All Cheat Sheets

Machine Learning, Deep Learning, Artificial Intelligence

BYSTANFORD UNIVERSITY AND MASSACHUSETTS INSTITUTE OF **TECHNOLOGY**





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Probability—the Science of Uncertainty and Data

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PROBABILITY

Probability models and axioms

Definition (Sample space) A sample space Ω is the set of all possible outcomes. The set's elements must be mutually exclusive, collectively exhaustive and at the right granularity.

Definition (Event) An event is a subset of the sample space. Probability is assigned to events.

Definition (Probability axioms) A probability law \mathbb{P} assigns probabilities to events and satisfies the following axioms:

Nonnegativity $\mathbb{P}(A) \geq 0$ for all events A.

Normalization $\mathbb{P}(\Omega) = 1$.

(Countable) additivity For every sequence of events A_1, A_2, \ldots such that $A_i \cap A_j = \emptyset$: $\mathbb{P}\left(\bigcup_i A_i\right) = \sum_i \mathbb{P}(A_i)$.

Corollaries (Consequences of the axioms)

- $\mathbb{P}(\emptyset) = 0$.
- For any finite collection of disjoint events A_1, \ldots, A_n , $\mathbb{P}\left(\bigcup_{i=1}^n A_i\right) = \sum_{i=1}^n \mathbb{P}(A_i).$
- $\mathbb{P}(A) + \mathbb{P}(A^c) = 1$.
- $\mathbb{P}(A) \leq 1$.
- If $A \subset B$, then $\mathbb{P}(A) \leq \mathbb{P}(B)$.
- $\mathbb{P}(A \cup B) = \mathbb{P}(A) + \mathbb{P}(B) \mathbb{P}(A \cap B)$.
- $\mathbb{P}(A \cup B) \leq \mathbb{P}(A) + \mathbb{P}(B)$.

Example (Discrete uniform law) Assume Ω is finite and consists of n equally likely elements. Also, assume that $A \subset \Omega$ with k elements. Then $\mathbb{P}(A) = \frac{k}{n}$.

Conditioning and Bayes' rule

Definition (Conditional probability) Given that event B has occurred and that P(B) > 0, the probability that A occurs is

$$\mathbb{P}(A|B) \stackrel{\triangle}{=} \frac{\mathbb{P}(A \cap B)}{\mathbb{P}(B)}.$$

Remark (Conditional probabilities properties) They are the same as ordinary probabilities. Assuming $\mathbb{P}(B) > 0$:

- $\mathbb{P}(A|B) \geq 0$.
- $\mathbb{P}(\Omega|B) = 1$
- $\mathbb{P}(B|B) = 1$.
- If $A \cap C = \emptyset$, $\mathbb{P}(A \cup C|B) = \mathbb{P}(A|B) + \mathbb{P}(C|B)$.

Proposition (Multiplication rule)

 $\mathbb{P}(A_1 \cap A_2 \cap \cdots \cap A_n) = \mathbb{P}(A_1) \cdot \mathbb{P}(A_2 | A_1) \cdots \mathbb{P}(A_n | A_1 \cap A_2 \cap \cdots \cap A_{n-1}).$

Theorem (Total probability theorem) Given a partition $\{A_1, A_2, \ldots\}$ of the sample space, meaning that $\bigcup_i A_i = \Omega$ and the events are disjoint, and for every event B, we have

$$\mathbb{P}(B) = \sum_{i} \mathbb{P}(A_i) \mathbb{P}(B|A_i)$$

Theorem (Bayes' rule) Given a partition $\{A_1,A_2,\ldots\}$ of the sample space, meaning that $\bigcup A_i=\Omega$ and the events are disjoint,

and if $\mathbb{P}(A_i) > 0$ for all i, then for every event B, the conditional probabilities $\mathbb{P}(A_i|B)$ can be obtained from the conditional probabilities $\mathbb{P}(B|A_i)$ and the initial probabilities $\mathbb{P}(A_i)$ as follows:

$$\mathbb{P}(A_i|B) = \frac{\mathbb{P}(A_i)\mathbb{P}(B|A_i)}{\sum_i \mathbb{P}(A_i)\mathbb{P}(B|A_i)}.$$

Independence

Definition (Independence of events) Two events are independent if occurrence of one provides no information about the other. We say that A and B are independent if

$$\mathbb{P}(A \cap B) = \mathbb{P}(A)\mathbb{P}(B).$$

Equivalently, as long as $\mathbb{P}(A) > 0$ and $\mathbb{P}(B) > 0$,

$$\mathbb{P}(B|A) = \mathbb{P}(B)$$
 $\mathbb{P}(A|B) = \mathbb{P}(A)$

Remarks

- The definition of independence is symmetric with respect to A and B.
- The product definition applies even if $\mathbb{P}(A) = 0$ or $\mathbb{P}(B) = 0$.

Corollary If A and B are independent, then A and B^c are independent. Similarly for A^c and B, or for A^c and B^c .

Definition (Conditional independence) We say that A and B are independent conditioned on C, where $\mathbb{P}(C) > 0$, if

$$\mathbb{P}(A \cap B|C) = \mathbb{P}(A|C)\mathbb{P}(B|C).$$

Definition (Independence of a collection of events) We say that events A_1, A_2, \ldots, A_n are independent if for every collection of distinct indices i_1, i_2, \ldots, i_k , we have

$$\mathbb{P}(A_{i_1} \cap \ldots \cap A_{i_k}) = \mathbb{P}(A_{i_1}) \cdot \mathbb{P}(A_{i_2}) \cdots \mathbb{P}(A_{i_k}).$$

Counting

This section deals with finite sets with uniform probability law. In this case, to calculate $\mathbb{P}(A)$, we need to count the number of elements in A and in Ω .

Remark (Basic counting principle) For a selection that can be done in r stages, with n_i choices at each stage i, the number of possible selections is $n_1 \cdot n_2 \cdots n_r$.

Definition (Permutations) The number of permutations (orderings) of n different elements is

$$n! = 1 \cdot 2 \cdot 3 \cdots n$$
.

Definition (Combinations) Given a set of n elements, the number of subsets with exactly k elements is

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

Definition (Partitions) We are given an n-element set and nonnegative integers n_1, n_2, \ldots, n_r , whose sum is equal to n. The number of partitions of the set into r disjoint subsets, with the ith subset containing exactly n_i elements, is equal to

$$\binom{n}{n_1,\ldots,n_r} = \frac{n!}{n_1!n_2!\cdots n_r!}.$$

Remark This is the same as counting how to assign n distinct elements to r people, giving each person i exactly n_i elements.

Discrete random variables

Probability mass function and expectation

Definition (Random variable) A random variable X is a function of the sample space Ω into the real numbers (or \mathbb{R}^n). Its range can be discrete or continuous.

Definition (Probability mass funtion (PMF)) The probability law of a discrete random variable X is called its PMF. It is defined as

$$p_X(x) = \mathbb{P}(X = x) = \mathbb{P}\left(\left\{\omega \in \Omega : X(\omega) = x\right\}\right).$$

Properties

 $p_X(x) \ge 0, \ \forall \ x.$

 $\sum_{x} p_X(x) = 1.$

Example (Bernoulli random variable) A Bernoulli random variable X with parameter $0 \le p \le 1$ ($X \sim \text{Ber}(p)$) takes the following values:

$$X = \begin{cases} 1 & \text{w.p. } p, \\ 0 & \text{w.p. } 1 - p. \end{cases}$$

An indicator random variable of an event (I_A = 1 if A occurs) is an example of a Bernoulli random variable.

Example (Discrete uniform random variable) A Discrete uniform random variable X between a and b with $a \le b$ ($X \sim \operatorname{Uni}[a,b]$) takes any of the values in $\{a,a+1,\ldots,b\}$ with probability $\frac{1}{b-a+1}$. Example (Binomial random variable) A Binomial random

variable X with parameters n (natural number) and $0 \le p \le 1$ $(X \sim \text{Bin}(n, p))$ takes values in the set $\{0, 1, \ldots, n\}$ with probabilities $p_X(i) = \binom{n}{i} p^i (1-p)^{n-i}$.

