

Section 1: Probability, Statistics, & Linear Algebra review

STATS 202: Data Mining and Analysis

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- ▶ Linear algebra
 - ▶ Basic concepts
 - ▶ Matrix multiplication
 - ▶ Operations and Properties
 - ▶ Matrix Calculus
- ▶ Probability
 - ▶ Sample space
 - ▶ Probability function
 - ▶ Probability space
 - ▶ Random variables
- ▶ Statistics
 - ▶ Expected value
 - ▶ Moments & Moment generating functions
 - ▶ Distributions



Linear algebra



Consider the following equations:

$$4x_1 - 5x_2 = -13 \quad (1)$$

$$-2x_1 + 3x_2 = 9 \quad (2)$$

Let's solve for x_1 and x_2 .



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$$4x_1 - 5x_2 = -13 \quad (1)$$

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Let's solve for x_1 and x_2 .

We can write this system of equations more compactly in matrix notation, e.g.

$$\mathbf{Ax} = \mathbf{b} \quad (3)$$

where $\mathbf{A} = \begin{bmatrix} 4 & -5 \\ -2 & 3 \end{bmatrix}$ and $\mathbf{b} = \begin{bmatrix} -13 \\ 9 \end{bmatrix}$



Some basic notation:

- ▶ We denote a matrix with m rows and n columns as $\mathbf{A} \in \mathbb{R}^{m \times n}$, where each entry in the matrix is a real number.
- ▶ We denote a vector with n entries as $\mathbf{x} \in \mathbb{R}^n$.
 - ▶ By convention, we typically think of a vector as a 1 column matrix.
- ▶ We denote the i^{th} element of a vector \mathbf{x} as x_i , e.g.

$$\mathbf{x} = \begin{bmatrix} x_1 \\ x_2 \\ \vdots \\ x_n \end{bmatrix} \quad (4)$$



Some basic notation:

- ▶ We denote each entry in a matrix \mathbf{A} by a_{ij} , corresponding to the i^{th} row and j^{th} column, e.g.

$$\mathbf{A} = \begin{bmatrix} a_{11} & a_{12} & \cdots & a_{1n} \\ a_{21} & a_{22} & \cdots & a_{2n} \\ \vdots & \vdots & \ddots & \vdots \\ a_{m1} & a_{m2} & \cdots & a_{mn} \end{bmatrix} \quad (5)$$

- ▶ We denote the *transpose* of a matrix as \mathbf{A}^{\top} , e.g.

$$\mathbf{A}^{\top} = \begin{bmatrix} a_{11} & a_{21} & \cdots & a_{m1} \\ a_{12} & a_{22} & \cdots & a_{m2} \\ \vdots & \vdots & \ddots & \vdots \\ a_{1n} & a_{2n} & \cdots & a_{mn} \end{bmatrix} \quad (6)$$



Some basic notation:

- ▶ We denote the j^{th} column of \mathbf{A} by \mathbf{a}_j or $\mathbf{A}_{.j}$, e.g.

$$\mathbf{A} = \begin{bmatrix} | & | & \cdots & | \\ \mathbf{a}_1 & \mathbf{a}_2 & & \mathbf{a}_n \\ | & | & & | \end{bmatrix} \quad (7)$$

- ▶ We denote the i^{th} row of \mathbf{A} by \mathbf{a}_i^{\top} or $\mathbf{A}_{i..}$.

$$\mathbf{A} = \begin{bmatrix} \text{---} & \mathbf{a}_1^{\top} & \text{---} \\ \text{---} & \mathbf{a}_2^{\top} & \text{---} \\ & \vdots & \\ \text{---} & \mathbf{a}_m^{\top} & \text{---} \end{bmatrix} \quad (8)$$

n.b. This isn't universal, though should be clear from its presentation and use.



Given two matrices $\mathbf{A} \in \mathbb{R}^{m \times n}$, $\mathbf{B} \in \mathbb{R}^{n \times p}$, we can multiply them by

$$\mathbf{C} = \mathbf{AB} \in \mathbb{R}^{m \times p} : \mathbf{C}_{ij} = \sum_{k=1}^n \mathbf{A}_{ik} \mathbf{B}_{kj} \quad (9)$$

n.b. The dimensions have to be compatible for matrix multiplication to be valid (e.g. the number of columns in \mathbf{A} must be equal to the number of rows in \mathbf{B}).



Given $\mathbf{x}, \mathbf{y} \in \mathbb{R}^n$, the quantity $\mathbf{x}^\top \mathbf{y} \in \mathbb{R}$ (aka *dot product* or *inner product*) is a scalar given by

$$\mathbf{x}^\top \mathbf{y} = [x_1 \quad x_2 \quad \cdots \quad x_n] \begin{bmatrix} y_1 \\ y_2 \\ \vdots \\ y_n \end{bmatrix} = \sum_{i=1}^n x_i y_i \quad (10)$$

Note: For vectors, we always have that $\mathbf{x}^\top \mathbf{y} = \mathbf{y}^\top \mathbf{x}$. This is not generally true for matrices.



Given $\mathbf{x} \in \mathbb{R}^m$, $\mathbf{y} \in \mathbb{R}^n$, the quantity $\mathbf{x}^\top \mathbf{y} \in \mathbb{R}^{m \times n}$ (aka *outer product*) is a matrix given by

$$\mathbf{xy}^\top = \begin{bmatrix} x_1 \\ x_2 \\ \vdots \\ x_n \end{bmatrix} \begin{bmatrix} y_1 & y_2 & \cdots & y_n \end{bmatrix} = \begin{bmatrix} x_1 y_1 & x_1 y_2 & \cdots & x_1 y_n \\ x_2 y_1 & x_2 y_2 & \cdots & x_2 y_n \\ \vdots & \vdots & \ddots & \vdots \\ x_m y_1 & x_m y_2 & \cdots & x_m y_n \end{bmatrix} \quad (11)$$



Example: Let $\mathbf{A} \in \mathbb{R}^{m \times n}$ be a matrix such that all columns are equal to some vector $\mathbf{x} \in \mathbb{R}^m$. Using outer products, we can represent \mathbf{A} compactly as

$$\mathbf{A} = \begin{bmatrix} | & | & \cdots & | \\ \mathbf{x} & \mathbf{x} & \cdots & \mathbf{x} \\ | & | & \cdots & | \end{bmatrix} = \begin{bmatrix} x_1 & x_1 & \cdots & x_1 \\ x_2 & x_2 & \cdots & x_2 \\ \vdots & \vdots & \ddots & \vdots \\ x_m & x_m & \cdots & x_m \end{bmatrix} \quad (12)$$

$$= \begin{bmatrix} x_1 \\ x_2 \\ \vdots \\ x_m \end{bmatrix} \begin{bmatrix} 1 & 1 & \cdots & 1 \end{bmatrix} \quad (13)$$

$$= \mathbf{x} \mathbf{1}^\top \quad (14)$$



Given $\mathbf{A} \in \mathbb{R}^{m \times n}$, $\mathbf{x} \in \mathbb{R}^n$, their product is a vector $\mathbf{y} = \mathbf{Ax} \in \mathbb{R}^m$.



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There are two ways of interpreting this:

$$\mathbf{y} = \mathbf{Ax} = \begin{bmatrix} \text{---} & \mathbf{a}_1^\top & \text{---} \\ \text{---} & \mathbf{a}_2^\top & \text{---} \\ & \vdots & \\ \text{---} & \mathbf{a}_m^\top & \text{---} \end{bmatrix} \mathbf{x} = \begin{bmatrix} \mathbf{a}_1^\top \mathbf{x} \\ \mathbf{a}_2^\top \mathbf{x} \\ \vdots \\ \mathbf{a}_m^\top \mathbf{x} \end{bmatrix} \quad (15)$$

$$= \begin{bmatrix} | & | & \cdots & | \\ \mathbf{a}_1 & \mathbf{a}_2 & \cdots & \mathbf{a}_n \\ | & | & & | \end{bmatrix} \begin{bmatrix} x_1 \\ x_2 \\ \cdots \\ x_n \end{bmatrix} \quad (16)$$

$$= \mathbf{a}_1 x_1 + \mathbf{a}_2 x_2 + \cdots + \mathbf{a}_n x_n \quad (17)$$



Example:

$$\text{Define } \mathbf{A} = \begin{bmatrix} 1 & 2 & 3 \\ 2 & 5 & 6 \\ 7 & 8 & 9 \\ 10 & 11 & 12 \end{bmatrix}, \mathbf{x} = \begin{bmatrix} -3 \\ -2 \\ -1 \end{bmatrix}.$$

Calculate $\mathbf{y} = \mathbf{Ax}$.

Matrix-matrix products



Given $\mathbf{A} \in \mathbb{R}^{m \times n}$, $\mathbf{B} \in \mathbb{R}^{n \times p}$, their product is a matrix $\mathbf{C} = \mathbf{AB} \in \mathbb{R}^{m \times p}$.



Given $\mathbf{A} \in \mathbb{R}^{m \times n}$, $\mathbf{B} \in \mathbb{R}^{n \times p}$, their product is a matrix $\mathbf{C} = \mathbf{AB} \in \mathbb{R}^{m \times p}$.

