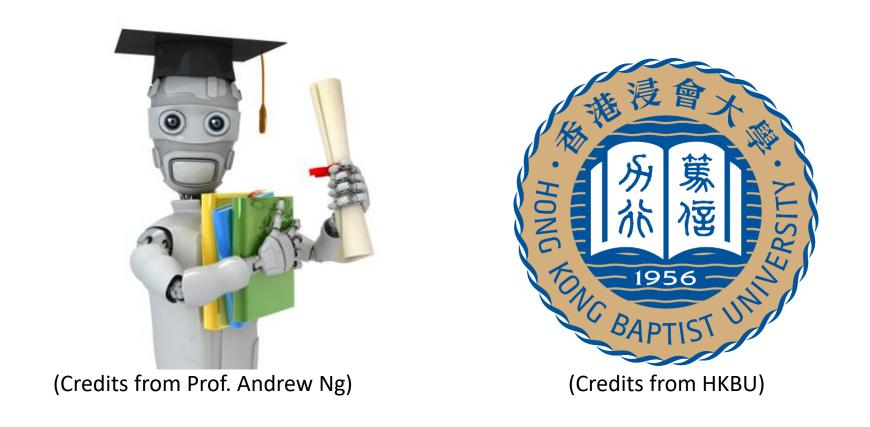
COMP7180: Quantitative Methods for DAAI



Course Instructors: Dr. Yang Liu and Dr. Bo Han Teaching Assistant: Mr. Minghao Li

Course Contents

- Continuous and Discrete Random Variables (Week 7)
- Conditional Probability and Independence (Week 8)
- Maximum Likelihood Estimation (Week 9)
- Mathematical Optimization (Week 10)
- Convex and Non-Convex Optimization (Week 11)
- Quiz and Course Review (Week 12) Our Focus

A probability can be regard as a function to estimate the value of every event.

As a function, we should have a domain (定义域). What is the domain?

Given a sample space S: set of all possible outcomes of an experiment. The domain consists of some subsets of S.

An element E in the domain is called event.

Example: Toss a coin (1 time). Then, the outcome is H or T, where H is the head of a coin and T is the tail of a coin.

Then S= { H, T};

The domain is $\{ \{H,T\}, \{H\}, \{T\}, \emptyset \}$.

{H,T}, {H}, {T}, Ø are called events.



Example: Toss a coin (1 time). Then, the outcome is H or T, where H is the head of a coin and T is the tail of a coin.

Then $S = \{H, T\}$; The domain is $\{\{H,T\}, \{H\}, \{T\}, \emptyset\}$. $\{H,T\}, \{H\}, \{T\}, \emptyset$ are called events.

- S and Ø should be event;
- $S^C = \emptyset$; $\{H\}^C = \{T\}$; $\{T\}^C = \{H\}$; $\emptyset^C = S$;
- $S \cap \emptyset = \emptyset$; $S \cap \{H\} = \{H\}$; $S \cap \{T\} = \{T\}$; $\{H\} \cap \{T\} = \emptyset$;
- $\{H\}\cup\{T\}=S$; $\{H\}\cup\emptyset=\{H\}$; $\{T\}\cup\emptyset=\{T\}$.

As a function, we should have a range (值域). What is the range?

Given an event E, a probability maps E into [0,1], that is $0 \le P(E) \le 1$.

If P(E)=0, then this event E will not occur.

If P(E)=1, then this event E occurs without uncertainty.

Example. Toss a coin (1 time). There are outcomes: H and T, where H is the head of a coin and T is the tail of a coin.

S= {H, T}; The domain is {{H,T}, {H}, {T}, \emptyset }.

{H,T}, {H}, {T}, Ø are called events.

 $P({H,T}) = 1; P({H}) = 0.5; P({T}) = 0.5; P(\emptyset) = 0.$



Probability is a special function, which should satisfy some properities:

• P(S)=1; $P(\emptyset)=0$; $0 \le P(E) \le 1$;

- If event E belongs to event F, then $P(E) \leq P(F)$;
- Given an event E, then $P(E^C) = 1-P(E)$;

• Given events E and F, then $P(E \cup F) = P(E) + P(F) - P(E \cap F)$.

Random Variables

Generally, it is very complex to represent an event;

• To deal with more complex events, researchers have developed random variables (随机变量).

Example. Toss a coin (1 time). In the sample space S={ H, T}, we design a function X: S→{1,-1} such that X(H)=1 and X(T)=-1. Then X is a random variable.

Moreover, $P(X=1) = P({H}) = 0.5$ and $P(X=-1) = P({T})=0.5$.

What are Random Variables

- A random variable is a variable that can take on different values randomly. We typically denote the random variable itself with an uppercase letter in plain typeface, and the values it can take on with lowercase letters.
- For vector-valued variables, we would write the random variable as X and one of its values as x.

 Random variables may be discrete or continuous. A discrete random variable is one that has a finite or countably infinite number of states. A continuous random variable is associated with a real value.

Discrete Variables and PMF

- A probability distribution over discrete variables may be described using a probability mass function (PMF, 概率质量函数)
- The probability mass function maps from a state of a random variable to the probability of that random variable taking on that state.
- $0 \le P(X = x) \le 1$
- $\sum_{x} P(X = x) = 1$. We refer to this property as being normalized

Discrete Variables and PMF: Examples

Discrete Random Variable with finite range:

Toss a coin (1 time).

In the sample space $S=\{H, T\}$, we design a random variable $X: S \rightarrow \{1,-1\}$ such that

X(H)=1 and X(T)=-1. Then X is a random variable with finite range.

The probability is P(X=1)=P(X=-1)=0.5.



Discrete Variables and PMF: Examples

Discrete Random Variable with infinite range:

Toss a coin (countably infinity times).

We design a random variable X: X = n means that the first head appears after throwing n times.



Then X is a random variable with countably infinite range.

The probability is $P(X=n) = 0.5^{n}$.

Discrete Variables and PMF: Examples

- Discrete uniform distribution (均匀分布) is one of the most important discrete distributions
- It is a finite discrete distribution
- Assume that the range is $x_1, x_2...x_n$, then

•
$$P(X = x_i) = \frac{1}{n}$$
; $\sum_i P(X = x_i) = \sum_i \frac{1}{n} = \frac{n}{n} = 1$

Continuous Variables and PDF

A continuous variable X is a function;

Range is not discrete and take values in real number;

- There is a probability density function (概率密度函数) $p_{\rm X}({
 m x})$ such that
- 1) $p_{X}(x) \geq 0$;
- 2) $P(a \le X \le b) = \int_a^b p_X(x) dx$;
- 3) $\int_{-\infty}^{+\infty} p_{\mathbf{X}}(\mathbf{x}) d\mathbf{x} = 1.$

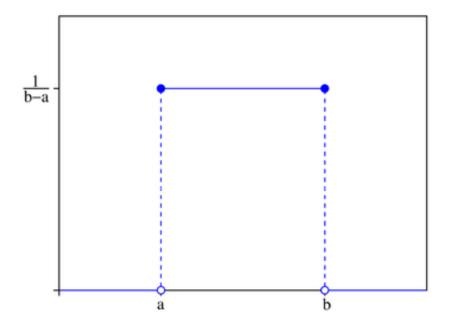
Continuous Variables and PDF

• In principle variables such as height, weight, and temperature are continuous, in practice the limitations of our measuring instruments restrict us to a discrete (though sometimes very finely subdivided) world.

 However, continuous models often approximate real-world situations very well, and continuous mathematics (calculus) is frequently easier to work with than mathematics of discrete variables and distributions.

Continuous Variables and PDF

- Continuous uniform distribution is one of the most important continuous distributions.
- The probability density function of continuous unflorm distribution can be written as $p(x; a, b) = \frac{1}{b-a}$



Joint Distribution

• In some practice case, we need to consider multiple randon variables

For example, it is clear that the weight is related to the height. So we are interested in knowing the joint distribution related to dog's weight and height.

 If random variables X, Y are discrete random variables, then the joint distribution of X and Y is

$$P(X=x, Y=y)$$

Joint Distribution

• If random variables X, Y are continuous random variables, then the joint distribution of X and Y is

$$P(a1 \le X \le b1, a2 \le Y \le b2)$$

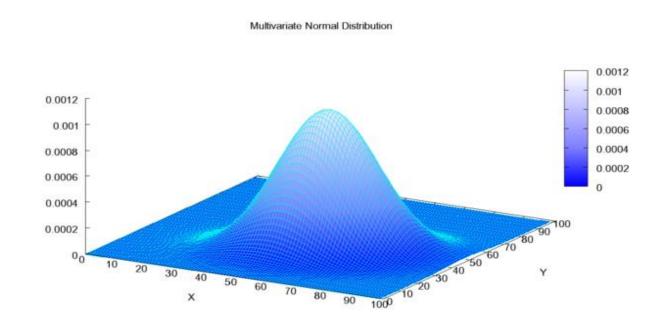
 In fact, when X, Y are continuous random variables, there exists a probability density function for the joint distribution:

P(a1
$$\leq$$
X \leq b1, a2 \leq Y \leq b2)= $\int_{a1}^{b1} \int_{a2}^{b2} p_{XY}(x,y) dxdy$,

where $p_{XY}(x,y)$ is the probability density function.

Joint Distribution

• If X represents the weight of dog and Y represents the height of dog, then the joint distribuion P(X,Y) is similar to a two-dimensional Gaussian distribution.



Joint Probability

• If random variable X is continuous random variable, and Y is discrete random variable, then the joint distribution of X and Y is

$$P(a1 \le X \le b1, Y=y)$$

In fact, when X, is continuous random variable, there exists a probability density function for the joint distribution:

$$P(a1 \le X \le b1, Y=y) = \int_{a1}^{b1} p_{XY}(x,y) dx,$$

where $p_{XY}(x,y)$ is continuous with respect to x, but discrete with respect to y.

Marginal Probability

• Using joint distribution, we can construct marginal distribution:

• If random variables X, Y are discrete random variables, then the marginal distributions are

$$P(X=x) = \sum_{y} P(X = x, Y = y)$$

$$P(Y=y) = \sum_{X} P(X = x, Y = y)$$

Marginal Probability

 If random variables X and Y are continuous random variables, then the marginal distribution with respect to X is

$$P(a \le X \le b) = \int_a^b \int_{-\infty}^{+\infty} p_{XY}(x, y) dx dy$$

and the density function of X is

$$p(x) = \int_{-\infty}^{+\infty} p_{XY}(x, y) dy$$

Similarly, we can obtain the marginal distribution with respect to Y and the density function of Y.