It represents the number of successes in n independent trials where each trial has a probability of success p. Therefore, it can also be seen as the sum of n independent Bernoulli random variables, each with parameter p.

Example (Geometric random variable) A Geometric random variable X with parameter $0 \le p \le 1$ ($X \sim \text{Geo}(p)$) takes values in the set $\{1,2,\ldots\}$ with probabilities $p_X(i) = (1-p)^{i-1}p$. It represents the number of independent trials until (and including)

the first success, when the probability of success in each trial is p.

Definition (Expectation/mean of a random variable) The expectation of a discrete random variable is defined as

$$\mathbb{E}[X] \stackrel{\triangle}{=} \sum_{x} x p_X(x).$$

assuming $\sum_{x} |x| p_X(x) < \infty$.

Properties (Properties of expectation)

- If $X \ge 0$ then $\mathbb{E}[X] \ge 0$.
- If $a \le X \le b$ then $a \le \mathbb{E}[X] \le b$.
- If X = c then $\mathbb{E}[X] = c$.

Example Expected value of know r.v.

- If $X \sim \operatorname{Ber}(p)$ then $\mathbb{E}[X] = p$.
- If $X = I_A$ then $\mathbb{E}[X] = \mathbb{P}(A)$.
- If $X \sim \text{Uni}[a, b]$ then $\mathbb{E}[X] = \frac{a+b}{2}$.
- If $X \sim \text{Bin}(n, p)$ then $\mathbb{E}[X] = np$.
- If $X \sim \text{Geo}(p)$ then $\mathbb{E}[X] = \frac{1}{p}$.

Theorem (Expected value rule) Given a random variable X and a Properties (Properties of joint PMF) function $q: \mathbb{R} \to \mathbb{R}$, we construct the random variable Y = q(X). Then

$$\sum_{y} y p_{Y}(y) = \mathbb{E}[Y] = \mathbb{E}[g(X)] = \sum_{x} g(x) p_{X}(x).$$

Remark (PMF of Y = q(X)) The PMF of Y = q(X) is $p_Y(y) = \sum_{x:g(x)=y} p_X(x).$

Remark In general $g(\mathbb{E}[X]) \neq \mathbb{E}[g(X)]$. They are equal if a(x) = ax + b.

Variance, conditioning on an event, multiple r.v.

Definition (Variance of a random variable) Given a random variable X with $\mu = \mathbb{E}[X]$, its variance is a measure of the spread of the random variable and is defined as

$$\operatorname{Var}(X) \stackrel{\triangle}{=} \mathbb{E}\left[(X - \mu)^2\right] = \sum_{x} (x - \mu)^2 p_X(x).$$

$$\sigma_X = \sqrt{\operatorname{Var}(X)}$$
.

Properties (Properties of the variance)

- $Var(aX) = a^2 Var(X)$, for all $a \in \mathbb{R}$.
- Var(X + b) = Var(X), for all $b \in \mathbb{R}$.
- $\operatorname{Var}(aX + b) = a^2 \operatorname{Var}(X)$.
- $\operatorname{Var}(X) = \mathbb{E}[X^2] (\mathbb{E}[X])^2$.

Example (Variance of known r.v.)

- If $X \sim \text{Ber}(p)$, then Var(X) = p(1-p).
- If $X \sim \text{Uni}[a,b]$, then $\text{Var}(X) = \frac{(b-a)(b-a+2)}{12}$.
- If $X \sim \text{Bin}(n, p)$, then Var(X) = np(1 p).
- If $X \sim \text{Geo}(p)$, then $\text{Var}(X) = \frac{1-p}{2}$

Proposition (Conditional PMF and expectation, given an event) Given the event A, with $\mathbb{P}(A) > 0$, we have the following

- $p_{X|A}(x) = \mathbb{P}(X = x|A)$.
- If A is a subset of the range of X, then: $p_{X|A}(x) \stackrel{\triangle}{=} p_{X|\{X \in A\}}(x) = \begin{cases} \frac{1}{P(A)} p_X(x), & \text{if } x \in A, \\ 0, & \text{otherwise.} \end{cases}$
- $\sum_{x} p_{X|A}(x) = 1$.
- $\mathbb{E}[X|A] = \sum_{x} x p_{X|A}(x)$.
- $\mathbb{E}[g(X)|A] = \sum_{x} g(x) p_{X|A}(x)$.

Proposition (Total expectation rule) Given a partition of disjoint events A_1, \ldots, A_n such that $\sum_i \mathbb{P}(A_i) = 1$, and $\mathbb{P}(A_i) > 0$,

$$\mathbb{E}[X] = \mathbb{P}(A_1)\mathbb{E}[X|A_1] + \dots + \mathbb{P}(A_n)\mathbb{E}[X|A_n].$$

Definition (Memorylessness of the geometric random variable)

When we condition a geometric random variable X on the event X > n we have memorylessness, meaning that the "remaining time" X-n, given that X>n, is also geometric with the same parameter. Formally,

$$p_{X-n|X>n}(i) = p_X(i).$$

Definition (Joint PMF) The joint PMF of random variables $X_1, X_2, ..., X_n$ is $p_{X_1,X_2,...,X_n}(x_1,...,x_n) = \mathbb{P}(X_1 = x_1,...,X_n = x_n).$

- $\bullet \sum_{x_1} \cdots \sum_{x_n} p_{X_1, \dots, X_n} (x_1, \dots, x_n) = 1.$
- $p_{X_1}(x_1) = \sum_{x_2} \cdots \sum_{x_n} p_{X_1,...,X_n}(x_1,x_2,...,x_n).$
- $p_{X_2,...,X_n}(x_2,...,x_n) = \sum p_{X_1,X_2,...,X_n}(x_1,x_2,...,x_n)$.

Definition (Functions of multiple r.v.) If $Z = g(X_1, \ldots, X_n)$, where $g: \mathbb{R}^n \to \mathbb{R}$, then $p_Z(z) = \mathbb{P}(g(X_1, \dots, X_n) = z)$.

Proposition (Expected value rule for multiple r.v.) Given

$$\mathbb{E}\left[g(X_1,\ldots,X_n)\right] = \sum_{x_1,\ldots,x_n} g(x_1,\ldots,x_n) p_{X_1,\ldots,X_n}(x_1,\ldots,x_n).$$

Properties (Linearity of expectations)

- $\mathbb{E}[aX + b] = a\mathbb{E}[X] + b$.
- $\mathbb{E}[X_1 + \dots + X_n] = \mathbb{E}[X_1] + \dots + \mathbb{E}[X_n].$

Conditioning on a random variable, independence

Definition (Conditional PMF given another random variable)

Given discrete random variables X, Y and y such that $p_Y(y) > 0$ we define

$$p_{X|Y}(x|y) \stackrel{\triangle}{=} \frac{p_{X,Y}(x,y)}{p_Y(y)}.$$

Proposition (Multiplication rule) Given jointly discrete random variables X, Y, and whenever the conditional probabilities are defined,

$$p_{X,Y}(x,y) = p_X(x)p_{Y|X}(y|x) = p_Y(y)p_{X|Y}(x|y).$$

Definition (Conditional expectation) Given discrete random variables X, Y and y such that $p_Y(y) > 0$ we define

$$\mathbb{E}[X|Y = y] = \sum_{x} x p_{X|Y}(x|y)$$

Additionally we have

$$\mathbb{E}\left[g(X)|Y=y\right] = \sum_{x} g(x) p_{X|Y}(x|y).$$

Theorem (Total probability and expectation theorems) If $p_Y(y) > 0$, then

$$p_X(x) = \sum_{y} p_Y(y) p_{X|Y}(x|y),$$

$$\mathbb{E}[X] = \sum_{y} p_{Y}(y) \mathbb{E}[X|Y = y].$$

Definition (Independence of a random variable and an event) A discrete random variable X and an event A are independent if $\mathbb{P}(X = x \text{ and } A) = p_X(x)\mathbb{P}(A), \text{ for all } x.$

Definition (Independence of two random variables) Two discrete random variables X and Y are independent if $p_{X,Y}(x,y) = p_X(x)p_Y(y)$ for all x,y.