Similar to before, we can think of this in two ways:

Interpretation # 1

$$\mathbf{C} = \mathbf{AB} = \begin{bmatrix} \text{---} & \mathbf{a}_1^\top & \text{---} \\ \text{---} & \mathbf{a}_2^\top & \text{---} \\ & \vdots & \\ \text{---} & \mathbf{a}_m^\top & \text{---} \end{bmatrix} \begin{bmatrix} | & | & \cdots & | \\ \mathbf{b}_1 & \mathbf{b}_2 & & \mathbf{b}_p \\ | & | & & | \end{bmatrix} \quad (18)$$

$$= \begin{bmatrix} \mathbf{a}_1^\top \mathbf{b}_1 & \mathbf{a}_1^\top \mathbf{b}_2 & \cdots & \mathbf{a}_1^\top \mathbf{b}_p \\ \mathbf{a}_2^\top \mathbf{b}_1 & \mathbf{a}_2^\top \mathbf{b}_2 & \cdots & \mathbf{a}_2^\top \mathbf{b}_p \\ \vdots & \vdots & \ddots & \vdots \\ \mathbf{a}_m^\top \mathbf{b}_1 & \mathbf{a}_m^\top \mathbf{b}_2 & \cdots & \mathbf{a}_m^\top \mathbf{b}_p \end{bmatrix} \quad (19)$$



Interpretation # 2

$$\mathbf{C} = \mathbf{AB} = \mathbf{A} \begin{bmatrix} | & | & \cdots & | \\ \mathbf{b}_1 & \mathbf{b}_2 & \cdots & \mathbf{b}_p \\ | & | & \cdots & | \end{bmatrix} \quad (20)$$

$$= \begin{bmatrix} | & | & \cdots & | \\ \mathbf{Ab}_1 & \mathbf{Ab}_2 & \cdots & \mathbf{Ab}_p \\ | & | & \cdots & | \end{bmatrix} \quad (21)$$

$$= \begin{bmatrix} \text{---} & \mathbf{a}_1^T & \text{---} \\ \text{---} & \mathbf{a}_2^T & \text{---} \\ & \vdots & \\ \text{---} & \mathbf{a}_m^T & \text{---} \end{bmatrix} \mathbf{B} = \begin{bmatrix} \text{---} & \mathbf{a}_1^T \mathbf{B} & \text{---} \\ \text{---} & \mathbf{a}_2^T \mathbf{B} & \text{---} \\ & \vdots & \\ \text{---} & \mathbf{a}_m^T \mathbf{B} & \text{---} \end{bmatrix} \quad (22)$$



- ▶ Associative: $(\mathbf{AB})\mathbf{C} = \mathbf{A}(\mathbf{BC})$
- ▶ Distributive: $\mathbf{A}(\mathbf{B} + \mathbf{C}) = \mathbf{AB} + \mathbf{AC}$
- ▶ Not commutative: $\mathbf{AB} \neq \mathbf{BA}$



Demonstrating *associativity*:

We just need to show that $((\mathbf{AB})\mathbf{C})_{ij} = (\mathbf{A}(\mathbf{BC}))_{ij}$:

$$((\mathbf{AB})\mathbf{C})_{ij} = \sum_{k=1}^p (\mathbf{AB})_{ik} \mathbf{C}_{kj} = \sum_{k=1}^p \left(\sum_{l=1}^n \mathbf{A}_{il} \mathbf{B}_{lk} \right) \mathbf{C}_{kj} \quad (23)$$

$$= \sum_{k=1}^p \left(\sum_{l=1}^n \mathbf{A}_{il} \mathbf{B}_{lk} \mathbf{C}_{kj} \right) = \sum_{l=1}^n \left(\sum_{k=1}^p \mathbf{A}_{il} \mathbf{B}_{lk} \mathbf{C}_{kj} \right) \quad (24)$$

$$= \sum_{l=1}^n \mathbf{A}_{il} \left(\sum_{k=1}^p \mathbf{B}_{lk} \mathbf{C}_{kj} \right) = \sum_{l=1}^n \mathbf{A}_{il} (\mathbf{BC})_{lj} \quad (25)$$

$$= (\mathbf{A}(\mathbf{BC}))_{ij} \quad (26)$$



The identity matrix:

The *identity matrix*, denoted $\mathbf{I} \in \mathbb{R}^{n \times n}$ is a square matrix with 1's in the diagonal and 0's everywhere else, i.e.

$$\mathbf{I}_{ij} = \begin{cases} 1 & i = j \\ 0 & i \neq j \end{cases} \quad (27)$$



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$$\mathbf{I}_{ij} = \begin{cases} 1 & i = j \\ 0 & i \neq j \end{cases} \quad (27)$$

It has the property

$$\mathbf{AI} = \mathbf{A} = \mathbf{IA} \quad \forall \mathbf{A} \in \mathbb{R}^{m \times n} \quad (28)$$

n.b. The dimensionality of \mathbf{I} is typically inferred (e.g. $n \times n$ vs $m \times m$)



The diagonal matrix: The *diagonal matrix*, denoted $\mathbf{D} = \text{diag}(d_1, d_2, \dots, d_n)$ is a matrix where all non-diagonal elements are 0, i.e.

$$\mathbf{D}_{ij} = \begin{cases} d_i & i = j \\ 0 & i \neq j \end{cases} \quad (29)$$

Clearly, $\mathbf{I} = \text{diag}(1, 1, \dots, 1)$.



The *transpose* of a matrix results from “flipping” the rows and columns, i.e.

$$(\mathbf{A}^\top)_{ij} = \mathbf{A}_{ji} \quad (30)$$

Consequently, for $\mathbf{A} \in \mathbb{R}^{m \times n}$ we have that $\mathbf{A}^\top \in \mathbb{R}^{n \times m}$.

Some properties:

- ▶ $(\mathbf{A}^\top)^\top = \mathbf{A}$
- ▶ $(\mathbf{AB})^\top = \mathbf{B}^\top \mathbf{A}^\top$
- ▶ $(\mathbf{A} + \mathbf{B})^\top = \mathbf{A}^\top + \mathbf{B}^\top$



A square matrix $\mathbf{A} \in \mathbb{R}^{n \times n}$ is *symmetric* if $\mathbf{A} = \mathbf{A}^\top$.

It is *anti-symmetric* if $\mathbf{A} = -\mathbf{A}^\top$.



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It is easy to show that $\mathbf{A} + \mathbf{A}^\top$ is symmetric and $\mathbf{A} - \mathbf{A}^\top$ is anti-symmetric. Consequently, we have that

$$\mathbf{A} = \frac{1}{2}(\mathbf{A} + \mathbf{A}^\top) + \frac{1}{2}(\mathbf{A} - \mathbf{A}^\top) \quad (31)$$



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Symmetric matrices tend to be denoted as $\mathbf{A} \in \mathbb{S}^n$.



The *trace* of a square matrix $\mathbf{A} \in \mathbb{R}^{n \times n}$, denoted $tr(\mathbf{A})$ or $tr\mathbf{A}$ is the sum of the diagonal elements, i.e.

$$tr\mathbf{A} = \sum_{i=1}^n \mathbf{A}_{ii} \quad (32)$$

The trace has the following properties:

- ▶ For $\mathbf{A} \in \mathbb{R}^{n \times n}$, $tr\mathbf{A} = tr\mathbf{A}^T$
- ▶ For $\mathbf{A}, \mathbf{B} \in \mathbb{R}^{n \times n}$, $tr(\mathbf{A} + \mathbf{B}) = tr\mathbf{A} + tr\mathbf{B}$
- ▶ For $\mathbf{A} \in \mathbb{R}^{n \times n}$, $c \in \mathbb{R}$, $tr(c\mathbf{A}) = c tr\mathbf{A}$
- ▶ For $\mathbf{A}, \mathbf{B} \in \mathbb{R}^{n \times n} \ni \mathbf{AB} \in \mathbb{R}^{n \times n}$, $tr\mathbf{AB} = tr\mathbf{BA}$
- ▶ For $\mathbf{A}, \mathbf{B}, \mathbf{C} \in \mathbb{R}^{n \times n} \ni \mathbf{ABC} \in \mathbb{R}^{n \times n}$,
 $tr\mathbf{ABC} = tr\mathbf{BCA} = tr\mathbf{CAB}$, and so on for more matrices



Example: Proving that $tr\mathbf{AB} = tr\mathbf{BA}$

$$tr\mathbf{AB} = \sum_{i=1}^m (\mathbf{AB})_{ii} = \sum_{i=1}^m \left(\sum_{j=1}^n \mathbf{A}_{ij} \mathbf{B}_{ji} \right) \quad (33)$$

$$= \sum_{i=1}^m \sum_{j=1}^n \mathbf{A}_{ij} \mathbf{B}_{ji} = \sum_{i=1}^m \sum_{j=1}^n \mathbf{B}_{ji} \mathbf{A}_{ij} \quad (34)$$

$$= \sum_{i=1}^m \left(\sum_{j=1}^n \mathbf{B}_{ji} \mathbf{A}_{ij} \right) = \sum_{j=1}^n (\mathbf{BA})_{jj} \quad (35)$$

$$= tr\mathbf{BA} \quad (36)$$



A *norm* of a vector \mathbf{x} , denoted $\|\mathbf{x}\|$ is a measure of the “length” of the vector. For example, the ℓ_2 -norm (aka Euclidean norm) is

$$\|\mathbf{x}\|_2 = \sqrt{\sum_{i=1}^n x_i^2} \quad (37)$$

n.b. $\|\mathbf{x}\|_2^2 = \mathbf{x}^\top \mathbf{x}$, i.e. the squared norm of a vector is the dot product with itself.



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Other norms:

- ▶ ℓ_1 -norm, i.e. $\|\mathbf{x}\|_1 = \sum_{i=1}^n |x_i|$.
- ▶ ℓ_∞ -norm, i.e. $\|\mathbf{x}\|_\infty = \max_i |x_i|$.
- ▶ ℓ_p -norm, i.e. $\|\mathbf{x}\|_p = (\sum_{i=1}^n |x_i|^p)^{1/p}$.



Formally, a norm is any function $f : \mathbb{R}^n \rightarrow \mathbb{R}$ satisfying four properties:

1. $\forall \mathbf{x} \in \mathbb{R}^n, f(\mathbf{x}) \geq 0$ (non-negativity).
2. $f(\mathbf{x}) = 0$ iff $\mathbf{x} = 0$ (definiteness).
3. $\forall \mathbf{x} \in \mathbb{R}^n, c \in \mathbb{R}, f(c\mathbf{x}) = |c|f(\mathbf{x})$ (homogeneity).
4. $\forall \mathbf{x}, \mathbf{y} \in \mathbb{R}^n, f(\mathbf{x} + \mathbf{y}) \leq f(\mathbf{x}) + f(\mathbf{y})$ (triangle inequality).



