

CS 181 Spring 2024 Section 2 Notes:

Probabilistic Regression, Classification

1 Probabilistic Regression (Review)

The idea behind probabilistic regression is to assume that there is a “story” for how the data were created. Informed assumptions allow for more rigorous model selection. For a model parameterized by θ , the **likelihood** of the data $D = \{(\mathbf{x}_i, y_i)\}_{i=1}^N$, $\mathbf{x}_i \in \mathbb{R}^m$, $y_i \in \mathbb{R}$ appearing, given the specific parameter θ is defined as:

$$L(\theta|D) = p(D|\theta), \text{ which is often also written as } f(D|\theta).$$

1. Note that θ may contain *multiple* elements!
2. The value of θ that maximizes the likelihood is called the maximum likelihood estimate or **MLE**. We find the MLE by taking the derivative of the likelihood (or log-likelihood, see below), setting it equal to 0, and solving for θ . Rigorously, we also should verify the second derivative to be negative, but in practice, this may be occasionally omitted.
3. For those of you who have taken Stat 110 and/or Stat 111, $f(D|\theta)$ is also known as the “joint PDF” of the data. However, in statistics and machine learning, there is a slight difference between the joint PDF $f(D|\theta)$ and the likelihood function $L(\theta|D)$, *even though they are mathematically identical*: the likelihood function is interpreted as a function of the parameter(s) θ , while the joint PDF is interpreted as a function of the observed data D .
4. If we want, we can further add a generative story for θ (and make it random). This is called the **prior distribution**, or just **prior** of θ : $p(\theta)$

Now, onto probabilistic regression itself:

1. For probabilistic regression, we modeled the conditional distribution, $p(y|\mathbf{x})$, with a target value y_i Normally distributed with mean $\mathbf{w}^\top \mathbf{x}_i$ and variance σ^2 .
2. In last week’s section, we showed that finding parameters \mathbf{w} that minimize the negative log-likelihood (see below for explanation) of labels \mathbf{y} given design matrix \mathbf{X} gives the same expression for the optimal parameters \mathbf{w}^* as from using ordinary least squares regression.
3. Later we will also see a “full Bayes” approach where we also reason about priors on the parameters θ .

A quick addendum: what (else) can we do with our generative model?

1. Based on the generative model, and using Bayes’ rule, we can find the **posterior distribution** for θ

$$p(\theta|D) = \frac{p(D|\theta)p(\theta)}{p(D)} \propto p(D|\theta)p(\theta)$$

Note that $L(\theta|D)$ is NOT the same thing as $p(\theta|D)$. This may seem confusing, but just remember that $L(\theta|D) = p(D|\theta)$! So, in other words, we have the following:

$$p(\theta|D) \propto L(\theta|D)p(\theta)$$

Again, this is confusing! Please read this line again. Some of you may know this distinction as the Bayesian vs. Frequentist schools of statistics.

2. The value θ that maximizes the **posterior distribution** is called the maximum a posteriori or **MAP** estimate.
3. A good way to remember the **posterior distribution** is the saying: “posterior is *proportional* to likelihood times prior.”

Remarks on Taking the Log:

1. When we are maximizing the likelihood or the posterior, often, we apply log on both sides of the equation so that we are then maximizing the **log-likelihood** or **log-posterior**.
2. Because log is a monotonically-increasing function, the value of θ that maximizes the log-likelihood (or log-posterior) also maximizes the original likelihood (or posterior, if we’re working with the posterior).
3. Practically, this is helpful because the *log* operation turns products into sums, which are easier to take the derivative of (because the product rule is not the most pleasant thing to do).
4. The **log-likelihood function** will often be denoted $\ell(\theta|D)$

2 What is Classification?

We previously explored ways of predicting a continuous, real-number target, whether it was through KNN, Kernelized Regression, or Ordinary Least-Squares (OLS.) In this section, we’re going to focus on a different problem: our target output is now going to be discrete-valued. Our “y’s” now will come as discrete classes, such as predicting star types or filtering emails for spam. This type of problem, one where we make a prediction by choosing between finite class options, is known as **classification**.

3 Linear Classification

3.0.1 Classification

1. Goal : Given an input vector x , assign it to one of K discrete classes C_k .
2. Strategy: Divide our input space into *disjoint* (i.e., no overlap) **decision regions** whose boundaries are called **decision boundaries** or **decision surfaces**.
3. Note: each decision region corresponds to being assigned to a certain class: there should be K decision regions if we are working with K discrete classes.

