Diskrete Mathematik für Informatik (SS 2025)

Martin Held

FB Informatik Universität Salzburg A-5020 Salzburg, Austria held@cs.sbg.ac.at

12. Juni 2025



Personalia

LVA-Leiter (VO+PS): Martin Held. Email-Adresse: held@cs.sbg.ac.at. Basis-URL: https://www.cosy.sbg.ac.at/~held. Büro: Universität Salzburg, FB Informatik, Zi. 1.20, Jakob-Haringer Str. 2, 5020 Salzburg-Itzling. Telefonnummer (Büro): (0662) 8044-6304. Telefonnummer (Sekr.): (0662) 8044-6300.





Personalia

LVA-Leiter (PS): Markus Flatz. Email-Adresse: mflatz@cs.sbg.ac.at. Telefonnummer (Sekr.): (0662) 8044-6300.



Personalia

LVA-Leiter (PS): Mara Grilnberger. Email-Adresse: mara.grilnberger@plus.ac.at. Büro: Universität Salzburg, FB Informatik, Zi. 2.34, Jakob-Haringer Str. 2, 5020 Salzburg-Itzling.

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Formalia

LVA-URL (VO+PS): https://www.cosy.sbg.ac.at/~held/teaching/diskrete_ mathematik/dm.html.

Allg. Information: Basis-URL/for_students.html.

PLUSonline: Bitte melden Sie sich unbedingt im PLUSonline zu VO/PS an!

Abhaltezeit der VO: Donnerstag 7⁴⁵–11⁰⁰, mit etwa 20–25 Minuten Pause.

Abhalteort der VO: T01, FB Informatik, Jakob-Haringer Str. 2.

Abhaltezeit des PS: Freitag 11⁴⁰–13⁴⁰.

Abhalteort des PS: T01+T02+T03, Jakob-Haringer Str. 2.

Tutorium: Andreas Auer und Jatin Kumar: Montag 16⁰⁰–18⁰⁰ (T06), Mittwoch 12³⁰–14³⁰ (T02); FB Informatik, Jakob-Haringer Str. 2.

Achtung — das Proseminar ist prüfungsimmanent!



Electronic Slides and Online Material

In addition to these slides, you are encouraged to consult the WWW home page of this lecture:

https://www.cosy.sbg.ac.at/~held/teaching/diskrete_mathematik/dm.html.

In particular, this WWW page contains up-to-date information on the course, plus links to online notes, slides and (possibly) sample code.





A Few Words of Warning

I hope that these slides will serve as a practice-minded introduction to various aspects of discrete mathematics which are of importance for computer science. I would like to warn you explicitly not to regard these slides as the sole source of information on the topics of my course. It may and will happen that I'll use the lecture for talking about subtle details that need not be covered in these slides! In particular, the slides won't contain all sample calculations, proofs of theorems, demonstrations of algorithms, or solutions to problems posed during my lecture. That is, by making these slides available to you I do not intend to encourage you to attend the lecture on an irregular basis.



Acknowledgments

These slides are a revised and extended version of a draft prepared by Kamran Safdar. Included is material written by Christian Alt, Caroline Atzl, Michael Burian, Peter Gintner, Bernhard Guillon, Yvonne Höller, Stefan Huber, Sandra Huemer, Christian Lercher, Sebastian Stenger, Alexander Zrinyi. I also benefited from comments and suggestions made by Stefan Huber and Peter Palfrader.

This revision and extension was carried out by myself, and I am responsible for all errors.

Salzburg, February 2025

Martin Held



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K.H. Rosen.

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📐 K.A. Ross, C.R.B. Wright. Discrete Mathematics Pearson Prentice Hall, 5th edition, Aug 2002; ISBN 9780130652478

🍆 C. Stein, R.L.S. Drysdale, K. Bogart. Discrete Mathematics for Computer Science Addison-Wesley, March 2010; ISBN 978-0132122719.





📎 J. O'Donnell, C. Hall, R. Page.

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M. Smid.

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Discrete Structures for Computer Science: Counting, Recursion, and Probability http://cglab.ca/~michiel/DiscreteStructures, 2019



📎 E. Lehman, F.T. Leighton, A.R. Meyer. Mathematics for Computer Science

https://courses.csail.mit.edu/6.042,2018



M M Eleck Building Blocks for Theoretical Computer Science http://mfleck.cs.illinois.edu/building-blocks/, 2017



Table of Content



- Propositional and Predicate Logic
- 3 Definitions and Theorem Proving
- Mumbers and Basics of Number Theory
- Principles of Elementary Counting and Combinatorics
- 6 Complexity Analysis and Recurrence Relations
- Graph Theory
- Cryptography





Introduction

- What is Discrete Mathematics?
- Motivation





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 - Elementary probability theory.



Applications of Discrete Mathematics

• DM forms the mathematical language of computer science. It is at the very heart of several other parts of computer science.



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- Applications of DM include but are not limited to
 - Algorithms and data structures,
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- We start with a set of sample problems; solutions for all problems will be worked out or, at least, sketched during this course.





Motivation



- Suppose that an algorithm needs $1 + 2 + 3 + \cdots + (n 1) + n$ many computational steps (of unit cost) to handle an input of size *n*.
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| 1 | = | 1 |
|---------------------------|---|----|
| 1 + 2 | = | 3 |
| 1 + 2 + 3 | = | 6 |
| 1 + 2 + 3 + 4 | = | 10 |
| 1 + 2 + 3 + 4 + 5 | = | 15 |
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• An inspection of the numbers on the right-hand side might let us suspect that

$$1+2+3+\cdots+(n-1)+n=\frac{n(n+1)}{2}.$$

 But is this indeed correct? And, by the way, what do the dots in this equation really mean??

 An answer can be established by means of number theory (natural numbers, induction). And we get indeed

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- Caution: Even after calculating this sum for all values of *n* between 1 and 500 one can not legitimately claim to know the sum for, say, *n* := 1000.
- Note: It would constitute a horrendous waste of CPU time to let a computer compute $1 + 2 + 3 + \cdots + (n 1) + n$ by successively adding numbers if we could simply obtain the result by evaluating $\frac{n(n+1)}{2}$.



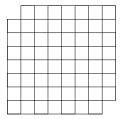
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500

• Consider an 8 × 8 chessboard

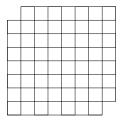


• Consider an 8 \times 8 chessboard with the upper-left and lower-right cells removed,





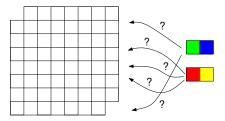
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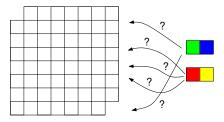


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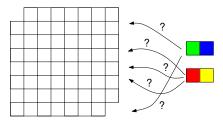
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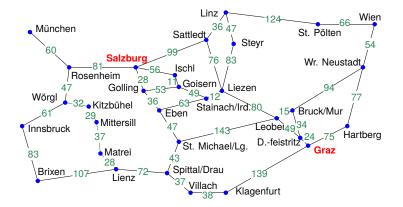


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- Caution: Simply trying out *all* possible placements of domino blocks hardly is an option for an 8 × 8 chessboard and definitely no option for an *n* × *n* board!



Sample Problem: Route Calculation

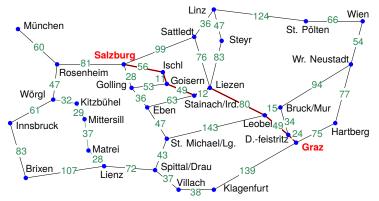
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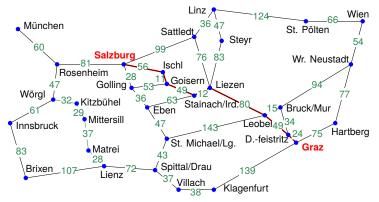
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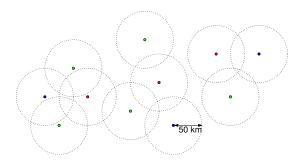


• Note: Simply trying all possible routes gets tedious! (How would you even guarantee that all possible routes have indeed been checked?)



Sample Problem: Channel Assignment

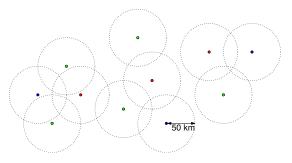
• Suppose that frequencies out of a set of *m* frequencies are to be assigned to *n* broadcast stations within Austria. We are told that the area serviced by a station lies within a disk with radius 50 kilometers. Obviously, no two different stations whose broadcast areas overlap may use the same frequency.





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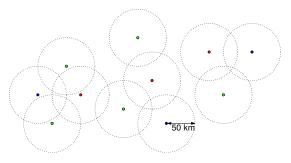
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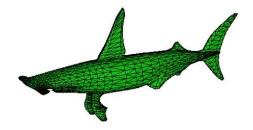
 The solution can be obtained by using techniques of computational geometry combined with graph coloring.

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Diskrete Mathematik (SS 2025)

Sample Problem: Memory Required for Storing a Polyhedron

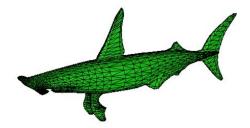
• Suppose that a polyhedral model has *n* vertices. How many edges and faces can it have at most? What is the storage complexity relative to *n*?





Sample Problem: Memory Required for Storing a Polyhedron

• Suppose that a polyhedral model has *n* vertices. How many edges and faces can it have at most? What is the storage complexity relative to *n*?



• Answer provided by graph theory: A polyhedron with *n* vertices has at most 3n - 6 edges and 2n - 4 faces.





input:

| | 100 |
|--|-----|
| | |



| input: | 100 |
|----------------|-----|
| after round 1: | 75 |



| input: | 100 |
|----------------|-----|
| after round 1: | 75 |
| after round 2: | 56 |



| input: | 100 |
|----------------|-----|
| after round 1: | 75 |
| after round 2: | 56 |
| after round 3: | 42 |



• Suppose that an algorithm is given *n* numbers as input and that it solves a problem by proceeding as follows: During one round of computation, it performs *n* computational steps. We know that during each round it discards at least 25% of the numbers. The algorithm executes one round after the other until only one number is left.

| input: | 100 |
|----------------|-----|
| after round 1: | 75 |
| after round 2: | 56 |
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• Question: How many rounds does the algorithm run in the worst case (depending on the input size *n*)? How many computational steps are carried out in the worst case?

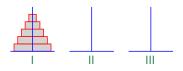


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- Question: How many rounds does the algorithm run in the worst case (depending on the input size *n*)? How many computational steps are carried out in the worst case?
- Answer provided by the theory of recurrence relations: The number of computational steps is linear in *n*, and the number of rounds is logarithmic in *n*.
- In asymptotic notation: O(n) and $O(\log n)$.

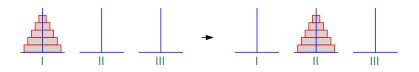


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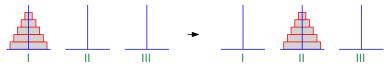


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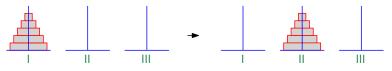


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One can also prove: Every(!) algorithm that solves ToH needs at least 2ⁿ - 1 moves.



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- Thus, the solution achieved by the recursive algorithm is optimal as far as the number of moves is concerned.





- Tower-of-Hanoi Problem (ToH): Given three pegs (labeled I,II,III) and a stack of *n* disks arranged on Peg I from largest at the bottom to smallest at the top, we are to move all disks to Peg II such that only one disk is moved at a time and such that no larger disk ever is placed on a smaller disk.
- Attributed to Édouard Lucas (1883). Supposedly based on an Indian legend about Brahmin priests moving 64 disks in the Great Temple of Benares; once they are finished, life on Earth will end.
- Goal: Find an algorithm that uses the minimum number of moves.



- One can prove: A (straightforward) recursive algorithm needs $2^n 1$ moves.
- One can also prove: Every(!) algorithm that solves ToH needs at least 2ⁿ 1 moves.
- Thus, the solution achieved by the recursive algorithm is optimal as far as the number of moves is concerned.
- [Buneman&Levy (1980)]: There exists a simple iterative solution that avoids an exponential-sized stack!

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and, in general, using the capital-sigma notation and geometric series,

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- [Sagan 1997]: "Exponentials can't go on forever, because they will gobble up everything".
- The "second half of the chessboard" is a phrase, coined by Kurzweil in 1999 to the point where exponential growth begins to have a significant impact.

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Sample Problem: Key Distribution and Message Encryption

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- By the way, how could Alice and Bob encrypt or decrypt messages once they have exchanged their key?
- Answer: This is yet another application of number theory!



Propositional and Predicate Logic

- Propositional Logic
- Predicate Logic
- Special Quantifiers





Propositional and Predicate Logic

- Propositional Logic
- Predicate Logic
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Propositional Logic

- Goal: specification of a language for formally expressing theorems and proofs.
- Aka: propositional calculus, logic of statements, statement logic;
- Dt.: Aussagenlogik.

Definition 1 (Proposition, Dt.: Aussage)

A proposition is a statement that is either true or false.

 Propositions can be *atomic*, like "The sun is shining",

or compound,

like "The sun is shining and the temperature is high".

 In the latter case, the proposition is a composition of atomic or compound propositions by means of logical junctors. (Junctors are also known as connectives or operators.)



Language of Propositional Logic

Definition 2 (Propositional formula, Dt.: aussagenlogische Formel)

A propositional formula is constructed inductively from a set of

- propositional variables (typically p, q, r or p₁, p₂,...);
- junctors: $\neg, \land, \lor, \Rightarrow, \Leftrightarrow;$
- parentheses: (,);
- constants (truth values): \bot , \top (or F, T);

based on the following rules:

- A propositional variable is a propositional formula.
- The constants \perp and \top are propositional formulas.
- If ϕ_1 and ϕ_2 are propositional formulas then so are the following:

 $(\neg \phi_1), (\phi_1 \land \phi_2), (\phi_1 \lor \phi_2), (\phi_1 \Rightarrow \phi_2), (\phi_1 \Leftrightarrow \phi_2).$



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Precedence Rules

 Precedence rules (Dt.: Vorrangregeln) are used frequently to avoid the burden of too many parentheses. From highest to lowest precedence, the following order is common.

 $\neg, \land, \lor, \overset{\Rightarrow}{\Leftrightarrow}$

- Unfortunately, different precedence rules tend to be used by different authors.
- Thus, make it clear which order you use, or in case of doubt, insert parentheses!
- It is common to represent the truth values of a proposition under all possible assignments to its variables by means of a *truth table*.
- In addition to the standard junctors we also define two other operators, NAND, denoted by ↑ (or sometimes by |), and NOR, denoted by ↓.



Truth Tables

| р | q | $\neg p$ | $p \land q$ | $p \lor q$ | $p \Rightarrow q$ | $p \Leftrightarrow q$ | $p \uparrow q$ | $p \downarrow q$ |
|---|---|----------|-------------|------------|-------------------|-----------------------|----------------|------------------|
| T | T | F | T | Т | T | Т | F | F |
| Т | F | F | F | Т | F | F | Т | F |
| F | T | T | F | Т | T | F | Т | F |
| F | F | T | F | F | T | Т | Т | Т |

- Common names for the junctors in natural language:
 - ¬p: Noт, negation;
 - $p \land q$: AND, conjunction;
 - $p \lor q$: OR, disjunction;
 - $p \Rightarrow q$: IMPLIES, conditional, if p then q, q if p, p sufficient for q, q necessary for p;
 - *p* ⇔ *q*: IFF, equivalence, biconditional, *p* if and only if *q*, *p* necessary and sufficient for *q*.
- Note: The truth table (Dt.: Wahrheitstabelle) of a formula with n variables has 2ⁿ rows.



Definition 3 (Tautology, Dt.: Tautologie)

A propositional formula is a *tautology* if it is true under all truth assignments to its variables.

Definition 4 (Contradiction, Dt.: Widerspruch)

A propositional formula is a *contradiction* if it is false under all truth assignments to its variables.

- Standard examples: $(p \lor \neg p)$ and $(p \land \neg p)$.
- Easy to prove: The negation of a tautology yields a contradiction, and vice versa.



Definition 5 (Logical equivalence, Dt.: logische Äquivalenz)

Two propositional formulas are *logically equivalent* if they have the same truth table. Logical equivalence of formulas ϕ_1, ϕ_2 is commonly denoted by $\phi_1 \equiv \phi_2$.

Theorem 6

Two propositional formulas ϕ_1, ϕ_2 are logically equivalent iff $\phi_1 \Leftrightarrow \phi_2$ is a tautology.

Definition 7 (Complete set of junctors, Dt.: vollständige Junktorenmenge)

A set S of junctors is said to be *complete* (or truth-functionally adequate/complete) if, for any given propositional formula, a logically equivalent one can be written using only junctors of S.

• Note: The sets $\{\uparrow\}$ and $\{\downarrow\}$ both are complete sets of junctors.



Laws for Logical Equivalence

Theorem 8

Let ϕ_1, ϕ_2 be propositional formulas. Then the following equivalences hold: Identity: $\phi_1 \wedge T \equiv \phi_1$ $\phi_1 \vee F \equiv \phi_1$ Domination: $\phi_1 \vee T \equiv T$ $\phi_1 \wedge F \equiv F$ Idempotence: $\phi_1 \lor \phi_1 \equiv \phi_1$ $\phi_1 \wedge \phi_1 \equiv \phi_1$ Double negation: $\neg \neg \phi_1 \equiv \phi_1$ Commutativity: $\phi_1 \wedge \phi_2 \equiv \phi_2 \wedge \phi_1$ $\phi_1 \lor \phi_2 \equiv \phi_2 \lor \phi_1$ $\phi_1 \Leftrightarrow \phi_2 \equiv \phi_2 \Leftrightarrow \phi_1$ $(\phi_1 \lor \phi_2) \land \phi_3 \equiv (\phi_1 \land \phi_3) \lor (\phi_2 \land \phi_3)$ Distributivity: $(\phi_1 \land \phi_2) \lor \phi_3 \equiv (\phi_1 \lor \phi_3) \land (\phi_2 \lor \phi_3)$ Associativity: $(\phi_1 \lor \phi_2) \lor \phi_3 \equiv \phi_1 \lor (\phi_2 \lor \phi_3)$ $(\phi_1 \land \phi_2) \land \phi_3 \equiv \phi_1 \land (\phi_2 \land \phi_3)$ De Morgan's laws: $\neg(\phi_1 \land \phi_2) \equiv \neg\phi_1 \lor \neg\phi_2$ $\neg(\phi_1 \lor \phi_2) \equiv \neg\phi_1 \land \neg\phi_2$ Trivial tautology: $\phi_1 \vee \neg \phi_1 \equiv T$ Trivial contradiction: $\phi_1 \wedge \neg \phi_1 \equiv F$ $\neg \phi_1 \Leftrightarrow \neg \phi_2 \equiv \phi_1 \Leftrightarrow \phi_2 \quad \neg \phi_2 \Rightarrow \neg \phi_1 \equiv \phi_1 \Rightarrow \phi_2$ Contraposition: Implication as Disj.: $\phi_1 \Rightarrow \phi_2 \equiv \neg \phi_1 \lor \phi_2$



Definition 9 (Logical implication, Dt.: logische Implikation)

A formula ϕ_1 *logically implies* ϕ_2 , denoted by $\phi_1 \models \phi_2$, if $\phi_1 \Rightarrow \phi_2$ is a tautology.

Definition 10 (Proof, Dt.: Beweis)

A proof of ψ based on premises ϕ_1, \ldots, ϕ_n is a finite sequence of propositions that ends in ψ such that each proposition is either a premise or a logical implication of the previous proposition.

- Note: Logical implication rather than logical equivalence!
- Thus,
 - note that it need not be possible to revert a proof!
 - pay close attention to which steps are actual equivalences if you intend to argue both ways!



Rules of Inference

- Aka: proof rules (Dt.: Schlußregeln).
- In addition to the following inference rules for propositional formulas φ₁, φ₂, all the equivalence rules apply: Each equivalence can be written as two inference rules since they are valid in both directions.

$$\begin{array}{c|c} \hline \phi_1 \wedge \phi_2 \\ \hline \phi_1 \\ \hline \phi_1 \\ \hline \phi_1 \\ \hline \phi_1 \\ \hline \phi_2 \\ \hline \phi_1 \\ \hline \phi_2 \\$$



φ

Definition 11 (Satisfiability, Dt.: Erfüllbarkeit)

A formula ϕ is *satisfiable* if there exists at least one truth assignment to the variables of ϕ that makes ϕ true.

Definition 12 (Satisfiability equivalent)

Two formulas are *satisfiability equivalent* if both formulas are either satisfiable or not satisfiable.



Conjunctive Normal Form

• In mathematics, normal forms are canonical representations of objects such that all equivalent objects have the same representation.

Definition 13 (Literal, Dt.: Literal)

A *literal* is a propositional variable or the negation of a propositional variable. A *clause* is a disjunction of literals.

• E.g., if p, q are variables then p and $\neg q$ are literals, and $(p \lor \neg q)$ is a clause.

Definition 14 (Conjunctive normal form, Dt.: konjunktive Normalform)

A propositional formula is in (general) *conjunctive normal form* (CNF) if it is a conjunction of clauses.

• E.g.,
$$\neg p_1 \land (p_2 \lor p_5 \lor \neg p_6) \land (\neg p_3 \lor p_4 \lor \neg p_6)$$
 is a CNF formula.

Definition 15 (k-CNF)

A CNF formula is a k-CNF formula if every clause contains at most k literals.



Conjunctive Normal Form

- Note: Some textbooks demand *exactly k literals* rather than at most k literals.
- Note: It is common to demand that no variable may appear more than once in a clause.
- Note: For k ≥ 3, a general CNF formula can easily be converted in polynomial time (in the number of literals) into a k-CNF formula with exactly k literals per clause such that no variable appears more than once in a clause and such that the two formulas are satisfiability equivalent.





- Propositional Logic
- Predicate Logic
- Special Quantifiers



Definition 16 (*n***-ary Relation, Dt.:** *n***-stellige Relation)**

Let A_1, A_2, \ldots, A_n be sets, for some $n \in \mathbb{N}$. An *n*-ary relation \mathcal{R} on A_1, A_2, \ldots, A_n is a subset of their Cartesian product, i.e., $\mathcal{R} \subseteq A_1 \times A_2 \times \cdots \times A_n$.

Definition 17 (*n***-ary Function, Dt.:** *n***-stellige Funktion)**

Let A_1, A_2, \ldots, A_n, B be sets, for some $n \in \mathbb{N}$. An *n*-ary function \mathcal{F} from $A_1 \times A_2 \times \cdots \times A_n$ to B is an (n + 1)-ary relation on A_1, A_2, \ldots, A_n, B such that for any $(a_1, a_2, \ldots, a_n) \in A_1 \times A_2 \times \cdots \times A_n$ there exists a unique $b \in B$ such that $(a_1, a_2, \ldots, a_n, b) \in \mathcal{F}$.

- It is common to write $y = \mathcal{F}(a_1, \ldots, a_n)$ for "pick y such that $(a_1, \ldots, a_n, y) \in \mathcal{F}$ ".
- The set $A_1 \times A_2 \times \cdots \times A_n$ is called the *domain* and the set *B* is called the codomain of \mathcal{F} .
- An *n*-ary relation/function over a set *A* is a relation/function where $A_1 = A_2 = \ldots = A_n = A$, i.e., $A_1 \times A_2 \times \cdots \times A_n = A^n$. It is also called an *n*-place relation/function.
- A 1-ary relation/function is called *unary*, and a 2-ary relation/function is called *unary*.

Definition 18 (Predicate, Dt.: Prädikat)

For an *n*-ary relation \mathcal{R} over A, an *n*-ary *predicate* over A is the *n*-ary function $f_{\mathcal{R}} : A^n \to \{T, F\}$, where

$$f_{\mathcal{R}}(a_1,\ldots,a_n) := \begin{cases} T & \text{if } (a_1,\ldots,a_n) \in \mathcal{R}, \\ F & \text{otherwise.} \end{cases}$$

- Thus, a predicate is a Boolean function.
- Note: This is a slight abuse of notation since the symbols ":" and " \rightarrow " in " $f: M \rightarrow N$ " actually form already a 3-ary predicate!
- An 1-ary predicate is called *unary*, and a 2-ary predicate is called *binary*.
- A sample unary predicate on \mathbb{R} is

"*x* is non-negative" := $\begin{cases} T & \text{if } x \ge 0, \\ F & \text{otherwise.} \end{cases}$

• Dt.: Prädikatenlogik.



Language of Predicate Logic

Definition 19 (Predicate vocabulary, Dt.: Symbolmenge)

A predicate vocabulary consists of

- a set C of constant symbols,
- a set \mathcal{F} of function symbols,
- a set \mathcal{V} of variables, typically $\{x_1, x_2, \ldots\}$ or $\{a, b, \ldots\}$,
- a set *P* of predicate symbols, including the 0-ary predicate symbols (truth values) ⊥, ⊤ or *F*, *T*,

together with

- logical junctors $\neg, \land, \lor, \Rightarrow, \Leftrightarrow$,
- quantifiers ∃, ∀,
- parentheses.



Definition 20 (Term)

A term over $(\mathcal{C},\mathcal{V},\mathcal{F})$ is defined inductively as follows:

- Every constant $c \in C$ is a term.
- Every variable $x \in \mathcal{V}$ is a term.
- If t_1, \ldots, t_n are terms and f is an *n*-ary function symbol then $f(t_1, \ldots, t_n)$ is a term.
- Note: Constants can be thought of as 0-ary function symbols. Thus, a set C of constants need not be considered when defining the language of predicate logic.



Definition 21 (Formulas)

The set of *formulas* over $(\mathcal{C}, \mathcal{V}, \mathcal{F}, \mathcal{P})$ is defined inductively as follows:

- \perp and \top are formulas.
- If t_1, \ldots, t_n are terms and $P \in \mathcal{P}$ is an *n*-ary predicate, then $P(t_1, \ldots, t_n)$ is a (so-called *atomic*) formula.
- If ϕ and ψ are formulas then $(\neg \phi)$, $(\phi \land \psi)$, $(\phi \lor \psi)$, $(\phi \Rightarrow \psi)$ and $(\phi \Leftrightarrow \psi)$ are formulas.
- If φ is a formula then (∀x φ) and (∃x φ) are formulas. In both cases, the scope of the quantifier is given by the formula φ to which the quantifier is applied.

Definition 22 (Quantifier-free formula, Dt.: quantorenfreie Formel)

A quantifier-free formula is a formula which does not contain a quantifier.



Quantifiers

Definition 23 (Universe of discourse, Dt.: Wertebereich, Universum)

The *universe of discourse* specifies the set of values that the variable *x* may assume in $(\forall x \ \phi)$ and $(\exists x \ \phi)$.

Definition 24 (Universal quantifier, Dt.: Allquantor)

- $(\forall x \ P(x))$ is the statement
 - "P(x) is true for all x (in the universe of discourse)".

Definition 25 (Existential quantifier, Dt.: Existenzquantor)

 $(\exists x \ P(x))$ is the statement

"there exists x (in the universe of discourse) such that P(x) is true".

- The notation (∃!*x P*(*x*)) is a convenience short-hand for *"there exists exactly one x such that P*(*x*) *is true",*
 - i.e., for denoting existence and uniqueness of a suitable *x*.



Precedence Rules for Quantified Formulas

- No universally accepted precedence rule exists.
- Thus, you have to make your specific order very clear.
- Even better, use parentheses or (significant!) spaces between coherent parts of the expression.
- First-order logic versus higher-order logic: In first-order predicate logic, predicate quantifiers or function quantifiers are not permitted, and variables are the only objects that may be quantified. Also, predicates are not allowed to have predicates as arguments.



Free Variables

Definition 26 (Free variables, Dt.: freie Variable)

The *free variables* of a formula ϕ or a term *t*, denoted by $FV(\phi)$ and FV(t), are defined inductively as follows:

 $FV(c) := \{\};$ For a constant $c \in C$: For a variable $x \in \mathcal{V}$: $FV(x) := \{x\}$: $FV(f(t_1,\ldots,t_n)) := FV(t_1) \cup \ldots \cup FV(t_n);$ For a term $f(t_1, \ldots, t_n)$: For a formula $P(t_1, \ldots, t_n)$: $FV(P(t_1, \ldots, t_n)) := FV(t_1) \cup \ldots \cup FV(t_n)$; Also. $FV(\perp) := \{\}.$ $FV(\top) := \{\};$ $FV((\neg \phi)) := FV(\phi),$ For formulas ϕ and ψ : $FV((\phi \land \psi)) := FV(\phi) \cup FV(\psi),$ $FV((\phi \lor \psi)) := FV(\phi) \cup FV(\psi),$ $FV((\phi \Rightarrow \psi)) := FV(\phi) \cup FV(\psi),$ $FV((\phi \Leftrightarrow \psi)) := FV(\phi) \cup FV(\psi)$: $FV((\forall x \ \phi)) := FV(\phi) \setminus \{x\},\$ For a formula ϕ : $FV((\exists x \ \phi)) := FV(\phi) \setminus \{x\}.$

Definition 27 (Bound variables, Dt.: gebundene Variable)

The *bound variables* of a formula ϕ or a term *t*, denoted by $BV(\phi)$ and BV(t), are defined inductively as follows:

For a constant $c \in C$: $BV(c) := \{\}:$ For a variable $x \in \mathcal{V}$: $BV(x) := \{\}:$ For a term $f(t_1, ..., t_n)$: $BV(f(t_1, ..., t_n)) := \{\};$ For a formula $P(t_1, ..., t_n)$: $BV(P(t_1, ..., t_n)) := \{\};$ Also. $BV(\perp) := \{\},\$ $BV(\top) := \{\};$ For formulas ϕ and ψ : $BV((\neg \phi)) := BV(\phi),$ $BV((\phi \land \psi)) := BV(\phi) \cup BV(\psi),$ $BV((\phi \lor \psi)) := BV(\phi) \cup BV(\psi),$ $BV((\phi \Rightarrow \psi)) := BV(\phi) \cup BV(\psi).$ $BV((\phi \Leftrightarrow \psi)) := BV(\phi) \cup BV(\psi)$: $BV((\forall x \ \phi)) := BV(\phi) \cup \{x\},\$ For a formula ϕ : $BV((\exists x \ \phi)) := BV(\phi) \cup \{x\}.$

Free and Bound Variables

- Note: Technically speaking, one variable symbol may denote both a free and a bound variable of a formula!
- However, common sense dictates to use a different symbol if a different variable is meant, even if not required by the syntax of predicate logic:
 - Do not use the same symbol for bound and free variables! E.g.,

 $(P(x) \Rightarrow (\forall x \ Q(x)))$

is syntactically correct but extremely difficult to parse for a human.

• Also, do not re-use symbols of bound variables inside nested quantifiers! E.g.,

 $(\forall x \ (P(x) \Rightarrow (\forall x \ Q(x))))$

is syntactically correct but horrible to parse.

Definition 28 (Sentence, Dt.: geschlossener Ausdruck)

A formula ϕ is a *sentence* if $FV(\phi) = \{\}$.



Substitutions

Definition 29 (Substitution, Dt.: Ersetzung)

For a formula ϕ , variable *x* and term *t*, we obtain the *substitution* of *x* by *t*, denoted as $\phi[t/x]$, by replacing each free occurrence of *x* in ϕ by *t*.

Definition 30 (Valid substitution, Dt.: gültige Ersetzung)

A substitution of *t* for *x* in a formula ϕ is *valid* if and only if no free variable of *t* ends up being bound in $\phi[t/x]$.

- Not a valid substitution of x: $\phi \equiv (\exists y \in \mathbb{N} \ y > 10 \ \land \ x < y)$ and t := 2y + 5.
- Again, it is very poor practice to substitute *x* by *t* if *t* contains any variable that also is a bound variable of *φ*!
 φ ≡ (∀*z* ∈ ℕ *z*² > 0) ∨ (∃*y* ∈ ℕ *y* > 10 ∧ *x* < *y*) and *t* := 2*z* + 5.



Theorem 31

Let *x* be a variable, and ϕ and ψ be formulas which normally contain *x* as a free variable. Then the following equivalences hold:

| De Morgan's laws: | $(\neg(\forall x \ \phi)) \equiv (\exists x \ (\neg\phi))$ $(\neg(\exists x \ \phi)) \equiv (\forall x \ (\neg\phi))$ |
|-------------------------------|---|
| Trivial conjunction: | $(\forall \mathbf{x} \ (\phi \land \psi)) \equiv ((\forall \mathbf{x} \ \phi) \land (\forall \mathbf{x} \ \psi))$ |
| Only if $x \notin FV(\psi)$: | |



Rules of Inference

 Let x, y be variables and φ, ψ be propositional formulas. The following inference rules allow to deduce new formulas.

$$\frac{((\forall x \ \phi) \lor (\forall x \ \psi))}{(\forall x \ (\phi \lor \psi))} \qquad \qquad \frac{(\exists x \ (\phi \land \psi))}{(\exists x \ \phi) \land (\exists x \ \psi)} \qquad \qquad \frac{(\exists x \ (\forall y \ \phi))}{(\forall y \ (\exists x \ \phi))}$$

- Note that the other direction does not hold for any of these inference rules!
- In addition to these three inference rules all the equivalence rules apply: Each
 equivalence can be written as two inference rules since they are valid in both
 directions.



Propositional and Predicate Logic

- Propositional Logic
- Predicate Logic
- Special Quantifiers



Special Quantifiers

• What is the syntactical meaning of

$$\sum_{i=m}^{n} f(i) \quad ?$$

• Apparently, this is the common short-hand notation for

$$\sum_{i=m}^{n} f(i) = \sum_{m \le i \le n} f(i) = \sum_{P(i,m,n)} f(i) = f(m) + f(m+1) + \dots + f(n-1) + f(n),$$

where f(i) is a term with the free variable *i* and $(m \le i \le n)$ is a formula with free variables *i*, *m*, *n*, and $P(i, m, n) :\Leftrightarrow [(i \ge m) \land (i \le n)]$.



Special Quantifiers

• Thus, the ∑-quantifier takes a predicate, *P*(*i*, *m*, *n*), and and a term, *f*(*i*), and converts it to the new term

$$(f(m) + f(m + 1) + f(m + 2) + \dots + f(n - 1) + f(n)),$$

By convention, the variable *i* is bound inside of $\sum_{i=m}^{n} f(i)$, while *m* and *n* remain free.

• Similarly,

$$\prod_{i=m}^n f(i) := f(m) \cdot f(m+1) \cdot f(m+2) \cdot \ldots \cdot f(n-1) \cdot f(n).$$

• Again, by convention, if *n* < *m* then

$$\sum_{i=m}^{n} f(i) := 0$$
 and $\prod_{i=m}^{n} f(i) := 1$.

Union (∪) and intersection (∩) of several sets are further examples of special quantifiers: ∪ⁿ_{i=1}A_i.



Special Quantifiers: Sets

- Standard notation for a set with a finite number of elements: $\{\ ,\ ,\ldots,\ \};$ e.g., $\{1,2,3,4\}.$
- Obvious disadvantage: explicit enumeration of all elements of a set allows to specify only finite sets!
- Infinite sets require us to give a statement *A* to specify a *characteristic property* of the set:

$$S := \{x : A\}$$
 or $S := \{f(x) : A\},$

where *S* shall contain those elements *x*, or those terms f(x), for some universe of discourse, for which the statement *A* holds.

- Typically, *x* will be a free variable of *A*.
- Thus, the three symbols "{" and ":" and "}" together act as a quantifier that binds *x*.



Convenient Short-Hand Notations

 The following short-hand notations are convenient for using the predicate x ∈ X in conjunction with sets or quantifiers:

 $\{x \in X : A(x)\}$ is a short-hand notation for $\{x : x \in X \land A(x)\}$

 $(\forall x \in X \ A(x))$ is a short-hand notation for $(\forall x \ (x \in X \Rightarrow A(x)))$

 $(\exists x \in X \ A(x))$ is a short-hand notation for $(\exists x \ (x \in X \land A(x)))$

 If x is a typed variable — e.g., a real number — and P is a "simple" unary predicate — e.g., P(x) :⇔ (x > 3) — then the following notations are also used commonly:

 $(\forall P(x) \ A(x))$ is a short-hand notation for $(\forall x \ (P(x) \Rightarrow A(x)))$

 $(\exists P(x) \ A(x))$ is a short-hand notation for $(\exists x \ (P(x) \land A(x)))$

• Another wide-spread notation is to drop the parentheses:

```
\forall x \ P(x) instead of (\forall x \ P(x))
```

Definitions and Theorem Proving

- Need for Rigorous Analysis
- Definitions
- Syntactical Proof Techniques
- Types of Proofs



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- Need for Rigorous Analysis
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3

- Syntactical Proof Techniques
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- Intuition: Start at 1 and scan the integers from 1 to 20, successively picking those integers which are compatible with all integers picked previously:

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• We get 9 compatible integers. Our selection scheme makes it plausible that this is indeed the maximum number of compatible integers within {1, 2, 3, ..., 19, 20}. Right?



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- Oops! Why should we believe that we can't find 11 or more compatible integers within {1,2,3,...,19,20}?
- The answer is provided by the pigeonhole principle (Thm. 147): Every subset of compatible integers of {1, 2, 3, ..., 19, 20} can contain at most one of each of the following 10 pairs:

| 1 | 2 | 3 | 4 | 5 | 11 | 12 | 13 | 14 | 15 |
|---|---|---|---|----|----|----|----|----|----|
| 6 | 7 | 8 | 9 | 10 | 16 | 17 | 18 | 19 | 20 |



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Even though proofs and a rigorous formal analysis might seem boring (difficult, mind-boggling, mind-numbing, unnecessary, ...) there is just no way around them if we want to be sure that our findings are correct!



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Even though proofs and a rigorous formal analysis might seem boring (difficult, mind-boggling, mind-numbing, unnecessary, \ldots) there is just no way around them if we want to be sure that our findings are correct!

So, be prepared for at least some boring (difficult, mind-boggling, mind-numbing, unnecessary, . . .) proofs!