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4. $\forall \mathbf{x}, \mathbf{y} \in \mathbb{R}^n, f(\mathbf{x} + \mathbf{y}) \leq f(\mathbf{x}) + f(\mathbf{y})$ (triangle inequality).

Norms can also be defined for matrices, e.g. The Frobenius norm,

$$\|\mathbf{A}\|^F = \sqrt{\sum_{i=1}^m \sum_{j=1}^n \mathbf{A}_{ij}^2} = \sqrt{\text{tr}(\mathbf{A}^\top \mathbf{A})} \quad (38)$$



A set of vectors $\{\mathbf{x}_1, \mathbf{x}_2, \dots, \mathbf{x}_n\} \in \mathbb{R}^m$ is *(linearly) dependent* if one of the vectors \mathbf{x}_i can be represented as a linear combination of the remaining vectors, i.e.

$$\mathbf{x}_n = \sum_{i=1}^{n-1} \alpha_i \mathbf{x}_i \quad (39)$$

for some scalar values $\alpha_1, \alpha_2, \dots, \alpha_{n-1} \in \mathbb{R}$



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Example: Let

$$\mathbf{x}_1 = \begin{bmatrix} 1 \\ 2 \\ 3 \end{bmatrix} \quad \mathbf{x}_2 = \begin{bmatrix} 4 \\ 1 \\ 5 \end{bmatrix} \quad \mathbf{x}_3 = \begin{bmatrix} 2 \\ -3 \\ -1 \end{bmatrix} \quad (40)$$

Is $\{\mathbf{x}_1, \mathbf{x}_2, \mathbf{x}_3\}$ linearly independent?



The *column rank* of $\mathbf{A} \in \mathbb{R}^{m \times n}$ is the largest subset of columns of \mathbf{A} that are linearly independent.

- ▶ The column rank is always $\leq n$.

The *row rank* of $\mathbf{A} \in \mathbb{R}^{m \times n}$ is the largest subset of rows of \mathbf{A} that are linearly independent.

- ▶ The row rank is always $\leq m$.



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- ▶ The row rank is always $\leq m$.

n.b. Column rank is always equal to row rank. Thus, we refer to both as the *rank* of the matrix.

- ▶ For $\mathbf{A} \in \mathbb{R}^{m \times n}$, if $\text{rank}(\mathbf{A}) = \min(m, n)$, then \mathbf{A} is said to be of *full rank*.
- ▶ For $\mathbf{A} \in \mathbb{R}^{m \times n}$, $\text{rank}(\mathbf{A}) = \text{rank}(\mathbf{A}^T)$.
- ▶ For $\mathbf{A} \in \mathbb{R}^{m \times n}$, $\mathbf{B} \in \mathbb{R}^{n \times p}$,
 $\text{rank}(\mathbf{AB}) \leq \min(\text{rank}(\mathbf{A}), \text{rank}(\mathbf{B}))$.
- ▶ For $\mathbf{A}, \mathbf{B} \in \mathbb{R}^{m \times n}$, $\text{rank}(\mathbf{A} + \mathbf{B}) \leq \text{rank}(\mathbf{A}) + \text{rank}(\mathbf{B})$



The *inverse* of a square matrix $\mathbf{A} \in \mathbb{R}^{n \times n}$ is denoted \mathbf{A}^{-1} , and is unique such that

$$\mathbf{A}^{-1}\mathbf{A} = \mathbf{I} = \mathbf{A}\mathbf{A}^{-1} \quad (41)$$



The *inverse* of a square matrix $\mathbf{A} \in \mathbb{R}^{n \times n}$ is denoted \mathbf{A}^{-1} , and is unique such that

$$\mathbf{A}^{-1}\mathbf{A} = \mathbf{I} = \mathbf{A}\mathbf{A}^{-1} \quad (41)$$

n.b. Not all matrices have inverses (e.g. $m \times n$ matrices).

Def:

A is *invertible* or *non-singular* if \mathbf{A}^{-1} exists.

Otherwise, it is *non-invertible* or *singular*.

1. $(\mathbf{A}^{-1})^{-1} = \mathbf{A}$
2. $(\mathbf{AB})^{-1} = \mathbf{B}^{-1}\mathbf{A}^{-1}$
3. $(\mathbf{A}^{-1})^\top = (\mathbf{A}^\top)^{-1}$

► This matrix is sometimes denoted $\mathbf{A}^{-\top}$



Def:

- ▶ A vector $\mathbf{x} \in \mathbb{R}^n$ is *normalized* if $\|\mathbf{x}\|_2 = 1$
- ▶ Two vectors $\mathbf{x}, \mathbf{y} \in \mathbb{R}^n$ are *orthogonal* if $\mathbf{x}^\top \mathbf{y} = 0$
- ▶ A square matrix $\mathbf{U} \in \mathbb{R}^{n \times n}$ is *orthogonal* or *orthonormal* if all its columns are:
 1. Orthogonal to each other
 2. Normalized

We therefore have that

$$\mathbf{U}^\top \mathbf{U} = \mathbf{I} = \mathbf{U} \mathbf{U}^\top \quad (42)$$



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We therefore have that

$$\mathbf{U}^\top \mathbf{U} = \mathbf{I} = \mathbf{U} \mathbf{U}^\top \quad (42)$$

Another nice property:

$$\|\mathbf{U}\mathbf{x}\|_2 = \|\mathbf{x}\|_2 \quad \forall \mathbf{x} \in \mathbb{R}^n, \mathbf{U} \in \mathbb{R}^{n \times n} \text{ orthogonal} \quad (43)$$



Def:

The *span* of a set of vectors $\{\mathbf{x}_1, \mathbf{x}_2, \dots, \mathbf{x}_n\}$ is

$$\text{span}(\{\mathbf{x}_1, \mathbf{x}_2, \dots, \mathbf{x}_n\}) = \left\{ \mathbf{v} : \mathbf{v} = \sum_{i=1}^n \alpha_i \mathbf{x}_i, \alpha_i \in \mathbb{R} \right\} \quad (44)$$

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n.b. If $\{\mathbf{x}_1, \mathbf{x}_2, \dots, \mathbf{x}_n\}$ is linearly independent, then $\text{span}(\{\mathbf{x}_1, \mathbf{x}_2, \dots, \mathbf{x}_n\}) = \mathbb{R}^n$.

Example:

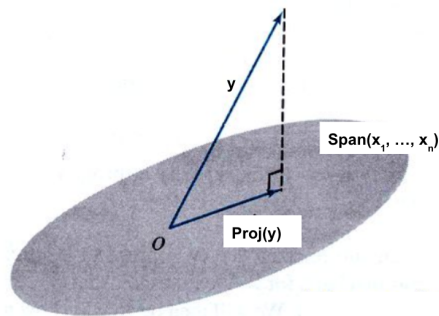
$$\mathbf{x}_1 = \begin{bmatrix} 1 \\ 0 \end{bmatrix} \quad \mathbf{x}_2 = \begin{bmatrix} 0 \\ 1 \end{bmatrix} \quad (45)$$



Def:

The *projection* of a vector $\mathbf{y} \in \mathbb{R}^m$ onto $\text{span}(\{\mathbf{x}_1, \mathbf{x}_2, \dots, \mathbf{x}_n\}) = \mathbb{R}^n$ is

$$\text{Proj}(\mathbf{y}; \{\mathbf{x}_1, \mathbf{x}_2, \dots, \mathbf{x}_n\}) = \underset{\mathbf{v} \in \text{span}(\{\mathbf{x}_1, \mathbf{x}_2, \dots, \mathbf{x}_n\})}{\text{arg min}} \|\mathbf{y} - \mathbf{v}\|_2 \quad (46)$$





Def:

The *range* of a matrix $\mathbf{A} \in \mathbb{R}^{m \times n}$, denoted $\mathcal{R}(\mathbf{A})$ is the span of the columns of \mathbf{A} , i.e.

$$\mathcal{R}(\mathbf{A}) = \{\mathbf{v} \in \mathbb{R}^m : \mathbf{v} = \mathbf{A}\mathbf{x}, \mathbf{x} \in \mathbb{R}^n\} \quad (47)$$

Assuming that \mathbf{A} is full rank and $n < m$, the projection of $\mathbf{y} \in \mathbb{R}^m$ onto $\mathcal{R}(\mathbf{A})$ is

$$\text{Proj}(\mathbf{y}; \mathbf{A}) = \arg \min_{\mathbf{v} \in \mathcal{R}(\mathbf{A})} \|\mathbf{v} - \mathbf{y}\|_2 \quad (48)$$

$$= \mathbf{A}(\mathbf{A}^\top \mathbf{A})^{-1} \mathbf{A}^\top \mathbf{y} \quad (49)$$



Def:

The *nullspace* of a matrix $\mathbf{A} \in \mathbb{R}^{m \times n}$, denoted $\mathcal{N}(\mathbf{A})$ is the set of all vectors that equal 0 when multiplied by \mathbf{A} , i.e.

$$\mathcal{N}(\mathbf{A}) = \{\mathbf{x} \in \mathbb{R}^n : \mathbf{A}\mathbf{x} = \mathbf{0}\} \quad (50)$$

Some properties:

- ▶ $\{w : w = u + v, u \in \mathcal{R}(\mathbf{A}^\top), v \in \mathcal{N}(\mathbf{A})\} = \mathbb{R}^n$
- ▶ $\mathcal{R}(\mathbf{A}^\top) \cap \mathcal{N}(\mathbf{A}) = \{\mathbf{0}\}$

This is referred to as *orthogonal complements*, denoted as $\mathcal{R}(\mathbf{A}^\top) = \mathcal{N}(\mathbf{A})^\perp$



Def:

The *determinant* of a square matrix $\mathbf{A} \in \mathbb{R}^{n \times n}$, denoted $|\mathbf{A}|$ or $\det \mathbf{A}$ is a function $\det: \mathbb{R}^{n \times n} \rightarrow \mathbb{R}$.