3.0.2 Binary Linear Classification

- We are working with two classes divided by a linear separator in our feature space. We will denote the two classes as -1 and 1 (note that in other situations, we might use 0 and 1).

- Note that linear in this sense is *not* limited to the 2D case. Formally, if each data point has D dimensions, then the linear separator dividing our two classes (also called a “hyperplane”) has $D - 1$ dimensions. For example, if each data point has 3 dimensions, then the linear separator / hyperplane is a 2D plane.

- **Discriminant function** : Function that directly assigns each vector to a specific class

$$\hat{y} = \text{sign}(h(\mathbf{x}; \mathbf{w}, w_0)) = \text{sign}(\mathbf{w}^\top \mathbf{x} + w_0)$$

*note: $\text{sign}(z) = 1$ if $z \geq 0$, and $\text{sign}(z) = -1$ if $z < 0$.

- \mathbf{w} is orthogonal to every point on the decision surface. It determines the orientation of the decision boundary.

3.0.3 Perceptron (Optional)

- Perceptron is a discriminative algorithm for binary classification that finds a linear decision boundary surface, if one exists.
- To define the loss, we use the hinge loss / rectified linear function, also called *ReLU*:

$$\text{ReLU}(z) = \max\{0, z\}$$

- We define the Perceptron loss function as follows, with $h(\mathbf{x}_i; \mathbf{w}, w_0) = \mathbf{w}^\top \mathbf{x}_i + w_0$ and correct data label y_i . Recall that for any misclassified point, the value of $-h(\mathbf{x}_i; \mathbf{w}, w_0)y_i$ will be positive, so minimizing the loss learns an optimal decision boundary.

$$\begin{aligned} \mathcal{L}(\mathbf{w}) &= \sum_{i=1}^n \text{ReLU}(-h(\mathbf{x}_i; \mathbf{w}, w_0)y_i) \\ &= - \sum_{i=1: y_i \neq \hat{y}_i}^n (\mathbf{w}^\top \mathbf{x}_i + w_0)y_i \end{aligned}$$

- We can find our optimal weights by updating using (stochastic) gradient descent. Below is the equation for updating the weights \mathbf{w} at time t using the i^{th} data point, where η is the learning rate and y_i is the correct class label. Note that Perceptron updates its weights one data point at a time.

$$\mathbf{w}^{(t+1)} = \mathbf{w}^{(t)} - \eta \frac{\partial}{\partial \mathbf{w}} \mathcal{L}^{(i)}(\mathbf{w}) = \mathbf{w}^{(t)} + \eta y_i \mathbf{x}_i,$$

- Note: This update rule is equivalent to setting our the mini-batch size to 1 in Stochastic Gradient Descent, we get this particular update rule.

3.1 Concept Question

Why do we choose the *ReLU* function over the 0/1 function when formulating the loss function?

Solution

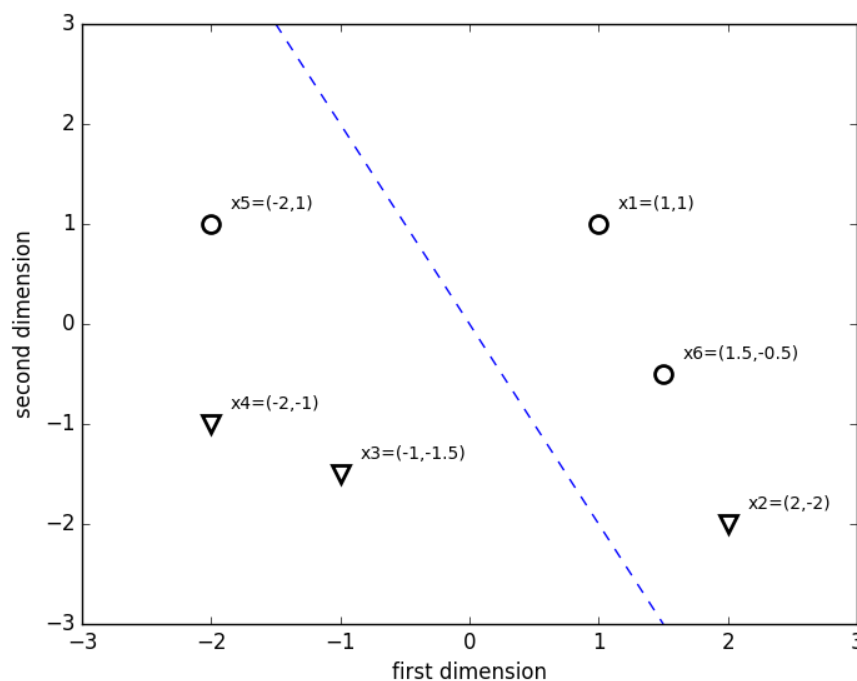
The gradient for the 0/1 function is uninformative. We are either right or wrong. 0/1 is not differentiable at 0 and its derivative is 0 at all other points.