Definitions and Theorem Proving

- Need for Rigorous Analysis
- Definitions
 - Basics of Definitions
 - Recursive Definitions
 - Fibonacci, Factorial, Sum, Product
 - Words
 - Caveats
- Syntactical Proof Techniques
- Types of Proofs



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- In the following slides on definitions and theorem proving we pre-suppose an "intuitive" understanding of natural numbers, integers, reals, etc.; e.g., as taught in school.
- We will later on put these number systems on slightly more formal grounds.



- We distinguish between *explicit* and *recursive* definitions.
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- Explicit definition of a function *f* with *n* arguments:

 $f(x_1, x_2, \ldots, x_n) := t,$

where the term *t* (normally) contains $x_1, x_2, ..., x_n$ as free variables.



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Warning

The definiendum does not occur in the definiens of an explicit definition of a function f or predicate P! That is, the symbols f and P do not appear on the right-hand side.

- It is common to use the special symbols := and :⇔ for definitions, where the symbol ":" appears on the side of the definiendum.
- Thus, one can also write =: or ⇔: to indicate that the definiendum is on the right-hand side.



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 - it forces the author to decide whether or not something is a consequence of prior knowledge or some newly introduced entity.
- However, if ":=" or ":⇔" are used once in a text then they have to be used for absolutely all definitions in that text!!



• Poster seen in a tutoring institute at Salzburg:

$$x^{2} + \rho x + q = 0$$

$$x_{1/2} = \frac{-\rho}{2} \pm \sqrt{\frac{\rho^{2}}{4}} - q$$

$$D = \frac{\rho^{2}}{4} - q$$

$$D > 0 \quad 2 \text{ Lösungen}$$

$$D = 0 \quad 1 \text{ Lösung}$$

$$D < 0 \quad \text{keine Lösung}$$



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- Can $x_{1/2}$ be derived?
- Can D be derived?

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- Better formalism:
 - If x_1, x_2 are the roots of the second-degree polynomial equation $x^2 + px + q = 0$, with $p, q \in \mathbb{R}$ and unknown $x \in \mathbb{R}$, then

$$x_{1/2} = -\frac{p}{2} \pm \sqrt{\frac{p^2}{4} - q}.$$



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• With $D := p^2 - 4q$ we get

$$D \left\{ \begin{array}{c} > \\ = \\ < \end{array} \right\} 0 : \left\{ \begin{array}{c} 2 \text{ distinct real roots,} \\ 1 \text{ real root,} \\ 0 \text{ real roots.} \end{array} \right.$$



- Aka: Inductive definition.
- How can we state

x is ancestor of *y* if *x* is parent of *y*, or if *x* is parent of parent of *y*, or if *x* is parent of parent of parent of *y*, or if ...

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Warning

To avoid infinite circles, the definiendum must not occur in the basis!



Definition 32 (Sum and product)

Consider *k* real numbers $a_1, a_2, ..., a_k \in \mathbb{R}$, together with some $m, n \in \mathbb{N}$ such that $1 \leq m, n \leq k$.



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$$\sum_{i=m}^{n} a_i := \begin{cases} 0 & \text{if } n < m, \\ a_m & \text{if } n = m, \\ (\sum_{i=m}^{n-1} a_i) + a_n & \text{if } n > m, \end{cases}$$



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• The definitions for *n* < *m* are convenience settings that have turned out to be useful in practice.



Definition 33 (Factorial, Dt.: Fakultät, Faktorielle)

For $n \in \mathbb{N}_0$,

$$n! := \begin{cases} 1 & \text{if } n \leq 1, \\ n \cdot (n-1)! & \text{if } n > 1. \end{cases}$$



| Defini | Definition 33 (Factorial, Dt.: Fakultät, Faktorielle) | | | | | | | | | | | | |
|--------|--|---|---|---|---|----|-----|-----|------|--------|---------|-----------|---|
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| | | | | | | | | | 7 | | 9 | 10 | _ |
| | n! | 1 | 1 | 2 | 6 | 24 | 120 | 720 | 5040 | 40 320 | 362 880 | 3 628 800 | |



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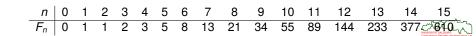
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Fibonacci Numbers

- The Fibonacci numbers are named after Leonardo da Pisa (1180?-1241?), aka "figlio di Bonaccio".
- The Fibonacci numbers have been studied extensively; they exhibit lots of interesting mathematical properties. For instance,

$$\lim_{n \to \infty} \frac{F_{n+1}}{F_n} = \phi, \quad \text{where } \phi := \frac{1 + \sqrt{5}}{2} \text{ is known as golden ratio.}$$



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 The Fibonacci numbers are also found in nature: E.g., the numbers of CW/CCW spirals of sunflower heads are given by subsequent Fibonacci numbers.



[Image credit: Wikipedia.]



73/406

500

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Definition 35 (Word)

Let Σ be a finite set. The set Σ^* of *words* over Σ is defined follows:

9 Base clause: The empty word, denoted by the Greek letter ϵ , belongs to Σ^* .



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- Pecursion clause: For all a ∈ Σ and all σ ∈ Σ*, the ordered pair (a, σ) belongs to Σ*.



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- E.g., $\Sigma := \{a, b, c, \dots, x, y, z\}$ or $\Sigma := \{0, 1\}$.

- Let Σ be a finite set. The set Σ^* of *words* over Σ is defined follows:
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 - Of course, in order to avoid confusion, *ϵ* is not allowed to be a character of Σ.
 - It is important to note that every element of Σ* is a finite sequence of zero or more characters (if we disregard the parentheses and commas) but that Σ* itself is an infinite set containing words of every possible finite length.

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Let Σ be a finite set. The *length* of a word σ over Σ is defined as follows:

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- In practice it is a convention to drop the ordered-pair notation and to write aσ rather than (a, σ). E.g., word rather than (w, (o, (r, (d, ε)))).
- Similarly, one writes word rather than wo rd. (This simplification is justified by the fact that the binary operator • is associative.)

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$$\frac{m}{n} \sharp \frac{p}{q} := \frac{m+p}{n+q}.$$

• Then
$$\frac{1}{1}\sharp_3^2 = \frac{3}{4}$$
, but $\frac{2}{2}\sharp_3^2 = \frac{4}{5}$.
• Since $\frac{1}{1} = \frac{2}{2}$, we conclude $\frac{4}{5} = \frac{3}{4}$, and, thus, $0 = 1$. Yikes!

3

Definitions and Theorem Proving

- Need for Rigorous Analysis
- Definitions
- Syntactical Proof Techniques
 - Syntax and Proofs
 - Equivalence Transformations
- Types of Proofs



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To *prove* a statement means to derive it from axioms (or postulates) and other previously established theorems by means of rules of logic.



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- Depending on the importance of the result, terms like *lemma* (Dt.: Lemma, Hilfssatz) or *corollary* (Dt.: Korollar) are also used instead of "theorem".
- A *conjecture* is a statement which has not yet been proved or disproved.
- The status of a conjecture may remain unknown for decades or even centuries: Fermat's Last Theorem was stated by Pierre de Fermat in 1637 and proved by Andrew Wiles (with the help of Richard Taylor) in 1993–1995.

- Syntactical proof techniques are proof techniques that are based on the analysis of the syntactical structure of a statement.
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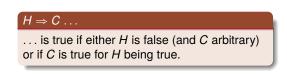
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Warning

In all the rules on this slide, A and B must not be part of a quantified formula. (Otherwise, get rid of the quantifier first!)

- If conclusion *C* is of the form $(\forall x \ A)$:
 - Proof technique: Let *x*₀ be arbitrary but fixed (Dt.: "beliebig aber fix"). From now on, *x*₀ can be treated as a constant!
 - It remains to prove $A[x_0/x]$ under the assumption *H*.



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- We are not allowed to make any assumptions on *x*₀ except for those that hold for all *x* in the universe of discourse.



- If conclusion *C* is of the form $(\exists x \ A)$:
 - Constructive Proof (Dt.: konstruktiver Beweis):
 - It "suffices" to find a suitable x_0 such that $A[x_0/x]$ if H.
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Proof (existential): We have p(2) = 5 > 0 and p(0) = -1 < 0. Since *p* is continuous on the closed interval [0,2], the Intermediate Value Theorem (Dt.: Zwischenwertsatz) tells us that there exists a real number *x* strictly between 0 and 2 such that p(x) = 0.



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 - Prove that such an x exists.
 - Prove its uniqueness.
- If hypothesis *H* is of the form $(\exists x \ A)$:
 - Let x_0 such that $A[x_0/x]$.
 - Add $A[x_0/x]$ to knowledge.
 - Again: x₀ must not occur anywhere else in H or C!



Natural-Language Synonyms of Formal Terms

- On many occasions a conjecture will not be stated in formal terms but by using a natural language.
- Then one has to *decode* the natural-language formulation and *translate* it into formal terms!



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 - A if and only if B, A genau dann wenn B,
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Equivalence Transformations

• First attempt to prove $(\forall n \in \mathbb{N} \quad \frac{2n+1}{n+1} \ge \frac{3}{2})$:

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• Correct proof of $(\forall n \in \mathbb{N} \mid \frac{2n+1}{n+1} \ge \frac{3}{2})$: Let $n \in \mathbb{N}$ be arbitrary but fixed. Then:

$$\begin{array}{ccc} & \frac{2n+1}{n+1} \geq \frac{3}{2} & |\cdot 2(n+1) \\ \Leftrightarrow & 2(2n+1) \geq 3(n+1) \\ \Leftrightarrow & 4n+2 \geq 3n+3 & |-(3n+2) \\ \Leftrightarrow & n \geq 1 \end{array}$$

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 $| \cdot a$
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 \Leftrightarrow

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Equivalence Transformations: Caveats

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$$\implies (x - \frac{9}{2})^2 = (y - \frac{9}{2})^2 \qquad | \sqrt{-}$$

$$\implies x - \frac{9}{2} = y - \frac{9}{2} \qquad | + \frac{9}{2}$$

$$\implies x = y$$

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Advice

In general, a relation a ∘ b may only be replaced by a new relation a' ∘ b' if one can argue that (a ∘ b) ⇔ (a' ∘ b').



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Advice

- In general, a relation a ∘ b may only be replaced by a new relation a' ∘ b' if one can argue that (a ∘ b) ⇔ (a' ∘ b').
- It is advisable to prove $a \circ b$, where $\circ \in \{=, <, >, \leqslant, \ge\}$, by constructing a chain $a_0 \circ a_1 \circ a_2 \circ \ldots \circ a_n$, with $a_0 = a$ and $a_n = b$, for some $n \in \mathbb{N}$.



3

Definitions and Theorem Proving

- Need for Rigorous Analysis
- Definitions
- Syntactical Proof Techniques
- Types of Proofs
 - Without Loss of Generality
 - Direct Enumeration
 - Case Analysis
 - Direct Proof
 - Proof by Contrapositive
 - Proof by Contradiction
 - Indirect Proof
 - Disproving Conjectures



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Warning

Do not use "w.l.o.g." unless *you could* indeed *explain* explicitly and in full detail how to carry on without that assumption!



Types of Proofs: Direct Enumeration

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• Note: Direct enumeration only works if the set given is finite!



Types of Proofs: Case Analysis

- Aka Proof by Exhaustion. Dt.: Fallunterscheidung.
- In order to prove $H \Rightarrow C$, it suffices to prove

 $A_1 \lor A_2 \lor \ldots \lor A_k$

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Warning

It is essential to guarantee that $A_1 \vee A_2 \vee \ldots \vee A_k$ holds, i.e., that no case is missing!



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 Proof: Factoring *n*⁷ − *n* yields

$$n^7 - n = n(n^6 - 1) = n(n^3 - 1)(n^3 + 1) = n(n - 1)(n^2 + n + 1)(n + 1)(n^2 - n + 1).$$



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- We want to prove $H \Rightarrow C$:
 - We build a chain of reasoning that starts at *H* and ends in *C*.
 - This approach is the classical example of deductive reasoning, where a logically valid sequence of steps establishes the truth of *C* under the assumption of *H*.



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Proof: (Direct Proof) Let $x_0, y_0 \in \mathbb{R}^+$ be arbitrary but fixed, with $x_0 < y_0$. We have $x_0 < y_0$, and therefore $x_0^2 = x_0 \cdot x_0 < y_0 \cdot x_0$.



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 - This approach is the classical example of deductive reasoning, where a logically valid sequence of steps establishes the truth of *C* under the assumption of *H*.
- Suppose we want to prove $(\forall x, y \in \mathbb{R}^+ \ (x < y) \Rightarrow (x^2 < y^2))$.

Proof: (Direct Proof) Let $x_0, y_0 \in \mathbb{R}^+$ be arbitrary but fixed, with $x_0 < y_0$. We have $x_0 < y_0$, and therefore $x_0^2 = x_0 \cdot x_0 < y_0 \cdot x_0$. Since $x_0 < y_0$ we know $y_0 \cdot x_0 < y_0^2$, and obtain $x_0^2 < y_0 \cdot x_0 < y_0^2$, which finally establishes $x_0^2 < y_0^2$:

$$x_0^2 = x_0 \cdot x_0 < y_0 \cdot x_0 < y_0 \cdot y_0 = y_0^2$$



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Types of Proofs: Proof by Contrapositive

- Dt.: Umkehrschluss, Kontraposition.
- We want to prove $H \Rightarrow C$:
 - In order to prove $H \Rightarrow C$ we build a (direct) proof for $(\neg C \Rightarrow \neg H)$.



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Warning

Make sure that the statements are negated correctly!



Types of Proofs: Proof by Contradiction

- Dt.: Widerspruchsbeweis.
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 - We assume $(H \land \neg C)$ as new hypothesis and prove $\neg H$.
 - This approach is correct since $(H \Rightarrow C) \equiv ((H \land \neg C) \Rightarrow \neg H)$.
- Warning: As when proving the contrapositive it is essential to check twice that the statements are indeed negated correctly!



- Aka Reductio ad absurdum.
- Dt.: indirekter Beweis.
- We want to prove $H \Rightarrow C$.
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 - Formally, $(H \land \neg C \land R) \Rightarrow \neg R$.
 - This is of the form $(A \Rightarrow B)$, and we have $(A \Rightarrow B) \equiv T$, where $B \equiv F$. Thus, $A \equiv F$.



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Note

Since an indirect proof is similar to a proof by contradiction, many textbooks treat it as one proof technique, or use the terms "reductio ad absurdum", "indirect proof", and "proof by contradiction" as synonyms.

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$$0 = \frac{p^3}{q^3} + \frac{p}{q} + 1$$
 and, thus, $0 = p^3 + pq^2 + q^3$.

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We do a case analysis, depending on whether p, q are even or odd:

Case p, q odd: Then $p^3 + pq^2 + q^3$ is odd, but 0 is even, yielding a contradiction to *R*.



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- Sometimes conjectures are false
- If the conjecture is of the form $(\forall x \ A)$:
 - Then we can disprove this conjecture by showing $(\exists x \neg A)$.
 - The latter is proved if we can come up with a *counterexample* (Dt.: Gegenbeispiel) to the original claim.



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 - Rather, to disprove this conjecture, we'd have to prove formally $(\forall x \neg A)$.



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4 Numbers and Basics of Number Theory

- Algebraic Structures
- Natural Numbers
- Integers
- Rational Numbers
- Real Numbers
- More Proof Techniques



4 Numbers and Basics of Number Theory

- Algebraic Structures
 - Operations
 - Properties of Operations
 - Group
 - Ring
 - Field

• Homomorphism and Isomorphism

- Natural Numbers
- Integers
- Rational Numbers
- Real Numbers
- More Proof Techniques



- An algebraic structure consists of a non-empty set together with one or more operations on it which satisfy certain identities ("*axioms*").
- The axioms tell us the properties of the operations.
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- Well-known example: \mathbb{R} with the standard addition "+".
- E.g., we have $(\sqrt{\pi} + 1) \sqrt{\pi} = 1$ because

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- Algebraic structures get their names based on the type of operations and axioms supported.
- Well-known structures include group, ring, field, and vector space. (Many more algebraic structures are studied in abstract algebra, though!)



Definition 40 (n-ary Operation, Dt.: n-stellig Verknüpfung)

Let *n* be a fixed non-negative integer and $X_1, X_2, ..., X_n$ be non-empty sets. An *n-ary* operation from $X_1, X_2, ..., X_n$ to another set *Y* is a function $\omega : X_1 \times X_2 \times ... \times X_n \to Y$.



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• An operation on a set *X* is also called an *internal operation* (Dt.: innere Verknüpfung).



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- An operation of arity zero is simply an element of the codomain *Y*, i.e., a constant.
- Note: The standard division ÷ is a binary operation neither on the natural numbers nor on the rational numbers.

• So, a binary operation on a set X is a function

 $\omega \colon X \times X \to X$ with $(x_1, x_2) \mapsto \omega(x_1, x_2)$ for $x_1, x_2 \in X$.



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- For binary operations it is customary to use symbols like ⋆, ∘, +, ·, ÷ rather than letters like ω.
- Furthermore, for binary operations it is common to use the infix notation

 $x_1 \star x_2$ or $x_1 + x_2$

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Definition 41 (Composition, Dt.: Hintereinanderausführung)

Consider two operations $f: A \to B$ and $g: B \to C$. The composition (Dt.: Komposition, Hintereinanderausführung) $g \circ f$ of f and g is defined as

 $(g \circ f)(x) := g(f(x))$ for all $x \in A$.

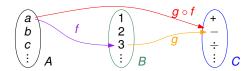


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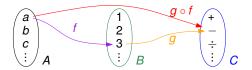
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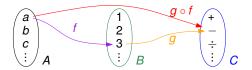
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- If A = B = C := X then \circ is a binary operation on operations from X to X.

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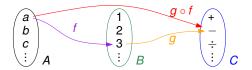
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- That is, the standard interpretation of *g f* is "carry out *f* followed by *g*".
- If A = B = C := X then \circ is a binary operation on operations from X to X.
- We will use the symbol
 o exclusively for denoting compositions of operations.

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Not all authors stick to the convention $(g \circ f)(x) := g(f(x)) \dots$





Properties of Operations: Associativity and Commutativity

Definition 42 (Associativity, Dt.: Assoziativität)

A binary operation * on a (non-empty) set G is associative if

 $\forall a, b, c \in G \quad (a \star b) \star c = a \star (b \star c).$



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- That is, the result does not change depending on whether the parentheses are associated with the first pair or the second pair of operands when the operation is applied to three operands.



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Definition 44 (Distributivity, Dt.: Distributivität)

A binary operation \cdot on a (non-empty) set *G* is *distributive over* a binary operation + on *G* if

$$\forall a, b, c \in G \quad a \cdot (b + c) = (a \cdot b) + (a \cdot c),$$

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- With the standard meaning of · and + over ℝ, multiplication distributes over addition, that is, when multiplying a sum by a factor we can distribute the factor over the summands.
- Note that addition does not distribute over multiplication (over \mathbb{R}).



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- Note that addition does not distribute over multiplication (over \mathbb{R}).
- Some textbooks prefer to split up the conditions of Def. 44 and say that · is *left-distributive* if

$$\forall a, b, c \in G \quad a \cdot (b + c) = (a \cdot b) + (a \cdot c),$$

and right-distributive if

$$\forall a, b, c \in G \quad (a+b) \cdot c = (a \cdot c) + (b \cdot c).$$



Properties of Operations: Neutral Element and Inverse Element

Definition 45 (Neutral element, Dt.: neutrales Element)

The element $n \in G$ is a *neutral element* (aka zero element, identity element) of a binary operation \star on a (non-empty) set *G* if

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• Hence, a neutral element of \star on *G* is an element in *G* that does not change the value of other elements when combined with them under the operation \star .



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- While addition over ℝ has zero as neutral element, subtraction does not have a neutral element: We get a 0 = a but, in general, 0 a ≠ a.



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Definition 46 (Inverse element, Dt.: inverses Element)

The element $b \in G$ is an *inverse element* of the element $a \in G$ for the binary operation \star on a (non-empty) set G if

 $a \star b = n = b \star a$,

where *n* denotes the neutral element of \star on *G*.



Properties of Operations: Uniqueness of Neutral Element

Lemma 47

A binary operation \star on a (non-empty) set *G* has at most one neutral element.



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A binary operation \star on a (non-empty) set G has at most one neutral element.

Proof: Assume that $n_1, n_2 \in G$ are neutral elements of \star on *G*. By Def. 45,

 $\forall a \in G \quad a \star n_1 = a = n_1 \star a \quad \text{and} \quad \forall a \in G \quad a \star n_2 = a = n_2 \star a.$



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These identities hold for all $a \in G$. Hence, in particular, they have to hold if $a := n_1$ and $a := n_2$:

 $n_2 \star n_1 = n_2 = n_1 \star n_2$ and $n_1 \star n_2 = n_1 = n_2 \star n_1$.



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Corollary 48

If a binary operation \star on a (non-empty) set G has a neutral element then it is unique.



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Corollary 48

If a binary operation \star on a (non-empty) set G has a neutral element then it is unique.

The neutral element is often denoted by 0 if + is used to denote the operation.

Lemma 49

An element $a \in G$ has at most one inverse element $b \in G$ for an associative binary operation \star on *G*.



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An element $a \in G$ has at most one inverse element $b \in G$ for an associative binary operation \star on G.

Proof: Assume that $b_1, b_2 \in G$ are inverse elements for $a \in G$ relative to an associative binary operation \star on G. Let $n \in G$ be the neutral element. By Def. 46,

 $a \star b_1 = n = b_1 \star a$ and $a \star b_2 = n = b_2 \star a$.



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If an element $a \in G$ has an inverse element relative to an associative binary operation \star on *G* then it is unique.



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If an element $a \in G$ has an inverse element relative to an associative binary operation \star on *G* then it is unique.

• Again, one may consider a left-inverse element and a right-inverse element.



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If an element $a \in G$ has an inverse element relative to an associative binary operation \star on *G* then it is unique.

- Again, one may consider a left-inverse element and a right-inverse element.
- The inverse element of *a* is often denoted by *a*⁻¹ if · or is used to denote the operation, and by *-a* if + denotes the operation.

Definition 51 (Group, Dt.: Gruppe)

A set G together with a binary operation \star on G defines a group if the following properties hold:



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- Since ★ is a binary operation on G, we know that G is closed under the application of ★. That is, if a, b ∈ G then a ★ b ∈ G.
- Note that a ★ b = b ★ a is not required for all a, b ∈ G. That is, commutativity need not hold!



Definition 52 (Abelian Group, Dt.: Abelsche Gruppe)



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A set *G* together with a binary operation \star on *G* defines an *Abelian group* (aka *commutative group*) if the following properties hold:

(G, \star) is a group.



Definition 52 (Abelian Group, Dt.: Abelsche Gruppe)

- **(** G, \star) is a group.
- **2** Commutativity: $\forall a, b \in G \quad a \star b = b \star a$.



Definition 52 (Abelian Group, Dt.: Abelsche Gruppe)

- $(G, \star) \text{ is a group.}$
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 - Sample (Abelian) groups: the integers Z under addition, non-zero rational numbers Q\{0} under multiplication.



Definition 52 (Abelian Group, Dt.: Abelsche Gruppe)

- $(G, \star) \text{ is a group.}$
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 - Sample (Abelian) groups: the integers Z under addition, non-zero rational numbers Q\{0} under multiplication.
 - Not a group: The integers under multiplication.



- A group (G, \star) is finite if G has a finite number of elements.
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- Multiplication tables for groups of orders two and three:

| * | n | а |
|---|---|---|
| n | n | а |
| а | а | п |

| * | n | а | b |
|---|---|---|---|
| n | n | а | b |
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- Up to renaming the elements of the groups, these are the only possible multiplication tables for groups of orders two and three.
- Again up to renaming, there are only two possible multiplication tables for groups with four elements, i.e., only two different groups.



• The *dihedral group* (Dt.: Diedergruppe) D_4 is formed by the clockwise rotations and reflections of a square which map the square onto itself:





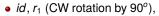
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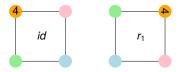
● *id*,





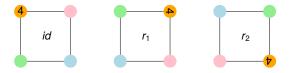
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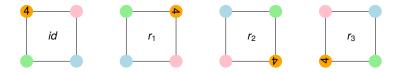


- The *dihedral group* (Dt.: Diedergruppe) D_4 is formed by the clockwise rotations and reflections of a square which map the square onto itself:
 - *id*, *r*₁ (CW rotation by 90°), *r*₂ (CW rotation by 180°),



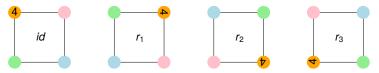


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 - *id*, *r*₁ (CW rotation by 90°), *r*₂ (CW rotation by 180°), *r*₃ (CW rotation by 270°);





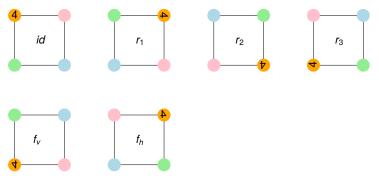
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 - *id*, r_1 (CW rotation by 90°), r_2 (CW rotation by 180°), r_3 (CW rotation by 270°);
 - f_v (vertical flip),





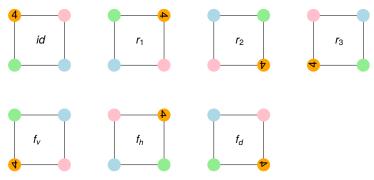


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 - *f_v* (vertical flip), *f_h* (horizontal flip),



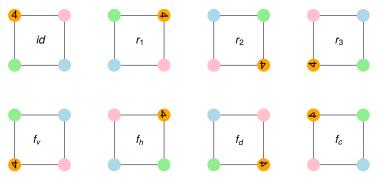


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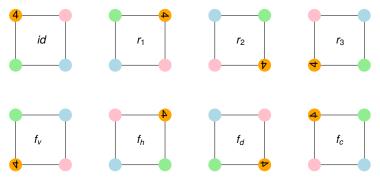


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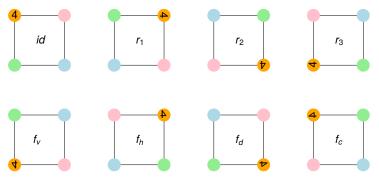
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• Does D₄ have eight elements? Or did we miss any element?



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• Does D₄ have eight elements? Or did we miss any element?

• No, we didn't!



- We denote the composition of functions by o.
- Multiplication table of *D*₄:

| 0 | id | <i>r</i> ₁ | <i>r</i> ₂ | <i>r</i> ₃ | f_V | f _h | f _d | f _c |
|-----------------------|-----------------------|-----------------------|-----------------------|-----------------------|-----------------------|-----------------------|-----------------------|-----------------------|
| id | id | <i>r</i> ₁ | <i>r</i> ₂ | <i>r</i> ₃ | f_v | f _h | f _d | f _c |
| <i>r</i> ₁ | <i>r</i> ₁ | r ₂ | <i>r</i> ₃ | id | f _c | f _d | f_v | f _h |
| <i>r</i> ₂ | <i>r</i> ₂ | r ₃ | id | <i>r</i> ₁ | f _h | f _v | f _c | f _d |
| <i>r</i> ₃ | <i>r</i> 3 | id | <i>r</i> ₁ | <i>r</i> 2 | f _d | f _c | f _h | f _v |
| f_v | f_{v} | f _d | f _h | f _c | id | <i>r</i> ₂ | <i>r</i> ₁ | <i>r</i> ₃ |
| f _h | f _h | f _c | f_{v} | f _d | <i>r</i> ₂ | id | r ₃ | <i>r</i> ₁ |
| f _d | f _d | f _h | f _c | f_{v} | r ₃ | <i>r</i> ₁ | id | <i>r</i> ₂ |
| f _c | f _c | f_{v} | f _d | f _h | <i>r</i> ₁ | <i>r</i> ₃ | <i>r</i> ₂ | id |

E.g., *f_d* ◦ *f_v*, which means flip vertically and then flip diagonally, corresponds to a (clockwise) rotation by 270°, i.e., to *r*₃.



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|-----------------------|-----------------------|-----------------------|-----------------------|-----------------------|-----------------------|-----------------------|-----------------------|-----------------------|
| id | id | <i>r</i> ₁ | <i>r</i> ₂ | <i>r</i> ₃ | f_v | f _h | f _d | f _c |
| <i>r</i> ₁ | <i>r</i> ₁ | r ₂ | <i>r</i> ₃ | id | f _c | f _d | f_v | f _h |
| <i>r</i> ₂ | <i>r</i> ₂ | r ₃ | id | <i>r</i> ₁ | f _h | f _v | f _c | f _d |
| <i>r</i> ₃ | <i>r</i> ₃ | id | <i>r</i> ₁ | <i>r</i> 2 | f _d | f _c | f _h | f _v |
| f_v | f_v | f _d | f _h | f _c | id | <i>r</i> ₂ | <i>r</i> ₁ | <i>r</i> ₃ |
| f _h | f _h | f _c | f_V | f _d | <i>r</i> ₂ | id | r ₃ | <i>r</i> ₁ |
| f _d | f _d | f _h | f _c | f_{v} | r ₃ | <i>r</i> ₁ | id | <i>r</i> ₂ |
| f _c | f _c | f_{v} | f _d | f _h | <i>r</i> ₁ | <i>r</i> ₃ | <i>r</i> ₂ | id |

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- Note: $f_d \circ f_v \neq f_v \circ f_d$. That is, D_4 is not commutative.



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|-----------------------|-----------------------|-----------------------|-----------------------|-----------------------|-----------------------|-----------------------|-----------------------|-----------------------|
| id | id | <i>r</i> ₁ | <i>r</i> ₂ | <i>r</i> ₃ | f_{v} | f _h | f _d | f _c |
| <i>r</i> ₁ | <i>r</i> ₁ | <i>r</i> ₂ | r ₃ | id | f _c | f _d | f_v | f _h |
| <i>r</i> ₂ | <i>r</i> ₂ | r ₃ | id | <i>r</i> ₁ | f _h | f _v | f _c | f _d |
| <i>r</i> ₃ | <i>r</i> ₃ | id | <i>r</i> ₁ | <i>r</i> 2 | f _d | f _c | f _h | f _v |
| f_{v} | f_{v} | f _d | f _h | f _c | id | <i>r</i> ₂ | <i>r</i> ₁ | <i>r</i> ₃ |
| f _h | f _h | f _c | f_{v} | f _d | <i>r</i> ₂ | id | <i>r</i> ₃ | <i>r</i> ₁ |
| f _d | f _d | f _h | f _c | f_{v} | r ₃ | <i>r</i> ₁ | id | <i>r</i> ₂ |
| f _c | f _c | f_v | f _d | f _h | <i>r</i> ₁ | <i>r</i> ₃ | <i>r</i> ₂ | id |

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- Note that each one of the transformations appears exactly once in each row and each column of the table: *Latin square*.



Real-World Application: Geometric Crystal Classes

- *D*₄ is one of the so-called *crystallographic point groups*, which describe sets of symmetry operations relative to a fixed point. Aka *geometric crystal class*.
- Each operation leaves the structure of the crystal unchanged. That is, the same types of atoms appear in similar positions as before the transformation induced by the operation.



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- Crystallographic point groups and their cousins, three-dimensional space groups, are studied and used by scientists such as crystallographers, mineralogists, and physicists.
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The Bauhinia flower has C_5 symmetry, and each star has D_5 symmetry.

This (color-inverted) snowflake has D₆ symmetry.

116/406 n Q Q C

Definition 53 (Ring, Dt.: Ring mit Eins)

A set *R* which possesses an "addition" $+ : R \times R \rightarrow R$ and a "multiplication"

 $\cdot : R \times R \rightarrow R$ defines a *(unit) ring* if the following conditions hold:



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Oistributivity:

 $\forall a, b, c \in R \quad a \cdot (b + c) = (a \cdot b) + (a \cdot c); \\ \forall a, b, c \in R \quad (a + b) \cdot c = (a \cdot c) + (b \cdot c).$



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 - Obstributivity: ∀a,b,c∈R a ⋅ (b + c) = (a ⋅ b) + (a ⋅ c); ∀a,b,c∈R (a + b) ⋅ c = (a ⋅ c) + (b ⋅ c).
 Neutral element: There exists an element 1 ∈ R ("one" element) such that ∀a ∈ R 1 ⋅ a = a = a ⋅ 1.



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• Note: The elements of a ring need not be numbers even though it is customary to use the terminology of arithmetic applied to numbers.



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- Note that $a \cdot b = b \cdot a$ for all $a, b \in R$ is not required. If commutativity holds then $(R, +, \cdot)$ forms a *commutative ring*.



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- Note that $a \cdot b = b \cdot a$ for all $a, b \in R$ is not required. If commutativity holds then $(R, +, \cdot)$ forms a *commutative ring*.
- Sample ring: The set of all continuous real-valued functions defined over an interval [α, β] ⊂ ℝ, with addition and multiplication of functions as operations, forms a ring.

Definition 54 (Field, Dt.: Körper)

A set *F* which possesses an "addition" $+ : F \times F \rightarrow F$ and a "multiplication" $\cdot : F \times F \rightarrow F$ defines a *field* if the following conditions hold:



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- **(**) For all $a \in F$ there exists an additive inverse $b \in F$, satisfying a + b = 0.
- **9** For all $a \in F \setminus \{0\}$ there exists a multiplicative inverse $b \in F$, satisfying $a \cdot b = 1$.



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- **Ommutativity:** $\forall a, b \in F \quad a \cdot b = b \cdot a$.
- **5** Distributivity: $\forall a, b, c \in F$ $a \cdot (b + c) = a \cdot b + a \cdot c$.
- **(**) Neutral element: There exists an element $0 \in F$ such that $\forall a \in F$ 0 + a = a.
- **2** Neutral element: There exists an element $1 \in F$ such that $\forall a \in F \mid 1 \cdot a = a$.
- **(**) For all $a \in F$ there exists an additive inverse $b \in F$, satisfying a + b = 0.
- For all *a* ∈ *F*\{0} there exists a multiplicative inverse *b* ∈ *F*, satisfying *a* ⋅ *b* = 1.
 0 ≠ 1.



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- For all $a \in F \setminus \{0\}$ there exists a multiplicative inverse $b \in F$, satisfying $a \cdot b = 1$. • $0 \neq 1$.
 - Again, the elements of *F* need not be numbers.
 - Note: The multiplication sign is often dropped if the meaning is clear within a specific context: It is common to write *ab* rather than *a* · *b*.

 In the sequel, we denote the additive neutral element of a field (*F*, +, ·) by 0 and its multiplicative neutral element by 1. Furthermore, we denote the inverse elements of *b* ∈ *F* by −*b* and *b*⁻¹.



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Lemma 57

Let $(F, +, \cdot)$ be a field. $\forall a \in F \quad a - a = 0$ and $\forall a \in F \setminus \{0\} \quad a \div a = 1.$

Theorem 58

Let $(F, +, \cdot)$ be a field. Then

-0 = 0 and $1^{-1} = 1$ and $\forall a \in F \quad 0 \cdot a = 0$.



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Proof: We have 0 = 0 + (-0) = -0.



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Theorem 59

Let $(F, +, \cdot)$ be a field. Then, for all $a, b \in F$, $(-1) \cdot a = -a$ and -(-a) = a and $(-a) \cdot b = -(a \cdot b) = a \cdot (-b)$ and $(-a) \cdot (-b) = a \cdot b$.



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Let $(F, +, \cdot)$ be a field. Then, for all $a, b \in F$,

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 $0=a^{-1}\cdot 0$



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Let $(F, +, \cdot)$ be a field. Then, for all $a, b \in F \setminus \{0\}$,

$$(a^{-1})^{-1} = a$$
 and $(a \cdot b)^{-1} = a^{-1} \cdot b^{-1}$.



Definition 62 (Fraction, Dt.: Bruch)

For $a \in F$, $b \in F \setminus \{0\}$, the *fraction* $\frac{a}{b}$ is defined as

$$\frac{a}{b} := a \div b.$$

We call a the enumerator (Dt.: Zähler) and b the denominator (Dt.: Nenner).



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Let $(F, +, \cdot)$ be a field. Then, for all $a, b, x, y \in F$ for which no denominator equals 0,

$$\frac{a}{b} \pm \frac{x}{y} = \frac{a \cdot y \pm b \cdot x}{b \cdot y}$$
 and $\frac{a}{b} \cdot \frac{x}{y} = \frac{a \cdot x}{b \cdot y}$ and $\frac{a}{b} \div \frac{x}{y} = \frac{a \cdot y}{b \cdot x}$.

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- E.g., $(\mathbb{R}, +)$ and (\mathbb{R}^+, \cdot) form groups. A group homomorphism from $(\mathbb{R}, +)$ to (\mathbb{R}^+, \cdot) is given by the exponential function $x \mapsto e^x$. (Recall that $e^{x+y} = e^x \cdot e^y$.)



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- An isomorphism is a bijective homomorphism.



Numbers and Basics of Number Theory

- Algebraic Structures
- Natural Numbers
 - Orders
 - Peano's Axioms for Introducing the Natural Numbers
 - The Principle of Mathematical Induction
 - Cardinality
- Integers
- Rational Numbers
- Real Numbers
- More Proof Techniques



How Shall We Define Natural Numbers or Real Numbers?

- Three options:
 - Ignore all formal details and presuppose an "intuitive" understanding of reals, integers, ...
 - Introduce the natural numbers, N, and then construct a hierarchy of number systems: N ⊂ Z ⊂ Q ⊂ R.
 - Set up the reals, \mathbb{R} , axiomatically and then define proper subsets for $\mathbb{N}, \mathbb{Z}, \mathbb{Q}$.



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 - Set up the reals, \mathbb{R} , axiomatically and then define proper subsets for $\mathbb{N}, \mathbb{Z}, \mathbb{Q}$.
- What is the best approach for a course on (applied) discrete mathematics? Much scholarly debate no consensus!
- We will start with introducing the natural numbers. However, since the gory details result in a lengthy discussion which provides little additional insight in \mathbb{N} and this is no course on number theory we base our introduction of \mathbb{N} on a simplified treatment of the so-called Peano axioms; see a book on number theory for a more formalized introduction of \mathbb{N} .



