

Series 1, Feb 21st, 2019
(Probability, Analysis, Linear Algebra)

We will publish sample solutions on Friday, Mar 8th.

Problem 1 (Sampling):

Knowing the CDF of a random variable X , enables one to draw samples from that distribution.

- (a) Show that if X has distribution function F , and $U \sim \text{Unif}(0, 1)$ is a uniform random number in the interval $(0, 1)$, then $F^{-1}(U)$ has the same distribution as X .

In situations where the inverse of F is not easy to compute, one can use the following method (known as the *rejection method*) for generating random variables with a density f . Suppose that γ be a function such that $\gamma(x) \geq f(x)$ for all $x \in \mathbb{R}$, and

$$\int_{-\infty}^{\infty} \gamma(x) dx = \alpha < \infty.$$

Then, $g(x) = \gamma(x)/\alpha$ is a probability density function. Suppose we generate a random variable X by the following algorithm:

- I. Generate a random variable T with density function g .
 - II. Generate a random variable $U \sim \text{Unif}(0, 1)$, independent of T . If $U \leq f(T)/\gamma(T)$ then set $X = T$; if $U > f(T)/\gamma(T)$ then repeat steps I and II.
- (b) Show that the generated random variable X has density f .
- (c) Show that the number of rejections before X is generated has a Geometric distribution. Give an expression for the parameter of this distribution.

Hints For part (a) note that

$$F^{-1}(u) = \inf\{x \mid F(x) \geq u\},$$

as F is right-continuous. For part (b), you need to evaluate

$$\mathbb{P}(T \leq x \mid U \leq f(T)/\gamma(T)).$$

Solution 1:

- (a) Define $Y = F^{-1}(U)$, where $U \sim \text{Unif}(0, 1)$ and F is a CDF. Computing the CDF of Y gives

$$F_Y(y) = \mathbb{P}(Y \leq y) = \mathbb{P}(F^{-1}(U) \leq y) = \mathbb{P}(U \leq F(y)) = F(y),$$

where we used the fact that $F^{-1}(u) \leq y \iff u \leq F(y)$. Thus, Y has F as its distribution function.

(b) Before continuing reading the solution, take a look at the first box in the solution of problem 2 (b).

First let us compute the probability

$$\mathbb{P}\left(U \leq \frac{f(T)}{\gamma(T)}\right).$$

As a reminder, this probability is defined on the joint distribution of U and T . As U and T are independent, the joint probability space is simply the product space defined over $(0, 1) \times \mathbb{R}$.

By conditioning on the value of T , and using the fact that T has density g , we get the following:

$$\begin{aligned}\mathbb{P}\left(U \leq \frac{f(T)}{\gamma(T)}\right) &= \int_{\mathbb{R}} \mathbb{P}\left(U \leq \frac{f(t)}{\gamma(t)} \mid T = t\right) g(t) dt \\ &= \int_{\mathbb{R}} \mathbb{P}\left(U \leq \frac{f(t)}{\gamma(t)}\right) g(t) dt \\ &= \int_{\mathbb{R}} \frac{f(t)}{\gamma(t)} g(t) dt \\ &= \int_{\mathbb{R}} \frac{1}{\alpha} f(t) dt = \frac{1}{\alpha},\end{aligned}$$

where in the second line, we used the fact that T and U are independent, thus we can remove the conditioning.

Also, we need to compute the following probability for $x \in \mathbb{R}$:

$$\begin{aligned}\mathbb{P}\left(T \leq x, U \leq \frac{f(T)}{\gamma(T)}\right) &= \int_{\mathbb{R}} \mathbb{P}\left(T \leq x, U \leq \frac{f(t)}{\gamma(t)} \mid T = t\right) g(t) dt \\ &= \int_{-\infty}^x \frac{f(t)}{\gamma(t)} g(t) dt = \frac{1}{\alpha} \int_{-\infty}^x f(t) dt.\end{aligned}$$

Now, by the definition of conditional probability, we have

$$\mathbb{P}\left(T \leq x \mid U \leq \frac{f(T)}{\gamma(T)}\right) = \frac{\mathbb{P}\left(T \leq x, U \leq \frac{f(T)}{\gamma(T)}\right)}{\mathbb{P}\left(U \leq \frac{f(T)}{\gamma(T)}\right)} = \frac{\frac{1}{\alpha} \int_{-\infty}^x f(t) dt}{1/\alpha} = \int_{-\infty}^x f(t) dt.$$

This means that if the choice of T and U resulted in an acceptance, the density of T is f .

(c) As computed in part (b), the probability of acceptance is $1/\alpha$. One can think of it as a coin with bias (probability of heads) $p = 1/\alpha$. Thus, the number of throws (rejections) until the first heads (the first acceptance) has a Geometric distribution.