3.2 Exercise: Small Perceptron Example

Let's train a perceptron on a small data set. Consider data $\{\mathbf{x}_i\}_{i=1}^N, \mathbf{x}_i \in \mathbb{R}^2$. Let the learning rate $\eta = 0.2$ and let the weights be initialized as:

$$\mathbf{w} = \begin{pmatrix} w_1 \\ w_2 \end{pmatrix} = \begin{pmatrix} 1 \\ 0.5 \end{pmatrix}, w_0 = 0$$

Let the circles have $y_i = 1$ and the triangles $y_i = -1$. The data and initial separation boundary (determined by \mathbf{w}) is illustrated below.



Proceed by iterating over each example until there are no more classification errors. When in doubt, refer to the notes above. We know a priori that we will be able to train the classifier and have no classification errors because one can see visually that the data is linearly separable (note: as mentioned above, if the data were not so obviously linearly separable, a new basis could make it so). How many updates do you have to make? Is this surprising?

Solution

1. Consider \mathbf{x}_1 : $\mathbf{w}^\top \mathbf{x}_1 + w_0 = 1 \cdot 1 + 0.5 \cdot 1 + 0 > 0$. This is a correct classification, so we take no action.
2. Consider \mathbf{x}_2 : $\mathbf{w}^\top \mathbf{x}_2 + w_0 = 1 \cdot 2 + 0.5 \cdot (-2) + 0 > 0$. This is an incorrect classification, so we need to update our weight parameters:

$$\begin{aligned}\mathbf{w} &\leftarrow \mathbf{w} + (0.2)(-1) \begin{pmatrix} 2 \\ -2 \end{pmatrix} = \begin{pmatrix} 0.6 \\ 0.9 \end{pmatrix} \\ w_0 &\leftarrow w_0 + 0.2(-1) = -0.2\end{aligned}$$

3. Consider \mathbf{x}_3 : $\mathbf{w}^\top \mathbf{x}_3 + w_0 = 0.6 \cdot (-1) + 0.9 \cdot (-1.5) + (-0.2) < 0$. This is a correct classification.
4. Consider \mathbf{x}_4 : Check that this is a correct classification.
5. Consider \mathbf{x}_5 : Check that this is an incorrect classification and our weight parameters are updated to:

$$\begin{aligned}\mathbf{w} &\leftarrow \begin{pmatrix} 0.2 \\ 1.1 \end{pmatrix} \\ w_0 &\leftarrow 0\end{aligned}$$

6. Consider \mathbf{x}_6 : Check that this is again an incorrect classification:

$$\begin{aligned}\mathbf{w} &\leftarrow \begin{pmatrix} 0.5 \\ 1 \end{pmatrix} \\ w_0 &\leftarrow 0.2\end{aligned}$$

All data points are correctly classified now. If the data is linearly separable, we are guaranteed to converge to a solution using the perceptron algorithm in a finite number of steps.

4 Probabilistic Classification

There are two distinct probabilistic approaches to modeling: discriminative or generative. We describe both below, and give an example for each.

4.1 Takeaways

4.1.1 Probabilistic Discriminative Model

1. In general, our goal with probabilistic discriminative modeling is to model $p(y|\mathbf{x})$.
2. Intuitively, and importantly, this means that we do not care about how \mathbf{x} is generated – we just care about the following: *given* \mathbf{x} , what is the distribution of y ?

A specific type of probabilistic discriminative modeling in the *binary case* is **logistic regression**.

4.1.2 Logistic Regression

1. In binary logistic regression, we only have *two* classes, which we will denote as 0 and 1. Note that we are *not* using -1 and 1 anymore!
2. We will model our probability distribution for the label of a certain data point y , given its features \mathbf{x} , as follows, for some weights \mathbf{w} and intercept w_0 :

$$p(y = 1|\mathbf{x}) = \sigma(\mathbf{w}^T \mathbf{x} + w_0)$$

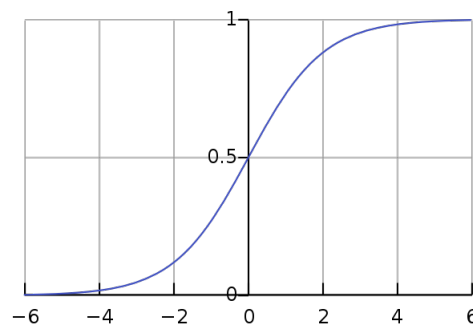
$$p(y = 0|\mathbf{x}) = 1 - \sigma(\mathbf{w}^T \mathbf{x} + w_0)$$

*In some texts, we might just see $\mathbf{w}^T \mathbf{x}$ instead of $\mathbf{w}^T \mathbf{x} + w_0$, because of the bias trick. They mean the same thing.