Remark (Independence of a collection of random variables) A collection X_1, X_2, \dots, X_n of random variables are independent if

$$p_{X_1,...,X_n}(x_1,...,x_n) = p_{X_1}(x_1) \cdots p_{X_n}(x_n), \forall x_1,...,x_n.$$

Remark (Independence and expectation) In general, $\mathbb{E}[g(X,Y)] \neq g(\mathbb{E}[X],\mathbb{E}[Y])$. An exception is for linear functions: $\mathbb{E}[aX + bY] = a\mathbb{E}[X] + b\mathbb{E}[Y].$

Proposition (Expectation of product of independent r.v.) If X and Y are discrete independent random variables,

$$\mathbb{E}[XY] = \mathbb{E}[X]\mathbb{E}[Y].$$

Remark If X and Y are independent, $\mathbb{E}[g(X)h(Y)] = \mathbb{E}[g(X)]\mathbb{E}[h(Y)].$

Proposition (Variance of sum of independent random variables) IF X and Y are discrete independent random variables.

$$Var(X + Y) = Var(X) + Var(Y).$$

Continuous random variables

PDF. Expectation. Variance. CDF

Definition (Probability density function (PDF)) A probability density function of a r.v. X is a non-negative real valued function f_X that satisfies the following

- $\bullet \int_{-\infty}^{\infty} f_X(x) dx = 1.$
- $\mathbb{P}(a \le X \le b) = \int_{a}^{b} f_X(x) dx$ for some random variable X.

Definition (Continuous random variable) A random variable X is continuous if its probability law can be described by a PDF f_X . Remark Continuous random variables satisfy:

- For small $\delta > 0$, $\mathbb{P}(a \le X \le a + \delta) \approx f_X(a)\delta$.
- $\mathbb{P}(X = a) = 0, \forall a \in \mathbb{R}.$

Definition (Expectation of a continuous random variable) The expectation of a continuous random variable is

$$\mathbb{E}[X] \stackrel{\triangle}{=} \int_{-\infty}^{\infty} x f_X(x) \mathrm{d}x.$$

assuming $\int_{0}^{\infty} |x| f_X(x) dx < \infty$.

Properties (Properties of expectation)

- If $X \ge 0$ then $\mathbb{E}[X] \ge 0$.
- If $a \le X \le b$ then $a \le \mathbb{E}[X] \le b$.
- $\mathbb{E}[g(X)] = \int_{-\infty}^{\infty} g(x) f_X(x) dx$.
- $\mathbb{E}[aX + b] = a\mathbb{E}[X] + b$.

Definition (Variance of a continuous random variable) Given a continuous random variable X with $\mu = \mathbb{E}[X]$, its variance is

$$\operatorname{Var}(X) = \mathbb{E}\left[(X - \mu)^2\right] = \int_{-\infty}^{\infty} (x - \mu)^2 f_X(x) dx.$$

It has the same properties as the variance of a discrete random variable.

Example (Uniform continuous random variable) A Uniform continuous random variable X between a and b, with a < b, $(X \sim \text{Uni}(a,b))$ has PDF

$$f_X(x) = \begin{cases} \frac{1}{b-a}, & \text{if } a < x < b, \\ 0, & \text{otherwise.} \end{cases}$$

We have $\mathbb{E}[X] = \frac{a+b}{2}$ and $\operatorname{Var}(X) = \frac{(b-a)^2}{12}$.

Example (Exponential random variable) An Exponential random variable X with parameter $\lambda > 0$ ($X \sim Exp(\lambda)$) has PDF

$$f_X(x) = \begin{cases} \lambda e^{-\lambda x}, & \text{if } x \ge 0, \\ 0, & \text{otherwise.} \end{cases}$$

We have $E[X] = \frac{1}{\lambda}$ and $Var(X) = \frac{1}{\lambda^2}$.

Definition (Cumulative Distribution Function (CDF)) The CDF of a random variable X is $F_X(x) = \mathbb{P}(X \le x)$.

In particular, for a continuous random variable, we have

$$F_X(x) = \int_{-\infty}^{x} f_X(x) dx,$$
$$f_X(x) = \frac{dF_X(x)}{dx}.$$

Properties (Properties of CDF)

- If $y \ge x$, then $F_X(y) \ge F_X(x)$.
- $\bullet \lim_{x \to -\infty} F_X(x) = 0.$
- $\lim_{x \to \infty} F_X(x) = 1$.

Definition (Normal/Gaussian random variable) A Normal random variable X with mean μ and variance $\sigma^2 > 0$ ($X \sim \mathcal{N}(\mu, \sigma^2)$) has PDF

$$f_X(x) = \frac{1}{\sqrt{2\pi\sigma^2}} e^{-\frac{1}{2\sigma^2}(x-\mu)^2}.$$

We have $E[X] = \mu$ and $Var(X) = \sigma^2$.

Remark (Standard Normal) The standard Normal is $\mathcal{N}(0,1)$.

Proposition (Linearity of Gaussians) Given $X \sim \mathcal{N}(\mu, \sigma^2)$, and if $a \neq 0$, then $aX + b \sim \mathcal{N}(a\mu + b, a^2\sigma^2)$.

Using this $Y = \frac{X - \mu}{\sigma}$ is a standard gaussian.

Conditioning on an event, and multiple continuous r.v.

Definition (Conditional PDF given an event) Given a continuous random variable X and event A with P(A) > 0, we define the conditional PDF as the function that satisfies

$$\mathbb{P}(X \in B|A) = \int_B f_{X|A}(x) dx.$$

Definition (Conditional PDF given $X \in A$) Given a continuous random variable X and an $A \subset \mathbb{R}$, with P(A) > 0:

$$f_{X|X\in A}(x) = \begin{cases} \frac{1}{\mathbb{P}(A)} f_X(x), & x \in A, \\ 0, & x \notin A. \end{cases}$$

Definition (Conditional expectation) Given a continuous random variable X and an event A, with P(A) > 0:

$$\mathbb{E}[X|A] = \int_{-\infty}^{\infty} f_{X|A}(x) dx.$$

Definition (Memorylessness of the exponential random variable)

When we condition an exponential random variable X on the event X > t we have memorylessness, meaning that the "remaining time" X - t given that X > t is also geometric with the same parameter i.e.,

$$\mathbb{P}(X-t>x|X>t)=\mathbb{P}(X>x).$$

Theorem (Total probability and expectation theorems) Given a partition of the space into disjoint events A_1, A_2, \ldots, A_n such that $\sum_i \mathbb{P}(A_i) = 1$ we have the following:

$$\begin{split} F_X(x) &= \mathbb{P}(A_1) F_{X|A_1}(x) + \dots + \mathbb{P}(A_n) F_{X|A_n}(x), \\ f_X(x) &= \mathbb{P}(A_1) f_{X|A_1}(x) + \dots + \mathbb{P}(A_n) f_{X|A_n}(x), \\ \mathbb{E}[X] &= \mathbb{P}(A_1) \mathbb{E}[X|A_1] + \dots + \mathbb{P}(A_n) \mathbb{E}[X|A_n]. \end{split}$$

Definition (Jointly continuous random variables) A pair (collection) of random variables is jointly continuous if there exists a joint PDF $f_{X,Y}$ that describes them, that is, for every set $B \subset \mathbb{R}^n$

$$\mathbb{P}\left((X,Y)\in B\right) = \iint_B f_{X,Y}(x,y) \mathrm{d}x\mathrm{d}y.$$

Properties (Properties of joint PDFs)

•
$$f_X(x) = \int_{-\infty}^{\infty} f_{X,Y}(x,y) dy$$
.