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Convention

In this course we adopt the following convention:

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- Caution: Read a text carefully to learn what an author means by "natural number". In particular, watch for clues such as terms like "positive natural numbers" (which indicates that zero is included) or statements like "n is a natural number, so it must be greater than zero" (which indicates that zero is not included).
- If one treats 0 as an element of $\mathbb N$ then $\{1,2,3,4,5,\ldots\}$ is often denoted by $\mathbb N^*.$



Definition 65 (Partial order, Dt.: Halbordnung)

A *partial order* on a set *S* is a binary relation \leq , i.e., a subset of $S \times S$, such that the following three properties hold for all $a, b, c \in S$:



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.

2 Transitivity:
$$(a < b \land b < c) \Rightarrow a < c$$
.



Definition 65 (Partial order, Dt.: Halbordnung)

A *partial order* on a set *S* is a binary relation \leq , i.e., a subset of $S \times S$, such that the following three properties hold for all $a, b, c \in S$:

- Reflexivity: $a \le a$.
- 2 Anti-symmetry: $(a \le b \land b \le a) \Rightarrow a = b$.
- **3** Transitivity: $(a \le b \land b \le c) \Rightarrow a \le c$.

If \leq is a partial order on S then (S, \leq) is called a *partially ordered set*, aka a *poset*.

Definition 66 (Strict partial order, Dt.: strikte Halbordnung)

A binary relation < on a set *S* forms a *strict partial order* on *S* if the following two properties hold for all $a, b, c \in S$:

- Irreflexivity: $\neg(a < a)$.
- 2 Transitivity: $(a < b \land b < c) \Rightarrow a < c$.

• A strict partial order is always *asymmetric*: If a < b then $\neg(b < a)$.

 $(a < b \land b < a) \stackrel{\text{trans.}}{\Rightarrow} a < a$, in contradiction to the irreflexivity: $\neg (a < a) \stackrel{\land \neg \uparrow}{\Rightarrow} a < a$,

Theorem 67

There is a one-to-one correspondence between non-strict and strict partial orders. Let *S* be a set and $a, b \in S$.

• If \leq is a non-strict partial order on *S* then the corresponding strict partial order "<" on *S* is the *reflexive reduction* given by

a < b : \Leftrightarrow $a \leq b$ and $a \neq b$.

If, on the other hand, < is a strict partial order on S then the corresponding non-strict partial order "≤" on S is the *reflexive closure* given by

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• As a notational convention, we omit the indication of an equality sign if we refer to a strict order, e.g., we write < rather than ≤ or ⊂ rather than ⊆.



• E.g., (\mathbb{Z}, \succeq) with (the non-strict order) \succeq as defined below forms a poset:

if *a* and *b* are even: $a \ge b$: $\Leftrightarrow a \ge b$

if *a* and *b* are odd: $a \ge b :\Leftrightarrow a \le b$



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Note that we do not know $a \ge b$ if one of a, b is even and the other one is odd. That is, if (S, \leq) is a poset then not all pairs of elements of S need to be comparable!



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• The subset relation, \subset , on the powerset $\mathcal{P}(X)$ of a set X is a strict partial order.



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• The subset relation, \subset , on the powerset $\mathcal{P}(X)$ of a set *X* is a strict partial order.

Definition 68 (Dual order, Dt.: duale Ordnung)

Let (S, \leq) resp. (S, <) be a (strict) poset. The *dual order* (or *reverse order*) on S, \geq resp. >, is defined as follows for $a, b \in S$:

 $a \ge b$: \Leftrightarrow $b \le a$ a > b : \Leftrightarrow b < a.



Let (S, \leq) be a poset and $T \subseteq S$. An element $a \in T$ is a *minimal element* of T if no $b \in T \setminus \{a\}$ exists such that $b \leq a$.



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Definition 70 (Least element, Dt.: kleinstes Element, Minimum)

Let (S, \leq) be a poset and $T \subseteq S$. An element $a \in T$ is a *least element* (or *minimum*) of T if $\forall b \in T \setminus \{a\} \ a \leq b$.



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Let (S, \leq) be a poset and $T \subseteq S$. An element $a \in T$ is a *greatest element* (or *maximum*) of T if $\forall b \in T \setminus \{a\} \ b \leq a$.



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Note: If a minimum or maximum exists then the anti-symmetry ensures that it is unique. Minimal or maximal elements need not be unique, though.

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Definition 73 (Total order, Dt.: totale Ordnung)



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A binary relation \leq on a set *S* forms a *total order* (or *linear order*) on *S* if the following three statements hold for all $a, b, c \in S$:

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Definition 74 (Well-order, Dt.: Wohlordnung)

A total order \leq on a set *S* forms a *well-order* if every non-empty subset of *S* has a least element.



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The set of all *natural numbers*, \mathbb{N} , together with an order relation \leq , is a totally ordered set defined as follows:



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 $\forall n \in \mathbb{N} \quad n+1 \in \mathbb{N} \quad \land \quad n < n+1.$

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\forall n \in \mathbb{N}, n \neq 1 \exists m \in \mathbb{N} \quad n = m + 1.
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The number n + 1 is called the *successor* of n, sometimes denoted by succ(n).

• N1 together with N2 establish the infinite sequence 1 < 2 < 3 < ...



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- The so-called *well-ordering principle*, N4, weeds out numbers like $\frac{1}{2}$ or π .
- One can show that the standard algebraic rules are compatible with the conditions imposed on \mathbb{N} , and that algebra and order interact smoothly within \mathbb{N} .
- One can also show that (up to a renaming of elements) there is only one set that fulfills all conditions of Def. 75. Hence, N is uniquely defined.

Definition 76 (Inductive)

A set $K \subseteq \mathbb{N}$ is *inductive* if

- $\bullet 1 \in K,$
- $2 \quad \forall k \in K \quad (k+1) \in K.$



Definition 76 (Inductive)

A set $K \subseteq \mathbb{N}$ is *inductive* if

- $\mathbf{0} \ \mathbf{1} \in K,$
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Theorem 78 (Weak Principle of Induction (W.P.I.))

Consider a predicate P over \mathbb{N} .

lf

P(1)

and if

 $\forall k \in \mathbb{N} \ (P(k) \Rightarrow P(k+1))$

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1 \in *K*, and 4 \forall *k* \in *K* (*k* + 1) \in *K*.

Thus, Thm. 77 is applicable and we conclude $K = \mathbb{N}$. That is, the predicate *P* holds for all natural numbers.



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- **1** *Induction basis* ("IB"): A basis step is done, i.e., *P*(1) is proved to be true.
- Induction hypothesis ("IH"): We assume P(k) to be true for an arbitrary but fixed k ∈ N.



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for some predicate P over \mathbb{N} by using induction:

- **1** *Induction basis* ("IB"): A basis step is done, i.e., *P*(1) is proved to be true.
- Induction hypothesis ("IH"): We assume P(k) to be true for an arbitrary but fixed k ∈ N.
- Inductive step ("IS"): We prove P(k + 1) based on the knowledge that P(k) is true.



• We claim that $\sum_{i=1}^{n} i = \frac{n \cdot (n+1)}{2}$ holds for all $n \in \mathbb{N}$.



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Proof: We use induction to prove our claim as follows:

• We define a suitable predicate *P*:

$$\forall n \in \mathbb{N} \left(P(n) :\Leftrightarrow \sum_{i=1}^{n} i = \frac{n \cdot (n+1)}{2} \right).$$



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• Induction basis (IB): We establish the truth of P(1):

$$\sum_{i=1}^{1} i = 1 = \frac{1 \cdot 2}{2}.$$



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Induction hypothesis (IH): Assume P(k) true for an arbitrary but fixed k ∈ N. That is, we assume (for this k ∈ N)

$$\sum_{i=1}^k i = \frac{k \cdot (k+1)}{2}.$$

Proof (cont'd):

• *Inductive step (IS)*: We have to prove P(k + 1) based on the induction hypothesis. That is, we have to prove

$$\sum_{i=1}^{k+1} i = \frac{(k+1) \cdot (k+2)}{2}$$



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We get

$$\sum_{i=1}^{k+1} i = \left(\sum_{i=1}^{k} i\right) + (k+1)$$



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Theorem 79 (Strong Principle of Induction (S.P.I.))

Consider a predicate P over \mathbb{N} . If

P(1)

and if

$$\forall k \in \mathbb{N} \left[\left(P(1) \land P(2) \land \ldots \land P(k) \right) \Rightarrow P(k+1) \right]$$

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Since

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• But W.P.I. and S.P.I. are equivalent, at least from a theoretical point of view.

Theorem 80 (S.P.I. with Larger Base)

Consider a predicate P over \mathbb{N} , and let $m \in \mathbb{N}$. If

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and apply the standard S.P.I.

• We could also carry out induction for smaller base values. That is, induction works for claims over \mathbb{N}_0 . (And even for negative base values!)

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• Thus, proving the base is mandatory!



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$$f(n) := \int_0^\infty \left(\prod_{i=0}^n \frac{\sin(\frac{x}{2i+1})}{\frac{x}{2i+1}}\right) \, dx.$$



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- We learn that f(n) is prime for all $0 \le n \le 40$. So, is f(n) always prime?
- No! For instance, f(41) is not prime.
- Thus, proving the inductive step is truly mandatory!



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 - IS: Consider a set *S* of *n* + 1 cats, and let *A* and *B* be two subsets of *S* such that

|A| = |B| = n and $A \cup B = S$.

Using the induction hypothesis and the transitivity of the equivalence, we conclude that all cats of the set S have the same color!



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As nature shows, this "proof" is seriously flawed



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$$2 \cdot (n+1) = 2 \cdot (k_1 + k_2) = 2 \cdot k_1 + 2 \cdot k_2 \stackrel{l.H.}{=} 0 + 0 = 0,$$

thus finishing the inductive "proof" ...



- Suppose that some limited and non-uniform resource has to be distributed fairly among *n* receivers, for some *n* ∈ N with *n* > 1.
- E.g., a cake (with fruits, whipped cream, chocolate crumbs, icing, etc.) might have to be distributed fairly among *n* kids. Aka: "Cake Cutting Problem".



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- The distribution should involve all kids such that each kid has to agree that it received a fair share of the cake by his/her preferences.



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Definition 81 (Fair distribution protocol)

A protocol for the distribution of a resource among *n* receivers is considered *fair* if each receiver gets at least 1/n-th of the resource (by his/her preferences), no matter what the preferences of the other receivers are and what the other receivers get.



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• How can we come up with a fair distribution protocol? Is there a general algorithm for fair cake cutting in the presence of *n* kids??



- If n = 2: Cut-and-choose distribution protocol.
 - Alice cuts the cake into two equal pieces (equal by her preferences).



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- If n > 2: Recursive application of the cut-and-choose distribution protocol.
 - The first n 1 kids cut the cake into n 1 pieces by applying the cut-and-choose distribution protocol recursively to n 2, n 3 etc. kids, thus each obtaining (hopefully) at least a fair 1/n-1 portion of the cake.



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 - **2** The *n*-th kid asks all other n 1 kids to cut his/her portion of the cake into *n* pieces such that the cutting is fair according to his/her preferences. (That is, according to each kid's preferences, each of the *n* pieces of his/her portion is equally desirable, for all of the first n 1 kids.)



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Theorem 82

The recursive cut-and-choose distribution protocol is fair.



Proof of Thm. 82 by induction: Assume that the total cake is worth 1 for each kid.

I.B.: n := 2 Alice cut the cake into two pieces that are equally desirable (according to her preferences) and, thus, both worth 1/2. Hence, she will get one half of the cake (by her preferences), no matter how Bob behaves.



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Hence, Bob can choose at least one half of the cake (according to his preferences), and both kids have no reason to complain about an unfair cutting.

I.H.: Assume that the recursive cut-and-choose cake cutting has been considered fair by the first k - 1 kids, for $k \ge 3$ arbitrary but fixed. Hence, each of the first k - 1 kids got a portion that is a least worth (according to the kid's preferences) $\frac{1}{k-1}$.



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$$\frac{w_1}{k} + \frac{w_2}{k} + \ldots + \frac{w_{k-1}}{k} = \frac{1}{k}(w_1 + w_2 + \ldots + w_{k-1}) = \frac{1}{k}.$$

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Theorem 85

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$$|\mathbb{N}| = |\mathbb{N} \times \mathbb{N}|$$
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Corollary 87

Consider three sets *A*, *B* and *C*. If $A \subseteq B \subseteq C$ and |A| = |C| then |A| = |B| = |C|.

Numbers and Basics of Number Theory

- Algebraic Structures
- Natural Numbers
- Integers
 - Construction of the Integers
 - Integral Powers
 - Divisibility and Prime Numbers
 - Quotient and Remainder
 - Congruences
 - Greatest Common Divisor
 - Chinese Remainder Theorem
- Rational Numbers
- Real Numbers
- More Proof Techniques



Integers: \mathbb{Z}

- Intuitive way to define the integers: $\mathbb{Z} := \mathbb{N}_0 \cup \{-n : n \in \mathbb{N}\}.$
- Thus, $\mathbb{Z} = \{0, \pm 1, \pm 2, \pm 3, \pm 4, \pm 5, \ldots\}.$
- $\bullet\,$ The blackboard-bold letter $\mathbb Z$ stands for the German word "Zahlen".



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- And how could we define a + b and $a \cdot b$ for $a, b \in \mathbb{Z}$??
- In order to put \mathbb{Z} on a more solid basis, we "extend" \mathbb{N} to obtain \mathbb{Z} .



• Let \cong_Z be a relation over \mathbb{N}_0 such that

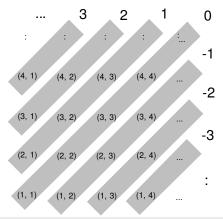
 $(a,b) \cong_Z (c,d) :\Leftrightarrow a+d=c+b.$



• Let \cong_Z be a relation over \mathbb{N}_0 such that

$$(a,b) \cong_Z (c,d) \quad :\Leftrightarrow \quad a+d=c+b.$$

• Easy to show: \cong_Z is an equivalence relation over \mathbb{N}_0 , with the equivalence classes shown below.



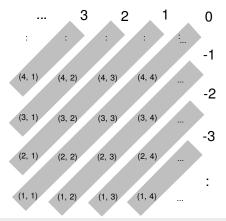


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• We interprete $[(a, b)]_{\cong_Z}$ as a - b.



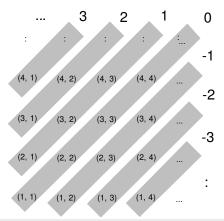


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Diskrete Mathematik (SS 2025)

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- We interprete $[(a, b)]_{\cong_Z}$ as a b.
- For n ∈ N, the equivalence classes [(n, 0)]_{≃z} form the natural numbers, while [(0, n)]_{≃z} form the negative integers.
- Zero is given by $[(0,0)]_{\cong_Z}$.



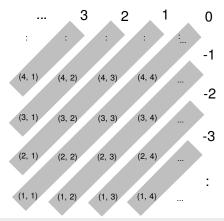


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Definition 88 (Integers)

The *integers* \mathbb{Z} are defined as $\mathbb{Z} := \{ [(a,b)]_{\cong_Z} : a, b \in \mathbb{N}_0 \}.$





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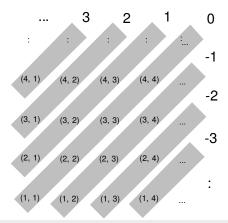
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• Furthermore, $\mathbb{Z}^+ := \mathbb{N}$ and $\mathbb{Z}_0^+ := \mathbb{N}_0$.





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It remains to define addition, multiplication and order on Z. For a, b, c, d ∈ N₀ we define an addition +_z, a multiplication ·_z and an order ≤_z as follows:

$$\begin{array}{lll} [(a,b)]_{\cong_{Z}} +_{Z} [(c,d)]_{\cong_{Z}} & := & [(a+c,b+d)]_{\cong_{Z}} \\ [(a,b)]_{\cong_{Z}} \cdot_{Z} [(c,d)]_{\cong_{Z}} & := & [(a\cdot c+b\cdot d,a\cdot d+b\cdot c)]_{\cong_{Z}} \\ [(a,b)]_{\cong_{Z}} \leqslant_{Z} [(c,d)]_{\cong_{Z}} & :\Leftrightarrow & a+d \leqslant b+c \end{array}$$



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An integer is *positive* if it is greater than zero and *negative* if it is less than zero; zero is neither positive nor negative.



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Theorem 90

 \mathbb{Z} is a countably infinite set. That is, $|\mathbb{N}| = |\mathbb{Z}|$.

Consider $x \in F$ for a field $(F, +, \cdot)$, with additive neutral element *e*. For $n \in \mathbb{N}_0$, we define integral powers of *x* as follows:

$$x^{n} := \begin{cases} 1 & \text{if } n = 0 \text{ and } x \neq e, \\ x & \text{if } n = 1, \\ x^{n-1} \cdot x & \text{if } n > 1. \end{cases}$$



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Lemma 92

X

Let $(F, +, \cdot)$ be a field. Then, for all $x, y \in F$ and all $m, n \in \mathbb{Z}$,

$$x^m \cdot x^n = x^{m+n}$$
 and $x^n \cdot y^n = (x \cdot y)^n$.

Definition 93 (Divisor, Dt.: Teiler, Faktor)

Let $a, b \in \mathbb{Z}$ with $a \neq 0$. Then *a divides b*, denoted by $a \mid b$, if there exists $c \in \mathbb{Z}$ such that $b = c \cdot a$.

 $a \mid b : \Leftrightarrow \exists c \in \mathbb{Z} \quad b = c \cdot a.$



Definition 93 (Divisor, Dt.: Teiler, Faktor)

Let $a, b \in \mathbb{Z}$ with $a \neq 0$. Then *a divides b*, denoted by $a \mid b$, if there exists $c \in \mathbb{Z}$ such that $b = c \cdot a$.

 $a \mid b : \Leftrightarrow \exists c \in \mathbb{Z} \ b = c \cdot a.$

In this case, we also say that *b* is a *multiple* of *a*, or *a* is a *divisor* or *factor* of *b*, or *b* is *divisible* by *a*. Otherwise we have $a \nmid b$. We have a *genuine divisor* if $a \mid b$ and $a \neq \pm 1$ and $a \neq \pm b$.



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Definition 94 (Even/odd, Dt.: gerade/ungerade)

A number $b \in \mathbb{Z}$ is said to be *even* if and only if $2 \mid b$; otherwise, b is *odd*.



Lemma 95

| $2 \forall a \in \mathbb{Z} \setminus \{0\} \ \forall b \in \mathbb{Z} \qquad a$ | $a \mid b \Rightarrow (\forall c \in \mathbb{Z} \ a \mid (b \cdot c)).$ |
|--|---|
| | $(a \mid b \land b \mid c) \Rightarrow a \mid c.$ |
| $ \ {\bf 0} \ \forall a \in \mathbb{Z} \backslash \{ 0 \} \ \forall b, c \in \mathbb{Z} $ | $(a \mid b \land a \mid c) \Rightarrow (\forall s, t \in \mathbb{Z} a \mid (b \cdot s + c \cdot t)).$ |
| | $a \mid b \Leftrightarrow (a \cdot c) \mid (b \cdot c).$ |
| $\bullet \forall a, b \in \mathbb{Z} \setminus \{0\} \qquad (a \mid b)$ | $\wedge b \mid a) \Rightarrow (a = b \lor a = -b).$ |



Lemma 95

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| $ 0 \forall \boldsymbol{a}, \boldsymbol{b} \in \mathbb{Z} \setminus \{0\} \; \forall \boldsymbol{c} \in \mathbb{Z} $ | $(a \mid b \land b \mid c) \Rightarrow a \mid c.$ |
| $ \ \ \textbf{0} \ \ \forall a \in \mathbb{Z} \backslash \{0\} \ \ \forall b, c \in \mathbb{Z} $ | $(a \mid b \land a \mid c) \Rightarrow (\forall s, t \in \mathbb{Z} a \mid (b \cdot s + c \cdot t)).$ |
| | $a \mid b \Leftrightarrow (a \cdot c) \mid (b \cdot c).$ |
| | $b \land b \mid a) \Rightarrow (a = b \lor a = -b).$ |

Lemma 96

For all $a, b, c \in \mathbb{Z}$ and all $k \in \mathbb{Z} \setminus \{0\}$,

 $(a = b + c \land k \mid b) \Rightarrow (k \mid a \Leftrightarrow k \mid c).$



Lemma 97

A number $a \in \mathbb{N}$ is divisible by

If its last digit is even, i.e., 0, 2, 4, 6 or 8;



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- if the hundreds digit is even and the number formed by the last two digits is divisible by eight, or if the hundreds digit is odd and the number formed by the last two digits plus four is divisible by eight;



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- If the sum of its digits is divisible by nine;



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- If the sum of its digits is divisible by nine;
- if its last digit is 0;
- if the alternating sum of its digits is divisible by eleven;



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- If it is divisible by three and four.
- There also exist divisibility rules for seven but all of them are a bit ackward



Proof of Lem. 97: We prove only the divisibility by three. Let $n \in \mathbb{N}$ and $a_0, a_1, \ldots, a_n \in \{0, 1, \ldots, 9\}$ such that

$$a=\sum_{i=0}^n a_i\cdot 10^i.$$



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Since

$$\mathbf{3} \mid \left(\sum_{i=0}^n a_i \cdot (\mathbf{10}^i - \mathbf{1})\right),$$

Lemma 96 implies that the number a is divisible by three if and only if

$$3 \mid \left(\sum_{i=0}^n a_i\right).$$

Definition 98 (Prime, Dt.: Primzahl)

A natural number $p \in \mathbb{N}$ is a *prime number*, or is *prime*, if $p \ge 2$ and if p is divisible only by 1 and p itself. All other natural numbers $p \ge 2$ are called *composite*.



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Definition 99 (Prime factor, Dt.: Primfaktor)

A natural number $p \in \mathbb{N}$ is a prime factor of $n \in \mathbb{N}$ if p is prime and $p \mid n$. If p is a prime factor of n then its *multiplicity* (Dt.: Vielfachheit) is the largest exponent k for which $p^k \mid n$.



Lemma 100

Let $k \in \mathbb{N}$ and $a_1, a_2, \ldots, a_k \in \mathbb{Z}$ and $p \in \mathbb{P}$. Then

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Theorem 102 (Fundamental Theorem of Arithmetic)

Every natural number n > 1 is representable uniquely in the form

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Corollary 103

There are infinitely many prime numbers.



- INT_MAX (in C/C++) is the eight Mersenne prime: $2147483647 = 2^{31} 1$.
- Mersenne primes are used by the *Mersenne twister*, a pseudo-random number generator developed in 1997 by Matsumoto and Nishimura.



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- Several unsolved problems related to Mersenne numbers:
 - Since $2^{11} 1 = 2047 = 23 \cdot 89$, not all Mersenne numbers are primes!
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Lemma 105

If $2^n - 1$ is prime for some $n \in \mathbb{N}$ then *n* is prime.

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- The strong version of this conjecture implies the weak version: If $n \in \mathbb{N}$, with $n \ge 7$, is odd then n' := n 3 is even with n' > 3. Hence, if n' can be written as the sum of two primes, then n can be written as the sum of three primes.



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- In November 2013, James Maynard reduced this bound to 600.
- This bound seems to have been further reduced to 246 by the Polymath project.



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For every natural number n > 2, the Diophantine equation $a^n + b^n = c^n$ has no solution $(a, b, c) \in \mathbb{N} \times \mathbb{N} \times \mathbb{N}$.

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- Finally proved by Andrew Wiles in 1993; a gap in the proof was fixed by Wiles and his former student Richard Taylor; the full proof was published in 1995.



Let $a \in \mathbb{N}$ and $b \in \mathbb{Z}$. Then there exist a unique *quotient* $q \in \mathbb{Z}$ and a unique *remainder* $r \in \mathbb{N}_0$ such that

```
b = a \cdot q + r and 0 \leq r < a.
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 We will use the operators div and mod for computing the quotient and remainder. That is, *q* and *r* of Lemma 110 are given by *q* := *b* div *a* and *r* := *b* mod *a*.



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- IEEE 754 defines a remainder function based on the round-to-nearest convention.

Warning

If one or both of a and b are allowed to be negative integers then the sign of the remainder may differ among different implementations!



• We know that $25 = (11001)_2$, i.e., $(11001)_2$ is the base-two representation of $25 = (25)_{10}$. (After all, $25 = 1 \cdot 2^4 + 1 \cdot 2^3 + 0 \cdot 2^2 + 0 \cdot 2^1 + 1 \cdot 2^0$.)



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- How can we represent an integer relative to an arbitrary base $b \in \mathbb{N} \setminus \{1\}$?
- Lemma 110 tells us that there exist unique $q_0, r_0 \in \mathbb{N}_0$ such that

 $n = b \cdot q_0 + r_0$ with $0 \leq r_0 < b$.



- We know that $25 = (11001)_2$, i.e., $(11001)_2$ is the base-two representation of $25 = (25)_{10}$. (After all, $25 = 1 \cdot 2^4 + 1 \cdot 2^3 + 0 \cdot 2^2 + 0 \cdot 2^1 + 1 \cdot 2^0$.)
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 The number r₀ becomes the rightmost digit of the base-b representation of n, and we seek q₁, r₁ such that

 $q_0 = b \cdot q_1 + r_1$ with $0 \leq r_1 < b$,

and so on until some $q_i = 0$.



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$$q_0 = b \cdot q_1 + r_1$$
 with $0 \leq r_1 < b$,

and so on until some $q_i = 0$.

• E.g.,

| 25 | = | $12 \cdot 2 + 1$ |
|----|---|------------------|
| 12 | = | $6 \cdot 2 + 0$ |
| 6 | = | $3 \cdot 2 + 0$ |
| 3 | = | $1 \cdot 2 + 1$ |
| 1 | = | $0\cdot 2+1$ |

and therefore $25 = (11001)_2$.

• Introduced by Carl Friedrich Gauss (1777–1855) in 1801.

Definition 111 (Congruence, Dt.: Kongruenz)

Let $a, b \in \mathbb{Z}$ and $m \in \mathbb{N}$. We say that *a* is *congruent* to *b modulo m*, and write

 $a \equiv_m b$,

if a - b is divisible by *m*. The term $a \equiv_m b$ is called a *congruence*.



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Lemma 112

For all $a, b \in \mathbb{Z}$ and $m \in \mathbb{N}$, we have $a \equiv_m b$ if and only if a and b have the same remainder after dividing by m, i.e., if and only if $a \mod m = b \mod m$.



| $38 \equiv_{12} 2$ | $even + even \equiv_2 even$ |
|--------------------|--|
| $-3 \equiv_5 2$ | $even + odd \equiv_2 odd$ |
| $0 \equiv_{3} 3$ | $\mathit{odd} + \mathit{odd} \equiv_2 \mathit{even}$ |
| $8 \equiv_3 2$ | even \cdot even \equiv_2 even |
| $7\equiv_3 1$ | $even \cdot odd \equiv_2 even$ |
| $7 \equiv_3 -8$ | $\mathit{odd}\cdot \mathit{odd} \equiv_2 \mathit{odd}$ |



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Lemma 113

For $m \in \mathbb{N}$, the relation \equiv_m is an equivalence relation on \mathbb{Z} , i.e., for all $a, b, c \in \mathbb{Z}$,

reflexivity $a \equiv_m a$,symmetryif $a \equiv_m b$ then $b \equiv_m a$, andtransitivityif $a \equiv_m b$ and $b \equiv_m c$ then $a \equiv_m c$ hold.

For $m \in \mathbb{N}$, the relation \equiv_m is a congruence relation on \mathbb{Z} , i.e., it respects addition, subtraction, and multiplication: Let $a, b, c, d \in \mathbb{Z}$ and $m \in \mathbb{N}$, and suppose that

$$a \equiv_m b$$
 and $c \equiv_m d$.

Then

 $a + c \equiv_m b + d$ and $a - c \equiv_m b - d$ and $a \cdot c \equiv_m b \cdot d$.



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Definition 115 (Residue, Dt.: Residuum, Restklasse)

Let $m \in \mathbb{N}$ with $m \ge 2$. The equivalence classes of \mathbb{Z} modulo m are called *residues* (or remainders) modulo m. For $a \in \mathbb{Z}$, its equivalence class modulo m is denoted by $[a]_m$. The set of residues modulo m is denoted by \mathbb{Z}_m or $\mathbb{Z}/m\mathbb{Z}$.



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Lemma 116

Let $m \in \mathbb{N}$ with $m \ge 2$. Then $\mathbb{Z}_m = \{ [a]_m : a \in \mathbb{N}_0 \land a < m \}$.



Let $m \in \mathbb{N}$ with $m \ge 2$, and $[a]_m, [b]_m \in \mathbb{Z}_m$. On \mathbb{Z}_m we define an addition $+_m$ and a multiplication \cdot_m as follows.

 $[a]_m +_m [b]_m := [a+b]_m$ $[a]_m \cdot_m [b]_m := [a \cdot b]_m$



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Lemma 118

Let $m \in \mathbb{N}$ with $m \ge 2$. Then addition $+_m$ and multiplication \cdot_m on \mathbb{Z}_m are well-defined. Furthermore, $(\mathbb{Z}_m, +_m, \cdot_m)$ forms a commutative ring.



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Often the notation [*a*]_m is simplified by omitting the modulus *m*, i.e., by writing [*a*], or even by simply writing *a* if it is clear that *a* ∈ Z_m. Similarly for +_m and ⋅_m.



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- It is also common to write

a mod *m*

instead of

[**a**]_m.



Theorem 119 (Fermat's Little Theorem)



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If $p \in \mathbb{N}$ is prime then $a^p \equiv_p a$ for every $a \in \mathbb{N}$.

If a is not divisible by p then this yields a^{p-1} ≡_p 1. In particular, this congruence holds for all 1 ≤ a ≤ p − 1.



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 - **1** Pick a random integer a with 1 < a < n 1.
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 - Otherwise, repeat the test for some other value of $a \in \{2, 3, ..., n-2\}$.



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One can prove that the probability for incorrectly classifying n as prime goes to zero (in most cases) as the number of tests is increased.



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- Linear congruential generators (LCG, [Lehmer 1954]) have been well studied, are easy to implement and used frequently.
- They generate a sequence of non-negative integers less than some specified modulus *m* ∈ N according to the following recursive definition:

 $x_{n+1} := (a \cdot x_n + c) \mod m,$

where

| $m \in \mathbb{N}$ | with | <i>m</i> > 1 | modulus, |
|------------------------|------|---------------------|-----------------|
| <i>a</i> ∈ ℕ | with | a < m | multiplier, |
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• E.g., m := 15, a := 1, c := 4 and $x_0 := 2$ yields the following sequence of numbers:

2 6 10 14 3 7 11 0 4 8 12 1 5 9 13 2 6 .



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• GCC/glibc: $m := 2^{31} - 1$, a := 1103515245, c := 12345.



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● GCC/glibc: *m* := 2³¹ − 1, *a* := 1103515245, *c* := 12345. More advanced pseudo-random number generators exist, e.g., Mersenne twister.



Greatest Common Divisor

Lemma 120

Let $a, b \in \mathbb{N}$. Then there exists a unique $n \in \mathbb{N}$ such that

 \bigcirc $n \mid a$ and $n \mid b$, and

2 for all $m \in \mathbb{N}$, if $m \mid a$ and $m \mid b$ then $m \leq n$.



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Definition 121 (Greatest common divisor, Dt.: größter gemeinsamer Teiler (ggT))

Let $a, b \in \mathbb{N}$. The unique number $n \in \mathbb{N}$ that exists according to Lem. 120 is called *greatest common divisor* of *a* and *b*, and is denoted by gcd(a, b).



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Definition 123 (Pairwise relatively prime)

A set *S* of natural numbers is called *pairwise relatively prime* (or *pairwise coprime* or *mutually coprime*) if all pairs of numbers *a* and *b* in *S*, with $a \neq b$, are relatively prime.

Let $a, b \in \mathbb{N}$. Then there exist $x, y \in \mathbb{Z}$ such that $gcd(a, b) = a \cdot x + b \cdot y$. Conversely, the smallest positive number $a \cdot x + b \cdot y$, for all $x, y \in \mathbb{Z}$, equals gcd(a, b).



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- For a, b, d ∈ Z given, the identity d = a ⋅ x + b ⋅ y over Z × Z is called a *linear* Diophantine equation in x and y.
- Lemma 124 was first stated by Étienne Bézout (1730–1783), and numbers $x, y \in \mathbb{Z}$ with $gcd(a, b) = a \cdot x + b \cdot y$ are called Bézout numbers.



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- Note: Bézout numbers are not unique! For instance, gcd(10, 15) = 5, and 10x + 15y = 5 has the solutions x := -1 and y := 1, and x := 2 and y := -1.



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Corollary 125

The numbers $a, b \in \mathbb{N}$ are relatively prime if and only if the linear Diophantine equation $a \cdot x + b \cdot y = 1$ has a solution, i.e., if and only if there exist $x, y \in \mathbb{Z}$ such that $a \cdot x + b \cdot y = 1$.



Theorem 126 (Euclidean Algorithm)

The following algorithm computes gcd(a, b) for $a, b \in \mathbb{N}_0$ with a > b.

```
function gcd(a,b)
precondition: a, b \in \mathbb{N}_0 with a > b.
postcondition: t = gcd(a,b)
while b > 0 do
t \leftarrow b
b \leftarrow a \mod b
a \leftarrow t
end while
t \leftarrow a
```



```
function gcd(a,b)

precondition: a, b \in \mathbb{N}_0 with a > b.

postcondition: t = gcd(a,b)

while b > 0 do

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end while

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```

• We want to compute the gcd of 78 and 99. Hence, b := 78 and $a := 99 = 1 \cdot 78 + 21$.



```
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precondition: a, b \in \mathbb{N}_0 with a > b.

postcondition: t = gcd(a,b)

while b > 0 do

t \leftarrow b

b \leftarrow a \mod b

a \leftarrow t

end while

t \leftarrow a
```

• We want to compute the gcd of 78 and 99. Hence, b := 78 and $a := 99 = 1 \cdot 78 + 21$. We get after different passes through the loop:

after 1st pass:
$$t = 78$$
, $b = 21$, $a = 78 = 3 \cdot 21 + 15$



```
function gcd(a,b)

precondition: a, b \in \mathbb{N}_0 with a > b.

postcondition: t = gcd(a,b)

while b > 0 do

t \leftarrow b

b \leftarrow a \mod b

a \leftarrow t

end while

t \leftarrow a
```

• We want to compute the gcd of 78 and 99. Hence, b := 78 and $a := 99 = 1 \cdot 78 + 21$. We get after different passes through the loop:

after 1st pass: t = 78, b = 21, $a = 78 = 3 \cdot 21 + 15$ after 2nd pass: t = 21, b = 15, $a = 21 = 1 \cdot 15 + 6$



```
function gcd(a,b)

precondition: a, b \in \mathbb{N}_0 with a > b.

postcondition: t = gcd(a,b)

while b > 0 do

t \leftarrow b

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end while

t \leftarrow a
```

• We want to compute the gcd of 78 and 99. Hence, b := 78 and $a := 99 = 1 \cdot 78 + 21$. We get after different passes through the loop:

| after 1st pass: | t = 78, | b = 21, | $a = 78 = 3 \cdot 21 + 15$ |
|-----------------|----------------|---------------|----------------------------|
| after 2nd pass: | <i>t</i> = 21, | b = 15, | $a = 21 = 1 \cdot 15 + 6$ |
| after 3rd pass: | <i>t</i> = 15, | <i>b</i> = 6, | $a = 15 = 2 \cdot 6 + 3$ |



```
function gcd(a,b)
precondition: a, b \in \mathbb{N}_0 with a > b.
postcondition: t = gcd(a,b)
while b > 0 do
t \leftarrow b
b \leftarrow a \mod b
a \leftarrow t
end while
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• We want to compute the gcd of 78 and 99. Hence, b := 78 and $a := 99 = 1 \cdot 78 + 21$. We get after different passes through the loop:

| after 1st pass: | <i>t</i> = 78, | b = 21, | $a = 78 = 3 \cdot 21 + 15$ |
|-----------------|----------------|----------------|----------------------------|
| after 2nd pass: | <i>t</i> = 21, | <i>b</i> = 15, | $a = 21 = 1 \cdot 15 + 6$ |
| after 3rd pass: | <i>t</i> = 15, | b = 6, | $a = 15 = 2 \cdot 6 + 3$ |
| after 4th pass: | <i>t</i> = 6, | b = 3, | $a=6=2\cdot 3+0$ |



```
function gcd(a,b)

precondition: a, b \in \mathbb{N}_0 with a > b.

postcondition: t = gcd(a,b)

while b > 0 do

t \leftarrow b

b \leftarrow a \mod b

a \leftarrow t

end while

t \leftarrow a
```

• We want to compute the gcd of 78 and 99. Hence, b := 78 and $a := 99 = 1 \cdot 78 + 21$. We get after different passes through the loop:

| after 1st pass: | t = 78, | <i>b</i> = 21, | $a = 78 = 3 \cdot 21 + 15$ |
|-----------------|----------------|----------------|----------------------------|
| after 2nd pass: | <i>t</i> = 21, | <i>b</i> = 15, | $a = 21 = 1 \cdot 15 + 6$ |
| after 3rd pass: | <i>t</i> = 15, | b = 6, | $a=15=2\cdot 6+3$ |
| after 4th pass: | <i>t</i> = 6, | b = 3, | $a=6=2\cdot 3+0$ |
| after 5th pass: | <i>t</i> = 3, | b = 0, | <i>a</i> = 3 |

• Hence, $t = 3 = \gcd(78, 99)$.



Theorem 127

Let $m \in \mathbb{N}$ with $m \ge 2$. An element $[a]_m \in \mathbb{Z}_m$ has a multiplicative inverse if and only if a is relatively prime to m.



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Corollary 128

Let $m \in \mathbb{N}$ with $m \ge 2$. The ring $(\mathbb{Z}_m, +_m, \cdot_m)$ forms a (finite) field if and only if *m* is prime.



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Let $m \in \mathbb{N}$ with $m \ge 2$. The ring $(\mathbb{Z}_m, +_m, \cdot_m)$ forms a (finite) field if and only if *m* is prime.

• If *m* is not prime then $(\mathbb{Z}_m, +_m, \cdot_m)$ may contain non-trivial zero divisors.



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• If *m* is not prime then $(\mathbb{Z}_m, +_m, \cdot_m)$ may contain non-trivial zero divisors.

Lemma 129

Let $m \in \mathbb{N}$ with $m \ge 2$ and $[a]_m \in \mathbb{Z}_m$ such that m and a are relatively prime. Let $x, y \in \mathbb{Z}$ such that $a \cdot x + m \cdot y = 1$. Then $[a]_m \cdot m [x]_m = [1]_m$, i.e., $[x]_m$ is the multiplicative inverse element for $[a]_m$.



Euclidean Algorithm Revisited

• Recursive formulation of the Euclidean Algorithm.