Let $\mathbf{A}_{\setminus i, \setminus j} \in \mathbb{R}^{(n-1) \times (n-1)}$ be the matrix that results from deleting the i^{th} row and j^{th} column. The general (recursive) formula for the determinant is

$$\begin{aligned} |\mathbf{A}| &= \sum_{i=1}^n (-1)^{i+j} a_{ij} |\mathbf{A}_{\setminus i, \setminus j}| \quad (\forall j \in 1, \dots, n) \\ &= \sum_{j=1}^n (-1)^{i+j} a_{ij} |\mathbf{A}_{\setminus i, \setminus j}| \quad (\forall i \in 1, \dots, n) \end{aligned} \quad (51)$$



Given a matrix

$$\mathbf{A} = \begin{bmatrix} \text{---} & \mathbf{a}_1^\top & \text{---} \\ \text{---} & \mathbf{a}_2^\top & \text{---} \\ & \vdots & \\ \text{---} & \mathbf{a}_n^\top & \text{---} \end{bmatrix} \quad (52)$$

and a set $\mathbf{S} \subset \mathbb{R}^n$,

$$\mathbf{S} = \{\mathbf{v} \in \mathbb{R}^n : \mathbf{v} = \sum_{i=1}^n \alpha_i \mathbf{a}_i \text{ where } 0 \leq \alpha_i \leq 1, i = 1, \dots, n\} \quad (53)$$

$|\mathbf{A}|$ is the volume of \mathbf{S} .



Example:

$$\mathbf{A} = \begin{bmatrix} 1 & 3 \\ 3 & 2 \end{bmatrix} \quad (54)$$



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The matrix rows are:

$$\mathbf{a}_1 = \begin{bmatrix} 1 \\ 3 \end{bmatrix} \quad \mathbf{a}_2 = \begin{bmatrix} 3 \\ 2 \end{bmatrix} \quad (55)$$

And $|\mathbf{A}| = -7$



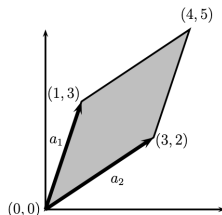
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Properties of determinants:

- ▶ For $\mathbf{A} \in \mathbb{R}^{n \times n}$, $|\mathbf{A}| = |\mathbf{A}^T|$
- ▶ For $\mathbf{A}, \mathbf{B} \in \mathbb{R}^{n \times n}$, $|\mathbf{AB}| = |\mathbf{A}||\mathbf{B}|$
- ▶ For $\mathbf{A} \in \mathbb{R}^{n \times n}$, $|\mathbf{A}| = 0$ iff \mathbf{A} is singular (i.e. non-invertible).
- ▶ For $\mathbf{A} \in \mathbb{R}^{n \times n}$ and \mathbf{A} non-singular, $|\mathbf{A}^{-1}| = 1/|\mathbf{A}|$



Given $\mathbf{A} \in \mathbb{R}^{n \times n}$ and a vector $\mathbf{x} \in \mathbb{R}^n$, the *quadratic form* is the scalar value

$$\mathbf{x}^\top \mathbf{A} \mathbf{x} = \sum_{i=1}^n x_i (\mathbf{A} \mathbf{x})_i = \sum_{i=1}^n x_i \left(\sum_{j=1}^n \mathbf{A}_{ij} x_j \right) = \sum_{i=1}^n \sum_{j=1}^n \mathbf{A}_{ij} x_i x_j \quad (56)$$



Some properties involving quadratic form:

- ▶ A symmetric matrix $\mathbf{A} \in \mathbb{S}^n$ is *positive definite* if for a non-zero $\mathbf{x} \in \mathbb{R}^n$, $\mathbf{x}^T \mathbf{A} \mathbf{x} > 0$
- ▶ A symmetric matrix $\mathbf{A} \in \mathbb{S}^n$ is *positive semi-definite* if for a non-zero $\mathbf{x} \in \mathbb{R}^n$, $\mathbf{x}^T \mathbf{A} \mathbf{x} \geq 0$
- ▶ A symmetric matrix $\mathbf{A} \in \mathbb{S}^n$ is *negative definite* if for a non-zero $\mathbf{x} \in \mathbb{R}^n$, $\mathbf{x}^T \mathbf{A} \mathbf{x} < 0$
- ▶ A symmetric matrix $\mathbf{A} \in \mathbb{S}^n$ is *negative semi-definite* if for a non-zero $\mathbf{x} \in \mathbb{R}^n$, $\mathbf{x}^T \mathbf{A} \mathbf{x} \leq 0$
- ▶ A symmetric matrix $\mathbf{A} \in \mathbb{S}^n$ is *indefinite* if it is neither positive nor negative semidefinite

n.b. Positive definite and negative definite matrices always have full rank.



Given $\mathbf{A} \in \mathbb{R}^{n \times n}$, $\lambda \in \mathbb{C}$ is an *eigenvalue* of \mathbf{A} with corresponding *eigenvector* $\mathbf{x} \in \mathbb{C}^n$ if

$$\mathbf{Ax} = \lambda\mathbf{x} : \mathbf{x} \neq \mathbf{0} \quad (57)$$

n.b. The eigenvector is (usually) normalized to have length 1



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We can write all of the eigenvector equations simultaneously as

$$\mathbf{AX} = \mathbf{X}\mathbf{\Lambda} \quad (58)$$

where

$$\mathbf{X} \in \mathbb{R}^{n \times n} = \begin{bmatrix} | & | & \cdots & | \\ \mathbf{x}_1 & \mathbf{x}_2 & \cdots & \mathbf{x}_n \\ | & | & \cdots & | \end{bmatrix}, \quad \mathbf{\Lambda} = \text{diag}(\lambda_1, \dots, \lambda_n) \quad (59)$$

This implies $\mathbf{A} = \mathbf{X}\mathbf{\Lambda}\mathbf{X}^{-1}$



Some properties:

- ▶ $\text{tr}\mathbf{A} = \sum_{i=1}^n \lambda_i$
- ▶ $|\mathbf{A}| = \prod_{i=1}^n \lambda_i$
- ▶ The rank of \mathbf{A} is equal to the number of non-zero eigenvalues of \mathbf{A} .
- ▶ If \mathbf{A} is non-singular, then $1/\lambda_i$ is an eigenvalue of \mathbf{A}^{-1} with corresponding eigenvector \mathbf{x}_i , i.e. $\mathbf{A}^{-1}\mathbf{x}_i = (1/\lambda_i)\mathbf{x}_i$
- ▶ The eigenvalues of a diagonal matrix $D = \text{diag}(d_1, \dots, d_n)$ are just its diagonal entries d_1, \dots, d_n



Example: For $\mathbf{A} \in \mathbb{S}^n$ with ordered eigenvalues $\lambda_1 \geq \lambda_2 \geq \dots \geq \lambda_n$,

$$\max_{\mathbf{x} \in \mathbb{R}^n} \mathbf{x}^T \mathbf{A} \mathbf{x} \text{ subject to } \|\mathbf{x}\|_2^2 = 1 \quad (60)$$

is solved with \mathbf{x}_1 corresponding to λ_1 . Similarly, it is solved with \mathbf{x}_n corresponding to λ_n .



Example:

Let $\mathbf{A} = \begin{bmatrix} 1 & 2 \\ 2 & 1 \end{bmatrix}$ Find the eigenvalues & eigenvectors.



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Example:

Let $\mathbf{A} = \begin{bmatrix} 1 & 2 \\ 2 & 1 \end{bmatrix}$ Find the eigenvalues & eigenvectors.

We want

$$(\mathbf{A} - \lambda\mathbf{I})\mathbf{x} = 0 \quad (61)$$

We want $\det(\mathbf{A} - \lambda\mathbf{I}) = 0$.

$$\det(\mathbf{A} - \lambda\mathbf{I}) = (1 - \lambda)^2 - 2^2 = \lambda^2 - 2\lambda - 3 \quad (62)$$

$$= (\lambda - 3)(\lambda + 1) \quad (63)$$

$\therefore \lambda = 3, -1$.



Finding the eigenvectors: calculating the null spaces of $(\mathbf{A} - \lambda \mathbf{I})$

$$\mathcal{N}(\mathbf{A} - 3\mathbf{I}) = \mathcal{N}\left(\begin{bmatrix} -2 & 2 \\ 2 & -2 \end{bmatrix}\right) = \begin{bmatrix} 1 \\ 1 \end{bmatrix} \quad (64)$$

$$\mathcal{N}(\mathbf{A} + \mathbf{I}) = \mathcal{N}\left(\begin{bmatrix} 2 & 2 \\ 2 & 2 \end{bmatrix}\right) = \begin{bmatrix} 1 \\ -1 \end{bmatrix} \quad (65)$$



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Thus:

$$\mathbf{X} = \begin{bmatrix} 1 & 1 \\ 1 & -1 \end{bmatrix}, \mathbf{\Lambda} = \begin{bmatrix} 3 & 0 \\ 0 & -1 \end{bmatrix} \quad (66)$$



SVD is a way of decomposing matrices.