Marginal Probability

• If random variable X is continuous random variable, and Y is a discrete random variable, then the marginal distribution with respect to X is

$$P(a \le X \le b) = \sum_{y} \int_{a}^{b} p_{XY}(x, y) dx$$

and the density function of X is

$$p(x) = \sum_{y} p_{XY}(x, y)$$

The marginal distribution with respect to Y is

$$P(Y=y) = \int_{-\infty}^{+\infty} p_{XY}(x,y) dx$$

Independence and Conditional Independence

• If events E and F are independent, then

$$P(E|F) = P(E \cap F)/P(F) = P(E)P(F)/P(F) = P(E).$$

So

$$P(E|F) = P(E)$$

We can also obtain that

$$P(F|E) = P(F)$$

Independence and Conditional Independence

• Two random variables X and Y are independent if their probability distribution can be expressed as a product of two factors:

$$P(X, Y) = P(X)P(Y)$$

If X and Y are both discrete random variables, then X and Y are independent, if for any x, y

$$P(X=x, Y=y) = P(X=x)P(Y=y)$$

We can also obtain that P(X=x|Y=y) = P(X=x) and P(Y=y|X=x)=P(Y=y).

Independence and Conditional Independence

If X and Y are continuous random variables, then X and Y are independent, if for any x, y,

$$p_{XY}(x,y)=p_X(x)p_Y(y)$$

where $p_{XY}(x,y)$ is the density function of the joint distribution, $p_X(x)$ is the density function with respect to random variable X, and $p_Y(x)$ is the density function with respect to random variable Y.

Conditional Independence to Random Variables

Two random variables X and Y are conditionally independent given a random variable Z if the conditional probability distribution

$$P(X,Y|Z) = P(X|Z)P(Y|Z)$$

When X, Y and Z are discrete variable variables, if X and Y are conditionally independent given a random variable Z: for any x, y, z

$$P(X=x,Y=y|Z=z) = P(X=x|Z=z)P(Y=y|Z=z)$$

Conditional Independence to Random Variables

When X, Y and Z are continous variable variables, if X and Y are conditionally independent given a random variable Z: for any x, y, z, then the probability density functions satisfy that

$$p_{XY|Z}(x,y|z) = p_{X|Z}(x|z)p_{Y|Z}(y|z)$$

So

$$\frac{p_{XYZ}(x,y,z)}{p_{Z}(z)} = \frac{p_{XZ}(x,z)}{p_{Z}(z)} \frac{p_{YZ}(y,z)}{p_{Z}(z)}$$

$$p_{XYZ}(x, y, z) = p_{XZ}(x, z)p_{YZ}(y, z)$$

• Bayes' theorem is an important tool in statistics and machine learning:

Given events A and B, then
$$P(A|B) = \frac{P(B|A)P(A)}{P(B)}$$
.

Proof:

$$P(A|B) = P(A \cap B)/P(B) = (P(B|A)P(A))/P(B).$$

How to use Baye's Rule?

Example: A bag I contains 4 white and 6 black balls while another bag II contains 4 white and 3 black balls.

One ball is drawn at random from one of the bags, and it is found to be black.

Find the probability that it was drawn from bag I.

Soultion: Let E1 be the event of choosing bag I, E2 the event of choosing bag II, and A be the event of drawing a black ball.

```
Then, P(E1) = P(E2) = 0.5.

P(A|E1) = P(drawing a black ball from Bag I) = 6/10 = 3/5.

P(A|E2) = P(drawing a black ball from Bag II) = 3/7

P(A) = P(A|E1)P(E1) + P(A|E2)P(E2) = 18/35
```

By using Bayes' theorem, the probability of drawing a black ball from bag I out of two bags,

$$P(E1|A) = P(A|E1)P(E1)/P(A)=0.6*0.5/(18/35) = 7/12$$

How to use Baye's Rule?

Example: A man is known to speak the truth 2 out of 3 times. He throws a dice and reports that the number obtained is a four. Find the probability that the number obtained is four.



Soultion: Let A be the event that the man reports that number four is obtained. Let E1 be the event that four is obtained and E2 be its complementary event.

Then, P(E1) = Probability that four occurs = 1/6.

P(E2) = Probability that four does not occur = 1- P(E1) = 1 - (1/6) = 5/6.

P(A|E1)= Probability that man reports four and it is actually a four = 2/3

 $P(A \mid E2)$ = Probability that man reports four and it is not a four = 1/3.

So P(A) = P(A|E1)P(E1)+P(A|E2)P(E2) = 1/9+5/18 = 7/18.

By using Bayes' theorem, probability that number obtained is actually a four, P(E1|A) = P(A|E1)P(E1)/P(A) = 2/18/(7/18) = 2/7.

Expectation (or Mean)

The expectation or mean of a random variable X is denoted by E[X] and defined as:

For discrete variable,

$$E[X] = \sum_{x} P(X = x)x$$

For continuous variable,

$$E[X] = \int_{-\infty}^{+\infty} p_X(x) x \, dx$$

In words, we are taking a weighted sum of the values that x can take on, where the weights are the probabilities of those respective values. The expected value has a physical interpretation as the "center of mass" of the distribution.

Expectation of Functions

The expectation or mean of f(X) (a function of random variable X) is denoted by E[f(X)] and defined as:

For discrete variable,

$$E[f(X)] = \sum_{x} P(X = x)f(x)$$

For continuous variable,

$$E[f(X)] = \int_{-\infty}^{+\infty} p_X(x) f(x) dx$$

Properties of Expectation

For any two random variables X and Y, functions f and g, and any constants a, $b \in R$, the following equations hold:

$$E[f(a)] = f(a)$$

•
$$E[X+Y] = E[X]+E[Y]$$

$$E[f(X)+g(Y)] = E[f(X)]+E[g(X)]$$

•
$$E[aX] = aE[X]$$

$$E[af(X)] = aE[f(X)]$$

$$E[af(X)+bg(Y)]=aE[g(X)]+bE[g(Y)]$$

Variance

• Expectation provides measure of the "center" of a distribution, but sometimes we are also interested in what the "spread" is about that center. Therefore, we define the variance Var(X) of a random variable X as follows:

$$Var(X) = E[(X - E[X])^2]$$

• In words, this is the average squared deviation of the values of X from the mean of X.

Properties of Variance

For any random variable X and any a, $b \in R$, the following equations hold:

$$Var(aX+b) = a^2Var(X)$$

$$Var(X) = E[X^2] - E[X]^2$$

Covariance

 The covariance gives some sense of how much two values are linearly related to each other, as well as the scale of these variables

The covariance of two random variables X and Y is denoted by Cov(X, Y) and defined as

$$Cov(X, Y) = E[(X - E[X])(Y - E[Y])].$$

Covariance of functions: given functions f and g, then

$$Cov(f(X), g(Y)) = E[(f(X) - E[f(X)])(g(Y) - E[g(Y)])].$$

Properties of Covariance

It is clear that

$$Cov(X,X) = Var(X);$$

$$Cov(f(X),f(X)) = Var(f(X)).$$

Properties of Covariance

If two random variables X and Y are independent, then

$$Cov(X, Y) = 0$$
.

Proof: From the definition of covariance, we have

$$Cov(X, Y) = E[(X - E[X])(Y - E[Y])]$$

- = E[XY E[X]Y XE[Y] + E[X]E[Y]]
- = E[XY] E[E[X]Y] E[XE[Y]] + E[E[X]E[Y]]
- = E[XY] E[X]E[Y] E[X]E[Y] + E[X]E[Y] = E[XY] E[X]E[Y] = 0.

Common Probability Distribution

Distribution	PDF or PMF	Mean	Variance
Bernoulli(p)	$\begin{cases} p, & \text{if } x = 1 \\ 1 - p, & \text{if } x = 0. \end{cases}$	р	p(1-p)
Binomial(n, p)	$\binom{n}{k} p^k (1-p)^{n-k}$ for $k = 0, 1,, n$	np	np(1-p)
Geometric(p)	$p(1-p)^{k-1}$ for $k=1,2,$	$\frac{1}{p}$	$\frac{1-p}{p^2}$
$Poisson(\lambda)$	$\frac{e^{-\lambda}\lambda^k}{k!}$ for $k=0,1,$	λ	λ
Uniform(a, b)	$\frac{1}{b-a}$ for all $x \in (a,b)$	<u>a+b</u> 2	$\frac{(b-a)^2}{12}$
Gaussian (μ, σ^2)	$\frac{1}{\sigma\sqrt{2\pi}}e^{-\frac{(x-\mu)^2}{2\sigma^2}} \text{ for all } x \in (-\infty, \infty)$	μ	σ^2
Exponential(λ)	$\lambda e^{-\lambda x}$ for all $x \ge 0, \lambda \ge 0$	$\frac{1}{\lambda}$	$\frac{1}{\lambda^2}$

Gaussian Distribution (Normal Distribution)

• Gaussian Distribution (Normal Distribution): It is the most widely used model for the distribution of continuous variables. For a single variable X, the Gaussian distribution can be represented as follows:

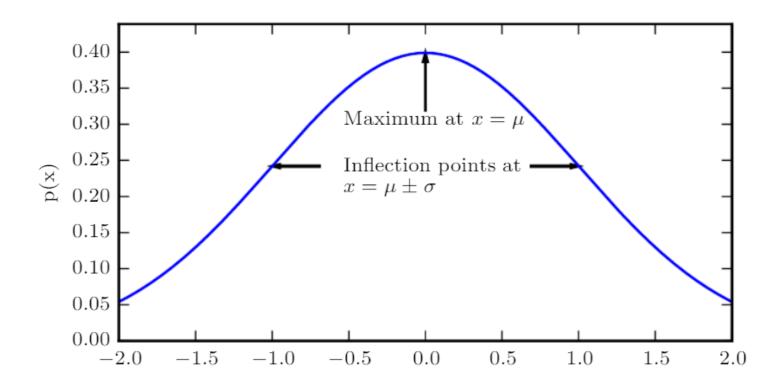
$$N(\mu, \sigma^2)$$
.

