-4 trees.
In the case where property (3) is not enforced, then it is possible to have a two red sibling edges. If we join all three nodes connected by these edges we get a 4-node (see part (b) of the same figure). Thus by enforcing property (3) we get something isomorphic to a $2-3$ tree, and by not enforcing it we get something isomorphic to $2-3-4$ trees. For example, the red-black tree given earlier could be expressed as the following (2-3)-tree.


Figure 85: Equivalent (2-3) tree.
Thus red-black trees are certainly similar to B-trees. One reason for studying red-black trees independently is that when storing B-trees in main memory (as opposed to the disk) red-black trees have the nice property that every node is fully utilized (since B-tree nodes can be as much as half empty). Another reason is to illustrate the fact that there are many ways to implement a single "concept", and these other implementations may be cleaner in certain circumstances. Unfortunately, red-black trees cannot replace B-trees when it comes to storing data on disks.
It follows that the height of a red-black tree of $n$ nodes is $O(\log n)$, since these trees are essentially the same as 2-3 or 2-3-4 trees, which also have logarithmic height.

Red-Black Insertion: It should be obvious from the discussion above that the methods used to
restore balance in red-black tree should be an easy extension of rebalancing in B-trees. We simply see what a B-tree would do in a particular circumstance, and figure out what the equivalent rebalancing in the red-black tree is needed. Here is where the similarity with AVL trees comes in, because it turns out that the fundamental steps needed to rebalance red-black trees are just rotations. Let's start by considering node insertion.

To insert an key $K$ into a red-black tree, we begin with the standard binary tree insertion. Search for the key until we either find the key, or until we fall out of the tree at some node $x$. If the key is found, generate an error message. Otherwise create a new node $y$ with the new key value $K$ and attach it as a child of $x$. Assign the color of the newly created edge $(y, x)$ to be red. Observe at this point that the black-depth to every extended leaf node is unchanged (since we have used a red edge for the insertion), but we may have have violated the red-edge constraints. Our goal will be to reestablish the red-edge constraints. We backtrack along the search path from the point of insertion to the root of the tree, performing rebalancing operations.

Our description of the algorithm will be nonrecursive, because (as was the case with splay trees) some of the operations require hopping up two levels of the tree at a time. The algorithm begins with a call to insert (), which is the standard (unbalanced) binary tree insertion. We assume that this routine returns a pointer y to the newly created node. In the absence of recursion, walking back to the root requires that we save the search path (e.g. in a stack, or using parent pointers). Here we assume parent pointers (but it could be done either way). We also assume we have two routines leftRotate() and rightRotate(), which perform a single left and single right rotation, respectively (and update parent pointers).
Whenever we arrive at the top of the while-loop we make the following assumptions:

- The black-depths of all external leaf nodes in the tree are equal.
- The edge to $y$ from its parent is red (that is, y.color $==$ red).
- All edges of the tree satisfy the red-conditions (2) and (3) above, except possibly for the edge to $y$ from its parent.

```
RBinsert(Object K, RBNode root) {
    y = insert(K, root) // standard insert; let y be new node
    y.color = red // y's incoming edge is red
    while (y != root) { // repeat until getting back to root
        x = y.parent // x is y's parent
        if (x.color == black) { // edge into x is black?
            z = otherChild(x,y) // z is x's other child
            if (z == null || z.color == black)
                return // no red violation - done
            else { // Case 1
                y.color = black // reverse colors
                    z.color = black
                x.color = red
                y = x // move up one level higher
            }
        }
        else { // edge into x is red
            w = x.parent // w is x's parent
            if (x == w.left) { // x is w's left child
                    if (y == x.right) { // Case 2(a): zig-zag path
                leftRotate(x) // make it a zig-zig path
                swap(x,y) // swap pointers x and y
            }
```

```
            rightRotate(w) // Case 2(b): zig-zig path
        }
        else { // x is w's right child
            ...this case is symmetric...
        }
        y.color = black // reverse colors
        w.color = black
        x.color = red
        y = x; // move up two levels
        }
    }
    root.color = black; // edge above root is black
}
```

Here are the basic cases considered by the algorithm. Recall that $y$ is the current node, whose incoming edge is red, and may violate the red-constraints. We let $x$ be $y$ 's parent, and $z$ be the other child of $x$.

Case 1: (Edge above $x$ is black:) There are two subcases. If $z$ is null or the edge above $z$ is black, then we are done, since edge $(x, y)$ can have no red neighboring edges, and hence cannot violate the red-constraints. Return.
Otherwise if $(x, z)$ exists and is red, then we have violated property (3). Change the color of edges above $y$ and $z$ to be black, change the color of the edge above $x$ to be red. (Note: This corresponds to a node split in a 2-3 tree.)


Figure 86: Case 1.
Case 2: (Edge above $x$ is red:) Since the root always has an incoming black edge we know $x$ is not the root, so let $w$ be $x$ 's parent. In this case we have two consecutive red edges on a path. If the path is a zig-zag path (Case 2(a)), we apply a single rotation at $x$ and swap $x$ and $y$ to make it a zig-zig path, producing Case $2(\mathrm{~b})$. We then apply a single rotation at $w$ to make the two red edges siblings of each other, and bringing $x$ to the top. Finally we change the lower edge colors to black, and change the color of the edge above $x$ to be red. (Note: This also corresponds to a node split in a 2-3 tree.)

We will not discuss red-black deletion, but suffice it to say that it is only a little more complicated, but operates in essentially the same way.