```
\begin{array}{l} \mbox{function } \gcd\_recursive(a,b) \\ \mbox{precondition: } a,b \in \mathbb{N} \mbox{ with } a > b. \\ \mbox{if } (a \bmod b) = 0 \mbox{ then} \\ \mbox{return } b \\ \mbox{else} \\ \mbox{return } \gcd\_recursive(b,a \bmod b) \\ \mbox{end if} \end{array}
```



Theorem 130 (Extended Euclidean Algorithm)

The following algorithm computes $x, y \in \mathbb{Z}$ and $d \in \mathbb{N}$ such that $gcd(a, b) = d = a \cdot x + b \cdot y$ for $a, b \in \mathbb{N}_0$ with a > b.

```
function gcd_extended(a,b)

precondition: a, b \in \mathbb{N}_0 with a > b.

postcondition: (d, x, y) \in \mathbb{N} \times \mathbb{Z} \times \mathbb{Z} such that gcd(a, b) = d = a \cdot x + b \cdot y

if (a \mod b) = 0 then

return (b, 0, 1)

else

(d, x, y) \leftarrow gcd_extended(b, a \mod b)

return (d, y, x - y \cdot (a \operatorname{div} b))

end if
```



```
function gcd_extended(a,b)

postcondition: (d, x, y) \in \mathbb{N} \times \mathbb{Z} \times \mathbb{Z} such that gcd(a, b) = d = a \cdot x + b \cdot y

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return (d, y, x - y \cdot (a \operatorname{div} b))

end if
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```
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if (a \mod b) = 0 then

return (b, 0, 1)

else

(d, x, y) \leftarrow gcd_extended(b, a \mod b)

return (d, y, x - y \cdot (a \operatorname{div} b))

end if
```

| а | b | a div b | a mod b | d | Х | У | |
|----|----|---------|---------|---|---|---|--|
| 99 | 78 | 1 | 21 | | | | |
| 78 | 21 | 3 | 15 | | | | |



```
function gcd_extended(a,b)

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end if
```

| а | b | a div <i>b</i> | a mod b | d | X | У | |
|----|----|----------------|---------|---|---|---|--|
| 99 | 78 | 1 | 21 | | | | |
| 78 | 21 | 3 | 15 | | | | |
| 21 | 15 | 1 | 6 | | | | |



```
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```

| а | b | a div b | a mod b | d | Х | У |
|----|----|---------|---------|---|---|---|
| 99 | 78 | 1 | 21 | | | |
| 78 | 21 | 3 | 15 | | | |
| 21 | 15 | 1 | 6 | | | |
| 15 | 6 | 2 | 3 | | | |

```
function gcd_extended(a,b)

postcondition: (d, x, y) \in \mathbb{N} \times \mathbb{Z} \times \mathbb{Z} such that gcd(a, b) = d = a \cdot x + b \cdot y

if (a \mod b) = 0 then

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```

| а | b | a div <i>b</i> | a mod b | d | Х | У |
|----|----|----------------|---------|---|---|---|
| 99 | 78 | 1 | 21 | | | |
| 78 | 21 | 3 | 15 | | | |
| 21 | 15 | 1 | 6 | | | |
| 15 | 6 | 2 | 3 | | | |
| 6 | 3 | _ | 0 | 3 | 0 | 1 |



```
function gcd_extended(a,b)

postcondition: (d, x, y) \in \mathbb{N} \times \mathbb{Z} \times \mathbb{Z} such that gcd(a, b) = d = a \cdot x + b \cdot y

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end if
```

| а | b | a div <i>b</i> | a mod b | d | X | У | |
|----|----|----------------|---------|---|---|----|--|
| 99 | 78 | 1 | 21 | | | | |
| 78 | 21 | 3 | 15 | | | | |
| 21 | 15 | 1 | 6 | | | | |
| 15 | 6 | 2 | 3 | 3 | 1 | -2 | |
| 6 | 3 | _ | 0 | 3 | 0 | 1 | |



```
function gcd_extended(a,b)

postcondition: (d, x, y) \in \mathbb{N} \times \mathbb{Z} \times \mathbb{Z} such that gcd(a, b) = d = a \cdot x + b \cdot y

if (a \mod b) = 0 then

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| а | b | a div <i>b</i> | a mod b | d | X | У |
|----|----|----------------|---------|---|----|----|
| 99 | 78 | 1 | 21 | | | |
| 78 | 21 | 3 | 15 | | | |
| 21 | 15 | 1 | 6 | 3 | -2 | 3 |
| 15 | 6 | 2 | 3 | 3 | 1 | -2 |
| 6 | 3 | _ | 0 | 3 | 0 | 1 |



```
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end if
```

| а | b | a div <i>b</i> | a mod b | d | X | У |
|----|----|----------------|---------|---|----|-----|
| 99 | 78 | 1 | 21 | | | |
| 78 | 21 | 3 | 15 | 3 | 3 | -11 |
| 21 | 15 | 1 | 6 | 3 | -2 | 3 |
| 15 | 6 | 2 | 3 | 3 | 1 | -2 |
| 6 | 3 | _ | 0 | 3 | 0 | 1 |



```
function gcd_extended(a,b)

postcondition: (d, x, y) \in \mathbb{N} \times \mathbb{Z} \times \mathbb{Z} such that gcd(a, b) = d = a \cdot x + b \cdot y

if (a \mod b) = 0 then

return (b, 0, 1)

else

(d, x, y) \leftarrow gcd_extended(b, a \mod b)

return (d, y, x - y \cdot (a \operatorname{div} b))

end if
```

| а | b | a div <i>b</i> | a mod b | d | Х | У |
|----|----|----------------|---------|---|-----|-----|
| 99 | 78 | 1 | 21 | 3 | -11 | 14 |
| 78 | 21 | 3 | 15 | 3 | 3 | -11 |
| 21 | 15 | 1 | 6 | 3 | -2 | 3 |
| 15 | 6 | 2 | 3 | 3 | 1 | -2 |
| 6 | 3 | _ | 0 | 3 | 0 | 1 |



```
function gcd_extended(a,b)

postcondition: (d, x, y) \in \mathbb{N} \times \mathbb{Z} \times \mathbb{Z} such that gcd(a, b) = d = a \cdot x + b \cdot y

if (a \mod b) = 0 then

return (b, 0, 1)

else

(d, x, y) \leftarrow gcd_extended(b, a \mod b)

return (d, y, x - y \cdot (a \operatorname{div} b))

end if
```

• We want to compute $x, y \in \mathbb{Z}$ and $d \in \mathbb{N}$ such that gcd(99, 78) = d = 99x + 78y.

| а | b | a div <i>b</i> | a mod b | d | X | У |
|----|----|----------------|---------|---|-----|-----|
| 99 | 78 | 1 | 21 | 3 | -11 | 14 |
| 78 | 21 | 3 | 15 | 3 | 3 | -11 |
| 21 | 15 | 1 | 6 | 3 | -2 | 3 |
| 15 | 6 | 2 | 3 | 3 | 1 | -2 |
| 6 | 3 | _ | 0 | 3 | 0 | 1 |
| | | | | | | |

• Hence, $gcd(99, 78) = -11 \cdot 99 + 14 \cdot 78 = -1089 + 1092 = 3$.



 Old Chinese folk tale: A Chinese Emperor used to count his army after a battle by ordering them to form groups of different sizes:



- Old Chinese folk tale: A Chinese Emperor used to count his army after a battle by ordering them to form groups of different sizes:
 - The soldiers should form groups of 3 and report back the number of soldiers that did not end up in a group consisting of 3 soldiers.



- Old Chinese folk tale: A Chinese Emperor used to count his army after a battle by ordering them to form groups of different sizes:
 - The soldiers should form groups of 3 and report back the number of soldiers that did not end up in a group consisting of 3 soldiers.
 - Then the soldiers should form groups of 5 and report back the number of soldiers that did not end up in a group consisting of 5 soldiers.



- Old Chinese folk tale: A Chinese Emperor used to count his army after a battle by ordering them to form groups of different sizes:
 - The soldiers should form groups of 3 and report back the number of soldiers that did not end up in a group consisting of 3 soldiers.
 - On the soldiers should form groups of 5 and report back the number of soldiers that did not end up in a group consisting of 5 soldiers.
 - Then the soldiers should form groups of 7 and report back the number of soldiers that could not join a group consisting of 7 soldiers.



- Old Chinese folk tale: A Chinese Emperor used to count his army after a battle by ordering them to form groups of different sizes:
 - The soldiers should form groups of 3 and report back the number of soldiers that did not end up in a group consisting of 3 soldiers.
 - On the soldiers should form groups of 5 and report back the number of soldiers that did not end up in a group consisting of 5 soldiers.
 - Then the soldiers should form groups of 7 and report back the number of soldiers that could not join a group consisting of 7 soldiers.
 - Then the soldiers should form groups of 11 and report back the number of soldiers that did not end up in a group consisting of 11 soldiers.



- Old Chinese folk tale: A Chinese Emperor used to count his army after a battle by ordering them to form groups of different sizes:
 - The soldiers should form groups of 3 and report back the number of soldiers that did not end up in a group consisting of 3 soldiers.
 - On the soldiers should form groups of 5 and report back the number of soldiers that did not end up in a group consisting of 5 soldiers.
 - Then the soldiers should form groups of 7 and report back the number of soldiers that could not join a group consisting of 7 soldiers.
 - Then the soldiers should form groups of 11 and report back the number of soldiers that did not end up in a group consisting of 11 soldiers.
 - 5 ...
- Based on this information he was able to figure out the number *n* of soldiers in his army.



- Old Chinese folk tale: A Chinese Emperor used to count his army after a battle by ordering them to form groups of different sizes:
 - The soldiers should form groups of 3 and report back the number of soldiers that did not end up in a group consisting of 3 soldiers.
 - On the soldiers should form groups of 5 and report back the number of soldiers that did not end up in a group consisting of 5 soldiers.
 - Then the soldiers should form groups of 7 and report back the number of soldiers that could not join a group consisting of 7 soldiers.
 - Then the soldiers should form groups of 11 and report back the number of soldiers that did not end up in a group consisting of 11 soldiers.
 - 5 ...
- Based on this information he was able to figure out the number *n* of soldiers in his army.
- Indeed, a mathematical solution was provided by the Chinese mathematician Sun Tzu sometime in the third to fifth century, and republished by Qin Jiushao in 1247!





 $n \mod 3 = 1$





| | 000 000 000 000 000 000 000 000 000 | |
|--|---|--|
| 999999999999999999999999999 89999999999 | © © | |
| <i>n</i> mod 3 = 1 | <i>n</i> mod 5 = <mark>2</mark> | |





| | 00 00 00 000 000 000 000 000 000 | 00 00 00 000 000 000 000 000 000 000 000 |
|--------------------|--|--|
| 8 | 00 | 88 |
| <i>n</i> mod 3 = 1 | <i>n</i> mod 5 = <mark>2</mark> | <i>n</i> mod 7 = <mark>2</mark> |



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Theorem 131 (Chinese Remainder Theorem, Dt.: Chinesischer Restsatz)

If, for some $k \in \mathbb{N}$, the numbers $m_1, m_2, \dots, m_k \in \mathbb{N}$ are pairwise relatively prime, then the following system of simultaneous congruences has an integer solution *b* for all $a_1, a_2, \dots, a_k \in \mathbb{Z}$ given:

$$\begin{array}{c} b \equiv_{m_1} a_1 \\ b \equiv_{m_2} a_2 \\ \vdots \\ b \equiv_{m_k} a_k \end{array} \right\} (*)$$

Furthermore, all solutions of (*) are congruent modulo $m := \prod_{i=1}^{k} m_i$. That is, the solution is unique if constrained to $\{1, 2, ..., m\}$.



Proof: We show the existence of an integer solution. Consider $i \in \mathbb{N}$ with $1 \le i \le k$. Since m_1, m_2, \dots, m_k are pairwise relatively prime, $gcd(\frac{m}{m_i}, m_i) = 1$. Using the extended Euclidean algorithm (Thm. 130), we can find integers x_i and y_i such that

$$x_i \cdot m_i + y_i \cdot \frac{m}{m_i} = 1. \tag{(\star)}$$



Proof: We show the existence of an integer solution. Consider $i \in \mathbb{N}$ with $1 \le i \le k$. Since m_1, m_2, \dots, m_k are pairwise relatively prime, $gcd(\frac{m}{m_i}, m_i) = 1$. Using the extended Euclidean algorithm (Thm. 130), we can find integers x_i and y_i such that

$$x_i \cdot m_i + y_i \cdot \frac{m}{m_i} = 1. \tag{(\star)}$$

Let $b_i := y_i \cdot \frac{m}{m_i}$. Equation (*) guarantees that the remainder of b_i when divided by m_i is 1. On the other hand, for $j \neq i$ every m_i divides b_i evenly. Thus,

 $b_i \equiv_{m_i} 1$ and $b_i \equiv_{m_i} 0$ for all j with $j \neq i$ and $1 \leq j \leq k$.



Proof: We show the existence of an integer solution. Consider $i \in \mathbb{N}$ with $1 \le i \le k$. Since m_1, m_2, \dots, m_k are pairwise relatively prime, $gcd(\frac{m}{m_i}, m_i) = 1$. Using the extended Euclidean algorithm (Thm. 130), we can find integers x_i and y_i such that

$$x_i \cdot m_i + y_i \cdot \frac{m}{m_i} = 1. \tag{(\star)}$$

Let $b_i := y_i \cdot \frac{m}{m_i}$. Equation (*) guarantees that the remainder of b_i when divided by m_i is 1. On the other hand, for $j \neq i$ every m_i divides b_i evenly. Thus,

 $b_i \equiv_{m_i} 1$ and $b_i \equiv_{m_i} 0$ for all j with $j \neq i$ and $1 \leq j \leq k$.

Since congruences respect multiplication, we get

 $a_i \cdot b_i \equiv_{m_i} a_i$ and $a_i \cdot b_i \equiv_{m_i} 0$ for all j with $j \neq i$ and $1 \leq j \leq k$.



Proof: We show the existence of an integer solution. Consider $i \in \mathbb{N}$ with $1 \le i \le k$. Since m_1, m_2, \dots, m_k are pairwise relatively prime, $gcd(\frac{m}{m_i}, m_i) = 1$. Using the extended Euclidean algorithm (Thm. 130), we can find integers x_i and y_i such that

$$x_i \cdot m_i + y_i \cdot \frac{m}{m_i} = 1. \tag{(\star)}$$

Let $b_i := y_i \cdot \frac{m}{m_i}$. Equation (*) guarantees that the remainder of b_i when divided by m_i is 1. On the other hand, for $j \neq i$ every m_i divides b_i evenly. Thus,

$$b_i \equiv_{m_i} 1$$
 and $b_i \equiv_{m_i} 0$ for all j with $j \neq i$ and $1 \leq j \leq k$.

Since congruences respect multiplication, we get

 $a_i \cdot b_i \equiv_{m_i} a_i$ and $a_i \cdot b_i \equiv_{m_i} 0$ for all j with $j \neq i$ and $1 \leq j \leq k$.

Thus, one solution of the simultaneous congruences is given by

$$b:=\sum_{i=1}^k a_i\cdot b_i.$$



- The Emperor collected the following information:
 - When the soldiers formed groups of 3, one soldier was left out.
 - When the soldiers formed groups of 5, two soldiers were left out.
 - When the soldiers formed groups of 7, again two soldiers were left out.



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 - When the soldiers formed groups of 3, one soldier was left out.
 - When the soldiers formed groups of 5, two soldiers were left out.
 - When the soldiers formed groups of 7, again two soldiers were left out.
- That is, since $a_1 = 1$, $a_2 = 2$, $a_3 = 2$ and $m_1 = 3$, $m_2 = 5$, $m_3 = 7$ and $m = 3 \cdot 5 \cdot 7 = 105$:

 $n \equiv_3 1$ $n \equiv_5 2$ $n \equiv_7 2$



- The Emperor collected the following information:
 - When the soldiers formed groups of 3, one soldier was left out.
 - When the soldiers formed groups of 5, two soldiers were left out.
 - When the soldiers formed groups of 7, again two soldiers were left out.
- That is, since $a_1 = 1$, $a_2 = 2$, $a_3 = 2$ and $m_1 = 3$, $m_2 = 5$, $m_3 = 7$ and $m = 3 \cdot 5 \cdot 7 = 105$:

$$n \equiv_3 1$$
 $n \equiv_5 2$ $n \equiv_7 2$

• Hence, we are to find $x_1, x_2, x_3, y_1, y_2, y_3 \in \mathbb{Z}$ such that

$$3x_1 + 35y_1 = 1$$
 $5x_2 + 21y_2 = 1$ $7x_3 + 15y_3 = 1$.



- The Emperor collected the following information:
 - When the soldiers formed groups of 3, one soldier was left out.
 - When the soldiers formed groups of 5, two soldiers were left out.
 - When the soldiers formed groups of 7, again two soldiers were left out.
- That is, since $a_1 = 1$, $a_2 = 2$, $a_3 = 2$ and $m_1 = 3$, $m_2 = 5$, $m_3 = 7$ and $m = 3 \cdot 5 \cdot 7 = 105$:

 $n \equiv_3 1$ $n \equiv_5 2$ $n \equiv_7 2$

• Hence, we are to find $x_1, x_2, x_3, y_1, y_2, y_3 \in \mathbb{Z}$ such that

 $3x_1 + 35y_1 = 1$ $5x_2 + 21y_2 = 1$ $7x_3 + 15y_3 = 1$.

• We have $x_1 := 12$, $y_1 := -1$, $x_2 := -4$, $y_2 := 1$, $x_3 := -2$, $y_3 := 1$ and, thus,

 $n = (35 \cdot (-1) \cdot 1 + 21 \cdot 1 \cdot 2 + 15 \cdot 1 \cdot 2) \mod 105 = 37 \mod 105 = 37.$



- Secret sharing refers to the distribution of information related to a secret (e.g., a number) among a group of receivers such that the secret can only be reconstructed if all or, at least, a large percentage of the receivers cooperate.
- Ideally, the information received by one individual receiver shall be of no (or very little) help for the receiver to obtain the secret without the help of the others.



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- Ideally, the information received by one individual receiver shall be of no (or very little) help for the receiver to obtain the secret without the help of the others.
- A secret sharing scheme is called a (*t*, *n*) threshold scheme, or *t*-out-of-*n* scheme, if at least *t* of the *n* receivers have to cooperate. (Of course, *t* ≤ *n*.)
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- Typically, *t* is large relative to *n* but not identical to *n*.
- Several different variants of schemes for secret sharing are used in practice.
- At least two published schemes rely on the Chinese Remainder Theorem 131.
- We sketch the very basic idea of a scheme based on the Chinese Remainder Theorem 131. (In our simple scheme we have *t* := *n*.)



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Note that

$$m := \prod_{i=1}^{5} m_i = 2 \cdot 3 \cdot 5 \cdot 7 \cdot 11 = 2310 > 1234.$$



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• Now consider $a_i := 1234 \mod m_i$, for $1 \le i \le 5$. This gives us the numbers

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- The numbers *m_i* and *a_i* are passed to the *i*-th receiver.
- Note that each individual receiver has gained little information about the secret b.
- Rather, in our simple approach, all five receivers need to cooperate in order to recover *b*: They have to solve the following set of five congruences:

$$b\equiv_2 0$$
 $b\equiv_3 1$ $b\equiv_5 4$ $b\equiv_7 2$ $b\equiv_{11} 2$



• The five receivers have to solve the following set of five congruences:

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• Since *a*₁ = 0, we need to solve only four congruences and get the following four Diophantine equations.

$$3x_2 + 770y_2 = 1$$
 $5x_3 + 462y_3 = 1$ $7x_4 + 330y_4 = 1$ $11x_5 + 210y_5 = 1$



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Hence, the secret sought is recovered as

 $b = (-1) \cdot 770 \cdot 1 + (-2) \cdot 462 \cdot 4 + 1 \cdot 330 \cdot 2 + 1 \cdot 210 \cdot 2 = -3386 \equiv_{2310} 1234.$



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 - **3** Represent an integer n < m by its k remainders n_1, n_2, \ldots, n_k upon division by m_1, m_2, \ldots, m_k .



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- This approach works as long as all intermediate results are less than *m*.
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 - One can run the computations for the different remainders in parallel, thus speeding up the computation.
- Standard choices for the modules are numbers of the form $2^i 1$:
 - One can prove $gcd(2^{i} 1, 2^{j} 1) = 2^{gcd(i,j)} 1$, which makes it easy to ensure that the modules are relatively prime.

- Suppose that we want to limit our arithmetic operations to numbers less than 12.
- We choose the five modules

 $m_1 := 2, \qquad m_2 := 3, \qquad m_3 := 5, \qquad m_4 := 7, \qquad m_5 := 11.$

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• Thus, *b* := 1234 + 1000 is uniquely determined as the solution of the following set of five congruences:

$$b \equiv_2 0$$
 $b \equiv_3 2$ $b \equiv_5 4$ $b \equiv_7 1$ $b \equiv_{11} 1$



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4 Numbers and Basics of Number Theory

- Algebraic Structures
- Natural Numbers
- Integers
- Rational Numbers
 - Construction of the Rational Numbers
 - Properties
- Real Numbers
- More Proof Techniques



Definition 132 (Rational equivalence)

On $\mathbb{Z} \times \mathbb{N}$ we define the binary relation \cong_Q as follows:

 $(p_1,q_1) \cong_Q (p_2,q_2) \quad :\Leftrightarrow \quad p_1 \cdot q_2 = p_2 \cdot q_1.$



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The rational numbers \mathbb{Q} are defined as

 $\mathbb{Q} := \{ [(p,q)]_{\cong_Q} : p \in \mathbb{Z}, q \in \mathbb{N} \}.$



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 $\mathbb{Q} := \{ [(\boldsymbol{p}, \boldsymbol{q})]_{\cong_{\boldsymbol{Q}}} : \boldsymbol{p} \in \mathbb{Z}, \boldsymbol{q} \in \mathbb{N} \}.$

The canonical representative of $[(p,q)]_{\cong_Q}$ is denoted by $\frac{p'}{q'}$, where p' := p div gcd(|p|,q) and q' := q div gcd(|p|,q).



• It is easy to define an addition $+_Q$, multiplication \cdot_Q and order \leq_Q on \mathbb{Q} that turns $(\mathbb{Q}, +, \cdot)$ into a totally ordered field. E.g.,

 $[(p_1, q_1)]_{\cong_Q} +_Q [(p_2, q_2)]_{\cong_Q} := [(p_1 \cdot q_2 + p_2 \cdot q_1, q_1 \cdot q_2)]_{\cong_Q}$



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• Of course, it is standard to simplify the notation and write

$$\frac{p}{q}$$
 instead of $[(p,q)]_{\cong_Q}$.

But keep in mind that fractions are equivalence classes. Thus,

$$(1,3) \cong_Q (3,9) \cong_Q (3000,9000)$$
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• In the sequel we resort to standard knowledge and deal with rational numbers as we learned in school. (However, this could be formalized based on Def. 134!)



Properties: \mathbb{Q} Is Not Complete

Theorem 135

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Proof : Suppose that there exists $x \in \mathbb{Q}$ such that $x^2 = 2$. Hence, there exist $p \in \mathbb{Z}$ and $q \in \mathbb{N}$ such that

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• Hence, $\sqrt{2} \notin \mathbb{Q}$.



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Lemma 136

There exists a rational number between any two distinct rational numbers.



Theorem 137

 \mathbb{Q} is a countably infinite set.



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 ${\mathbb Q}$ is a countably infinite set.

Proof by Cantor: Construct a bijection between \mathbb{N} and $\mathbb{Z} \times \mathbb{N}$ (as a "superset" of \mathbb{Q}).

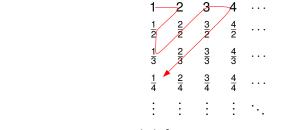
| 1 | 2 | 3 | 4 | • • • |
|---|---------------|----------------------------------|---------------|-------|
| <u>1</u> 2 | 2 2 2 | $\frac{3}{2}$ | 4 4 2 | |
| $\frac{1}{2}$ $\frac{1}{3}$ $\frac{1}{4}$ | <u>2</u> 3 | 3 32 33 33 34 | $\frac{4}{3}$ | |
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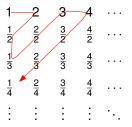
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This gives the sequence $1, 2, \frac{1}{2}, \frac{1}{3}, \frac{2}{2}, 3, \ldots$ Now start with zero and include every number's negative number, thus obtaining a systematic enumeration of $\mathbb{Z} \times \mathbb{N}$:

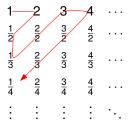
$$0 \quad 1 \quad -1 \quad 2 \quad -2 \quad \frac{1}{2} \quad -\frac{1}{2} \quad \frac{1}{3} \quad -\frac{1}{3} \quad \frac{2}{2} \quad -\frac{2}{2} \quad 3 \quad -3 \quad \dots$$

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Numbering this sequence yields a bijection from \mathbb{N} onto $\mathbb{Z} \times \mathbb{N}$, and Cor. 87 implies the claim.

Numbers and Basics of Number Theory

- Algebraic Structures
- Natural Numbers
- Integers

4

- Rational Numbers
- Real Numbers
 - Decimal Notation
 - Properties and Cardinality
- More Proof Techniques



- Intuitively, the reals comprise both rational and irrational numbers like $\sqrt{2}$ or π .
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- Convenient notations for intervals of real numbers:

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• Note: Some authors prefer to denote the open interval]*a*, *b*[by (*a*, *b*).



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 $\begin{array}{ll} \forall a,b \in \mathbb{R} \quad [a,b] := \{x \in \mathbb{R} : a \leqslant x \leqslant b\}; \\ \forall a,b \in \mathbb{R} \quad]a,b[:= \{x \in \mathbb{R} : a \leqslant x \leqslant b\}; \\ \forall a,b \in \mathbb{R} \quad [a,b[:= \{x \in \mathbb{R} : a \leqslant x \leqslant b\}; \\ \forall a,b \in \mathbb{R} \quad]a,b] := \{x \in \mathbb{R} : a < x \leqslant b\}. \end{array}$

- Note: Some authors prefer to denote the open interval]*a*,*b*[by (*a*,*b*).
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 $|x| := \max\{k \in \mathbb{Z} \colon k \leq x\},\$

 $[x] := \min\{k \in \mathbb{Z} \colon k \ge x\}.$



- Intuitively, the reals comprise both rational and irrational numbers like $\sqrt{2}$ or π .
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 Gauß introduced the square-bracket notation [x] ("Gaussklammer") in 1808. The names "floor" and "ceiling" and the corresponding notations were introduced by Iverson in 1962 in his book on APL.

• We have
$$[x] = \lfloor x \rfloor$$
 for all $x \in \mathbb{R}$.



Decimal Notation

Definition 138 (Decimal representation, Dt.: Dezimalzahl)

A real number $x \in \mathbb{R}_0^+$ is in *decimal representation* (or a *decimal number*) if it is represented as a sum of (negative) powers of ten:

$$x = x_0 + \sum_{i=1}^{\infty} \frac{x_i}{10^i}$$
, with an integer part $x_0 \in \mathbb{N}_0$ and with $0 \le x_i \le 9$ for all $i \in \mathbb{N}$.



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The decimal representation is *finite* if, for some $n_0 \in \mathbb{N}_0$, we have $x_i = 0$ for all $i \ge n_0$.



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• It is straightforward to extend Def. 138 to negative reals.



Definition 139 (Recurring decimal, Dt.: periodische Dezimalzahl)

A decimal representation of a real number is a *recurring decimal* (or *repeating decimal*) if it becomes periodic at some point: a finite subsequence of the digits after the decimal separator is repeated indefinitely.



4

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$$\frac{1}{3} = 0.333\cdots$$
$$\frac{1}{7} = 0.142857142857142857\cdots$$

are written as $0.\overline{3}$ or 0.3, and $0.\overline{142857}$. (The horizontal line is known as *vinculum*.)



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- Note: The decimal representation is not unique: we have 1.0 = 0.9 = 0.9999..., where the ellipsis "..." represents an infinite sequence of the digit 9.
- In fact, every non-zero, finitely represented decimal number has an alternate representation with trailing 9s, such as 123.4567 as 123.45669.

or

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Definition 141 (Irrational)

A number $x \in \mathbb{R} \setminus \mathbb{Q}$ is called *irrational*.

Definition 142 (Decimal separator)

The decimal separator is a symbol which is used to mark the boundary between the integer part and the fractional part of a number in decimal representation.



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Warning

A least two symbols are in wide-spread use for the decimal separator!

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- Most of Europe, most of South America and French Canada use the comma, while the UK, USA, Australia, English Canada and several Asian countries use a dot ("period", "full stop"). The dot also prevails in English-language publications.
- Dots or commas are frequently used to group three digits into groups within the integer part. However, this practice is discouraged by ISO!



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Well-Order "Theorem"

Every set can be well-ordered.

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Well-Order "Theorem"

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- In 1883, Georg Cantor stated that the Well-Order Theorem is a "fundamental law of thought". This statement started a mathematical flame war!
- In any case, this "theorem" can only be taken as an axiom, since it has been proved that it does not follow from any of the other commonly accepted axioms of set theory.
- In first-order logic, the Well-Order Theorem is equivalent to the Axiom of Choice (Dt.: Auswahlaxiom) and to Zorn's Lemma, in the sense that either one of them together with the Zermelo-Fraenkel Axioms allows to deduce the other onescience.

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Proof by Cantor (1891): Suppose to the contrary that there exists a bijection $a : \mathbb{N} \to \mathbb{R}$. We show that we can construct a number *r* which is not in the list a_1, a_2, a_3, \ldots :



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If $d_k = 1$ then $r_k := 2$ else $r_k := 1$. Now regard r_k as the *k*-th digit of a number $r \in \mathbb{R}$: we have $r = 0.r_1r_2r_3r_4...$ Since at least the *k*-th digit of *r* differs from the *k*-th digit of a_k , we conclude that $r \neq a_n$ for all $n \in \mathbb{N}$.



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• Hence, $|\mathbb{N}| < |\mathbb{R}|$.



For every $x \in \mathbb{R}$, every arbitrarily small neighborhood of x contains a rational number.

Sketch of proof: Let $x \in \mathbb{R}$ and $\varepsilon \in \mathbb{R}^+$ be arbitrary but fixed. W.l.o.g, 0 < x < 1. Let $k \in \mathbb{N}$ such that $10^{-k} < \varepsilon$.



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We define the rational number P/q as follows:

 $p := \left\lfloor x \cdot 10^k \right\rfloor \qquad \qquad q := 10^k$



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• E.g.,
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- E.g., $\pi \approx 3.1415 = \frac{31\,415}{10\,000}$ with $|\pi \frac{31\,415}{10\,000}| \leq \frac{1}{10\,000}$.
- Thus, we can approximate a real number by a rational number *P*/*q*.
- If the denominator *q* is a power of 10 then we can guarantee the error to be at most ¹/_q. Otherwise, if we allow an arbitrary integer *q* as denominator, we can guarantee the error to be at most ¹/_{q²}.

No (non-empty) set *A* has the same cardinality as its power set $\mathcal{P}(A)$.



Theorem 145

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• This implies that $|\mathbb{N}| < |\mathcal{P}(\mathbb{N})| < |\mathcal{P}(\mathcal{P}(\mathbb{N}))| < \dots$



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- The continuum hypothesis started out as a conjecture, until it was shown to be consistent with the usual axioms of the reals (by Gödel in 1940), and independent of those axioms (by Cohen in 1963).



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- The continuum hypothesis started out as a conjecture, until it was shown to be consistent with the usual axioms of the reals (by Gödel in 1940), and independent of those axioms (by Cohen in 1963).
- Under this hypothesis, the cardinality of \mathbb{R} equals \aleph_1 , and we have $2^{\aleph_0} = \aleph_1$.
- Furthermore, $|\mathcal{P}(\mathbb{R})| =: \aleph_2$, etc.



Numbers and Basics of Number Theory

- Algebraic Structures
- Natural Numbers
- Integers

4

- Rational Numbers
- Real Numbers
- More Proof Techniques
 - Pigeonhole Principle
 - Well-founded Induction
 - Structural Induction



The Pigeonhole Principle

- In 1834, Johann Dirichlet noted that if there are five objects in four drawers then there is a drawer with two or more objects.
- Pigeonhole Principle: If n letters are posted to m pigeonholes, then
 - at least one pigeonhole receives more than one letter if *n* > *m*.
 - at least one pigeonhole remains empty if n < m.
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Theorem 147 (Pigeonhole Principle, Dt.: Schubfachschluss)

Consider two finite sets A and B. If A has more elements then B then every mapping from A to B will cause at least one element of B to be the target of two or more elements of A.

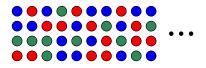


Consider a rectangular grid of points which consists of four rows and 100 columns.

0000000000 0000000000 0000000000

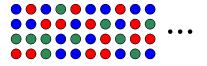


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Consider a rectangular grid of points which consists of four rows and 100 columns. Each point is colored with a color which is picked randomly among red, green and blue. Prove that there always exist four points of the same color that form the corners of a rectangle (with sides parallel to the grid), no matter how the coloring is done.





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Proof: A column pattern is the top-to-bottom sequence of colors assigned to the four points of a column of the grid. There are exactly $3^4 = 81$ different column patterns.

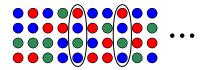


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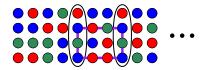
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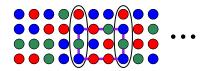
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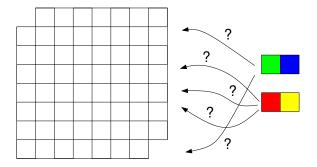
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Proof: A column pattern is the top-to-bottom sequence of colors assigned to the four points of a column of the grid. There are exactly $3^4 = 81$ different column patterns. Since there are more than 81 columns, we are guaranteed to have at least two columns with the same column pattern. Consider two such columns. Since there are four rows but only three colors, we conclude that two of the rows have the same color, thus giving us the four corners of the rectangle sought.

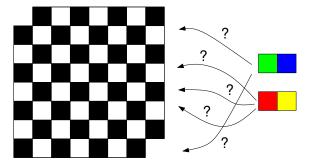
• Note: Just 19 columns suffice to guarantee the existance of such a rectangle

• Question: Can our modified chessboard be covered completely by 31 domino blocks of arbitrary color combinations?





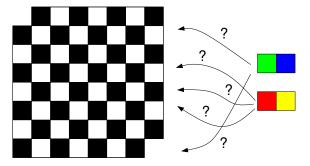
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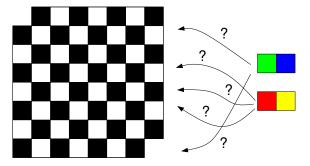
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- We observe that every permissible domino placement covers exactly one black square and one white square of the chessboard.
- Thus, all domino placements would establish a one-to-one mapping between black and white squares. However, there are 32 black squares and only 30 white squares! We conclude that our chessboard cannot be covered completely domino blocks.

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- Hence, every compression algorithm will increase the size of at least some file, or whether the sizes of all files unchanged.

A strict partial order < on M is called *well-founded* if every $X \subseteq M$, with $X \neq \emptyset$, has at least one minimal element relative to <. A poset (M, <) is called a well-founded poset if < is well-founded.



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Lemma 150

The poset (M, \prec) is well-founded if and only if no infinite strictly decreasing sequence in *M* exists, i.e., if an $a : \mathbb{N} \to M$ with $a_{i+1} < a_i$ for all $i \in \mathbb{N}$ does not exist.



Lexicographical Order

Definition 151

Let $(M_1, <_1)$ and $(M_2, <_2)$ be two posets. The *lexicographical ordering* $(<_1, <_2)_{lex}$ on $M_1 \times M_2$ is defined as

 $(a_1,b_1) \ (\prec_1,\prec_2)_{\textit{lex}} \ (a_2,b_2) \quad :\Leftrightarrow \quad \big((a_1\prec_1 a_2) \lor ((a_1=a_2) \land (b_1\prec_2 b_2))\big),$

where $(a_1, b_1), (a_2, b_2) \in M_1 \times M_2$.



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Lemma 152

Let $(M_1, <_1)$ and $(M_2, <_2)$ be two posets. Then $M_1 \times M_2$ together with the lexicographical order $(<_1, <_2)_{lex}$ is a poset.

• Similarly for a non-strict partial order \leq .



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Lemma 153

The posets (M_1, \prec_1) and (M_2, \prec_2) are well-founded if and only if $(M_1 \times M_2, (\prec_1, \prec_2)_{lex})$ is well-founded.



• Consider a predicate *P* over \mathbb{N} and recall the Strong Induction Principle (Thm 79): If *P*(1) and if

$$\forall k \in \mathbb{N} \left[\left(\forall (m \in \mathbb{N}, m \leq k) \ P(m) \right) \Rightarrow P(k+1) \right]$$

then

 $\forall n \in \mathbb{N} \ P(n).$



Induction Revisited

• And yet another version with "implicit" base: If

$$\forall k \in \mathbb{N} \left[\left(\forall (m \in \mathbb{N}, m < k) \ P(m) \right) \Rightarrow P(k) \right]$$

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• Note: The base case was not lost! Rather, it is included since we have to prove P(1) using the "helpful knowledge" that P(m) holds for all $m \in \mathbb{N}$ with m < 1.