• It is a continuous distribution with the probability density function

$$\sqrt{\frac{1}{2\pi\sigma^2}}\exp(-\frac{1}{2\sigma^2}(x-\mu)^2)$$

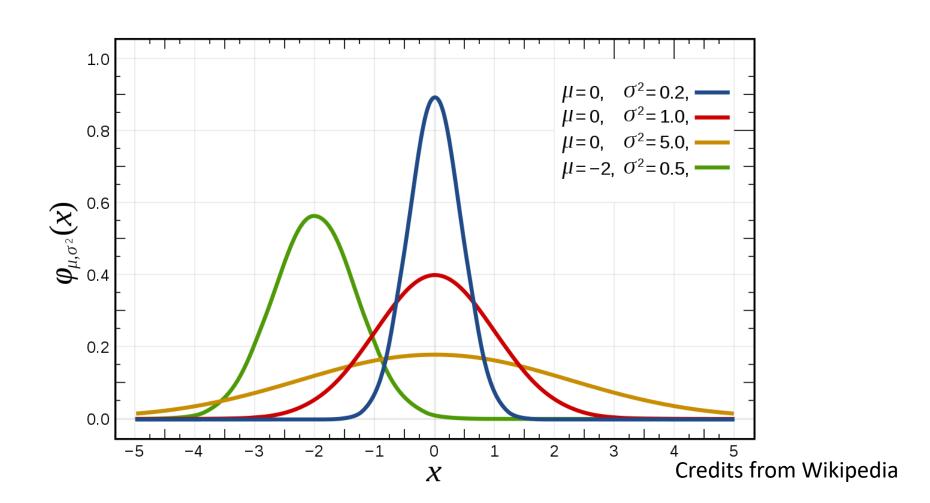
Gaussian Distribution (Normal Distribution)

Gaussian distribution depends on two parameters μ , σ . Then the expectation is μ and the variance is σ^2 .



Gaussian Distribution (Normal Distribution)

Gaussian distributions with different expectations and variances



Consider following table of counts that are obtained from an observed sample of individuals both males and females, who had taken Covid-19. X, Y, Z are random variables. Z represents the gender, X represents whether individuals recover from Covid-19, and Y represents whether individuals have been treated.



Z=1 Male; Z=0 Female; X=1 Recovery (康复); X=0 No Recovery; Y=1 Treatment (治疗); Y=0 No Treatment;

	Z=1, X=1	Z=1, X=0	Z=0, X=1	Z=0, X=0
Y=1	245	105	315	735
Y=0	630	420	70	280

245, 105, 315, 735, 630,420,70, 280 represent the number of individuals corresponding to different values of X, Y and Z.

Please compute that



	Z=1, X=1	Z=1, X=0	Z=0, X=1	Z=0, X=0
Y=1	245	105	315	735
Y=0	630	420	70	280

Solution:

$$P(X=1|Y=1,Z=1)=P(X=1,Y=1,Z=1)/P(Y=1,Z=1)=P(X=1,Y=1,Z=1)/(P(X=1,Y=1,Z=1)+P(X=0,Y=1,Z=1))=245/(245+105)=0.7$$

$$P(X=1|Y=0,Z=1)=P(X=1,Y=0,Z=1)/P(Y=0,Z=1)=P(X=1,Y=0,Z=1)/(P(X=1,Y=0,Z=1)+P(X=0,Y=0,Z=1))=630/(630+420)=0.6$$

P(X=1|Y=1,Z=1) means the the recovery probability for individuals who are male and have been treated.

P(X=1|Y=0,Z=1) means the the recovery probability for individuals who are male and have not been treated.

Solution:

$$P(X=1|Y=1,Z=0)=P(X=1,Y=1,Z=0)/P(Y=1,Z=0)=P(X=1,Y=1,Z=0)/(P(X=1,Y=1,Z=0)+P(X=0,Y=1,Z=0))=315/(315+735)=0.3$$

$$P(X=1|Y=0,Z=0)=P(X=1,Y=0,Z=0)/P(Y=0,Z=0)=P(X=1,Y=0,Z=0)/(P(X=1,Y=0,Z=0)+P(X=0,Y=0,Z=0))=70/(70+280)=0.2$$

P(X=1|Y=1,Z=0) means the the recovery probability for individuals who are female and have been treated.

P(X=1|Y=0,Z=0) means the the recovery probability for individuals who are female and have not been treated.

Solution:

$$P(X=1|Y=1)=P(X=1,Y=1)/P(Y=1)=(245+315)/(245+105+315+735)=0.4$$

$$P(X=1|Y=0)=P(X=1,Y=0)/P(Y=0)=(630+70)/(630+420+70+280)=0.5$$

P(X=1|Y=1) means the the recovery probability for individuals who have been treated.

P(X=1|Y=0) means the the recovery probability for individuals who have not been treated.

$$P(X=1|Y=1,Z=1)=0.7>P(X=1|Y=0,Z=1)=0.6,$$

 $P(X=1|Y=1,Z=0)=0.3>P(X=1|Y=0,Z=0)=0.2.$



- The recovery probability for individuals who are male and have been treated is larger than the recovery probability for individuals who are male and have not been treated.
- The recovery probability for individuals who are female and have been treated is larger than the recovery probability for individuals who are female and have not been treated.

- The recovery probability for individuals who are male and have been treated is larger than the recovery probability for individuals who are male and have not been treated.
- The recovery probability for individuals who are female and have been treated is larger than the recovery probability for individuals who are female and have not been treated.
- Does this mean that the treatment can make postive affects? That is: is the recovery probability for individuals who have been treated larger than the recovery probability for individuals who have not been treated?

To answer this question, we need to compute P(X=1|Y=1) and P(X=1|Y=0)

Then, we need to compare them.

We discover that P(X=1|Y=1)=0.4 < P(X=1|Y=0)=0.5.

 That is: the recovery probability for individuals who have been treated is smaller than the recovery probability for individuals who have not been treated.

It seems that

Why does this happen?

It is called Simpson's paradox.

- The recovery probability for individuals who are male and have been treated is larger than the recovery probability for individuals who are male and have not been treated; and
- the recovery probability for individuals who are female and have been treated is larger than the recovery probability for individuals who are female and have not been treated.

But

 the recovery probability for individuals who have been treated is smaller than the recovery probability for individuals who have not been treated

Note that

$$P(Z=1|Y=1)P(X=1|Y=1,Z=1)+P(Z=0|Y=1)P(X=1|Y=1,Z=0)=P(X=1|Y=1) \\ P(Z=1|Y=0)P(X=1|Y=0,Z=1)+P(Z=0|Y=0)P(X=1|Y=0,Z=0)=P(X=1|Y=0) \\ P(Z=1|Y=0)P(X=1|Y=0,Z=1)+P(Z=0|Y=0)P(X=1|Y=0,Z=0)=P(X=1|Y=0) \\ P(Z=1|Y=0)P(X=1|Y=0,Z=1)+P(Z=0|Y=0)P(X=1|Y=0,Z=0)=P(X=1|Y=0) \\ P(Z=1|Y=0)P(X=1|Y=0,Z=0)=P(X=1|Y=0)P(X=1|Y=0,Z=0)=P(X=1|Y=0) \\ P(Z=1|Y=0)P(X=1|Y=0,Z=0)=P(X=1|Y=0)P(X=1|Y=0,Z=0)=P(X=1|Y=0) \\ P(Z=1|Y=0)P(X=1|Y=0,Z=0)=P(X=1|Y=0)P(X=1|Y=0,Z=0)=P(X=1|Y=0) \\ P(Z=1|Y=0)P(X=1|Y=0,Z=0)=P(X=1|Y=0)P(X=1|Y=0,Z=0)=P(X=1|Y=0) \\ P(Z=1|Y=0,Z=0)=P(X=1|Y=0)P(X=1|Y=0,Z=0)=P(X=1|Y=0)P(X=1|Y=0) \\ P(Z=1|Y=0,Z=0)=P(X=1|Y=0)P(X=1|Y=0,Z=0)=P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|Y=0)P(X=1|$$

Although P(X=1|Y=1,Z=1)>P(X=1|Y=0,Z=1), P(X=1|Y=1,Z=0)>P(X=1|Y=0,Z=0), the conditional probabilities P(Z=1|Y=1), P(Z=0|Y=1), P(Z=1|Y=0) and P(Z=0|Y=0) can affect the values of P(X=1|Y=1) and P(X=1|Y=0).

That is the basic reason why the Simpson's paradox happens.

Detailly,

let
$$u = P(Z=1|Y=1)$$
 and $v = P(Z=1|Y=0)$, then if we hope that $P(X=1|Y=1) > P(X=1|Y=0)$

We need the following inequality:

$$0.7u + 0.3(1 - u) > 0.6v + 0.2(1 - v)$$
.

Whether the inequality can success depends on the values of u and v.

$$0.7u+0.3(1-u)>0.6v+0.2(1-v)$$
 if and only if $v-u<0.25$.

But, u=P(Z=1|Y=1)=0.25 and v=P(Z=1|Y=0)=0.75. So v-u=0.5>0.25. That is the reason why P(X=1|Y=1)< P(X=1|Y=0).

The selected distribution $P(X;\alpha^*)$ is the most possible distribution sampling data S=(x1,x2,...,xn), i.i.d..

Understanding above sentence, we can formulate it as follows:

$$\underset{\alpha \in \Delta}{\operatorname{argmax}} P(x1, x2, \dots, xn; \alpha)$$

here we assume P(X; a) is a discrete distribution.

• $\max_{\alpha \in \Delta} P(x1, x2, ..., xn; \alpha)$

means the largest probability for $P(X; \alpha)$ that S is observed.