Problem 2 (Multivariate Normal Distribution):

Recall the following fact about characteristic functions:

Fact 1. For a random vector X in \mathbb{R}^d , define its characteristic function φ_X as

$$\varphi_X(\mathbf{t}) = \mathbb{E}[\exp(i\mathbf{t}^\top X)], \quad \text{for all } \mathbf{t} \in \mathbb{R}^d.$$

The characteristic function completely identifies a distribution. For a multivariate Normal distribution $\mathcal{N}(\mu, \Sigma)$, one has

$$\varphi(\mathbf{t}) = \exp(i\mathbf{t}^\top \mu - \frac{1}{2}\mathbf{t}^\top \Sigma \mathbf{t}).$$

- (a) Let $X = (X_1, \dots, X_d)$ be a d -dimensional standard Gaussian random vector, that is, $X \sim \mathcal{N}_d(0, I)$. Define $Y = AX + \mu$, where A is a $d \times d$ matrix and $\mu \in \mathbb{R}^d$. What is the distribution of Y ? If B is an $r \times d$ matrix, what is the distribution of BY ?
- (b) Let X be a bivariate Normal random variable (taking on values in \mathbb{R}^2) with mean $\mu = (1, 1)$ and covariance matrix $\Sigma = \begin{pmatrix} 3 & 1 \\ 1 & 2 \end{pmatrix}$. Find the conditional distribution of $Y = X_1 + X_2$ given $Z = X_1 - X_2 = 0$.
- (c*) For $Y \sim \mathcal{N}_d(0, I)$, we say that the random variable $V = \|Y\|^2$ has the χ^2 (chi-square) distribution with d degrees of freedom ($V \sim \chi^2(d)$). Assume that X_1, \dots, X_n are i.i.d. samples from the Normal distribution $\mathcal{N}(\mu, \sigma^2)$. One way to estimate σ^2 from these samples is to look at the sample variance:

$$S^2 = \frac{1}{n-1} \sum_{i=1}^n (X_i - \bar{X})^2,$$

where $\bar{X} = \frac{1}{n}(X_1 + \dots + X_n)$.

Prove that $\frac{(n-1)}{\sigma^2} S^2$ has a chi-square distribution with $n - 1$ degrees of freedom.

Hint: Can you write S^2 as the norm-squared of a vector? Which vector? Take care of the dimensions.

Solution 2:

(a) Let us compute the characteristic function of Y . Define $\mathbf{s} = A^\top \mathbf{t}$. We have

$$\begin{aligned} \varphi_Y(\mathbf{t}) &= \mathbb{E}[\exp(i\mathbf{t}^\top Y)] \\ &= \mathbb{E}[\exp(i\mathbf{t}^\top AX) \cdot \exp(i\mathbf{t}^\top \mu)] \\ &= \mathbb{E}[\exp(i\mathbf{s}^\top X)] \cdot \exp(i\mathbf{t}^\top \mu) \\ &= \varphi_X(\mathbf{s}) \cdot \exp(i\mathbf{t}^\top \mu) \\ &= \exp(-\frac{1}{2}\mathbf{s}^\top \mathbf{s} + i\mathbf{t}^\top \mu) \\ &= \exp(i\mathbf{t}^\top \mu - \frac{1}{2}\mathbf{t}^\top AA^\top \mathbf{t}), \end{aligned}$$

which means that $Y \sim \mathcal{N}(\mu, AA^\top)$. With the same argument as above, one gets $BY \sim \mathcal{N}(B\mu, BAA^\top B^\top)$.

(b) First, take a look at the following facts:

Let A, B be events. The definition of conditional probability $\mathbb{P}(A \mid B)$ assumes that $\mathbb{P}(B) \neq 0$. So one essentially cannot condition on events of zero probability in the usual way. The following is a workaround to this issue.

Let X, Y be random variables with joint density f and joint CDF F . For $\varepsilon > 0$ and $x, y \in \mathbb{R}$, we compute

$$\begin{aligned}\mathbb{P}(X \leq x \mid Y \in [y, y + \varepsilon]) &= \frac{\mathbb{P}(X \leq x, Y \in [y, y + \varepsilon])}{\mathbb{P}(Y \in [y, y + \varepsilon])} \\ &= \frac{F(x, y + \varepsilon) - F(x, y)}{F_Y(y + \varepsilon) - F_Y(y)} \\ &= \frac{[F(x, y + \varepsilon) - F(x, y)]/\varepsilon}{[F_Y(y + \varepsilon) - F_Y(y)]/\varepsilon}.\end{aligned}$$

Now if $\varepsilon \rightarrow 0$, the right hand side has the limit $\frac{\partial_y F(x, y)}{f_Y(y)}$, and the left hand side can be regarded as $\mathbb{P}(X \leq x \mid Y = y)$. Taking derivative with respect to x gives the conditional density

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

One can use this density to compute probabilities like $\mathbb{P}(X \in A \mid Y = y) = \int \int_A \frac{f(x, y)}{f_Y(y)} dx dy$.

We present two approaches for this exercise:

APPROACH 1. Note that $Z = 0$ implies $X_1 = X_2$. Furthermore by the definition of Y , we have $X_1 = X_2 = Y/2$ given $Z = 0$. Hence the marginal density of Y given $Z = 0$ is proportional to

$$f_{Y|Z}(y \mid 0) = \frac{f_{Y,Z}(y, 0)}{f_Z(0)} \propto f_{Y,Z}(y, 0) \propto f_X \left[\begin{pmatrix} y/2 \\ y/2 \end{pmatrix} \right].$$

