

ENGR 421/DASC 521: Introduction to Machine Learning
Fall 2021 Final – Solution Key

Question 1:

We have two classes that are assumed to have one-dimensional Gaussian distributions with different means and variances: $p(x|y = 1) \sim N(\mu_1, \sigma^2)$ and $p(x|y = 2) \sim N(\mu_2, 4\sigma^2)$. Derive the position(s) of the intersection of the two posterior probabilities.

We would like to find x that satisfy $P(y = 1|x) = P(y = 2|x)$.

$$\begin{aligned} p(x|y = 1)P(y = 1) &= p(x|y = 2)P(y = 2) \\ \log p(x|y = 1) + \log P(y = 1) &= \log p(x|y = 2) + \log P(y = 2) \\ -\frac{\log(2\pi\sigma^2)}{2} - \frac{(x - \mu_1)^2}{2\sigma^2} + \log P(y = 1) &= -\frac{\log(8\pi\sigma^2)}{2} - \frac{(x - \mu_2)^2}{8\sigma^2} + \log P(y = 2) \\ -\frac{3}{8\sigma^2}x^2 + \frac{4\mu_1 - \mu_2}{4\sigma^2}x - \frac{(2\mu_1 - \mu_2)(2\mu_1 + \mu_2)}{8\sigma^2} + \log \frac{P(y = 1)}{P(y = 2)} + \log 2 &= 0 \end{aligned}$$

Question 2:

Consider a data set of N data points, in which each data point has one real-valued positive input x_i and the corresponding real-valued output y_i , i.e., $\{(x_i, y_i)\}_{i=1}^N$. We use the following model to fit the data, which has an unknown parameter w (the variance is known in advance and is set to 1).

$$p(y_i|x_i) \sim N(\log(wx_i), 1) \quad \forall i$$

- (a) Describe a maximum likelihood approach to infer w and write down the log-likelihood objective for this problem.
- (b) Find the maximum likelihood solution for w .

- (a) By assuming the data points are independent from each other, we can write down the likelihood function as follows:

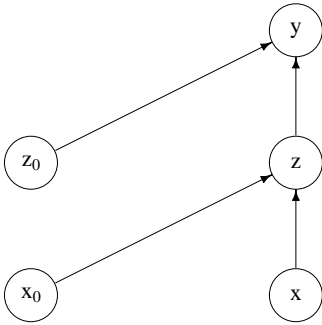
$$\begin{aligned} \text{likelihood} &= \prod_{i=1}^N p(y_i|x_i) = \prod_{i=1}^N N(\log(wx_i), 1) \\ &= \prod_{i=1}^N \frac{1}{\sqrt{2\pi(1)^2}} \exp\left(-\frac{(y_i - \log(wx_i))^2}{2(1)^2}\right) \\ \text{log-likelihood} &= \sum_{i=1}^N \left(-\frac{1}{2} \log(2\pi) - \frac{1}{2}(y_i - \log(wx_i))^2\right) \end{aligned}$$

- (b) To maximize log-likelihood, we need to minimize $\sum_{i=1}^N (1/2)(y_i - \log(wx_i))^2$ with respect to w .

$$\begin{aligned} \frac{\partial \sum_{i=1}^N (1/2)(y_i - \log(wx_i))^2}{\partial w} &= \sum_{i=1}^N \frac{\partial (1/2)(y_i - \log(wx_i))^2}{\partial w} \\ &= \sum_{i=1}^N (2(1/2)(y_i - \log(wx_i)))(-(x_i/(wx_i))) = 0 \\ \sum_{i=1}^N (y_i - (\log(w^*) + \log(x_i))) &= 0 \Rightarrow N \log(w^*) = \sum_{i=1}^N (y_i - \log(x_i)) \\ w^* &= \exp \left(\left(\sum_{i=1}^N (y_i - \log(x_i)) \right) / N \right) \end{aligned}$$

Question 3:

Given a multilayer perceptron with one input, one tanh hidden unit, and one sigmoid output unit, derive the weight update equations to minimize the cross-entropy using gradient-descent.