Because (x1,...,xn), are Independent and identically distributed,

$$\underset{\alpha \in \Delta}{\operatorname{argmax}} P(x1, x2, \dots, xn; \alpha)$$

is equal to

$$\underset{\alpha \in \Delta}{\operatorname{argmax}} \prod_{i=1}^{n} P(X = xi; \alpha)$$

To reduce the affects of Multiplier operator $\prod_{i=1}^n$, we take a small trick (we use the property of log function to help us):

$$\log \prod_{i=1}^{n} ai = \sum_{i=1}^{n} \log ai$$

Step 1. We take log function.

$$\underset{\alpha \in \Delta}{\operatorname{argmax}} \prod_{i=1}^{n} P(X = xi; \alpha) = \underset{\alpha \in \Delta}{\operatorname{argmax}} \log \prod_{i=1}^{n} P(X = xi; \alpha)$$

$$\underset{\alpha \in \Delta}{\operatorname{argmax}} \prod_{i=1}^{n} P(X = xi; \alpha) = \underset{\alpha \in \Delta}{\operatorname{argmax}} \log \prod_{i=1}^{n} P(X = xi; \alpha)$$

Step 2. Using the property of log function:

$$\log \prod_{i=1}^{n} P(X = xi; a) = \sum_{i=1}^{n} \log P(X = xi; a)$$

Therefore,

$$\underset{\alpha \in \Delta}{\operatorname{argmax}} \prod_{i=1}^{n} P(X = xi; \alpha) = \underset{\alpha \in \Delta}{\operatorname{argmax}} \sum_{i=1}^{n} \log P(X = xi; \alpha)$$

Step 3. We need to optimize

$$\underset{\alpha \in \Delta}{\operatorname{argmax}} \sum_{i=1}^{n} \log P(X = xi; \alpha) \quad (1)$$

and obtain the optimal solution.

The solution of Eq. 1 is called Maximum Likelihood Estimation.

If the distribution class consists of continuous distributions, that is $P(X;\alpha)$ is a continuous distribution with respect to all $\alpha \in \Delta$.

Then the Maximum Likelihood Estimation is

$$\underset{\alpha \in \Delta}{\operatorname{argmax}} \sum_{i=1}^{n} \log p_X(xi; \alpha) \quad (2)$$

where $p_X(x; a)$ is the density function of P(X; a).

How to obtain the solution of

$$\underset{\alpha \in \Delta}{\operatorname{argmax}} \sum_{i=1}^{n} \log P(X = xi; \alpha) ?$$

- This is related to optimization problem.
- Generally, there are no unviersal approaches to give soultions to all Maximum Likelihood (ML) Estimation.
- The approaches are case by case.

In this class, we introduce a common used approach.

This approach is based on a simple theorem:

- If 1) a function f(x1,x2,...,xd) is differentiable,
 - 2) $x^* = (x1^*, x2^*, ..., xd^*)$ is the maximum point of f, then

$$\frac{\partial f}{\partial x_i}(x1^*, x2^*, \dots, xd^*) = 0.$$

Using this theorem, if $\sum_{i=1}^{n} log P(X = xi; a)$ is differentiable, then

Let a = (a1, a2, ..., ad),

$$\frac{\partial \sum_{i=1}^{n} \log P(X=xi;a)}{\partial aj} = 0, \text{ for } j=1,...,d$$

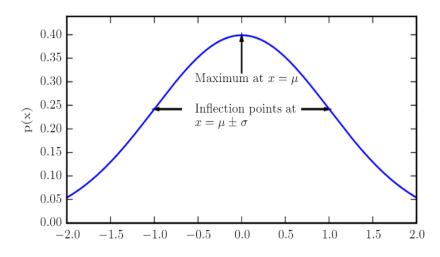
Then, addressing above equations.

Check that you've found a maximum rather than a minimum or saddle-point, and be careful if α belongs to Δ .

Suppose you have x1,x2,...,xn (i.i.d) $N(\mu,\sigma^2)$

$$\sqrt{\frac{1}{2\pi\sigma^2}}\exp(-\frac{1}{2\sigma^2}(x-\mu)^2)$$

- Assume that you know σ^2
- But you don't know μ



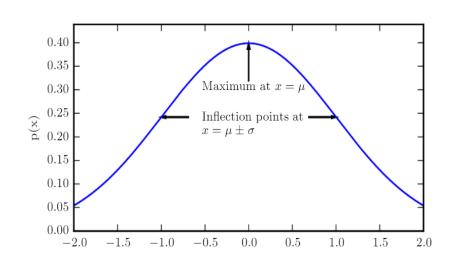
MLE: For which μ is x1, x2, ..., xn most likely?

Compute the MLE $\underset{\mu \in R}{\operatorname{argmax}} \sum_{i=1}^{n} \log p_X(xi; \mu)$

$$\underset{\mu \in R}{\operatorname{arg \, max}} \frac{1}{\sqrt{2\pi}} \sum_{\sigma=1}^{n} -\frac{(xi - \mu)^2}{2\sigma^2}$$

$$= \arg\min \sum_{\mu \in R}^{n} (xi - \mu)^{2}$$

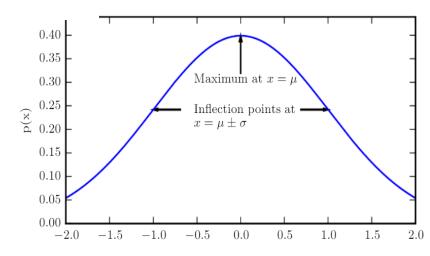
$$\mu \in R \quad i=1$$



Noticed that σ is known, we can ignore it when finding the proper μ

Derivation the equation $\underset{\mu \in R}{\arg\min} \sum_{i=1}^{n} (xi - \mu)^2$

$$\frac{d\sum_{i=1}^{n} (xi - \mu)^2}{d\mu} = 2\sum_{i=1}^{n} (xi - \mu) = 0$$

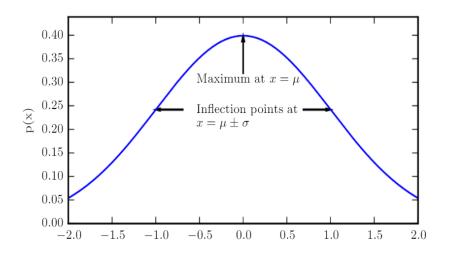


So the solution is

$$\mu = \frac{\sum_{i=1}^{n} xi}{n}$$

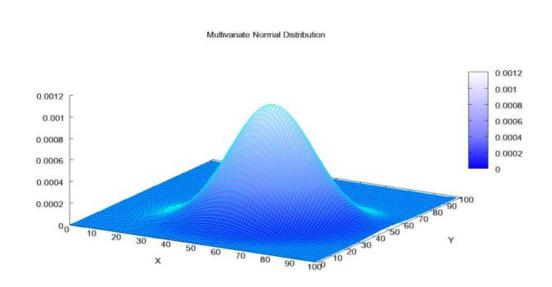
•In conclusion, the best estimate of the mean of a gaussian distribution is the mean of the sample!

$$\mu = \frac{\sum_{i=1}^{n} xi}{n}$$



Exercises: MLE for high-dimensional Gaussian Distribution

• Given a 2 \times 2 positive semi-definite matrix (半正定) Σ and a 2 \times 1 vector μ , a three dimensional normal distribution $N(\mu, \Sigma)$ can be represented as follows: the density function of this distribution is



$$p_{XY}(x, y; \mu, \Sigma) = \frac{1}{\sqrt{(2\pi)^2 |\Sigma|}} e^{-\frac{1}{2}(x-\mu)^T \Sigma^{-1}(x-\mu)}$$

where $|\Sigma|$ is the determinant (行列式) of Σ and $\mathbf{x} = (x, y)^{\mathrm{T}}$.

Definition: Positive Semi-Definite Matrix

- Definite Matrix:
 - An $n \times n$ symmetric real matrix A is said to be positive definite if $x^T A x > 0$ for all non-zero $x \in \mathbb{R}^n$
 - An $n \times n$ symmetric real matrix A is said to be negative definite if $x^T A x < 0$ for all non-zero $x \in \mathbb{R}^n$
- Semi-Definite Matrix:
 - An $n \times n$ symmetric real matrix A is said to be positive semi-definite if $x^T A x \ge 0$ for all non-zero $x \in \mathbb{R}^n$
 - An $n \times n$ symmetric real matrix A is said to be negative semi-definite if $x^T A x \leq 0$ for all non-zero $x \in \mathbb{R}^n$

Useful Properties

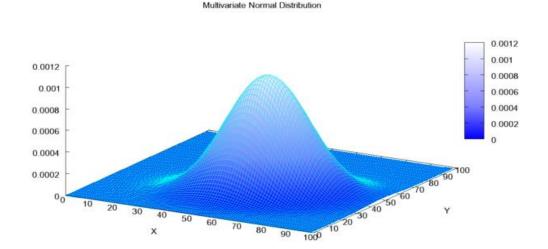
- For Positive Definite Matrix, all the eigenvalues (特征根) $\lambda > 0$.
- For Negative Definite Matrix, all the eigenvalues $\lambda < 0$.
- For Positive Semi-Definite Matrix, all the eigenvalues $\lambda \geq 0$.
- For Negative Semi-Definite Matrix, all the eigenvalues $\lambda \leq 0$.

 These properties are useful to check whether the matrix is definite, semi-definite or not.