•
$$F_{X,Y}(x,y) = \mathbb{P}(X \le x, Y \le y) = \int_{-\infty}^{x} \left[\int_{-\infty}^{y} f_{X,Y}(u,v) dv \right] du.$$

•
$$f_{X,Y}(x) = \frac{\partial^2 F_{X,Y}(x,y)}{\partial x \partial y}$$

Example (Uniform joint PDF on a set S) Let $S \subset \mathbb{R}^2$ with area s > 0, then the random variable (X, Y) is uniform over S if it has PDF

$$f_{X,Y}(x,y) = \begin{cases} \frac{1}{s}, & (x,y) \in S, \\ 0, & (x,y) \notin S. \end{cases}$$

Conditioning on a random variable, independence, Bayes' rule

Definition (Conditional PDF given another random variable)

Given jointly continuous random variables X, Y and a value y such that $f_Y(y) > 0$, we define the conditional PDF as

$$f_{X|Y}(x|y) \stackrel{\triangle}{=} \frac{f_{X,Y}(x,y)}{f_Y(y)}.$$

Additionally we define $\mathbb{P}(X \in A|Y = y) \int_A f_{X|Y}(x|y) dx$. Proposition (Multiplication rule) Given jointly continuous random variables X, Y, whenever possible we have

$$f_{X,Y}(x,y) = f_X(x)f_{Y|X}(y|x) = f_Y(y)f_{X|Y}(x|y).$$

Definition (Conditional expectation) Given jointly continuous random variables X, Y, and y such that $f_Y(y) > 0$, we define the conditional expected value as

$$\mathbb{E}[X|Y=y] = \int_{-\infty}^{\infty} x f_{X|Y}(x|y) dx.$$

Additionally we have

$$\mathbb{E}[g(X)|Y=y] = \int_{-\infty}^{\infty} g(x) f_{X|Y}(x|y) dx.$$

Theorem (Total probability and total expectation theorems)

$$f_X(x) = \int_{-\infty}^{\infty} f_Y(y) f_{X|Y}(x|y) dy,$$
$$\mathbb{E}[X] = \int_{-\infty}^{\infty} f_Y(y) \mathbb{E}[X|Y = y] dy.$$

Definition (Independence) Jointly continuous random variables X, Y are independent if $f_{X,Y}(x,y) = f_X(x)f_Y(y)$ for all x, y.

Proposition (Expectation of product of independent r.v.) If X and Y are independent continuous random variables,

$$\mathbb{E}[XY] = \mathbb{E}[X]\mathbb{E}[Y].$$

Remark If X and Y are independent, $\mathbb{E}[g(X)h(Y)] = \mathbb{E}[g(X)]\mathbb{E}[h(Y)].$

Proposition (Variance of sum of independent random variables) If X and Y are independent continuous random variables,

$$Var(X + Y) = Var(X) + Var(Y)$$
.

Proposition (Bayes' rule summary)

- For X, Y discrete: $p_{X|Y}(x|y) = \frac{p_X(x)p_{Y|X}(y|x)}{p_Y(y)}$.
- For X, Y continuous: $f_{X|Y}(x|y) = \frac{f_X(x)f_{Y|X}(y|x)}{f_Y(y)}$.
- For X discrete, Y continuous: $p_{X|Y}(x|y) = \frac{p_X(x)f_{Y|X}(y|x)}{f_Y(y)}$.
- For X continuous, Y discrete: $f_{X|Y}(x|y) = \frac{f_X(x)p_{Y|X}(y|x)}{p_Y(y)}$.

Derived distributions

Proposition (Discrete case) Given a discrete random variable X and a function g, the r.v. Y = g(X) has PMF

$$p_Y(y) = \sum_{x:g(x)=y} p_X(x).$$

Remark (Linear function of discrete random variable) If g(x) = ax + b, then $p_Y(y) = p_X\left(\frac{y-b}{a}\right)$.

Proposition (Linear function of continuous r.v.) Given a continuous random variable X and Y = aX + b, with $a \neq 0$, we have

$$f_Y(y) = \frac{1}{|a|} f_X\left(\frac{y-b}{a}\right).$$

Corollary (Linear function of normal r.v.) If $X \sim \mathcal{N}(\mu, \sigma^2)$ and Y = aX + b, with $a \neq 0$, then $Y \sim \mathcal{N}(a\mu + b, a^2\sigma^2)$.

Example (General function of a continuous r.v.) If X is a continuous random variable and g is any function, to obtain the pdf of Y = g(X) we follow the two-step procedure:

- 1. Find the CDF of Y: $F_Y(y) = \mathbb{P}(Y \le y) = \mathbb{P}(q(X) \le y)$.
- 2. Differentiate the CDF of Y to obtain the PDF: $f_Y(y) = \frac{dF_Y(y)}{dy}$.

Proposition (General formula for monotonic g) Let X be a continuous random variable and g a function that is monotonic wherever $f_X(x) > 0$. The PDF of Y = g(X) is given by

$$f_Y(y) = f_X(h(y)) \left| \frac{\mathrm{d}h}{\mathrm{d}y}(y) \right|.$$

where $h = g^{-1}$ in the interval where g is monotonic.

Sums of independent r.v., covariance and correlation

Proposition (Discrete case) Let X, Y be discrete independent random variables and Z = X + Y, then the PMF of Z is

$$p_Z(z) = \sum_x p_X(x) p_Y(z-x).$$

Proposition (Continuous case) Let X, Y be continuous independent random variables and Z = X + Y, then the PDF of Z is

$$f_Z(z) = \int_{-\infty}^{\infty} f_X(x) f_Y(z-x) dx.$$

Proposition (Sum of independent normal r.v.) Let $X \sim \mathcal{N}(\mu_x, \sigma_x^2)$ and $Y \sim \mathcal{N}(\mu_y, \sigma_y^2)$ independent. Then $Z = X + Y \sim \mathcal{N}(\mu_x + \mu_y, \sigma_x^2 + \sigma_y^2)$.

Definition (Covariance) We define the covariance of random variables X, Y as

$$Cov(X,Y) \stackrel{\triangle}{=} \mathbb{E}[(X - \mathbb{E}[X])(Y - \mathbb{E}[Y])].$$

Properties (Properties of covariance)

- If X, Y are independent, then Cov(X, Y) = 0.
- Cov(X, X) = Var(X)
- Cov(aX + b, Y) = a Cov(X, Y).
- Cov(X, Y + Z) = Cov(X, Y) + Cov(X, Z).
- $Cov(X, Y) = \mathbb{E}[XY] \mathbb{E}[X]\mathbb{E}[Y].$

Proposition (Variance of a sum of r.v.)

$$\operatorname{Var}(X_1+\cdots+X_n)=\sum_i\operatorname{Var}(X_i)+\sum_{i\neq j}\operatorname{Cov}(X_i,X_j).$$

Definition (Correlation coefficient) We define the correlation coefficient of random variables X, Y, with $\sigma_X, \sigma_Y > 0$, as

$$\rho(X,Y) \stackrel{\triangle}{=} \frac{\mathrm{Cov}(X,Y)}{\sigma_X \sigma_Y}.$$

Properties (Properties of the correlation coefficient)

- $-1 \le \rho \le 1$.
- If X, Y are independent, then $\rho = 0$.
- $|\rho| = 1$ if and only if $X \mathbb{E}[X] = c(Y \mathbb{E}[Y])$.
- $\rho(aX + b, Y) = \operatorname{sign}(a)\rho(X, Y)$.

Conditional expectation and variance, sum of random number of r.v.