Given $\mathbf{A} \in \mathbb{R}^{m \times n}$ with rank r , \exists
 $\Sigma \in \mathbb{R}^{m \times n}$, $\mathbf{U} \in \mathbb{R}^{m \times m}$, $\mathbf{V} \in \mathbb{R}^{n \times m}$ \ni

$$\mathbf{A} = \mathbf{U}\Sigma\mathbf{V}^T \quad (67)$$

Notes:

- ▶ Σ is a diagonal matrix with entries $\sigma_1, \dots, \sigma_r > 0$ known as *singular values*.
- ▶ \mathbf{U} and \mathbf{V} are orthogonal matrices.
- ▶ Common uses:
 - ▶ Least squares models
 - ▶ Range, rank, null space
 - ▶ Moore-Penrose inverse



Some intuition:

$\mathbf{A} \in \mathbb{R}^{m \times n}$ can be thought of as a linear transformation, such that for $\mathbf{x} \in \mathbb{R}^n$,

$$f(\mathbf{x}) = \mathbf{A}\mathbf{x} \quad (68)$$

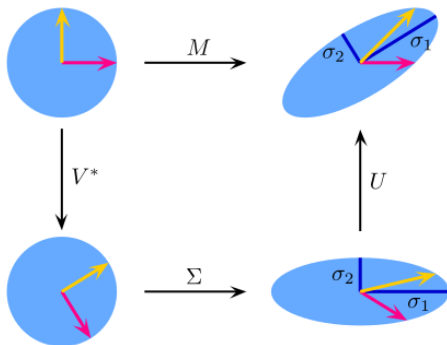


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SVD can be thought of as breaking this into individual steps:





Given $f : \mathbb{R}^{m \times n} \rightarrow \mathbb{R}$, the *gradient* of f wrt $\mathbf{A} \in \mathbb{R}^{m \times n}$ is

$$\nabla_{\mathbf{A}} f(\mathbf{A}) \in \mathbb{R}^{m \times n} = \begin{bmatrix} \frac{\partial f(\mathbf{A})}{\partial \mathbf{A}_{11}} & \frac{\partial f(\mathbf{A})}{\partial \mathbf{A}_{12}} & \cdots & \frac{\partial f(\mathbf{A})}{\partial \mathbf{A}_{1n}} \\ \frac{\partial f(\mathbf{A})}{\partial \mathbf{A}_{21}} & \frac{\partial f(\mathbf{A})}{\partial \mathbf{A}_{22}} & \cdots & \frac{\partial f(\mathbf{A})}{\partial \mathbf{A}_{2n}} \\ \vdots & \vdots & \ddots & \vdots \\ \frac{\partial f(\mathbf{A})}{\partial \mathbf{A}_{m1}} & \frac{\partial f(\mathbf{A})}{\partial \mathbf{A}_{m2}} & \cdots & \frac{\partial f(\mathbf{A})}{\partial \mathbf{A}_{mn}} \end{bmatrix} \quad (69)$$

Some properties

- ▶ $\nabla_{\mathbf{x}}(f(\mathbf{x}) + g(\mathbf{x})) = \nabla_{\mathbf{x}}f(\mathbf{x}) + \nabla_{\mathbf{x}}g(\mathbf{x})$
- ▶ For $c \in \mathbb{R}$, $\nabla_{\mathbf{x}}(c f(\mathbf{x})) = c \nabla_{\mathbf{x}}(f(\mathbf{x}))$



Given $f : \mathbb{R}^n \rightarrow \mathbb{R}$, the *Hessian* of f wrt $\mathbf{x} \in \mathbb{R}^n$ is

$$\nabla_{\mathbf{x}}^2 f(\mathbf{x}) \in \mathbb{R}^{n \times n} = \begin{bmatrix} \frac{\partial^2 f(\mathbf{x})}{\partial x_1^2} & \frac{\partial^2 f(\mathbf{x})}{\partial x_1 \partial x_2} & \cdots & \frac{\partial^2 f(\mathbf{x})}{\partial x_1 \partial x_n} \\ \frac{\partial^2 f(\mathbf{x})}{\partial x_2 \partial x_1} & \frac{\partial^2 f(\mathbf{x})}{\partial x_2^2} & \cdots & \frac{\partial^2 f(\mathbf{x})}{\partial x_2 \partial x_n} \\ \vdots & \vdots & \ddots & \vdots \\ \frac{\partial^2 f(\mathbf{x})}{\partial x_n \partial x_1} & \frac{\partial^2 f(\mathbf{x})}{\partial x_n \partial x_2} & \cdots & \frac{\partial^2 f(\mathbf{x})}{\partial x_n^2} \end{bmatrix} \quad (70)$$

n.b. The Hessian is always symmetric, since $\frac{\partial^2 f(\mathbf{x})}{\partial x_i \partial x_j} = \frac{\partial^2 f(\mathbf{x})}{\partial x_j \partial x_i}$



Given $\mathbf{A} \in \mathbb{R}^{m \times n}$, $\mathbf{b} \in \mathbb{R}^m \ni \mathbf{b} \notin \mathcal{R}(\mathbf{A})$, we want to find $\mathbf{x} \in \mathbb{R}^n$ as close as possible to \mathbf{b} (via the Euclidean norm),

$$\|\mathbf{Ax} - \mathbf{b}\|_2^2 = (\mathbf{Ax} - \mathbf{b})^\top (\mathbf{Ax} - \mathbf{b}) \quad (71)$$

$$= \mathbf{x}^\top \mathbf{A}^\top \mathbf{Ax} - 2\mathbf{b}^\top \mathbf{Ax} + \mathbf{b}^\top \mathbf{b} \quad (72)$$



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Taking the gradient wrt \mathbf{x} , we have

$$\begin{aligned} \nabla_{\mathbf{x}}(\mathbf{x}^\top \mathbf{A}^\top \mathbf{Ax} - 2\mathbf{b}^\top \mathbf{Ax} + \mathbf{b}^\top \mathbf{b}) &= \nabla_{\mathbf{x}} \mathbf{x}^\top \mathbf{A}^\top \mathbf{Ax} - \nabla_{\mathbf{x}} 2\mathbf{b}^\top \mathbf{Ax} + \nabla_{\mathbf{x}} (\mathbf{b}^\top \mathbf{b}) \\ &= \mathbf{A}^\top \mathbf{Ax} - 2\mathbf{A}^\top \mathbf{b} \end{aligned} \quad (74)$$



Given $\mathbf{A} \in \mathbb{R}^{m \times n}$, $\mathbf{b} \in \mathbb{R}^m \ni \mathbf{b} \notin \mathcal{R}(\mathbf{A})$, we want to find $\mathbf{x} \in \mathbb{R}^n$ as close as possible to \mathbf{b} (via the Euclidean norm),

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Setting this expression equal to zero and solving for \mathbf{x} gives the normal equations,

$$\mathbf{x} = (\mathbf{A}^\top \mathbf{A})^{-1} \mathbf{A}^\top \mathbf{b} \quad (75)$$



Some textbooks on linear algebra:

- ▶ *Linear Algebra (Jim Hefferon)*
- ▶ *Introduction to Applied Linear Algebra (Boyd & Vandenberghe)*
- ▶ *Linear Algebra (Cherney, Denton et al.)*
- ▶ *Linear Algebra (Hoffman & Kunze)*
- ▶ *Fundamentals of Linear Algebra (Carrell)*
- ▶ *Linear Algebra (S. Friedberg A. Insel L. Spence)*



Probability



The set of all possible values is called the *sample space* S .

- ▶ It's the space where realizations can be produced.



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Example: Tossing a coin

$$S = \{Heads, Tails\} \quad (76)$$



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Example: Tossing a coin

$$S = \{Heads, Tails\} \quad (76)$$

More notation:

- ▶ \emptyset is the *empty set*. Can be denoted as $\emptyset = \{\}$.
- ▶ $\bigcup_{i=1}^{\infty} B_i$ is the union of sets B_i . Formally,
 - ▶ $\bigcup_{i=1}^{\infty} B_i = \{s \in S : s \in B_i \forall i\}$
- ▶ $B \subseteq S$ means B is a *subset* of the sample space.
- ▶ *Heads*, without curly braces, is an *element* of set B .
- ▶ $B^C = S \setminus B$ is the complement of set B



A *probability function* is a function $P : \mathcal{B} \rightarrow [0, 1]$, where

- ▶ $P(S) = 1$
- ▶ $P(\cup_{i=1}^{\infty} B_i) = \sum_{i=1}^{\infty} P(B_i)$ when B_1, B_2, \dots are disjoint



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n.b. We can define the domain \mathcal{B} many ways, e.g. $\mathcal{B} = 2^S$



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n.b. We can define the domain \mathcal{B} many ways, e.g. $\mathcal{B} = 2^S$

Example: For flipping a coin, we have

$$\mathcal{B} = 2^S = \{\emptyset, \{Heads\}, \{Tails\}, \{Heads, Tails\}\} \quad (77)$$

This implies that

$$P(B) = \begin{cases} 1 & B = \{Heads, Tails\} \\ \frac{1}{2} & B = \{Heads\} \\ \frac{1}{2} & B = \{Tails\} \\ 0 & B = \emptyset \end{cases} \quad (78)$$

n.b. The power set is a 'set of sets'



Problem: Power sets don't work well for \mathbb{R} .



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Solution: Define the domain using σ -algebra:

- ▶ $\emptyset \in \mathcal{B}$
- ▶ $B \in \mathcal{B} \Rightarrow B^C \in \mathcal{B}$
- ▶ $B_1, B_2, \dots \in \mathcal{B} \Rightarrow \bigcup_{i=1}^{\infty} B_i \in \mathcal{B}$



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Example:

- ▶ The *discrete* σ -algebra:
 $\mathcal{B} = 2^S = \{\emptyset, \{Heads\}, \{Tails\}, \{Heads, Tails\}\}$
- ▶ The *trivial* σ -algebra: $\mathcal{B} = \emptyset \cup S = \{\emptyset, \{Heads, Tails\}\}$

n.b. For uncountable sets, we use the *Borel* σ -algebra.

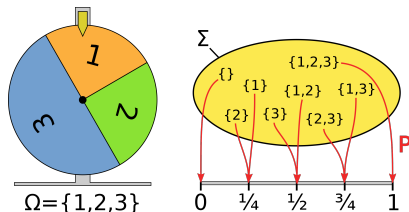


Def:

A *probability space* is a triple (S, \mathcal{B}, P) .