• If μ =(a1,a2) and Σ is a diagonal matrix with eigenvalues λ 1, λ 2 (λ 1 > 0, λ 2 > 0),

$$\Sigma = \begin{bmatrix} \lambda 1 & 0 \\ 0 & \lambda 2 \end{bmatrix}$$

then f(x, y) can be writeen as:



$$p_{XY}(x,y;\mu,\Sigma) = \frac{1}{\sqrt{(2\pi)^2 \lambda 1 * \lambda 2}} e^{-\frac{1}{2\lambda 1}(x-a_1)^2 - \frac{1}{2\lambda 2}(y-a_2)^2}$$

Two Dimension Gaussian Distribution

- How to get $p_{XY}(x,y;\mu,\Sigma) = \frac{1}{\sqrt{(2\pi)^2\lambda 1*\lambda 2}}e^{-\frac{1}{2\lambda 1}(x-a_1)^2-\frac{1}{2\lambda 2}(y-a_2)^2}$?
- Consider one dimension version:

•
$$p_X(x; \mu, \sigma) = \frac{1}{\sqrt{2\pi\lambda_1}} e^{-\frac{1}{2\lambda_1}(x-a_1)^2}$$
 (we have $\sigma = \lambda 1, \mu = a_1$)

- $\Sigma = \begin{bmatrix} \lambda 1 & 0 \\ 0 & \lambda 2 \end{bmatrix}$ tells us X and Y are independent.
- Then we have:

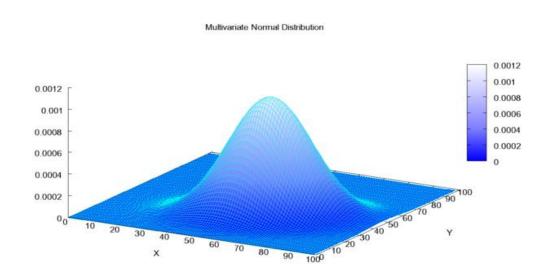
•
$$p_{XY}(x, y; \mu, \Sigma) = \frac{1}{\sqrt{2\pi\lambda 1}} e^{-\frac{1}{2\lambda 1}(x-a_1)^2} \times \frac{1}{\sqrt{2\pi\lambda 2}} e^{-\frac{1}{2\lambda 2}(x-a_2)^2}$$

= $\frac{1}{\sqrt{(2\pi)^2\lambda 1*\lambda 2}} e^{-\frac{1}{2\lambda 1}(x-a_1)^2 - \frac{1}{2\lambda 2}(y-a_2)^2}$

• If

$$\Sigma = \begin{bmatrix} \lambda 1 & 0 \\ 0 & \lambda 2 \end{bmatrix}$$

and we have n data (x1,y1),...,(xn,yn) sampled from a two-dimensional Gaussian Distribution $N(\mu, \Sigma)$, i.i.d., calculate μ by the maximum likelihood estimation method.



Maximum Likelihood (ML) Estimation:

$$\underset{\mu}{\operatorname{argmax}} \sum_{i=1}^{n} \log p_{XY}(xi, yi; \mu, \Sigma)$$

It is equal to

$$\underset{\text{a1,a2}}{\operatorname{argmin}} \sum_{i=1}^{n} \left(\frac{1}{2\lambda 1} (xi - a1)^2 + \frac{1}{2\lambda 2} (yi - a2)^2 \right)$$

Derivation the equation $G(a1,a2) = \sum_{i=1}^{n} (\frac{1}{2\lambda 1}(xi-a1)^2 + \frac{1}{2\lambda 2}(yi-a2)^2)$

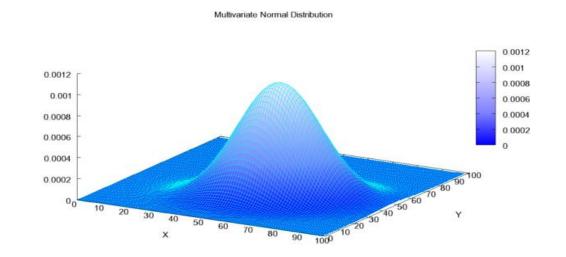
$$\frac{\partial G}{\partial a1} = \sum_{i=1}^{n} \frac{a1 - xi}{\lambda 1} = 0, \qquad \frac{\partial G}{\partial a2} = \sum_{i=1}^{n} \frac{a2 - yi}{\lambda 2} = 0$$

So
$$a1 = \frac{1}{n} \sum_{i=1}^{n} xi$$
, $a2 = \frac{1}{n} \sum_{i=1}^{n} yi$

(As we only need to estimate μ , so the λ will not influence the estimation.)

•In conclusion, the best estimate of the mean of a two-dimensional gaussian distribution is the mean of the sample!

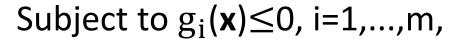
In fact, it also holds for highdimensional gaussian distribution.



Convex Optimization

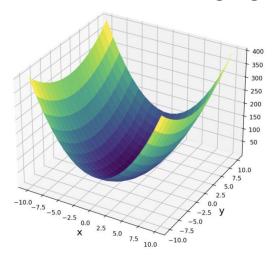
Definition of Convex Optimization Problem:

Minimize f(x)



$$h_j(x) = 0, j = 1,...,n.$$

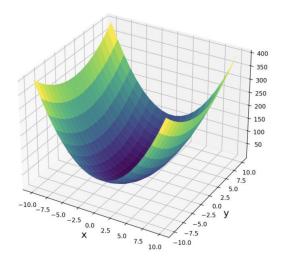
- $g_i(x)$ is convex function, i=1,...,m
- $h_j(x)$ is linear function A_jx+b_j , j=1,...,n
- f(x) is a convex function



Convex Optimization

Examples of Convex Optimization:

Linear optimization belongs to Convex Optimization.
 Because linear functions are also convex function.



• Minimize $f(x) = x^T Mx + bx$ is a Convex Optimization problem without constraints, if M is a positive semi-definite matrix.

We first introduce how to address convex optimization without constraints: that is

Minimize f(x), where f(x) is convex

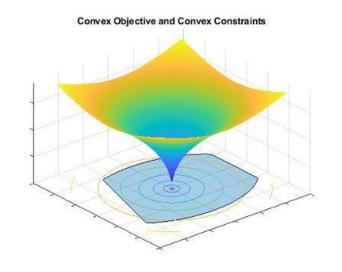
We want to ask some issues:

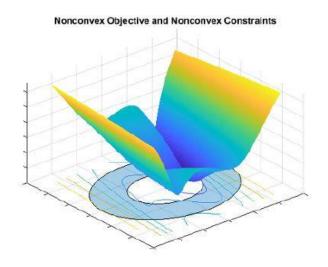
Issue 1. Whether we can find a solution to this issue?

Issue 2. Whether the solution is unique.

Exercise:
$$\mathbf{M} = \begin{pmatrix} 2 & -1 & 0 \\ -1 & 2 & -1 \\ 0 & -1 & 2 \end{pmatrix}, \mathbf{b} = (1,0,0)$$

what is the solution of Minimize $f(x) = x^T Mx + bx$?





Solution: Firstly, we note that **M** is an inverse matrix.

What is the inverse matrix of M?

The inverse matrix \mathbf{M}^{-1} of \mathbf{M} satisfies that

Inverse Inverse

$$\mathbf{M}^{-1}\mathbf{M} = \mathbf{I}$$
, where \mathbf{I} is the identity matrix.

How to compute the inverse matrix?

In matlab, we can use code inv(M) to compute the inverse matrix of M.

To understand how to compute the inverse matrix, we need to know the determinant (行列式) of M.

$$\det(\mathbf{M}) = \begin{vmatrix} m_{11} m_{12} \dots m_{1d} \\ m_{21} m_{22} \dots m_{2d} \\ \vdots & \vdots & \vdots \\ m_{d1} m_{d2} \dots m_{dd} \end{vmatrix} = \sum_{i=1}^{d} (-1)^{i+j} m_{ij} \det(\mathbf{M}_{ij})$$

where M_{ij} is the (d-1)×(d-1) submatrix obtained by deleting row i and column j from M.

$$\det(\mathbf{M}) = \begin{vmatrix} 2 & -1 & 0 \\ -1 & 2 & -1 \\ 0 & -1 & 2 \end{vmatrix} = \sum_{i=1}^{d} (-1)^{i+j} m_{ij} \det(\mathbf{M}_{ij})$$

$$= 2 \begin{vmatrix} 2 & -1 \\ -1 & 2 \end{vmatrix} + 1 \begin{vmatrix} -1 & -1 \\ 0 & 2 \end{vmatrix} + 0 \begin{vmatrix} -1 & 2 \\ 0 & -1 \end{vmatrix} + 1 \begin{vmatrix} -1 & 0 \\ -1 & 2 \end{vmatrix}$$

$$+ 2 \begin{vmatrix} 2 & 0 \\ 0 & 2 \end{vmatrix} + 1 \begin{vmatrix} 2 & -1 \\ 0 & -1 \end{vmatrix} + 0 \begin{vmatrix} -1 & 0 \\ 2 & -1 \end{vmatrix} + 1 \begin{vmatrix} 2 & 0 \\ 2 & -1 \end{vmatrix} + 2 \begin{vmatrix} 2 & -1 \\ -1 & 2 \end{vmatrix}$$

Now we need to compute the determinant of 2×2 matrix using folloing equations

$$\begin{vmatrix} a & b \\ c & d \end{vmatrix} = ad - bc$$

$$\begin{vmatrix} 2 & -1 \\ -1 & 2 \end{vmatrix} = 2*2 - (-1)* (-1) = 3, \begin{vmatrix} -1 & -1 \\ 0 & 2 \end{vmatrix} = (-1)*2 - 0*(-1) = -2$$

We omit the computing process of other 7 matrics.

$$\det(\mathbf{M}) = \begin{vmatrix} 2 & -1 & 0 \\ -1 & 2 & -1 \\ 0 & -1 & 2 \end{vmatrix} = 4$$

Before computing the inverse matrix, we need to know how to compute the adjoint matrix (伴随矩阵).

What is adjoint matrix?

$$adj(M) = \begin{pmatrix} A_{11}A_{12} \cdots A_{1d} \\ A_{21}A_{22} \cdots A_{2d} \\ \vdots & \ddots & \vdots \\ A_{d1}A_{d2} \cdots A_{dd} \end{pmatrix}^{T}$$

where $A_{ij} = (-1)^{i+j}$

 $\mathbf{M}_{ij,}$

here M_{ij} is the (d-1)×(d-1) submatrix obtained by deleting row i and column j from M.