$$\begin{aligned} \frac{\partial \tanh(a)}{\partial a} &= (1 - \tanh(a))^2 \\ \frac{\partial \text{sigmoid}(a)}{\partial a} &= \text{sigmoid}(a)(1 - \text{sigmoid}(a)) \\ z_i &= \tanh(wx_i + w_0) \\ \hat{y}_i &= \text{sigmoid}(vz_i + v_0) \\ \text{Error}_i &= -y_i \log \hat{y}_i - (1 - y_i) \log(1 - \hat{y}_i) \end{aligned}$$

$$\begin{aligned} \Delta v &= -\eta \frac{\partial \text{Error}_i}{\partial \hat{y}_i} \frac{\partial \hat{y}_i}{\partial v} \\ &= \eta (y_i - \hat{y}_i) z_i \\ \Delta v_0 &= -\eta \frac{\partial \text{Error}_i}{\partial \hat{y}_i} \frac{\partial \hat{y}_i}{\partial v_0} \\ &= \eta (y_i - \hat{y}_i) \\ \Delta w &= -\eta \frac{\partial \text{Error}_i}{\partial \hat{y}_i} \frac{\partial \hat{y}_i}{\partial z_i} \frac{\partial z_i}{\partial w} \\ &= \eta (y_i - \hat{y}_i) v (1 - z_i)^2 x_i \\ \Delta w_0 &= -\eta \frac{\partial \text{Error}_i}{\partial \hat{y}_i} \frac{\partial \hat{y}_i}{\partial z_i} \frac{\partial z_i}{\partial w_0} \\ &= \eta (y_i - \hat{y}_i) v (1 - z_i)^2 \end{aligned}$$

Question 4:

We know that an $N \times N$ symmetric real matrix \mathbf{K} is said to be positive semidefinite if $\mathbf{a}^\top \mathbf{K} \mathbf{a} \geq 0$ for all \mathbf{a} in \mathbb{R}^N .

- (a) Show that $k(\mathbf{x}_i, \mathbf{x}_j) = \mathbf{x}_i^\top \mathbf{x}_j$ produces a positive semidefinite kernel matrix on a given set of N data points, i.e., $\mathcal{X} = \{\mathbf{x}_i \in \mathbb{R}^D\}_{i=1}^N$.
- (b) Show that $k(\mathbf{x}_i, \mathbf{x}_j) = \frac{\mathbf{x}_i^\top \mathbf{x}_j}{\|\mathbf{x}_i\|_2 \|\mathbf{x}_j\|_2}$ produces a positive semidefinite kernel matrix on a given set of N data points, i.e., $\mathcal{X} = \{\mathbf{x}_i \in \mathbb{R}^D\}_{i=1}^N$.

- (a) $k(\mathbf{x}_i, \mathbf{x}_j) = \mathbf{x}_i^\top \mathbf{x}_j$ produces the following kernel matrix:

$$\mathbf{K} = \begin{bmatrix} \mathbf{x}_1^\top \mathbf{x}_1 & \mathbf{x}_1^\top \mathbf{x}_2 & \dots & \mathbf{x}_1^\top \mathbf{x}_N \\ \mathbf{x}_2^\top \mathbf{x}_1 & \mathbf{x}_2^\top \mathbf{x}_2 & \dots & \mathbf{x}_2^\top \mathbf{x}_N \\ \vdots & \vdots & \ddots & \vdots \\ \mathbf{x}_N^\top \mathbf{x}_1 & \mathbf{x}_N^\top \mathbf{x}_2 & \dots & \mathbf{x}_N^\top \mathbf{x}_N \end{bmatrix} = \underbrace{\begin{bmatrix} \mathbf{x}_1^\top \\ \mathbf{x}_2^\top \\ \vdots \\ \mathbf{x}_N^\top \end{bmatrix}}_{\mathbf{X}} \underbrace{\begin{bmatrix} \mathbf{x}_1^\top \\ \mathbf{x}_2^\top \\ \vdots \\ \mathbf{x}_N^\top \end{bmatrix}^\top}_{\mathbf{X}^\top}$$

$$\begin{aligned} \mathbf{a}^\top \mathbf{K} \mathbf{a} &\geq 0 \\ \mathbf{a}^\top \mathbf{X} \mathbf{X}^\top \mathbf{a} &\geq 0 \\ (\mathbf{X}^\top \mathbf{a})^\top \mathbf{X}^\top \mathbf{a} &\geq 0 \\ \|\mathbf{X}^\top \mathbf{a}\|_2^2 &\geq 0 \end{aligned}$$

