

Design & Analysis of Algorithms

Algorithm: An algorithm is a finite set of instructions that is to be carried out in order to solve a particular task (input each instruction tells what action is to be performed).

- An algorithm is a tool for solving a well-specified computational problem
- An algorithm is a sequence of unambiguous instructions for solving a problem

An algorithm is any well-defined computational procedure that takes some value or set of values as input and produces some value or set of values as output (An algorithm is thus a sequence of computational steps that transform the input into the output).

Ex:- for sorting problem

Input: A sequence of n numbers $\langle a_1, a_2, \dots, a_n \rangle$

Output: $\langle a'_1, a'_2, \dots, a'_n \rangle$ in increasing order.

Procedure → It is the code in an algorithm & that of program. That means both code & program.
 In general, the coding part of an algorithm is known as program.
 In general, an instance of a problem consists of input needed to compute a solⁿ to the problem.

→ Study of algorithm provides insight into the intrinsic of the problem as well as possible solution technique independent of ^{prior understanding}

- of:
- PL
 - programming paradigm
 - computer hardware
 - or any other implementation aspects.

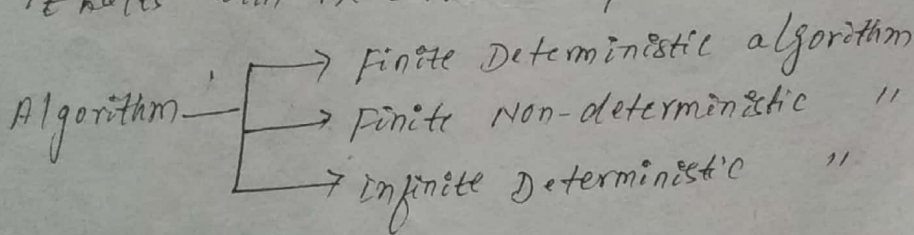
- An algorithm can be specified in any language.
- Normally algorithms are written in a pseudocode.
- pseudocode conveys the structure of the algorithm clearly enough that programmes can implement it in the language of his choice.

- The prefix pseudo is used to give the information that the code is not meant to be compiled and executed on a computer.
- It is easy to understand an algorithm by using pseudocode.
- The pseudocode hides the implementation details and thus one can easily focus on the computational aspects of an algorithm.
- Pseudo-code is a mixture of natural language and high-level programming constructs that describe the main idea behind a general implementation of algorithm.

Which algorithm is best for a given application depends on many factors

- time
- space
- no. of items to be sorted
- the extent to which the item is ^{somehow} already sorted
- possible restrictions on the item values
- kind of storage device used (main memory / disk / tapes)

→ An algorithm is said to be correct if, for every input sequence, it halts with the correct output.

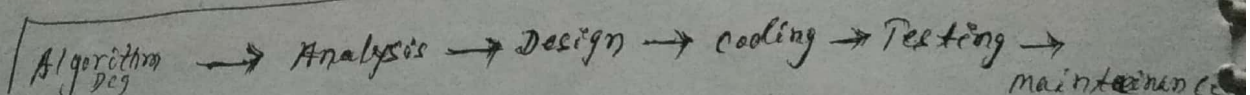


Finite Deterministic Algo Algo terminate in a finite amount of time. It always gives result that is uniquely dependent on the input.
Ex:- Finding root of a quadratic equation, square of a root.

Finite Non-deterministic Algo:- Algo terminates within a finite amount of time but output may not be unique.
Ex:- To generate a random number.

Infinite Deterministic Algo:- Algo which do not terminate either because a terminating condition was not satisfied for a given set of input.

Ex:- Task of monitoring the temperature in a nuclear reactor is an infinite algorithm.



Ex: criterion 5 of algorithm is that each opⁿ is effective; each step must be such that it can, at least in principle, be done by a person using pencil and paper in a finite amount of time. Ex: performing arithmetic on integers is an example of an effective opⁿ, but arithmetic with real numbers is not since some values may be expressible only by infinitely long decimal expansion. Adding two such numbers would violate the effectiveness property.

- algorithms that are definite and effective are also called computational procedures.
- To help us achieve the criteria of definiteness, algorithms are written in a programming language (PL).
- such languages are designed so that each legitimate sentence has a unique meaning.
- A program is the expression of an algorithm in a PL.

Group 1a The study of algorithms includes many important and active areas of research. There are four distinct areas of study we can identify:

- (1) How to devise algorithms: creating an algorithm is an art which may never be fully automated. Various design techniques are there. By mastering these design strategies, it will become easier for you to devise new and useful algorithms.
- (2) How to validate algorithms: - once an algorithm is devised, it is necessary to show that it computes the correct answer for all possible legal inputs. We refer to this process as algorithm validation.

The purpose of the validation is to assure us that this algorithm will work correctly independently of the issues concerning the PL it will eventually be written in. Once the validity of the method has been shown, a program can be written and a second phase begins. This phase is referred to as program verification.

(3) How to analyze algorithms:- As an algorithm is executed, it uses the computer's CPU to perform operations and its memory to hold the program + data.

Analysis of algorithm or performance analysis refers to the task of determining how much computing time and storage an algorithm requires.

→ This is a challenging area which sometimes requires great mathematical skill.

→ An important result of this study is that it allows you to make quantitative judgements about the value of one algorithm over another.

(4) How to test a program:- Testing a program consists of two phases:-

- Debugging
- Profiling (or performance measurement)

Debugging is the process of executing programs on sample data gets to determine whether faulty results occur and, if so, to correct them.

However as E. Dijkstra has pointed out "debugging can only point to the presence of errors but not to their absence".

Profiling or performance measurement is the process of executing a correct program on datasets and measuring the time and space it takes to compute the results.

Extra on pseudocode it is not typically concerned with issues of software engineering.

→ issues of data abstraction, modularity, error handling are ignored in order to convey the essence of algo more precisely.

Testing is the process of executing the program with an intention of finding errors, faults and failures.

Debugging is the process of fixing the bugs by the testers and then submitted to the developer.

Validation evaluates the product itself. Validation ensures the functionality intended behavior of the product (reconfiguration takes place before testing).

Verification evaluates the documents, plans, codes, requirements & specifications.

Basic properties of Algorithm has 5 parts

- ① Input :- Zero or more ip data are externally supplied. That is no. of quantities are provided initially before the algo begins.
- ② Output :- There must be at least one op produced.
- ③ Definiteness :- Each instruction must be clear and unambiguous. i.e., the processing rule specified in the algorithm must be precise and unambiguous and lead to a specific action.

Ex:- ADD 50 to p \leftarrow clearly defined 50 will be added to variable p .

ADD 50 to p or q \leftarrow does not clearly defined. So there is ambiguity because we are not sure

Ex:- "compute $5/0$ " is not clear. \leftarrow is added to p or q .

- ④ Finiteness :- Algo should terminate in finite amount of time. i.e. The total time to carry out all the steps in the algorithm must be finite.

- ⑤ Effectiveness :- Each instruction must be sufficiently basic so that a person using paper and pencil can carry it out in a finite time.

Algorithm structure :-

- \Rightarrow Iterative \rightarrow executes action f in loop
- \Rightarrow Recursive \rightarrow re apply action to subproblem(s)

problem type :-

- \Rightarrow satisfying \rightarrow Find any satisfactory solution.
- \Rightarrow optimization \rightarrow Find best solution (V , cost metric)

Analyzing Algorithms:- Analyzing an algorithm has come to mean predicting the resources that the algorithm requires.

- occasionally resources such as
 - memory
 - com/n bandwidth
 - & computer h/w are of primary concern.
- but most often it is computational time that we want to measure.
- So analysis of algorithm refers to task of determining how much computing time and storage an algorithm requires.
- In general, the time taken by an algorithm grows with the size of the input, so it is traditional to describe the running time of a program as a function of the size of its input.
- The running time of an algorithm on a particular input is the number of primitive operations or steps executed.

Running time is affected by the

- h/w environment (processor, clock rate, memory, disk etc)
- SW environment (OS, PL, compiler, interpreter)

Kind of Analysis:-

- worst case (usually) - upper bound on the running time for any i/p.
 - $T(n)$ = maximum time of algo on any input of size n .
 - = Algo runs longest among all possible inputs of size n .
- Average case (sometimes)
 - $T(n)$ = expected time of algo over all i/p of size n .
- Best case
 - Algo runs the fastest among all possible inputs of size n .

NOTE Knowing the worst-case it guarantees that algorithm will not take any longer.

performance analysis of Algo: - To judge the performance of an algorithm there are basically two criteria:

- (i) Time complexity: (T_n) Amount of time it needs to run to completion as a function of its input size.
- (ii) Space complexity: (S_n) Amount of memory it needs to run to completion as a function of its input size.

Ex: 1 [Analysis of an algo provides background info that gives us a general idea of how long an algo will take for a given problem set.]

- The purpose of analysis is not to give a formula that will tell us exactly how many seconds or computer cycles needed.
- Analysis of an algo is to predict the resources that the algo requires and such analysis is based on individual computational model.

How to calculate running time of an algorithm?

- By running the implementation of the algo on a computer. Alternatively we can calculate by using a technique called algo analysis.
- we can estimate by counting the no. of basic opⁿ required by the algo to process an i/p of a certain size.

Basic opⁿ: The time to complete basic opⁿ does not depend on the particular values of its operands. So it takes a constant amount of time.
Ex: - add, sub, mult, div, bitlen (AND, OR, NOT), comparison, modulo & floor operation.

Ex: 2 Execution time depends on many parameters which are independent of the particular algo: m/c clock rate, quality of code produced by compiler, whether or not the computer is multi-processed etc.
→ Exe time of an algo is a function of the values of its i/p parameters.

GROWTH OF FUNCTIONS: The order of growth of the running time of an algorithm gives a simple characterization of the algorithm's efficiency and also allows us to compare the relative performance of alternative algorithms.

- When our input size n is going to be large enough in that case we have to take the help of asymptotic notations.
- For large enough input the multiplicative constants and lower order terms of an exact running time are dominated by the effects of input size.
- Notation which describe the algorithm efficiency and performance in a meaningful way is asymptotic notation.
- The growth rate of an algo is the rate at which running time of the algo grows as the size of the input grows.

$$\log n < n < n^2 < n^3 < 2^n$$

(n is very large)

Asymptotic Notations: - "Asymptotic means" - what will happen if n (i.e. input size) is very large.

- The asymptotic running time of an algorithm is defined in terms of functions.
- Asymptotic analysis studies how the values of functions compare as their arguments converge to infinity.

~~Three~~ ^{five} notations are used to compare orders of growth of algorithms.

- (i) O → big oh → upper bound → worst case
- (ii) Ω → big omega → lower bound → best case
- (iii) Θ → Theta → average case
- (iv) o → small oh / little oh
- (v) ω → small omega / little omega.

1) Big Oh (O) notation: upper bound for the function f if provided by the big-oh notation.

The function $f(n) = O(g(n))$ (read as " $f(n)$ is big-oh of $g(n)$ ") if and only if there exist a positive constant $c > 0$ and n_0 such that $f(n) \leq c * g(n)$ for all $n, n > n_0$.

(or) we can write $0 \leq f(n) \leq c * g(n)$ for all $n > n_0$

→ So $f(n)$ always lies below $c * g(n)$ $\forall n, n_0$
 → class of function $f(n)$ that grows no faster than $g(n)$

OR $O(g(n)) = \{ f(n) : \exists \text{ positive constants } c \text{ and } n_0 \text{ such that } \forall n > n_0, \text{ we have } 0 \leq f(n) \leq c * g(n) \}$

Intuitively: The set of all functions whose rate of growth is the same as or lower than that of $g(n)$.

So, $g(n)$ is an asymptotic upper bound for $f(n)$.

Ex:-

Ex:- $3n+3 = O(n)$ $\forall c > 0, \forall n > 3$
 as $3n+3 \leq 4n$

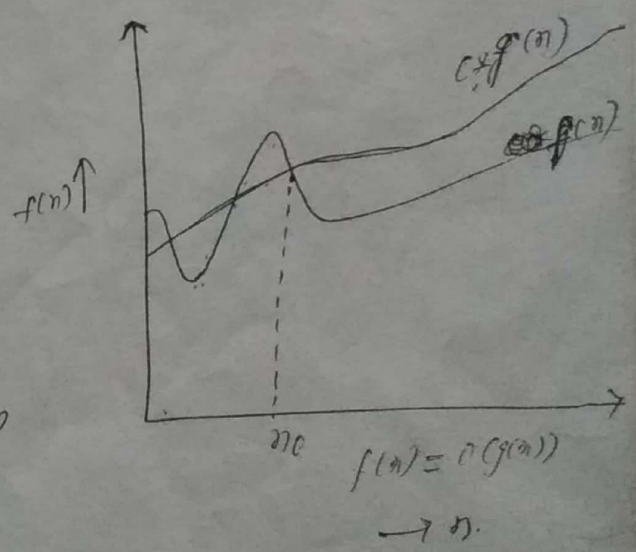
Ex:- $3n+2 = O(n)$
 as $3n+2 \leq 4n, \forall n > 2$

Ex:- $f(n) = 7n+5$ so $f(n) = O(n)$
 as $7n+5 \leq 8n, \forall n > 5$

Ex:- $10n^r + 4n + 2 = O(n^r)$
 as $10n^r + 4n + 2 \leq 11n^r, \forall n > 5$

Ex:- $3n+2 = O(n^r)$ as $3n+2 \leq 2n^r, \forall n > 2$
 If we take any upper bound $O(n^3), O(n^4), \dots$ if satisfied.

But not $O(1)$ so $3n+2 = O(n)$
 $= O(n^r)$
 $= O(n^3)$
 $\neq O(1)$



Ex: $3n+2 \neq O(1)$ as $3n+2$ is not less than or equal to c for any constant c and all n, n_0 .

Ex: $10n^2 + 4n + 2 \neq O(n) \neq O(1)$

Ex: Let $f(n) = 3n+1$

- $3n+1 \leq 4n$
 - $3n+1 \leq 5n$
 - $3n+1 \leq n^2$
 - $3n+1 \leq n^3$
 - $3n+1 \leq 2^n$
- } Growth function

$$f(n) \leq c \cdot g(n), \quad n > n_0$$

Incorrect bound

- $7n+5 \neq O(1)$
- $2n+3 \neq O(1)$
- $10n^2+7 \neq O(n)$

Loose bounds

- $2n+3 = O(n^2)$
- $4n^2+5n+6 = O(n^4)$
- $3n^3+4n^2+n = O(n^6)$

Ex: Let $f(n) = n+2$, then which condition is not satisfied?

- (i) $O(n)$
- (ii) $O(n^2)$
- (iii) $O(n^3)$
- (iv) $O(1)$

If $f(n)$ satisfies any lower function $g(n)$ then it will satisfy the higher function $g(n)$.

NOTE Since O notation describes an upper bound, we use it to bound the worst-case running time of an algorithm.

Ex: $f(n) = 5n-2 \neq O(n)$ for $c=5, n_0=1$

TASK $100n^3+7n+5 \neq O(n^3)$ for $c=7, 22, n_0=1$
 $3 \log n + \log \log n \neq O(\log n)$ for $c=7, 4, n_0=2$

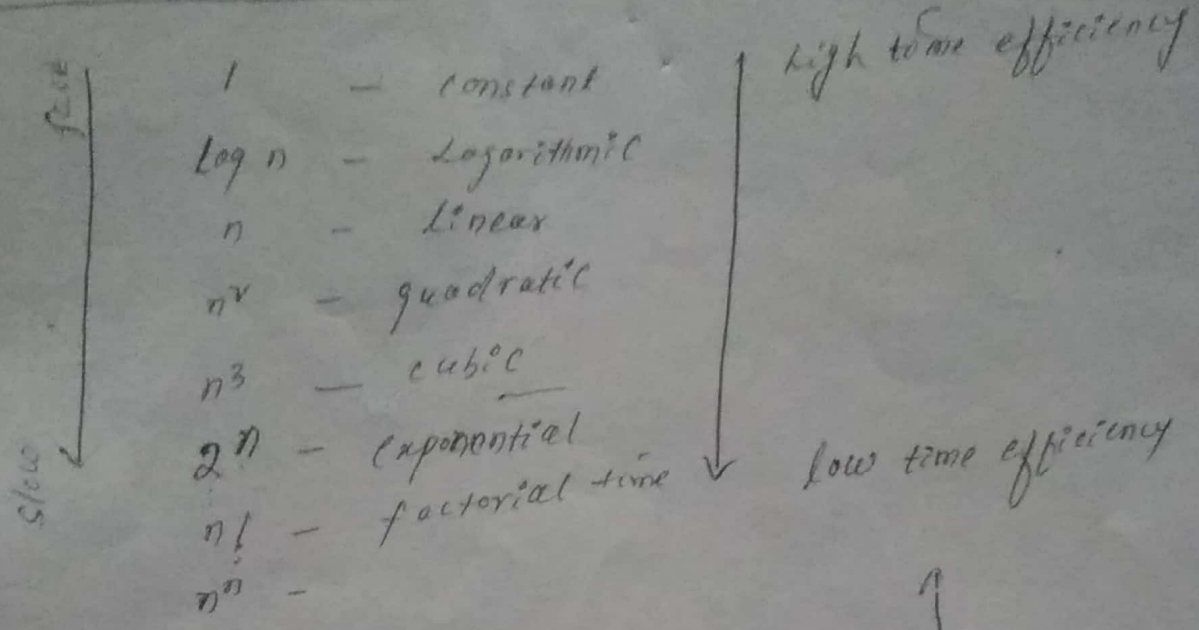
Ex: $2^{100} \neq O(1)$ since $2^{100} \leq 2^{100+1}$

Ex: $7/n \neq O(1/n)$

Ex: $4n^3 + O(n^2) = O(n^3)$

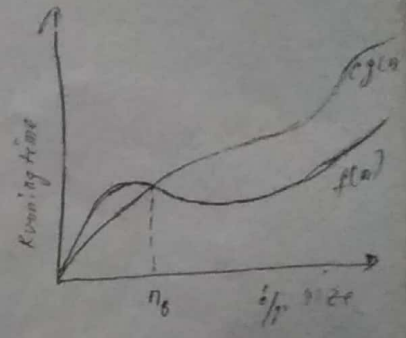
NOTE We write $O(1)$ to mean computing time of constant.

Some common functions



NOTE:- Algo taking $O(2^n)$ time should never be considered efficient.

$$\log n < \log^2 n < \sqrt{n} < n < n \log n < n^2 < n^3 < 2^n$$



Big-O notation

- $O((\log n)^k)$ - poly logarithmic time
- $n \cdot \log n$ - linearithmic function
- $n \log^k n$ - quasi linear time
- $O(n^k)$ = polynomial time.

$$O(1) < \log n < n < n \log n < n^2 < n^3 < 2^n$$

(2) Big-omega notation This gives asymptotic lower bound
 The function $f(n) = \Omega(g(n))$ (pronounced big-omega of g of n) or sometimes just omega of g of n) if and only if there exist the positive c and n_0 such that $f(n) \geq c \cdot g(n)$ for all n where $n > n_0$

$$\Omega(g(n)) = \left\{ f(n) : \text{there exist positive constants } c \text{ and } n_0 \text{ such that } 0 \leq c \cdot g(n) \leq f(n) \text{ for all } n > n_0 \right\}$$

Ex - $3n+3 = \Omega(n)$
 as $3n+3 \geq 3n \quad \forall n > 1$

(2) $3n+2 = \Omega(n)$
 as $3n+2 \geq 3n \quad \forall n > 1$

(3) $3n+3 = \Omega(n) \neq \Omega(n^2)$
 But it can be

$$3n+3 = \Omega(1)$$

(4) $6 \cdot 2^n + n^r \neq \Omega(3^n)$
 $6 \cdot 2^n + n^r = \Omega(n^r)$
 $= \Omega(n)$ or $\Omega(2^n)$ as $6 \cdot 2^n + n^r \geq 2^n$ for $n > 1$

(5) $10n^r + 4n + 2 = \Omega(n^r)$ as $10n^r + 4n + 2 \geq n^r$
 $= \Omega(n) = \Omega(1) \neq \Omega(n^3)$

(6) $6 \cdot 2^n + n^r = \Omega(2^n)$ as $6 \cdot 2^n + n^r \geq 2^n$ for $n > 1$. or $\Omega(1)$

(7) $3 \log n + \log \log n \neq \Omega(\log n)$
 $3 \log n + \log \log n \geq 3 \log n$ for $n > 2$

(8) $f(n) = 3n+1$
 $3n+1 \geq 3n$ here $f(n) = 3n+1$ and $c = 3$, $g(n) = n$

So $n_0 = 1$ so $3n+1 = \Omega(n)$

(9) prove that $3n+1 = \Omega(1)$

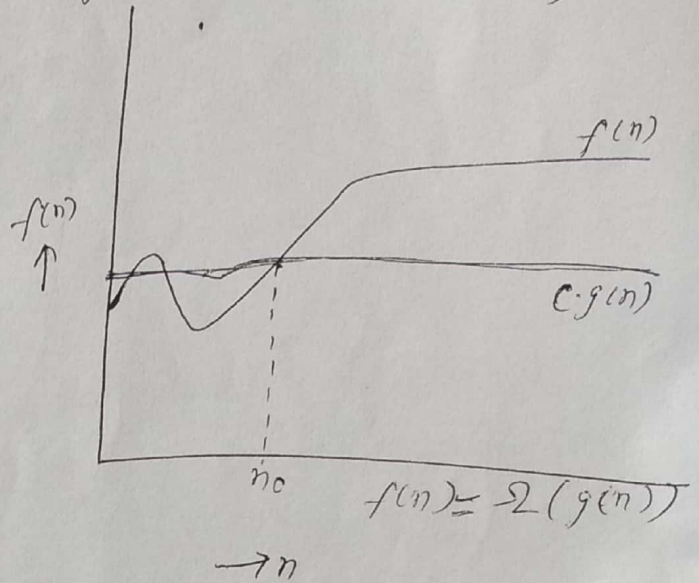
$$f(n) \geq c \cdot g(n)$$

$$f(n) = 3n+1$$

$$3n+1 \geq 1 \cdot 1$$

$$n > 1$$

So here $c = 1$, $g(n) = 1$ so $3n+1 = \Omega(1)$ proved



③ Theta (Θ) notation: - the function $f(n) = \Theta(g(n))$ if there exist constants $c_1 > 0$ and $c_2 > 0$ and an integer constant $n_0 > 1$ such that $\boxed{c_1 g(n) \leq f(n) \leq c_2 g(n)}$ for every integer $n > n_0$

$$\Theta(g(n)) = \left\{ f(n) : \exists \text{ positive constants } c_1, c_2 \text{ and } n_0 \text{ such that } \forall n > n_0, 0 \leq c_1 g(n) \leq f(n) \leq c_2 g(n) \right\}$$

→ If $f(n) \not\leq c_1 g(n)$ and $f(n) \not\geq c_2 g(n)$ then $\underline{f(n) = \Theta(g(n))}$

→ $\Theta(g(n)) = O(g(n)) \cap \Omega(g(n))$

→ Θ -notation is tight bound.