The last equality is due to the fact that the linear map $(x_1, x_2) \mapsto (x_1 + x_2, x_1 - x_2)$ has constant determinant of -2 . Thus, by a change of variables formula, the density changes by a constant factor. We then have

$$\begin{aligned}f_X \left[\begin{pmatrix} y/2 \\ y/2 \end{pmatrix} \right] &\propto \exp \left(-\frac{1}{2} \begin{pmatrix} \frac{y}{2} - 1 \\ \frac{y}{2} - 1 \end{pmatrix}^T \begin{pmatrix} 3 & 1 \\ 1 & 2 \end{pmatrix}^{-1} \begin{pmatrix} \frac{y}{2} - 1 \\ \frac{y}{2} - 1 \end{pmatrix} \right) \\ &= \exp \left(-\frac{1}{2} \begin{pmatrix} \frac{y}{2} - 1 \\ \frac{y}{2} - 1 \end{pmatrix}^T \frac{1}{5} \begin{pmatrix} 2 & -1 \\ -1 & 3 \end{pmatrix} \begin{pmatrix} \frac{y}{2} - 1 \\ \frac{y}{2} - 1 \end{pmatrix} \right) \\ &= \exp \left(-\frac{1}{2} \frac{(y - 2)^2}{\frac{20}{3}} \right).\end{aligned}$$

Clearly, the conditional distribution of Y given $Z = 0$ is hence Normal with mean 2 and variance $\frac{20}{3}$.

In this problem, we used the following trick which prevents a lot of computational headaches. If one is trying to derive the density of a random variable X at x , that is, $f_X(x)$, it is easier to neglect all *multiplicative* terms that does not include x . The reason is simply because $\int_{\mathbb{R}} f_X(x) dx = 1$.

Two important examples are single variable Normal random variables and multivariate Gaussian vectors. In the first case, following the trick above, we conclude that if a density function is of the form

$$f(x) \propto \exp(-ax^2 + bx)$$

for $a > 0$ and $b \in \mathbb{R}$, by completing the squares, we obtain

$$-ax^2 + bx = -a\left(x - \frac{b}{2a}\right)^2 + \frac{b^2}{4a},$$

and thus, by removing the terms that does not depend on x , we get

$$f(x) \propto \exp\left(-\frac{(x - \frac{b}{2a})^2}{1/a}\right),$$

meaning that the distribution is a Normal distribution with mean $\frac{b}{2a}$ and variance $1/a$.

The situation for multivariate normal distribution is the same. One needs only to create a proper quadratic form in the exponent to get the familiar multivariate Gaussian density.

APPROACH 2. We define the random variable R as

$$R = \begin{pmatrix} Y \\ Z \end{pmatrix} = \underbrace{\begin{pmatrix} 1 & 1 \\ 1 & -1 \end{pmatrix}}_{=A} X.$$

Notice that R is a linear transformation of a Gaussian vector, and by part (a), it is a Gaussian vector. Thus, we only need to compute its mean and covariance matrix. By linearity of expectation, the mean μ_R of R is

$$\mathbb{E}[R] = A\mathbb{E}[X] = A\mu = \begin{pmatrix} 2 \\ 0 \end{pmatrix}.$$

The covariance matrix Σ_R of R is also given by part (a):

$$\Sigma_R = A\Sigma A^\top = \begin{pmatrix} 1 & 1 \\ 1 & -1 \end{pmatrix} \begin{pmatrix} 3 & 1 \\ 1 & 2 \end{pmatrix} \begin{pmatrix} 1 & 1 \\ 1 & -1 \end{pmatrix} = \begin{pmatrix} 7 & 1 \\ 1 & 3 \end{pmatrix}$$

The conditional density of Y given $Z = 0$ is then given by

$$\begin{aligned} f_{Y|Z}(y | 0) &= \frac{f_{Y,Z}(y, 0)}{f_Z(0)} \propto f_{Y,Z}(y, 0) \\ &\propto \exp\left(-\frac{1}{2} \begin{pmatrix} y-2 \\ 0 \end{pmatrix}^\top \begin{pmatrix} 7 & 1 \\ 1 & 3 \end{pmatrix}^{-1} \begin{pmatrix} y-2 \\ 0 \end{pmatrix}\right) \\ &= \exp\left(-\frac{1}{2} \begin{pmatrix} y-2 \\ 0 \end{pmatrix}^\top \frac{1}{20} \begin{pmatrix} 3 & -1 \\ -1 & 7 \end{pmatrix} \begin{pmatrix} y-2 \\ 0 \end{pmatrix}\right) \\ &= \exp\left(-\frac{1}{2} \frac{(y-2)^2}{\frac{20}{3}}\right). \end{aligned}$$

Clearly, the conditional distribution of Y given $Z = 0$ is hence Normal with mean 2 and variance $\frac{20}{3}$.

(c*) Before bringing the solution, let us introduce three important linear transformations.

Let v be a unit vector in \mathbb{R}^d . The *orthogonal projector* on the direction of v is the linear transformation described by the matrix vv^\top (the outer product of v by itself). The orthogonal projection on the hyperplane defined by v is then $I - vv^\top$. Also the *reflection* about the hyperplane defined by v is $I - 2vv^\top$ (verify these by drawing a picture). Sometimes, the last transformation is called a *Householder Reflector*.

If one is searching for a unitary matrix that maps u to v , one possible way is to consider the Householder reflector about the hyperplane defined by $(v - u)/\|v - u\|$.