- (b) $k(\mathbf{x}_i, \mathbf{x}_j) = \frac{\mathbf{x}_i^\top \mathbf{x}_j}{\|\mathbf{x}_i\|_2 \|\mathbf{x}_j\|_2}$ produces the following kernel matrix:

$$\mathbf{K} = \begin{bmatrix} \frac{\mathbf{x}_1^\top \mathbf{x}_1}{\|\mathbf{x}_1\|_2 \|\mathbf{x}_1\|_2} & \frac{\mathbf{x}_1^\top \mathbf{x}_2}{\|\mathbf{x}_1\|_2 \|\mathbf{x}_2\|_2} & \dots & \frac{\mathbf{x}_1^\top \mathbf{x}_N}{\|\mathbf{x}_1\|_2 \|\mathbf{x}_N\|_2} \\ \frac{\mathbf{x}_2^\top \mathbf{x}_1}{\|\mathbf{x}_2\|_2 \|\mathbf{x}_1\|_2} & \frac{\mathbf{x}_2^\top \mathbf{x}_2}{\|\mathbf{x}_2\|_2 \|\mathbf{x}_2\|_2} & \dots & \frac{\mathbf{x}_2^\top \mathbf{x}_N}{\|\mathbf{x}_2\|_2 \|\mathbf{x}_N\|_2} \\ \vdots & \vdots & \ddots & \vdots \\ \frac{\mathbf{x}_N^\top \mathbf{x}_1}{\|\mathbf{x}_N\|_2 \|\mathbf{x}_1\|_2} & \frac{\mathbf{x}_N^\top \mathbf{x}_2}{\|\mathbf{x}_N\|_2 \|\mathbf{x}_2\|_2} & \dots & \frac{\mathbf{x}_N^\top \mathbf{x}_N}{\|\mathbf{x}_N\|_2 \|\mathbf{x}_N\|_2} \end{bmatrix} = \underbrace{\begin{bmatrix} \frac{\mathbf{x}_1^\top}{\|\mathbf{x}_1\|_2} \\ \frac{\mathbf{x}_2^\top}{\|\mathbf{x}_2\|_2} \\ \vdots \\ \frac{\mathbf{x}_N^\top}{\|\mathbf{x}_N\|_2} \end{bmatrix}}_{\mathbf{U}} \underbrace{\begin{bmatrix} \frac{\mathbf{x}_1^\top}{\|\mathbf{x}_1\|_2} \\ \frac{\mathbf{x}_2^\top}{\|\mathbf{x}_2\|_2} \\ \vdots \\ \frac{\mathbf{x}_N^\top}{\|\mathbf{x}_N\|_2} \end{bmatrix}^\top}_{\mathbf{U}^\top}$$

$$\begin{aligned} \mathbf{a}^\top \mathbf{K} \mathbf{a} &\geq 0 \\ \mathbf{a}^\top \mathbf{U} \mathbf{U}^\top \mathbf{a} &\geq 0 \\ (\mathbf{U}^\top \mathbf{a})^\top \mathbf{U}^\top \mathbf{a} &\geq 0 \\ \|\mathbf{U}^\top \mathbf{a}\|_2^2 &\geq 0 \end{aligned}$$

Question 5:

The XOR problem is given by:

\mathbf{x}_i	x_{i1}	x_{i2}	y_i
\mathbf{x}_1	-1	-1	-1
\mathbf{x}_2	+1	-1	+1
\mathbf{x}_3	-1	+1	+1
\mathbf{x}_4	+1	+1	-1

If we use the following mapping function $\Phi(\cdot)$ for the support vector machine formulation, can we solve this binary classification problem successfully? Justify your answer.

$$\Phi(\mathbf{x}_i) = \begin{bmatrix} x_{i1}^2 - x_{i2}^2 \\ x_{i1}x_{i2} \\ x_{i1}^2 + x_{i2}^2 \end{bmatrix}$$

The mapping function produces the following representations: $\Phi(\mathbf{x}_1) = [0 \ 1 \ 2]^\top$, $\Phi(\mathbf{x}_2) = [0 \ -1 \ 2]^\top$, $\Phi(\mathbf{x}_3) = [0 \ -1 \ 2]^\top$, and $\Phi(\mathbf{x}_4) = [0 \ 1 \ 2]^\top$.