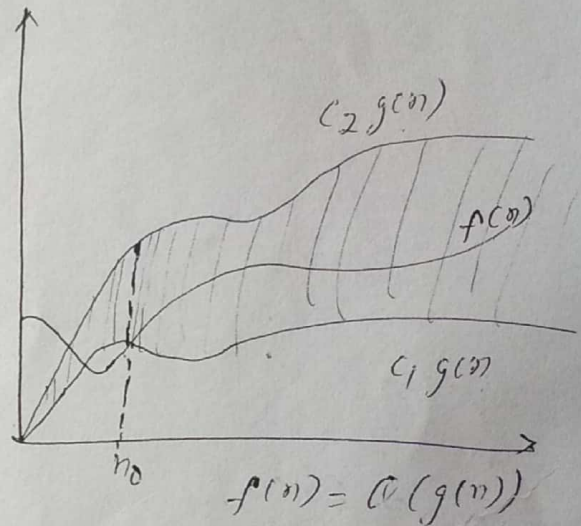
→ $T_{avg}(n) \leq T_{worst}(n)$

→ $f(n) = O(g(n)) \Rightarrow f(n) = \Omega(g(n))$

$\Theta(g(n)) \subseteq \Omega(g(n))$

→ "Running time of $O(f(n))$ " \Rightarrow worst case of $O(f(n))$

→ "Running time of $\Omega(f(n))$ " \Rightarrow Best case of $\Omega(f(n))$



Ex: $f(n) = 3n^r + 5$

$3n^r + 5 \leq 4n^r$

Here $f(n) = 3n^r + 5$, $c_1 = 4$, $g(n) = n^r$

So $f(n) = 3n^r + 5 = O(n^r)$ ————— (1)

Again $f(n) = 3n^r + 5$

$3n^r + 5 \geq 3n^r$

$f(n) = 3n^r + 5$, $c_2 = 3$, $g(n) = n^r$

So $f(n) = 3n^r + 5 = \Omega(n^r)$ ————— (2)

from (1) and (2), we can conclude that

$\boxed{f(n) = \Theta(n^r)}$

Theorem-1

For any two functions $f(n)$ and $g(n)$, we have $f(n) = \Theta(g(n))$ if and only if $f(n) = O(g(n))$ and $f(n) = \Omega(g(n))$

Ex: The function $5n+2$ is $O(n)$
 As $5n+2 \geq 5n$ for all $n \geq 1$ and $5n+2 \leq 6n$ for all $n \geq 2$
 So $c_1=5$, $c_2=6$ and $n_0=2$ for which we can verify that
 $5n+2 = O(n)$.

Ex: $10n^2+4n+2 = O(n^2)$ as $5n^2 \leq 10n^2+4n+2 \leq 11n^2$, $\forall n \geq 5$

Ex: $3n+2 = O(n)$ as $3n \leq 3n+2 \leq 4n$, $\forall n \geq 2$. Here $c_1=3$, $c_2=4$
 and $n_0=2$.

Ex: $f(n) = 10n^2+7$ as $10n^2 < 10n^2+7$ for all n , $c_1=10$
 also $10n^2+7 \leq 11n^2$, $c_2=11$, $n_0=7$

Thus $10n^2 < 10n^2+7 \leq 11n^2$, $c_1=10$, $c_2=11$, $n \geq n_0=7$.

Incorrect bound

$$7n+5 \neq O(1)$$

$$7n+5 \neq O(n^2)$$

$$2n+3 \neq O(1)$$

$$10n^2+7 \neq O(n^3)$$

$$10n^2+7 \neq O(1)$$

$$10n^2+4n+2 \neq O(n)$$

Ex: $6 \times 2^n + n^r = O(2^n)$ as $6 \times 2^n \leq 6 \times 2^n + n^r \leq 7 \times 2^n$, $n \geq 2$.

$$6 \times 2^n + n^r \neq O(n^r)$$

$$6 \times 2^n + n^r \neq O(n^{100})$$

$$6 \times 2^n + n^r \neq O(1)$$

Theorem-1 If $f(n) = a_m n^m + a_{m-1} n^{m-1} + \dots + a_1 n + a_0$, then
 $f(n) = \Omega(n^m)$ (or) $f(n) = \Theta(n^m)$ (or) $f(n) = O(n^m)$

(4) Small-oh (o) notations: The function $f(n) = o(g(n))$ if there exist a real constant $c > 0$ and an integer constant $n_0 > 0$ such that $f(n) < c g(n)$ for every integer $n > n_0$.

(*) $O(g(n)) = \{ f(n) : \text{for any positive constants } c > 0, \text{ there exist a constant } n_0 > 0 \text{ such that } \boxed{0 \leq f(n) < c g(n)} \text{ for all } n > n_0 \}$

Ex: $2n = o(n^2)$ but $2n^2 \neq o(n^2)$ ← small-oh

→ The asymptotic upper bound provided by O -notation may or may not be asymptotically tight.

→ The bound $2n^2 = O(n^2)$ is asymptotically tight but the bound $2n = O(n^2)$ is not.

→ we use o -notation to denote an upper bound that is not asymptotically tight.

→ The defⁿ of O -notation and o -notation are similar.

→ The main difference is that in $f(n) = O(g(n))$, the bound $0 \leq f(n) \leq c g(n)$ holds for some constant $c > 0$.

but in $f(n) = o(g(n))$, the bound $0 \leq f(n) < g(n)$ holds for all constants $c > 0$.

→ Intuitively, in the o -notation, the function $f(n)$ becomes insignificant relative to $g(n)$ as n approaches infinity; that is

$$\boxed{\lim_{n \rightarrow \infty} \frac{f(n)}{g(n)} = 0}$$

Ex: $f(n) = 3n+2 = o(n^2)$ as $\lim_{n \rightarrow \infty} \frac{3n+2}{n^2} = 0$ (8)

$\lim_{n \rightarrow \infty} \frac{3n+2}{n} = \lim_{n \rightarrow \infty} \frac{3}{1} = 3 \neq 0$

Ex: But $f(n) = 3n+2 \neq o(n)$ as $\lim_{n \rightarrow \infty} \frac{3n+2}{n} = \lim_{n \rightarrow \infty} 3 + \frac{2}{n} \neq 0$

So $\boxed{3n+2 \neq o(n) \text{ but } 3n+2 = o(n^2)}$

Ex: $3n+2 = o(n \log n)$ as $\lim_{n \rightarrow \infty} \frac{3n+2}{n \log n} = \lim_{n \rightarrow \infty} \left(\frac{3}{\log n} + \frac{2}{n \log n} \right) = 0$

NOTE: $g(n)$ is an upper bound for $f(n)$ that is not asymptotically tight.
 (2) For any higher term growth of big-oh notation, the little-oh is satisfied.

(5) Little omega notation (ω): We use ω -notation to denote a lower bound that is not asymptotically tight.

$\omega(g(n)) = \left\{ f(n) : \text{for any positive constant } c > 0, \text{ there exists a constant } n_0 > 0 \text{ such that } 0 < c g(n) < f(n) \text{ for all } n > n_0 \right\}$

\rightarrow if $f(n) \in \omega(g(n))$ if and only if $g(n) \in o(f(n))$

Ex: $n^{1/2} = \omega(n)$ but $n^{1/2} \neq \omega(n^2)$

The relation $f(n) = \omega(g(n))$ implies that

$\lim_{n \rightarrow \infty} \frac{f(n)}{g(n)} = \infty$ or $\lim_{n \rightarrow \infty} \frac{g(n)}{f(n)} = 0$

That if $f(n)$ becomes arbitrarily large relative to $g(n)$ as n approaches infinity

Alternate defⁿ The function $f(n) = \omega(g(n))$ iff $\lim_{n \rightarrow \infty} \frac{g(n)}{f(n)} = 0$

Ex: $f(n) = 3n+2 \neq \omega(n)$ as $\lim_{n \rightarrow \infty} \frac{n}{3n+2} = \frac{1}{3}$ (L'Hospital rule)
 $\frac{f(n)}{g(n)} = \omega$ $\frac{3n+2}{n} = 3 + \frac{2}{n} = 3 + 0 = 3$

Ex: But $f(n) = 3n+2 = \omega(\log n)$ as $\lim_{n \rightarrow \infty} \frac{\log n}{3n+2} = \lim_{n \rightarrow \infty} \frac{1}{3n} = 0$

Ex: But $f(n) = 3n+2 \neq \omega(n^r)$ as $\lim_{n \rightarrow \infty} \frac{n^r}{3n+2} = \lim_{n \rightarrow \infty} \frac{n}{\frac{3}{2}}$

So $3n+2 \neq \omega(n^r)$ because n^r is larger than $3n+2$

Ex: - $3n+2 = \omega(1)$
 as $\lim_{n \rightarrow \infty} \frac{1}{3n+2} = 0$ So $3n+2 = \omega(1)$

NOTE: For any lower growth of the omega notation (ω) little omega notation condition is satisfied.

- (2) for small-oh $f(n)$ is less than $g(n)$
- (3) for small-omega $f(n)$ is larger than $g(n)$.

Correct Bounds (little-oh)

$\rightarrow f(n) = 3n+5 = o(n^r)$ as $3n+5 = O(n^r)$
 $\rightarrow f(n) = 3n^r+4n = o(n^4)$ as $3n^r+4n = O(n^4)$
 $\rightarrow f(n) = 4n^3+2n+3 = o(n^4)$ as $4n^3+2n+3 = O(n^4)$

Incorrect Bounds (little-oh)

$\rightarrow f(n) = 2n+3 \neq o(n)$
 $\rightarrow f(n) = 27n^r+16n \neq o(n^r)$ as $\lim_{n \rightarrow \infty} \frac{27n^r+16n}{n^r} = 27 \neq 0$
 $\rightarrow f(n) = 10n^r+7 \neq o(n^r)$ as $\lim_{n \rightarrow \infty} \frac{10n^r+7}{n^r} = 10 \neq 0$
 $\rightarrow f(n) = 3n^3+4n \neq o(n^3)$

Ex of little-oh or small-oh

$3n+2 = o(n \log \log n)$ but $3n+2 \neq o(n)$
 $6 \times 2^n + n^r = o(2^n)$
 $= o(2^{n \log n})$
 $\neq o(2^n)$

Recurrences - The fact of occurring again

A recurrence is an equation or inequality that describes a function in terms of its value on smaller inputs.

- Recurrence relations are recursive defⁿ of mathematical functions or sequences
- Recurrence relations arises when we analyze the running time of iterative or recursive algorithm.
- The running time of a recursive algorithm can be obtained by a recurrence.

Ex:- $T(n) = 2 * T(n/2) + n$ example of recurrence relation

Ex:- The complexity of many Divide and Conquer algorithms is given by recurrences of the form:-

$$T(n) = \begin{cases} T(1) & , n=1 \\ a T(n/b) + f(n) & , n > 1 \end{cases}$$

where 'a' is the no. of subproblems and 'b' is the size of each subproblems. i.e a and b are known constants.

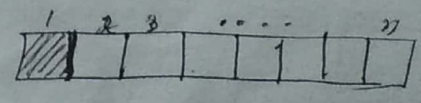
We assume that T(1) is known and n is power of b (i.e $n = b^k$)

Methods of Solving recurrence relation:-

- (1) Substitution method → Guess the form of the solⁿ. That will work when solⁿ is easy to guess.
- (2) Iteration "
- (3) Recursion tree. "
- (4) Master method

Consider the linear search.

Recursively look at one element then search remaining elements



so $T(n) = T(n-1) + c$

" The cost of searching n elements is the cost of searching looking at 1 element, plus the cost of searching n-1 elements."

(1) Substitution method :- (guess method)

In this method we guess a bound and then use mathematical induction to prove our guess correct.

So the substitution method for solving recurrences entails two steps:

- ① \rightarrow Guess the form of the solution
- ② \rightarrow Use the mathematical induction to find the constants and show that the solution works.

This method is powerful but it is obviously can be applied only in cases when it is easy to guess the form of the answer.

- This method can be used to establish either upper or lower bounds on the recurrence.
- Unfortunately, there is no general way to guess the correct solutions to recurrences.
- Guessing a solⁿ takes experience and occasionally, creativity.

The most general method:

- guess the form of solⁿ
- verify by induction
- solve for constants

Mathematical induction :-
is a method of mathematical proof typically used to establish a given result for all natural numbers.

✓ Ex-1 $T(n) = T(n/2) + 1$ Determine upper bound

Let's guess that the solution is $T(n) = O(\log n)$

So our method is to prove that $T(n) \leq c \log n$ for any constant $c > 0$.

$$\text{So } T(n) \leq c \cdot \log n \quad \text{--- (2)}$$

$$\begin{aligned} T(n/2) &\leq c \cdot \log(n/2) \\ &\leq c \cdot \log n - c \cdot \log 2 \end{aligned}$$

Now, eqn (2) becomes,

$$\begin{aligned} T(n) &\leq c \cdot \log n - c \cdot \log 2 + 1 \\ &\leq c \cdot \log n - c + 1 \\ &\leq c \cdot \log n \\ \therefore T(n) &= O(\log n) // \text{proven} \end{aligned}$$

Recursion Tree method :- A recursion tree is generated by tracing the execution of a recursive algorithm. It is a pictorial representation of recursion method.

A recursion tree models the costs (time) of a recursive execution of an algorithm.

→ In a recursion tree, each node represents the cost of a single subproblem somewhere in the set of recursive function invocations.

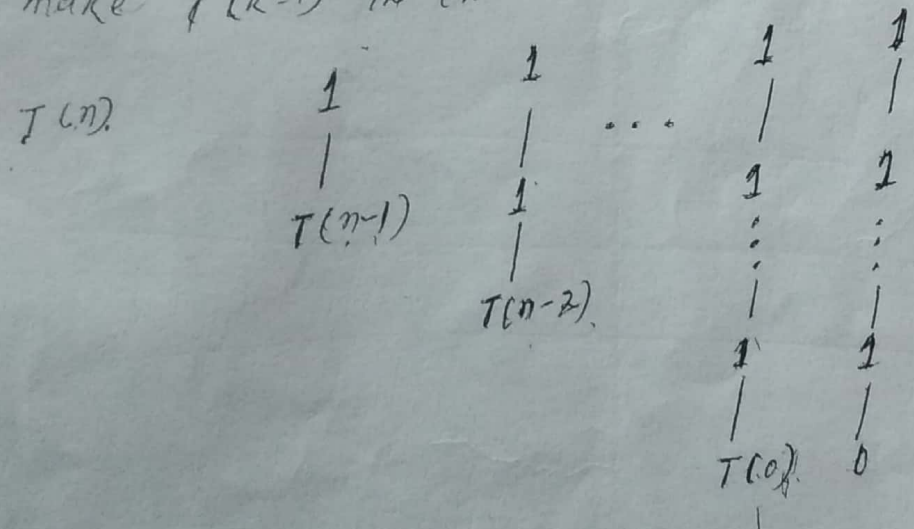
→ We sum the costs within each level of the tree to obtain a set of per-level costs and we sum all the per-level costs to determine the total cost of all levels of the recursion.

→ Recursion trees are useful when the recurrence describes the running time of a D-A-C algorithm's Digital to analog Conversion Divide and Conquer algorithm which based on multi. This method is good for generating guesses for the substitution method.

branch recursion

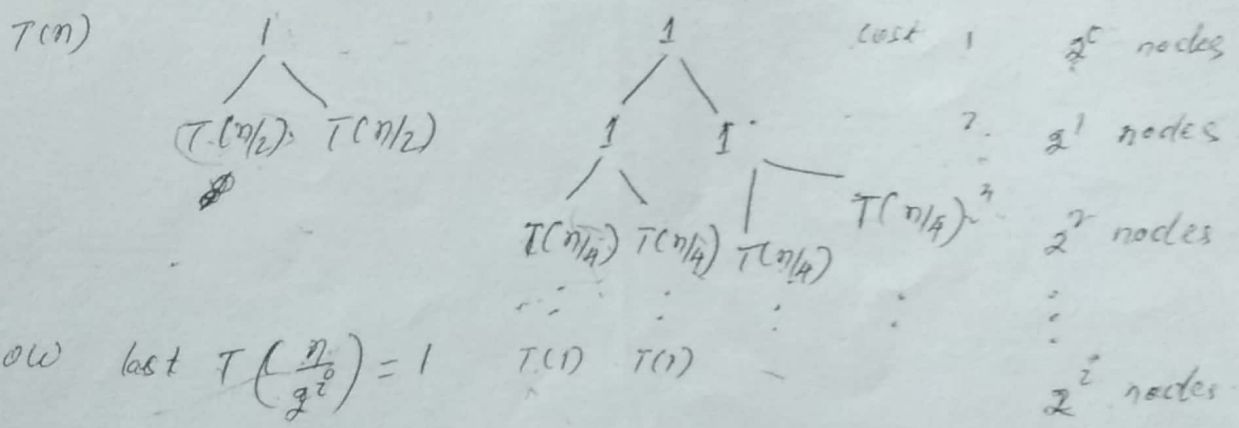
Ex - $T(n) = T(n-1) + 1$ every parent replaced by constant term

At each level of tree, replace the parent of that subtree by the constant term of the recurrence relation $T(n) = T(n-1) + 1$ and make $T(k-1)$ the child.



$T(n)$ is computed by finding the sum of elements at each level of the tree. Therefore $T(n) = O(n)$

Ex-2 $T(n) = 2T(\frac{n}{2}) + 1$ Here the no. of child = 2.



Now last $T(\frac{n}{2^i}) = 1$

so $\frac{n}{2^i} = 1 \Rightarrow n = 2^i \Rightarrow 2^i = n$ so $i = \log_2 n$

so height = $\log_2 n$ and level = $\log_2 n + 1$

But last step contains 2^i nodes

so total = $1 + 2 + 4 + \dots + (2^i)$

$$= 1 \left(\frac{2^{i+1} - 1}{2 - 1} \right) = 2^{i+1} - 1 = 2 \cdot 2^i - 1 = 2n - 1 \leq 3n$$

so $T(n) = O(n)$

Master method :- The master method applies to recurrences of the form $T(n) = aT(n/b) + f(n)$ where $a \geq 1$ and $b > 1$ are constants $f(n)$ is asymptotically +ve function

The above recurrence describes the running time of an algorithm that divides a problem of size n into a subproblem each of size n/b , where a and b are positive constants.

The master method depends on the following theorem

Master theorem :- Let $a \geq 1$ and $b > 1$ be constants, let $f(n)$ be a function and let $T(n)$ be defined by the recurrence.

$$T(n) = aT(n/b) + f(n)$$

Then $T(n)$ can be bounded asymptotically as follows :-

1. if $f(n) = O(n^{\log_b a - \epsilon})$ for some constant $\epsilon > 0$ then $T(n) = \Theta(n^{\log_b a})$
2. if $f(n) = \Theta(n^{\log_b a})$ then $T(n) = \Theta(n^{\log_b a} \lg n)$
3. if $f(n) = \Omega(n^{\log_b a + \epsilon})$ then $T(n) = \Theta(f(n))$ and if $a f(n/b) \leq c f(n)$ for some constants $c < 1$ and all sufficiently large n , then

$$T(n) = \Theta(f(n))$$

NOTE :- It means compare $f(n)$ with $n^{\log_b a}$

Three common cases :-

1. $f(n)$ grows polynomially slower than $n^{\log_b a}$ then $T(n) = \Theta(n^{\log_b a})$
2. $f(n)$ and $n^{\log_b a}$ grows at similar rate then $T(n) = \Theta(n^{\log_b a} \lg n)$
3. $f(n)$ grows polynomially faster than $n^{\log_b a}$ and $f(n)$ satisfies the regularity condition then $T(n) = \Theta(f(n))$

Ex - $T(n) = 4T(n/2) + n$

Here $a=4, b=2$ so $n^{\log_b a} = n^2$ but $f(n) = n$

Case-1 $f(n) = O(n^{r-\epsilon})$ for $\epsilon=1$ // case-1 is satisfied

$\therefore T(n) = O(n^2)$

$\log_2 4 = \log_2 2^2 = 2$
 $= 2 / \log_2 2 = 2$
 $= 2$
 $f(n) < n^{\log_2 4}$
 so case-1 is satisfied

Ex - $T(n) = 4T(n/2) + n^2$

Here $a=4, b=2$ so $n^{\log_b a} = n^2$ and $f(n) = n^2$

so case-2 is satisfied

$f(n) = O(n^{\log_b a} \lg n) = O(n^2 \lg n)$

Ex - $T(n) = 4T(n/2) + n^3$

Here $a=4, b=2$ so $n^{\log_b a} = n^2$ but $f(n) = n^3$

so case-3 is satisfied

$f(n) = \Omega(n^{r+\epsilon})$ for $\epsilon=1$

and $4(n/2)^3 \leq c n^3$ for $c=1/2$

$\therefore T(n) = O(n^3)$

$f(n) = n^3 > n^2$
 $a f(n/b) \leq c f(n)$
 $4 f(n/2) \leq c n^3$
 $4 (n/2)^3 \leq c n^3$
 $4 n^3 / 8 \leq c n^3$
 $n^3 \leq c n^3$
 $c \geq 1/2$

Ex - $T(n) = 9T(n/3) + n$

Here $a=9, b=3, f(n)=n$

$n^{\log_b a} = n^{\log_3 9} = n^2$. Since $f(n) = n$, case-1 is satisfied.