Well-founded Induction

Theorem 154 (Principle of Well-founded Induction, Dt.: wohlfundierte Induktion)

Let (M, \prec) be well-founded and *P* be a predicate on *M*. If

$$\forall k \in M \left[\left(\forall (m \in M, m < k) \ P(m) \right) \Rightarrow P(k) \right]$$

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 That is, as inductive step we have to prove that the predicate holds for k if it holds for all predecessors m of k relative to ≺.

Proof: Let $X := \{m \in M : P(m) \text{ is false}\}$, and suppose $X \neq \emptyset$. Since (M, \prec) is well-founded, *X* has a minimal element *n*. Thus, $\forall m \in M$ with m < n the predicate P(m) holds. The inductive step

 $(\forall (m \in M, m < n) \ P(m)) \Rightarrow P(n)$

yields that P(n) holds, in contradiction to $n \in X$.

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Hence, both m_1 and m_2 are predecessors of k. By the inductive hypothesis, we know that m_1 is either prime or has a prime factorization; same for m_2 . Thus, also k has a prime factorization, which establishes the inductive step.



Partial Order on Recursive Structures

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- E.g., for a ∈ Σ and σ, σ' ∈ Σ*, if σ = aσ' then we could regard σ' to be "smaller" than σ.
- More generally, σ' <_Σ σ if and only if σ can be obtained from σ' and other words over Σ by applying constructors finitely often. (Hence, in this case σ' is a *sub-string* of σ.)
- Easy to prove: $<_{\Sigma}$ is a well-founded partial order on Σ^* .



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Structural induction can be seen as a special case of a well-founded induction.



Sample Structural Induction

Lemma 156

Let Σ be a finite set. For every $\sigma \in \Sigma^*$ we have $\sigma \bullet \epsilon = \epsilon \bullet \sigma = \sigma$.

Proof: Def. 37 immediately gives $\epsilon \bullet \sigma = \sigma$ for all $\sigma \in \Sigma^*$.



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The empty word ϵ is the only minimal element stated in the base case of the definition of $\Sigma^*,$ and we have

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Theorem 157 (Functional completeness of NAND)

The NAND junctor, \uparrow , is (truth-functionally) complete.

- Thus, every formula of propositional logic has a logically equivalent formula that uses only NAND junctors.
- Hence, any digital circuit can be realized by using only one type of gate: NAND gates. (This is also true for the NOR inverter.)



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Lemma 158

Let *p*, *q* denote two Boolean variables. The following logical equivalences hold:

$$\neg p \equiv (p \uparrow p) \qquad (p \land q) \equiv ((p \uparrow q) \uparrow (p \uparrow q)) \qquad (p \lor q) \equiv ((p \uparrow p) \uparrow (q \uparrow q))$$
$$(p \Rightarrow q) \equiv (\neg p \lor q) \qquad (p \Leftrightarrow q) \equiv ((p \Rightarrow q) \land (q \Rightarrow p))$$
$$\top \equiv (p \uparrow (p \uparrow p)) \qquad \bot \equiv (\top \uparrow \top)$$



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Proof of Thm. 157: Recall Def. 2: Propositional formulas (over some fixed set of *n* propositional variables p_1, p_2, \ldots, p_n) follow a rigid recursive construction scheme.



Proof of Thm. 157: Recall Def. 2: Propositional formulas (over some fixed set of *n* propositional variables $p_1, p_2, ..., p_n$) follow a rigid recursive construction scheme. Hence, we may use structural induction:

The minimal elements of the base case are given by the variables p₁, p₂,..., p_n and the constants ⊥ and ⊤. Lem. 158 tells us that ⊥ and ⊤ can be expressed using NAND junctors.



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- The minimal elements of the base case are given by the variables p_1, p_2, \ldots, p_n and the constants \perp and \top . Lem. 158 tells us that \perp and \top can be expressed using NAND junctors.
- Onsider an arbitrary but fixed propositional formula \u03c6₀ that contains at least one junctor.



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- The minimal elements of the base case are given by the variables p_1, p_2, \ldots, p_n and the constants \perp and \top . Lem. 158 tells us that \perp and \top can be expressed using NAND junctors.
- Consider an arbitrary but fixed propositional formula φ₀ that contains at least one junctor. By the construction scheme of propositional formulas, the formula φ₀ is of the form (¬φ₁) or (φ₁ # φ₂), for suitable propositional formulas φ₁, φ₂ and where # is one of the junctors ∧, ∨, ⇔, ⇒.



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Assume as inductive hypothesis that ϕ_1, ϕ_2 can be expressed using only NAND junctors (or are simply variables).



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By using the scheme outlined in Lem. 158, also ϕ_0 can be expressed using only NAND junctors.



Principles of Elementary Counting and Combinatorics

- Sum and Product Rule
- Inclusion-Exclusion Principle
- Binomial Coefficient
- Permutations
- Ordered Selection (Variation)
- Unordered Selection (Combination)



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Theorem 159 (Sum rule, Dt.: Additionsprinzip)

Let *A*, *B* be two finite sets with $A \cap B = \emptyset$. Then

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



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Corollary 160

For $n \in \mathbb{N}$, let A_1, A_2, \ldots, A_n be *n* finite sets that are pairwise disjoint. Then

$$|\mathbf{A}_1 \cup \mathbf{A}_2 \cup \ldots \cup \mathbf{A}_n| = \sum_{i=1}^n |\mathbf{A}_i|.$$



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Theorem 161 (Product rule, Dt.: Multiplikationsprinzip)

Let A, B be two finite sets. Then

 $|A \times B| = |A| \cdot |B|.$



Proof of Theorem 161:

• We observe that

$$A \times B = \bigcup_{b \in B} (A \times \{b\}), \quad \text{with } (A \times \{b_1\}) \cap (A \times \{b_2\}) = \emptyset \text{ if } b_1 \neq b_2.$$



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 There exists a bijective mapping between A and A × {b} for every b ∈ B. Thus, |A| = |A × {b}|, and the theorem is a consequence of Corollary 160.



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Corollary 162

For $n \in \mathbb{N}$, let A_1, A_2, \ldots, A_n be *n* finite sets. Then

$$|A_1 \times A_2 \times \ldots \times A_n| = \prod_{i=1}^n |A_i|.$$



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Corollary 163

For a propositional formula that contains n variables, 2^n evaluations are necessary in order to test all possible combinations of truth assignments to its variables.

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Characteristic Function and Cardinality of Power Set

Definition 164 (Characteristic function, Dt.: Indikatorfunktion)

Let *A* be a finite set, and $B \subseteq A$. The *characteristic function* $1_B : A \to \{0, 1\}$ indicates membership of an element of *A* in *B*:

$$1_B(a) := \begin{cases} 1 & \text{if } a \in B, \\ 0 & \text{if } a \notin B. \end{cases}$$



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Lemma 165

A finite set *A* has $2^{|A|}$ many different subsets. That is, $|\mathcal{P}(A)| = 2^{|A|}$.



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Proof: We observe that every subset of *A*, including \emptyset and *A* itself, has a one-to-one correspondance to a characteristic function. Thus, every subset of *A* corresponds to a sequence of *n* 0's and 1's, where n := |A|. We conclude that the power set $\mathcal{P}(A)$ has 2^n members.



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Lemma 166

Let *A* be a finite set, and $B \subseteq A$. Then $|B| = \sum_{a \in A} 1_B(a)$.

• How many 3-element strings *s* can be formed over the standard Latin alphabet — 26 lower-case letters — such that every string contains at least one *x*?



- How many 3-element strings *s* can be formed over the standard Latin alphabet 26 lower-case letters such that every string contains at least one *x*?
- Obviously such a 3-element string *s* is in exactly one of the following sets:

 $\begin{aligned} A_1 &:= \{s: \text{ first } x \text{ in first place of } s\}, \\ A_2 &:= \{s: \text{ first } x \text{ in second place of } s\}, \\ A_3 &:= \{s: \text{ first } x \text{ in third place of } s\}. \end{aligned}$



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• Since A₁, A₂, A₃ are pairwise disjoint, the Sum Rule 159 implies

$$|A_1 \cup A_2 \cup A_3| = |A_1| + |A_2| + |A_3| = 26 \cdot 26 + 25 \cdot 26 + 25 \cdot 25 = 1951.$$

- Suppose that passwords are limited to strings of six to eight characters, where each character is one of the 26 uppercase letters or a digit. Every password has to contain at least one digit.
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- By the Product Rule 161, the total number of six-character strings (over the 26 letters and the 10 digits) is 36⁶, with 26⁶ of them containing no digit at all. Hence,

 $N_6 = 36^6 - 26^6 = 1\,867\,866\,560.$



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$$N_7 = 36^7 - 26^7 = 70\,332\,353\,920$$

and

$$N_8 = 36^8 - 26^8 = 2\,612\,282\,842\,880.$$



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Hence, by the Sum Rule 159,

$$N = N_6 + N_7 + N_8 = 2\,684\,483\,063\,360.$$



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5

- Ordered Selection (Variation)
- Unordered Selection (Combination)



Theorem 167 (Inclusion-exclusion principle, Dt.: Siebprinzip, Poincaré-Formel)

Let A_1, A_2, \ldots, A_n be finite sets. Then

$$|\bigcup_{i=1}^{''} A_i| = \sum_{\substack{l \neq \emptyset \\ l \subseteq \{1, \dots, n\}}} (-1)^{|l|+1} |\bigcap_{i \in I} A_i|.$$



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$$\sum_{1 \le i \le n} (-1)^{1+1} |A_i| = \sum_{i=1}^n |A_i|$$



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$$\sum_{1 \leq i \leq n} (-1)^{1+1} |A_i| = \sum_{i=1}^n |A_i|. \qquad \sum_{1 \leq i < j \leq n} (-1)^{2+1} |A_i \cap A_j| = -\sum_{1 \leq i < j \leq n} |A_i \cap A_j|.$$



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• For |*I*| = 1: • For |*I*| = 2:

$$\sum_{\leqslant i \leqslant n} (-1)^{1+1} |A_i| = \sum_{i=1}^n |A_i|. \qquad \sum_{1 \leqslant i < j \leqslant n} (-1)^{2+1} |A_i \cap A_j| = -\sum_{1 \leqslant i < j \leqslant n} |A_i \cap A_j|.$$

• In particular:

1

$$|A_1 \cup A_2| = |A_1| + |A_2| - |A_1 \cap A_2|$$
$$|A_1 \cup A_2 \cup A_3| = |A_1| + |A_2| + |A_3| - |A_1 \cap A_2| - |A_1 \cap A_3| - |A_2 \cap A_3| + |A_1 \cap A_2 \cap A_3|$$

• How many bit strings of length eight either start with 1 as first bit or end in 00 as the two last bits? (This is a non-exclusive or!)



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- By the Product Rule 161,

 $|A_1| = 2^7 = 128$ and $|A_2| = 2^6 = 64$ and $|A_1 \cap A_2| = 2^5 = 32$.



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• Hence, by the Inclusion-Exclusion Principle (Thm. 167),

$$|A_1 \cup A_2| = |A_1| + |A_2| - |A_1 \cap A_2| = 128 + 64 - 32 = 160.$$



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Definition 168 (Binomial coefficient, Dt.: Binomialkoeffizient)

Let $n \in \mathbb{N}_0$ and $k \in \mathbb{Z}$. The binomial coefficient $\binom{n}{k}$ of n and k is defined as follows:

$$\binom{n}{k} := \begin{cases} 0 & \text{if } k < 0, \\ \frac{n!}{k! \cdot (n-k)!} & \text{if } 0 \le k \le n, \\ 0 & \text{if } k > n. \end{cases}$$

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- Recall *k*! := 1 for *k* := 0.
- The binomial coefficient $\binom{n}{k}$ is pronounced as "*n* choose *k*"; Dt.: "*n* über *k*".



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Lemma 169

Let $n \in \mathbb{N}_0$ and $k \in \mathbb{Z}$.

$$\binom{n}{0} = \binom{n}{n} = 1 \qquad \qquad \binom{n}{1} = \binom{n}{n-1} = n \qquad \qquad \binom{n}{k} = \binom{n}{n-k}$$

• The following table contains the non-zero values of $\binom{n}{k}$ for $0 \le n, k \le 6$.

| | | | | k | | | |
|--------|---|---|----|----|----|---|---|
| n | 0 | 1 | 2 | 3 | 4 | 5 | 6 |
| 0 | 1 | | | | | | |
| 1 | 1 | 1 | | | | | |
| 2 3 | 1 | 2 | 1 | | | | |
| 3 | 1 | 3 | 3 | 1 | | | |
| 4 | 1 | 4 | 6 | 4 | 1 | | |
| 5 | 1 | 5 | 10 | 10 | 5 | 1 | |
| 6 | 1 | 6 | 15 | 20 | 15 | 6 | 1 |



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| 1 | 1 | 1 | | | | | |
| 2 | 1 | 2 | 1 | | | | |
| 2 3 | 1 | 3 | 3 | 1 | | | |
| 4 | 1 | 4 | 6 | 4 | 1 | | |
| 4 5 6 | 1 | 5 | 10 | 10 | 5 | 1 | |
| 6 | 1 | 6 | 15 | 20 | 15 | 6 | 1 |

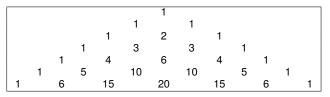
• Trivial to observe:

- Each row begins and ends with 1.
- Initially each row contains increasing numbers till its middle but then the numbers start to decrease.
- Each row's first half is exactly the mirror image of its second half.



Binomial Coefficients: Pascal's Triangle

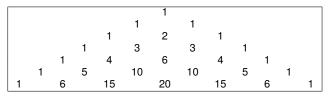
• A simple rearrangement of the previous table yields what is known as *Pascal's Triangle* in the Western world (Blaise Pascal, 1623–1662). But it was already studied in India in the 10th century, and discussed by Omar Khayyam (1048–1131)!





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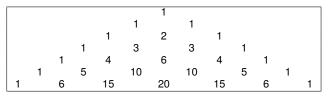


• All entries in this triangle, except for the left-most and right-most entries per row, are the sum of the two entries above them in the previous row.



Binomial Coefficients: Pascal's Triangle

• A simple rearrangement of the previous table yields what is known as *Pascal's Triangle* in the Western world (Blaise Pascal, 1623–1662). But it was already studied in India in the 10th century, and discussed by Omar Khayyam (1048–1131)!



• All entries in this triangle, except for the left-most and right-most entries per row, are the sum of the two entries above them in the previous row.

Theorem 170 (Khayyam, Yang Hui, Tartaglia, Pascal)

For $n \in \mathbb{N}$ and $k \in \mathbb{Z}$,

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

Binomial Theorem

• We know: $(a + b)^2 = a^2 + 2ab + b^2$ and $(a + b)^3 = a^3 + 3a^2b + 3ab^2 + b^3$.



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Theorem 171 (Binomial Theorem, Dt.: Binomischer Lehrsatz)

For all $n \in \mathbb{N}_0$ and $a, b \in \mathbb{R}$,

$$(a+b)^n = \binom{n}{0}a^n + \binom{n}{1}a^{n-1}b + \dots + \binom{n}{n}b^n$$

or, equivalently,

$$(\mathbf{a}+\mathbf{b})^n = \sum_{i=0}^n \binom{n}{i} \mathbf{a}^{n-i} \mathbf{b}^i.$$



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$$(\mathbf{a}+\mathbf{b})^n = \sum_{i=0}^n \binom{n}{i} \mathbf{a}^{n-i} \mathbf{b}^i.$$

Corollary 172

For all $n \in \mathbb{N}$ and all $x \in \mathbb{R}$:

$$\sum_{i=0}^{n} \binom{n}{i} x^{i} = (1+x)^{n} \qquad \sum_{i=0}^{n} \binom{n}{i} = 2^{n} \qquad \sum_{i=0}^{n} (-1)^{i} \binom{n}{i} = 0$$

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Principles of Elementary Counting and Combinatorics

- Sum and Product Rule
- Inclusion-Exclusion Principle
- Binomial Coefficient
- Permutations

5

- Ordered Selection (Variation)
- Unordered Selection (Combination)



Definition 173 (Permutation)

Let A be a finite set. A *permutation* of A is a bijective function from A to A.



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- Many encryption schemes used in cryptography can be seen as permutations.
- Standard notation for a permutation π of $\{1, 2, ..., n\}$:

$$\left(\begin{array}{cccc} 1 & 2 & 3 & \dots & n \\ \pi(1) & \pi(2) & \pi(3) & \dots & \pi(n) \end{array}\right)$$



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$$\begin{pmatrix} 1 & 2 & 3 & \dots & n \\ \pi(1) & \pi(2) & \pi(3) & \dots & \pi(n) \end{pmatrix}$$

$$\begin{pmatrix} 1 & 2 & 3 & 4 \\ 2 & 3 & 1 & 4 \end{pmatrix}$$



• E.g., for n := 4:

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Definition 174 (Product of permutations)

 $\left(\begin{array}{cccc} 1 & 2 & 3 & \dots & n \\ \pi(1) & \pi(2) & \pi(3) & \dots & \pi(n) \end{array}\right)$

Let A be a finite set together with two permutations α, β . Then the *product* (or *composition*) $\alpha \circ \beta$ is the function

$$\alpha \circ \beta : A \to A$$
 with $(\alpha \circ \beta)(a) := \alpha(\beta(a))$ for all $a \in A$.



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- The product of two permutations is not commutative.

$$\alpha := \begin{pmatrix} 1 & 2 & 3 & 4 \\ 1 & 3 & 4 & 2 \end{pmatrix} \qquad \beta := \begin{pmatrix} 1 & 2 & 3 & 4 \\ 2 & 1 & 3 & 4 \end{pmatrix}$$
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Lemma 175

For all $n \in \mathbb{N}$ and all finite sets A with n = |A|, the set of all permutations, S_n , over A together with \circ as operation forms a group, the so-called *symmetric group*.



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Lemma 176

For all $n \in \mathbb{N}$ and all finite sets A with n = |A|, the group (S_n, \circ) is a finite group with exactly n! members.

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Definition 177 (Cycle, Dt.: Zyklus)

Let *A* be a finite set of cardinality *n*. A permutation π of *A* is a *cycle of length* $k \le n$ if there exists a set $B \subseteq A$ with |B| = k such that, with $B := \{b_1, b_2, \dots, b_k\}$,

$$\pi(b_1) = b_2, \ \pi(b_2) = b_3, \ \dots, \ \pi(b_{k-1}) = b_k, \ \pi(b_k) = b_1,$$

and $\pi(a) = a$ for all $a \in A \setminus B$. In this case this *k*-cycle is written as

 $(b_1 \ b_2 \ \dots \ b_k)$ or as $b_1 \mapsto b_2 \mapsto \dots \mapsto b_k \mapsto b_1$.

A cycle is *non-trivial* if $k \ge 2$.



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A transposition is a cycle of length two, aka 2-cycle.



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Definition 178 (Transposition)

A transposition is a cycle of length two, aka 2-cycle.

Lemma 179

Every permutation (of two or more elements) can be written as

- (1) a product of cycles,
- (2) a product of transpositions.

Lemma 180

If two different products of transpositions correspond to the same permutation then both products consist of either an even or an odd number of transpositions.



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Definition 181 (Signature, Dt.: Signum)

The *signature* of a permutation is +1, and the permutation is *even*, if it consists of an even number of transpositions. Otherwise, the signature is -1 and the permutation is *odd*.



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Definition 183 (Inversion, Dt.: Inversion, Fehlstand)

A permutation $\pi \in S_n$ has an *inversion* (i,j) if $\pi(i) > \pi(j)$ for $1 \le i < j \le n$.

Principles of Elementary Counting and Combinatorics

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- Ordered Selection (Variation)
- Unordered Selection (Combination)



Ordered Selection

Definition 184 (Ordered selection without repetition, Dt.: Variation ohne Zurücklegen, Variation ohne Wiederholung)

Let $n \in \mathbb{N}$ and $k \in \mathbb{N}_0$, with $k \leq n$, and A be a finite set of cardinality n. An ordered selection without repetition of k elements from A is a k-tuple

 (a_1, a_2, \dots, a_k) with $a_i \in A$ for $i = 1, 2, \dots, k$ and $a_i \neq a_j$ for $1 \leq i < j \leq k$.



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Let $n \in \mathbb{N}$ and $k \in \mathbb{N}_0$, with $k \leq n$, and A be a finite set of cardinality n. There exist

$$V_k^n := \frac{n!}{(n-k)!}$$

many different ordered selections without repetition of k elements from A.



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- Convention: $V_k^n := 0$ for k > n.
- V_k^n is the number of injective functions from I_k to A.
- Sometimes, V(n,k) is written instead of Vⁿ_k. English-language textbooks often speak of a k-permutation rather than of an ordered selection without repetition

Definition 186 (Ordered selection with repetition, Dt.: Variation mit Zurücklegen, Variation mit Wiederholung)

Let $n \in \mathbb{N}$ and $k \in \mathbb{N}_0$, and A be a finite set of cardinality n. An ordered selection with repetition of k elements from A is a k-tuple

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Let $n \in \mathbb{N}$ and $k \in \mathbb{N}_0$, and A be a finite set of cardinality n. There exist

 $^{r}V_{k}^{n}:=n^{k}$

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• Note: ${}^{r}V_{k}^{n} = |A^{k}|.$

• Sometimes, $V_r(n,k)$ is written instead of ${}^rV_k^n$.



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- Ordered Selection (Variation)
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Definition 188 (Unordered selection without repetition, Dt.: Kombination ohne Zurücklegen, Kombination ohne Wiederholung)

Let $n \in \mathbb{N}$ and $k \in \mathbb{N}_0$, with $k \leq n$, and A be a finite set of cardinality n. An unordered selection without repetition of k elements from A is a set B such that

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- Convention: $C_k^n := 0$ for k > n. Sometimes, C(n, k) is written instead of C_k^n .
- Lemma 189 yields an alternate proof of |P(A)| = 2ⁿ. It also implies that there exist (ⁿ_k) different binary sequences where exactly k elements are 1.

Definition 190 (Unordered selection with repetition, Dt.: Kombination mit Zurücklegen, Kombination mit Wiederholung)

Let $n \in \mathbb{N}$ and $k \in \mathbb{N}_0$, and A be a finite set of cardinality n. An unordered selection with repetition of k elements from A is a k-element multiset, i.e., a set $B \subseteq A$ together with a multiplicity function, mult: $A \to \mathbb{N}_0$, such that

 $\operatorname{mult}(a) = 0$ for all $a \in A \setminus B$ and $\operatorname{mult}(b) > 0$ for all $b \in B$ and $\sum \operatorname{mult}(b) = k$.



b∈B

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Lemma 191

Let $n \in \mathbb{N}$ and $k \in \mathbb{N}_0$, and A be a finite set of cardinality n. There exist

$${}^{r}C_{k}^{n}:=\binom{n+k-1}{k}$$

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• Sometimes, $C_r(n,k)$ is written instead of C_k^n .



Proofs of Lemmas 185–191

Proof of Lemma 185: We have *n* options for a_1 , leaving n - 1 options for a_2 , etc. Thus, we have $n \cdot (n - 1) \cdot \ldots \cdot (n - k + 1) = \frac{n!}{(n-k)!}$ options.

Proof of Lemma 187: We have *n* options for every selection. Thus, we have n^k options in total.

Proof of Lemma 189: We know that $V_k^n = \frac{n!}{(n-k)!}$. There are k! many different ordered selections that correspond to the same unordered selection. Thus, $C_k^n = V_k^n/k! = \frac{n!}{(n-k)!k!} = \binom{n}{k}$.



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Proof of Lemma 191: Let a_1, \ldots, a_n be the *n* elements of *A*, and $k \in \mathbb{N}_0$. We encode such an unordered selection with repetition of *k* elements from *A* as a sequence of length n + k - 1 of *k* crosses × which are separated by n - 1 vertical bars |, where *i* crosses between the *j*-th vertical bar and the (j + 1)-st vertical bar, for $1 \le j \le n - 2$, indicate that element a_{j+1} was chosen with multiplicity *i*.



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$$C_k^{n+k-1} = \binom{n+k-1}{k}$$

ways to choose the positions of the k crosses within this sequence.



Real-World Application: Elementary Probability

• What is the probability to win in the Austrian "6-aus-45" lottery after choosing one set of six numbers?



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- As usual, we define the probability of an event among (finitely many) equally-likely outcomes as the number of favorable outcomes divided by the total number of possible outcomes.



- What is the probability to win in the Austrian "6-aus-45" lottery after choosing one set of six numbers?
- As usual, we define the probability of an event among (finitely many) equally-likely outcomes as the number of favorable outcomes divided by the total number of possible outcomes.
- Assuming that the lottery is fair and, thus, that all combinations are equally likely to win, we get

$$\frac{1}{C_6^{45}} = \frac{1}{\binom{45}{6}} = \frac{1}{8\,145\,060} \approx 1.22774 \cdot 10^{-7}$$

as probability for having all six numbers right.



A standard deck of cards contains 52 cards grouped into four suits (Dt.: Farben)
 — diamonds (Dt.: Schelle, Karo), clubs (Dt.: Eichel, Kreuz), hearts (Dt.: Herz), and spades (Dt.: Laub, Pik) — with 13 cards in each suit (ace, 2, 3, 4, 5, 6, 7, 8, 9, 10, jack, queen, king).



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- What is the probability that all hearts appear in consecutive (but arbitrary) order after a decent shuffling of the deck?
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- There are 40! different permutations of the block of 13 hearts and the other 39 cards, and 13! many permutations of the 13 hearts within that block.



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- There are 52! different permutations of the 52 cards.
- There are 40! different permutations of the block of 13 hearts and the other 39 cards, and 13! many permutations of the 13 hearts within that block.
- Hence, the probability that all hearts are consecutive is given by

 $\frac{40!\cdot 13!}{52!}\approx 6.29908\cdot 10^{-11}.$



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Complexity Analysis and Recurrence Relations

- Growth Rates
- Bachmann-Landau (Asymptotic) Notation
- Recurrence Relations
- Master Theorem



Complexity Analysis and Recurrence Relations

- Growth Rates
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• Algorithms/codes tend to process inputs of small sizes instantaneously. Therefore we are most interested in how an algorithm performs as the input size *n* gets large: *asymptotic complexity analysis*.



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- Determine the dominating term in the complexity function it gives the order of magnitude of the asymptotic behavior.

1, $\log n$, $\log^2 n$, \sqrt{n} , n, $n \log n$, $n \log^2 n$, $n^{\frac{7}{6}}$, n^2 , n^3 , ..., 2^n , 3^n , $2^{(2^n)}$, ...



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Convention regarding logarithms

In this course, log *n* will always denote the logarithm of *n* to the base 2, i.e., $\log n := \log_2 n$.

• Recall that
$$\log_{\alpha} n = \frac{1}{\log_2 \alpha} \log_2 n$$
.



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 $g(n) = 9n + 20 \leq 9n + n = 10n = 10f(n),$



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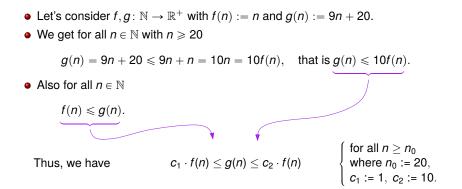
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Thus, we have

 $c_1 \cdot f(n) \leq g(n) \leq c_2 \cdot f(n)$ $\begin{cases} \text{for all } n \geq n_0 \\ \text{where } n_0 \coloneqq 20, \\ c_1 \coloneqq 1, \ c_2 \coloneqq 10. \end{cases}$ $g \text{ grows at most as fast as } c_2 \cdot f$ f is an asymptotic upper bound on g $\text{we'll say that } g \in O(f)$



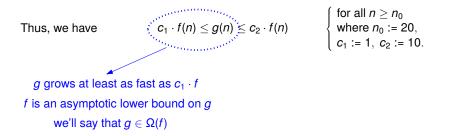
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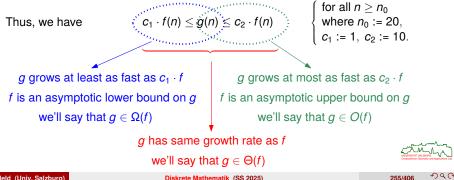


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Diskrete Mathematik (SS 2025)

Complexity Analysis and Recurrence Relations

- Growth Rates
- Bachmann-Landau (Asymptotic) Notation
 - Bachmann-Landau Symbols
 - Limit of a Sequence
 - Basic Facts
 - Conditional Asymptotic Notation and Smoothness Rule
- Recurrence Relations
- Master Theorem



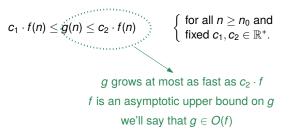
$$c_1 \cdot f(n) \leq g(n) \leq c_2 \cdot f(n) \qquad \left\{ \begin{array}{l} \text{for all } n \geq n_0 \text{ and} \\ \text{fixed } c_1, c_2 \in \mathbb{R}^+. \end{array} \right.$$

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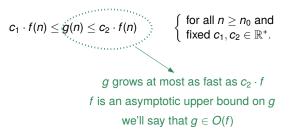


Definition 192 (Big-O, Dt.: Groß-O)

Let $f : \mathbb{N} \to \mathbb{R}^+$. Then the set O(f) is defined as

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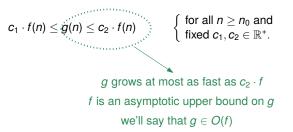
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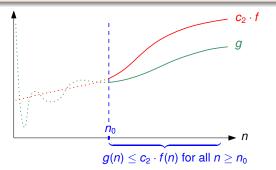
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- Some authors prefer to use the symbol \mathcal{O} instead of O.
- Note: O(f) is a set of functions! Definitions of the form $O(f(n)) := \{g : \mathbb{N} \to \mathbb{R}^+ | \exists c_2 \in \mathbb{R}^+ \exists n_0 \in \mathbb{N} \forall n \ge n_0 \quad g(n) \le c_2 \cdot f(n) \}$ are a (wide-spread) formal nonsense.

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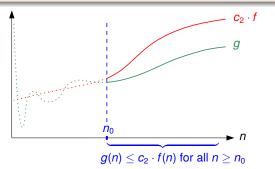




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• Equivalent definition used by some authors:

$$O(f) := \left\{g \colon \mathbb{N} \to \mathbb{R}^+ \mid \exists c_2 \in \mathbb{R}^+ \exists n_0 \in \mathbb{N} \quad \forall n \ge n_0 \quad \frac{g(n)}{f(n)} \le c_2 \right\}.$$

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Diskrete Mathematik (SS 2025)

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Compute(i, j)
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- We get

$$g(n) = n + (n - 1) + \dots + 2 + 1$$

= $\frac{n(n + 1)}{2} = \frac{1}{2}n^2 + \frac{1}{2}n.$



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- Consider $f : \mathbb{N} \to \mathbb{R}^+$ with $f(n) := n^2$.
- Let's compare the growth rates of *f* and *g* when we double *n*:

| n | g (n) | <i>f</i> (<i>n</i>) |
|----|-----------------------|-----------------------|
| 5 | 15 | 25 |
| 10 | 55 | 100 |
| 20 | 210 | 400 |
| 40 | 820 | 1600 |
| 80 | 3240 | 6400 |



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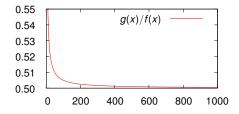
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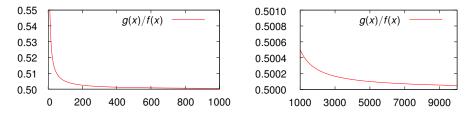
• Doubling *n* causes both *f*(*n*) and *g*(*n*) to (roughly) quadruple!





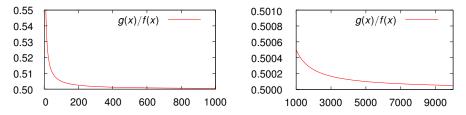






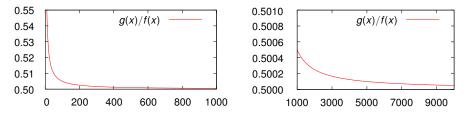


• We plot the growth ratio $\frac{g(n)}{f(n)}$ for $f,g: \mathbb{N} \to \mathbb{R}^+$ with $f(n) := n^2$ and $g(n) := \frac{1}{2}n^2 + \frac{1}{2}n$.



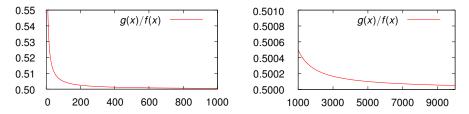
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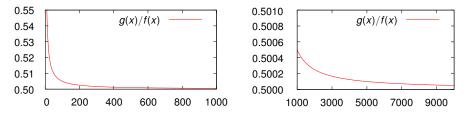
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- The plots also suggest $\frac{g(n)}{f(n)} \ge \frac{1}{2}$, which would imply $g \in \Omega(f)$.
- Hence $g(n) \approx \frac{1}{2}f(n)$, which would imply $g \in \Theta(f)$.



Asymptotic Notation: Big-Omega

 $c_1 \cdot f(n) \leq g(n) \leq c_2 \cdot f(n)$ *g* grows at least as fast as $c_1 \cdot f$ *f* is an asymptotic lower bound on *g* we'll say that $g \in \Omega(f)$

 $\begin{cases} \text{ for all } n \ge n_0 \text{ and} \\ \text{ fixed } c_1, c_2 \in \mathbb{R}^+. \end{cases}$



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$$g \text{ grows at least as fast as } c_1 \cdot f$$

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Definition 193 (Big-Omega, Dt.: Groß-Omega)

Let $f: \mathbb{N} \to \mathbb{R}^+$. Then the set $\Omega(f)$ is defined as

 $\Omega(f) := \{g \colon \mathbb{N} \to \mathbb{R}^+ \mid \exists c_1 \in \mathbb{R}^+ \exists n_0 \in \mathbb{N} \forall n \ge n_0 \quad c_1 \cdot f(n) \le g(n) \}.$



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Equivalently,

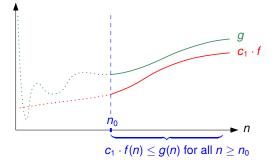
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(

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Asymptotic Notation: Big-Theta

$$c_1 \cdot f(n) \leq g(n) \leq c_2 \cdot f(n)$$

 $\left\{ \begin{array}{l} \text{for all } n \geq n_0 \text{ and} \\ \text{fixed } c_1, c_2 \in \mathbb{R}^+. \end{array} \right.$

g has same growth rate as f we'll say that $g \in \Theta(f)$



Asymptotic Notation: Big-Theta

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Let $f : \mathbb{N} \to \mathbb{R}^+$. Then the set $\Theta(f)$ is defined as

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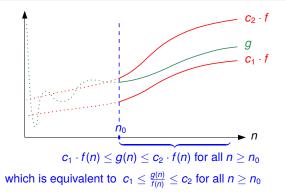


Graphical Illustration of $\Theta(f)$

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$$\begin{aligned} \Theta(f) &:= \left\{ g \colon \mathbb{N} \to \mathbb{R}^+ \mid \exists c_1, \, c_2 \in \mathbb{R}^+ \quad \exists n_0 \in \mathbb{N} \quad \forall n \ge n_0 \\ c_1 \cdot f(n) \leqslant g(n) \leqslant c_2 \cdot f(n) \right\}. \end{aligned}$$





Diskrete Mathematik (SS 2025)

• We prove
$$g \in \Theta(f)$$
 for $f(n) := n^2$ and $g(n) := \frac{1}{2}n^2 + \frac{1}{2}n$.



- We prove $g \in \Theta(f)$ for $f(n) := n^2$ and $g(n) := \frac{1}{2}n^2 + \frac{1}{2}n$. *Proof*:
 - We get, for all $n \in \mathbb{N}$,

$$g(n) = \frac{1}{2}n^2 + \frac{1}{2}n \leq \frac{1}{2}n^2 + \frac{1}{2}n^2 = n^2 = f(n),$$



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• Thus,
$$g \in O(f)$$
 with $c_2 := 1$ and $n_0 := 1$.



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 - We get, for all $n \in \mathbb{N}$,

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• Thus, $g \in \Omega(f)$ with $c_1 := \frac{1}{2}$ and $n_0 := 1$. Def. 194 or Lemma 201 yield $g \in \Theta(f)$.



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 - We get, for all $n \in \mathbb{N}$,

$$g(n) = \frac{1}{2}n^2 + \frac{1}{2}n \leq \frac{1}{2}n^2 + \frac{1}{2}n^2 = n^2 = f(n), \text{ that is } g(n) \leq f(n).$$

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• Now we prove $g \in \Omega(f)$ and get, again for all $n \in \mathbb{N}$,

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• Thus, $g \in \Omega(f)$ with $c_1 := \frac{1}{2}$ and $n_0 := 1$. Def. 194 or Lemma 201 yield $g \in \Theta(f)$.

Don't be overly zealous!

There is no need to try to obtain the "best-possible" values for n_0 and c_1, c_2 !