Let us now define $X = (X_1, \dots, X_n)$. The fact that X_i are i.i.d. implies that $X \sim \mathcal{N}(\mu, \sigma^2 I_n)$. Consider the unit vector $v = (1/\sqrt{n}, \dots, 1/\sqrt{n})$. The projection of X on the direction of v is

$$vv^\top X = \begin{pmatrix} 1/n & \cdots & 1/n \\ \vdots & \ddots & \vdots \\ 1/n & \cdots & 1/n \end{pmatrix} X = \begin{pmatrix} \bar{X} \\ \vdots \\ \bar{X} \end{pmatrix}.$$

Thus, the projection on the hyperplane defined by v is the vector $Y = (I - vv^\top)X = (X_1 - \bar{X}, \dots, X_n - \bar{X})$. Note here that $\|Y\|^2 = (n-1)S^2$. Note that Y is a Gaussian vector, as it is a linear function of X . Also notice that the transformation $I - vv^\top$ is of rank $n-1$. Hence, it is better to transform Y in a way that one component becomes zero, while keeping the norm of Y fixed. That is, we need a unitary map that maps v to $w = (1, 0, \dots, 0)$. Using Householder reflectors, this map is indeed $I - 2uu^\top$, where $u = (v - w)/\|v - w\|$.

Denote by $Z = (I - 2uu^\top)Y$. Observe that Z is a Gaussian vector. It is easy to verify that the mean of Z is zero. The covariance matrix can be computed using part (a):

$$\begin{aligned}\Sigma_Z &= (I - 2uu^\top)(I - vv^\top)(\sigma^2 I)(I - vv^\top)^\top(I - 2uu^\top)^\top \\ &= \sigma^2(I - 2uu^\top)(I - vv^\top)(I - 2uu^\top) \\ &= \sigma^2 \begin{pmatrix} 1 & 0 & \cdots & 0 \\ 0 & 1 & \cdots & 0 \\ \vdots & \vdots & \ddots & \vdots \\ 0 & 0 & \cdots & 0 \end{pmatrix}\end{aligned}$$

Thus, Z/σ is a Gaussian random vector, that is supported on a $(n-1)$ -dimensional space, with mean 0 and covariance I_{n-1} . That is, it is a standard Gaussian vector in \mathbb{R}^{n-1} . Hence, $\frac{1}{\sigma^2}\|Z\|^2$ has chi-square distribution with $(n-1)$ degrees of freedom. But $(n-1)S^2 = \|Y\|^2 = \|Z\|^2$. Thus,

$$\frac{(n-1)}{\sigma^2}S^2 \sim \chi^2(n-1).$$

Problem 3 (Linear Regression and Ridge Regression):

Let $D = \{(\mathbf{x}_1, y_1), (\mathbf{x}_2, y_2), \dots, (\mathbf{x}_n, y_n)\}$ where $\mathbf{x}_i \in \mathbb{R}^d$ and $y_i \in \mathbb{R}$ be the training data that you are given. As you have to predict a continuous variable, one of the simplest possible models is linear regression, i.e. to predict y as $\mathbf{w}^T \mathbf{x}$ for some parameter vector $\mathbf{w} \in \mathbb{R}^d$.¹ We thus suggest minimizing the following loss

$$\operatorname{argmin}_{\mathbf{w}} \hat{R}(\mathbf{w}) = \operatorname{argmin}_{\mathbf{w}} \sum_{i=1}^n (y_i - \mathbf{w}^T \mathbf{x}_i)^2. \quad (1)$$

Let us introduce the $n \times d$ matrix $\mathbf{X} \in \mathbb{R}^{n \times d}$ with the \mathbf{x}_i as rows, and the vector $\mathbf{y} \in \mathbb{R}^n$ consisting of the scalars y_i . Then, (1) can be equivalently re-written as

$$\operatorname{argmin}_{\mathbf{w}} \|\mathbf{X}\mathbf{w} - \mathbf{y}\|^2.$$

In this exercise, $\|\cdot\|$ is always the Euclidean norm. We refer to any \mathbf{w}^* that attains the above minimum as a solution to the problem.

- (a) Show that if $\mathbf{X}^T \mathbf{X}$ is invertible, then there is a unique \mathbf{w}^* that can be computed as $\mathbf{w}^* = (\mathbf{X}^T \mathbf{X})^{-1} \mathbf{X}^T \mathbf{y}$.
- (b) Show for $n < d$ that (1) does not admit a unique solution. Intuitively explain why this is the case.
- (c) Consider the case $n \geq d$. Under what assumptions on \mathbf{X} does (1) admit a unique solution \mathbf{w}^* ? Give an example with $n = 3$ and $d = 2$ where these assumptions do not hold.