$f(n) = O(n^{\log_3 9 - \epsilon})$ where $\epsilon=1 = O(n^{2-1})$

$\therefore T(n) = O(n^2)$

Ex - $T(n) = 2T(n/3) + 1$

Here $a=2, b=3/2$

$f(n)=1$ but $\log_b a = \log_{3/2} 2 = \log_{3/2} 2 = 1$

so case-2 is satisfied

$f(n) = O(n^{\log_b a}) = O(1)$

so $T(n) = O(\lg n)$

$O(n^{\log_{3/2} 2 - \epsilon}) < 1$

Algorithm Classifications: Brute force

→ This algorithm simply tries all possibilities until a satisfactory solution is found.

• Such an algorithm can be either optimizing or ~~satisfactory~~/satisficing

Optimizing: - Find the best solution.

This may require finding all solutions, or if a value for the best solution is known, it may stop when any best sol.ⁿ is found.

- Example: Finding a best path for a Travelling salesman problem. i.e. cover all the cities & the total path cost is minimum.

Satisficing: - Stop as soon as a solution is found that is good enough

- Example: Finding travelling salesman path which is within 10% of optimal.

Improvements: -

- * Often Brute-force algorithms require exponential time
- * Various heuristics and optimizations can be used.

• Heuristic: A rule of thumb that helps you decide which possibilities to look at first.

Explain \Rightarrow we guess that this is the previous value using that previous value which we should explore to next. Optimization: In this case, a way to

eliminate certain possibilities without fully exploring them.

ϕ That means we are not going to explore all the possibilities.

If we go on doing such then the complexity of algorithm will be too high. So that the optimization we ~~can~~ may not explore all the possibilities with ~~the~~ fullest extent. We may leave them i.e. we may not execute all the possibilities we may not test all the possibilities, eliminate certain possibilities without fully exploring them.

Module - 11

Selection Sort

Assume that we have the data

These are the Location →

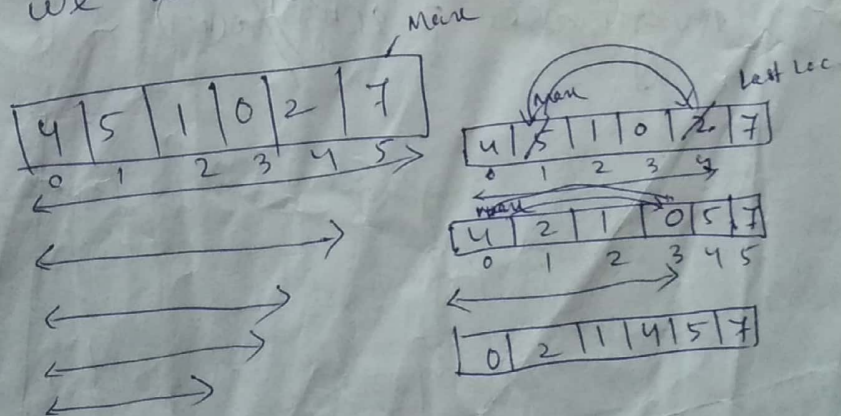
4	5	1	0	2	7
0	1	2	3	4	5

In selection sort we are going to do is in this complete data initially we have find which is the maximum element → then put the maximum element to the last location.

→ That means we have to put in 5th location ^{this array from}
→ Next time we are going to find 0 to 4th location, then find the maximum element and put that last location.

Example

Assume we are searching by linear search



Selection sort (int A[], int n)

```

int i, j;
int max;
for (i = n-1; i > 0; i--)

```

```

{
    max = A[0];
    for (j = 1; j < i; j++)

```

```

{
    if (A[j] > max A[max]);

```

after this
inner for
loop will have
to stop

```

{
    max = j;
}
}

```

```

int temp = A[max];

```

```

A[max] = A[i]

```

```

A[i] = temp

```

```

}
}

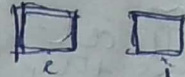
```

Best case = $O(n^2)$, $\omega(n^2)$ general book
 worst case = $O(n^2)$
 Avg case = $O(n^2) / O(n^2)$

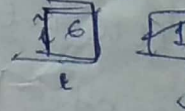
Assume this is the

array

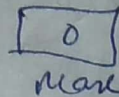
0	1	1	6	7	8	9
0	1	2	3	4	5	6



total array = 6



location is actually



How when A[i] is

GREEDY METHOD

(3)

✓ A greedy algorithm always makes the choice that looks best at the moment. That is, it makes a locally optimal choice in the hope that this choice will lead to globally optimal solution.

→ ~~but~~ greedy algorithms don't always yield optimal solⁿ but for many problems they do.

✓ → once a decision is made it is never revoked.

✓ An essential point of greedy solution is that we never have to revise our greedy decisions and they lead to fast algorithm provided.

✓ → The solution obtained may or may not be optimal solution, but it is the best solution obtained.

→ Instead of finding optimal solⁿ to subproblems and then making choice greedy algorithm, first make choice - the choice that looks best at the time and then solve a resulting subproblem.

→ The greedy method is also applied to optimization problems.

→ Greedy algorithm ascends to a local optimum which is a part of the global optimum.

The steps for achieving greedy algorithms are:

• feasible: Here we check whether it satisfies all the possible problem constraints or not, so as to obtain at least one solⁿ to our problem.

• Local optimal choice: in this choice should be the optimum which is selected from the current available feasible choices.

• unalterable: once the choice is made, at any subsequent step that choice is not altered.

Greedy (i) may or may not give optimal solⁿ (ii) Top-down design (iii) In greedy we find a feasible solⁿ but not in DP (iv) In DP we make choice at each step, but the choice

NOTE Greedy and DP are methods for solving optimization problems. Greedy makes choice then solve subproblems but in DP solve subproblems first then make the first choice.

Algorithm Greedy (1, n)

```
Solution =  $\emptyset$  (initialize the solution)
for  $i = 1$  to  $n$ 
     $x = \text{select}(a)$ ;
    if feasible (Solution,  $x$ ) then
        Solution = Union (Solution,  $x$ );
return Solution
```

- The function select an x from $a[i]$ and remove it.
- feasible is a boolean valued function that determines whether x can be included into solution vector.
- The function Union combines x with the solution and updates the objective function.

Example of Greedy :

- activity selection problem
- Huffman coding
- knapsack
- Spanning tree (e.g. MST :- Prim's algo, - Kruskal)
- Shortest path (e.g. single source shortest path - Dijkstra's algo)

feasible solⁿ Suppose we have a problem with n -inputs that require a solⁿ satisfying some constraints. Any subset that satisfies these constraints is called a feasible solution.

- We need to find a feasible solⁿ that either maximizes or minimizes the given objective function.
- A feasible solⁿ that does this is called an optimal solⁿ.

Activity selection problem The problem of scheduling several competing activities that require exclusive use of common resources, with the goal of selecting a maximum-size set of mutually compatible activities.

→ Suppose we have a set $S = \{a_1, a_2, \dots, a_n\}$ of n proposed activities that wish to use a resource, such as lecture halls, which can be used by only one activity at a time.

→ Each activity a_i has start time s_i and finish time f_i

where $0 \leq s_i < f_i < \infty$

→ Activities a_i and a_j are compatible if the interval $[s_i, f_i)$ and $[s_j, f_j)$ do not overlap (i.e. a_i and a_j are compatible if $s_i > f_j$ or $s_j > f_i$). So for non-interfering $a_i \cap a_j = \emptyset$.

→ The activity selection problem is to select maximum-size subset of mutually compatible activities.

Ex Consider following set S

i	1	2	3	4	5	6	7	8	9	10	11
s_i	1	3	0	5	3	5	6	8	8	2	12
f_i	4	5	6	7	8	9	10	11	12	13	14

NOTE
 $[2, 3)$
 less than or equal to;
 but less than 3.

→ For this example, the subset $\{a_3, a_9, a_{11}\}$ consists of mutually compatible activities.

→ but it is not maximal subset

→ The subset $\{a_1, a_4, a_8, a_{11}\}$ is largest subset of mutually compatible activities.

→ another is $\{a_2, a_4, a_9, a_{11}\}$.

Greedy choice property - a locally optimal choice is globally optimal

But - DP make decision based on all the decisions made in the previous stage, and may re-consider the previous stage's algorithm path to find.

But - Greedy never re-consider its choices.

Greedy-activity-selector algorithm It assumes that the input activities are ordered by monotonically increasing finish time.

GREEDY-ACTIVITY-SELECTOR (S, f)

1. $n = \text{length}[S];$
 2. $A = \{a_1\}$
 3. $i = 1$
 4. for $m \leftarrow 2$ to n
 5. do if $s_m > f_i$
 6. then $A \leftarrow A \cup \{a_m\}$
 7. $i \leftarrow m$
- return A .

→ It collects selected activities into a set A and returns this set when it is done.

→ Since the activities are considered in order of monotonically increasing finish time, f_i is always the maximum finish time of any activity in A . That is

$$f_i = \max \{ f_k : a_k \in A \}$$

Time complexity Above algorithm schedules a set of n -activities in $O(n)$ time, assuming that the activity were already initially by their finish time.

If the activity is not sorted then it takes

$$T(n) = O(n \log n) + O(n)$$

↑
any sorting also (merge, heap, quick sort)

Recursive Greedy algo This algorithm purely greedy, top-down fashion we give here recursive soln for activity selector. It returns maximum-size set of mutually compatible activities in S_{ij} . We assume that the n -input activities are ordered by monotonically increasing finish time.

REC-ACT-SEL (S, f, i, j)

1. $m \leftarrow i+1$ $\rightarrow 12$
2. while $m < j$ and $S_m < f_i$ // find the first activity in S_{ij}
3. do $m \leftarrow m+1$
4. if $m < j$
5. then return $\{a_m\} \cup \text{REC-ACT-SEL}(S, f, m, j)$
6. else return \emptyset .

→ The while loop of line 2-3 look for the first activity in S_{ij} which can be included.

→ The loop examine $a_{i+1}, a_{i+2}, \dots, a_{j-1}$, until it finds the first activity a_m that is compatible with a_i ; such an activity has $S_m > f_i$.

→ loop terminates when it finds an activity.

→ Assuming the activities are sorted by finish time this algorithm takes $O(n)$ time because each activity is examined exactly once in the while loop part of line-2.

→ Assume a_0 activity finishes at time 0. then the initial call $\text{REC-ACT-SEL}(S, f, 0, 12)$ is made then a_1 is selected.

→ In each recursive call, the activities that have already been selected are shaded, and activity shown in whole is being considered.

→ The arrow points directly up or to the right.

→ To the right arrow is selected.

Knapsack problem: We are given a set S of n items and a knapsack or bag. Each item has a +ve benefit p_i and a +ve weight w_i . The maximum capacity of the knapsack is m .

- If a fraction x_i , $0 \leq x_i \leq 1$, of object i is placed into the knapsack then a profit of $p_i x_i$ is earned.
- The objective is to obtain a filling of the knapsack that maximizes the total profit earned.
 - Since the knapsack capacity is m , we require the total weight of all chosen objects to be at most m .

Formally the problem can be defined as

$$\text{maximize } \sum_{1 \leq i \leq n} p_i x_i \quad \text{--- (1)}$$

$$\text{Subject to } \left[\sum_{1 \leq i \leq n} w_i x_i \leq m \right] \quad \text{--- (2)}$$

$$\text{and } 0 \leq x_i \leq 1, 1 \leq i \leq n \quad \text{--- (3)}$$

for 0-1 knapsack
 $x_i \in \{0, 1\}$

- The profits and weights are +ve numbers.
- The feasible solution is any set (x_1, x_2, \dots, x_n) satisfying above condition.
- An optimal solution is a feasible soln for which eqn (1) is maximized.
- We have to pack the knapsack (or bag) in such a manner as to get the maximum total value.

The knapsack problem has two variants:

0/1 knapsack

- In this items are individual; either we take full one item or discard it.
- This problem solved using DP not by GP.

Fractional knapsack we take any fraction of an item. Solvable by GP.

→ 0-1 knapsack restricts the number x_j to be zero or one.

- Greedy strategy solve the problem by putting items into bag one-by-one.
- This approach is greedy because once an item has put into the bag it is never removed.

Consider the following instances of the knapsack problem

$$n = 3, m = 20 \quad (P_1, P_2, P_3) = (25, 24, 15) \text{ and}$$

$$(w_1, w_2, w_3) = (18, 15, 10)$$

x_1	x_2	x_3	$\sum w_i x_i$	$\sum P_i x_i$
1	$\frac{1}{2}$	0	$18 \times 1 + 15 \times \frac{1}{2} + 10 \times 0 = 25.5$ (not feasible) because > 20	
$\frac{1}{2}$	$\frac{1}{3}$	$\frac{1}{4}$	$18 \times \frac{1}{2} + 15 \times \frac{1}{3} + 10 \times \frac{1}{4} = 16.5$ (feasible sol ⁿ)	$12.5 + 8 + 3.5 = 24.25$
1	$\frac{2}{15}$	0	$18 \times 1 + 15 \times \frac{2}{15} + 0 \times 10 = 20$	$25 + 24 \times \frac{2}{15} + 0 \times 15 = 28.2$
0	$\frac{2}{3}$	1	$18 \times 0 + \frac{2}{3} \times 15 + 1 \times 10 = 20$	$25 \times 0 + \frac{2}{3} \times 24 + 15 \times 1 = 16 + 15 = 31$
0	1	$\frac{1}{2}$	$0 \times 18 + 15 \times 1 + 1 \times 10 = 20$	31.5

Out of these four feasible solution, solution 4 (i.e. last) yields the maximum profit.

Ex The greedy strategy does not work for 0-1 knapsack problem.

Suppose there are 3 items and knapsack can hold 50 pounds.

$$w_1 = 10 \text{ pounds}, P_1 = 60 \text{ dollars}$$

$$w_2 = 20, P_2 = 100 "$$

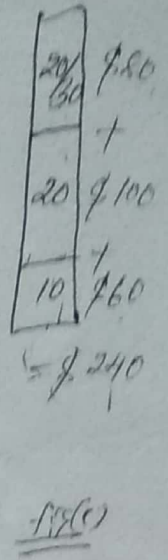
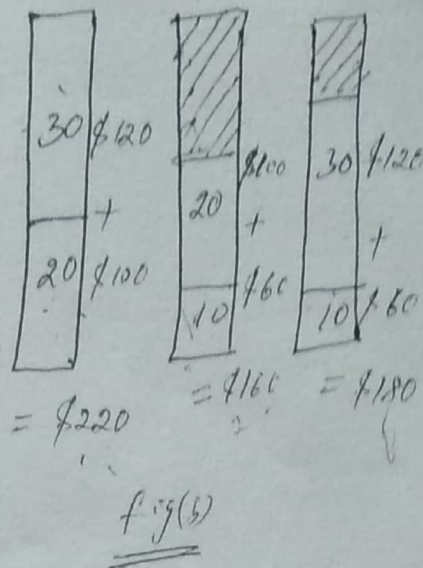
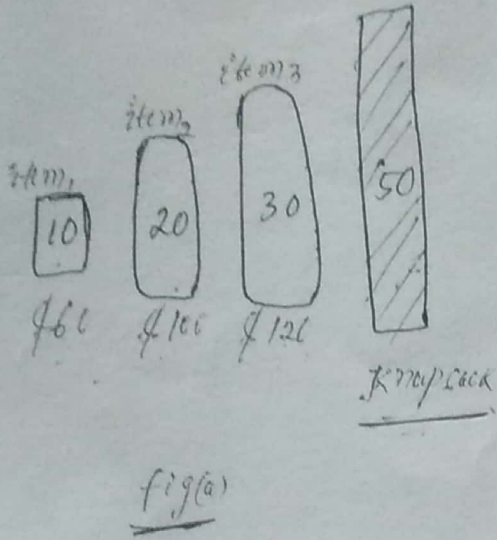
$$w_3 = 30, P_3 = 120 "$$

Thus the value per pound of item 1 is 6 dollars per pounds which is greater than item 2 (5 dollars/pounds) or item 3 (4 \$/p)

→ The greedy strategy therefore would take item 1 first.

→ However optimal solution takes items 2 and 3 leaving behind 1.

→ The two possible solⁿ that involve item 1 are both suboptimal.



- The optimal subset includes items 2 and 3.
- my solⁿ with item 1 is suboptimal. even though item 1 has the greatest value per pound.
- Fig(c) shows the fraction knapsack taking the items in order of greatest value per pound yields an optimal solution.

The methods used to obtain the solution is termed as greedy method because at each step we choose an object which would increase the objective function value the most. However greedy method does not yield an optimal solution.

- To get the optimal solⁿ, at each step we include the object which has the maximum profit / unit capacity used.
- This means objects are arranged by P_i / W_i order.

How do we select the next item to be put into the knapsack?

There are several possibilities:-

- Greedy by profit: At each step select from the remaining items with the highest profit. This approach tries to maximize the profit by choosing the most profitable items first.
- Greedy by weight: with one with least weight - it tries to maximize the profit by putting as many items into the knapsack as possible.
- Greedy by profit density: with the largest profit density, P_i / W_i . it tries to maximize the profit by choosing items with the largest profit per units of weight.

(Example shown in next paper)

Huffman codes ^{→ factor up numbers} are widely used and very effective technique for compressing data, depending on the characteristics of data being compressed.

- Huffman's greedy algorithm uses a table of the frequencies of occurrence of the characters to build up an optimal way of representing each character as a binary string.
- Used for security purpose where data is transferred from sender to receiver.
- There are many ways to represent such a file of information.

Suppose we have a 100,000-character data file that we wish to store compactly. Only six characters appear and the character 'a' occurs 45,000 times.

→ Huffman coding is a lossless data compression also. EO ~~Header~~ OR

	a	b	c	d	e	f
Frequency (in thousands)	45	13	12	16	9	5
Fixed-length code word	000	001	010	011	100	101
Variable-length code word	0	101	100	111	1101	1100

- A character-coding problem.
- A datafile of 100,000 characters contains only the characters a-f.
- If each character is assigned 3-bit code word, the file can be encoded in 300,000 bits.
- Using the variable-length codes shown, the file can be encoded in 224,000 bits.

It is of two types :-
 — Fixed length coding
 — variable length "

- If we use fixed-length code we need 3 bits to represent six characters. This requires 300,000 bits.
- A variable length code can do considerably better than a fixed length code. This coding requires

$$(45 \cdot 1 + 13 \cdot 3 + 12 \cdot 3 + 16 \cdot 3 + 9 \cdot 4 + 5 \cdot 4) 1000 = 224,000 \text{ bits}$$

→ So it saves 28%.

~~prefix codes are widely used in applications that compress data including JPEG for images and MP3 for music.~~

Prefix codes (No codeword is a prefix of some other codeword. Such codes are called prefix code (or prefix free codes).)

→ The standard solⁿ for unique decoding is to insist that the code-word be prefix free.

→ This means that if $a, b \in \Sigma$, $a \neq b$ then code word for a is not a prefix of the code-word for b vice versa.

→ Encoding is always simple for any binary character code, we just concatenate the codewords representing each character of the file.

→ Suppose we code 3-character file abc as $0 \cdot 101 \cdot 100 = 0101100$

→ prefix codes are desirable because they simplify decoding.

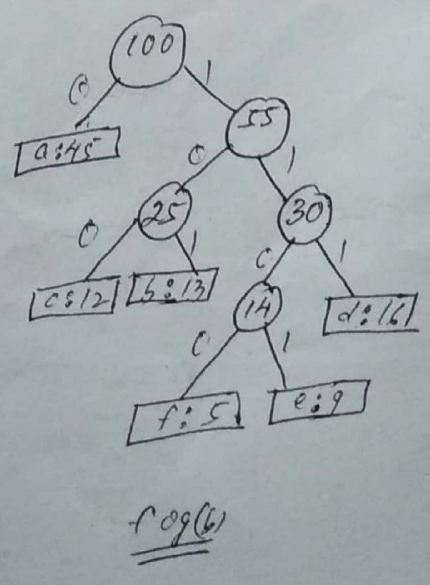
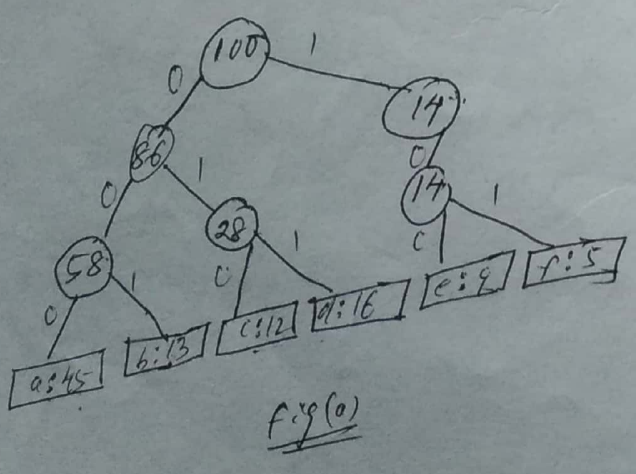
→ Since no codeword is a prefix of any other, the codeword that begin an encoded file is unambiguous.