- ▶ S is the set of possible singleton events
- ▶ \mathcal{B} is the set of questions to ask P
- ▶ P maps sets into probabilities

n.b. They represent the ingredients needed to talk about probabilities





Some properties of $P(\cdot)$

- ▶ $P(B) = 1 - P(B^C)$
- ▶ $P(\emptyset) = 0$, since $P(\emptyset) = 1 - P(S)$
- ▶ $P(A \cup B) = P(A) + P(B) - P(A \cap B)$, implying that
 - ▶ $P(A \cup B) \leq P(A) + P(B)$
 - ▶ $P(A \cap B) \geq P(A) + P(B) - 1$



For events A and B where $P(B) > 0$, the *conditional probability* of A given B (denoted $P(A|B)$) is

$$P(A|B) = \frac{P(A \cap B)}{P(B)} \quad (79)$$

Example: In an agricultural region with 1000 farms, we want to know if the farm has vineyards or cork trees.

		Cork Trees	
		Yes	No
Vineyard	Yes	200	50
	No	150	600

Table: Frequency counts



Example: In an agricultural region with 1000 farms, we want to know if the farm has vineyards or cork trees.

		Cork Trees	
		Yes	No
Vineyard	Yes	20%	5%
	No	15%	60%

Table: Joint probabilities

Questions:

- ▶ What is the probability of seeing cork trees in a farm with vineyards?
- ▶ Among farms with cork trees or vineyards, what is the probability of having both?



Let's assume the following joint probabilities

		Cork Trees	
		Yes	No
Vineyard	Yes	25%	25%
	No	25%	25%

We have that $P(A \cap B) = P(A) \cdot P(B)$, meaning that they are *independent*



Let $B_1, B_2, \dots, B_k \in \mathcal{B}$ and $P(B_i) > 0 : i = 1, \dots, k$. The *law of total probability* states that

$$P(A) = \sum_{i=1}^k P(B_i)P(A|B_i) \quad (80)$$



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The *conditional law of total probability* states that

$$P(A|C) = \sum_{i=1}^k P(B_i|C)P(A|B_i, C) \quad (81)$$



Let $B_1, B_2, \dots, B_k \in \mathcal{B}$, $P(B_i) > 0 : i = 1, \dots, k$, and $P(A) > 0$.
Then Bayes' Theorem states that for $i = 1, \dots, k$

$$P(B_i|A) = \frac{P(B_i)P(A|B_i)}{\sum_{j=1}^k P(B_j)P(A|B_j)} \quad (82)$$

n.b. Can be proven using the def of conditional probability



Example: You test positive for disease X , which has 90% sensitivity and a FPR of 10%. Past genetic screening has indicated that you have a 1 in 10,000 chance of having the disease. What is the probability of having disease X ?



Example: You test positive for disease X , which has 90% sensitivity and a FPR of 10%. Past genetic screening has indicated that you have a 1 in 10,000 chance of having the disease. What is the probability of having disease X ?

$$P(B_1|A) = \frac{P(A|B_1)P(B_1)}{P(A|B_1)P(B_1) + P(A|B_2)P(B_2)} \quad (83)$$

$$= \frac{(0.9)(0.0001)}{(0.9)(0.0001) + (0.1)(0.9999)} = 0.0009 \quad (84)$$



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Notes:

- ▶ $P(B_1)$ is often referred to as the *prior* probability
- ▶ $P(B_1|A)$ is often referred to as the *posterior* probability



A *random variable* is a (Borel measurable) function

$$X : \mathcal{S} \rightarrow \mathbb{R}$$

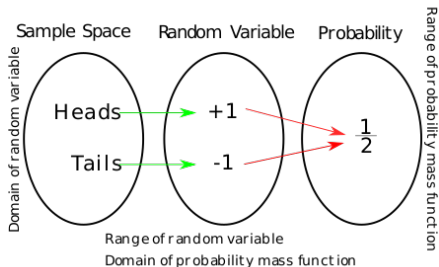


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Example: For coin tossing, we have $X : \{Heads, Tails\} \rightarrow \mathbb{R}$, where

$$X(s) = \begin{cases} 1 & \text{if } s = \text{Heads} \\ 0 & \text{if } s = \text{Tails} \end{cases} \quad (85)$$



Cumulative distribution function



The *cumulative distribution function* (cdf) of a random variable X is the function $F_X : \mathbb{R} \rightarrow [0, 1]$.



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$$F_X(x) = \begin{cases} 0 & \text{if } x < 0 \\ \frac{1}{2} & \text{if } 0 \leq x < 1 \\ 1 & \text{if } x \geq 1 \end{cases} \quad (87)$$



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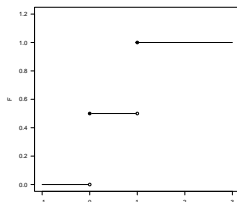
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n.b. We have two ways of thinking about probabilities:

1. Probability functions
2. Cumulative distribution functions

Question: Which one should we use?



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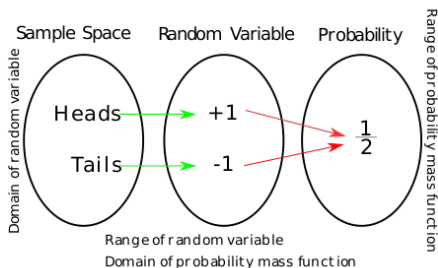
The Correspondence Theorem: Let $P_X(\cdot)$ and $P_Y(\cdot)$ be probability functions and $F_X(\cdot)$ and $F_Y(\cdot)$ be their associated cdfs. Then

$$P_X(\cdot) = P_Y(\cdot) \iff F_X(\cdot) = F_Y(\cdot) \quad (88)$$



Some properties for cdfs:

- ▶ $\lim_{x \rightarrow -\infty} F(x) = 0$
- ▶ $\lim_{x \rightarrow \infty} F(x) = 1$
- ▶ $F(\cdot)$ is non-decreasing
- ▶ $F(\cdot)$ is right-continuous





Let X be a continuous rv and one-to-one over the the possible values of X . Then

$$F^{-1}(p) = \inf\{x \in \mathbb{R} : p \leq F(x)\} \quad (89)$$

Is the quantile function of X .

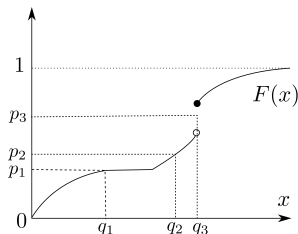


Let X be a continuous rv and one-to-one over the the possible values of X . Then

$$F^{-1}(p) = \inf\{x \in \mathbb{R} : p \leq F(x)\} \quad (89)$$

Is the quantile function of X . Let X be a *discrete* rv and one-to-one over the the possible values of X . Then $F^{-1}(p)$ states that we take the smallest value of x .

Example:





A random variable X is

- ▶ **Discrete** if $\exists f_X : \mathbb{R} \rightarrow [0, 1] \ni F_X(x) = \sum_{t \leq x} f_X(t), x \in \mathbb{R}$
 - ▶ f_X is referred to as the probability mass function (pmf)
- ▶ **Continuous** if $\exists f_X : \mathbb{R} \rightarrow \mathbb{R}_+ \ni F_X(x) = \int_{-\infty}^x f_X(t) dt, x \in \mathbb{R}$
 - ▶ f_X is referred to as the probability density function (pdf).
 - ▶ n.b. We can have multiple pdf's consistent with the same cdf.
 - ▶ n.b. For any specific value of a continuous random variable, its probability is 0, i.e. $P(\{x\}) = 0 \forall x \in \mathbb{R}$.



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n.b. pmf's and pdf's sum to 1, i.e.

- ▶ $f : \mathbb{R} \rightarrow [0, 1]$ is the pmf of a discrete RV iff $\sum_{x \in \mathbb{R}} f(x) = 1$
- ▶ $f : \mathbb{R} \rightarrow \mathbb{R}_+$ is the pdf of a continuous RV iff $\int_{-\infty}^{\infty} f(x) dx = 1$



Example #1: Coin tossing

$$F_X(x) = \begin{cases} 0 & \text{if } x < 0 \\ \frac{1}{2} & \text{if } 0 \leq x < 1 \\ 1 & \text{if } x \geq 1 \end{cases} \quad (90)$$

Here, F_X is a step function with pmf

$$f_X(x) = \begin{cases} \frac{1}{2} & x \in \{0, 1\} \\ 0 & \text{otherwise} \end{cases} \quad (91)$$



Example #2: Uniform distribution on $(0,1)$

$$F_X(x) = \begin{cases} 0 & \text{if } x < 0 \\ x & \text{if } 0 \leq x < 1 \\ 1 & \text{if } x \geq 1 \end{cases} \quad (92)$$

Here, F_X is a continuous function. Two consistent pdfs include

$$f_X(x) = \begin{cases} 1 & x \in [0, 1] \\ 0 & \text{otherwise} \end{cases} \quad (93)$$

$$f_X(x) = \begin{cases} 1 & x \in (0, 1) \\ 0 & \text{otherwise} \end{cases} \quad (94)$$



Suppose $Y = g(X)$, where $g : \mathbb{R} \rightarrow \mathbb{R}$ and X is a *discrete* rv with cdf F_X .



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Since the function is applied to a rv, Y is also a random variable with probability function

$$f_Y(y) = P_Y(g(X) = y) = \sum_{x:g(x)=y} f_X(x) \quad (95)$$



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Example:

Let X be a uniform random variable on $\{-n, -n + 1, \dots, n - 1, n\}$. Then $Y = |X|$ has mass function

$$f_Y(y) = \begin{cases} \frac{1}{2n+1} & \text{if } x = 0 \\ \frac{2}{2n+1} & \text{if } x \neq 0 \end{cases} \quad (96)$$



Suppose $Y = g(X)$, where $g : \mathbb{R} \rightarrow \mathbb{R}$ and rv X with cdf F_X .



Suppose $Y = g(X)$, where $g : \mathbb{R} \rightarrow \mathbb{R}$ and rv X with cdf F_X .