Because
$$\mathbf{M} = \begin{pmatrix} 2 & -1 & 0 \\ -1 & 2 & -1 \\ 0 & -1 & 2 \end{pmatrix}$$
 adj $(\mathbf{M}) = \begin{pmatrix} A_{11} & A_{12} & A_{13} \\ A_{21} & A_{22} & A_{23} \\ A_{31} & A_{32} & A_{33} \end{pmatrix}^{T}$

$$A_{11} = (-1)^2 \begin{vmatrix} 2 & -1 \\ -1 & 2 \end{vmatrix} = 3, \qquad A_{12} = (-1)^3 \begin{vmatrix} -1 & -1 \\ 0 & 2 \end{vmatrix} = 2, A_{13} = (-1)^4 \begin{vmatrix} -1 & 2 \\ 0 & -1 \end{vmatrix} = 1,$$

$$A_{21} = (-1)^3 \begin{vmatrix} -1 & 0 \\ -1 & 2 \end{vmatrix} = 2, \qquad A_{22} = (-1)^4 \begin{vmatrix} 2 & 0 \\ 0 & 2 \end{vmatrix} = 4, A_{23} = (-1)^5 \begin{vmatrix} 2 & -1 \\ 0 & -1 \end{vmatrix} = 2,$$

$$A_{31} = (-1)^4 \begin{vmatrix} -1 & 0 \\ 2 & -1 \end{vmatrix} = 1, \qquad A_{32} = (-1)^5 \begin{vmatrix} 2 & 0 \\ -1 & -1 \end{vmatrix} = 2, A_{33} = (-1)^6 \begin{vmatrix} 2 & -1 \\ -1 & 2 \end{vmatrix} = 3$$

Then, the adjoint matrix is
$$adj(M) = \begin{pmatrix} 3 & 2 & 1 \\ 2 & 4 & 2 \\ 1 & 2 & 3 \end{pmatrix}$$

Finally, the inverse matrix is

$$M^{-1} = adj(M)/det(M)$$

Therefore, the inverse matrix

$$\mathbf{M}^{-1} = \begin{pmatrix} 0.75 & 0.5 & 0.25 \\ 0.5 & 1 & 0.5 \\ 0.25 & 0.5 & 0.75 \end{pmatrix}$$

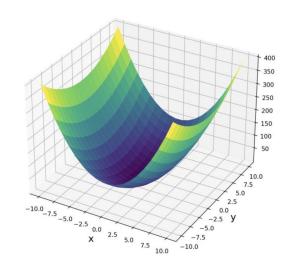
More details about the inverse matrix can be found in *Lecture 1: Introduction to Linear Algebra: Vectors and Matrices* and *Lecture 2: Linear Independence, Rank, and Orthogonality*

Solution: Firstly, we note that **M** is an inverse matrix, and the inverse is

$$\mathbf{M}^{-1} = \begin{pmatrix} 0.75 & 0.5 & 0.25 \\ 0.5 & 1 & 0.5 \\ 0.25 & 0.5 & 0.75 \end{pmatrix}$$

Secondly, we need to compute the gradient $\nabla f(\mathbf{x})$

How to compute the gradient?

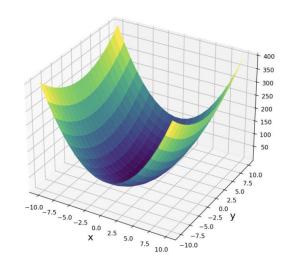


How to compute the gredient?

We introduce the Matrix derivatives:

$$\frac{\partial \mathbf{x}^{\mathrm{T}} \mathbf{M} \mathbf{x}}{\partial \mathbf{x}} = (\mathbf{M} + \mathbf{M}^{\mathrm{T}}) \mathbf{x},$$

$$\frac{\partial \mathbf{b} \mathbf{x}}{\partial \mathbf{x}} = \mathbf{b}^{\mathrm{T}},$$

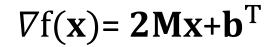


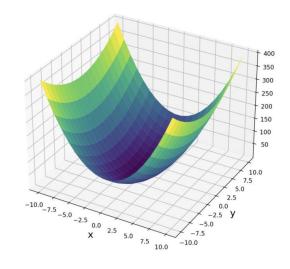
More details can be found in https://cloud.tencent.com/developer/article/1551901

Because **M** is sysmetric, then

$$\frac{\partial \mathbf{x}^{\mathrm{T}} \mathbf{M} \mathbf{x}}{\partial \mathbf{x}} = (\mathbf{M} + \mathbf{M}^{\mathrm{T}}) \mathbf{x} = 2 \mathbf{M} \mathbf{x}.$$

So we obtain that





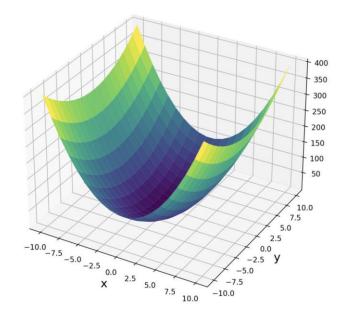
Then according to Theorem 4, the optimal solution sholud satisfy that $2Mx_0+b^T=0$

Then according to Theorem 4, the optimal solution sholud satisfy that

$$2\mathbf{M}\mathbf{x_0} + \mathbf{b}^{\mathrm{T}} = \mathbf{0}$$

The solution of
$$2Mx_0 + b^T = 0$$
 is $\frac{-M^{-1}b^T}{2} = \frac{-1}{2} \begin{pmatrix} 0.75 & 0.5 & 0.25 \\ 0.5 & 1 & 0.5 \\ 0.25 & 0.5 & 0.75 \end{pmatrix} \begin{pmatrix} 1 \\ 0 \\ 0 \end{pmatrix}$

So the optimal solution is $\begin{pmatrix} -3/8 \\ -1/4 \\ -1/\Omega \end{pmatrix}$



We next introduce convex optimization with constraints: that is

Minimize
$$f(x)$$
,

Subject to
$$g_i(\mathbf{x}) \leq 0$$
, $i=1,...,m$,

$$h_j(x) = 0, j = 1,...,n.$$

We want to ask an issue:

Issue. Whether we can find a solution to this issue? (解的存在性)

Issue. Whether we can find a solution to this issue? (解的存在性)

Following theorem gives the answer:

Theorem 8. Assume that $f(\mathbf{x})$ is differential, then \mathbf{x}_0 is the optimal solution of the Convex optimization problem with constraints if and only if

$$\nabla f(\mathbf{x}_0)^{\mathrm{T}}(\mathbf{y} - \mathbf{x}_0) \ge 0$$
,

for all **y** satisfy the constraints.

The proof can be found in Section 4.2.3 in https://web.stanford.edu/~boyd/cvxbook/bv_cvxbook.pdf

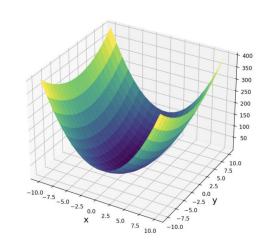
We consider a conic form problem

Minimize **bx**

Subject to $\mathbf{x}^{\mathsf{T}}\mathbf{M}\mathbf{x} - c \leq 0$,

where

$$\mathbf{M} = \begin{pmatrix} 2 & -1 & 0 \\ -1 & 2 & -1 \\ 0 & -1 & 2 \end{pmatrix}, \mathbf{b} = (1,0,0), \mathbf{c} = 1$$



We firstly transform this problem to a simple form

Note that **M** is positive definite. According to the property of positive definite matrix, **M** can be decomposed as folloing form:

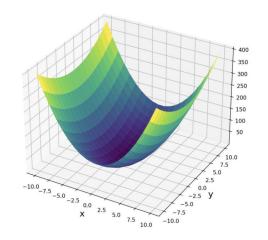
$$\mathbf{M} = \mathbf{U}^{\mathrm{T}} \mathbf{D} \mathbf{U},$$

where U is the orthogonal matrix and D is the diagonal matrix whose elements in the diagonal elements are M's eigenvalues.

This decomposition is called singular value decomposition, see https://en.wikipedia.org/wiki/Singular_value_decomposition

D can be written as follows:

$$\begin{pmatrix} a_1 & \cdots & 0 \\ \vdots & \ddots & \vdots \\ 0 & \cdots & a_d \end{pmatrix},$$



where $a_i > 0$ (i=1,...,d), because **M** is positive definite.

D can be written as follows:

$$\sqrt{D}^{T}\sqrt{D}$$
,

where \sqrt{D} is the diagonal matrix whose elements in the square roots of diagonal elements

$$\sqrt{\mathbf{D}} = \begin{pmatrix} \sqrt{a_1} & \cdots & 0 \\ \vdots & \ddots & \vdots \\ 0 & \cdots & \sqrt{a_d} \end{pmatrix}$$

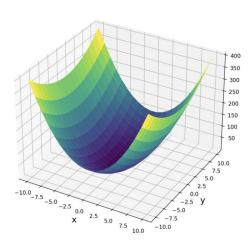
Then if we set (We transform this problem to a simple form as follows)

$$y = \sqrt{D} U x$$
, so $U^T \sqrt{D}^{-1} y = x$ and $x^T M x = y^T y$

Then, the problem will be transformed into

Minimize **b** $U^T \sqrt{D}^{-1} y$

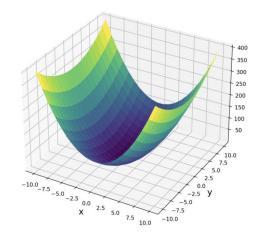
Subject to $\mathbf{y}^{\mathsf{T}}\mathbf{y} - \mathbf{c} \leq 0$.



Now we address this simpler issue

Minimize **b**
$$U^T \sqrt{D}^{-1}$$
 y

Subject to $\mathbf{y}^{\mathsf{T}}\mathbf{y} - \mathbf{c} \leq 0$.

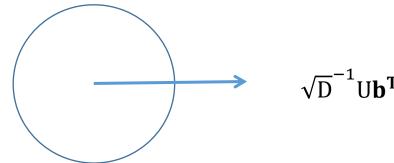


Using **Theorem 8**, we obtain that: if y_0 is the optimal solution, then

$$\nabla f(\mathbf{y}_0)^T(\mathbf{y} - \mathbf{y}_0) \ge 0$$
, for all \mathbf{y} satisfy $\mathbf{y}^T\mathbf{y} - c \le 0$.

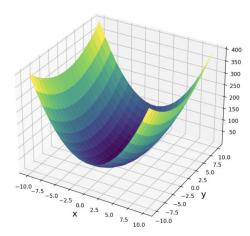
which is equal to
$$\mathbf{b} \mathbf{U}^{\mathrm{T}} \sqrt{\mathbf{D}}^{-1} \mathbf{y} \ge \mathbf{b} \mathbf{U}^{\mathrm{T}} \sqrt{\mathbf{D}}^{-1} \mathbf{y}_{0}$$

Note that $y^Ty - c \le 0$ means a ball with radius \sqrt{c} .



$$\mathbf{b} \ \mathbf{U}^{\mathrm{T}} \sqrt{\mathbf{D}}^{-1} \ \mathbf{y} \ge \mathbf{b} \ \mathbf{U}^{\mathrm{T}} \sqrt{\mathbf{D}}^{-1} \mathbf{y}_{0}$$

means that inner product between the optimal solution \mathbf{y}_0 and $\mathbf{b} \ \mathbf{U}^T \sqrt{\mathbf{D}}^{-1}$ should be smallest in the ball.