→ we can simply identify the initial codeword, translate it back to the original character and repeat the decoding process on the remainder of the encoded file.

→ In our example the string 00101101 parses uniquely as $0 \cdot 0 \cdot 101 \cdot 1101$ which decodes to 0abe.

The decoding process needs a convenient representation for the prefix code so that the initial codeword can be easily peeled off.

→ A binary tree whose leaves are the given characters provides one such representation. (0 means left child and 1 means right)



Constructing a Huffman code: Huffman invented a greedy algorithm that constructs an optimal prefix code called a Huffman code.

- C is a set of n characters and that each character $c \in C$ is an object with a defined frequency $f[c]$.
- The algorithm builds the tree T corresponding to the optimal code in bottom-up manner.
- It begins with a set of $|C|$ leaves and performs a sequence of $|C|-1$ "merging" operations to create the final tree.
- A min-priority queue Q , keyed on f , is used to identify the two least-frequent objects to merge together.

HUFFMAN(C)

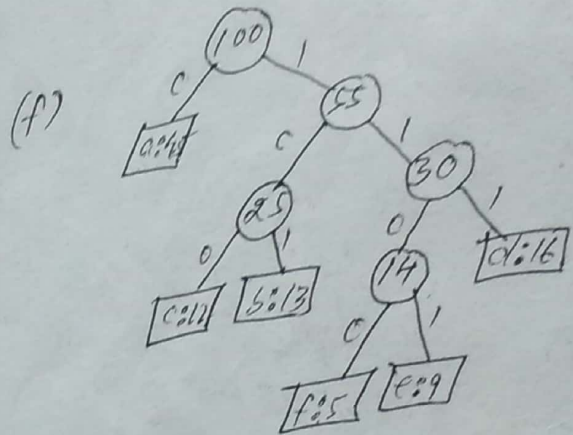
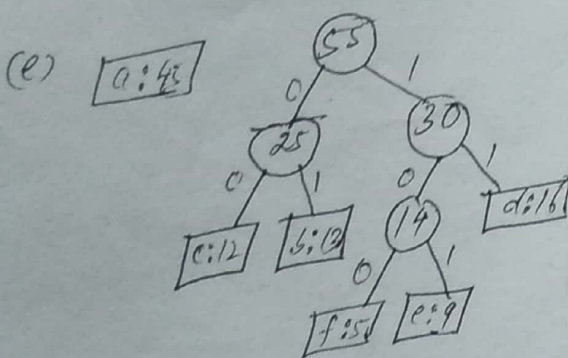
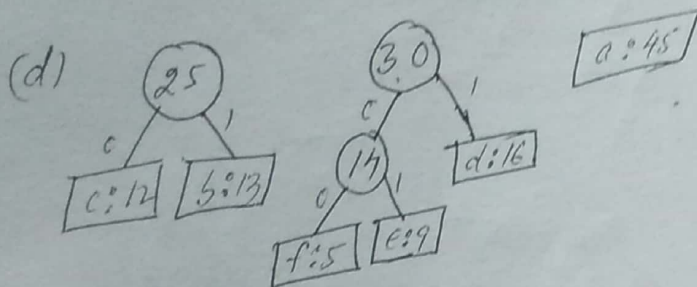
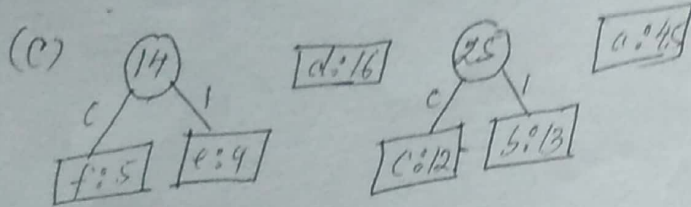
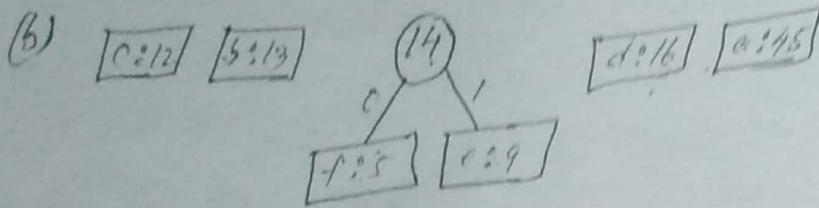
1. $n \leftarrow |C|$
 2. $Q \leftarrow C$ $O(n)$
 3. for $i \leftarrow 1$ to $n-1$
 4. do allocate a new node Z
 5. left[Z] $\leftarrow x \leftarrow \text{EXTRACT-MIN}(Q)$
 6. right[Z] $\leftarrow y \leftarrow \text{EXTRACT-MIN}(Q)$
 7. $f[Z] \leftarrow f[x] + f[y]$
 8. INSERT(Q, Z)
- } 3-8 executed exactly $n-1$ times.
and each heap op. requires $O(\lg n)$
Total $O(n \lg n)$.
9. return EXTRACT-MIN(Q) // returns the root of the tree.

→ Since there are 6 letters in the alphabet, the initial queue size is $n=6$, and 5 merge steps are required to build tree.

→ The final tree represents the optimal prefix code.

→ The codeword for a letter is the sequence of edge labels on the path from the root to the letter.

(a) $f:5$ $e:9$ $c:12$ $b:13$ $d:16$ $a:45$



→ leaves are shown as rectangle.
 → internal node as circle

Ex Build a binary tree using Huffman code for $P = \langle 29, 25, 20, 12, 5, 09 \rangle$

Time complexity The analysis of the running time of Huffman's algo assumes that Q is implemented as a binary min-heap.

→ For a set E of n -characters the initialization of Q in line 2 takes $O(n)$ time using the BUILD-MIN-HEAP procedure.

→ The for loop in lines 3-8 is executed exactly $n-1$ times and since each heap operation requires time $O(\lg n)$, the loop contributes $O(n \lg n)$ to the running time.

Dynamic Programming

Dynamic programming typically applies to optimization problems in which a set of choices must be made in order to arrive at an optimal solution.

It is generally used with optimization problems where a sequence of solutions are available, each solution has a value and we require to find the optimal value, each solution is called an optimal solution (one way to solve optimization problem to try all possible solⁿ and then pick out the best).

It is similar to DANC method, which solves problems by combining the solution to subproblems. But DANC algorithm partition the problem into independent subproblems, solve subproblems recursively and then combine their solutions to solve the original problem. But in DP, the subproblems are not independent. DP solves every subproblem once and then save its answer in a table, thereby avoiding the work of recomputing the answer every time the subproblem is encountered. (save in table at the cost of space)

DANC

- ① Each part is independent
- ② There is no need of table
- ③ Each part solⁿ saving is not required
- ④ Time complexity is more
- ⑤ we may or may not find optimal solution
- ⑥ subproblems of DANC are non-overlap
- ⑦ DANC solves the sub-problem top-down

DP

- ① Each part is dependent
- ② Need of table
- ③ Required
- ④ Less
- ⑤ we must find the optimal solution guaranteed to find solution.
- ⑥ Overlap
- ⑦ DP is bottom-up technique.

NOTE ① DP solves the subproblems once and then gave answer in a table but DANC call repeatedly the subproblems in many times so complexity is more.
 ② in DP we try to find redundancies and reduce the space for search.

- DP is applicable when the subproblems are not independent. that is when subproblems share subsubproblems. So
- DP does more work than necessary (repeatedly solving the common subproblems).

The development of a DP algorithm can be broken into a sequence of four steps.

1. Characterizing the structure of an optimal solution.
2. Recursively define the value of an optimal solution.
3. Compute the optimal solution in a bottom up fashion.
4. Construct an optimal solution from computed information.

bottom-up means:

- ① Start with the smallest subproblems.
 - ② Combining their solⁿ obtain the solⁿ to subproblems of increasing size.
 - ③ Until arrive at the solⁿ of the original solution.
- Try to solve solve subproblem first and use their solⁿ to build-on and arrive at solⁿ to bigger subproblems.

EXTRA ① The word "programming" in DP has no connection to computer program and instead come from the term "mathematical programming", a synonym for optimization.

- ② DP is mainly used to tackle optimization problems that are solvable in polynomial time.

Matrix-chain Multiplication is an example of DP.

We are given a sequence (chain) $\langle A_1, A_2, \dots, A_n \rangle$ of n -matrices to be multiplied and we wish to compute the product $A_1 A_2 \dots A_n$. We can evaluate the expⁿ once we have parenthesized it to resolve all ambiguity in how the matrices are multiplied together.

→ matrix multiplication is associative and so all parenthesizations yield the same product.

Ex: - If the chain of matrices is $\langle A_1, A_2, A_3, A_4 \rangle$, the product $A_1 A_2 A_3 A_4$ can be fully parenthesized in 5 distinct ways:

- $(A_1 (A_2 (A_3 A_4)))$
- $(A_1 ((A_2 A_3) A_4))$
- $((A_1 A_2) (A_3 A_4))$
- $((A_1 (A_2 A_3)) A_4)$
- $((A_1 A_2) A_3) A_4$

→ The way we parenthesize a chain of matrices can have a dramatic impact on the cost of evaluating the product.

Example: - Suppose the dimensions of matrices $\langle A_1, A_2, A_3 \rangle$ are 10×100 , 100×5 and 5×50

i.e. $(A_1)_{10 \times 100}$ $(A_2)_{100 \times 5}$ $(A_3)_{5 \times 50}$

We can multiply the above chain of matrices according to the parenthesizations:

- (i) $(A_1 A_2) A_3$ or
- (ii) $A_1 (A_2 A_3)$

$(A_1 \times A_2) = (A)_{10 \times 5}$

$(A \times A_3) = (B)_{10 \times 50}$

No. of scalar multiplications $10 \times 100 \times 5 = 5000$
 " " " " " " $10 \times 5 \times 50 = 2500$

total = 7500 ops

need for $(A_1 A_2) A_3$

If we multiply according to (ii) then we have

$(A_2 A_3) = (C)_{100 \times 50}$

$(A_1 \times C) = (D)_{10 \times 50}$

No. of scalar multiplications $100 \times 5 \times 50 = 25000$
 " " " " " " $10 \times 100 \times 50 = 50,000$

75,000 ops

- Thus computing the product according to first parenthesis is 10 times faster
- Note that in the matrix-chain multiplication (MEM), we are not actually multiplying matrices.
- Our goal is only to determine an order for multiplying matrices that has the lowest cost.

MEM can be stated as follows: given a chain (A_1, A_2, \dots, A_n) of n matrices where for $i = 1, 2, \dots, n$ matrix A_i has dimension $p_{i-1} \times p_i$. Fully parenthesize the product $A_1 A_2 \dots A_n$ in a way that minimize the no. of scalar multiplications.

- Any parenthesization of the product $A_i A_{i+1} \dots A_j$ must split the product between A_k and A_{k+1} for some constant k , $i \leq k \leq j$.
- The cost of the parenthesization is thus the cost of computing the matrix $A_i \dots A_k$ + cost of computing $A_{k+1} \dots A_j$ + (plus) the cost of multiplying them together.

Let us define the cost of an optimal solution recursively in terms of the optimal solutions to subproblems. For MEM we pick as our subproblems the problems of determining the minimum cost of a parenthesization of $A_i A_{i+1} \dots A_j$ for $1 \leq i \leq j \leq n$.

Let $m[i, j]$ = minimum no. of scalar multiplications needed to compute the matrix $A_i \dots A_j$.

i.e. $(A_i A_{i+1} \dots A_k) (A_{k+1} \dots A_j)$
 $p_{i-1} \times p_k \quad p_k \times p_j$

i.e. $A_i \dots A_j = (A_i \dots A_k) (A_{k+1} \dots A_j)$
 $p_{i-1} \times p_k \quad p_k \times p_j$

→ If $i = j$, the problem is trivial so $m[i, i] = 0$.

$$m[i, j] = m[i, k] + m[k+1, j] + p_{i-1} p_k p_j$$

↓
↓
↓

cost of computing subproduct $A_i \dots A_k$ cost of computing $A_{k+1} \dots A_j$ cost of computing $(A_i \dots A_k) (A_{k+1} \dots A_j)$

→ This recursive eqⁿ assumes that we know the value of k , which don't. There are only $j-i$ possible values for k .

$$k \in \{i, i+1, \dots, j-1\}$$

→ To find the best we use one of these values for k .

$$m[3,6] = \min \begin{cases} m[3,3] + m[4,6] + p_2 p_2 p_6 \\ m[3,4] + m[5,6] + p_2 p_4 p_6 \\ m[3,5] + m[6,6] + p_2 p_5 p_6 \end{cases}$$

$$= \min \begin{cases} 0 + 3500 + 15 \times 5 \times 25 = 5375 \\ 750 + 5000 + 15 \times 10 \times 25 = 7500 \\ 2500 + 0 + 15 \times 20 \times 25 = 10000 \end{cases}$$

$$s[3,6] = 3$$

$$m[1,5] = \min \begin{cases} m[1,1] + m[2,5] + p_0 p_1 p_5 \\ m[1,2] + m[3,5] + p_0 p_2 p_5 \\ m[1,3] + m[4,5] + p_0 p_3 p_5 \\ m[1,4] + m[5,5] + p_0 p_4 p_5 \end{cases}$$

$$= \min \begin{cases} 0 + 7125 + 30 \times 35 \times 20 = 28125 \\ 15750 + 2500 + 30 \times 15 \times 20 = 27250 \\ 7875 + 1000 + 30 \times 5 \times 20 = 11875 \\ 9375 + 0 + 30 \times 10 \times 20 = 15375 \end{cases}$$

$$s[1,5] = 3$$

$$m[2,6] = \min \begin{cases} m[2,2] + m[3,6] + p_1 p_2 p_6 \\ m[2,3] + m[4,6] + p_1 p_3 p_6 \\ m[2,4] + m[5,6] + p_1 p_4 p_6 \\ m[2,5] + m[6,6] + p_1 p_5 p_6 \end{cases}$$

$$= \min \begin{cases} 0 + 5375 + 35 \times 15 \times 25 = 18500 \\ 2625 + 3500 + 35 \times 5 \times 25 = 10500 \\ 4375 + 5000 + 35 \times 10 \times 25 = 15125 \\ 7125 + 0 + 35 \times 20 \times 25 = 24625 \end{cases}$$

$$s[2,6] = 3$$

$$m[1,6] = \min \begin{cases} m[1,1] + m[2,6] + p_0 p_1 p_6 \\ m[1,2] + m[3,6] + p_0 p_2 p_6 \\ m[1,3] + m[4,6] + p_0 p_3 p_6 \\ m[1,4] + m[5,6] + p_0 p_4 p_6 \\ m[1,5] + m[6,6] + p_0 p_5 p_6 \end{cases}$$

$$= \min \begin{cases} 0 + 10500 + 30 \times 35 \times 25 = 36750 \\ 15750 + 5375 + 30 \times 15 \times 25 = 32375 \\ 7875 + 3500 + 30 \times 5 \times 25 = 15125 \\ 9375 + 5000 + 30 \times 10 \times 25 = 21875 \\ 11875 + 0 + 30 \times 20 \times 25 = 26875 \end{cases}$$

$$\therefore s[1,6] = 3$$

Although MATRIX-CHAIN-ORDER determines the optimal no. of scalar multiplications needed to compute a matrix-chain product, it does not directly show how to multiply the matrices. It is not difficult to construct an optimal solution from the computed information stored in the table $S[1..n, 1..n]$.

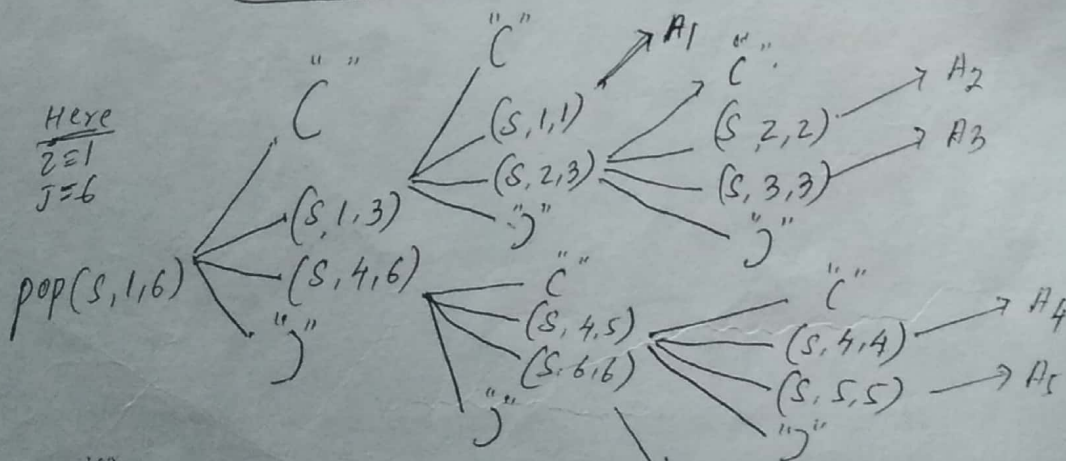
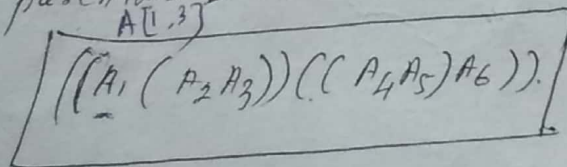
The following recursive procedure prints an optimal parenthesization of $\langle A_i, A_{i+1}, \dots, A_j \rangle$.

PRINT-OPTIMAL-PARENS (S, i, j)

1. if $i = j$
2. then print " A_i "
3. else print "("
4. PRINT-OPTIMAL-PARENS($S, i, S[i, j]$)
5. PRINT-OPTIMAL-PARENS($S, S[i, j] + 1, j$)
6. print ")"

Initially $i=1, j=6$
 $\rightarrow S, 1, (1,6)$
 value $[1,6]$
 $\rightarrow (S, 1, 3)$
 Now $i=1, j=3$
 Now $(1,1)$
 (A_1)
 $S[i, j]$ stored in S
 $(C=2)$
 (A_2)

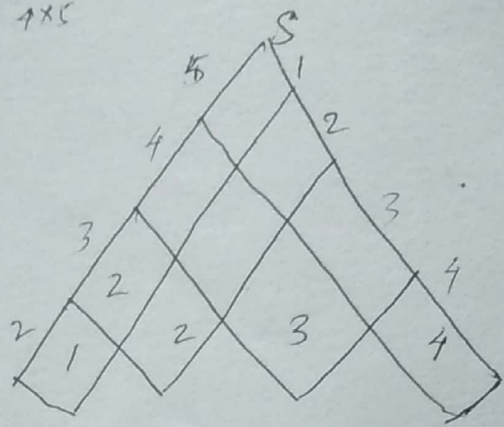
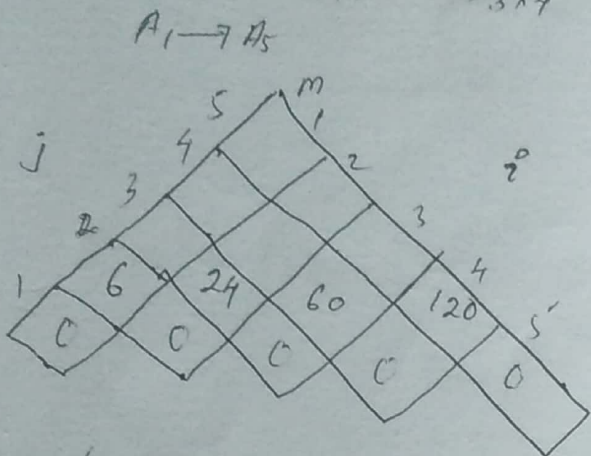
In the above example the call PRINT-OPTIMAL-PARENS($S, 1, 6$) prints the parenthesization



Remember When $i = j$, it will print the matrix " A_i "
 when $i \neq j$ the problem follows 4 steps

- print "("
- pop($S, i, S[i, j]$)
- pop($S, S[i, j] + 1, j$)
- print ")"

Ex $(A_1)_{1 \times 2}$ $(A_2)_{2 \times 3}$ $(A_3)_{3 \times 4}$ $(A_4)_{4 \times 5}$ $(A_5)_{5 \times 6}$



put $l=2$, $i \rightarrow 1$ to 4 , for $i=1, j=2$ $m[1,2]$, $m[2,3]$
 $m[3,4]$ $m[4,5]$

for $k=1$; $m[1,2] = m[1,1] + m[2,2] + p_0 p_1 p_2$
 for $k=2$; $m[2,3] = m[2,2] + m[3,3] + p_1 p_2 p_3$
 for $k=3$; $m[3,4] = m[3,3] + m[4,4] + p_2 p_3 p_4$
 for $k=4$; $m[4,5] = m[4,4] + m[5,5] + p_3 p_4 p_5$

for $l=3$, $i \rightarrow 1$ to 3 $j=3$ $m[1,3]$, $m[2,4]$, $m[3,5]$

for $k=1$ $m[1,3] = m[1,1] + m[2,3] + p_0 p_1 p_3 = 32$
 for $k=2$ $m[1,3] = m[1,2] + m[3,3] + p_0 p_2 p_3 = 6 + 0 + 12 = 18$

Ex - find the optimal parenthesis of $5, 10, 12, 2, 15, 16, 11$

Ans: - $A_1 = 5 \times 10$ $A_3 = 12 \times 2$ $A_5 = 15 \times 16$
 $A_2 = 10 \times 12$ $A_4 = 2 \times 15$ $A_6 = 16 \times 11$

$p_0 p_1 p_2 p_3 p_4 p_5 p_6$ are given in equation

Ex - $A_1 = 2 \times 3$ $p_0 = 2$ find (i) $m[1,4]$
 $A_2 = 3 \times 4$ $p_1 = 3$ (ii) $m[1,3]$
 $A_3 = 4 \times 5$ $p_2 = 4$ (iii) draw the parenthesis
 $A_4 = 5 \times 6$ $p_3 = 5$
 $p_4 = 6$

$m[1,2] = 24$ ($k=1$) $m[1,3] = 64$ ($k=2$)
 $m[2,3] = 60$ ($k=2$) $m[2,4] = 150$ ($k=3$)
 $m[3,4] = 120$ ($k=3$) $m[1,4] = 124$ ($k=3$)

optimal = $(((A_1 A_2) A_3) A_4)$

Longest Common Subsequence :-

→ A subsequence of a given sequence is just the given sequence with zero or more elements left out.