The *ridge regression* optimization problem with parameter $\lambda > 0$ is given by

$$\operatorname{argmin}_{\mathbf{w}} \hat{R}_{\text{ridge}}(\mathbf{w}) = \operatorname{argmin}_{\mathbf{w}} \left[\sum_{i=1}^n (y_i - \mathbf{w}^T \mathbf{x}_i)^2 + \lambda \mathbf{w}^T \mathbf{w} \right]. \quad (2)$$

- (d) Show that \hat{R}_{ridge} is convex with respect to \mathbf{w} . You can use the fact that a twice differentiable function $f : \mathbb{R}^d \rightarrow \mathbb{R}$ is convex if and only if for any $\mathbf{x} \in \mathbb{R}^d$ its Hessian $D^2 f(\mathbf{x}) \in \mathbb{R}^{d \times d}$ is positive semi-definite.
- (e) Derive the closed form solution $\mathbf{w}_{\text{ridge}}^* = (\mathbf{X}^T \mathbf{X} + \lambda I_d)^{-1} \mathbf{X}^T \mathbf{y}$ to (2), where I_d denotes the identity matrix of size $d \times d$.
- (f) A continuously differentiable function $f : \mathbb{R}^d \rightarrow \mathbb{R}$ is called α -strongly convex for some $\alpha > 0$, if for any points $\mathbf{x}, \mathbf{y} \in \mathbb{R}^d$ one has

$$f(\mathbf{y}) \geq f(\mathbf{x}) + \nabla f(\mathbf{x})^T (\mathbf{y} - \mathbf{x}) + \frac{\alpha}{2} \|\mathbf{y} - \mathbf{x}\|^2.$$

If f is twice differentiable, an equivalent condition is that for any point $\mathbf{x} \in \mathbb{R}^d$, one has

$$D^2 f(\mathbf{x}) \succeq \alpha I,$$

which means $D^2 f(\mathbf{x}) - \alpha I$ is always positive semi-definite. Prove that a strongly convex function admits a unique minimizer in \mathbb{R}^d . *Hint: prove that $f(\mathbf{x}) \rightarrow \infty$ as $\|\mathbf{x}\| \rightarrow \infty$.*

- (g) Show that (2) admits the unique solution $\mathbf{w}_{\text{ridge}}^*$ for any matrix \mathbf{X} . Show that this even holds for the cases in (b) and (c) where (1) does not admit a unique solution \mathbf{w}^* .
- (h) What is the role of the term $\lambda \mathbf{w}^T \mathbf{w}$ in \hat{R}_{ridge} ? What happens to $\mathbf{w}_{\text{ridge}}^*$ as $\lambda \rightarrow 0$ and $\lambda \rightarrow \infty$? You do not need to give a complete proof, only an intuitive answer suffice.

¹Without loss of generality, we assume that both \mathbf{x}_i and y_i are centered, i.e. they have zero empirical mean. Hence we can neglect the otherwise necessary bias term.

Solution 3:

- (a) Note that $\hat{R} : \mathbb{R}^d \rightarrow \mathbb{R}$ and

$$\hat{R}(\mathbf{w}) = \|\mathbf{X}\mathbf{w} - \mathbf{y}\|^2 = (\mathbf{X}\mathbf{w} - \mathbf{y})^T (\mathbf{X}\mathbf{w} - \mathbf{y}) = \mathbf{w}^T \mathbf{X}^T \mathbf{X} \mathbf{w} - 2\mathbf{w}^T \mathbf{X}^T \mathbf{y} + \mathbf{y}^T \mathbf{y}.$$

The gradient of this function is equal to (see the recap slides; also note that the gradient is a vector in \mathbb{R}^d)

$$\nabla \hat{R}(\mathbf{w}) = 2\mathbf{X}^T \mathbf{X} \mathbf{w} - 2\mathbf{X}^T \mathbf{y}.$$

Because $\hat{R}(\mathbf{w})$ is convex (formally proven in (d)), its optima (if they exist) are exactly those points that have a zero gradient, i.e., those \mathbf{w}^* that satisfy $\mathbf{X}^T \mathbf{X} \mathbf{w}^* = \mathbf{X}^T \mathbf{y}$. Under the given assumption, the unique minimizer is indeed equal to $\mathbf{w}^* = (\mathbf{X}^T \mathbf{X})^{-1} \mathbf{X}^T \mathbf{y}$.

- (b) Consider the *singular value decomposition* $\mathbf{X} = \mathbf{U}\Sigma\mathbf{V}^T$ where \mathbf{U} is an unitary $n \times n$ matrix, \mathbf{V} is a unitary $d \times d$ matrix and Σ is a diagonal $n \times d$ matrix with the singular values of \mathbf{X} on the diagonal. We then have

$$\operatorname{argmin}_{\mathbf{w}} \hat{R}(\mathbf{w}) = \operatorname{argmin}_{\mathbf{w}} [\mathbf{w}^T \mathbf{V} \Sigma^2 \mathbf{V}^T \mathbf{w} - 2\mathbf{y}^T \mathbf{U} \Sigma \mathbf{V}^T \mathbf{w}]$$

Note that $\mathbf{y}^T \mathbf{U} \Sigma \mathbf{V}^T \mathbf{w} \in \mathbb{R}$ is a number. Thus, we have the equality $\mathbf{y}^T \mathbf{U} \Sigma \mathbf{V}^T \mathbf{w} = \mathbf{w}^T \mathbf{V} \Sigma^T \mathbf{U}^T \mathbf{y}$. Also, for brevity, we write Σ^2 instead of $\Sigma^T \Sigma \in \mathbb{R}^{d \times d}$.