## Lecture X02: BB-trees

(Supplemental)
Reading: This is not covered in our readings.
BB-trees: We introduced B-trees earlier, and mentioned that they are good for storing data on disks. However, B-trees can be stored in main memory as well, and are an alternative to AVL


Figure 87: Case 2.
trees, since both guarantee $O(\log n)$ wost case time for dictionary operations (not amortized, not randomized, not average case). Unfortunately, implementing B-trees is quite a messy programming task in general. BB-trees are a special binary tree implementation of 2-3 trees, and one of their appealing aspects is that they are very easy to code (arguably even easier than AVL trees).
Recall that in a $2-3$ tree, each node has either 2 or 3 children, and if a node has $j$ children, then it has $j-1$ key values. We can represent a single node in a 2-3 tree using either 1 or 2 nodes in a binary tree, as illustrated below.


Figure 88: Binary representation of 2-3 tree nodes.
When two nodes of a BB-tree are used to represent a single node of a 2-3 tree, this pair of nodes is called pseudo-node. When we represent a $2-3$ tree using this binary representation, we must be careful to keep straight which pairs of vertices are pseudo-nodes. To do this, we create an additional field in each node that contains the level of the node in the 2-3 tree. The leaves of the 2-3 tree are at level 1, and the root is at the highest level. Two adjacent nodes in a BB-tree (parent and right child) that are of equal level form a single pseudo-node.
The term "BB-tree" stands for "binary B-tree". Note that in other textbooks there is a data structure called a "bounded balance" tree, which also goes by the name BB-tree. Be sure you are aware of the difference. BB-trees are closely related to red-black trees, which are a binary representation of 2-3-4 trees. However, because of their additional structure the code for BB-trees is quite a bit simpler than the corresponding code for red-black trees.
As is done with skip-lists, it simplifies coding to create a special sentinel node called nil. Rather than using null pointers, we store a pointer to the node nil as the child in each leaf. The level of the sentinel node is 0 . The left and right children of nil point back to nil.

The figure below illustrates a 2-3 tree and the corresponding BB-tree.
Note that the edges of the tree can be broken into two classes, vertical edges that correspond to 2-3 tree edges, and horizontal edges that are used to join the two nodes that form a pseudonode.


Figure 89: BB-tree corresponding to a 2-3 tree.

BB-tree operations: Since a BB-tree is essentially a binary search tree, find operations are no different than they are for any other binary tree. Insertions and deletions operate in essentially the same way they do for AVL trees, first insert or delete the key, and then retrace the search path and rebalance as you go. Rebalancing is performed using rotations. For BB-trees the two rotations go under the special names skew() and split. Their intuitive meanings are:
skew ( p ) : replace any horizontal left edge with a horizontal right edge by a right rotation at p.
$\operatorname{split}(p)$ : if a pseudo-node is too large (i.e. more than two consecutive nodes at the same level), then split it by increasing the level of every other node. This is done by making left rotations along a right path of horizontal edges.

Here is their implementation. Split only splits one set of three nodes.
BB-tree Utilities skew and split

```
BBNode skew(BBNode p) {
    if (p.left.level == p.level) {
            q = p.left
            p.left = q.right
            q.right = p
            return q
    }
    else return p
}
BBNode split(BBNode p) {
    if (p.right.right.level == p.level) {
        q = p.right
        p.right = q.left
        q.left = p
        q.level++
        return q
    }
    else return p
}
```

Insertion: Insertion performs in the usual way. We walk down the tree until falling out, and insert the new key at the point we fell out. The new node is at level 1. We return up the search path and rebalance. At each node along the search path it suffices to perform one skew and one split.


Figure 90: BB-tree insertion.

```
BBNode insert(int x, BBNode t) {
    if (t == nil) { // empty tree
        t = new BBNode(x,1,nil,nil) // new node a level 1
    }
    else {
        if (x < t.data)
            t.left = insert(x, t.left) // insert on left
        else if (x > t.data)
            t.right = insert(x, t.right) // insert on right
        else
            ...error: duplicate key...
        t = skew(t)
        t = split(t)
        }
        return t
}
```

Deletion: As usual deletion is more of a hassle. We first locate the node to be deleted. We replace it's key with an appropriately chose key from level 1 (which need not be a leaf) and proceed to delete the node with replacement key. We retrace the search path towards the root rebalancing along the way. At each node $p$ that we visit there are a number of things that can happen.
(a) If $p$ 's child is 2 levels lower, then $p$ drops a level. If $p$ is part of a pseudo-node of size 2 , then both nodes drop a level.
(b) Apply a sufficient number of skew operations to align the nodes at this level. In general 3 may be needed: one at $p$, one at $p$ 's right child, and one at $p$ 's right-right grandchild.
(c) Apply a sufficient number of splits to fix things up. In general, two may be needed: one at $p$, and one at $p$ 's right child.

Important note: We use 3 global variables: nil is a pointer to the sentinel node at level 0 , del is a pointer to the node to be deleted, repl is the replacement node. Before calling this routine from the main program, initialize del $=$ nil.

BB-Tree Deletion

```
BBNode delete(int x, BBNode t) {
    if (t != nil) { // search tree for repl and del
        repl = t
        if (x < t.data)
            t.left = delete(x, t.left)
        else {
            del = t
            t.right = delete(x, t.right)
        }
        if (t == repl) { // if at bottom remove item
            if ((del != nil) && (x == del.data)) {
            del.data = t.data
            del = nil
            t = t.right // unlink replacement
            delete repl // destroy replacement
            }
```



Figure 91: BB-tree deletion.

```
            else
            ...error: deletion of nonexistent key...
        }
                                    // lost a level?
            else if ((t.left.level < t.level - 1) ||
            (t.right.level < t.level - 1)) {
        t.level-- // drop down a level
        if (t.right.level > t.level) {
            t.right.level = t.level
        }
        t = skew(t)
        t.right = skew(t.right)
        t.right.right = skew(t.right.right)
        t = split(t)
        t.right = split(t.right)
        }
    }
return t
}
```


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