Cauchy inequality: https://en.wikipedia.org/wiki/Cauchy%E2%80%93Schwarz_inequality

 $x^{T} y \ge - ||x|| ||y|| (||*|| is L2 norm)$

and $\mathbf{x}^T \mathbf{y} = -\|\mathbf{x}\| \|\mathbf{y}\|$ if and only if $\mathbf{x} = -k\mathbf{y}$, where k is any positive

constant.

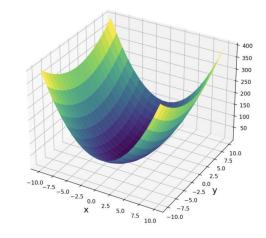
Using Cauchy inequality

$$\mathbf{b} \ \mathbf{U}^{\mathrm{T}} \sqrt{\mathbf{D}}^{-1} \mathbf{y} \ge \mathbf{b} \ \mathbf{U}^{\mathrm{T}} \sqrt{\mathbf{D}}^{-1} \mathbf{y}_{0} \ge - \left\| \mathbf{b} \ \mathbf{U}^{\mathrm{T}} \sqrt{\mathbf{D}}^{-1} \right\| \|\mathbf{y}_{0}\| \ge - \sqrt{c} \left\| \mathbf{b} \ \mathbf{U}^{\mathrm{T}} \sqrt{\mathbf{D}}^{-1} \right\|$$

So

$$\mathbf{b} \ \mathbf{U}^{\mathrm{T}} \sqrt{\mathbf{D}}^{-1} \mathbf{y} \geq -\sqrt{\mathbf{c}} \ \left\| \mathbf{b} \ \mathbf{U}^{\mathrm{T}} \sqrt{\mathbf{D}}^{-1} \right\|$$

If we take
$$\mathbf{y}_0 = -\sqrt{c}\sqrt{D}^{-1}U\mathbf{b}^T/\left\|\mathbf{b}\ U^T\sqrt{D}^{-1}\right\|$$
 ,



Then, it is clear that
$$\mathbf{b} \mathbf{U}^{\mathrm{T}} \sqrt{\mathbf{D}}^{-1} \mathbf{y}_0 = -\sqrt{\mathbf{c}} \| \mathbf{b} \mathbf{U}^{\mathrm{T}} \sqrt{\mathbf{D}}^{-1} \|.$$

So
$$\mathbf{y}_0 = -\sqrt{c}\sqrt{D}^{-1}U\mathbf{b}^{\mathrm{T}}/\left\|\mathbf{b} \mathbf{U}^{\mathrm{T}}\sqrt{D}^{-1}\right\|$$
 is the optimal solution.

So
$$\mathbf{x}_0 = \mathbf{U}^{\mathsf{T}} \sqrt{\mathbf{D}}^{-1} \mathbf{y}_0 = -\sqrt{\mathbf{c}} \, \mathbf{U}^{\mathsf{T}} \sqrt{\mathbf{D}}^{-1} \sqrt{\mathbf{D}}^{-1} \mathbf{U} \mathbf{b}^{\mathsf{T}} / \left\| \mathbf{b} \, \mathbf{U}^{\mathsf{T}} \sqrt{\mathbf{D}}^{-1} \right\|$$
$$= -\sqrt{\mathbf{c}} \, \mathbf{M}^{-1} \mathbf{b}^{\mathsf{T}} / \left\| \mathbf{b} \, \mathbf{U}^{\mathsf{T}} \sqrt{\mathbf{D}}^{-1} \right\|$$

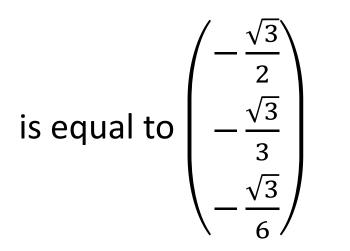
$$\left\| \mathbf{b} \sqrt{\mathbf{D}}^{-1} \mathbf{U}^{\mathrm{T}} \right\| = \sqrt{\mathbf{b} \ \mathbf{U}^{\mathrm{T}} \sqrt{\mathbf{D}}^{-1} \sqrt{\mathbf{D}}^{-1} \mathbf{U} \mathbf{b}^{\mathrm{T}}} = \sqrt{\mathbf{b} \mathbf{M}^{-1} \mathbf{b}^{\mathrm{T}}}$$

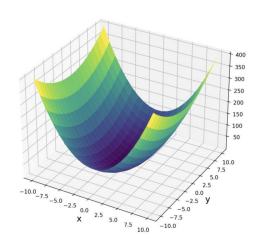
So, the optimal solution should be $-\sqrt{c} \mathbf{M}^{-1} \mathbf{b}^{\mathrm{T}} / \sqrt{\mathbf{b} \mathbf{M}^{-1} \mathbf{b}^{\mathrm{T}}}$

$$\mathbf{M}^{-1} = \begin{pmatrix} 0.75 & 0.5 & 0.25 \\ 0.5 & 1 & 0.5 \\ 0.25 & 0.5 & 0.75 \end{pmatrix}, \mathbf{b} = (1,0,0), \mathbf{c} = 1$$

So the solution

$$-\sqrt{c} \mathbf{M}^{-1} \mathbf{b}^{\mathrm{T}} / \sqrt{\mathbf{b} \mathbf{M}^{-1} \mathbf{b}^{\mathrm{T}}}$$

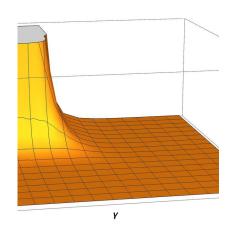


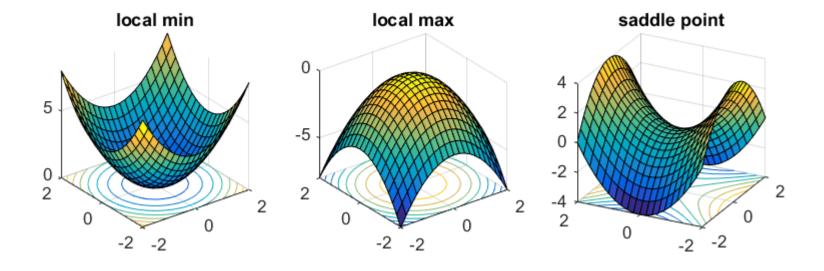


Non-Convex Optimization

Why is non-convex optimization hard?

- Potentially many local minimal points
- Saddle points
- Very flat regions



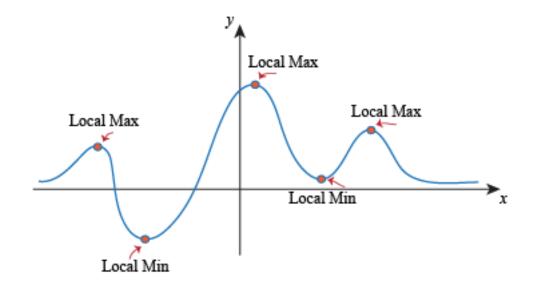


Global minimum point:

A real-valued function f defined on a domain Ω has a global maximum point at x*, if $f(x*) \leq f(x)$ for all x in Ω .

Local minimum point:

A real-valued function f defined on a domain Ω has a local maximum point at x*, if there exists some $\varepsilon > 0$ such that $f(x*) \le f(x)$ for all x in Ω within distance ε of x*

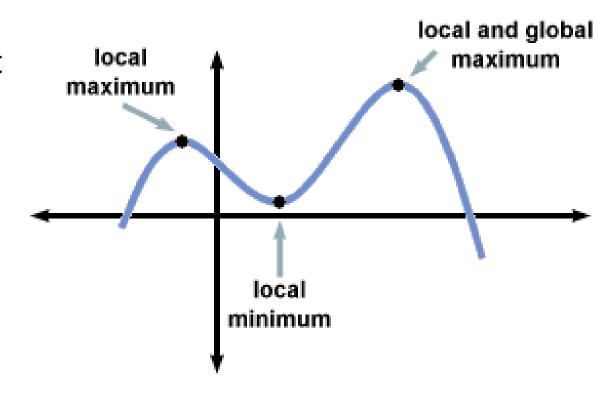


Global maximum point:

A real-valued function f defined on a domain Ω has a global maximum point at x*, if $f(x*) \ge f(x)$ for all x in Ω .

Local maximum point:

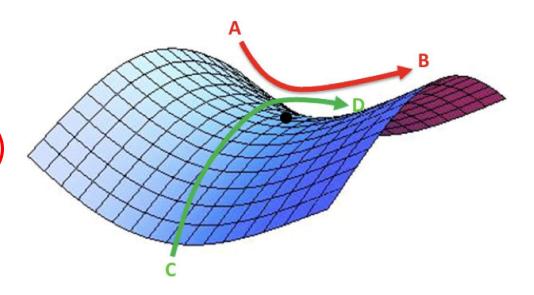
A real-valued function f defined on a domain Ω has a local maximum point at x*, if there exists some $\varepsilon > 0$ such that $f(x*) \ge f(x)$ for all x in Ω within distance ε of x*



Saddle Point:

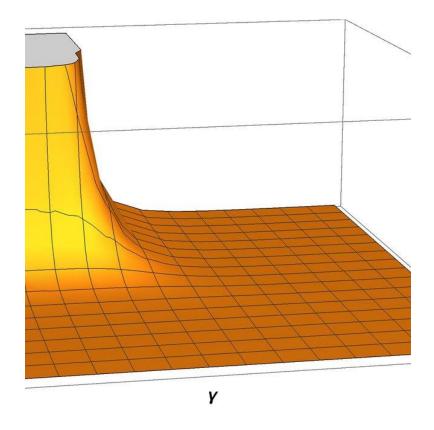
A saddle point or minimax point is a point on the surface of the graph of a function where the slopes (derivatives) in orthogonal directions are all zero (a critical point), but which is not a local extremum (local minimal point or local maximal point) of the function.