Since \mathbf{V} is unitary (and hence it is a bijection), we may rotate \mathbf{w} using \mathbf{V} to $\mathbf{z} = \mathbf{V}^T \mathbf{w}$ and formulate the optimization problem in terms of \mathbf{z} , i.e.

$$\operatorname{argmin}_{\mathbf{z}} [\mathbf{z}^T \Sigma^2 \mathbf{z} - 2\mathbf{y}^T \mathbf{U} \Sigma \mathbf{z}] = \operatorname{argmin}_{\mathbf{z}} \sum_{i=1}^d [z_i^2 \sigma_i^2 - 2(\mathbf{U}^T \mathbf{y})_i z_i \sigma_i]$$

where σ_i is the i th entry in the diagonal of Σ . Note that this problem gets decomposed into d independent optimization problems of the form

$$z_i = \operatorname{argmin}_z [z^2 \sigma_i^2 - 2(\mathbf{U}^T \mathbf{y})_i z \sigma_i]$$

for $i = 1, 2, \dots, d$. Since each problem is quadratic with positive coefficient and thus convex we may obtain the solution by finding the root of the first derivative. For $i = 1, 2, \dots, d$ we require that z_i satisfies

$$z_i \sigma_i^2 - (\mathbf{U}^T \mathbf{y})_i \sigma_i = 0.$$

For all $i = 1, 2, \dots, d$ such that $\sigma_i \neq 0$, the solution z_i is thus given by

$$z_i = \frac{(\mathbf{U}^T \mathbf{y})_i}{\sigma_i}.$$

For the case $n < d$, however, \mathbf{X} has at most rank n as it is a $n \times d$ matrix and hence at most n of its singular values are nonzero.

We use the fact that the rank of a matrix A is equal to the number of nonzero singular values of A .

This means that there is at least one index j such that $\sigma_j = 0$ and hence any $z_j \in \mathbb{R}$ is a solution to the optimization problem. As a result, the set of optimal solutions for \mathbf{z} is a linear subspace of at least one dimension. By rotating this subspace back using \mathbf{V} , i.e., $\mathbf{w} = \mathbf{V}\mathbf{z}$, it is evident that the optimal solution to the optimization problem in terms of \mathbf{w} is also a linear subspace of at least one dimension and that thus no unique solution exists. Furthermore, since \mathbf{X} has at most rank n , $\mathbf{X}^T \mathbf{X}$ is not of full rank (for a proof, look at the SVD of $\mathbf{X}^T \mathbf{X}$). As a result $(\mathbf{X}^T \mathbf{X})^{-1}$ does not exist and \mathbf{w}^* is ill-defined.

The intuition behind these results is that the “linear system” $\mathbf{X}\mathbf{w} \approx \mathbf{y}$ is underdetermined as there are less data points than parameters that we want to estimate.

- (c) We showed in (b) that the optimization problem admits a unique solution only if all the singular values of \mathbf{X} are nonzero. For $n \geq d$, this is the case if and only if \mathbf{X} is of full rank, i.e., all the columns of \mathbf{X} are linearly independent. As an example for a matrix not satisfying these assumptions, any matrix with linearly dependent columns suffices, e.g.,

$$\mathbf{X}_{\text{degenerate}} = \begin{pmatrix} 1 & -2 \\ 0 & 0 \\ -2 & 4 \end{pmatrix}.$$

- (d) Because convex functions are closed under addition, we will show that each term in the objective is convex, from which the claim will follow. Each data term $(y_i - \mathbf{w}^T \mathbf{x}_i)^2$ has the Hessian $\mathbf{x}_i \mathbf{x}_i^T$, which is positive semi-definite because for any $\mathbf{w} \in \mathbb{R}^d$ we have $\mathbf{w}^T \mathbf{x}_i \mathbf{x}_i^T \mathbf{w} = (\mathbf{x}_i^T \mathbf{w})^2 \geq 0$ (note that $\mathbf{x}_i^T \mathbf{w} = \mathbf{w}^T \mathbf{x}_i$ are scalars).

The regularizer $\lambda \mathbf{w}^T \mathbf{w}$ has the identity matrix λI_d as a Hessian, which is also positive semi-definite because for any $\mathbf{w} \in \mathbb{R}^d$ we have $\mathbf{w}^T (\lambda I_d) \mathbf{w} = \lambda \|\mathbf{w}\|^2 \geq 0$, and this completes the proof.

- (e) The gradient of $\hat{R}_{\text{ridge}}(\mathbf{w})$ with respect to \mathbf{w} is given by

$$\nabla \hat{R}_{\text{ridge}}(\mathbf{w}) = 2\mathbf{X}^T(\mathbf{X}\mathbf{w} - \mathbf{y}) + 2\lambda\mathbf{w}.$$

Similar to (a), because $\hat{R}_{\text{ridge}}(\mathbf{w})$ is convex, we only have to find a point $\mathbf{w}_{\text{ridge}}^*$ such that

$$\nabla \hat{R}_{\text{ridge}}(\mathbf{w}_{\text{ridge}}^*) = 2\mathbf{X}^T(\mathbf{X}\mathbf{w}_{\text{ridge}}^* - \mathbf{y}) + 2\lambda\mathbf{w}_{\text{ridge}}^* = 0.$$

This is equivalent to

$$(\mathbf{X}^T \mathbf{X} + \lambda I_d) \mathbf{w}_{\text{ridge}}^* = \mathbf{X}^T \mathbf{y}$$

which implies the required result

$$\mathbf{w}_{\text{ridge}}^* = (\mathbf{X}^T \mathbf{X} + \lambda I_d)^{-1} \mathbf{X}^T \mathbf{y}.$$