Definition (Conditional expectation as a random variable) Given random variables X, Y the conditional expectation $\mathbb{E}[X|Y]$ is the random variable that takes the value $\mathbb{E}[X|Y=y]$ whenever Y=y. Theorem (Law of iterated expectations)

$$\mathbb{E}\left[\mathbb{E}[X|Y]\right] = \mathbb{E}[X].$$

Definition (Conditional variance as a random variable) Given random variables X, Y the conditional variance Var(X|Y) is the random variable that takes the value Var(X|Y=y) whenever Y=y.

Theorem (Law of total variance)

$$Var(X) = \mathbb{E}\left[Var(X|Y)\right] + Var\left(\mathbb{E}[X|Y]\right).$$

Proposition (Sum of a random number of independent r.v.)

Let N be a nonnegative integer random variable. Let X, X_1, X_2, \ldots, X_N be i.i.d. random variables. Let $Y = \sum_i X_i$. Then

$$\mathbb{E}[Y] = \mathbb{E}[N]\mathbb{E}[X],$$

$$Var(Y) = \mathbb{E}[N]Var(X) + (\mathbb{E}[X])^{2}Var(N).$$

Convergence of random variables

Inequalities, convergence, and the Weak Law of Large Numbers

Theorem (Markov inequality) Given a random variable $X \ge 0$ and, for every a > 0 we have

$$\mathbb{P}(X \ge a) \le \frac{\mathbb{E}[X]}{a}.$$

Theorem (Chebyshev inequality) Given a random variable X with $\mathbb{E}[X] = \mu$ and $\text{Var}(X) = \sigma^2$, for every $\epsilon > 0$ we have

$$\mathbb{P}\left(\left|X-\mu\right| \geq \epsilon\right) \leq \frac{\sigma^2}{\epsilon^2}.$$

Theorem (Weak Law of Large Number (WLLN)) Given a sequence of i.i.d. random variables $\{X_1, X_2, ...\}$ with $\mathbb{E}[X_i] = \mu$ and $\text{Var}(X_i) = \sigma^2$, we define

$$M_n = \frac{1}{n} \sum_{i=1}^n X_i,$$

for every $\epsilon > 0$ we have

$$\lim_{n\to\infty} \mathbb{P}\left(|M_n - \mu| \ge \epsilon\right) = 0.$$

Definition (Convergence in probability) A sequence of random variables $\{Y_i\}$ converges in probability to the random variable Y if

$$\lim_{n\to\infty} \mathbb{P}\left(|Y_i - Y| \ge \epsilon\right) = 0,$$

for every $\epsilon > 0$.

Properties (Properties of convergence in probability) If $X_n \to a$ and $Y_n \to b$ in probability, then

- $X_n + Y_n \rightarrow a + b$.
- If q is a continuous function, then $q(X_n) \to q(a)$.
- $\mathbb{E}[X_n]$ does not always converge to a.

The Central Limit Theorem

Theorem (Central Limit Theorem (CLT)) Given a sequence of independent random variables $\{X_1, X_2, ...\}$ with $\mathbb{E}[X_i] = \mu$ and $\text{Var}(X_i) = \sigma^2$, we define

$$Z_n = \frac{1}{\sigma\sqrt{n}} \sum_{i=1}^n (X_i - \mu).$$

Then, for every z, we have

$$\lim_{n\to\infty} \mathbb{P}(Z_n \le z) = \mathbb{P}(Z \le z),$$

where $Z \sim \mathcal{N}(0,1)$.

Corollary (Normal approximation of a binomial) Let $X \sim Bin(n,p)$ with n large. Then S_n can be approximated by $Z \sim \mathcal{N}(np, np(1-p))$.

Remark (De Moivre-Laplace 1/2 approximation) Let $X \sim Bin$, then $\mathbb{P}(X=i) = \mathbb{P}\left(i-\frac{1}{2} \leq X \leq i+\frac{1}{2}\right)$ and we can use the CLT to approximate the PMF of X.

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1 Supervised Learning

1.1 Introduction to Supervised Learning

Given a set of data points $\{x^{(1)},...,x^{(m)}\}$ associated to a set of outcomes $\{y^{(1)},...,y^{(m)}\}$, we want to build a classifier that learns how to predict y from x.

 \square Type of prediction – The different types of predictive models are summed up in the table below:

	Regression	Classifier	
Outcome	Continuous	Class	
Examples Linear regression		Logistic regression, SVM, Naive Bayes	

☐ **Type of model** – The different models are summed up in the table below:

	Discriminative model	Generative model	
Goal	Directly estimate $P(y x)$	Estimate $P(x y)$ to deduce $P(y x)$	
What's learned	Decision boundary	Probability distributions of the data	
Illustration			
Examples	Regressions, SVMs	GDA, Naive Bayes	

1.2 Notations and general concepts

 \square Hypothesis – The hypothesis is noted h_{θ} and is the model that we choose. For a given input data $x^{(i)}$, the model prediction output is $h_{\theta}(x^{(i)})$.

□ Loss function – A loss function is a function $L:(z,y) \in \mathbb{R} \times Y \longmapsto L(z,y) \in \mathbb{R}$ that takes as inputs the predicted value z corresponding to the real data value y and outputs how different they are. The common loss functions are summed up in the table below:

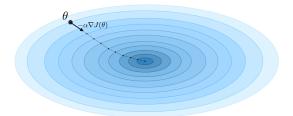
Least squared	Logistic	Hinge	Cross-entropy
$\frac{1}{2}(y-z)^2$	$\log(1 + \exp(-yz))$	$\max(0,1-yz)$	$-\left[y\log(z)+(1-y)\log(1-z)\right]$
$y\in\mathbb{R}$	y = -1 $y = 1$	y = -1 z y = 1	y = 0 $y = 1$ $y = 1$
Linear regression	Logistic regression	SVM	Neural Network

 \square Cost function – The cost function J is commonly used to assess the performance of a model, and is defined with the loss function L as follows:

$$J(\theta) = \sum_{i=1}^{m} L(h_{\theta}(x^{(i)}), y^{(i)})$$

□ Gradient descent – By noting $\alpha \in \mathbb{R}$ the learning rate, the update rule for gradient descent is expressed with the learning rate and the cost function J as follows:

$$\theta \longleftarrow \theta - \alpha \nabla J(\theta)$$



Remark: Stochastic gradient descent (SGD) is updating the parameter based on each training example, and batch gradient descent is on a batch of training examples.

□ Likelihood – The likelihood of a model $L(\theta)$ given parameters θ is used to find the optimal parameters θ through maximizing the likelihood. In practice, we use the log-likelihood $\ell(\theta) = \log(L(\theta))$ which is easier to optimize. We have:

$$\theta^{\text{opt}} = \underset{\theta}{\text{arg max } L(\theta)}$$

□ Newton's algorithm – The Newton's algorithm is a numerical method that finds θ such that $\ell'(\theta) = 0$. Its update rule is as follows:

$$\theta \leftarrow \theta - \frac{\ell'(\theta)}{\ell''(\theta)}$$

Remark: the multidimensional generalization, also known as the Newton-Raphson method, has the following update rule:

$$\theta \leftarrow \theta - \left(\nabla_{\theta}^2 \ell(\theta)\right)^{-1} \nabla_{\theta} \ell(\theta)$$

1.3 Linear models

1.3.1 Linear regression

We assume here that $y|x;\theta \sim \mathcal{N}(\mu,\sigma^2)$

 \square Normal equations – By noting X the matrix design, the value of θ that minimizes the cost function is a closed-form solution such that:

$$\theta = (X^T X)^{-1} X^T y$$

 \square LMS algorithm – By noting α the learning rate, the update rule of the Least Mean Squares (LMS) algorithm for a training set of m data points, which is also known as the Widrow-Hoff learning rule, is as follows:

$$\forall j, \quad \theta_j \leftarrow \theta_j + \alpha \sum_{i=1}^m \left[y^{(i)} - h_{\theta}(x^{(i)}) \right] x_j^{(i)}$$

Remark: the update rule is a particular case of the gradient ascent.