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$$F_Y(y) = P(Y \leq y) = P(g(X) \leq y) = \int_{x : g(x) \leq y} f_X(x) dx \quad (97)$$

We can get the probability function by taking the derivative

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Example:

Let X be a uniform rv on $[-1, 1]$. Then $Y = X^2$ has cdf

$$\begin{aligned} F_Y(y) &= P_Y(Y \leq y) = P_X(X^2 \leq y) = P_X(-y^{1/2} \leq X \leq y^{1/2}) \\ &= \int_{-y^{1/2}}^{y^{1/2}} f(x) dx = y^{1/2} \end{aligned} \quad (99)$$

and $f_Y(y) = \frac{\partial}{\partial y} F_Y(y) = \frac{1}{2y^{1/2}}$



Suppose $Y = g(X) = aX + b$, $a > 0$, $b \in \mathbb{R}$. Then

$$P(Y \leq y) = P(aX + b \leq y) = P\left(X \leq \frac{y - b}{a}\right) = F_X\left(\frac{y - b}{a}\right) \quad (100)$$



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In general, as long as the transformation $Y = g(X)$ is monotonic, then

$$f_Y(y) = f_X(g^{-1}(y)) \left| \frac{\partial}{\partial y} g^{-1}(y) \right| \quad (102)$$



- ▶ Grinstead & Snell Chapters 1,2,4
- ▶ DeGroot & Schervish Chapters 1,2,3



Statistics



The *expected value* of rv X is defined as

$$\mathbb{E}[X] = \begin{cases} \sum_x x f_X(x) & \text{if } x \text{ is discrete} \\ \int x f_X(x) dx & \text{if } x \text{ is continuous} \end{cases} \quad (103)$$

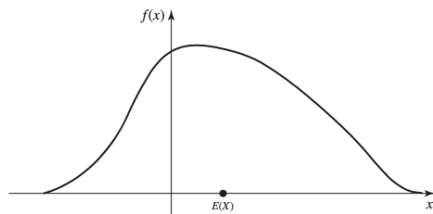
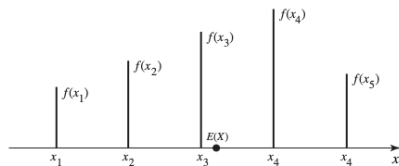
For functions g of X ,

$$\mathbb{E}[g(X)] = \begin{cases} \sum_x g(x) f_X(x) & \text{if } x \text{ is discrete} \\ \int g(x) f_X(x) dx & \text{if } x \text{ is continuous} \end{cases} \quad (104)$$

n.b. In general, $\mathbb{E}[g(X)] \neq g(\mathbb{E}[X])$



Examples:





Important: Expectations might not exist!

Example: Suppose $f_X(x) = \frac{1}{x^2}$, defined on $[1, \infty]$. Then

$$\mathbb{E}[X] = \int xf_X(x)dx = \int x\frac{1}{x^2}dx = \int \frac{1}{x}dx = \infty \quad (105)$$



Important: Expectations might not exist!

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Some properties of expectations:

- ▶ Linearity: $\mathbb{E}[ag(X) + bh(X)] = \mathbb{E}[ag(X)] + \mathbb{E}[bh(X)]$
- ▶ Order preserving:
 $g(X) \leq h(X), \forall x \in \mathbb{R} \Rightarrow \mathbb{E}[g(X)] \leq \mathbb{E}[h(X)]$



The *variance* of rv X is defined as

$$\text{var}(X) = \mathbb{E}[(X - \mu)^2] : \mu = \mathbb{E}[X] \quad (106)$$



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Some notes:

- ▶ If $\mathbb{E}[X]$ doesn't exist then $\text{var}(X)$ doesn't exist.
- ▶ $\text{var}(X)$ can be infinite.
- ▶ The standard deviation σ of X is $\sqrt{\text{var}(X)}$.



With some algebra, we see that

$$\text{var}(X) = \mathbb{E}[(X - \mu)^2] \quad (107)$$

$$= \mathbb{E}[X^2 - 2X\mu + \mu^2] \quad (108)$$

$$= \mathbb{E}[X^2] - \mathbb{E}[2X\mu] + \mathbb{E}[\mu^2] \quad (109)$$

$$= \mathbb{E}[X^2] - \mu^2 \quad (110)$$

$$= \mathbb{E}[X^2] - \mathbb{E}[X]^2 \quad (111)$$



Some properties:

- ▶ If X is bounded, then $\text{var}(X)$ exists and is finite.
- ▶ $\text{var}(X) = 0 \iff P(X = c) = 1$ for some constant c .
- ▶ $\text{var}(cX) = c^2 \text{var}(X)$ for some constant c .
- ▶ variance is linear, i.e. $\text{var}(X_1 + X_2) = \text{var}(X_1) + \text{var}(X_2)$.



The k^{th} *moment* of rv X is defined as

$$\mathbb{E}[X^k] = \mu'_k : k \in \mathbb{N} \quad (112)$$

The k^{th} *central/centered moment* of rv X is defined as

$$\mathbb{E}[(X - \mu)^k] = \mu_k : k \in \mathbb{N} \quad (113)$$



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Notes:

- ▶ μ'_k exists if and only if $\mathbb{E}[|X|^k] < \infty$.
- ▶ If μ'_k exists, then for all $j < k$, μ'_j also exists.
- ▶ Variance is μ_2 .
- ▶ *Skewness* is μ_3/σ^2 .
- ▶ *Kurtosis* is μ_4/σ^4 .



Example: Suppose $X \sim N(0, 1) \ni f_X(x) = \frac{1}{\sqrt{2\pi}} \exp\left(-\frac{x^2}{2}\right)$.

$$\mu_1^i = \mathbb{E}[X] = \int xf_X(x)dx = f_X(x)|_{-\infty}^{\infty} = 0 \quad (114)$$

n.b. For the normal distribution, $xf_X(x) = -\frac{\partial}{\partial x} f_X(x)$.



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n.b. For the normal distribution, $x f_X(x) = -\frac{\partial}{\partial x} f_X(x)$.

$$\mu_2 = \mathbb{E}[(X - \mu)^2] = \mathbb{E}[(X - 0)^2] = \mathbb{E}[X^2] = \int x^2 f_X(x) dx \quad (115)$$

using integration by parts, we get

$$\int x^2 f_X(x) dx = \underbrace{-x f_X(x) \Big|_{-\infty}^{\infty}}_{=0} + \underbrace{\int_{-\infty}^{\infty} f_X(x) dx}_{=1} = 1 \quad (116)$$



Moment generating functions (mgf) are used to calculate the moments of a rv. The mgf of a rv X is a function $M_X : \mathbb{R} \Rightarrow \mathbb{R}_+$ such that

$$M_X(t) = \mathbb{E}[e^{tX}] : t \in \mathbb{R} \quad (117)$$



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Notes:

- ▶ The mgf is a function of t ; X is integrated out by \mathbb{E} .
- ▶ The mgf only applies if the moments of the rv exists.
- ▶ If two rv X, Y have the same mgf (i.e. $M_X(t) = M_Y(t)$), then they have the same distribution.
- ▶ Even if a rv has moments, the mgf may yield infinity (e.g. log-normal distribution).



Taking the derivative of the mgf, we see that

$$\frac{\partial}{\partial t} M_X(t) = \frac{\partial}{\partial t} \int e^{tx} f_X(x) dx = \int x \cdot e^{tx} f_X(x) dx \quad (118)$$

What happens when $t = 0$?



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What happens when $t = 0$ for the k^{th} derivative?

$$\frac{\partial}{\partial t^k} M_X(t) = \int x^k \cdot e^{tx} f_X(x) dx \quad (120)$$

At $t = 0$, we get $\frac{\partial}{\partial t^k} M_X(t)|_{t=0} = \mathbb{E}[X^k]$

Evaluating the k^{th} derivative at $t = 0$ gives us the k^{th} moment of X .



Example: The standard normal distribution

$$M_X(t) = \mathbb{E}[e^{tX}] = \int e^{tX} f_X(x) dx \quad (121)$$

$$= \int e^{tX} \frac{1}{\sqrt{2\pi}} \exp\left(-\frac{x^2}{2}\right) dx \quad (122)$$

$$= \int \frac{1}{\sqrt{2\pi}} \exp\left(-\frac{(x-t)^2}{2}\right) \exp\left(\frac{t^2}{2}\right) dx \quad (123)$$

$$= \exp\left(\frac{t^2}{2}\right) \int \frac{1}{\sqrt{2\pi}} \exp\left(-\frac{(x-t)^2}{2}\right) dx \quad (124)$$

$$= \exp\left(\frac{t^2}{2}\right) \quad (125)$$



The mgf for *affine transformations* is straight forward, e.g. If $Y = aX + b$, then $M_Y(t) = e^{bt} M_X(at)$.

Example: Let $X = \mu + \sigma Z : Z \sim N(0, 1)$. Then

$$M_X(t) = M_{\mu + \sigma Z}(t) = e^{\mu t} M_Z(\sigma t) = e^{\mu t} e^{\frac{1}{2}\sigma^2 t^2} = e^{\mu t + \frac{1}{2}\sigma^2 t^2} \quad (126)$$



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Another example:

Let $X_1, \dots, X_n \stackrel{iid}{\sim} P_0$ and $Y = \sum_{i=1}^n X_i$. Then

$$M_Y(t) = \mathbb{E}[e^{tY}] = \mathbb{E}[e^{t(X_1 + \dots + X_n)}] = \mathbb{E} \left[\prod_{i=1}^n e^{tX_i} \right] \quad (127)$$

$$= \prod_{i=1}^n \mathbb{E} \left[e^{tX_i} \right] = \prod_{i=1}^n M_{X_i}(t) \quad (128)$$



Most useful distributions have names, e.g.

- ▶ Normal distribution
- ▶ Uniform distribution
- ▶ Bernoulli distribution
- ▶ Binomial distribution
- ▶ Poisson distribution
- ▶ Gamma distribution



A rv X follows a *Normal distribution*, denoted as $X \sim N(\mu, \sigma^2)$ with mean μ and variance σ^2 , if X is continuous with pdf

$$f_X(x) = \frac{1}{\sqrt{2\pi\sigma^2}} \exp\left(-\frac{(x-\mu)^2}{2\sigma^2}\right) : x \in \mathbb{R} \quad (129)$$

Note:

If $Z \sim N(0, 1)$ then $X = \mu + \sigma Z \sim N(\mu, \sigma^2)$. It follows that

- ▶ $\mathbb{E}[X] = \mathbb{E}[\mu + \sigma Z] = \mu + \sigma \mathbb{E}[Z] = \mu$.
- ▶ $\text{var}(X) = \text{var}(\mu + \sigma Z) = \sigma^2 \text{var}(Z) = \sigma^2$.