Let $f : \mathbb{N} \to \mathbb{R}^+$. Then the set o(f) is defined as

$$o(f) := \{g \colon \mathbb{N} \to \mathbb{R}^+ \mid \forall c \in \mathbb{R}^+ \exists n_0 \in \mathbb{N} \forall n \ge n_0 \qquad g(n) \leqslant c \cdot f(n) \}.$$



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Mind the difference

$$\begin{aligned} O(f) &:= \left\{ g \colon \mathbb{N} \to \mathbb{R}^+ \mid \exists c \in \mathbb{R}^+ \exists n_0 \in \mathbb{N} \quad \forall n \ge n_0 \qquad g(n) \leqslant c \cdot f(n) \right\} \\ o(f) &:= \left\{ g \colon \mathbb{N} \to \mathbb{R}^+ \mid \forall c \in \mathbb{R}^+ \exists n_0 \in \mathbb{N} \quad \forall n \ge n_0 \qquad g(n) \leqslant c \cdot f(n) \right\} \end{aligned}$$



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| $O(f) := \{g \colon \mathbb{N} \to \mathbb{R}^+ \mid$ | $\exists c \in \mathbb{R}^+$ | $\exists n_0 \in \mathbb{N}$ | $\forall n \ge n_0$ | $g(n) \leq c \cdot f(n)$ |
|--|---|------------------------------|---------------------|-----------------------------|
| $o(f) := \{g \colon \mathbb{N} \to \mathbb{R}^+ \mid f \}$ | $\forall \boldsymbol{c} \in \mathbb{R}^+$ | $\exists n_0 \in \mathbb{N}$ | $\forall n \ge n_0$ | $g(n) \leq c \cdot f(n) \}$ |

• Similarly, $\omega(f)$ can be defined relative to $\Omega(f)$.



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| $O(f) := \{g \colon \mathbb{N} \to \mathbb{R}^+ \mid$ | $\exists c \in \mathbb{R}^+$ | $\exists n_0 \in \mathbb{N}$ | $\forall n \ge n_0$ | $g(n) \leqslant c \cdot f(n) \}$ |
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- Similarly, $\omega(f)$ can be defined relative to $\Omega(f)$.
- It is trivial to extend Definitions 192–195 such that \mathbb{N}_0 rather than \mathbb{N} is taken as the domain.
- We can also replace the codomain \mathbb{R}^+ by \mathbb{R}^+_0 (or even \mathbb{R}) provided that all functions are eventually positive.
- The same comments apply to the subsequent slides.



Definition 196 (Sequence, Dt.: Folge)

A (real) sequence is a function from \mathbb{N} (or \mathbb{N}_0) to \mathbb{R} . For $x : \mathbb{N} \to \mathbb{R}$ it is common to write the sequence as $(x_n)_{n \in \mathbb{N}}$ or $\langle x_n \rangle_{n \in \mathbb{N}}$, or simply (x_n) or $\langle x_n \rangle$.



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Definition 197 (Limit, Dt. Grenzwert)

The value $\bar{x} \in \mathbb{R}$ is the limit of the (real) sequence (x_n) , denoted by $\lim_{n\to\infty} x_n = \bar{x}$, if

 $\forall \varepsilon \in \mathbb{R}^+ \ \exists n_0 \in \mathbb{N} \ \forall n \ge n_0 \quad |x_n - \bar{x}| < \varepsilon.$



Definition 196 (Sequence, Dt.: Folge)

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Lemma 198

If $z_n = x_n + y_n$ for three sequences $(x_n), (y_n), (z_n)$ and if $\lim_{n\to\infty} x_n$ and $\lim_{n\to\infty} y_n$ exist, then $\lim_{n\to\infty} z_n$ exists and we have $\lim_{n\to\infty} z_n = \lim_{n\to\infty} x_n + \lim_{n\to\infty} y_n$.



Theorem 199 (Squeeze theorem, Dt.: Einschnürungssatz)

Consider three real sequences $(x_n), (y_n), (z_n)$ and suppose that $x_n \leq y_n \leq z_n$ for all $n \geq n_0$ for some $n_0 \in \mathbb{N}$. If the limits of (x_n) and (z_n) exist such that

 $\lim_{n\to\infty}x_n=\lim_{n\to\infty}z_n,$

then the limit of (y_n) exists with

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 $\lim_{n\to\infty} x_n = \lim_{n\to\infty} y_n = \lim_{n\to\infty} z_n.$

• For $z_n := \frac{8}{n}$ it is easy to see that $\lim_{n\to\infty} z_n = 0$.



Theorem 199 (Squeeze theorem, Dt.: Einschnürungssatz)

Consider three real sequences (x_n) , (y_n) , (z_n) and suppose that $x_n \le y_n \le z_n$ for all $n \ge n_0$ for some $n_0 \in \mathbb{N}$. If the limits of (x_n) and (z_n) exist such that

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Now consider the following sequences:

$$x_n := 0$$
 $y_n := \frac{\log n + 7\sqrt{n} - 10}{n^2}$ $z_n := \frac{8}{n}$



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Now consider the following sequences:

$$x_n := 0$$
 $y_n := \frac{\log n + 7\sqrt{n} - 10}{n^2}$ $z_n := \frac{8}{n}$

• We have for all $n \in \mathbb{N} \setminus \{1, 2, 3\}$

$$x_n \leqslant y_n \leqslant z_n$$
 and $\lim_{n\to\infty} x_n = 0 = \lim_{n\to\infty} z_n$.

Thus, $\lim_{n\to\infty} y_n = 0$.

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• The following theorem (by Guillaume de l'Hôpital, 1661–1704) allows to handle limits that involve indeterminate terms of the form

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Consider two real functions f and g, and a real value c. If

$$Iim_{x\to c} f(x) = 0 = Iim_{x\to c} g(x) \text{ or } Iim_{x\to c} f(x) = \pm \infty = Iim_{x\to c} g(x),$$



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$$g'(x) \neq 0 \text{ for all } x \in I \setminus \{c\}, \text{ and i}$$

$$Iim_{x\to c} \frac{f'(x)}{g'(x)} exists,$$



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2 f and g are differentiable in an open interval I with $c \in I$, except possibly at c itself,

$$Iim_{x\to c} \frac{f'(x)}{g'(x)} exists,$$

then

$$\lim_{x \to c} \frac{f(x)}{g(x)} = \lim_{x \to c} \frac{f'(x)}{g'(x)}$$

Lemma 201

Let $f_1, f_2, g_1, g_2 : \mathbb{N} \to \mathbb{R}^+$, and $c \in \mathbb{R}^+$. Then the following relations hold: **(** $g_1 \in O(f_1) \land g_2 \in O(f_2)$) $\Rightarrow g_1 + g_2 \in O(f_1 + f_2)$



Lemma 201

- $(g_1 \in O(f_1) \land g_2 \in O(f_2)) \Rightarrow g_1 + g_2 \in O(f_1 + f_2)$
- $(g_1 \in O(f_1) \land g_2 \in O(f_2)) \Rightarrow g_1 \cdot g_2 \in O(f_1 \cdot f_2)$



Lemma 201

- $\bigcirc (g_1 \in O(f_1) \land g_2 \in O(f_2)) \Rightarrow g_1 + g_2 \in O(f_1 + f_2)$
- $(g_1 \in O(f_1) \land g_2 \in O(f_2)) \Rightarrow g_1 \cdot g_2 \in O(f_1 \cdot f_2)$



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- $(g_1 \in O(f_1) \land g_2 \in O(f_2)) \Rightarrow g_1 + g_2 \in O(f_1 + f_2)$
- $(g_1 \in O(f_1) \land g_2 \in O(f_2)) \Rightarrow g_1 \cdot g_2 \in O(f_1 \cdot f_2)$
- $O(\boldsymbol{c} \cdot \boldsymbol{f}_1) = \boldsymbol{O}(\boldsymbol{f}_1)$



Lemma 201

- $(g_1 \in O(f_1) \land g_2 \in O(f_2)) \Rightarrow g_1 + g_2 \in O(f_1 + f_2)$
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- $O(\boldsymbol{c} \cdot \boldsymbol{f}_1) = O(\boldsymbol{f}_1)$



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- $(g_1 \in O(f_1) \land g_2 \in O(f_2)) \Rightarrow g_1 + g_2 \in O(f_1 + f_2)$
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- $O(\boldsymbol{c} \cdot \boldsymbol{f}_1) = O(\boldsymbol{f}_1)$
- $\Theta(f_1) = O(f_1) \cap \Omega(f_1)$



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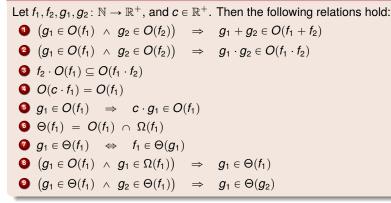
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- $(g_1 \in O(f_1) \land g_1 \in \Omega(f_1)) \Rightarrow g_1 \in \Theta(f_1)$



Lemma 201





Lemma 202

Let $f, g: \mathbb{N} \to \mathbb{R}^+$ and $c \in \mathbb{R}^+$. Then:

$$\lim_{n\to\infty}\frac{g(n)}{f(n)}=c\quad\Rightarrow\quad g\in\Theta(f).$$



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Let $f, g \colon \mathbb{N} \to \mathbb{R}^+$ and $c \in \mathbb{R}^+$. Then:

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• For example, let $f, g, h: \mathbb{N} \to \mathbb{R}^+$ with $f(n) := n^2 - 7n$, $g(n) := 3n^2 + 5n\sqrt{n}$ and $h(n) := n^2$.



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$$\lim_{n \to \infty} \frac{f(n)}{h(n)} = \lim_{n \to \infty} \frac{n^2 - 7n}{n^2} = \lim_{n \to \infty} \left(1 - \frac{7}{n}\right) = 1$$



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Sac

Lemma 202

Let $f, g \colon \mathbb{N} \to \mathbb{R}^+$ and $c \in \mathbb{R}^+$. Then:

$$\lim_{n\to\infty}\frac{g(n)}{f(n)}=c \quad \Rightarrow \quad g\in\Theta(f).$$

and

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$$\lim_{n \to \infty} \frac{g(n)}{h(n)} = \lim_{n \to \infty} \frac{3n^2 + 5n\sqrt{n}}{n^2} = \lim_{n \to \infty} \left(3 + \frac{5}{\sqrt{n}}\right) = 3$$



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Diskrete Mathematik (SS 2025)

Asymptotic Notation: Wide-spread Notational Abuse

• It is convenient to be a bit sloppy and write, e.g.,

$$g(n) = O(n^2)$$
 or $g \in O(n^2)$

rather than to resort to the λ -quantifier and write $g \in O(\lambda n.n^2)$, or

$$g \in O(f)$$
 with $f : \mathbb{N} \to \mathbb{R}^+$, $n \mapsto n^2$.



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 with $f : \mathbb{N} \to \mathbb{R}^+$, $n \mapsto n^2$.

• Similarly,

$$g(n) = h(n) + O(n^3)$$

means

 $|g-h| \in O(f)$ with $f: \mathbb{N} \to \mathbb{R}^+$, $n \mapsto n^3$.



Asymptotic Notation: Wide-spread Notational Abuse

• It is convenient to be a bit sloppy and write, e.g.,

$$g(n) = O(n^2)$$
 or $g \in O(n^2)$

rather than to resort to the λ -quantifier and write $g \in O(\lambda n.n^2)$, or

$$g \in O(f)$$
 with $f : \mathbb{N} \to \mathbb{R}^+$, $n \mapsto n^2$.

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$$g(n) = h(n) + O(n^3)$$

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$$|g-h| \in O(f)$$
 with $f: \mathbb{N} \to \mathbb{R}^+$, $n \mapsto n^3$.

• Furthermore,

$$g(n) = n^{O(1)}$$

indicates that

$$g \in O(f)$$
 with $f: \mathbb{N} \to \mathbb{R}^+, n \mapsto n^c$

for some constant $c \in \mathbb{R}^+$.



Asymptotic Notation: Wide-spread Notational Abuse — Caveats!

Warning

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- So, keep in mind that an *is-element-of* or *subset relation* is meant even if an equality sign is used!
- Unfortunately, several textbooks are fuzzy about this important distinction



Consider a function $f : \mathbb{N} \to \mathbb{R}^+$

 $O(f) := \{g \colon \mathbb{N} \to \mathbb{R}^+ \mid \exists c \in \mathbb{R}^+ \exists n_0 \in \mathbb{N} \quad \forall n \ge n_0 :$

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



Consider a function $f : \mathbb{N} \to \mathbb{R}^+$ and a predicate $P : \mathbb{N} \to \{F, T\}$.

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$$P(n) :\Leftrightarrow n \equiv_2 0$$
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Diskrete Mathematik (SS 2025)

Consider a function $f : \mathbb{N} \to \mathbb{R}^+$ and a predicate $P : \mathbb{N} \to \{F, T\}$.

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$$\begin{split} \Theta(f \mid P) &:= \left\{ g \colon \mathbb{N} \to \mathbb{R}^+ \mid \exists c_1, c_2 \in \mathbb{R}^+ \exists n_0 \in \mathbb{N} \quad \forall n \ge n_0 : \\ P(n) &\Rightarrow c_1 \cdot f(n) \le g(n) \le c_2 \cdot f(n) \right\}. \end{split}$$

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A function $f : \mathbb{N} \to \mathbb{R}^+$ is eventually non-decreasing exactly if

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Definition 205 (b-smooth, Dt.: b-glatt)

A function $f : \mathbb{N} \to \mathbb{R}^+$ is *b-smooth* for some integer $b \ge 2$ exactly if *f* is eventually non-decreasing and if

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Lemma 207

If $f : \mathbb{N} \to \mathbb{R}^+$ is b'-smooth for some integer $b' \ge 2$ then it is smooth.

Theorem 208 (Smoothness Rule)

Let $f, g \colon \mathbb{N} \to \mathbb{R}^+$, and consider an integer $b \ge 2$.



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- Similarly for $\Omega(f)$ and $\Theta(f)$.
- Again, it is trivial to extend the definitions and lemmas such that \mathbb{N}_0 rather than \mathbb{N} is taken as the base set. Similarly, we can replace \mathbb{R}^+ by \mathbb{R}^+_0 or even by \mathbb{R} provided that all functions are eventually positive.
- The same comments apply to the subsequent slides.



• For
$$a, b \in \mathbb{R}^+_0$$
 we define $g \colon \mathbb{N} \to \mathbb{R}^+_0$ as

$$g(n) := \begin{cases} a & \text{if } n = 1, \\ 4 \cdot g\left(\left\lceil \frac{n}{2} \right\rceil\right) + b \cdot n & \text{otherwise.} \end{cases}$$



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• Standard application in computer science: Solving the recurrence relation

$$T(n) = T\left(\left\lceil \frac{n}{2} \right\rceil\right) + T\left(\left\lfloor \frac{n}{2} \right\rfloor\right) + b \cdot n,$$

e.g., as derived when analyzing the complexity of merge sort.



6 Complexity Analysis and Recurrence Relations

- Growth Rates
- Bachmann-Landau (Asymptotic) Notation
- Recurrence Relations
 - Heuristics for Solving Recurrences
 - Solving Linear Recurrence Relations
- Master Theorem



• Sample sequence $t: \mathbb{N}_0 \to \mathbb{R}$: (1, 2, 4, 8, 16, 32, 64, 128, 256, ...)



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A *recurrence relation* for a sequence *t* is an equation that relates elements of *t*. It is of order *k*, for some $k \in \mathbb{N}$, if t_n can be expressed in terms of *n* and $t_{n-1}, t_{n-2}, \ldots, t_{n-k}$, i.e., if t_n is of the form $t_n = f(t_{n-1}, t_{n-2}, \ldots, t_{n-k}, n)$ for $f : \mathbb{R}^k \times \mathbb{N} \to \mathbb{R}$ (or for $f : \mathbb{R}^k \times \mathbb{N}_0 \to \mathbb{R}$).



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• Recurrence relation (of order 1) for the sample sequence given above:

$$t_n := \begin{cases} 1 & \text{if } n = 0, \\ 2 \cdot t_{n-1} & \text{if } n > 0. \end{cases}$$



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Note

We will freely mix the notations t_k and t(k) for denoting the *k*-th element of a sequence $(t_n)_{n \in \mathbb{N}}$ or $(t_n)_{n \in \mathbb{N}_0}$.



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- Also according to legend, the priests apply a recursive algorithm, thereby moving
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- A solution of this recurrence relation tells us when life on Earth might end ...
- So, is it already time for an apocalyptic mood?
- We start with heuristics for solving recurrence relations.



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 - First "guess" a solution.
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Note

All heuristics require induction to prove that the result obtained is indeed correct!



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$$\stackrel{a:=2}{=} 2(n^2 + 2n + 1)$$



- Solve the recurrence relation $t_n = t_{n-1} + n$, with $t_0 := 0$.
- Guess: $t \in O(f)$ for $f(n) := n^2$.
- Our guess could be verified by showing t_n ≤ a ⋅ n² for all n ∈ N₀ for a suitable (but yet unknown) a ∈ ℝ⁺.
- If we assume $t_n \leq a \cdot n^2$ then we get

$$t_{n+1} = t_n + (n+1)$$

$$\leq a \cdot n^2 + (n+1)$$

$$\leq a \cdot n^2 + 4n + 2$$

$$= 2(\frac{a}{2} \cdot n^2 + 2n + 1)$$

$$\stackrel{a := 2}{=} 2(n^2 + 2n + 1)$$

$$= 2(n+1)^2.$$

• Now use standard induction to show that $t_n \leq 2n^2$ is indeed correct for all $n \in \mathbb{N}_0$.

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Heuristics for Solving Recurrences: Cascading

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$$\vdots$$

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$$\frac{t_{1} = t_{0} + 1}{t_{n} = t_{0} + 1 + 2 + \dots + (n-2) + (n-1) + n}$$



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$$= 0 + 1 + 2 + \dots + (n-2) + (n-1) + n$$

This indicates that

$$t_n=\sum_{i=0}^n i=\frac{n(n+1)}{2}\in \Theta(n^2),$$

which is proved by induction.

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- Solve the recurrence relation $t_n = t_{n-1} + n$, with $t_0 := 0$.
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= $(t_{n-4} + (n-3)) + ((n-2) + (n-1) + n)$



- Solve the recurrence relation $t_n = t_{n-1} + n$, with $t_0 := 0$.
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$$t_{n} = t_{n-1} + n$$

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$$= (t_{n-4} + (n-3)) + ((n-2) + (n-1) + n)$$

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• Again, this indicates that

$$t_n=\sum_{i=0}^n i=\frac{n(n+1)}{2}\in\Theta(n^2),$$

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= $2^{2}(2^{1}T(n-3) + 2^{0}) + 2^{1} + 2^{0} = 2^{3}T(n-3) + 2^{2} + 2^{1} + 2^{0}$
:
= $2^{n-1}T(n-(n-1)) + 2^{n-2} + \dots + 2^{2} + 2^{1} + 2^{0}$
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= $2^{n} - 1$



• We have the Tower-of-Hanoi recurrence relation

T(n) = 2T(n-1) + 1 with T(1) := 1.

• Iteration yields the following identities:

$$T(n) = 2T(n-1) + 1 = 2^{1}T(n-1) + 2^{0}$$

= $2(2^{1}T(n-2) + 2^{0}) + 2^{0} = 2^{2}T(n-2) + 2^{1} + 2^{0}$
= $2^{2}(2^{1}T(n-3) + 2^{0}) + 2^{1} + 2^{0} = 2^{3}T(n-3) + 2^{2} + 2^{1} + 2^{0}$
:
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= $2^{n-1} + 2^{n-2} + \dots + 2^{2} + 2^{1} + 2^{0}$
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 Hence, if the priests manage to move one disk per second then we would have to expect the end of Earth 2⁶⁴ – 1 seconds after they started, i.e., roughly within 5 · 10¹¹ years ...

Definition 210 (Homogeneous recurrence, Dt.: homogene Rekurrenz)

A recurrence relation of order *k* is *homogeneous* if it is satisfied by the zero sequence.



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A homogeneous recurrence relation of order k is *linear* if $t_n = \sum_{i=1}^k a_i(n) \cdot t_{n-i}$, where $a_i : \mathbb{N} \to \mathbb{R}$ for i = 1, 2, ..., k.



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A linear homogeneous recurrence relation of order *k* has *constant coefficients* if $t_n = \sum_{i=1}^k a_i \cdot t_{n-i}$, where $a_1, a_2, \ldots, a_k \in \mathbb{R}$.



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• E.g.,
$$t_n := 2 \cdot t_{n-1} + 3 \cdot t_{n-2}$$
.



Lemma 213

Consider the recurrence relation $a_0t_n + a_1t_{n-1} + \cdots + a_kt_{n-k} = 0$, with $a_i \in \mathbb{R}$. If (f_n) and (g_n) satisfy the recurrence relation then $(\alpha f_n + \beta g_n)$ satisfies the recurrence relation for all $\alpha, \beta \in \mathbb{R}$.



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Proof: Suppose that

$$a_0 f_n + a_1 f_{n-1} + \dots + a_k f_{n-k} = \sum_{i=0}^k a_i f_{n-i} = 0$$
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for all $n \ge k$. Let $\alpha, \beta \in \mathbb{R}$ arbitrary but fixed and consider $(\alpha f_n + \beta g_n)$. We get

$$\sum_{i=0}^k a_i (\alpha f_{n-i} + \beta g_{n-i}) = \alpha \sum_{i=0}^k a_i f_{n-i} + \beta \sum_{i=0}^k a_i g_{n-i} = 0.$$



- So, consider $a_0t_n + a_1t_{n-1} + \cdots + a_kt_{n-k} = 0$
- Guess $t_n = x^n$ for some unknown $x \in \mathbb{R}$.



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- So, consider $a_0t_n + a_1t_{n-1} + \cdots + a_kt_{n-k} = 0$
- Guess $t_n = x^n$ for some unknown $x \in \mathbb{R}$.
- Then $a_0x^n + a_1x^{n-1} + \cdots + a_kx^{n-k} = 0$.
- Further $x^{n-k}(a_0x^k + a_1x^{k-1} + \dots + a_k) = 0.$



• So, consider
$$a_0t_n + a_1t_{n-1} + \cdots + a_kt_{n-k} = 0$$

• Guess $t_n = x^n$ for some unknown $x \in \mathbb{R}$.

• Then
$$a_0x^n + a_1x^{n-1} + \cdots + a_kx^{n-k} = 0.$$

- Further $x^{n-k}(a_0x^k + a_1x^{k-1} + \dots + a_k) = 0.$
- If we ignore the trivial solution *x* := 0 then we get

 $a_0x^k+a_1x^{k-1}+\cdots+a_k=0$

as the so-called characteristic equation of the recurrence relation

$$a_0t_n + a_1t_{n-1} + \cdots + a_kt_{n-k} = 0.$$

• Hence, any root *r* of this equation serves as a partial solution of the recurrence relation, with $t_n := r^n$.



• Suppose that the characteristic equation has k distinct roots r_1, \ldots, r_k such that all roots are real numbers. I.e., the characteristic equation is given as

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• The constants *c_i* are determined based on the initial condition(s).



• Consider the *Fibonacci* sequence (over N₀)

$$F_n := \begin{cases} 0 & \text{if } n = 0, \\ 1 & \text{if } n = 1, \\ F_{n-1} + F_{n-2} & \text{if } n \ge 2. \end{cases}$$



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• This characteristic equation has the roots

$$r_1 := \frac{1+\sqrt{5}}{2}$$
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- This yields

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$$n := 0: \qquad F_0 = 0 = c_1 + c_2$$

$$n := 1: \qquad F_1 = 1 = c_1 \cdot \frac{1 + \sqrt{5}}{2} + c_2 \cdot \frac{1 - \sqrt{5}}{2}$$



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By solving this linear system we obtain c₁ = −c₂ = ¹/_{√5}.
Hence,

$$F_n = \frac{1}{\sqrt{5}} \cdot \left(\frac{1+\sqrt{5}}{2}\right)^n - \frac{1}{\sqrt{5}} \cdot \left(\frac{1-\sqrt{5}}{2}\right)^n.$$



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• *Multiple roots*: Suppose that the characteristic equation has *s* distinct roots r_1, \ldots, r_s of multiplicities m_1, \ldots, m_s such that all roots are real numbers. I.e., the characteristic equation is given as

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• Then we have

$$t_n = \sum_{i=1}^s \sum_{j=0}^{m_i-1} c_{ij} \cdot n^j \cdot r_i^n,$$

for constants $c_{ij} \in \mathbb{R}$.



• *Multiple roots*: Suppose that the characteristic equation has *s* distinct roots r_1, \ldots, r_s of multiplicities m_1, \ldots, m_s such that all roots are real numbers. I.e., the characteristic equation is given as

$$\prod_{i=1}^{s} (x-r_i)^{m_i} = 0.$$

• Then we have

$$t_n = \sum_{i=1}^s \sum_{j=0}^{m_i-1} c_{ij} \cdot n^j \cdot r_i^n,$$

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• E.g., for the characteristic equation $(x - 1) \cdot (x - 2)^2 = 0$ we have $s = 2, r_1 = 1, r_2 = 2, m_1 = 1, m_2 = 2$,

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• E.g., for the characteristic equation $(x - 1) \cdot (x - 2)^2 = 0$ we have $s = 2, r_1 = 1, r_2 = 2, m_1 = 1, m_2 = 2$, and get

$$t_n = c_{10} \cdot n^0 \cdot 1^n + c_{20} \cdot n^0 \cdot 2^n + c_{21} \cdot n^1 \cdot 2^n = c_{10} + c_{20} \cdot 2^n + c_{21} \cdot n \cdot 2^n.$$



• Assume we have an inhomogeneous recurrence relation of the following form:

$$\mathbf{a}_0 \cdot t_n + \mathbf{a}_1 \cdot t_{n-1} + \cdots + \mathbf{a}_k \cdot t_{n-k} = \mathbf{b}_1^n \cdot \mathbf{p}_1(n) + \mathbf{b}_2^n \cdot \mathbf{p}_2(n) + \cdots + \mathbf{b}_t^n \cdot \mathbf{p}_t(n),$$

where $t \in \mathbb{N}_0$ and b_i is constant and p_i is a polynomial in n of degree $d_i \in \mathbb{N}_0$ for each $1 \leq i \leq t$.



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where $t \in \mathbb{N}_0$ and b_i is constant and p_i is a polynomial in n of degree $d_i \in \mathbb{N}_0$ for each $1 \leq i \leq t$.

• Then the characteristic polynomial is

$$(a_0 \cdot x^k + a_1 \cdot x^{k-1} + \dots + a_k) \cdot \prod_{i=1}^t (x - b_i)^{d_i+1} = 0.$$

• Now proceed as in the homogeneous case.



Theorem 214

Consider the linear inhomogeneous recurrence relation

$$a_0t_n + a_1t_{n-1} + \cdots + a_kt_{n-k} = \sum_{i=1}^l b_i^n \cdot p_i(n),$$

where $t \in \mathbb{N}_0$, and b_i is constant and p_i is a polynomial in n of degree $d_i \in \mathbb{N}_0$ for each $1 \le i \le t$,

DAG

Theorem 214

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where $t \in \mathbb{N}_0$, and b_i is constant and p_i is a polynomial in n of degree $d_i \in \mathbb{N}_0$ for each $1 \le i \le t$, and suppose that its characteristic equation

$$(a_0x^k + a_1x^{k-1} + \dots + a_k) \cdot \prod_{i=1}^{l} (x - b_i)^{d_i+1} = 0$$

has *s* distinct roots r_1, \ldots, r_s of multiplicities m_1, \ldots, m_s such that all roots are real numbers.

Sac

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has *s* distinct roots r_1, \ldots, r_s of multiplicities m_1, \ldots, m_s such that all roots are real numbers. Then the general solution of the recurrence relation is given by

$$t_n = \sum_{i=1}^s \sum_{j=0}^{m_i-1} c_{ij} \cdot n^j \cdot r_i^n,$$

for constants $c_{ij} \in \mathbb{R}$.

Sac

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• Hence, relative to Thm. 214, we get

$$k = 1$$
 $a_0 = 1$ $a_1 = -2$ $t = 2$
 $b_1 = 1$ $p_1(n) = n$ $d_1 = 1$ $b_2 = 2$ $p_2(n) = 1$ $d_2 = 0$.



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$$t_n = c_{10} \cdot n^0 \cdot 1^n + c_{11} \cdot n^1 \cdot 1^n + c_{20} \cdot n^0 \cdot 2^n + c_{21} \cdot n^1 \cdot 2^n$$

= $c_{10} + c_{11} \cdot n + c_{20} \cdot 2^n + c_{21} \cdot n \cdot 2^n$.

So, we know that

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• The constants $c_{10}, c_{11}, c_{20}, c_{21}$ are determined by resorting to the initial conditions:

$$n := 0: \qquad 0 = c_{10} + c_{11} \cdot 0 + c_{20} \cdot 2^0 + c_{21} \cdot 0 \cdot 2^0 = c_{10} + c_{20}$$

$$n := 1: \qquad 3 = c_{10} + c_{11} + 2 \cdot c_{20} + 2 \cdot c_{21}$$

$$n := 2: \qquad 12 = c_{10} + 2 \cdot c_{11} + 4 \cdot c_{20} + 8 \cdot c_{21}$$

$$n := 3: \qquad 35 = c_{10} + 3 \cdot c_{11} + 8 \cdot c_{20} + 24 \cdot c_{21}$$



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• Solving this system of four linear equations for $c_{10}, c_{11}, c_{20}, c_{21}$ yields

$$c_{10} = -2,$$
 $c_{11} = -1,$ $c_{20} = 2,$ $c_{21} = 1.$



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We conclude that

$$t_n = -2 - n + 2 \cdot 2^n + n \cdot 2^n$$
, i.e., $t_n = -2 - n + 2^{n+1} + n \cdot 2^n$.



6 Complexity Analysis and Recurrence Relations

- Growth Rates
- Bachmann-Landau (Asymptotic) Notation
- Recurrence Relations
- Master Theorem



Master Theorem

Theorem 215 (Master theorem, Dt.: Hauptsatz der Laufzeitfunktionen)

Consider constants $c \in \mathbb{R}^+$, $k, n_0 \in \mathbb{N}$ and $a, b \in \mathbb{N}$ with $b \ge 2$, and let $T : \mathbb{N} \to \mathbb{R}_0^+$ be an eventually non-decreasing function such that

$$T(n) = \mathbf{a} \cdot T\left(\frac{n}{b}\right) + \mathbf{c} \cdot n^{k}$$

for all $n \in \mathbb{N}$ with $n \ge n_0$, where we interpret $T(\frac{n}{b})$ as (a combination of) $T(\lceil \frac{n}{b} \rceil)$ or $T(\lfloor \frac{n}{b} \rfloor)$.



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$$T \in \begin{cases} \Theta(n^k) & \text{if } a < b^k, \\ \Theta(n^k \log n) & \text{if } a = b^k, \\ \Theta(n^{\log_b a}) & \text{if } a > b^k. \end{cases}$$



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• E.g., we get $T \in \Theta(n \log n)$ for T defined as follows:

$$T(n) = T\left(\left\lceil \frac{n}{2} \right\rceil\right) + T\left(\left\lfloor \frac{n}{2} \right\rfloor\right) + c \cdot n.$$



Theorem 216

Consider constants $k, n_0 \in \mathbb{N}$ and $a, b \in \mathbb{N}$ with $b \ge 2$, and a function $f \colon \mathbb{N} \to \mathbb{R}_0^+$ with $f \in \Theta(n^k)$. Let $T \colon \mathbb{N} \to \mathbb{R}_0^+$ be an eventually non-decreasing function such that

$$T(n) = \mathbf{a} \cdot T\left(\frac{n}{b}\right) + f(n)$$

for all $n \in \mathbb{N}$ with $n \ge n_0$, where we interpret $T(\frac{n}{b})$ as (a combination of) $T(\lceil \frac{n}{b} \rceil)$ or $T(\lfloor \frac{n}{b} \rfloor)$. Then we have

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Master Theorem (Refined Asymptotic Version)

Theorem 217

Consider constants $n_0 \in \mathbb{N}$ and $a \in \mathbb{N}$, $b \in \mathbb{R}$ with b > 1, and a function $f \colon \mathbb{N} \to \mathbb{R}_0^+$. Let $T \colon \mathbb{N} \to \mathbb{R}_0^+$ be an eventually non-decreasing function such that

$$T(n) = \mathbf{a} \cdot T\left(\frac{n}{b}\right) + f(n)$$

for all $n \in \mathbb{N}$ with $n \ge n_0$, where we interpret $T(\frac{n}{b})$ as (a combination of) $T(\lceil \frac{n}{b} \rceil)$ or $T(\lfloor \frac{n}{b} \rfloor)$.



Theorem 217

Consider constants $n_0 \in \mathbb{N}$ and $a \in \mathbb{N}$, $b \in \mathbb{R}$ with b > 1, and a function $f \colon \mathbb{N} \to \mathbb{R}_0^+$. Let $T \colon \mathbb{N} \to \mathbb{R}_0^+$ be an eventually non-decreasing function such that

$$T(n) = a \cdot T\left(\frac{n}{b}\right) + f(n)$$

for all $n \in \mathbb{N}$ with $n \ge n_0$, where we interpret $T(\frac{n}{b})$ as (a combination of) $T(\lceil \frac{n}{b}\rceil)$ or $T(\lfloor \frac{n}{b}\rceil)$. Then we have

$$T \in \begin{cases} \Theta(f) & \text{if} \begin{cases} f \in \Omega(n^{(\log_b a) + \varepsilon}) \text{ for some } \varepsilon \in \mathbb{R}^+, \\ \text{and if the following regularity condition holds} \\ \text{for some } 0 < s < 1 \text{ and all sufficiently large } n: \\ a \cdot f(n/b) \leqslant s \cdot f(n), \\ \Theta(n^{\log_b a}) & \text{if } f \in \Theta(n^{\log_b a}), \\ \Theta(n^{\log_b a}) & \text{if } f \in O(n^{(\log_b a) - \varepsilon}) \text{ for some } \varepsilon \in \mathbb{R}^+. \end{cases}$$

• This is a simplified version of the Akra-Bazzi Theorem [Akra&Bazzi 1998].



Real-World Application: Analysis of Fast Integer Multiplication

 The standard multiplication of two integers *a*, *b* represented as binary numbers with 2*n* bits each requires Θ(n²) many additions and shifts of bits.



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- [Karatsuba (1960–1963)]: Let

 $(a_{2n-1}a_{2n-2}\cdots a_1a_0)_2$ and $(b_{2n-1}b_{2n-2}\cdots b_1b_0)_2$

be the 2*n*-bit binary representations of *a* and *b*. Hence, $a = \sum_{i=0}^{2n-1} a_i 2^i$ and $b = \sum_{i=0}^{2n-1} b_i 2^i$.



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be the 2*n*-bit binary representations of *a* and *b*. Hence, $a = \sum_{i=0}^{2n-1} a_i 2^i$ and $b = \sum_{i=0}^{2n-1} b_i 2^i$. • We have

 $a \sim 2^n A_1 + A_0$ and $b \sim 2^n B_1 + B_0$

with

$$\begin{aligned} &A_1 := (a_{2n-1}a_{2n-2}\cdots a_{n+1}, a_n)_2, \quad A_0 := (a_{n-1}a_{n-2}\cdots a_1a_0)_2, \\ &B_1 := (b_{2n-1}b_{2n-2}\cdots b_{n+1}, b_n)_2, \quad B_0 := (b_{n-1}b_{n-2}\cdots b_1b_0)_2. \end{aligned}$$



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with

$$A_{1} := (a_{2n-1}a_{2n-2}\cdots a_{n+1}, a_{n})_{2}, \quad A_{0} := (a_{n-1}a_{n-2}\cdots a_{1}a_{0})_{2},$$

$$B_{1} := (b_{2n-1}b_{2n-2}\cdots b_{n+1}, b_{n})_{2}, \quad B_{0} := (b_{n-1}b_{n-2}\cdots b_{1}b_{0})_{2}.$$

We get

$$a \cdot b \sim 2^{2n}A_1 \cdot B_1 + 2^n(A_1 \cdot B_0 + A_0 \cdot B_1) + A_0 \cdot B_0.$$



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which can be rewritten as

 $a \cdot b \sim (2^{2n} + 2^n)A_1 \cdot B_1 + 2^n(A_1 - A_0) \cdot (B_0 - B_1) + (2^n + 1)A_0 \cdot B_0.$



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• Thus, the multiplication of two 2*n*-bit binary numbers can be carried out recursively by computing

• three multiplications of *n*-bit binary numbers, plus

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• Thus, the multiplication of two 2*n*-bit binary numbers can be carried out recursively by computing

three multiplications of n-bit binary numbers, plus

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- Hence, if *T*(*n*) denotes the total number of bit operations used by this recursive algorithm for *n*-bit binary numbers, then

$$T(n) = 3T\left(\frac{n}{2}\right) + f(n)$$
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The asymptotic version of the Master Theorem 216 allows us to conclude that

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 [Schönhage&Strassen (1971), Fürer (2007)]: Faster methods based on Fast Fourier Transform.

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 $T\in \Theta(n^{\log_2 3}), \qquad \text{i.e., that} \ T\in \Theta(n^{1.58496\dots}) \ \text{and, thus,} \ T\in o(n^2).$

- [Schönhage&Strassen (1971), Fürer (2007)]: Faster methods based on Fast Fourier Transform.
- [Harvey&van der Hoeven (2021)]: Achieved $O(n \log n)$.



Graph Theory

- What is a (Directed) Graph?
- Paths
- Trees
- Special Graphs
- Graph Coloring



Graph Theory

- What is a (Directed) Graph?
 - Undirected and Directed Graph
 - Applications: Hasse Diagram and Precedence Graphs
 - Adjacency and Degree
 - Euler's Handshaking Lemma
- Paths
- Trees
- Special Graphs
- Graph Coloring



For $n \in \mathbb{N}$ and $m \in \mathbb{N}_0$, a (simple finite undirected) graph $\mathcal{G} := (V, E)$ with *n* vertices (aka nodes) and *m* edges consists of a vertex set $V := \{v_1, v_2, \ldots, v_n\}$ and an edge set $E := \{e_1, e_2, \ldots, e_m\}$, where $V \cap E = \emptyset$ and each edge is an unordered pair of distinct vertices:



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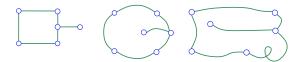


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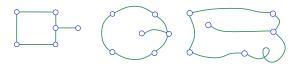
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- If we allow multiple edges between two vertices then we get a multigraph.

- Graphical representation of a graph:
 - Denote the vertices by markers of the same form (circles, dots, squares, ...).
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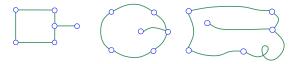


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- Use arrows to denote directed edges.