□ LWR - Locally Weighted Regression, also known as LWR, is a variant of linear regression that weights each training example in its cost function by $w^{(i)}(x)$, which is defined with parameter $\tau \in \mathbb{R}$ as:

$$w^{(i)}(x) = \exp\left(-\frac{(x^{(i)} - x)^2}{2\tau^2}\right)$$

1.3.2Classification and logistic regression

 \square Sigmoid function – The sigmoid function q, also known as the logistic function, is defined as follows:

$$\forall z \in \mathbb{R}, \quad g(z) = \frac{1}{1 + e^{-z}} \in]0,1[$$

 \square Logistic regression – We assume here that $y|x;\theta \sim \text{Bernoulli}(\phi)$. We have the following form:

$$\phi = p(y = 1|x; \theta) = \frac{1}{1 + \exp(-\theta^T x)} = g(\theta^T x)$$

Remark: there is no closed form solution for the case of logistic regressions.

□ Softmax regression – A softmax regression, also called a multiclass logistic regression, is used to generalize logistic regression when there are more than 2 outcome classes. By convention, we set $\theta_K = 0$, which makes the Bernoulli parameter ϕ_i of each class i equal to:

$$\phi_i = \frac{\exp(\theta_i^T x)}{\sum_{j=1}^K \exp(\theta_j^T x)}$$

1.3.3Generalized Linear Models

□ Exponential family – A class of distributions is said to be in the exponential family if it can be written in terms of a natural parameter, also called the canonical parameter or link function, η , a sufficient statistic T(y) and a log-partition function $a(\eta)$ as follows:

$$p(y;\eta) = b(y) \exp(\eta T(y) - a(\eta))$$

Remark: we will often have T(y) = y. Also, $\exp(-a(\eta))$ can be seen as a normalization parameter that will make sure that the probabilities sum to one.

Here are the most common exponential distributions summed up in the following table:

Distribution	η	T(y)	$a(\eta)$	b(y)
Bernoulli	$\log\left(\frac{\phi}{1-\phi}\right)$	y	$\log(1 + \exp(\eta))$	1
Gaussian	μ	y	$\frac{\eta^2}{2}$	$\frac{1}{\sqrt{2\pi}}\exp\left(-\frac{y^2}{2}\right)$
Poisson	$\log(\lambda)$	y	e^{η}	$\frac{1}{y!}$
Geometric	$\log(1-\phi)$	y	$\log\left(\frac{e^{\eta}}{1-e^{\eta}}\right)$	1

□ Assumptions of GLMs – Generalized Linear Models (GLM) aim at predicting a random variable y as a function fo $x \in \mathbb{R}^{n+1}$ and rely on the following 3 assumptions:

(1)
$$y|x; \theta \sim \text{ExpFamily}(\eta)$$
 (2) $h_{\theta}(x) = E[y|x; \theta]$ (3) $\eta = \theta^{T}x$

(2)
$$h_{\theta}(x) = E[y|x;$$

Remark: ordinary least squares and logistic regression are special cases of generalized linear models.

1.4 Support Vector Machines

The goal of support vector machines is to find the line that maximizes the minimum distance to the line.

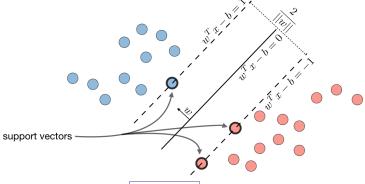
 \Box Optimal margin classifier – The optimal margin classifier h is such that:

$$h(x) = \operatorname{sign}(w^T x - b)$$

where $(w,b) \in \mathbb{R}^n \times \mathbb{R}$ is the solution of the following optimization problem:

$$\min \frac{1}{2}||w||^2$$
 such that

such that
$$y^{(i)}(w^Tx^{(i)}-b) \geqslant 1$$



Remark: the line is defined as $w^T x - b = 0$

☐ **Hinge loss** – The hinge loss is used in the setting of SVMs and is defined as follows:

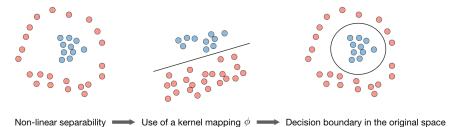
$$L(z,y) = [1 - yz]_{+} = \max(0,1 - yz)$$

4

 \square Kernel – Given a feature mapping ϕ , we define the kernel K to be defined as:

$$K(x,z) = \phi(x)^T \phi(z)$$

In practice, the kernel K defined by $K(x,z)=\exp\left(-\frac{||x-z||^2}{2\sigma^2}\right)$ is called the Gaussian kernel and is commonly used.



Remark: we say that we use the "kernel trick" to compute the cost function using the kernel because we actually don't need to know the explicit mapping ϕ , which is often very complicated. Instead, only the values K(x,z) are needed.

□ Lagrangian – We define the Lagrangian $\mathcal{L}(w,b)$ as follows:

$$\mathcal{L}(w,b) = f(w) + \sum_{i=1}^{l} \beta_i h_i(w)$$

Remark: the coefficients β_i are called the Lagrange multipliers.

1.5 Generative Learning

A generative model first tries to learn how the data is generated by estimating P(x|y), which we can then use to estimate P(y|x) by using Bayes' rule.

1.5.1 Gaussian Discriminant Analysis

 \square Setting – The Gaussian Discriminant Analysis assumes that y and x|y=0 and x|y=1 are such that:

$$\boxed{y \sim \text{Bernoulli}(\phi)}$$

$$\boxed{x|y=0 \sim \mathcal{N}(\mu_0, \Sigma)} \quad \text{and} \quad \boxed{x|y=1 \sim \mathcal{N}(\mu_1, \Sigma)}$$

 \square **Estimation** – The following table sums up the estimates that we find when maximizing the likelihood:

$\widehat{\phi}$	$\widehat{\mu_j}$ $(j=0,1)$	$\widehat{\Sigma}$
$\frac{1}{m} \sum_{i=1}^{m} 1_{\{y^{(i)}=1\}}$	$\frac{\sum_{i=1}^{m} 1_{\{y^{(i)}=j\}} x^{(i)}}{\sum_{i=1}^{m} 1_{\{y^{(i)}=j\}}}$	$\frac{1}{m} \sum_{i=1}^{m} (x^{(i)} - \mu_{y^{(i)}}) (x^{(i)} - \mu_{y^{(i)}})^{T}$

1.5.2 Naive Bayes

 \square **Assumption** – The Naive Bayes model supposes that the features of each data point are all independent:

$$P(x|y) = P(x_1, x_2, ...|y) = P(x_1|y)P(x_2|y)... = \prod_{i=1}^{n} P(x_i|y)$$

 $\hfill \Box$ Solutions – Maximizing the log-likelihood gives the following solutions, with $k \in \{0,1\},$ $l \in [\![1,L]\!]$

$$P(y=k) = \frac{1}{m} \times \#\{j|y^{(j)} = k\} \quad \text{and} \quad P(x_i = l|y = k) = \frac{\#\{j|y^{(j)} = k \text{ and } x_i^{(j)} = l\}}{\#\{j|y^{(j)} = k\}}$$

Remark: Naive Bayes is widely used for text classification and spam detection.

1.6 Tree-based and ensemble methods

These methods can be used for both regression and classification problems.

 \square CART – Classification and Regression Trees (CART), commonly known as decision trees, can be represented as binary trees. They have the advantage to be very interpretable.

□ Random forest – It is a tree-based technique that uses a high number of decision trees built out of randomly selected sets of features. Contrary to the simple decision tree, it is highly uninterpretable but its generally good performance makes it a popular algorithm.

Remark: random forests are a type of ensemble methods.