For example, x^2 - y^2 (see the figure). (0,0) is a saddle point, because the gradient at (0,0) is zero, but it is not the local extremum

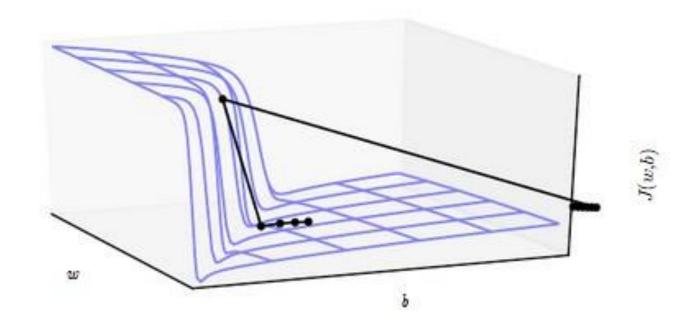


Very Flat Regions:

Very Flat Regions is an area on the surface of the graph of a function where the gredients (derivatives) in orthogonal directions are very close to zero, but which is not a local extremum of the function.



Cliffs and Exploding Gradients



Neural networks with many layers will have cliffs and exploding gradients. Therefore, gradient clipping is useful

How to solve non-convex problems?

- Gradient descent
- Stochastic gradient descent https://en.wikipedia.org/wiki/Stochastic_gradient_descent
- Adaptive gradient algorithm https://conferences.mpi-inf.mpg.de/adfocs/material/alina/adaptive-L1.pdf
- RMSprop https://optimization.cbe.cornell.edu/index.php?title=RMSProp
- Momentum https://en.wikipedia.org/wiki/Momentum

Consider the following non-convex optimization problem:

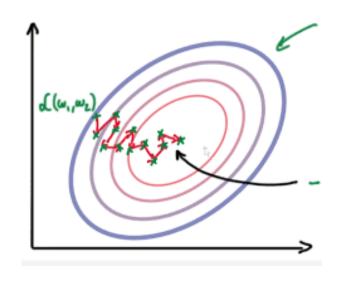
Minimize
$$f(x)$$
,

where f is non-convex and is defined in R^d.

Given an inital point x_0 (which is also called initial weight)

Then the next updated point should be

$$\mathbf{x_1} = \mathbf{x_0} - \mathbf{t} \nabla \mathbf{f}(\mathbf{x_0})$$



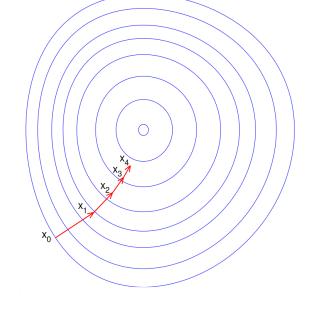
where t (t>0) is called learning rate.

After we obtain the n-th weight x_n , then the (n+1)-th weight x_{n+1} is

$$\mathbf{x_{n+1}} = \mathbf{x_n} - \mathsf{t} \nabla f(\mathbf{x_n})$$

Motivation 1. We hope the final point \mathbf{x} will get close to the a critical point $\nabla f(\mathbf{x}) = 0$.

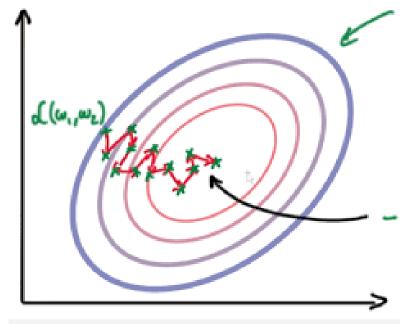
Motivation 2. to take repeated steps in the opposite direction of the gradient (or approximate gradient) of the function at the current point, because this is the direction of steepest descent (下降最快的方向).



Doesn't necessarily go towards optimal point.

Exercises

• Minimize f(x), where $f(x) = x^3$. Initial weight $x_0 = 1$, learning rate t = 1/3. Then what is the convergent point?



Solution

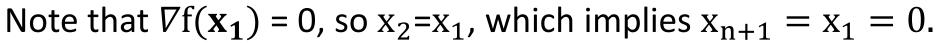
• Minimize f(x), where $f(x) = x^3$.

Then the grendent descent formula is

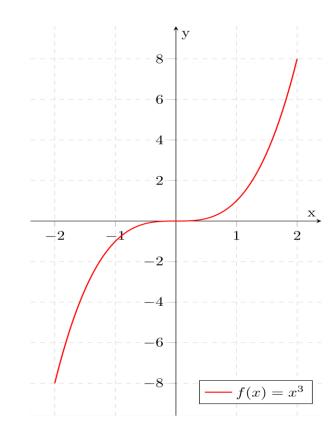
$$x_{n+1} = x_n - 3tx_n^2$$

Because $x_0 = 1$, t = 1/3, then

$$x_1 = 0$$



So 0 is the convergent point. But 0 is a saddle point of f(x).



Stochastic Gradient Descent

```
Algorithm 8.1 Stochastic gradient descent (SGD) update

Require: Learning rate schedule \epsilon_1, \epsilon_2, \ldots

Require: Initial parameter \theta

k \leftarrow 1

while stopping criterion not met do

Sample a minibatch of m examples from the training set \{x^{(1)}, \ldots, x^{(m)}\} with corresponding targets y^{(i)}.

Compute gradient estimate: \hat{g} \leftarrow \frac{1}{m} \nabla_{\theta} \sum_{i} L(f(x^{(i)}; \theta), y^{(i)})

Apply update: \theta \leftarrow \theta - \epsilon_k \hat{g}

k \leftarrow k + 1

end while
```

Sample i uniformly from $\{1, \dots, n\}$, and update θ by

$$\theta = \theta - \epsilon \nabla_{\theta} J^{(i)}(\theta)$$

AdaGrad

```
Algorithm 8.4 The AdaGrad algorithm

Require: Global learning rate \epsilon

Require: Initial parameter \theta

Require: Small constant \delta, perhaps 10^{-7}, for numerical stability

Initialize gradient accumulation variable r=0

while stopping criterion not met do

Sample a minibatch of m examples from the training set \{x^{(1)}, \dots, x^{(m)}\} with corresponding targets y^{(i)}.

Compute gradient: g \leftarrow \frac{1}{m} \nabla \theta \sum_i L(f(x^{(i)}; \theta), y^{(i)}).

Accumulate squared gradient: r \leftarrow r + g \odot g.

Compute update: \Delta \theta \leftarrow \underbrace{\epsilon}_{\delta + \sqrt{r}} \odot g. (Division and square root applied element-wise)

Apply update: \theta \leftarrow \theta + \Delta \theta.

end while
```

$$g_t = \nabla_{\theta_t} J(\theta_t)$$
 AdaGrad: $\theta_{t+1} = \theta_t - \frac{\eta}{\sqrt{G_t + \epsilon}} \odot g_t$

RMSProp

```
Algorithm 8.5 The RMSProp algorithm Require: Global learning rate \epsilon, decay rate \rho Require: Initial parameter \theta Require: Small constant \delta, usually 10^{-6}, used to stabilize division by small numbers Initialize accumulation variables r=0 while stopping criterion not met \mathbf{do} Sample a minibatch of m examples from the training set \{\boldsymbol{x}^{(1)},\dots,\boldsymbol{x}^{(m)}\} with corresponding targets \boldsymbol{y}^{(i)}. Compute gradient: \boldsymbol{g} \leftarrow \frac{1}{m} \nabla_{\boldsymbol{\theta}} \sum_{i} L(f(\boldsymbol{x}^{(i)};\boldsymbol{\theta}),\boldsymbol{y}^{(i)}). Accumulate squared gradient: \boldsymbol{r} \leftarrow \rho \boldsymbol{r} + (1-\rho)\boldsymbol{g}\odot\boldsymbol{g} Compute parameter update: \Delta \boldsymbol{\theta} = -\frac{\epsilon}{\sqrt{\delta+r}}\odot\boldsymbol{g}. (\frac{1}{\sqrt{\delta+r}} applied element-wise) Apply update: \boldsymbol{\theta} \leftarrow \boldsymbol{\theta} + \Delta \boldsymbol{\theta}. end while
```

Changing the gradient accumulation into an exponentially weighted moving average

Adam

```
Algorithm 8.7 The Adam algorithm
Require: Step size \epsilon (Suggested default: 0.001)
Require: Exponential decay rates for moment estimates, \rho_1 and \rho_2 in [0,1).
  (Suggested defaults: 0.9 and 0.999 respectively)
Require: Small constant \delta used for numerical stabilization (Suggested default:
  10^{-8})
Require: Initial parameters \theta
   Initialize 1st and 2nd moment variables s = 0, r = 0
   Initialize time step t = 0
   while stopping criterion not met do
     Sample a minibatch of m examples from the training set \{x^{(1)}, \dots, x^{(m)}\} with
     corresponding targets y^{(i)}.
     Compute gradient: \mathbf{g} \leftarrow \frac{1}{m} \nabla_{\boldsymbol{\theta}} \sum_{i} L(f(\mathbf{x}^{(i)}; \boldsymbol{\theta}), \mathbf{y}^{(i)})
     t \leftarrow t + 1
     Update biased first moment estimate. s \leftarrow \rho_1 s + (1 - \rho_1)g
     Update biased second moment estimate: r \leftarrow \rho_2 r + (1 - \rho_2) g \odot g
     Correct bias in first moment: \hat{s} \leftarrow \frac{s}{1}
     Correct bias in second moment: \hat{r} \leftarrow \frac{r}{1-r}
     Compute update: \Delta \theta = -\epsilon \frac{\hat{s}}{\sqrt{\hat{r}} + \delta} (operations applied element-wise)
     Apply update: \theta \leftarrow \theta + \Delta \theta
   end while
```

$$\begin{split} g_t &= \nabla_{\theta_t} J(\theta_t) \\ \text{Adam: } \theta_{t+1} &= \theta_t - \frac{\eta}{\sqrt{\hat{v}_t + \epsilon}} \, \widehat{m}_t \text{ ,} \\ \text{where } \widehat{m}_t &= \frac{m_t}{1 - \beta_1^t} \\ \text{and } \widehat{v}_t &= \frac{v_t}{1 - \beta_2^t} \\ m_t &= \beta_1 m_{t-1} + (1 - \beta_1) g_t \\ \text{and } v_t &= \beta_2 v_{t-1} + (1 - \beta_2) g_t^2 \end{split}$$

Using Adam when you are unfamiliar about optimization; while using momentum SGD when you are familiar about optimization

Thank You!