✓ A subsequence is any subset of the elements of a sequence that maintains the same relative order.

→ Given a sequence $x = \langle x_1, x_2, \dots, x_n \rangle$, another sequence $z = \langle z_1, z_2, \dots, z_k \rangle$ is a subsequence of x if there exists a strictly increasing order sequence $\langle i_1, i_2, \dots, i_k \rangle$ of indices of x such that for all $j = 1, 2, \dots, k$, we have

$$x_{i_j} = z_j$$

✓ Ex - $z = \langle B, C, D, B \rangle$ is a subsequence of $x = \langle A, B, C, B, D, A, B \rangle$ with corresponding index subsequence $\langle 2, 3, 5, 7 \rangle$

✓ → Given two sequences x and y , we say that a sequence z is a common subsequence of x and y if z is subsequence of both x and y .

Ex - If $x = \langle A, B, C, B, D, A, B \rangle$

$y = \langle B, D, C, A, B, A \rangle$

- For position A of x
(i) A B A is a common subsequence of x, y
- For position B of x
(ii) B C A B is a " " " "
- For position C of x
(iii) C A B " " " "
- For position D of x
(iv) D A B " " " "
- For position A of y
(v) A B " " " "
- For position B of y
(vi) B " " " "

→ In the longest common subsequence (LCS) problem, also called LPS, we are given two sequences and wish to find a maximum-length common subsequence of X and Y.
 One way to solve the LCS problem is to enumerate all subsequences of X and check the longest one that is also a subsequence of Y.

Several motivation applications

- (1) molecular biology
- (2) file comparison
- (3) Screen display

When we look for a good solⁿ to a problem, we examine the two key ingredients that an optimization problem must have in order for DP to apply.

Elements of Dynamic programming

The problem has the following properties:

I. optimal substructure. A solⁿ is optimal only if \bar{Z} subproblems sub-structure to the subproblem is optimal.

[exists optimal solⁿ in the problem contains an optimal solⁿ to subproblems]

Overlapping subproblems

- The same subproblem is visited over and over again.
- The # of distinct subproblems is polynomial.

OR A problem is said to have overlapping subproblems if the problem can be broken down into subproblems which are re-used several times or a recursive algo for the problem solves the same subproblem over and over rather than always generating new subproblems.

OR
 when a recursive algo revisits the same problems repeatedly we say that the optimization problem has overlapping subproblems.

Algo. ^{fall} 1. Characterizing a LCS: - A brute-force approach to solving the LCS is to enumerate all subsequences of X and check each subsequence to see if it is also a subsequence of Y , keeping track of the LCS found.

The LCS problem has an optimal-substructure property

Optimal substructure of an LCS

Let $X = \langle x_1, x_2, \dots, x_m \rangle$ and $Y = \langle y_1, y_2, \dots, y_n \rangle$ be sequences and let $Z = \langle z_1, z_2, \dots, z_k \rangle$ be any LCS of X and Y .

1. If $x_m = y_n$, then $z_k = x_m = y_n$ and Z_{k-1} is an LCS of X_{m-1} and Y_{n-1} .
2. If $x_m \neq y_n$, then $z_k \neq x_m$ implies that Z is an LCS of X_{m-1} and Y .
3. If $x_m \neq y_n$, then $z_k \neq y_n$ implies that Z is an LCS of X and Y_{n-1} .

2. Recursive solution previous section implies that if $x_m = y_n$ we must find an LCS of X_{m-1} and Y_{n-1} . Appending $x_m = y_n$ to this LCS yields an LCS of X and Y .

If $x_m \neq y_n$ the we must solve two subproblems:

- finding an LCS of X_{m-1} and Y
- " " " " X and Y_{n-1}

whichever of these two LCS's is longer is an LCS of X and Y .

- Let $c[i, j]$ length of an LCS of the sequences X_i and Y_j .

$$c[i, j] = \begin{cases} 0 & \text{if } i=0 \text{ or } j=0 \\ c[i-1, j-1]+1 & \text{if } i, j > 0 \text{ and } x_i = y_j \\ \max(c[i, j-1], c[i-1, j]) & \text{if } i, j > 0 \text{ and } x_i \neq y_j \end{cases}$$

Example: Let $X = \langle x_1, x_2, \dots, x_n \rangle$

$Y = \langle \quad \quad \quad \rangle$ no elements.

So length of $X = L(X) = n$

" " $Y = L(Y) = 0$

// Brute force approach

Algorithm to build the LCS: (computing the length of LCS)

LCS-length (X, Y)

1. $m \leftarrow \text{length}[X]$
2. $n \leftarrow \text{length}[Y]$
3. for $i \leftarrow 1$ to m } $O(m)$
do $c[i, 0] \leftarrow 0$
- 4.
5. for $j \leftarrow 1$ to n } $O(m)$
do $c[0, j] \leftarrow 0$
- 6.
7. for $i \leftarrow 1$ to m } $O(m \times n)$
do for $j \leftarrow 1$ to n
8. do if $x_i = y_j$
9. then $c[i, j] \leftarrow c[i-1, j-1] + 1$
10. $b[i, j] \leftarrow \leftarrow$
11. else if $c[i-1, j] \geq c[i, j-1]$
12. then $c[i, j] \leftarrow c[i-1, j]$
13. $b[i, j] \leftarrow \uparrow$
14. else $c[i, j] \leftarrow c[i, j-1]$
15. $b[i, j] \leftarrow \leftarrow$
- 16.
17. return c and b

find LCS for $X = \langle A, B, C, B, D, A, B \rangle$
 $Y = \langle B, D, C, A, B, A \rangle$

It stores the $c[i, j]$ value in a table $c[0..m, 0..n]$
whose entries are computed in row-major order.
It also maintains the table $b[1..m, 1..n]$ to help in
construction of an optimal solution.

x_i	y_j	0	1	2	3	4	5	6
0	0	0	0	0	0	0	0	0
1	A	0	0↑	0↑	0↑	1	1←	1
2	B	0	0	1←	1←	1↑	2←	2
3	C	0	1↑	1↑	2	2	2↑	2
4	B	0	1←	1↑	2↑	2↑	3	3
5	D	0	1↑	2	2↑	2↑	3↑	3↑
6	A	0	1↑	2↑	2↑	3	3	4
7	B	0	1	2↑	2↑	3	4	4

B C B A

The running time of the procedure is $O(mn)$, hence each table entry $O(1)$ time to compute. [or $T(n) = O(m) + O(n) + O(m \times n)$]

Constructing in LCS

The b table returned by `LCS-LENGTH`

can be used to quickly construct an LCS of x and y .

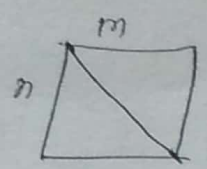
→ we simply begin at $b[m, n]$ and trace through the table following the arrows

→ whenever we encounter a "K" in entry $b[i, j]$ it implies that $x_i = y_j$ is an element of the LCS.

→ the elements of procedure for the LCS are encountered in reverse order by this method.

PRINT-LCS (b, x, i, j) // find the optimal value

1. if $i=0$ or $j=0$ then return ""
2. if $b[i, j] \neq \uparrow$ then PRINT-LCS (b, x, i-1, j-1)
3. print x_i
4. else if $b[i, j] = \uparrow$ then PRINT-LCS (b, x, i-1, j)
5. else PRINT-LCS (b, x, i, j-1)



→ This procedure prints "B C B A".
 → This procedure takes time $O(m+n)$, since at least one of i and j is decremented in each stage of the recursion.

EXERCISE Determine an LCS of $A = \{1, 0, 0, 1, 0, 1, 0, 1\}$ $m=8$
 $B = \{0, 1, 0, 1, 1, 0, 1, 0\}$ $n=9$

Minimum Spanning Trees

Suppose $G = (V, E)$ is a connected undirected graph and for each edge $(u, v) \in E$, we have a weight $w(u, v)$ specifying the cost to connect u and v .

→ We then wish to find an acyclic subset $T \subseteq E$ that connects all of the vertices and whose total weight

$$w(T) = \sum_{(u,v) \in T} w(u,v) \text{ is minimized.}$$

→ Since T is acyclic and connects all of the vertices, it must form a tree, which we call a spanning tree. Since it spans the graph G , we call the problem of determining the tree T the minimum-spanning-tree problem.

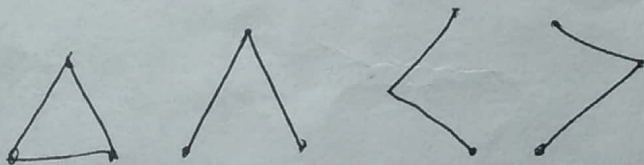
OR

→ Let a graph $G = (V, E)$ if T is a sub-graph and contains all the vertices but no cycle then T is said to be a spanning tree.
→ A connected graph is a tree if and only if it has n -vertices and $n-1$ edges.

OR

Spanning tree of G is a graph H if the spanning tree of the graph G , if

- (i) H is the subgraph of G
- (ii) H is a tree (connected tree)
- (iii) H contains all the vertices of G .

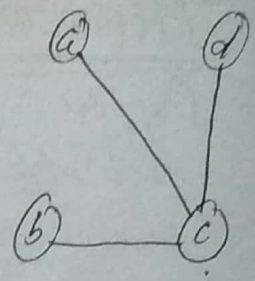
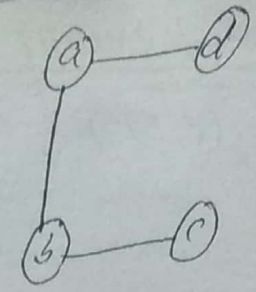
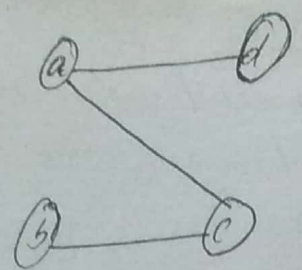
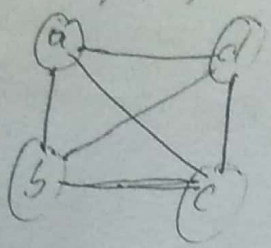


(Spanning tree of G)

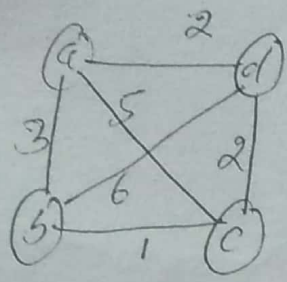
→ If the graph is a complete graph with n -vertices, then it can have n^{n-2} spanning trees

$$\text{No. of spanning trees} = n^{n-2}$$

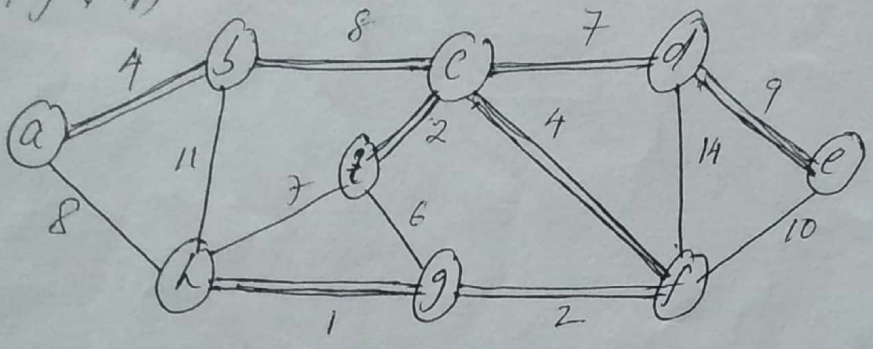
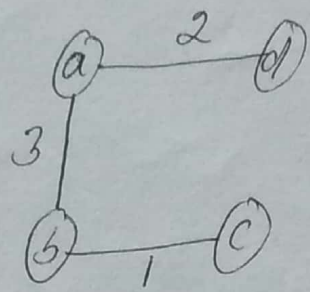
A graph G



(Tree of many possible spanning tree of G)



(a weighted graph G)



→ The total weight of the tree shown is 37

→ This MST is not unique:
 removing the edge (b,c) and replacing it with the edge (a,h) yields another tree with weight 37.

Growing a minimum spanning tree: We have a connected, undirected graph $G = (V, E)$ with a weight function $w: E \rightarrow \mathbb{R}$ and we wish to find a MST for G .

→ Two algorithms we consider in this case a greedy approach to the problem, although they differ in how they apply this approach
 → This greedy strategy is captured by the following "generic" algorithm which grows the minimum spanning tree one edge at a time.

GENERIC-MST (G, w)

1. $A \leftarrow \emptyset$
2. while A does not form a spanning tree
3. do find an edge (u, v) that is safe for A
4. $A \leftarrow A \cup \{(u, v)\}$
5. return A

Cut :- A cut $(S, V-S)$ of an undirected graph $G = (V, E)$ is a partition of V .

Cut vertex :- From the connected graph G when we remove a vertex v from the graph and $G-v$ becomes a graph which is not connected, then the vertex v is called cut vertex.

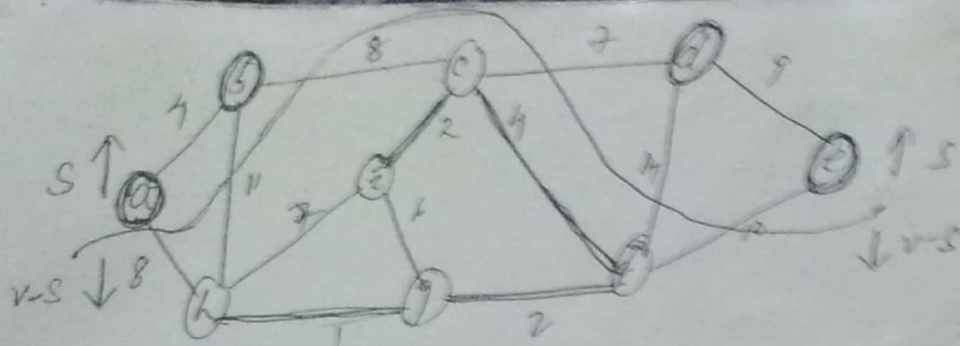
Cut set :- The cutset is a set of edges whose removal from the graph G , makes the graph disconnected then that set of edges are called cutset. (It is minimal no. of edges whose removal from G destroys all paths between these two sets of vertices).

Cross edge An edge $(u, v) \in E$ crosses the cut $(S, V-S)$ if one of its endpoints is in S and the other is in $V-S$.

⊗ cross edge are the edges which belong to the cutset.

Light edge :- Among the cross edges, whose weight is minimum is called light edge. It is the minimum weight among the cross edge.

An edge is a light edge crossing a cut if its weight is the minimum of any edge crossing the cut.



- The vertices in the set S are shown in black and those in $v-S$ are shown in white.
- The edges crossing the cut are those connecting white vertices with black vertices.
- The edge (d, c) is the unique light edge crossing the cut.
- A subset A of the edges is shaded, since no edge of A crosses the cut.

Single-Source Shortest paths: In a shortest-path problem we are given a weighted, directed graph $G=(V, E)$ with weight function $w: E \rightarrow \mathbb{R}$ mapping edges to real-valued weights.

→ The weight of path $P = \langle v_0, v_1, v_2, \dots, v_k \rangle$ is the sum of the weights of its constituent edges:

$$w(P) = \sum_{i=1}^k w(v_{i-1}, v_i)$$

We define the shortest-path weight from u to v by
$$s(u, v) = \begin{cases} \min \{ w(P) : u \overset{P}{\rightsquigarrow} v \} & \text{if there is a path from } u \text{ to } v. \\ \infty & \text{otherwise} \end{cases}$$

- A shortest path from vertex u to vertex v is then defined as any path p with weight $w(p) = s(u, v)$
- Edge weight used to represent time, cost, penalties, loss or any other quantity.

Variants

- (1) Single-source shortest-path: From the given graph we want to find a shortest path from a given source vertex $s \in V$ to each vertex $v \in V$.
- (2) Single-destination shortest-path: Find a shortest path to a given destination vertex from each vertex v . By reversing the direction of each edge in the graph, we can reduce this problem to a single-source problem.
- (3) Single-pair shortest-path problem: Find a shortest path from u to v for given vertices u and v . If we solve the single-source problem with source vertex u we solve this problem also.
- (4) All-pairs shortest-path problem: Find a shortest path from u to v for every pair of vertices u and v .

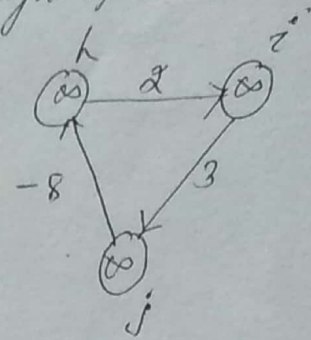
Negative-weight edges In some instances of the single-source shortest path problems, there may be edges whose weights are negative.

Negative-weight cycle: cycle whose sum of weights are $-ve$. Here vertices h, i and j also form a $-ve$ weight cycle.

→ They are reachable from source s so their shortest path weights are ∞ .

$$\boxed{f(s, h) = f(s, i) = f(s, j) = \infty}$$

even though they lie on $-ve$ weight cycle.



→ Let us discuss the given figure below

→ There is only one path from s to a (the path $\langle s, a \rangle$).

$$f(s, a) = w(s, a) = 3.$$

→ Similarly there is only one path from s to b and so

$$f(s, b) = w(s, e) + w(e, b) = 3 + (-4) = -1.$$

→ There are infinitely many paths from s to c :

$\langle s, c \rangle, \langle s, c, d, c \rangle, \langle s, c, d, c, d, c \rangle$ and so on.

Because the cycle $\langle c, d, c \rangle$ has weight $6 + (-3) = \boxed{3 > 0}$, the shortest path from s to c is $\langle s, c \rangle$ with $f(s, c) = 5$.

→ Similarly the shortest path from s to d is $\langle s, c, d \rangle$ with weight $f(s, d) = w(s, c) + w(c, d) = 11$.

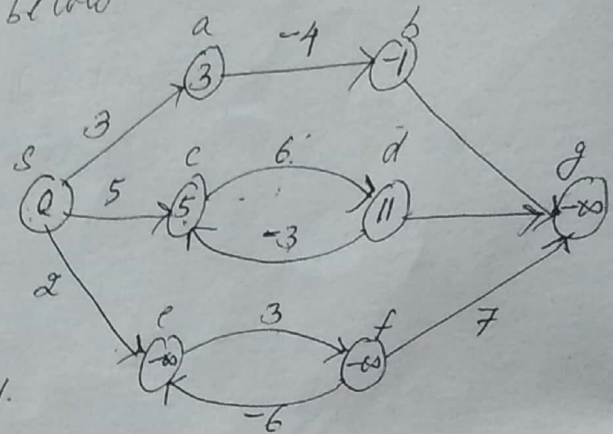
→ But there are infinitely many paths from s to e :

$\langle s, e \rangle, \langle s, e, f, e \rangle, \langle s, e, f, e, f, e \rangle$ and so on.

Since the cycle $\langle e, f, e \rangle$ has weight $3 + (-6) = \boxed{-3 < 0}$, however there is not shortest path from s to e . By traversing the $-ve$ weight cycle $\langle e, f, e \rangle$ many times we can find path from $\langle s, e \rangle$

$$\boxed{f(s, e) = -\infty} \quad \text{Similarly} \quad \boxed{f(s, f) = -\infty}$$

→ vertices h, i, j also form $-ve$ weight cycle. they are not reachable from s so $f(s, h) = f(s, i) = f(s, j) = \infty$.



→ If there is a -ve weight cycle and reachable from source s then $\delta(s, \text{that vertex}) = -\infty$.

→ If there is a -ve weight cycle but not reachable from s then $\delta(s, \text{that vertex}) = \infty$.