□ Boosting – The idea of boosting methods is to combine several weak learners to form a stronger one. The main ones are summed up in the table below:

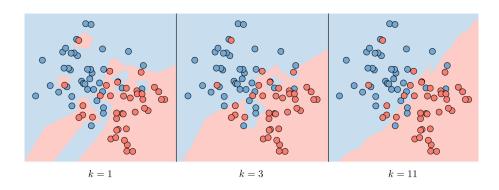
Adaptive boosting	Gradient boosting
- High weights are put on errors to improve at the next boosting step - Known as Adaboost	- Weak learners trained on remaining errors

1.7 Other non-parametric approaches

 \square k-nearest neighbors – The k-nearest neighbors algorithm, commonly known as k-NN, is a non-parametric approach where the response of a data point is determined by the nature of its k neighbors from the training set. It can be used in both classification and regression settings.

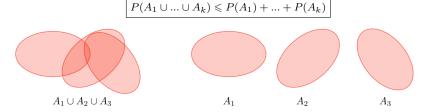
Remark: The higher the parameter k, the higher the bias, and the lower the parameter k, the higher the variance.

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1.8 Learning Theory

 \square Union bound – Let $A_1,...,A_k$ be k events. We have:



□ Hoeffding inequality – Let $Z_1,..,Z_m$ be m iid variables drawn from a Bernoulli distribution of parameter ϕ . Let $\widehat{\phi}$ be their sample mean and $\gamma > 0$ fixed. We have:

$$P(|\phi - \widehat{\phi}| > \gamma) \le 2\exp(-2\gamma^2 m)$$

Remark: this inequality is also known as the Chernoff bound.

 \square Training error – For a given classifier h, we define the training error $\widehat{\epsilon}(h)$, also known as the empirical risk or empirical error, to be as follows:

$$\widehat{\epsilon}(h) = \frac{1}{m} \sum_{i=1}^{m} 1_{\{h(x^{(i)}) \neq y^{(i)}\}}$$

- □ Probably Approximately Correct (PAC) PAC is a framework under which numerous results on learning theory were proved, and has the following set of assumptions:
 - the training and testing sets follow the same distribution
 - the training examples are drawn independently
- \square Shattering Given a set $S = \{x^{(1)},...,x^{(d)}\}$, and a set of classifiers \mathcal{H} , we say that \mathcal{H} shatters S if for any set of labels $\{y^{(1)},...,y^{(d)}\}$, we have:

$$\exists h \in \mathcal{H}, \quad \forall i \in [1,d], \quad h(x^{(i)}) = y^{(i)}$$

□ Upper bound theorem – Let \mathcal{H} be a finite hypothesis class such that $|\mathcal{H}| = k$ and let δ and the sample size m be fixed. Then, with probability of at least $1 - \delta$, we have:

$$\widehat{\epsilon(h)} \leqslant \left(\min_{h \in \mathcal{H}} \epsilon(h)\right) + 2\sqrt{\frac{1}{2m}\log\left(\frac{2k}{\delta}\right)}$$

 \square VC dimension – The Vapnik-Chervonenkis (VC) dimension of a given infinite hypothesis class \mathcal{H} , noted VC(\mathcal{H}) is the size of the largest set that is shattered by \mathcal{H} .

Remark: the VC dimension of $\mathcal{H} = \{ set \ of \ linear \ classifiers \ in \ 2 \ dimensions \}$ is 3.



□ Theorem (Vapnik) – Let \mathcal{H} be given, with VC(\mathcal{H}) = d and m the number of training examples. With probability at least $1 - \delta$, we have:

$$\widehat{\epsilon(h)} \leqslant \left(\min_{h \in \mathcal{H}} \epsilon(h)\right) + O\left(\sqrt{\frac{d}{m}\log\left(\frac{m}{d}\right) + \frac{1}{m}\log\left(\frac{1}{\delta}\right)}\right)$$

2 Unsupervised Learning

2.1 Introduction to Unsupervised Learning

□ Motivation – The goal of unsupervised learning is to find hidden patterns in unlabeled data $\{x^{(1)},...,x^{(m)}\}.$

 \square Jensen's inequality – Let f be a convex function and X a random variable. We have the following inequality:

$$E[f(X)] \geqslant f(E[X])$$

2.2 Clustering

2.2.1 Expectation-Maximization

 \Box Latent variables – Latent variables are hidden/unobserved variables that make estimation problems difficult, and are often denoted z. Here are the most common settings where there are latent variables:

Setting	Latent variable z	x z	Comments
Mixture of k Gaussians	$\operatorname{Multinomial}(\phi)$	$\mathcal{N}(\mu_j, \Sigma_j)$	$\mu_j \in \mathbb{R}^n, \phi \in \mathbb{R}^k$
Factor analysis	$\mathcal{N}(0,I)$	$\mathcal{N}(\mu + \Lambda z, \psi)$	$\mu_j \in \mathbb{R}^n$

 \square Algorithm – The Expectation-Maximization (EM) algorithm gives an efficient method at estimating the parameter θ through maximum likelihood estimation by repeatedly constructing a lower-bound on the likelihood (E-step) and optimizing that lower bound (M-step) as follows:

• E-step: Evaluate the posterior probability $Q_i(z^{(i)})$ that each data point $x^{(i)}$ came from a particular cluster $z^{(i)}$ as follows:

$$Q_i(z^{(i)}) = P(z^{(i)}|x^{(i)};\theta)$$

• M-step: Use the posterior probabilities $Q_i(z^{(i)})$ as cluster specific weights on data points $\overline{x^{(i)}}$ to separately re-estimate each cluster model as follows:

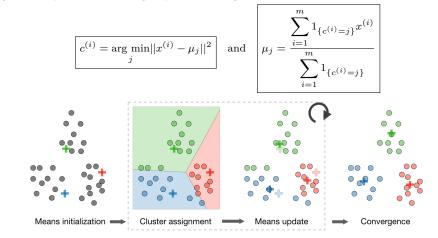
$$\theta_i = \underset{\theta}{\operatorname{argmax}} \sum_i \int_{z^{(i)}} Q_i(z^{(i)}) \log \left(\frac{P(x^{(i)}, z^{(i)}; \theta)}{Q_i(z^{(i)})} \right) dz^{(i)}$$

Gaussians initialization Expectation step Maximization step Convergence

2.2.2 k-means clustering

We note $c^{(i)}$ the cluster of data point i and μ_i the center of cluster j.

 \square Algorithm – After randomly initializing the cluster centroids $\mu_1, \mu_2, ..., \mu_k \in \mathbb{R}^n$, the k-means algorithm repeats the following step until convergence:



□ Distortion function – In order to see if the algorithm converges, we look at the distortion function defined as follows:

$$J(c,\mu) = \sum_{i=1}^{m} ||x^{(i)} - \mu_{c(i)}||^2$$

2.2.3 Hierarchical clustering

 \square Algorithm – It is a clustering algorithm with an agglomerative hierarchical approach that build nested clusters in a successive manner.

□ Types – There are different sorts of hierarchical clustering algorithms that aims at optimizing different objective functions, which is summed up in the table below:

Ward linkage	Average linkage	Complete linkage	
Minimize within cluster distance	Minimize average distance between cluster pairs	Minimize maximum distance of between cluster pairs	

2.2.4 Clustering assessment metrics

In an unsupervised learning setting, it is often hard to assess the performance of a model since we don't have the ground truth labels as was the case in the supervised learning setting.

 \square Silhouette coefficient – By noting a and b the mean distance between a sample and all other points in the same class, and between a sample and all other points in the next nearest cluster, the silhouette coefficient s for a single sample is defined as follows:

$$s = \frac{b - a}{\max(a, b)}$$

 \square Calinski-Harabaz index – By noting k the number of clusters, B_k and W_k the between and within-clustering dispersion matrices respectively defined as

$$B_k = \sum_{i=1}^k n_{c(i)} (\mu_{c(i)} - \mu) (\mu_{c(i)} - \mu)^T, \qquad W_k = \sum_{i=1}^m (x^{(i)} - \mu_{c(i)}) (x^{(i)} - \mu_{c(i)})^T$$

the Calinski-Harabaz index s(k) indicates how well a clustering model defines its clusters, such that the higher the score, the more dense and well separated the clusters are. It is defined as follows:

$$s(k) = \frac{\operatorname{Tr}(B_k)}{\operatorname{Tr}(W_k)} \times \frac{N-k}{k-1}$$

2.3 Dimension reduction

2.3.1 Principal component analysis

It is a dimension reduction technique that finds the variance maximizing directions onto which to project the data.