- (f) First, let us prove that the function f is *coercive*, i.e., $\lim_{\|\mathbf{x}\| \rightarrow \infty} f(\mathbf{x}) = \infty$. In the definition of strong convexity, by putting $\mathbf{x} = 0$ we get

$$f(\mathbf{y}) \geq f(0) + \nabla f(0)^T \mathbf{y} + \frac{\alpha}{2} \|\mathbf{y}\|^2 \geq f(0) - \|\nabla f(0)\| \cdot \|\mathbf{y}\| + \frac{\alpha}{2} \|\mathbf{y}\|^2,$$

where we used the Cauchy-Schwartz inequality: $\nabla f(0)^T \mathbf{y} \geq -\|\nabla f(0)\| \cdot \|\mathbf{y}\|$. The right-hand side of the equation above is a quadratic function of $\|\mathbf{y}\|$ with a positive coefficient for second degree term. Thus, it goes to infinity as $\|\mathbf{y}\| \rightarrow \infty$. Hence, f also goes to infinity.

Next, we prove that f has a global minimum. Denote by $s = \inf_{\mathbf{x} \in \mathbb{R}^d} f(\mathbf{x}) < \infty$. Then, there exists a sequence $\mathbf{x}_1, \mathbf{x}_2, \dots$ such that $f(\mathbf{x}_n) \rightarrow s$. We claim that this sequence is bounded: otherwise, there was a subsequence that $\|\mathbf{x}_{n_i}\| \rightarrow \infty$. But as f is coercive, $f(\mathbf{x}_{n_i}) \rightarrow \infty$, contradicting $f(\mathbf{x}_{n_i}) \rightarrow s < \infty$. Hence, the sequence $\mathbf{x}_1, \mathbf{x}_2, \dots$ is inside some bounded set. By compactness, we obtain that there exists a convergent subsequence. As f is continuous, the f value of this subsequence converges as well, meaning that the infimum is attained. That is, $\exists \mathbf{x}_\infty : f(\mathbf{x}_\infty) = s = \inf f(\mathbf{x})$.

Finally, we prove uniqueness. If \mathbf{x} and \mathbf{y} were two distinct global minima for f , then, by strong convexity, we have

$$f\left(\frac{\mathbf{x} + \mathbf{y}}{2}\right) < \frac{1}{2}(f(\mathbf{x}) + f(\mathbf{y})) = \min f,$$

a contradiction.

- (g) Note that $\mathbf{X}^T \mathbf{X}$ is a positive semi-definite matrix, since $\forall \mathbf{w} \in \mathbb{R}^d : \mathbf{w}^T \mathbf{X}^T \mathbf{X} \mathbf{w} = \|\mathbf{X}\mathbf{w}\|^2 \geq 0$, which implies that it has non-negative eigenvalues. But then, $\mathbf{X}^T \mathbf{X} + \lambda I_d$ has eigenvalues bounded from below by $\lambda > 0$, which means that it is invertible and thus the optimum is uniquely defined.

Note. Since $\mathbf{X}^T \mathbf{X}$ is symmetric, all of its eigenvalues are real, and it is clear that μ is an eigenvalue of $\mathbf{X}^T \mathbf{X}$ if and only if $\mu + \lambda$ is an eigenvalue of $\mathbf{X}^T \mathbf{X} + \lambda I$. Also note that if a linear function is injective, then its kernel is $\{\mathbf{0}\}$, meaning that it does not have a zero eigenvalue. The converse is also true.

Another way to state the result is that $\mathbf{X}^T \mathbf{X} + \lambda I_d \succeq \lambda I_d$, which means that \hat{R}_{ridge} is λ -strongly convex. Thus, the claim follows.

- (h) The term $\lambda \mathbf{w}^T \mathbf{w}$ “biases” the solution towards the origin, i.e., there is a quadratic penalty for solutions \mathbf{w} that are far from the origin. The parameter λ determines the extend of this effect: As $\lambda \rightarrow 0$, $\hat{R}_{\text{ridge}}(\mathbf{w})$ converges to $\hat{R}(\mathbf{w})$. As a result the optimal solution $\mathbf{w}_{\text{ridge}}^*$ approaches the solution of (1). As $\lambda \rightarrow \infty$, only the quadratic penalty $\mathbf{w}^T \mathbf{w}$ is relevant and $\mathbf{w}_{\text{ridge}}^*$ hence approaches the null vector $(0, 0, \dots, 0)$.

One can also show this interesting property (however, the proof is involved): Assume $n < d$ (as the situation discussed in (b)). Then \mathbf{w}^* for linear regression is not unique. Denote by \mathbf{w}_λ^* the *unique* solution to the Ridge regression problem for $\lambda > 0$. Then the limit $\lim_{\lambda \rightarrow 0} \mathbf{w}_\lambda^*$ exists, and the limit point falls inside the space of solutions to linear regression problem. One can further show that this solution is the one with the minimum norm.

Problem 4 (Normal Random Variables):

Let X be a Normal random variable with mean $\mu \in \mathbb{R}$ and variance $\tau^2 > 0$, i.e. $X \sim \mathcal{N}(\mu, \tau^2)$. Furthermore, the random variable Y given $X = x$ is normally distributed with mean x and variance σ^2 , i.e. $Y|_{X=x} \sim \mathcal{N}(x, \sigma^2)$.