→ Remember: → Dijkstra's algorithm assume that all edge weights in the input graph are non-negative. (+ve)

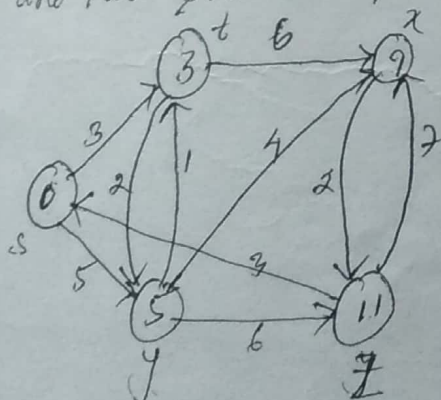
→ Bellman-Ford algorithm allow negative-weight edges in the input graph and produce a correct answer as long as no -ve weight cycles are reachable from the source.

Shortest-path tree: rooted at s is a directed graph $G' = (V', E')$ where $V' \subseteq V$ and $E' \subseteq E$, such that

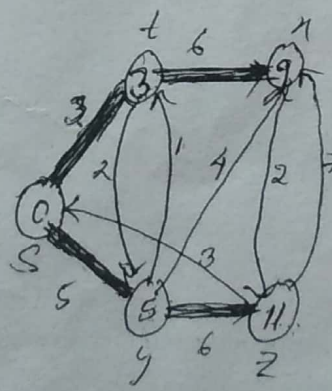
1. V' is the set of vertices reachable from s in G .
2. G' forms a rooted tree with root s , and
3. for all $v \in V'$, the unique simple path from s to v in G' is a shortest path from s to v in G .

Shortest paths are not necessarily unique and neither are shortest-path trees.

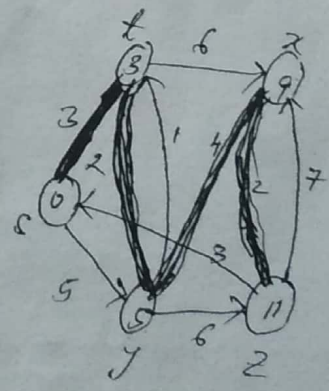
For example below diagram shows a weighted, directed graph and two shortest-path trees with the same root.



(a)



(b)



(c)

→ The shaded edges form a shortest-path tree rooted at the source s .

→ In (b), another shortest path tree with the same root.

Relaxation: - The process of relaxing an edge (u, v) consists of testing whether we can improve the shortest path to v found so far by going through u and if so, updating $d[v]$ and $\pi(v)$.

→ For each vertex $v \in V$, we maintain an attribute $d[v]$, which is an upper bound on the weight of a shortest path from s to v .

→ We call $d[v]$ a shortest-path estimate.

We initialize the shortest-path estimates and predecessors by the following $O(V)$ -time procedure.

✓ INITIALIZE-SINGLE-SOURCE (G, s)

1. for each vertex $v \in V[G]$
2. do $d[v] \leftarrow \infty$ // distance of vertex is ∞
3. $\pi[v] \leftarrow \text{NIL}$ // Parent of v is nil i.e. each is a child.
4. $d[s] \leftarrow 0$.

Here $d[s] = 0$ for $v \in V - \{s\}$

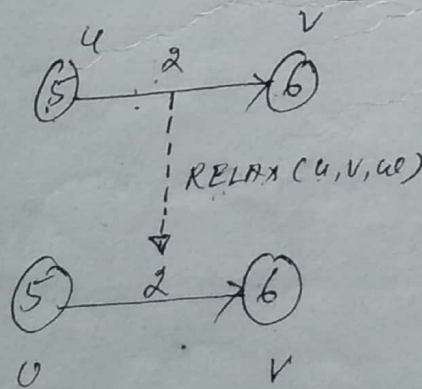
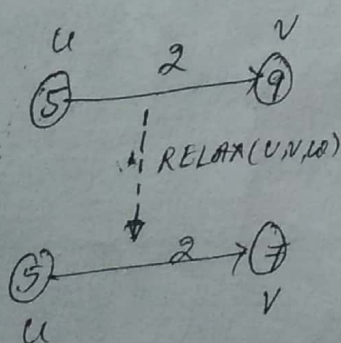
→ A relaxation step may decrease the value of the shortest-path estimate $d[v]$ and update v 's predecessor field $\pi(v)$.

→ the following code performs a relaxation step on edge (u, v)

RELAX (u, v, w)

1. if $d[v] > d[u] + w(u, v)$
2. then $d[v] \leftarrow d[u] + w(u, v)$
3. $\pi[v] \leftarrow u$

NOTE (1) In Dijkstra and shortest-path algo for directed acyclic graph each edge is relaxed exactly once.
 (2) In Bellman-Ford algo each edge is relaxed many times.



→ Relaxation of an edge (u, v) with weight $w(u, v) = 2$.

→ In fig(a) $d[v] > d[u] + w(u, v)$ prior to relaxation, the value of $d[v]$ decreases. But in fig(b) $d[v]$ is unchanged by relaxation.

Breadth First Search (BFS) :- BFS used for searching a graph.

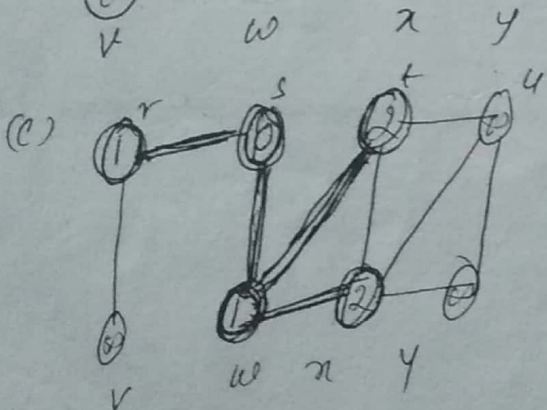
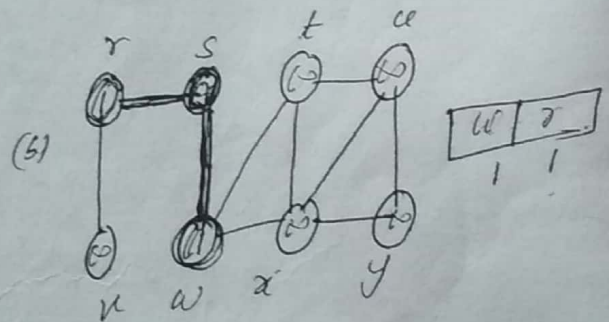
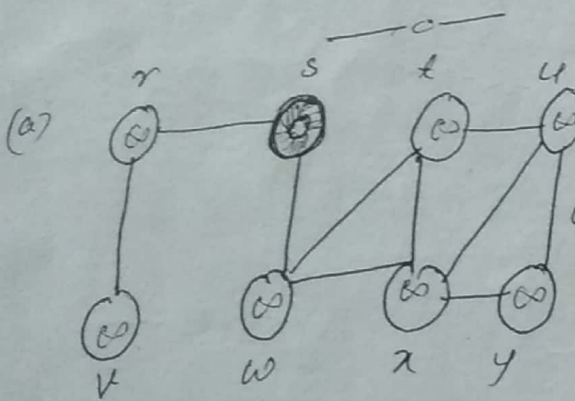
- BFS is named because it expands the ^{border} frontier between discovered and undiscovered vertices uniformly across the breadth of the frontier.
- Given a graph $G=(V, E)$ and a source vertex s , BFS asymptotically explores the edges of G to "discover" every vertex that is reachable from s .
- It computes the distance (smallest no. of edges) from s to each reachable vertex.
- It also produces a "BFS tree" with root s that contains all reachable vertices.
- This algorithm works on both directed and undirected graph.

To keep track of progress, BFS colors each vertex white, gray or black. Initially all vertices are white.

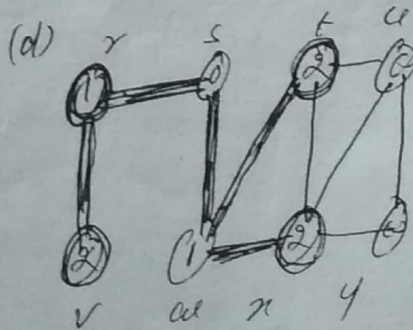
- When the vertex is discovered first time, it becomes non-white.
- If $(u, v) \in E$ and vertex u is black, then vertex v is either gray or black; i.e. all vertices adjacent to black vertices have been discovered.
- Gray vertices may have some adjacent white vertices; they represent the frontier (border) between discovered and undiscovered.
- Since a vertex is discovered at most once, it has at most one parent.
- BFS assumes input graph G is represented using adjacency lists.

BFS (G, s)

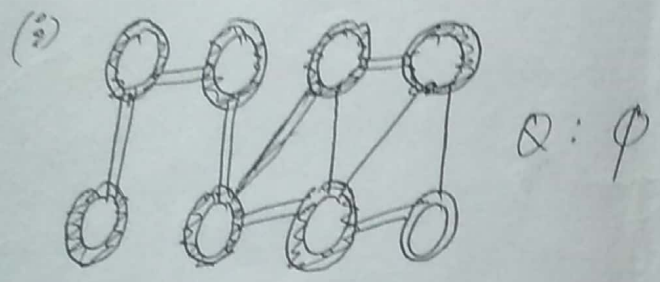
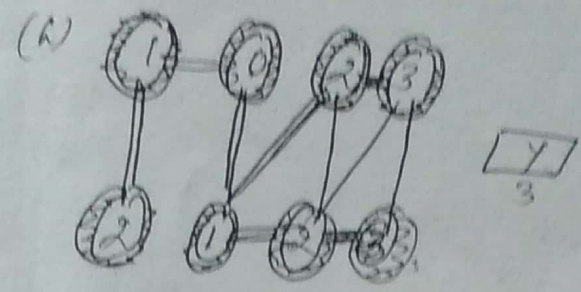
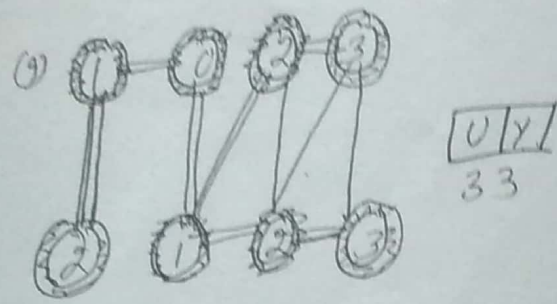
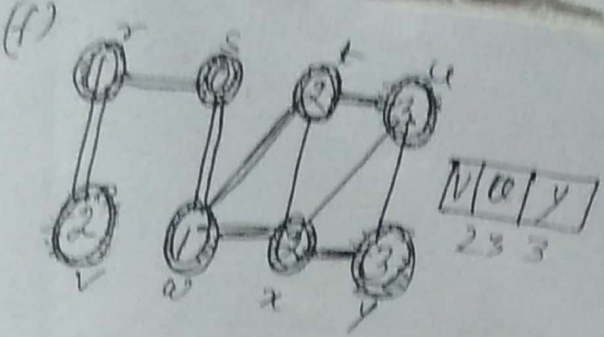
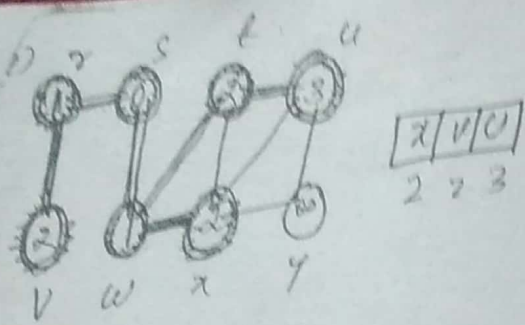
1. for each vertex $u \in V(G) - \{s\}$
2. do $color[u] \leftarrow white$ // initially all vertices white
3. $d[u] \leftarrow \infty$ // distance from s to u
4. $\pi[u] \leftarrow NIL$ // predecessor of u if not there.
5. $color[s] \leftarrow GRAY$ // each step takes $O(1)$
6. $d[s] \leftarrow 0$
7. $\pi[s] \leftarrow NIL$
8. $Q \leftarrow \emptyset$ // $O(V)$
9. **ENQUEUE**(Q, s)
10. while $Q \neq \emptyset$
11. do $u \leftarrow$ **DEQUEUE**(Q).
12. for each $v \in Adj[u]$
13. do if $color[v] = WHITE$
14. then $color[v] \leftarrow GRAY$
15. $d[v] \leftarrow d[u] + 1$
16. $\pi[v] \leftarrow u$
17. **ENQUEUE**(Q, v)
18. $color[u] \leftarrow BLACK$



r	t	x
1	2	2



t	x	v
2	2	2



Time complexity

→ The opⁿ of enqueueing and dequeueing take $O(1)$ time since line-13 executes. So total time devoted to queue opⁿ is $O(V)$.

→ Because the adjacency list of each vertex is scanned only when the vertex is dequeued, each adjacency list is scanned at most once.

Since the sum of the lengths of all the adjacency lists is $O(E)$, the total time spent in scanning adjacency list is $O(E)$.

→ The overhead for initialization is $O(V)$

→ Thus total time of BFS is $O(V+E)$

OK
Step 1 and 2 : $O(V)$
Step 3 ~~step~~ : $O(V)$
Step 4 : $O(V)$
Step 5 : $O(V)$
Step 6 : $O(V)$

Step 7 : $O(V)$
Step 8 : $|V|$
Step 9 to 11 : $O(V+E)$
Step 12 : $O(V)$ or $|V|$
Step 13 : $O(V)$
Step 14 : $O(V)$

So time complexity

$$T(n) = O(V+E)$$

Application of BFS

- (1) Finding all nodes which are in one connected component.
- (2) Finding shortest path between two nodes u and v .
- (3) Testing a graph for bipartiteness.

Depth-first search (DFS): - As its name implies it search "deeper" in the graph whenever possible. (6)

→ Edges are explored out of the most recently discovered vertex v that still has unexplored edges leaving it.

→ when all of v 's edges have been explored, the search "backtracks" to explore edges leaving the vertex from which v was discovered. The process continues.

→ DFS timestamps each vertex. Each vertex u has two timestamps: $d[u]$ records when u is first discovered (and grayed); $f[u]$ " " the search finishes examining u .

v 's adjacency list (and blocks v).

→ $d[u]$ records when it discovers vertex u

→ $f[u]$ " " finishes vertex u

DFS(G)

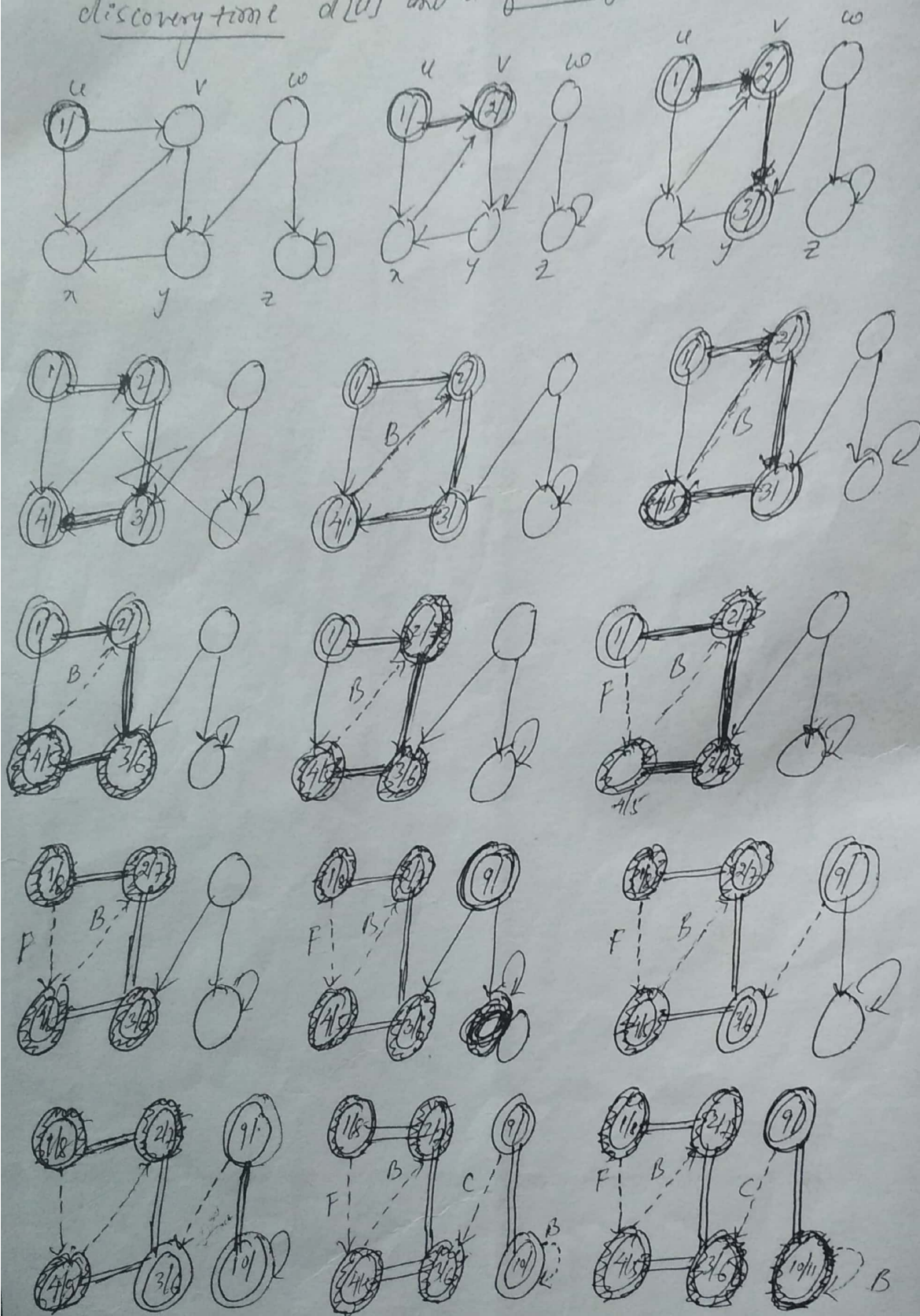
1. for each vertex $u \in V[G]$
2. do $color[u] \leftarrow WHITE$
3. $\pi[u] \leftarrow NIL$
4. $time \leftarrow 0$
5. for each vertex $u \in V[G]$
6. do if $color[u] = WHITE$
7. then $DFS-VISIT(u)$

for each vertex u ,
 $d[u] < f[u]$.
 vertex u is WHITE before time $d[u]$, GRAY between time $d[u]$ and time $f[u]$ and BLACK thereafter.

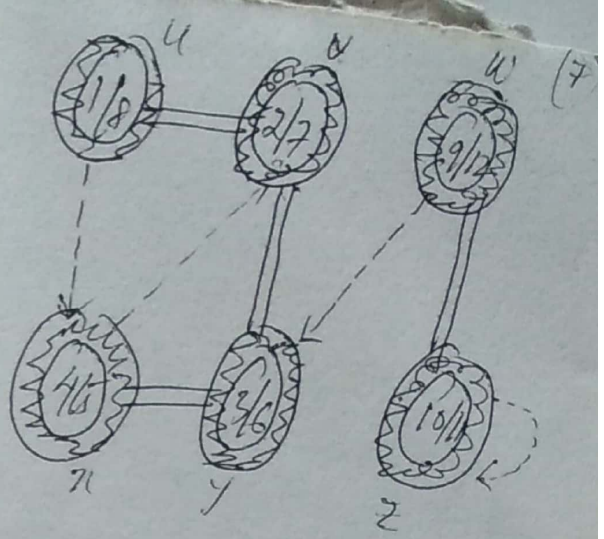
DFS-VISIT(u)

1. $color[u] \leftarrow GRAY$ // white vertex u has just been discovered
2. $time \leftarrow time + 1$
3. $d[u] \leftarrow time$
4. for each $v \in Adj[u]$ // explore edge (u, v)
5. do if $color[v] = WHITE$
6. then $\pi[v] \leftarrow u$
7. $DFS-VISIT(v)$
8. $color[u] \leftarrow BLACK$ // blacken u ; it is finished.

→ Every time DFS-VISIT (w) is called in line 7, vertex w becomes the root of a new tree in the DF forest.
 → When DFS returns, every vertex u has been assigned a discovery time $d[u]$ and a finishing time $f[u]$.



- Nontree edges are classified B, C or F according to whether they are back, cross or forward edges.
- Vertices are timestamped by discovery time / finishing time.



Time complexity

- Loops on lines 1-3 and lines 5-7 of DFS take time $\Theta(V)$
- DFS-VISIT called exactly once for each vertex $v \in V$, since DFS-VISIT is invoked only on white vertices and the first thing it does is paint the vertex gray.

→ During an execution of DFS-VISIT(v), the loop on lines 4-7 is executed $|Adj[v]|$ times since.

$$\sum_{v \in V} |Adj[v]| = \Theta(E)$$

The total cost of executing lines 4-7 of DFS-VISIT is $\Theta(E)$.
 The running time of DFS is therefore $\Theta(V+E)$

1 Tractable and Intractable Problems

So far, almost all of the problems that we have studied have had complexities that are *polynomial*, i.e. whose running time $T(n)$ has been $O(n^k)$ for some fixed value of k . Typically k has been small, 3 or less. We will let \mathbf{P} denote the class of all problems whose solution can be computed in polynomial time, i.e. $O(n^k)$ for some fixed k , whether it is 3, 100, or something else. We consider all such problems efficiently solvable, or *tractable*. Notice that this is a very relaxed definition of tractability, but our goal in this lecture and the next few ones is to understand which problems are *intractable*, a notion that we formalize as *not being solvable in polynomial time*. Notice how *not being in \mathbf{P}* is certainly a strong way of being intractable.