The kernel matrix can be constructed as follows:

$$\mathbf{K} = \begin{bmatrix} 5 & 3 & 3 & 5 \\ 3 & 5 & 5 & 3 \\ 3 & 5 & 5 & 3 \\ 5 & 3 & 3 & 5 \end{bmatrix}$$

$$\text{maximize} \quad J = \alpha_1 + \alpha_2 + \alpha_3 + \alpha_4 - \frac{1}{2} \begin{bmatrix} \alpha_1 \\ \alpha_2 \\ \alpha_3 \\ \alpha_4 \end{bmatrix}^\top \begin{bmatrix} 5 & -3 & -3 & 5 \\ -3 & 5 & 5 & -3 \\ -3 & 5 & 5 & -3 \\ 5 & -3 & -3 & 5 \end{bmatrix} \begin{bmatrix} \alpha_1 \\ \alpha_2 \\ \alpha_3 \\ \alpha_4 \end{bmatrix}$$

with respect to $\alpha_1, \alpha_2, \alpha_3, \alpha_4$

subject to $-\alpha_1 + \alpha_2 + \alpha_3 - \alpha_4 = 0$

$\alpha_i \geq 0 \quad i \in \{1, 2, 3, 4\}$

$$\begin{aligned} \frac{\partial J}{\partial \alpha_1} &= \frac{\partial J}{\partial \alpha_4} = 1 - 5\alpha_1 + 3\alpha_2 + 3\alpha_3 - 5\alpha_4 = 0 \\ \frac{\partial J}{\partial \alpha_2} &= \frac{\partial J}{\partial \alpha_3} = 1 + 3\alpha_1 - 5\alpha_2 - 5\alpha_3 + 3\alpha_4 = 0 \end{aligned}$$

For example, $\alpha_1^* = 1/2$, $\alpha_2^* = 1/2$, $\alpha_3^* = 0$, and $\alpha_4^* = 0$ is one of the optimum solutions.

Question 6:

Recall the error function for k -means clustering with K clusters, data points $\mathbf{x}_1, \dots, \mathbf{x}_N$, and centers $\hat{\boldsymbol{\mu}}_1, \dots, \hat{\boldsymbol{\mu}}_K$:

$$E = \sum_{i=1}^N \sum_{k=1}^K b_{ik} \|\mathbf{x}_i - \hat{\boldsymbol{\mu}}_k\|_2^2$$

where b_{ik} is equal to 1 if data point \mathbf{x}_i is closer to center $\hat{\boldsymbol{\mu}}_k$ than to any other center and to 0 otherwise.

- Instead of updating $\{\hat{\boldsymbol{\mu}}_k\}_{k=1}^K$ by computing the means, let us minimize E with batch gradient descent while holding $\{b_{ik}\}_{i=1, k=1}^{N, K}$ fixed. Derive the update formula for $\hat{\boldsymbol{\mu}}_1$ with learning rate η .
- Derive the update formula for $\hat{\boldsymbol{\mu}}_1$ with gradient descent on a single data point \mathbf{x}_i . Use learning rate η .

- (c) Recall that in the update step of the standard algorithm, we assign each cluster center to the mean of the data points closest to that center. It turns out that a particular choice of the learning rate η , which may be different for each cluster, makes the two algorithms (batch gradient descent and the standard k -means algorithm) have identical update steps. Let us focus on the update for the first cluster, with center $\hat{\mu}_1$. Calculate the value of η so that both algorithms perform the same update for $\hat{\mu}_1$.

(a)

$$\frac{\partial E}{\partial \hat{\mu}_1} = \frac{\partial \left(\sum_{i=1}^N b_{i1} \|\mathbf{x}_i - \hat{\mu}_1\|_2^2 + \text{constant} \right)}{\partial \hat{\mu}_1} = -2 \sum_{i=1}^N b_{i1} (\mathbf{x}_i - \hat{\mu}_1)$$

Therefore the update formula is $\hat{\mu}_1^{(t+1)} = \hat{\mu}_1^{(t)} + 2\eta \sum_{i=1}^N b_{i1} (\mathbf{x}_i - \hat{\mu}_1^{(t)})$

(b)

$$\frac{\partial E_i}{\partial \hat{\mu}_1} = \frac{\partial (b_{i1} \|\mathbf{x}_i - \hat{\mu}_1\|_2^2 + \text{constant})}{\partial \hat{\mu}_1} = -2b_{i1} (\mathbf{x}_i - \hat{\mu}_1)$$