- (a) Derive the *marginal distribution* of Y , i.e. compute the density $f_Y(y)$.
- (b) Use Bayes' theorem to derive the *conditional distribution* of X given $Y = y$.

Hint: For both tasks derive the density up to a constant factor and use this to identify the distribution.

Solution 4:

Before starting calculations, it is good to mention that one can easily compute the following integral for $a > 0$ by creating complete squares:

$$\begin{aligned} \int_{\mathbb{R}} e^{-(ax^2+2bx+c)} dx &= \int_{\mathbb{R}} \exp \left(-a \left[\left(x + \frac{b}{a} \right)^2 - \frac{b^2 - ac}{a^2} \right] \right) dx \\ &= \exp \left(\frac{b^2 - ac}{a} \right) \cdot \int_{\mathbb{R}} \exp \left(-\frac{1}{2} \frac{\left(x + \frac{b}{a} \right)^2}{1/2a} \right) dx \\ &= \exp \left(\frac{b^2 - ac}{a} \right) \sqrt{\pi/a} \end{aligned}$$

As a prelude to both (a) and (b) we consider the joint density function $f_{X,Y}(x, y)$ of X and Y

$$f_{X,Y}(x, y) = f_{Y|X}(y|x)f_X(x) = \frac{1}{2\pi\sigma\tau} \exp \left(-\frac{1}{2} \underbrace{\left[\frac{(x-\mu)^2}{\tau^2} + \frac{(y-x)^2}{\sigma^2} \right]}_{(A)} \right).$$

For brevity, let us define

$$\begin{aligned} a &:= \frac{\sigma^2 + \tau^2}{2\sigma^2\tau^2}, \\ b &:= -\frac{\sigma^2\mu + \tau^2y}{2\sigma^2\tau^2}, \\ c &:= \frac{\sigma^2\mu^2 + \tau^2y^2}{2\sigma^2\tau^2}. \end{aligned}$$

Using simple algebraic operations, we obtain that $(A) = ax^2 + 2bx + c$.

- (a) The marginal density of Y is given by

$$f_Y(y) = \int_{\mathbb{R}} f_{X,Y}(x, y) dx = \int_{\mathbb{R}} f_{Y|X}(y|x)f_X(x) dx.$$

Using the formula discussed at the beginning of the solution, we can compute this integral by just putting

in the values of a, b and c :

$$\begin{aligned}
f_Y(y) &= \int_{\mathbb{R}} f_{X,Y}(x,y) dx \\
&= \int_{\mathbb{R}} \frac{1}{2\pi\sigma\tau} e^{-(ax^2+2bx+c)} dx \\
&= \frac{1}{2\pi\sigma\tau} \exp\left(\frac{b^2-ac}{a}\right) \sqrt{\pi/a} \\
&\propto \exp\left(\frac{b^2-ac}{a}\right) \quad (a \text{ does not depend on } y)
\end{aligned}$$

Now we try to write $(b^2 - ac)/a$ as a complete square:

$$\begin{aligned}
\frac{b^2 - ac}{a} &= \frac{1}{a} \left\{ \left(\frac{\sigma^2\mu + \tau^2y}{2\sigma^2\tau^2} \right)^2 - \frac{(\sigma^2 + \tau^2)(\sigma^2\mu^2 + \tau^2y^2)}{(2\sigma^2\tau^2)^2} \right\} \\
&= -\frac{1}{a} \cdot \frac{1}{(2\sigma^2\tau^2)^2} \cdot (\sigma^2\tau^2y^2 - 2\tau^2\sigma^2\mu y + \sigma^2\tau^2\mu^2) \\
&= -\frac{1}{a} \cdot \frac{\sigma^2\tau^2}{(2\sigma^2\tau^2)^2} \cdot ((y - \mu)^2 + \dots) \\
&= -\frac{1}{2} \frac{1}{(\sigma^2 + \tau^2)} \cdot ((y - \mu)^2 + \dots)
\end{aligned}$$

Putting everything together yields

$$f_Y(y) \propto \exp\left[-\frac{1}{2} \frac{(y - \mu)^2}{(\sigma^2 + \tau^2)}\right],$$

meaning that Y has a Gaussian distribution with mean μ and variance $\sigma^2 + \tau^2$.

(b) The conditional density of X given $Y = y$ is proportional to the joint density function, i.e.

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

By the discussion at the beginning of the solution, $f_{X,Y}(x,y) \propto \exp(-(ax^2 + 2bx + c))$. Since c does not depend on x (and y is considered as fixed/given), we can say :

$$f_{X|Y}(x|y) \propto \exp\left(-\frac{1}{2} \frac{\left(x + \frac{b}{a}\right)^2}{1/2a}\right)$$

So the mean would be $-b/a$ and the variance will be $1/2a$. Concretely:

$$\text{mean} = -\frac{b}{a} = \frac{\sigma^2\mu + \tau^2y}{\sigma^2 + \tau^2} = \frac{\sigma^2}{\sigma^2 + \tau^2}\mu + \frac{\tau^2}{\sigma^2 + \tau^2}y$$

Note that the mean is a convex combination of μ and the observation y . Also

$$\text{variance} = \frac{1}{2a} = \frac{\sigma^2\tau^2}{\sigma^2 + \tau^2}.$$