We will focus on a class of problems, called the *NP-complete problems*, which is a class of very diverse problems, that share the following properties: we only know how to solve those problems in time much larger than polynomial, namely *exponential time*, that is $2^{O(n^k)}$ for some k ; and if we could solve one NP-complete problem in polynomial time, then there is a way to solve *every* NP-complete problem in polynomial time.

There are two reasons to study NP-complete problems. The practical one is that if you recognize that your problem is NP-complete, then you have three choices:

1. you can use a known algorithm for it, and accept that it will take a long time to solve if n is large;
2. you can settle for *approximating* the solution, e.g. finding a nearly best solution rather than the optimum; or

3. you can change your problem formulation so that it is in P rather than being NP-complete, for example by restricting to work only on a subset of simpler problems.

Most of this material will concentrate on recognizing NP-complete problems (of which there are a large number, and which are often only slightly different from other, familiar, problems in P) and on some basic techniques that allow to solve some NP-complete problems in an *approximate* way in polynomial time (whereas an exact solution seems to require exponential time).

The other reason to study NP-completeness is that one of the most famous open problem in computer science concerns it. We stated above that “we *only know* how to solve NP-complete problems in time much larger than polynomial” not that we *have proven* that NP-complete problems require exponential time. Indeed, this is the million dollar question,¹ one of the most famous open problems in computer science, the question whether “ $P = NP$?”, or whether the class of NP-complete problems have polynomial time solutions. After

¹This is not a figure of speech. See <http://www.claymath.org/prizeproblems>

decades of research, everyone believes that $P \neq NP$, i.e. that no polynomial-time solutions for these very hard problems exist. But no one has proven it. If you do, you will be very famous, and moderately wealthy.

So far we have not actually defined what NP-complete problems are. This will take some time to do carefully, but we can sketch it here. First we define the larger class of problems called NP: these are the problems where, if someone hands you a potential solution, then you can *check* whether it is a solution in polynomial time. For example, suppose the problem is to answer the question “Does a graph have a simple path of length $|V|$?”. If someone hands you a path, i.e. a sequence of vertices, and you can *check* whether this sequence of vertices is indeed a path and that it contains all vertices in polynomial time, then the problem is in NP. It should be intuitive that any problem in P is also in NP, because we are all familiar with the fact that checking the validity of a solution is easier than coming up with a solution. For example, it is easier to get jokes than to be a comedian, it is easier to have average taste in books than to write a best-seller, it is easier to read a textbook in a math or theory course than to come up with the proofs of all the theorems by yourself. For all these reasons (and more technical ones) people believe that $P \neq NP$, although nobody has any clue how to prove it. (But once it will be proved, it will probably not be too hard to understand the proof.)

The NP-complete problems have the interesting property that if you can solve any one of them in polynomial time, then you can solve *every* problem in NP in polynomial time. In other words, they are at least as hard as any other problem in NP; this is why they are called *complete*. Thus, if you could show that *any one* of the NP-complete problems that we will study *cannot* be solved in polynomial time, then you will have not only shown that $P \neq NP$, but also that none of the NP-complete problems can be solved in polynomial time.

Conversely, if you find a polynomial-time algorithm for just one NP-complete problem, you will have shown that $P=NP$.²

2 Decision Problems

To simplify the discussion, we will consider only problems with Yes-No answers, rather than more complicated answers. For example, consider the *Traveling Salesman Problem* (TSP) on a graph with nonnegative integer edge weights. There are two similar ways to state it:

1. Given a weighted graph, what is the minimum length cycle that visits each node exactly once? (If no such cycle exists, the minimum length is defined to be ∞ .)
2. Given a weighted graph and an integer K , is there a cycle that visits each node exactly once, with weight at most K ?

Question 1 above seems more general than Question 2, because if you could answer Question 1 and find the minimum length cycle, you could just compare its length to K to answer Question 2. But Question 2 has a Yes/No answer, and so will be easier for us to consider. In

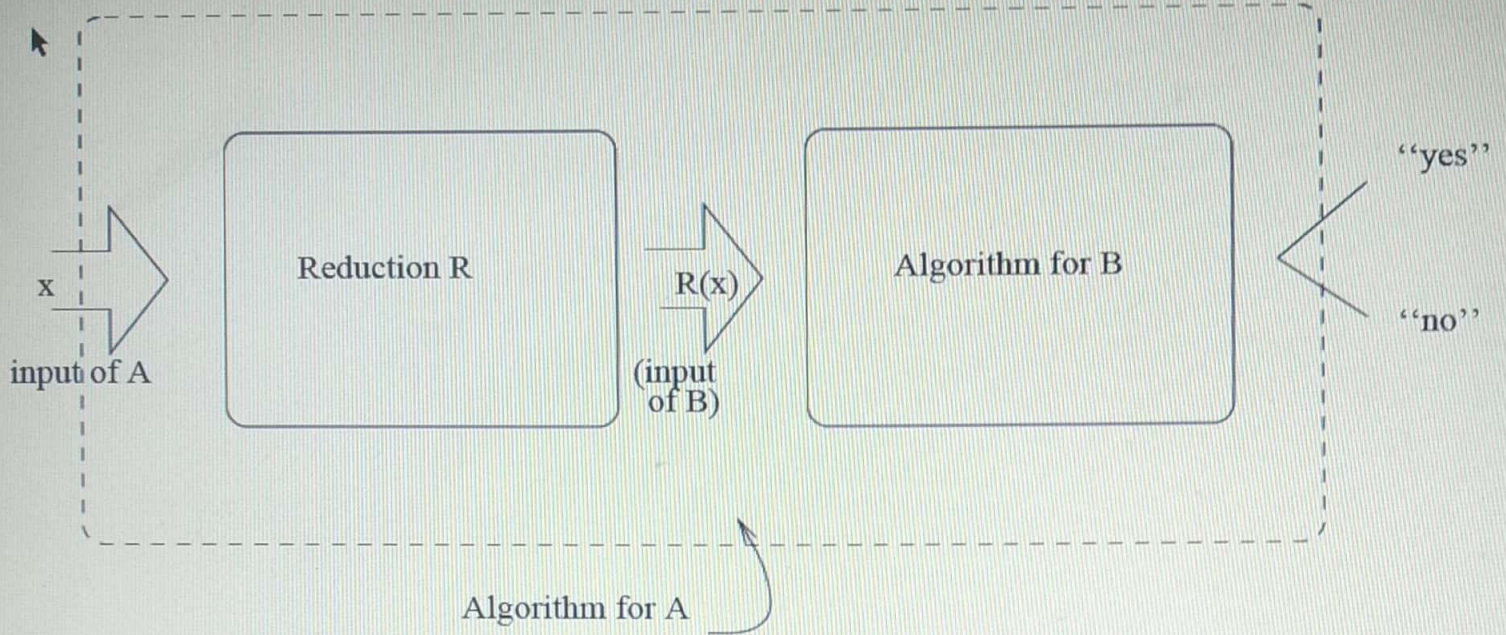


Figure 1: A reduction.

particular, if we show that Question 2 is NP-complete (it is), then that means that Question 1 is at least as hard, which will be good enough for us.³

Another example of a problem with a Yes-No answer is *circuit satisfiability* (which we abbreviate CSAT). Suppose we are given a Boolean circuit with n Boolean inputs x_1, \dots, x_n connected by AND, OR and NOT gates to one output x_{out} . Then we can ask whether there is a set of inputs (a way to assign True or False to each x_i) such that $x_{out} = \text{True}$. In

there is a set of inputs (a way to assign True or False to each x_i) such that $x_{out} = \text{True}$. In particular, we will not ask what the values of x_i are that make x_{out} True.

If A is a Yes-No problem (also called a *decision* problem), then for an input x we denote by $A(x)$ the right Yes-No answer.

3 Reductions

Let A and B be two problems whose instances require Yes/No answers, such as TSP and CSAT. A *reduction* from A to B is a polynomial-time algorithm R which transforms inputs of A to equivalent inputs of B . That is, given an input x to problem A , R will produce an input $R(x)$ to problem B , such that x is a “yes” input of A if and only if $R(x)$ is a “yes” input of B . In a compact notation, if R is a reduction from A to B , then for every input x we have $A(x) = B(R(x))$.

A reduction from A to B , together with a polynomial time algorithm for B , constitute a polynomial algorithm for A (see Figure 1). For any input x of A of size n , the reduction R takes time $p(n)$ —a polynomial—to produce an equivalent input $R(x)$ of B . Now, this input $R(x)$ can have size at most $p(n)$ —since this is the largest input R can conceivably construct in $p(n)$ time. If we now submit this input to the assumed algorithm for B , running in time $q(m)$ on inputs of size m , where q is another polynomial, then we get the right answer of x , within a total number of steps at most $p(n) + q(p(n))$ —also a polynomial!

I We have seen many reductions so far, establishing that problems are easy (e.g., from matching to max-flow). In this part of the class we shall use reductions in a more sophisticated and counterintuitive context, in order to prove that certain problems are hard. If we reduce A to B , we are essentially establishing that, *give or take a polynomial*, A is no harder than B . We could write this as

$$A \leq B$$

an inequality between the complexities of the two problems. If we know B is easy, this establishes that A is easy. If we know A is hard, this establishes B is hard. It is this latter implication that we shall be using soon.

4 Definition of Some Problems

Before giving the formal definition of NP and of NP-complete problem, we define some problems that are NP-complete, to get a sense of their diversity, and of their similarity to some polynomial time solvable problems.

In fact, we will look at pairs of very similar problems, where in each pair a problem is solvable in polynomial time, and the other is presumably not.

- minimum spanning tree: Given a weighted graph and an integer K , is there a tree that connects all nodes of the graph whose total weight is K or less?
- travelling salesman problem: Given a weighted graph and an integer K , is there a cycle that visits all nodes of the graph whose total weight is K or less?

Notice that we have converted each one of these familiar problems into a decision problem, a “yes-no” question, by supplying a goal K and asking if the goal can be met. Any optimization problem can be so converted

If we can solve the optimization problem, we can certainly solve the decision version (actually, the converse is in general also true). Therefore, proving a negative complexity result about the decision problem (for example, proving that it cannot be solved in polynomial time) immediately implies the same negative result for the optimization problem.

By considering the decision versions, we can study optimization problems side-by-side with decision problems (see the next examples). This is a great convenience in the theory of complexity which we are about to develop.

- **Eulerian graph:** Given a directed graph, is there a closed path that visits each edge of the graph exactly once?
- **Hamiltonian graph:** Given a directed graph, is there a closed path that visits each *node* of the graph exactly once?

A graph is Eulerian if and only if it is strongly connected and each node has equal in-degree and out-degree; so the problem is squarely in \mathbf{P} . There is no known such characterization—or algorithm—for the Hamilton problem (and notice its similarity with the TSP).

- **circuit value:** Given a Boolean circuit, and its inputs, is the output T?
- **circuit SAT:** Given a Boolean circuit, is there a way to set the inputs so that the output is T? (Equivalently: If we are given *some* of its inputs, is there a way to set the remaining inputs so that the output is T.)

We know that **circuit value** is in **P**: also, the naïve algorithm for that evaluates all gates bottom-up is polynomial. How about **circuit SAT**? There is no obvious way to solve this problem, sort of trying all input combinations for the unset inputs—and this is an exponential algorithm.

General circuits connected in arbitrary ways are hard to reason about, so we will consider them in a certain standard form, called *conjunctive normal form (CNF)*: Let x_1, \dots, x_n be the input Boolean variables, and x_{out} be the output Boolean variable. Then a Boolean expression for x_{out} in terms of x_1, \dots, x_n is in CNF if it is the AND of a set of *clauses*, each of which is the OR of some subset of the set of *literals* $\{x_1, \dots, x_n, \neg x_1, \dots, \neg x_n\}$. (Recall that “conjunction” means “and”, whence the name CNF.) For example,

$$x_{out} = (x_1 \vee \neg x_1 \vee x_2) \wedge (x_3 \vee x_2 \vee \neg x_1) \wedge (x_1 \vee x_2) \wedge (x_3)$$

is in CNF. This can be translated into a circuit straightforwardly, with one gate per logical operation. Furthermore, we say that an expression is in **2-CNF** if each clause has two distinct literals. Thus the above expression is not 2-CNF but the following one is:

$$(x_1 \vee \neg x_1) \wedge (x_3 \vee x_2) \wedge (x_1 \vee x_2)$$

3-CNF is defined similarly, but with 3 distinct literals per clause:

$$(x_1 \vee \neg x_1 \vee x_4) \wedge (x_3 \vee x_2 \vee x_1) \wedge (x_1 \vee x_2 \vee \neg x_3)$$

- **2SAT:** Given a Boolean formula in **2-CNF**, is there a satisfying truth assignment to the input variables?
- **3SAT:** Given a Boolean formula in **3-CNF** is there a satisfying truth assignment to the input variables?

2SAT can be solved by graph-theoretic techniques in polynomial time. For **3SAT**, no such techniques are available, and the best algorithms known for this problems are exponential in the worst case, and they run in time roughly $(1.4)^n$, where n is the number of variables. (Already a non-trivial improvement over 2^n , which is the time needed to check all possible assignments of values to the variables.)

- **matching:** Given a boys-girls compatibility graph, is there a complete matching?
- **3D matching:** Given a boys-girls-homes compatibility relation (that is, a set of boy-girl-home “triangles”), is there a complete matching (a set of disjoint triangles that covers all boys, all girls, and all homes)?

We know that matching can be solved by a reduction to max-flow. For **3D matching** there is a reduction too. Unfortunately, the reduction is *from 3SAT to 3D matching*—and this is bad news for **3D matching**...

- **unary knapsack:** Given integers a_1, \dots, a_n , and another integer K in unary, is there a subset of these integers that sum exactly to K ?
- **knapsack:** Given integers a_1, \dots, a_n , and another integer K in binary, is there a subset of these integers that sum exactly to K ?

unary knapsack is in **P**—simply because the input is represented so wastefully, with about $n + K$ bits, so that a $O(n^2K)$ dynamic programming algorithm, which would be exponential *in the length of the input* if K were represented in binary, is bounded by a polynomial in the length of the input. There is no polynomial algorithm known for the real knapsack problem. This illustrates that you have to represent your input in a sensible way, binary instead of unary, to draw meaningful conclusions.

5 NP, NP-completeness

Intuitively, a problem is in **NP** if it can be formulated as the problem of whether there is a solution

- They are *small*. In each case the solution would never have to be longer than a polynomial in the length of the input.
- They are *easily checkable*. In each case there is a polynomial algorithm which takes as inputs the input of the problem and the alleged solution, and checks whether the solution is a valid one for this input. In the case of **3SAT**, the algorithm would just check that the truth assignment indeed satisfies all clauses. In the case of *Hamilton*

cycle whether the given closed path indeed visits every node once. And so on.

- Every “yes” input to the problem has at least one solution (possibly many), and each “no” input has none.

Not all decision problems have such certificates. Consider, for example, the problem **non-Hamiltonian graph**: Given a graph G , is it true that there is no Hamilton cycle in G ? How would you prove to a suspicious person that a given large, dense, complex graph has *no* Hamilton cycle? Short of listing all cycles and pointing out that none visits all nodes once (a certificate that is certainly not succinct)?

These are examples of problems in NP:

- Given a graph G and an integer k , is there a simple path of length at least k in G ?
- Given a set of integers a_1, \dots, a_n , is there a subset S of them such that $\sum_{a \in S} a = \sum_{a \notin S} a$?

We now come to the formal definition.

DEFINITION 1 *A problem A is NP if there exist a polynomial p and a polynomial-time algorithm $V()$ such that x is a YES-input for problem A if and only if there exists a solution y , with $\text{length}(y) \leq p(\text{length}(x))$ such that $V(x, y)$ outputs YES.*

I We also call \mathbf{P} the set of decision problems that are solvable in polynomial time. Observe every problem in \mathbf{P} is also in \mathbf{NP} .

We say that a problem A is \mathbf{NP} -hard if for every N in \mathbf{NP} , N is reducible to A , and that a problem A is \mathbf{NP} -complete if it is \mathbf{NP} -hard *and* it is contained in \mathbf{NP} . As an exercise to understand the formal definitions, you can try to prove the following simple fact, that is one of the fundamental reasons why \mathbf{NP} -completeness is interesting.

LEMMA 1

If A is \mathbf{NP} -complete, then A is in \mathbf{P} if and only if $\mathbf{P} = \mathbf{NP}$.

So now, if we are dealing with some problem A that we can prove to be \mathbf{NP} -complete, there are only two possibilities:

- A has no efficient algorithm.
- All the infinitely many problems in \mathbf{NP} , including factoring and all conceivable optimization problems are in \mathbf{P} .

If $\mathbf{P} = \mathbf{NP}$, then, given the statement of a theorem, we can find a proof in time polynomial in the number of pages that it takes to write the proof down.

If it was so easy to find proof, why do papers in mathematics journal have theorems *and* proofs, instead of just having theorems. And why theorems that had reasonably short proofs have been open questions for centuries? Why do newspapers publish solutions for crossword puzzles? If $\mathbf{P} = \mathbf{NP}$, whatever exists can be found efficiently. It is too bizarre to be true.

In conclusion, it is safe to assume $P \neq NP$, or at least that the contrary will not be proved by anybody in the next decade, and it is *really* safe to assume that the contrary will not be proved by us in the next month. So, if our short-term plan involves finding an efficient algorithm for a certain problem, and the problem turns out to be NP-hard, then we should change the plan.

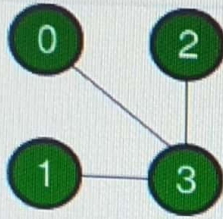
Vertex Cover Problem | Set 1 (Introduction and Approximate Algorithm)

A vertex cover of an undirected graph is a subset of its vertices such that for every edge (u, v) of the graph, either 'u' or 'v' is in vertex cover. Although the name is Vertex Cover, the set covers all edges of the given graph. **Given an undirected graph, the vertex cover problem is to find minimum size vertex cover.**

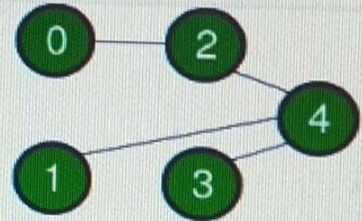
Following are some examples.



Minimum vertex cover is $\{\}$



Minimum vertex cover is $\{3\}$



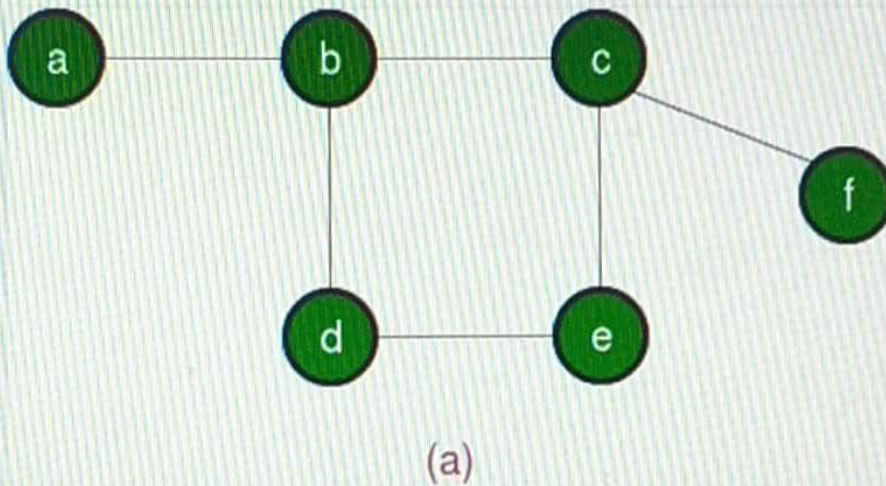
Minimum vertex cover is $\{4, 2\}$ or $\{4, 0\}$

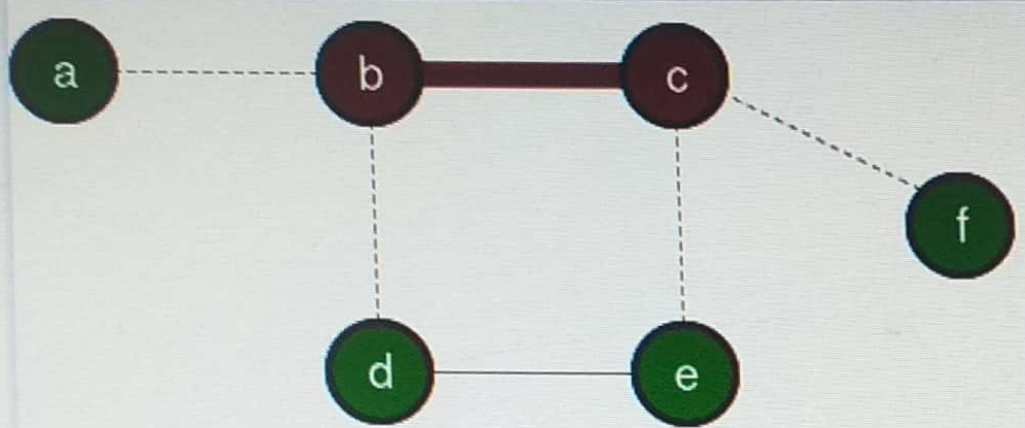
Vertex Cover Problem is a known NP Complete problem, i.e., there is no polynomial time solution for this unless $P = NP$. There are approximate polynomial time algorithms to solve the problem though. Following is a simple approximate algorithm adapted from CLRS book.