□ Eigenvalue, eigenvector – Given a matrix $A \in \mathbb{R}^{n \times n}$, λ is said to be an eigenvalue of A if there exists a vector $z \in \mathbb{R}^n \setminus \{0\}$, called eigenvector, such that we have:

$$Az = \lambda z$$

□ Spectral theorem – Let $A \in \mathbb{R}^{n \times n}$. If A is symmetric, then A is diagonalizable by a real orthogonal matrix $U \in \mathbb{R}^{n \times n}$. By noting $\Lambda = \operatorname{diag}(\lambda_1, ..., \lambda_n)$, we have:

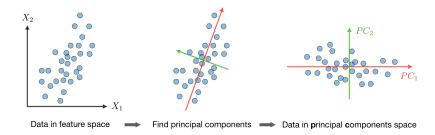
$$\exists \Lambda \text{ diagonal}, \quad A = U \Lambda U^T$$

Remark: the eigenvector associated with the largest eigenvalue is called principal eigenvector of matrix A.

- \square Algorithm The Principal Component Analysis (PCA) procedure is a dimension reduction technique that projects the data on k dimensions by maximizing the variance of the data as follows:
 - Step 1: Normalize the data to have a mean of 0 and standard deviation of 1.

$$\boxed{x_{j}^{(i)} \leftarrow \frac{x_{j}^{(i)} - \mu_{j}}{\sigma_{j}}} \quad \text{where} \quad \boxed{\mu_{j} = \frac{1}{m} \sum_{i=1}^{m} x_{j}^{(i)}} \quad \text{and} \quad \boxed{\sigma_{j}^{2} = \frac{1}{m} \sum_{i=1}^{m} (x_{j}^{(i)} - \mu_{j})^{2}}$$

- Step 2: Compute $\Sigma = \frac{1}{m} \sum_{i=1}^{m} x^{(i)} x^{(i)}^T \in \mathbb{R}^{n \times n}$, which is symmetric with real eigenvalues.
- Step 3: Compute $u_1, ..., u_k \in \mathbb{R}^n$ the k orthogonal principal eigenvectors of Σ , i.e. the orthogonal eigenvectors of the k largest eigenvalues.
- Step 4: Project the data on $\operatorname{span}_{\mathbb{R}}(u_1,...,u_k)$. This procedure maximizes the variance among all k-dimensional spaces.



2.3.2 Independent component analysis

It is a technique meant to find the underlying generating sources.

□ Assumptions – We assume that our data x has been generated by the n-dimensional source vector $s = (s_1, ..., s_n)$, where s_i are independent random variables, via a mixing and non-singular matrix A as follows:

$$x = As$$

The goal is to find the unmixing matrix $W = A^{-1}$ by an update rule.

 \Box Bell and Sejnowski ICA algorithm – This algorithm finds the unmixing matrix W by following the steps below:

• Write the probability of $x = As = W^{-1}s$ as:

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$$p(x) = \prod_{i=1}^{n} p_s(w_i^T x) \cdot |W|$$

• Write the log likelihood given our training data $\{x^{(i)}, i \in [\![1,m]\!]\}$ and by noting g the sigmoid function as:

$$l(W) = \sum_{i=1}^{m} \left(\sum_{j=1}^{n} \log \left(g'(w_{j}^{T} x^{(i)}) \right) + \log |W| \right)$$

Therefore, the stochastic gradient ascent learning rule is such that for each training example $x^{(i)}$, we update W as follows:

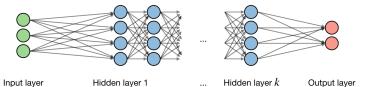
$$W \longleftarrow W + \alpha \begin{pmatrix} \begin{pmatrix} 1 - 2g(w_1^T x^{(i)}) \\ 1 - 2g(w_2^T x^{(i)}) \\ \vdots \\ 1 - 2g(w_n^T x^{(i)}) \end{pmatrix} x^{(i)^T} + (W^T)^{-1} \end{pmatrix}$$

3 Deep Learning

3.1 Neural Networks

Neural networks are a class of models that are built with layers. Commonly used types of neural networks include convolutional and recurrent neural networks.

 $\hfill \square$ Architecture – The vocabulary around neural networks architectures is described in the figure below:



By noting i the i^{th} layer of the network and j the j^{th} hidden unit of the layer, we have:

$$z_{j}^{[i]} = w_{j}^{[i]} x + b_{j}^{[i]}$$

where we note w, b, z the weight, bias and output respectively.

□ Activation function – Activation functions are used at the end of a hidden unit to introduce non-linear complexities to the model. Here are the most common ones:

Sigmoid	Tanh	\mathbf{ReLU}	Leaky ReLU
$g(z) = \frac{1}{1 + e^{-z}}$	$g(z) = \frac{e^z - e^{-z}}{e^z + e^{-z}}$	$g(z) = \max(0, z)$	$g(z) = \max(\epsilon z, z)$ with $\epsilon \ll 1$
$\begin{array}{c c} 1 \\ \hline \frac{1}{2} \\ \hline -4 & 0 \end{array}$	1 -4 0 4	0 1	0 1

 \square Cross-entropy loss – In the context of neural networks, the cross-entropy loss L(z,y) is commonly used and is defined as follows:

$$L(z,y) = -\left[y\log(z) + (1-y)\log(1-z)\right]$$

 \square Learning rate – The learning rate, often noted η , indicates at which pace the weights get updated. This can be fixed or adaptively changed. The current most popular method is called Adam, which is a method that adapts the learning rate.

 \square Backpropagation – Backpropagation is a method to update the weights in the neural network by taking into account the actual output and the desired output. The derivative with respect to weight w is computed using chain rule and is of the following form:

$$\frac{\partial L(z,y)}{\partial w} = \frac{\partial L(z,y)}{\partial a} \times \frac{\partial a}{\partial z} \times \frac{\partial z}{\partial w}$$

As a result, the weight is updated as follows:

$$w \longleftarrow w - \eta \frac{\partial L(z, y)}{\partial w}$$

□ Updating weights – In a neural network, weights are updated as follows:

- Step 1: Take a batch of training data.
- Step 2: Perform forward propagation to obtain the corresponding loss.
- Step 3: Backpropagate the loss to get the gradients.
- Step 4: Use the gradients to update the weights of the network.

 \square Dropout – Dropout is a technique meant at preventing overfitting the training data by dropping out units in a neural network. In practice, neurons are either dropped with probability p or kept with probability 1-p.

3.2 Convolutional Neural Networks

 \square Convolutional layer requirement – By noting W the input volume size, F the size of the convolutional layer neurons, P the amount of zero padding, then the number of neurons N that fit in a given volume is such that:

$$N = \frac{W - F + 2P}{S} + 1$$

 \square Batch normalization – It is a step of hyperparameter γ, β that normalizes the batch $\{x_i\}$. By noting μ_B, σ_B^2 the mean and variance of that we want to correct to the batch, it is done as follows:

$$x_i \longleftarrow \gamma \frac{x_i - \mu_B}{\sqrt{\sigma_B^2 + \epsilon}} + \beta$$

It is usually done after a fully connected/convolutional layer and before a non-linearity layer and aims at allowing higher learning rates and reducing the strong dependence on initialization.

3.3 Recurrent Neural Networks

☐ Types of gates – Here are the different types of gates that we encounter in a typical recurrent neural network:

Input gate	Forget gate	Output gate	Gate
Write to cell or not?	Erase a cell or not?	Reveal a cell or not?	How much writing?

 $\hfill \square$ LSTM – A long short-term memory (LSTM) network is a type of RNN model that avoids the vanishing gradient problem by adding 'forget' gates.