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- ▶ $\text{var}(X) = \text{var}(\mu + \sigma Z) = \sigma^2 \text{var}(Z) = \sigma^2$.

Most well known distribution due to:

1. Good mathematical properties
2. Often (approximately) observed in the real world (e.g. heights, weights, etc.)
3. Central limit theorem



Let $X_1, \dots, X_n \stackrel{iid}{\sim} P_0$, where $\mathbb{E}[X_i] = \mu$ and $\text{var}(X_i) = \sigma^2$.

Then

$$\lim_{n \rightarrow \infty} P \left(\frac{n^{1/2}(\bar{X}_n - \mu)}{\sigma} \leq x \right) = \Phi(x) \quad (130)$$

where $\Phi(x)$ is the cdf for the standard normal distribution.



Let $X_1, \dots, X_n \stackrel{iid}{\sim} P_0$, where $\mathbb{E}[X_i] = \mu$ and $\text{var}(X_i) = \sigma^2$.

Then

$$\lim_{n \rightarrow \infty} P \left(\frac{n^{1/2}(\bar{X}_n - \mu)}{\sigma} \leq x \right) = \Phi(x) \quad (130)$$

where $\Phi(x)$ is the cdf for the standard normal distribution.

Example: The sample mean

$$\bar{X}_n = \frac{1}{n} \sum_{i=1}^n X_i \quad (131)$$

The 95% CI: $\bar{X}_n \pm z_{\alpha/2} \hat{\text{se}}_n$



A rv X follows a Uniform distribution $U(a, b)$ if X is continuous with pdf

$$f_X(x) = \begin{cases} \frac{1}{b-a} & x \in [a, b] \\ 0 & \text{otherwise} \end{cases} \quad (132)$$

Under $U(a, b)$, all observations are “*equally likely*”

$$\mathbb{E}[X] = \frac{a+b}{2}, \text{ var}(X) = \frac{(b-a)^2}{12}, \text{ and } M_X(t) = \frac{e^{bt} - e^{at}}{(b-a)t}.$$



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Note: if $X \sim U(a, b)$, then $X = (b - a)\tilde{X} + a : \tilde{X} \sim U(0, 1)$
and

$$f_{\tilde{X}}(x) = \begin{cases} 1 & x \in [0, 1] \\ 0 & \text{otherwise} \end{cases} \quad (133)$$



A rv X follows a Bernoulli distribution $Ber(p)$ if X is discrete with pmf

$$f_X(x) = \begin{cases} p & \text{if } x = 1 \\ 1 - p & \text{if } x = 0 \\ 0 & \text{otherwise} \end{cases} \quad (134)$$

$\mathbb{E}[X] = p$, $\text{var}(X) = p(1 - p)$, and $M_X(t) = e^t p + (1 - p)$.



A rv X follows a Binomial distribution $Bin(n, p)$ if X is discrete with pmf

$$f_X(x) = \begin{cases} \binom{n}{x} p^x (1-p)^{n-x} & \text{if } x \in \{0, 1, \dots, n\} \\ 0 & \text{otherwise} \end{cases} \quad (135)$$

$\mathbb{E}[X] = np$, $\text{var}(X) = np(1-p)$, and

$M_X(t) = (e^t p + (1-p))^n$.

If $X_1, \dots, X_n \stackrel{iid}{\sim} Ber(p)$, then $Y = X_1 + \dots + X_n$ follows $B(n, p)$.



A rv X follows a Negative Binomial distribution $NB(r, p)$ if X is discrete with pmf

$$f_X(x) = \begin{cases} \binom{r+x-1}{x} p^x (1-p)^r & \text{if } x \in \{0, 1, \dots, n\} \\ 0 & \text{otherwise} \end{cases} \quad (136)$$

$\mathbb{E}[X] = \frac{r(1-p)}{p}$, $\text{var}(X) = \frac{r(1-p)}{p^2}$, and

$$M_X(t) = \left(\frac{p}{1-qe^t} \right)^r : t < \log \left(\frac{1}{q} \right).$$

When $r = 1$, we refer to it as the *Geometric distribution*.

- ▶ It has a *memoryless* property.



A rv X follows a Poisson distribution $Pois(\lambda)$ if X is discrete with pmf

$$f_X(x) = \begin{cases} e^{-\lambda} \frac{\lambda^x}{x!} & x \in \mathbb{N} \\ 0 & \text{otherwise} \end{cases} \quad (137)$$

$\mathbb{E}[X] = \lambda$, $var(X) = \lambda$, and $M_X(t) = e^{\lambda(e^t-1)}$.

Some notes:

- ▶ $Bin(n, p) \approx Pois(np)$ when n is large and np is small.
- ▶ “Poisson Processes” are typically used to model rates, e.g. mortality rates
 1. The number of events in each fixed time interval t has a Poisson distribution with mean λt .
 2. The number of events in each time interval is independent.



A rv X follows a Gamma distribution $\text{Gamma}(\alpha, \beta)$ if X is continuous with pdf

$$f_X(x) = \begin{cases} \frac{1}{\Gamma(\alpha)\beta^\alpha} x^{\alpha-1} e^{-\frac{x}{\beta}} & x > 0 \\ 0 & \text{otherwise} \end{cases} \quad (138)$$

where $\Gamma(x) = \int_0^\infty t^{x-1} e^{-t} dt : x > 0$.

$\mathbb{E}[X] = \alpha\beta$, $\text{var}(X) = \alpha\beta^2$, and

$M_X(t) = \left(1 - \frac{t}{\beta}\right)^{-\alpha} : t < \beta$.



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$\mathbb{E}[X] = \alpha\beta$, $\text{var}(X) = \alpha\beta^2$, and

$M_X(t) = \left(1 - \frac{t}{\beta}\right)^{-\alpha} : t < \beta$.

Notes:

- ▶ $\frac{1}{\Gamma(\alpha)\beta^\alpha}$ is often referred to as the '*normalizing constant*'.
- ▶ When $\alpha = 1$, we get the exponential distribution.



A rv X follows a Beta distribution $Beta(\alpha, \beta)$ if X is continuous with pdf

$$f_X(x) = \begin{cases} \frac{\Gamma(\alpha+\beta)}{\Gamma(\alpha)\Gamma(\beta)} x^{\alpha-1} (1-x)^{\beta-1} & 0 < x < 1 \\ 0 & \text{otherwise} \end{cases} \quad (139)$$

$\mathbb{E}[X] = \frac{\alpha}{\alpha+\beta}$, $\text{var}(X) = \frac{\alpha\beta}{(\alpha+\beta)^2(\alpha+\beta+1)}$, and

$$M_X(t) = \frac{\Gamma(\alpha+\beta)}{\Gamma(\alpha)\Gamma(\beta)} \int_0^1 x^{\alpha+k-1} (1-x)^{\beta-1} dx.$$

n.b. Very popular distribution in Bayesian statistics.



Suppose rv $\mathbf{X} = (X_1, \dots, X_k)$ represents counts of k different classes. Then it follows a Multinomial distribution $Multi(p_1, \dots, p_k)$ if it has pdf

$$f_{\mathbf{X}}(\mathbf{x}) = \begin{cases} \binom{n}{x_1, \dots, x_k} p_1^{x_1} \cdots p_k^{x_k} & x_1 \geq 0, \dots, x_k \geq 0 \\ 0 & \text{otherwise} \end{cases} \quad (140)$$

where $n = \sum_{i=1}^k X_i$.

$\mathbb{E}[X_i] = np$, $\text{var}(X_i) = np_i(1 - p_i)$, and
 $\text{Cov}(X_i, X_j) = -np_i p_j$.



While not technically a pdf, often used for e.g. mixture of discrete distributions

The Dirac delta function is defined as $\delta : \mathbb{R} \rightarrow \mathbb{R} \cup \infty \ni$

$$\delta(x) = \begin{cases} +\infty & x = 0 \\ 0 & \text{otherwise} \end{cases} \quad (141)$$

and $\int_{-\infty}^{\infty} \delta(x) dx = 1$

The sifting property:

$$\int f(x) \delta(x - a) dx = f(a) \quad (142)$$



Example: Let

$$Y = \begin{cases} 1 & \text{w.p. } \alpha \\ U(0, 1) & \text{w.p. } 1 - \alpha \end{cases} \quad (143)$$

Then $f_Y(y) = \alpha\delta(y - 1) + (1 - \alpha)\mathbb{I}(y \in [0, 1])$



Example: Let

$$Y = \begin{cases} 1 & \text{w.p. } \alpha \\ U(0, 1) & \text{w.p. } 1 - \alpha \end{cases} \quad (143)$$

Then $f_Y(y) = \alpha\delta(y - 1) + (1 - \alpha)\mathbb{I}(y \in [0, 1])$

$$\mathbb{E}[Y] = \int_{-\infty}^{\infty} y(\alpha\delta(y - 1) + (1 - \alpha)\mathbb{I}(y \in [0, 1]))dy \quad (144)$$

$$= \alpha \int_{-\infty}^{\infty} y(\delta(y - 1))dy + (1 - \alpha) \int_0^1 ydy \quad (145)$$

$$= \alpha + (1 - \alpha) \frac{y^2}{2} \Big|_0^1 \quad (146)$$

$$= \alpha + \frac{1 - \alpha}{2} \quad (147)$$

$$= \frac{1 + \alpha}{2} \quad (148)$$



- ▶ DeGroot & Schervish Chapters 4.1-4.5,5.1-5.9
- ▶ Grinstead & Snell Chapters 5,6