• Which of the following drawings show simple graphs?



multigraph



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not a simple graph: loop!



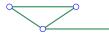
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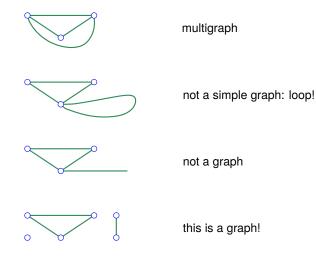
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not a graph



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Definition 219 (Directed graph, Dt.: (schlichter endlicher) gerichteter Graph)

For $n \in \mathbb{N}$ and $m \in \mathbb{N}_0$, a (simple finite) directed graph, or digraph, $\mathcal{G} := (V, E)$ with *n* vertices (aka nodes) and *m* edges consists of a vertex set $V := \{v_1, v_2, \ldots, v_n\}$ and an edge set $E := \{e_1, e_2, \ldots, e_m\}$, where $V \cap E = \emptyset$ and each edge is an ordered pair of distinct vertices:





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- For a digraph, *uv* indicates the edge (*u*, *v*), i.e., an edge where *u* is the *tail* and *v* is the *head*.
- In this lecture we will always specify a directed graph explicitly; that is, the term "graph" without the qualifier "directed" shall mean "undirected graph".



Basic Definitions: How to Deal with $V = \emptyset$

- There is no consensus on whether or not to allow V = Ø in the definition of a graph. (Of course, if V = Ø then E = Ø.)
- And, indeed, there are pros and cons of allowing $V = \emptyset$.



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- And, indeed, there are pros and cons of allowing $V = \emptyset$.
- Furthermore, if *V* = Ø is allowed then there is little consensus on how to call such a graph:
 - Common terms are *order-zero graph*, *K*₀, and *null graph*.
 - Some authors also use the term *empty graph* to indicate V = Ø while other authors prefer to reserve this term for a graph with E = Ø but V ≠ Ø.



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Convention

We will always assume that every (directed) graph has at least one node.



No common terminology

The terminology in graph theory lacks a rigorous standardization, both in the German and in the English literature.



No common terminology

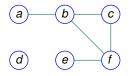
The terminology in graph theory lacks a rigorous standardization, both in the German and in the English literature.

- In several cases the meanings of different terms coincide for simple undirected graphs, which seems to serve as a justification for authors to freely mix and match terms.
- Thus, always make sure to check how some author defines standard terms of graph theory . . .



Undirected Graphs as Directed Graphs

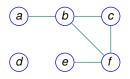
• It is straightforward to represent an undirected graph as a directed graph.

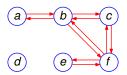




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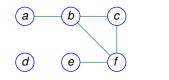


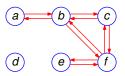




Undirected Graphs as Directed Graphs

- It is straightforward to represent an undirected graph as a directed graph.
- Hence, undirected graphs can be seen as a special case of directed graphs, and most algorithms that work for directed graphs are applicable to undirected graphs, too.







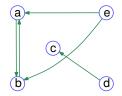
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- E.g., the relation R on the set $\{a, b, c, d, e\}$, with

$$R := \{(a,b), (b,a), (d,c), (e,a), (e,b)\},\$$

corresponds to the following directed graph:

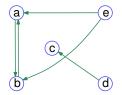




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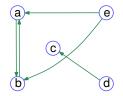
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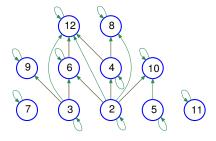
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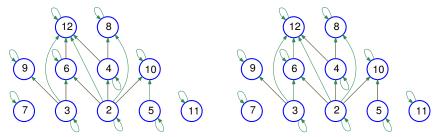
- Hence, statements about relations can be translated to statements about digraphs, and vice versa.
- Note, though, that the digraph corresponding to a relation
 - need not be simple but might contain loops,
 - need not have a finite vertex set.
- Simplified representation of the digraph of an order relation: Hasse diagram





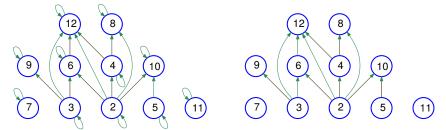


• Consider the poset (S, R), where $S := \{n \in \mathbb{N} : 1 < n \leq 12\}$ and R denotes the partial order of divisibility on S. (That is, for $a, b \in S$, we have a R b iff $a \mid b$.)



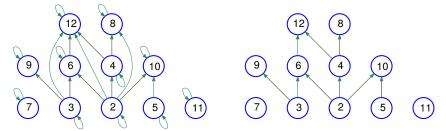
Redraw the digraph such that all oriented (non-loop) edges point upwards.





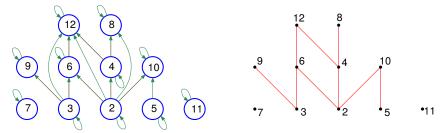
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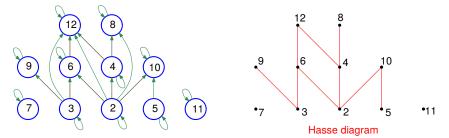




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Definition 220 (Hasse diagram)

The graph obtained after carrying out Steps (1)–(4) is the *Hasse diagram* of the poset.

• Typically, some statements of a computer program could be executed in parallel.

(1) a := 1
(2) b := 2
(3) c := 3
(4) d := a + 2
(5) e := 2a + b
(6) f := d + c
(7) g := c + e
(8) h := d + e + f



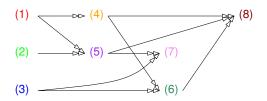
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- A *precedence graph* is a directed graph that models dependences. E.g., the dependence of statements of a computer program on other statements:
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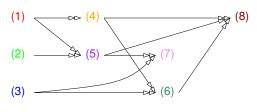






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- Precedence graphs are used in all sorts of scheduling tasks: E.g., job scheduling, concurrency control and instruction scheduling, resolving linker dependencies, data serialization, automated parallelization of sequential code.
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The *degree* (aka *valence*) of a vertex u of a graph $\mathcal{G} := (V, E)$ is the number of edges incident to u. It is denoted by deg(u).



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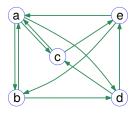
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Definition 223 (Subgraph, Dt.: Teilgraph)

A graph $\mathcal{G}' := (V', E')$ is a *subgraph* of a (directed) graph $\mathcal{G} := (V, E)$ if $V' \subseteq V$ and $E' \subseteq E$ such that all edges of E' are formed by vertices of V'.

The *adjacency matrix* of a (directed) graph $\mathcal{G} := (V, E)$ is an $n \times n$ matrix **M**, where n := |V| and

$$m_{ij} := \left\{ egin{array}{cc} 1 & ext{if } v_i v_j \in E, \ 0 & ext{otherwise.} \end{array}
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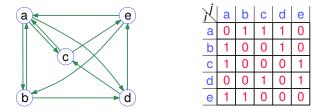


| Ķ | a | b | С | d | e |
|---|---|---|---|---|---|
| а | 0 | 1 | 1 | 1 | 0 |
| b | 1 | 0 | 0 | 1 | 0 |
| С | 1 | 0 | 0 | 0 | 1 |
| d | 0 | 0 | 1 | 0 | 1 |
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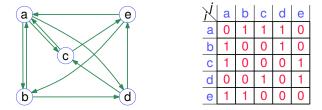


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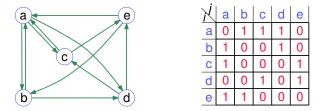


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- Note: Storing **M** (as an $n \times n$ array) requires $\Theta(n^2)$ memory!
- Adjacency lists (and their variants) help to preserve memory if $|E| \ll |V|^2$.

Definition 225 (Regular graph, Dt.: regulärer Graph)

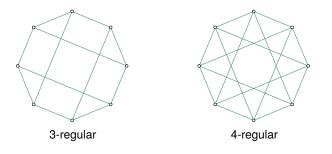
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• A 3-regular graph is known as a cubic graph, and a 4-regular graph is known as a quartic graph.

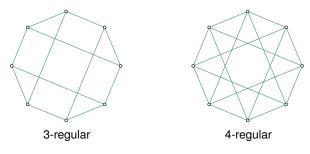




Definition 225 (Regular graph, Dt.: regulärer Graph)

A graph G is *regular* if every vertex of G has the same degree. A regular graph with vertices of degree k is called a k-regular graph or regular graph of degree k.

- A 3-regular graph is known as a cubic graph, and a 4-regular graph is known as a quartic graph.
- For directed regular graphs it is common to demand that the in-degree and the out-degree of each vertex is identical.





The sum over all degrees of vertices of a graph $\mathcal{G} := (V, E)$ equals twice the number of its edges, i.e., $\sum_{\nu \in V} \text{deg}(\nu) = 2|E|$.



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Corollary 227 (Euler's Handshaking Lemma, Dt.: Handschlag-Lemma)

In every graph the number of vertices of odd degree is even.



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- Simple application of Euler's Handshaking Lemma:
 - Suppose that a party is attended by 15 guests. Is it possible that every guest at the party knows all others except for precisely one guest?



The sum over all degrees of vertices of a graph $\mathcal{G} := (V, E)$ equals twice the number of its edges, i.e., $\sum_{\nu \in V} \text{deg}(\nu) = 2|E|$.

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In every graph the number of vertices of odd degree is even.

- Simple application of Euler's Handshaking Lemma:
 - Suppose that a party is attended by 15 guests. Is it possible that every guest at the party knows all others except for precisely one guest?
 - No: Consider a graph with 15 nodes (guests) where two nodes are linked by an edge if the corresponding guests do not know each other. Hence, we would get 15 nodes of degree one, in contradiction to Cor. 227.



Graph Theory

- What is a (Directed) Graph?
- Paths
 - Walks
 - Connectedness
 - Euler Tour and Hamilton Cycle
- Trees
- Special Graphs
- Graph Coloring



Walks

Definition 228 (Walk, Dt.: Wanderung, Kantenfolge)

A *walk* of length *k*, with $k \in \mathbb{N}_0$, on $\mathcal{G} := (V, E)$ is an alternating sequence

 $V_0 e_1 V_1 e_2 V_2 \dots e_k V_k$

of k + 1 vertices $v_0, v_1, \ldots, v_k \in V$ and k edges $e_1, \ldots, e_k \in E$ such that

 $\forall (1 \leq i \leq k) \ e_i = v_{i-1}v_i.$



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Definition 229 (Closed walk, Dt.: geschlossene Wanderung)

A walk is called *closed* if the start vertex and the end vertex are identical. A closed walk of length k is called *trivial* if $k \leq 2$.

A *trail* in a (directed) graph \mathcal{G} is a walk in which all edges are distinct.



A *trail* in a (directed) graph G is a walk in which all edges are distinct.

Definition 231 (Path, Dt.: Pfad)

A path in a (directed) graph G is a walk in which all vertices are distinct.



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Definition 231 (Path, Dt.: Pfad)

A path in a (directed) graph G is a walk in which all vertices are distinct.

Definition 232 (Tour, Dt.: Tour)

A *tour* in a (directed) graph G is a closed trail.



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A *tour* in a (directed) graph G is a closed trail.

Definition 233 (Cycle, Dt.: Zyklus, Kreis)

A cycle in a (directed) graph \mathcal{G} is a non-trivial closed walk in which all but the start and the end vertices are distinct.



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• Note: Distinct vertices implies distinct edges; i.e., every path is a trail and every cycle is a tour.



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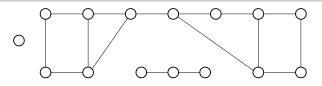
Definition 233 (Cycle, Dt.: Zyklus, Kreis)

A *cycle* in a (directed) graph \mathcal{G} is a non-trivial closed walk in which all but the start and the end vertices are distinct.

- Note: Distinct vertices implies distinct edges; i.e., every path is a trail and every cycle is a tour.
- Note that some authors prefer to use the terms "path", "simple path", "cycle" and "simple cycle" instead of "trail", "path", "tour" and "cycle" ...

Definition 234 (Connected component, Dt.: Zusammenhangskomponente)

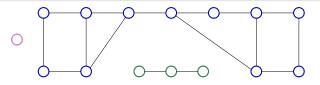
A connected component of a graph $\mathcal{G} := (V, E)$ is a maximal subgraph $\mathcal{G}' := (V', E')$ of \mathcal{G} such that for every unordered pair $\{u, v\}$, with $u, v \in V'$ and $u \neq v$, there exists a path between u and v within \mathcal{G}' .





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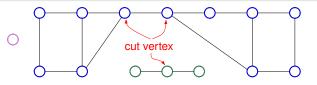
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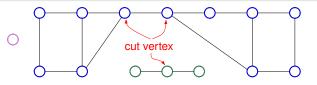
Definition 235 (Cut vertex, Dt.: Artikulationspunkt, Schnittknoten)

A *cut vertex* of a graph $\mathcal{G} := (V, E)$ is a vertex $v \in V$ such that the removal of v and of all edges incident to v would increase the number of connected components.



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Definition 236 (Connected, Dt.: zusammenhängend)

A graph is *connected* if it contains only one connected component.



Definition 237 (Weakly connected, Dt.: schwach zusammenhängend)

A directed graph is *weakly connected* if replacing all its directed edges by undirected edges results in a connected (undirected) graph.



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Definition 238 (Strong component, Dt.: starke Zusammenhangskomponente)

A strong component (aka strongly connected component) of a directed graph $\mathcal{G} := (V, E)$ is a maximal subgraph $\mathcal{G}' = (V', E')$ of \mathcal{G} such that for every ordered pair (u, v), with $u, v \in V'$ and $u \neq v$, there exists a path from u to v within \mathcal{G}' .



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A directed graph is *weakly connected* if replacing all its directed edges by undirected edges results in a connected (undirected) graph.

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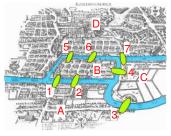
Definition 239 (Strongly connected, Dt.: stark zusammenhängend)

A directed graph $\mathcal{G} := (V, E)$ is *strongly connected* if it consists of only one strong component, i.e., if for every ordered pair (u, v), with $u, v \in V$ and $u \neq v$, there exists a path from u to v.



Seven Bridges of Königsberg

• Early 18th century: Does there exist a trail (or even a tour) through the city of Königsberg that crosses every of its seven bridges exactly once? (Of course, every bridge had to be crossed fully, and no other means to get across the river Pregel were allowed.)

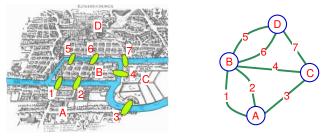


[Image credit for background image: Wikipedia.]



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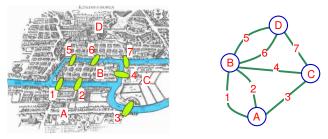


[Image credit for background image: Wikipedia.]

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[Image credit for background image: Wikipedia.]

- In 1736, Leonhard Euler (1707–1783) treated this problem as a graph problem and proved, using a parity argument, that such a trail or tour does not exist.
- His solution is generally regarded as the first theorem of graph theory.



An Euler trail is a trail that contains all edges of a graph exactly once.



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Definition 241 (Euler tour, Dt.: Eulersche Tour)

An *Euler tour* is a tour that contains all edges of a graph exactly once. A graph is an *Eulerian graph* if it has an Euler tour.



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Definition 241 (Euler tour, Dt.: Eulersche Tour)

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A Hamilton path is a path that passes through all vertices of a graph exactly once.

Definition 243 (Hamilton cycle, Dt.: Hamiltonscher Kreis)

A Hamilton cycle is a cycle that passes through all vertices of a graph exactly once.

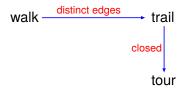


walk

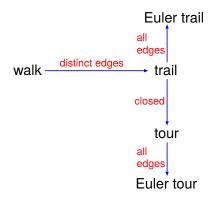


walk <u>distinct edges</u> trail

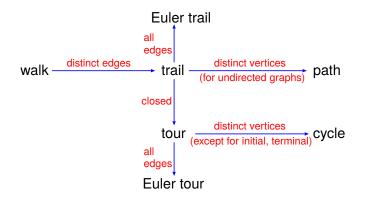




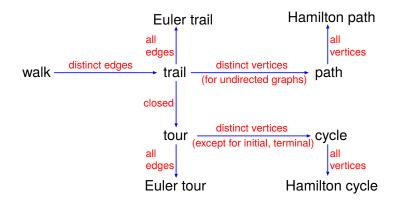














Euler Tour

Theorem 244

Suppose that every node of a graph \mathcal{G} has degree at least one. Then \mathcal{G} has an Euler tour if and only if \mathcal{G} is connected and every vertex of \mathcal{G} has even degree.



Euler Tour

Theorem 244

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Theorem 245

Suppose that every node of a graph \mathcal{G} has degree at least one. Then \mathcal{G} has an Euler trail (but no Euler tour) if and only if \mathcal{G} is connected and exactly two vertices of \mathcal{G} have odd degrees.



Euler Tour

Theorem 244

Suppose that every node of a graph \mathcal{G} has degree at least one. Then \mathcal{G} has an Euler tour if and only if \mathcal{G} is connected and every vertex of \mathcal{G} has even degree.

Theorem 245

Suppose that every node of a graph \mathcal{G} has degree at least one. Then \mathcal{G} has an Euler trail (but no Euler tour) if and only if \mathcal{G} is connected and exactly two vertices of \mathcal{G} have odd degrees.

Corollary 246

An Euler tour or trail in a graph $\mathcal{G} := (V, E)$ can be determined in O(|E|) time, if it exists. Otherwise, again in O(|E|) time, we can determine that neither an Euler tour nor an Euler trail exists in \mathcal{G} .



Constructive Proof of Theorem 244

Sketch of proof of Theorem 244: Let $\mathcal{G} := (V, E)$ be a graph such that every node of a graph \mathcal{G} has degree at least one.



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Suppose that \mathcal{G} has an Euler tour \mathcal{T} . It is obvious that \mathcal{G} is connected.



Constructive Proof of Theorem 244

Sketch of proof of Theorem 244: Let $\mathcal{G} := (V, E)$ be a graph such that every node of a graph \mathcal{G} has degree at least one.

Suppose that \mathcal{G} has an Euler tour T. It is obvious that \mathcal{G} is connected. Every occurrence of a vertex $v \in V$ in T is preceded and followed by an edge. Thus, each time T passes through v, two of the edges incident to v are consumed. Since T does neither start nor end in v, it is necessary that deg(v) is even.



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Now suppose that every vertex of \mathcal{G} has even degree, and, of course, that \mathcal{G} is connected. We give a constructive proof that \mathcal{G} admits an Euler tour.



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We realize that, eventually, T will get us back to v. (We cannot be stuck in some other vertex w since w has even degree.)



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Now suppose that every vertex of \mathcal{G} has even degree, and, of course, that \mathcal{G} is connected. We give a constructive proof that \mathcal{G} admits an Euler tour. Pick any vertex v to start with and trace out a trail T. Every edge that is being traversed is marked. As above, we observe that passing through a vertex that is neither the start nor the end vertex of T consumes two edges.

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This process continues until no unmarked edges remain. At the end the trails are spliced together appropriately.

Hamilton Cycle

Theorem 247

It is $\mathcal{NP}\text{-}complete$ to determine whether a Hamilton cycle or Hamilton path exists in a general graph.

- Informally, Theorem 247 says that no (deterministic sequential) algorithm is known which determines the existence of a Hamilton cycle or path in an *n*-vertex graph in a time that is a polynomial function of *n*.
- Even worse, an efficient (polynomial-time) algorithm will never be found unless $\mathcal{P} = \mathcal{NP}$ holds, which seems rather unlikely.



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Theorem 248 (Dirac, 1952)

If the degree of every vertex of an *n*-vertex graph \mathcal{G} , with $n \ge 3$, is at least $\lceil \frac{n}{2} \rceil$ then \mathcal{G} has a Hamilton cycle.



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Theorem 249 (Ore, 1960)

If the sum of the degrees of every pair of non-adjacent vertices of an *n*-vertex graph \mathcal{G} , with $n \ge 3$, is at least *n* then \mathcal{G} has a Hamilton cycle.

Graph Theory

- What is a (Directed) Graph?
- Paths

Trees

- Basic Definitions
- Elementary Properties
- Binary Trees
- Balance and Height
- Spanning Trees
- Recursion Trees
- Special Graphs
- Graph Coloring



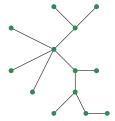
A graph is called acyclic if it contains no cycles.



A graph is called *acyclic* if it contains no cycles.

Definition 251 (Tree, Dt.: Baum)

A tree is an undirected graph that is acyclic and connected.



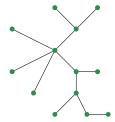


A graph is called *acyclic* if it contains no cycles.

Definition 251 (Tree, Dt.: Baum)

A tree is an undirected graph that is acyclic and connected.

• For trees most authors prefer to speak about nodes rather than vertices.



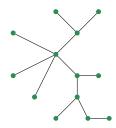


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Definition 251 (Tree, Dt.: Baum)

A tree is an undirected graph that is acyclic and connected.

- For trees most authors prefer to speak about nodes rather than vertices.
- Unless explicitly stated otherwise, we will only deal with trees that have at least one node. (Some authors call a tree with V = E = Ø a null tree.)

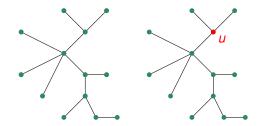




Definition 252 (Rooted tree, Dt.: Baum mit Wurzel, Wurzelbaum)

A rooted tree is a directed graph with a node u such that

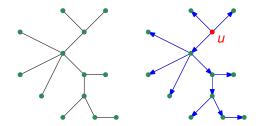
• the graph contains *u* as node ("*root*"),





Definition 252 (Rooted tree, Dt.: Baum mit Wurzel, Wurzelbaum)

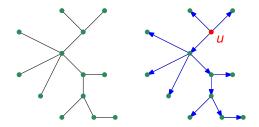
- the graph contains u as node ("root"),
- 2 paths from u to all other nodes of the graph exist,





Definition 252 (Rooted tree, Dt.: Baum mit Wurzel, Wurzelbaum)

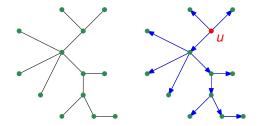
- the graph contains u as node ("root"),
- 2 paths from u to all other nodes of the graph exist,
- the in-degree of u is zero,





Definition 252 (Rooted tree, Dt.: Baum mit Wurzel, Wurzelbaum)

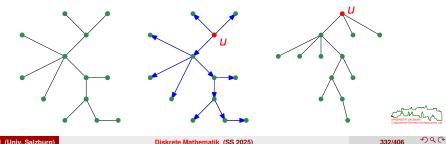
- the graph contains u as node ("root"),
- 2 paths from u to all other nodes of the graph exist,
- the in-degree of u is zero,
- the in-degree of every other node of the graph is one.





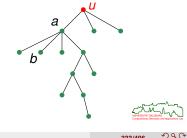
Definition 252 (Rooted tree, Dt.: Baum mit Wurzel, Wurzelbaum)

- the graph contains u as node ("root").
- 2 paths from u to all other nodes of the graph exist.
- the in-degree of u is zero,
- the in-degree of every other node of the graph is one.
- It is common practice to draw rooted trees from the root downwards such that the (downwards) orientations of the edges are implied by the positions of the nodes.



Definition 253 (Child and parent, Dt.: Kind und Eltern)

For a rooted tree $\mathcal{T} := (V, E)$ and nodes $a, b \in V$, the node b is a *child* of the node a, and a is the parent of b, if the edge ab belongs to E. Siblings are nodes which share the same parent.

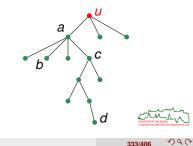


Definition 253 (Child and parent, Dt.: Kind und Eltern)

For a rooted tree $\mathcal{T} := (V, E)$ and nodes $a, b \in V$, the node *b* is a *child* of the node *a*, and *a* is the *parent* of *b*, if the edge *ab* belongs to *E*. *Siblings* are nodes which share the same parent.

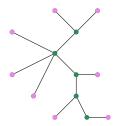
Definition 254 (Descendant and ancestor, Dt.: Nachfahre und Vorfahre)

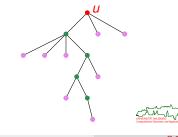
In a rooted tree $\mathcal{T} := (V, E)$, with $c, d \in V$, a node d is a *descendant* of a node c, and c is an *ancestor* of d, if $c \neq d$ and if the path from the root to d contains c.



Definition 255 (Leaf, Dt.: Blatt)

A *leaf* of a rooted tree is a node without children. For a tree (that is not rooted) a leaf is a node with degree 1. All non-leaf nodes of a (rooted) tree are called *inner nodes*.

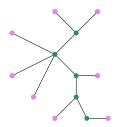


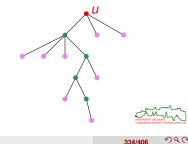


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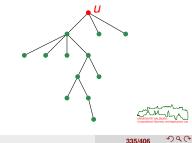
• Of course, the root of a rooted tree \mathcal{T} may also be the (only) leaf of \mathcal{T} .





Definition 256 (Subtree, Dt.: Teilbaum)

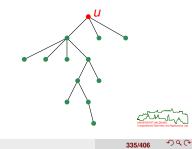
A tree $\mathcal{T}' := (V', E')$ is a *subtree* of a tree $\mathcal{T} := (V, E)$ rooted at the node *u* if \mathcal{T}' is a subgraph of \mathcal{T} ,



Definition 256 (Subtree, Dt.: Teilbaum)

A tree $\mathcal{T}' := (V', E')$ is a *subtree* of a tree $\mathcal{T} := (V, E)$ rooted at the node *u* if

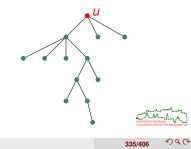
- \mathcal{T}' is a subgraph of \mathcal{T} ,
- 2 \mathcal{T}' is rooted at a node v that is a descendant of u,



Definition 256 (Subtree, Dt.: Teilbaum)

A tree $\mathcal{T}' := (V', E')$ is a *subtree* of a tree $\mathcal{T} := (V, E)$ rooted at the node *u* if

- $\ \, \bullet \ \, {\mathcal T}' \text{ is a subgraph of } {\mathcal T},$
- 2 \mathcal{T}' is rooted at a node v that is a descendant of u, and
- **9** \mathcal{T}' contains all descendants of v in \mathcal{T} , together with the appropriate edges of E.

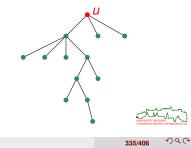


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A subtree rooted at v is called a proper subtree if v is a child of u.



Definition 256 (Subtree, Dt.: Teilbaum)

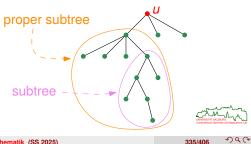
A tree $\mathcal{T}' := (V', E')$ is a *subtree* of a tree $\mathcal{T} := (V, E)$ rooted at the node *u* if

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A subtree rooted at v is called a *proper subtree* if v is a child of u.

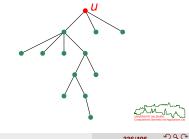
Warning

Some authors do not make the distinction between the node v being a child of u or some arbitrary descendant of u.



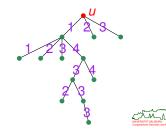
Definition 257 (Ordered tree, Dt.: geordneter Baum)

An *ordered tree* is a rooted tree T such that the children of every node of T are arranged in some specific order, e.g., by means of a numbering scheme.



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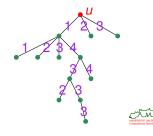


Definition 257 (Ordered tree, Dt.: geordneter Baum)

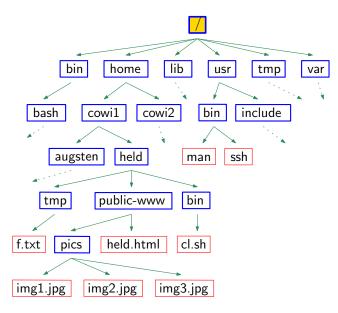
An *ordered tree* is a rooted tree \mathcal{T} such that the children of every node of \mathcal{T} are arranged in some specific order, e.g., by means of a numbering scheme.

Definition 258 (Forest, Dt.: Wald)

A forest is a graph such that all its connected components are trees.

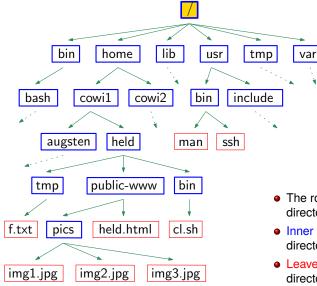


Real-World Application: File System as a Rooted Tree





Real-World Application: File System as a Rooted Tree



- The root of the tree is the root directory /.
- Inner nodes are (non-empty) directories.
- Leaves are files (or empty directories).



Theorem 259

Every pair of nodes in a tree is connected by exactly one path.



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Theorem 260

In a rooted tree there exists exactly one path from the root to any node.



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In a rooted tree there exists exactly one path from the root to any node.

Lemma 261

Removing an edge from a (rooted) tree results in a graph with two connected components, each of which is a (rooted) tree.



Theorem 262

For every (rooted) tree $\mathcal{T} := (V, E)$ we get |E| = |V| - 1.



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Proof of Theorem 262 for rooted trees : We use structural induction relative to proper subtrees.



Theorem 262

For every (rooted) tree $\mathcal{T} := (V, E)$ we get |E| = |V| - 1.

Proof of Theorem 262 for rooted trees : We use structural induction relative to proper subtrees. Obviously, the claim holds for the minimal elements, i.e., for trees that contain no proper subtrees and, thus, have only a root and no edges.



For every (rooted) tree
$$\mathcal{T} := (V, E)$$
 we get $|E| = |V| - 1$.

Proof of Theorem 262 for rooted trees : We use structural induction relative to proper subtrees. Obviously, the claim holds for the minimal elements, i.e., for trees that contain no proper subtrees and, thus, have only a root and no edges. Now consider an arbitrary but fixed rooted tree $\mathcal{T} := (V, E)$ and suppose that the equality claimed holds for all its k > 0 proper subtrees $(V_1, E_1), \ldots, (V_k, E_k)$. (We do not need to assume explicitly that it holds for all subtrees of \mathcal{T} .)



For every (rooted) tree
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$$|E| = k + \sum_{i=1}^{k} |E_i|$$



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$$|E| = k + \sum_{i=1}^{k} |E_i| = k + \sum_{i=1}^{k} (|V_i| - 1) = k + (-k) + \sum_{i=1}^{k} |V_i| = \sum_{i=1}^{k} |V_i|$$



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= |V| - 1,

thus establishing the claim also for $\mathcal{T} = (V, E)$.



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Proof of Theorem 262 for rooted trees : We use structural induction relative to proper subtrees. Obviously, the claim holds for the minimal elements, i.e., for trees that contain no proper subtrees and, thus, have only a root and no edges. Now consider an arbitrary but fixed rooted tree $\mathcal{T} := (V, E)$ and suppose that the equality claimed holds for all its k > 0 proper subtrees $(V_1, E_1), \ldots, (V_k, E_k)$. (We do not need to assume explicitly that it holds for all subtrees of \mathcal{T} .) We get

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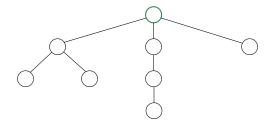
= |V| - 1,

thus establishing the claim also for T = (V, E).

Corollary 263

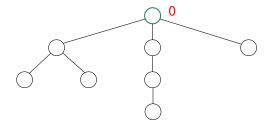
If |V| > 1 holds for a (rooted) tree $\mathcal{T} := (V, E)$, then \mathcal{T} has at least one leaf.

Definition 264 (Depth, Dt.: Tiefe)



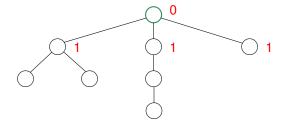


Definition 264 (Depth, Dt.: Tiefe)



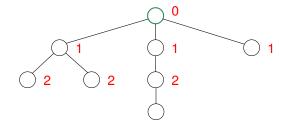


Definition 264 (Depth, Dt.: Tiefe)



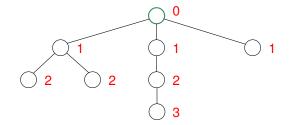


Definition 264 (Depth, Dt.: Tiefe)





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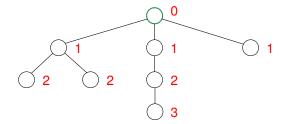


Definition 264 (Depth, Dt.: Tiefe)

The *depth* of the root *u* of a rooted tree $\mathcal{T} := (V, E)$ is 0, and the depth of a node $v \neq u$ of \mathcal{T} is *k* if the depth of the parent of *v* is k - 1, for all $v \in V$.

Warning

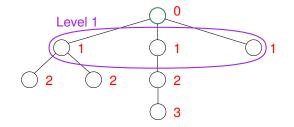
Some authors prefer to regard the root as a node at depth 1. Hence, make sure to check how depth is defined in a textbook prior to using the results stated!





Definition 265 (Level, Dt.: Niveau)

A *level* of a rooted tree \mathcal{T} comprises all nodes of \mathcal{T} which have the same depth.



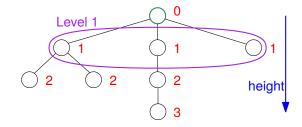


Definition 265 (Level, Dt.: Niveau)

A *level* of a rooted tree \mathcal{T} comprises all nodes of \mathcal{T} which have the same depth.

Definition 266 (Height, Dt.: Höhe)

The *height* of a rooted tree T is the maximum depth of nodes of T.





Definition 267 (Binary tree, Dt.: Binärbaum)

A *binary tree* is an ordered tree T with a root node u and at most two proper subtrees that are called *left subtree*, L, and *right subtree*, R. If T has a left (right, resp.) subtree then L(R, resp.) is in turn a binary tree rooted in the left (right, resp.) child of u.



Binary Tree

Definition 267 (Binary tree, Dt.: Binärbaum)

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Definition 268 (Complete binary tree, Dt.: vollständiger Binärbaum)

A *complete binary tree* is a binary tree in which every level, except possibly the last level, is completely filled, and the last level is filled from left to right.



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• E.g., a (binary) heap is a complete binary tree.



Binary Tree

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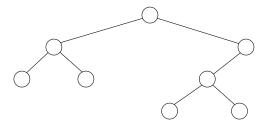
• E.g., a (binary) heap is a complete binary tree.

Definition 269 (Perfect binary tree, Dt.: perfekter Binärbaum)

A *perfect binary tree* is a binary tree that has the maximum number of nodes (relative to its height).



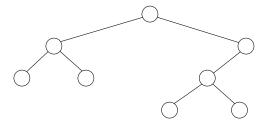
Definition 270 (Binary search tree, Dt.: binärer Suchbaum)





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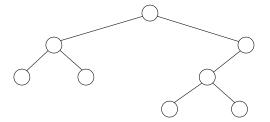
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 - all values of nodes in L are less than the root value,





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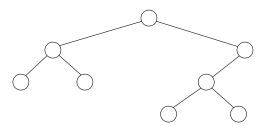
- if it has a proper left subtree L then
 - all values of nodes in L are less than the root value,
 - 2 L is a binary search tree itself,





Definition 270 (Binary search tree, Dt.: binärer Suchbaum)

- if it has a proper left subtree L then
 - **1** all values of nodes in *L* are less than the root value,
 - 2 L is a binary search tree itself,
- if it has a proper right subtree R then
 - 3 all values of nodes in *R* are greater than the root value,
 - Is a binary search tree itself.

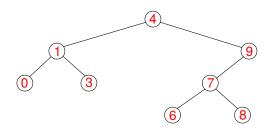




Definition 270 (Binary search tree, Dt.: binärer Suchbaum)

A *binary search tree* is a binary tree T which has distinct values associated with its nodes such that (relative to some total order)

- if it has a proper left subtree L then
 - all values of nodes in L are less than the root value,
 - 2 L is a binary search tree itself,
- if it has a proper right subtree R then
 - 3 all values of nodes in *R* are greater than the root value,
 - Is a binary search tree itself.





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Definition 271 (k-balanced tree, Dt.: k-balanzierter Baum)

A binary tree is *height-balanced* with balance factor k if it either has no proper subtrees



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- it has two proper subtrees and the heights of both subtrees differ by not more than *k*, or if
- 2 it has one proper subtree of height at most k 1,

and if

- 3 all proper subtrees are height-balanced with balance factor *k*.
 - E.g., for *k* := 1: AVL tree.
 - Trees with balance factor 1 are simply called balanced or self-balancing.



Definition 272 (Perfectly balanced binary tree, Dt.: perfekt balanz. Binärbaum)

A binary tree T is *perfectly balanced* if all inner nodes of T, except possibly on the second-last level, have exactly two children.



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A binary tree T is *perfectly balanced* if all inner nodes of T, except possibly on the second-last level, have exactly two children.

• E.g., a (binary) heap is a perfectly balanced binary tree.

Lemma 273

A complete binary tree is perfectly balanced.



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• E.g., a (binary) heap is a perfectly balanced binary tree.

Lemma 273

A complete binary tree is perfectly balanced.

Lemma 274

A perfectly balanced binary tree has leaves only at its two bottom-most levels.



Height-Related Properties of Binary Trees

Lemma 275

For $i \in \mathbb{N}_0$, level *i* of a binary tree contains at most 2^i nodes.



For $i \in \mathbb{N}_0$, level *i* of a binary tree contains at most 2^i nodes.

Sketch of proof by induction: The claim holds for i := 0. If we have at most 2^k nodes on level *k* then we have at most $2 \cdot 2^k = 2^{k+1}$ nodes on level k + 1.



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Sketch of proof by induction : The claim holds for i := 0. If we have at most 2^k nodes on level k then we have at most $2 \cdot 2^k = 2^{k+1}$ nodes on level k + 1.

Lemma 276

Let *h* be the height and *n* be the number of nodes of a binary tree. Then $h \ge \lfloor \log(n+1) \rfloor - 1$, i.e., $h \in \Omega(\log n)$.



For $i \in \mathbb{N}_0$, level *i* of a binary tree contains at most 2^i nodes.