Approximate Algorithm for Vertex Cover:

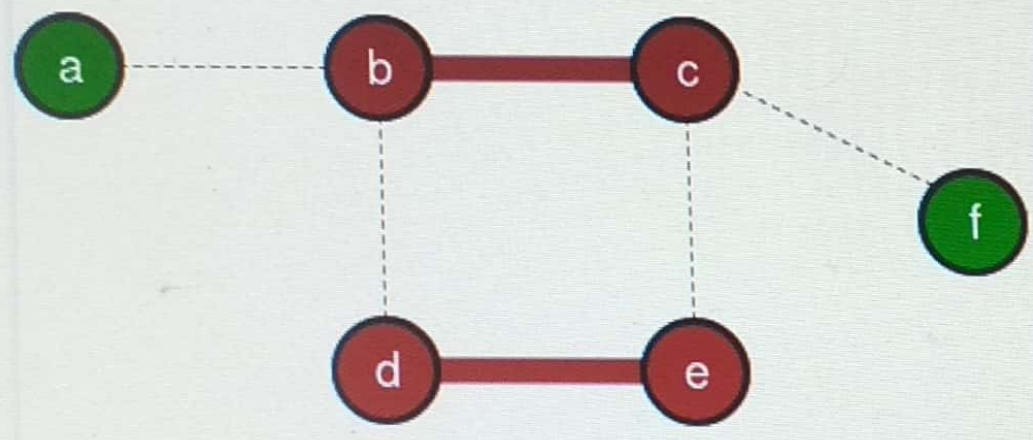
- 1) Initialize the result as $\{\}$
- 2) Consider a set of all edges in given graph. Let the set be E .
- 3) Do following while E is not empty
 - ...a) Pick an arbitrary edge (u, v) from set E and add ' u ' and ' v ' to result
 - ...b) Remove all edges from E which are either incident on u or v .
- 4) Return result

Below diagram to show execution of above approximate algorithm:





(b)



(c)

Minimum Vertex Cover is {b, c, d} or {b, c, e}

THE TRAVELING-SALESMAN PROBLEM

Given a graph $G = (V, E)$ non negative edge weight $C(e)$, and an integer c , is there a Hamiltonian cycle whose total cost is at most C ?

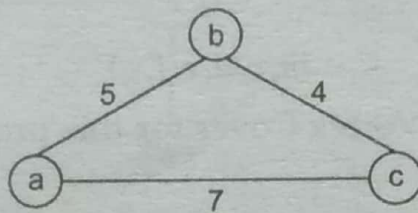
In other words we can also define TSP problem as follows : A Salesperson is required to visit a number of cities during a trip, given the distance between the cities, in what order should he travel so as to visit every city precisely once and returned home, with the minimum mileage traveled ?

Special Case of the Traveling Salesperson Problem

Here first of all we define OPT-TSP as : "Given a complete, weighted graph, find a cycle of minimum cost that visits each vertex". It is also noticeable that OPT-TSP is NP-hard.

Special Case : Edge weights satisfy the triangle inequality (which is common in many applications)

$$w(a, b) + w(b, c) \geq w(a, c)$$



Here,

⇒

⇒

⇒

⇒

$$w(a, b) = 5, \quad w(b, c) = 4 \quad \& \quad w(a, c) = 7$$

$$w(a, b) + w(b, c)$$

$$5 + 4$$

$$9 > w(a, c)$$

$$9 > 7$$

That is here the condition for triangle inequality is satisfied.

A Traveling-Salesman Problem with the Triangle Inequality

For developing an approximate algorithm for traveling salesman problem is impossible unless $P = NP$. We will first compute a minimum spanning tree we then crawl around

each node following the edges and finally returns back to the first vertex. This is a preorder walk then edges are joined in this order to obtain approx tour walk. Cost of this tour is no more than twice that of MST Cost as long as cost function satisfies the triangle inequality.

TSP-APPROX (G)

Input : Weighted complete graph G satisfying triangle inequality.

Output : ATSP tour T for G.

1. $M \leftarrow$ a minimum spanning tree for G
2. $P \leftarrow$ an Euler tour traversal of M, starting at some vertex s.
3. $T \leftarrow$ Empty list
4. for each vertex v in P (in traversal order)
5. if this is V 's first appearance in P
6. then T. insert Last (V)
7. T. insert Last (S)
8. return T

Example : Show the operation of TSP-APPROX (G) for the given set of points.

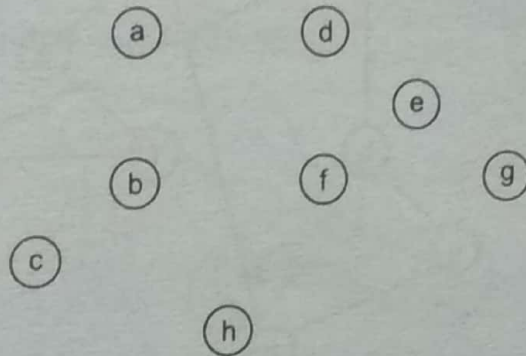


Fig 29.1 [Given set of Points]

Solution :

Step 1 : Find a minimum spanning tree for the given set of points, as :

$M \leftarrow$ a minimum spanning tree of G.

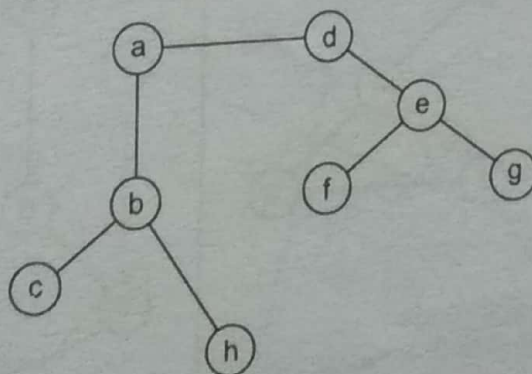


Fig 29.2 [Minimum Spanning Tree]

Step 2 : Find ordered list of vertices in "preorder walk" of M (MST) by using line 2 of TSP-APPROX (G) as :

$P \leftarrow$ an Euler tour traversal of M, starting at some vertex s.

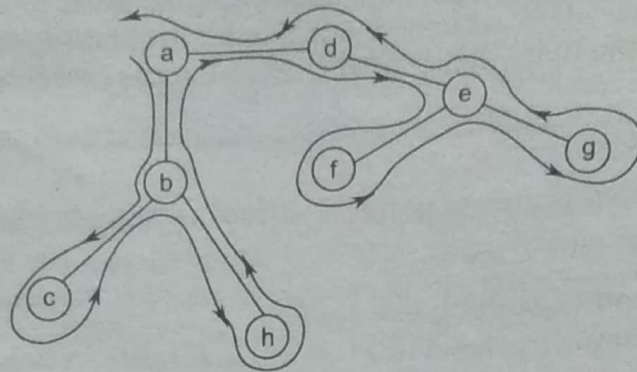
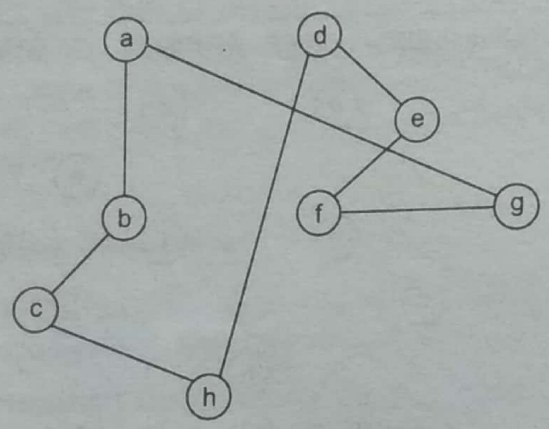


Fig 29.3 [Preorder Traversal Full walk]

Therefore preorder traversal full walk P is :

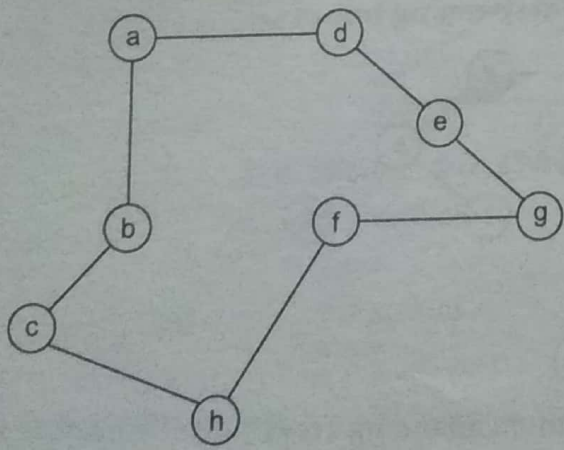
$a b c b h b a d e f e g e d a$

Step 3 : Find the cycle that visits the vertices in the order L.

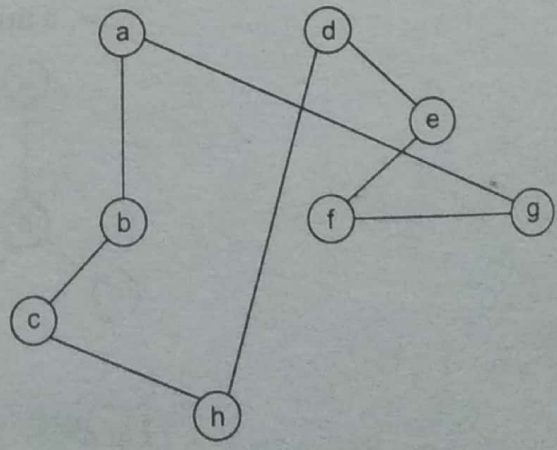


Therefore, Hamiltonian cycle is : $a b c h d e f g a$

Step 4 : Here we will show optimal tour in comparison to Hamiltonian cycle.



An optimal tour : 14.715



Hamiltonian cycle H : 19.074

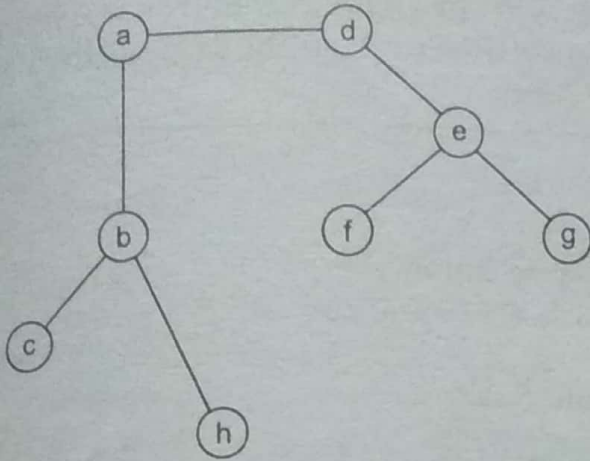
Here in both of the Cases we are assuming Euclidean distance.

Theorem 29.1 : Approx-TSP is a 2-approximation algorithm for Δ -TSP.

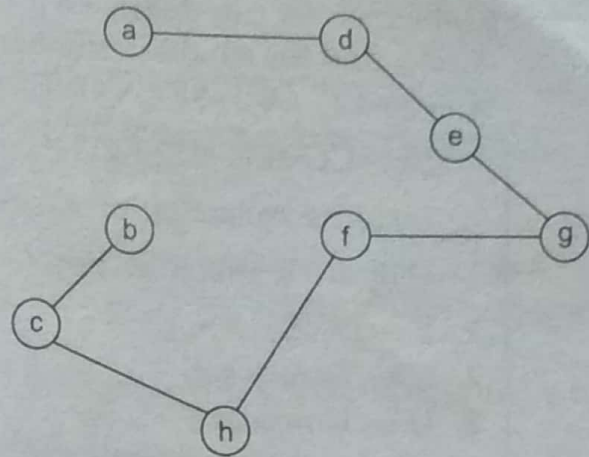
Proof : Let H^* denote an optimal tour. We want to show

$$C(H) \leq 2C(H^*)$$

(i) The optimal tour is a spanning tour hence $C(T) \leq C(H^*)$ since we obtain spanning tree by deleting any edge from optimal tour.



[MST T]

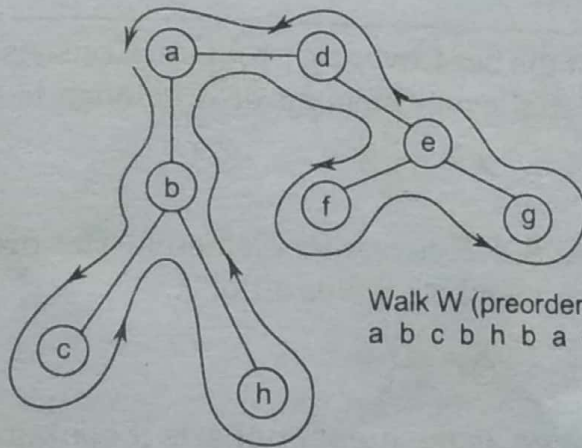


[An Optimal Tour]

(ii) The Euler tour $C(w)$ visits each edge of $C(T)$ twice hence

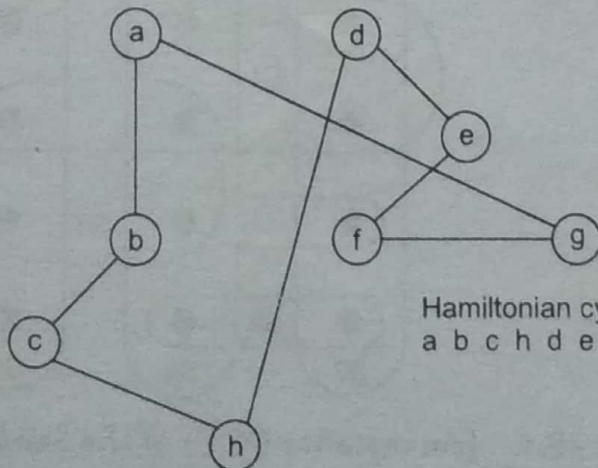
$$C(w) = 2C(T) \Rightarrow C(w) \leq 2C(H^*)$$

Since every edge visited exactly twice.



Walk W (preorder traversal)
a b c b h b a d e f e g e d a

(iii) Each time we shortcut a vertex in the Euler Tour we will not increase the total length, by the triangle inequality ($w(a, b) + w(b, c) \geq w(a, c)$) and hence $C(H) \leq C(w)$.



Hamiltonian cycle H
a b c h d e f g a

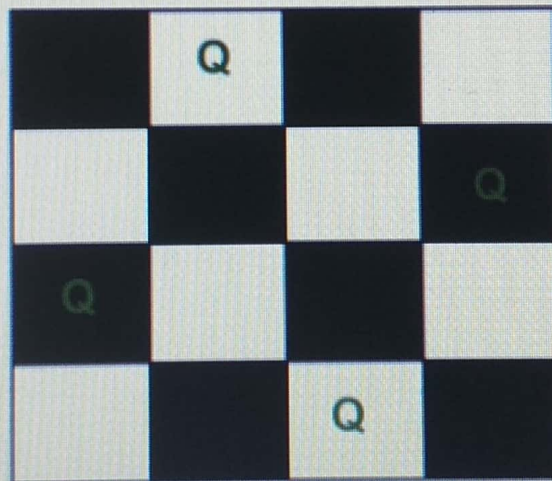
Therefore, from (ii) and (iii) points given above we can write :

$$C(H) \leq 2C(H^*)$$

N Queen Problem | Backtracking-3

We have discussed Knight's tour and Rat in a Maze problems in Set 1 and Set 2 respectively. Let us discuss N Queen as another example problem that can be solved using Backtracking.

The N Queen is the problem of placing N chess queens on an $N \times N$ chessboard so that no two queens attack each other. For example, following is a solution for 4 Queen problem.



QUICK SORT

Quicksort is a sorting algorithm whose worst case running time is $O(n^2)$ & Best case & average case is $O(n \log n)$

→ It is one of the fastest sorting algorithms.

→ Like Mergesort QUICK SORT is based on DAC paradigm.

Approaches

Divide :- The array $A[p \dots r]$ is partitioned into two nonempty subarrays $A[p \dots q]$ and $A[q+1 \dots r]$ such that each element of $A[p \dots q]$ is less than or equal to $A[q+1 \dots r]$.

The index q is computed as part of this partitioning procedure. Conquer :- The two subarrays $A[p \dots q]$ & $A[q+1 \dots r]$ are sorted by recursive calls to quicksort.

Combine :- ~~The~~ ^{Since} two subarrays are sorted in place, no work is needed to combine them. The entire array $A[p \dots r]$ is now sorted.

partition takes place in quick sort in 3 methods —

(i) Find the pivot, if the pivot element is the first element then it is the worst case.

(ii) If pivot element is the middle element, then it is the best case.

(iii) If the pivot element is chosen any element of the array then it is average case.

Basic Idea.

1. Pick one element in the array which will be the pivot.

2. Make one pass through the array called the partition step rearrange the entries so that:

→ the pivot is in its proper place.

→ entries smaller than the pivot are to its left.

→ entries larger than the pivot are to its right.

3. Recursively call (apply) Quicksort to the part of the array.

Algorithm of Quick sort
QUICKSORT (A, P, r)

// P - 1st pos of the array
// r - length of the array

- 1) if $P < r$
- 2) then $q \leftarrow$ PARTITION(A, P, r)
- 3) QUICKSORT (A, P, q-1)
- 4) QUICKSORT (A, q+1, r)

PARTITION (A, P, r)

- 1) $x \leftarrow A[P]$
- 2) $i \leftarrow P-1$
- 3) $j \leftarrow r$
- 4) While TRUE
- 5) do repeat $J \leftarrow J-1$
- 6) Until $A[J] \leq x$
- 7) repeat $i \leftarrow i+1$
- 8) Until $A[i] > x$
- 9) if $i < j$
- 10) Then exchange $A[i] \leftrightarrow A[j]$
- 11) else return j .

// x - pivot

Example $\Theta = n \log n$ Avg. Case.
 $O(n^2)$ = worst case
 $O(n^2) =$

Appl^y - Language Library

Ex: - [7 2 1 6 8 5 3 6]

[2 | 1 | 3 | 4 | 8 | 5 | 7 | 6]

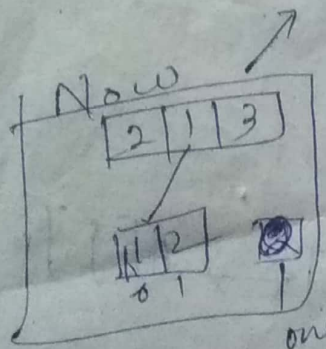
less than pivot left to the pivot
 & greater " right to the pivot

so we have to sort left to
 the pivot & sort right to the pivot
 partitioned 2 subproblem sublist part.

[2 | 1 | 3]
 0 1 2
 ↑
 pivot

[8 | 5 | 7 | 6]
 4 5 6 7

[2 | 1]
 ↑
 pivot



one element
 stop recursion

Performance of Quicksort

The running time of quick sort depends on whether partitioning is balanced or unbalanced and this in turn depends on which elements are used for partitioning.

→ If partitioning is balanced, the algorithm runs as fast as merge sort.

→ If unbalanced it can run as slow.

Best Case Partitioning

Equal size of subpartitioning is the best case Q.A.L. $n/2$ $n/2$

The recurrence for the running time

$$T(n) \leq 2T(n/2) + O(n)$$

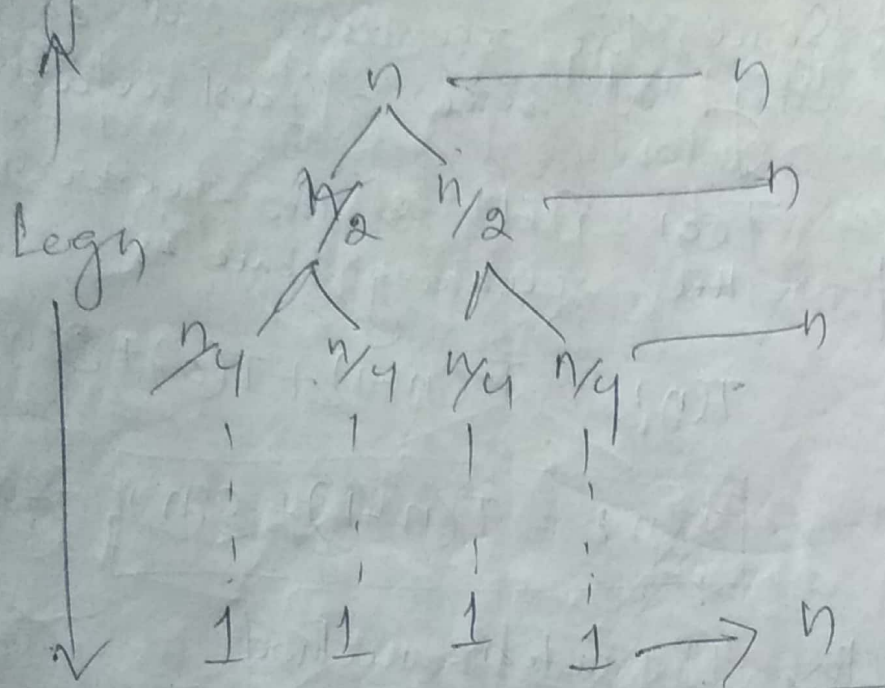
By using master theorem Case-2

it has the solⁿ

$$T(n) = O(n \log n)$$

Thus the equal balancing of the two sides of the partitions at

every level of the recursion produces an asymptotically faster algorithm.



total = $O(n \log n)$

Thus resulting running time is $O(n \log n)$

Worst-case partitioning

The worst-case behavior for QUICK SORT occurs when the partitioning routine produces one subproblem with $n-1$ elements and one with 0 elements.

→ let us assume that this unbalanced partitioning arises in each recursive call.

→ The partitioning costs $O(n)$ time