Therefore the update formula is $\hat{\mu}_1^{(t+1)} = \hat{\mu}_1^{(t)} + 2\eta b_{i1} (\mathbf{x}_i - \hat{\mu}_1^{(t)})$

- (c) In the standard k -means algorithm, we assign $\hat{\mu}_1^{(t+1)} = \frac{\sum_{i=1}^N b_{i1} \mathbf{x}_i}{\sum_{i=1}^N b_{i1}}$.

$$\frac{\sum_{i=1}^N b_{i1} \mathbf{x}_i}{\sum_{i=1}^N b_{i1}} = \hat{\mu}_1^{(t)} + 2\eta \sum_{i=1}^N b_{i1} (\mathbf{x}_i - \hat{\mu}_1^{(t)}) \Rightarrow \eta = \frac{1}{2 \sum_{i=1}^N b_{i1}}$$

Question 7:

One may be concerned that the randomness introduced in random forests may cause trouble, for instance, some features or samples may not be considered at all.

- (a) Consider N training samples in a feature space of D dimensions. Consider building a random forest with T binary trees, each having exactly H internal nodes. Let F be the number of features randomly selected at each node. In order to simplify our calculations, we will let $F = 1$. For this setting, compute the probability that a certain feature (say, the first feature) is never considered for splitting.
- (b) Now let us investigate the concern regarding the random selection of the samples. Suppose each tree employs N bootstrapped training samples. Compute the probability that a particular sample (say, the first sample) is never considered in any of the trees.
- (c) Compute the values of the probabilities you obtained in the previous two parts for the case when there are $N = 2$ training samples, $D = 2$ dimensions, $T = 10$ trees of depth $H = 4$. What conclusions can you draw from your answer with regard to the concern mentioned at the beginning of the question?

- (a) The probability that it is not considered for splitting in a particular node of a particular tree is $(1 - 1/D)$. The subsampling of $F = 1$ features at each node is independent of all others. There are a total of TH nodes and hence the final answer is $(1 - 1/D)^{TH}$.
- (b) The probability that it is not considered in one of the trees is $(1 - 1/N)^N$. Since the choice for every tree is independent, the probability that it is not considered in any of the trees is $(1 - 1/N)^{NT}$.
- (c) $(1/2)^{40}$ and $(1/2)^{20}$. It is quite unlikely that a feature or a sample will be missed.

Question 8:

It is suggested to run the k -means clustering algorithm multiple times with different initializations. Explain the reasoning behind this suggestion.

The k -means clustering algorithm is heavily affected by the initial configuration. That is why it is always a good idea to run the algorithm multiple times and pick the best solution.

Question 9:

Let a configuration of the k -means clustering algorithm correspond to the k way partition (on the set of instances to be clustered) generated by the clustering at the end of each iteration. Is it possible for the k -means algorithm to revisit a configuration? Justify your answer.

It is guaranteed that the objective function value of the k -means clustering algorithm is monotonically decreasing during successive iterations. That is why it is not possible to revisit a configuration once you move to another but a better configuration.

Question 10:

What are the distance and linkage functions used in hierarchical clustering algorithms? How do they differ in their roles?

Distance function: Distance measure between two instances \mathbf{x}_i and \mathbf{x}_j .

Linkage function: Distance measure between two groups \mathcal{G}_A and \mathcal{G}_B .

Question 11:

What is an unstable learner? Why does bagging (i.e., bootstrap aggregating) rely on having an unstable learner as the base classifier?

A learning algorithm is unstable if small changes in the training set causes a large difference in the generated learner.

Bagging trains base-learners on similar data sets and, in order to generate diverse learners on these data sets, we need to have unstable learners.

Question 12:

What are the advantages of K -fold cross-validation and 5×2 cross-validation on each other?

K -fold cross-validation:

- + Training set size is large
- Overlap between training sets is large

5×2 cross-validation:

- + Overlap between training sets is small
- Training set size is small

Question 13:

For classification algorithms, there are different performance metrics such as classification accuracy, area under the ROC curve, precision, recall, etc. Explain when it is not a good idea to use the classification accuracy.

For heavily imbalanced data sets, it is not a good idea to use the classification accuracy as the performance metric since returning the majority class for all test samples will give a very high classification accuracy.