Sketch of proof by induction : The claim holds for i := 0. If we have at most 2^k nodes on level k then we have at most $2 \cdot 2^k = 2^{k+1}$ nodes on level k + 1.

Lemma 276

Let *h* be the height and *n* be the number of nodes of a binary tree. Then $h \ge \lceil \log(n+1) \rceil - 1$, i.e., $h \in \Omega(\log n)$.

Proof: Lemma 275 implies that a binary tree with height *h* contains at most

$$\sum_{i=0}^{h} 2^{i} = 2^{h+1} - 1$$

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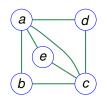
Theorem 277

If \mathcal{T} is a balanced binary tree with *n* nodes and height *h* then $h \in \Theta(\log n)$.

Definition 278 (Spanning tree, Dt.: spannender Baum)

A spanning tree of a connected graph ${\mathcal G}$ is a subgraph of ${\mathcal G}$ that

- is a tree,
- 2 includes all vertices of \mathcal{G} .

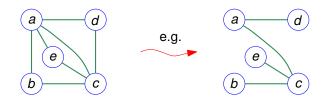




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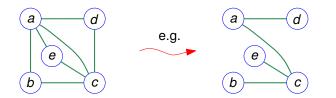
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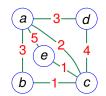
Every connected graph \mathcal{G} contains a spanning tree.





Definition 280 (Weighted graph, Dt.: gewichteter Graph)

An (*edge-*)weighted graph is a graph in which every edge is assigned a (non-negative) real number, the so-called weight or cost.



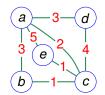


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A minimum spanning tree (MST) of a weighted connected graph \mathcal{G} is a spanning tree \mathcal{T} of \mathcal{G} such that the sum of the weights of the edges of \mathcal{T} is minimum over all spanning trees of \mathcal{G} .





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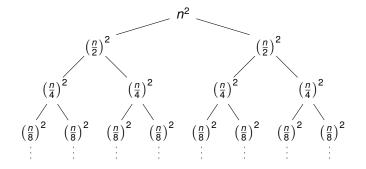
• A *recursion tree* visualizes the recursive calls of a function and the work done at each call, as given by a recurrence relation.



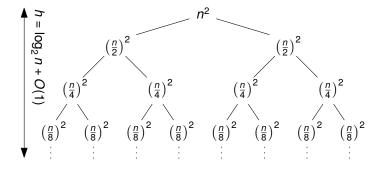
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- E.g., consider $T(n) = 2T(n/2) + n^2$.



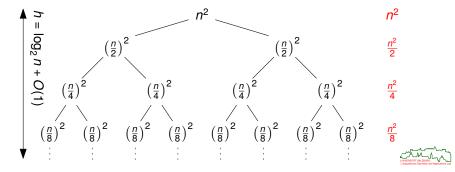
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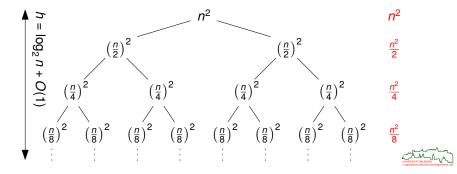
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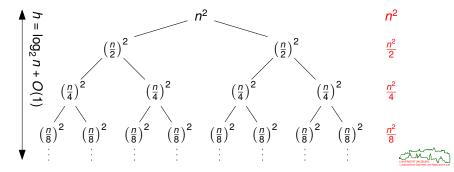
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- Summing across every level gives the total work done per level.



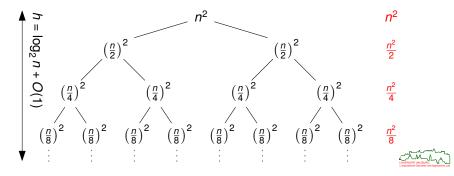
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- Master Theorem 215: We have a = b = k = 2 and, thus, $a < b^k$.



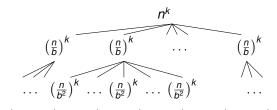
 Note that in this case the height of the tree does not really matter: The amount of work done at every level decreases so quickly that the total work is only a constant factor more than the work done at the root.



• For the recurrence relation $T(n) = a \cdot T\left(\frac{n}{b}\right) + n^k$ we get an *a*-ary recursion tree:

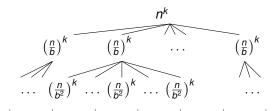


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 - The problem size at level *i* is n/b^i .



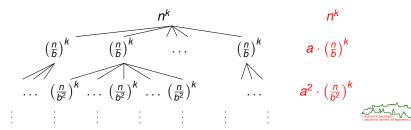


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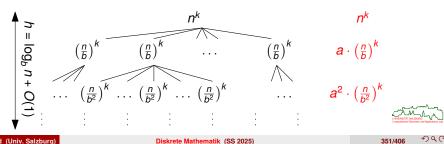




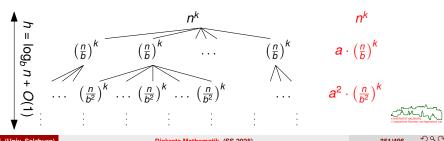
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 - The tree has $\log_b n + O(1)$ levels, i.e., a height of $O(\log n)$.

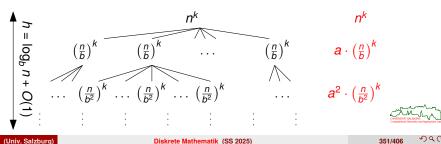


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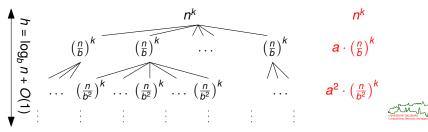
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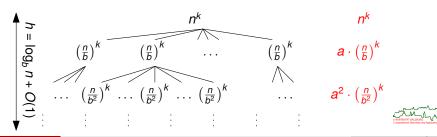
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 - Total work:

$$T(n) = \sum_{0 \leq i < \log_b n} a^i \cdot \left(\frac{n}{b^i}\right)^k + O(n^{\log_b a}) = \sum_{0 \leq i < \log_b n} n^k \cdot \left(\frac{a}{b^k}\right)^i + O(n^{\log_b a}).$$



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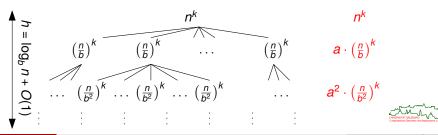
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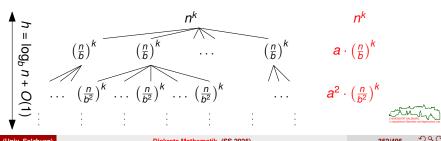
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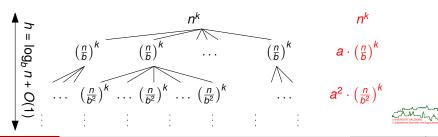
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 Hence, the same order of work is done on every level, and since the tree has $O(\log n)$ levels, we get $T \in \Theta(n^{\log_b a} \log n)$; recall Thm. 215.



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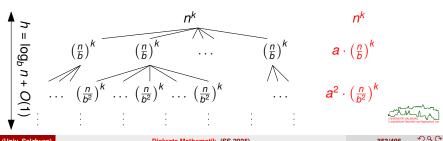
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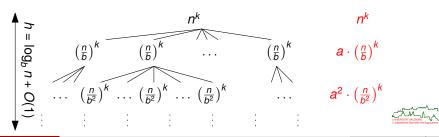
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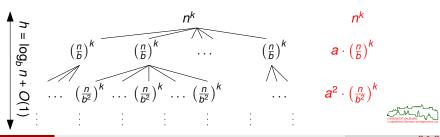
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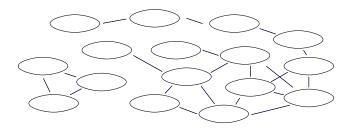
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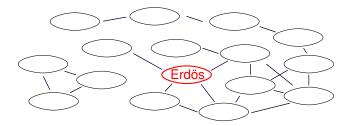


• A *collaboration graph* for a set of *n* scientists is a graph with *n* vertices such that two vertices are connected by an edge if the corresponding scientists have at least one joint publication.



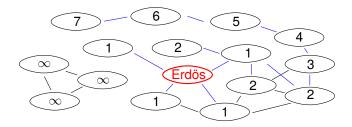


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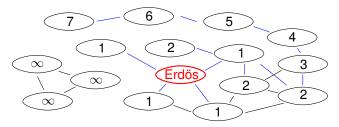


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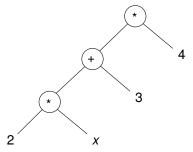
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- One's Erdös number can be obtained by computing minimum-weight paths on a collaboration graph.





Real-World Application: Algebraic Expression Trees

- An algebraic expression tree is a rooted tree that corresponds to an expression.
- E.g., an in-order traversal of the tree

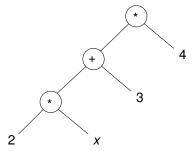


produces the standard (infix) expression $(2x + 3) \cdot 4$.



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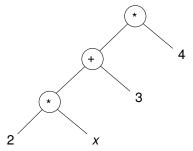
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A post-order traversal yields the postfix expression 2 x · 3 + 4 · , while a pre-order traversal yields the prefix expression · (+(·(2 x) 3) 4).



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- The analysis of expression trees is a central task for the simplification and parallel evaluation of an expression.



7 G

Graph Theory

- What is a (Directed) Graph?
- Paths
- Trees
- Special Graphs
 - Complete and Bipartite Graphs
 - Hypercube
 - Isomorphic Graphs
 - Planar Graphs
- Graph Coloring



Definition 282 (Complete graph, Dt.: vollständiger Graph)

For $n \in \mathbb{N}$, the *complete graph* on *n* vertices, commonly denoted by K_n , is an undirected graph with *n* vertices in which every pair of vertices is adjacent.



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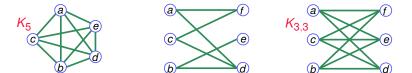
Definition 283 (Bipartite graph, Dt.: bipartiter Graph)

An undirected graph $\mathcal{G} := (V, E)$ is a *bipartite graph* if V can be partitioned into two (non-empty) subsets V_1, V_2 such that $E \subseteq \{\{v_1, v_2\} : v_1 \in V_1, v_2 \in V_2\}$.



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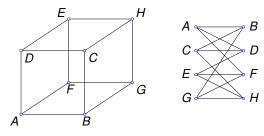
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Definition 284 (Complete bipartite graph, Dt.: vollständig-bipartiter Graph)

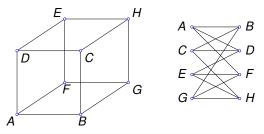
An undirected graph $\mathcal{G} := (V, E)$ is a *complete bipartite graph* if *V* can be partitioned into two (non-empty) subsets V_1 , V_2 such that $E = \{ \{v_1, v_2\} : v_1 \in V_1, v_2 \in V_2 \}$. If $n := |V_1|$ and $m := |V_2|$ then \mathcal{G} is denoted by $K_{n,m}$.

• The edges and corners of a cube can be interpreted as a bipartite graph.

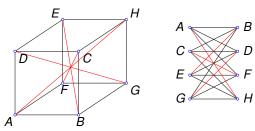




• The edges and corners of a cube can be interpreted as a bipartite graph.



• If we add all diagonals that cross the cube then we get $K_{4,4}$.





Lemma 285

Let $\mathcal{G} := (V, E)$ be a bipartite graph and let V_1, V_2 be the partition of V according to Def. 283. Then

$$\sum_{v_1 \in V_1} \deg(v_1) = \sum_{v_2 \in V_2} \deg(v_2) = |\boldsymbol{E}|.$$



Lemma 285

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$$\sum_{v_1 \in V_1} \deg(v_1) = \sum_{v_2 \in V_2} \deg(v_2) = |E|.$$

Proof:

• As each edge has exactly one vertex from V₁, we can write

$$\sum_{v_1\in V_1} \deg(v_1) = |E|.$$

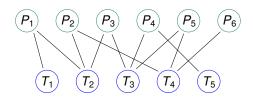
Similarly,

$$\sum_{v_2 \in V_2} \deg(v_2) = |E|.$$

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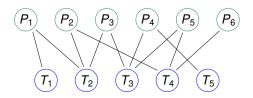


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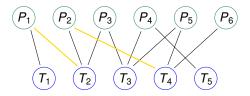




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Definition 286 (Matching, Dt.: Paarung)

• A *matching* in a simple graph $\mathcal{G} := (V, E)$ is a subset E' of E such that no two edges of E' are incident upon the same vertex of V.

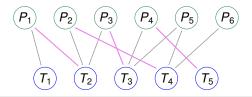




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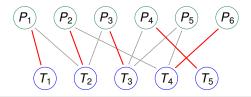




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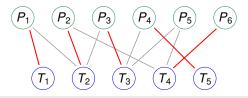




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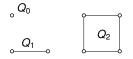




Definition 287 (Hypercube)

For $n \in \mathbb{N}_0$, the hypercube Q_n is defined recursively as follows:

- Q_0 is a single vertex;
- Q_{n+1} is obtained by taking two disjoint copies of Q_n and linking each vertex in one copy of Q_n to the corresponding vertex in the other copy of Q_n.

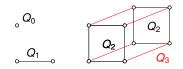




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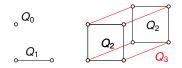


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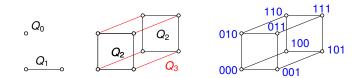
For $n \in \mathbb{N}_0$, the hypercube Q_n is a regular graph of degree n with 2^n vertices and $n \cdot 2^{n-1}$ edges; it is bipartite for $n \ge 1$.



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 We could also obtain Q_n by labeling 2ⁿ vertices with distinct n-bit binary strings, and by connecting those vertices by edges whose strings differ in exactly one bit.

Definition 289 (Gray code)

A (cyclic) Gray code of an ordered sequence of 2^n entities, for $n \in \mathbb{N}$, is a sequence of *n*-bit binary strings such that the encodings of two neighboring entities have Hamming distance one, i.e., differ by exactly one bit.

• Gray codes are widely used in position encoders and for error detection and correction in digital communication.



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For $n \in \mathbb{N}$ with $n \ge 2$, the number of different *n*-bit cyclic Gray codes equals the number of different Hamilton cycles in Q_n .



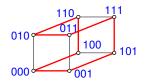
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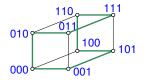
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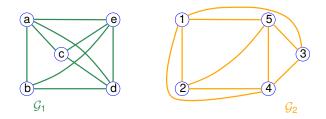




Definition 291 (Isomorphic, Dt.: isomorph)

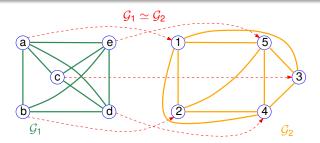


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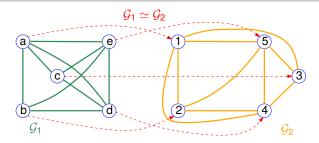


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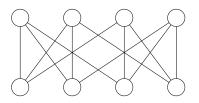


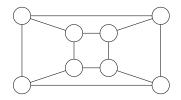
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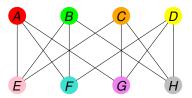
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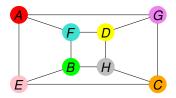






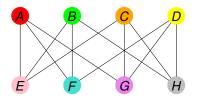
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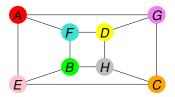






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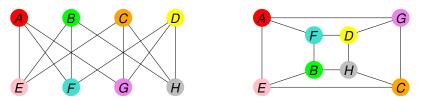




 Necessary (but not sufficient) conditions for two graphs to be isomorphic: same numbers of vertices and edges, same degrees.



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- Necessary (but not sufficient) conditions for two graphs to be isomorphic: same numbers of vertices and edges, same degrees.
- The complexity of the graph isomorphism problem for general *n*-vertex graphs is unknown. No polynomial-time algorithm is known, but the problem is also not known to be \mathcal{NP} -complete. In December 2015, Babai announced a deterministic algorithm that runs in time $2^{O(\log^c n)}$ time for some positive constant *c*, i.e., in quasi-polynomial time. In 2017, Helfgott claimed that one can take *c* := 3.
- Practically efficient algorithms for graph canonical labeling are known, though.

CTM MA

Real-World Application: Non-Isomorphic Trees Represent Molecules

- [Cayley 1857]: Molecules can be represented as graphs, where atoms are represented by vertices and bonds are represented by edges.
- Saturated hydrocarbons, C_nH_{2n+2}, are given by trees where each carbon atom is represented by a degree-four vertex and each hydrogen atom is a leaf.



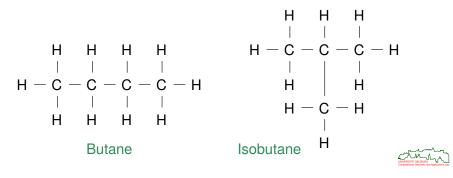
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- How many different isomers can exist for *n* := 4?
- We have exactly two non-isomorphic trees of this type and, thus, two different isomers of C₄H₁₀, namely butane and isobutane.



Planar Graphs

Definition 293 (Planar graph, Dt.: planarer oder plättbarer Graph)

A *planar graph* is a graph which can be drawn in the plane without edge crossings. A suitable drawing is called a *(planar) embedding* (Dt.: planare Einbettung).

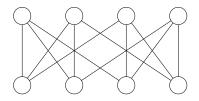


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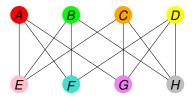


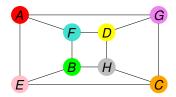
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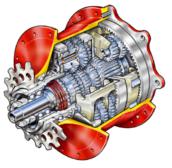
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- VLSI circuits are easier/cheaper to manufacture if their connections live in fewer layers.
- A scheme for a planetary gearset is compatible if and only if a suitably designed graph is planar.

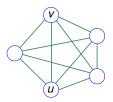


[Image credit: Rohloff AG, http://www.rohloff.de/]



Definition 294 (Subdivision, Dt.: Unterteilung)

An *edge subdivision* of the edge $uv \in E$ by means of the vertex $w \notin V$ transforms the graph $\mathcal{G} := (V, E)$ into the graph $\mathcal{G}' = (V', E')$, where $V' = V \cup \{w\}$ and $E' = (E \setminus \{uv\}) \cup \{uw, wv\}$.





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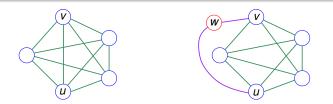


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A graph \mathcal{G}' is a subdivision graph of \mathcal{G} if \mathcal{G}' is obtained from \mathcal{G} via a finite sequence of edge subdivisions.





A graph is planar if and only if it does not contain a subgraph that is isomorphic to a subdivision graph of K_5 or $K_{3,3}$.



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If a graph contains K_5 or $K_{3,3}$ as a subgraph then it is not planar.

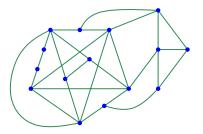


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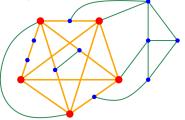


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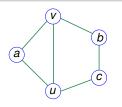
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• Is the following graph planar? No: It contains a subdivision graph of K_5 as a subgraph. Hence, it is not planar.



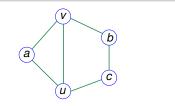


In a graph $\mathcal{G} := (V, E)$, the *contraction* of an edge $e \in E$, with e = uv for some $u, v \in V$, replaces u and v by a new vertex $w \notin V$ such that edges incident to w are all edges other than e that were incident with u or v. All other nodes and edges are preserved.





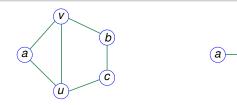
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Theorem 299 (Wagner (1937))

A graph is planar if and only if it does not contain a subgraph that can be contracted to K_5 or $K_{3,3}$ via a finite sequence of edge contractions.



Testing whether a given graph with *n* vertices is planar can be done in O(n) time.



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Theorem 301 (Wagner (1936), Fáry (1948), Stein (1951))

Any planar graph can be embedded into the plane without edge crossings such that all its edges are straight-line segments:





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A *face* of a PSLG embedding of a planar graph is a maximal connected region of the plane that is disjoint from all edges. The embedding of the graph together with the collection of faces induced is called *planar subdivision*.



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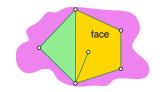
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• Note that one of the faces of a planar subdivision is unbounded: *outer face*.



Consider a planar subdivision induced by a connected planar graph G. We denote

- the number of its vertices by v,
- the number of its edges by e,
- the number of its faces by f.

Then

v - e + f = 2.



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Proof: Suppose that \mathcal{G} is connected but no tree. Therefore \mathcal{G} contains a cycle, and we may remove an edge from \mathcal{G} without destroying its connectivity. The removal of one edge of a cycle decreases both e and f by one, implying that the value of v - e + f does not change.



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• Euler's Formula generalizes to v - e + f = 1 + c for a planar graph with c connected components.

Let v, e, f for a connected planar graph G as defined in Theorem 303. If $v \ge 3$ then

$$e \leq 3\nu - 6$$
 and $f \leq 2\nu - 4$ and $f \leq \frac{2}{3}e$.



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If every vertex of ${\mathcal G}$ has a degree of three or greater then we get

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Furthermore, every planar graph contains one node with degree at most five.



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 and $e \leq 3f - 6$ and $v \leq 2f - 4$.

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Euler's Formula for Planar Graphs

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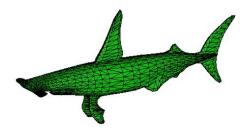
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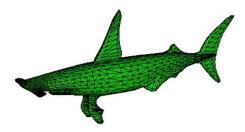
Proof: $K_{3,3}$ is triangle-free and has six vertices and nine edges. If it were planar then, by Cor. 307, it could have at most $2 \cdot 6 - 4 = 8$ edges. Thus, $K_{3,3}$ is non-planar.

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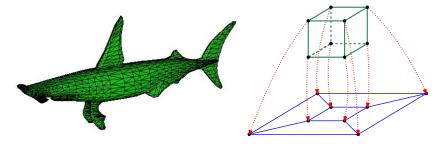


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The vertices and edges of a simple (bounded) polyhedron form a planar graph.



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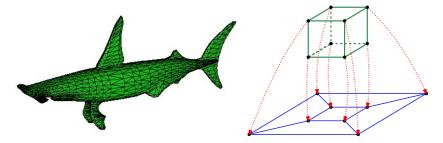


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A simple (bounded) polyhedron with *n* vertices has at most 3n - 6 edges and 2n - 4 faces.

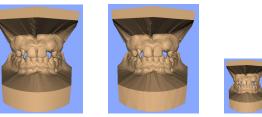
Real-World Application: Reducing the Face Count

- Recent improvements in laser rangefinder technology allow the digitization of the shapes of physical objects at extremely high resolutions.
- The resulting polyhedral models are huge: E.g., a 0.25 mm model of Michelangelo's 5-meter statue of David contains about 1 billion polygonal faces!



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- E.g., the left dental model has 424 376 faces, while the other two models have only a few thousand faces.



[Image credit: Michael Garland, Eurographics'99 STAR]



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- Edge contraction is one of the techniques used for reducing the face count.







[Image credit: Michael Garland, Eurographics'99 STAR]





Graph Theory

- What is a (Directed) Graph?
- Trees
- Special Graphs
- Graph Coloring



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Lemma 314

The chromatic number of a graph \mathcal{G} is two if and only if \mathcal{G} is bipartite.



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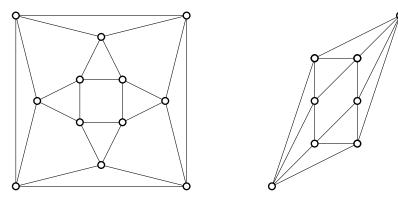
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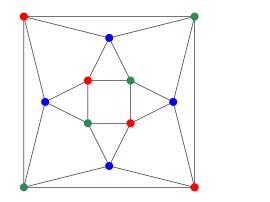
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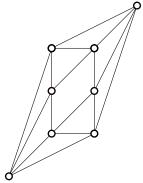
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- In 1996, Robertson et alii reduced the number of computer-checked cases to 633.
- In 2005, Werner and Gonthier used a general-purpose proof assistant ("Coq") to prove the theorem.



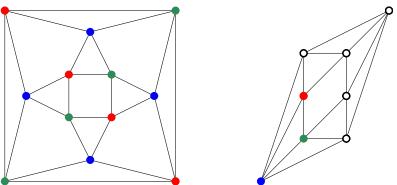




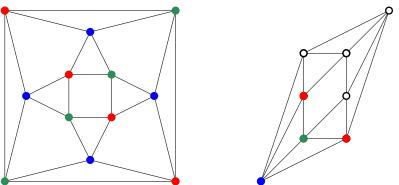




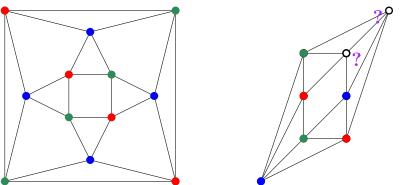






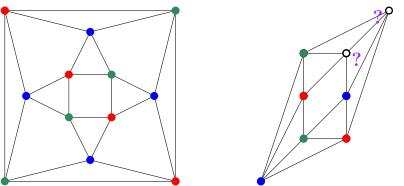






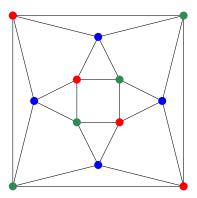


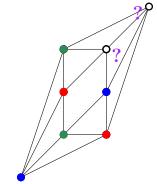
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- However, fairly efficient heuristics exist for approximate graph coloring.









[Image credit: Wikipedia]

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If every entity of a topographic map is a connected area then four colors suffice to color the map such that no two entities that share a common border (other than a common point) are colored with the same color.





[Image credit: Wikipedia]

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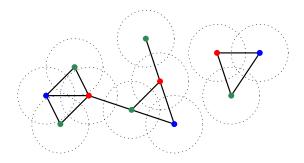
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If every entity of a topographic map is a connected area then four colors suffice to color the map such that no two entities that share a common border (other than a common point) are colored with the same color.

Note that this result holds only in the plane! E.g., on the surface of a torus seven colors are sufficient and may be necessary.

Real-World Application: Channel Assignment

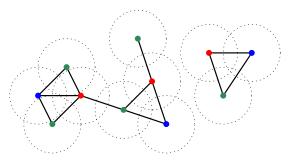
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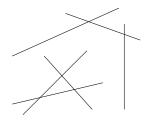
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- Obviously, the chromatic number of that graph equals the minimum number of frequencies needed.





Real-World Application: Minimum Plane Partition

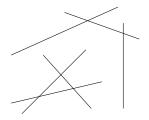
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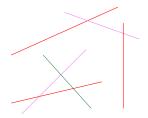
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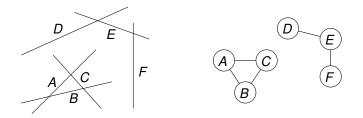
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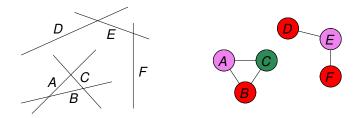
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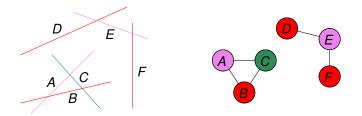
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- Other applications of graph coloring:
 - Scheduling consumer-producer interactions to allow concurrency.
 - Sudoku puzzles.



Cryptography

- Introduction
- Symmetric-Key Cryptography
- Public-Key Cryptography





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Symmetric-Key Cryptography (SKC): The same secret key is used for both encryption and decryption; aka secret-key cryptography.

Public-Key Cryptography (PKC): Different keys are used for encryption and decryption, with some keys being known publicly; aka asymmetric-key cryptography.

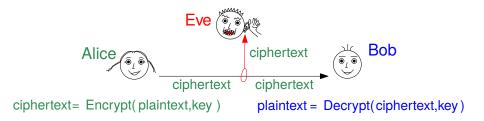


Basic Terms

- Plaintext original message.
- Ciphertext encoded/encrypted message.
- Encryption generating ciphertext from plaintext.
- Decryption / Deciphering generating plaintext from ciphertext.
- Cryptanalysis trying to break the encryption by applying various methods.
- Adversary, Spy the message thief.
- Eavesdropper a secret listener who listens to private conversations.
- Authentication the process of proving one's identity.
- *Privacy* ensuring that the message is read only by the intended receiver. (GnuPG: "Privacy is not a crime!")

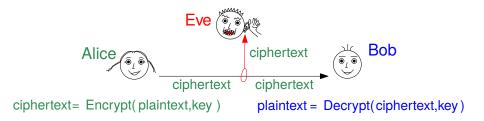


Eavesdropper Attacks





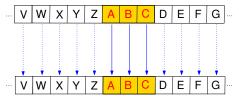
Eavesdropper Attacks



- Eve might attempt to
 - break the encryption,
 - replay the encrypted message (e.g., login) without breaking the encryption,
 - modify the message,
 - block the message,
 - fabricate a new message.

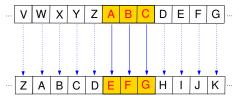


• According to Suetonius, Caesar (100–44 BCE) used an encryption scheme (for communication with his generals) that shifted the alphabet of the plaintext by some fixed position value *n*.



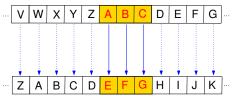


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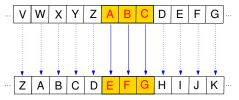
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- With *n* := 4:
 - Plaintext: alea iacta est
 - Ciphertext: epie megxe iwx



 According to Suetonius, Caesar (100–44 BCE) used an encryption scheme (for communication with his generals) that shifted the alphabet of the plaintext by some fixed position value n.



With n := 4:

Plaintext: alea iacta est Ciphertext: epie megxe iwx

- Suppose that the (Roman) letters are mapped to the numbers 0, 1, ..., 25.
- Then Caesar's encryption and decryption with shift *n* can be computed as follows:

ciphertext := $Encrypt_n(plaintext)$ = $(plaintext + n) \mod 26$

plaintext := $\text{Decrypt}_n(\text{ciphertext}) = (\text{ciphertext} - n) \mod 26$



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- It is broken easily by means of frequency analysis and brute-force attacks it offers no security by today's standards!
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- However, it is used within more complex systems, e.g., the Vigenère cipher.



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- Nevertheless, Caesar's cipher with *n* := 13, aka ROT13, has been (mis-)used for serious applications even rather recently.
- However, it is used within more complex systems, e.g., the Vigenère cipher.
- On a Unix machine, the tr utility can be used for carrying out Caesar's cipher. E.g.,

```
echo "alea iacta est" | tr 'A-Za-z' 'E-ZA-De-za-d' yields
```

```
epie megxe iwx,
```

and

```
echo "epie megxe iwx" | tr 'E-ZA-De-za-d' 'A-Za-z' recovers the original text
```

alea iacta est.



8 Cryptography

- Introduction
- Symmetric-Key Cryptography
- Public-Key Cryptography



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Simple example: Suppose that Alice wants to encrypt a bit string *A*. Then Alice and Bob could choose a secret key *B* and apply a bit-wise XOR (exclusive OR, ⊕) — an output bit is 1 if exactly one of the two input bits is 1 — in order to transmit *A* ⊕ *B*. Then Bob would compute (*A* ⊕ *B*) ⊕ *B* and, thus, retrieve *A*.



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| 0 | 1 | 1 | 0 |
| 1 | 0 | 1 | 1 |
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- The key distribution problem is a major roadblock on the road to secure communication among folks who do not meet regularly.
- A second big disadvantage is the need for multiple keys in order to encrypt messages intended for different receivers.



Cryptography

- Introduction
- Symmetric-Key Cryptography
- Public-Key Cryptography
 - Diffie-Hellman Algorithm
 - RSA



Public-Key Cryptography

• A pair of keys is used to encrypt and decrypt the messages, with one key being public.



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 - multiplication versus factorization ("factorization problem"):
 - If f(a,b) := a ⋅ b, then
 f(a,b) = 533 for a := 13, b := 41;
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- exponentiation versus logarithms ("discrete log problem"):
 - If $f(a,b) := a^b$, then
 - f(a,b) = 243 for a := 3, b := 5;
 - Again, finding *a* and *b* such that $\log_a 243 = b$ is considerably more difficult.



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 - If $f(a, b) := a^b$, then
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 - Again, finding *a* and *b* such that $\log_a 243 = b$ is considerably more difficult.
- Diffie and his advisor Hellman were the first to publish a PKC scheme in 1976 (They were the recipients of the 2015 ACM Turing Award.)

Sac

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$$\begin{array}{c|c} Alice & Bob \\\hline (1) & selects s with 1 < s < p - 1 & selects t with 1 < t < p - 1 \\ \end{array}$$



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| (1) | selects s with $1 < s < p - 1$ | selects t with $1 < t < p - 1$ |
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$$T^s = (g^t \mod p)^s$$



Alice and Bob share two public numbers: a (large) prime number p ∈ P and a so-called generator g ∈ {2,3,...,p-1} such that for every n ∈ {1,2,...,p-1} there exists a k ∈ N with n = g^k mod p.

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Hence, $k := T^s \mod p = S^t \mod p$ can be used as a common key by Alice and Bob.



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• In general, the public information is p, g, S and T, while s and t are secret.



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 Same for t. At present, nobody knows how to solve this problem efficiently.



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- To find *s*, Eve could attempt to solve the discrete log problem $S = g^s \mod p$. Same for *t*. At present, nobody knows how to solve this problem efficiently.
- Diffie-Hellman key exchange is used by the Tor system to set-up secure communication links with onion routers.
- The Diffie-Hellman key exchange is vulnerable to man-in-the-middle attacks.

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• Alice and Bob make *p* := 13 and *g* := 2 public.



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| | | |

• Alice and Bob make p := 13 and g := 2 public. The number 2 is indeed a generator modulo 13 because the following powers of two taken modulo 13 yield the integers 1, 2, ..., 12: $2^{12}, 2^1, 2^4, 2^2, 2^9, 2^5, 2^{11}, 2^{15}, 2^8, 2^{10}, 2^7, 2^6$.



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- Finally, $T^s \mod p = 12^5 \mod 13 = (12140 \cdot 13 + 12) \mod 13 = 12$, and $S^t \mod p = 6^6 \mod 13 = (3588 \cdot 13 + 12) \mod 13 = 12$.
- Hence, Alice and Bob have managed to exchange 12 as a master key for their future communication.



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No toy numbers!

Of course, in practice considerably larger values are chosen for p!!

Number Theory and Cryptography

Lemma 317

Let $a, b \in \mathbb{N}$ such that gcd(a, b) = 1. Then there exists $x \in \mathbb{Z}$ such that $a \cdot x \equiv_b 1$.



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Proof: Since gcd(a,b) = 1, Cor. 125 tells us that there exist $x, y \in \mathbb{Z}$ such that $a \cdot x + b \cdot y = 1$. Hence, $a \cdot x = 1 - b \cdot y \equiv_b 1$.



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Definition 318 (Euler's Totient Function, Dt.: Eulersche φ -Funktion)

Euler's totient function $\varphi : \mathbb{N} \to \mathbb{N}$ is defined as

 $\varphi(n) := |U_n|, \text{ with } U_n := \{x \in \mathbb{N} \colon 1 \leq x \leq n \land \gcd(x, n) = 1\}.$

The set U_n is called the group of units of n.



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• Hence, $\varphi(n)$ is the number of integers among 1, 2, ..., *n* that are coprime to *n*.

• More generally, $\varphi(p) = p - 1$ for every $p \in \mathbb{P}$.



Number Theory and Cryptography

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Let $p, q \in \mathbb{P}$. If $p \neq q$ then $\varphi(pq) = (p-1)(q-1)$.



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Proof: There are *q* multiples of *p* and *p* multiples of *q* within $\{1, 2, ..., pq\}$, and the only common multiple of both *p* and *q* is *pq*. Hence, by the Inclusion-Exclusion Principle (Thm. 167), $\varphi(pq) = pq - p - q + 1 = (p - 1)(q - 1)$.



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Lemma 320 (Fermat/Euler)

Let $n \in \mathbb{N}$ and $x \in U_n$. Then $x^{\varphi(n)} \equiv_n 1$.



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Corollary 321

Let $n \in \mathbb{N}$ and $x \in U_n$. If n = pq, with $p, q \in \mathbb{P}$ and $p \neq q$, then $x^{(p-1)(q-1)} \equiv_n 1$.



 The RSA system [Rivest, Shamir, Adleman 1977] makes use of Lemma 320 and of the fact that state-of-the-art factorization methods take far too long for products of numbers with several hundred digits each.



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 - Select two distinct prime numbers p and q (each of which, in practice, has at least 150 digits) and compute $n = p \cdot q$.
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 - Select an integer $e \in \mathbb{N} \setminus \{1\}$ such that $gcd(e, \varphi(n)) = 1$.
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 - The numbers *n* and *e* are published (Bob's *public key*).
 - Compute a number d which is the inverse of e in Z_{φ(n)}, i.e., such that
 - $d \cdot e \equiv_{\varphi(n)} 1$. (Such a number exists due to Lem. 317.)
 - The number *d* is called Bob's *private key* and is kept secret.



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| | $7^{13} = 96889010407 = 1761618371 \cdot 55 + 2 \equiv_{55} 2 =: z$ |
| Caesar: | $2^9 = 512 = 9 \cdot 55 + 17 \equiv_{55} 17 =: y$ |
| | $17^9 = 118587876497 = 2156143209 \cdot 55 + 2 \equiv_{55} 2 =: z$ |



- Note that there are $\varphi(n)$ many messages that can be sent for *n* given.
- Since

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- An eavesdropper who only knows *n*, *e*, and *y* cannot do much with this information. In particular, no efficient algorithm is known to factor *n* into *p*, *q* as a simple means to obtain φ(*n*).
- It is also important to ensure that x^e > n, i.e., that y is obtained by exponentiation and then by a reduction modulo n.
 - If $x^e < n$ then one could simply recover x by taking the *e*-th root of y. (After all, *e* is known publicly!)
 - Hence, it is wise to select *e* such that $2^e > n$.



The End!

I hope that you enjoyed this course, and I wish you all the best for your future studies.