3. The σ denotes the **sigmoid function**, which is defined as:

$$\sigma(z) = \frac{1}{1 + \exp(-z)}$$

The **sigmoid function** is important because it takes any value z on the real line (i.e., \mathbb{R}) and returns an output on $(0, 1)$. See below for an image. This is very important because probabilities must be between 0 and 1. For clarity, $\exp(-z) = e^{-z}$.



4. To find the best weights \mathbf{w} , we need to set up a loss function. We will use what is called the **negative log-likelihood** loss function.

- (a) Intuitively, we want to find \mathbf{w} that maximizes the *likelihood* of our data.
- (b) However, unlike ordinary least squares linear regression from last week, there is no clean-cut analytical solution. Thus, we have to use gradient descent.
- (c) The problem is – gradient descent is used to minimize a function. Well, minimizing the negative likelihood (i.e., likelihood $\times -1$) is the same thing as maximizing the likelihood. Furthermore, we know that minimizing the log-likelihood is the same thing as minimizing the likelihood, but more mathematically tractable.
- (d) By definition, a “loss function” is something that we want to minimize when trying to find our optimal weights \mathbf{w} . Thus, we use the negative log-likelihood as our loss function.
- (e) Power trick: Suppose we have 1 data point with a label $y \in \{0, 1\}$. Its likelihood could be expressed as $p^y(1 - p)^{1-y}$, which evaluates to p if $y = 1$ and $(1 - p)$ if $y = 0$. So, the likelihood function essentially captures the probability of the correct class for each data point. You can employ this trick to construct likelihood expressions.

$$p(y|\mathbf{x}) = p(y = 1|\mathbf{x})^y \cdot p(y = 0|\mathbf{x})^{1-y}$$

- (f) With a training data set of N points of the form (\mathbf{x}_i, y_i) , our negative log-likelihood loss (which, fun fact, is also sometimes called the “cross-entropy loss”) is defined as follows:

$$\mathcal{L}(\theta) = - \sum_{n=1}^N (y_n \ln p(y_n = 1|\mathbf{x}_n; \theta) + (1 - y_n) \ln p(y_n = 0|\mathbf{x}_n; \theta))$$

- (g) After fitting our model, if we want to predict the class y^* for a new data point \mathbf{x}^* , we will calculate the following class probabilities, and assign this new data point to which ever class has the higher probability.

$$p(y^* = 1|\mathbf{x}^*; \mathbf{w})$$

$$p(y^* = 0|\mathbf{x}^*; \mathbf{w}) = 1 - p(y^* = 1|\mathbf{x}^*; \mathbf{w})$$

- (h) Because of the $\mathbf{w}^T \mathbf{x}$, logistic regression has *linear* decision boundaries! Yes, the sigmoid function isn’t a straight line/hyperplane, but the $\mathbf{w}^T \mathbf{x}$ ensures that we have a linear decision boundary!
- (i) Remarks on Notation:

- i. In some texts, you may see a \hat{y}_i term. Don’t be scared! In the context of logistic regression,

$$\hat{y}_i = p(y_i = 1|\mathbf{x}_i) = \sigma(\mathbf{w}^T \mathbf{x} + w_0), \text{ or with the bias trick, just } \sigma(\mathbf{w}^T \mathbf{x}).$$

- ii. Instead of $y_i = 1$, in some course materials, you might sometimes see $y_i = C_1$. They mean literally the same thing. Analogously, C_2 corresponds to class 0.
- iii. In the expression for the negative log-likelihood loss above, θ refers to the parameters of our model. In the context of logistic regression, θ and \mathbf{w} (with maybe w_0) mean the same thing.
- iv. \mathcal{L} refers to the **loss function**, while L refers to the likelihood, and ℓ refers to the log-likelihood. Be sure to check the context in which these symbols are used!
- (j) So far, we've been classifying inputs into 2 classes. We can also extend logistic regression to the multi-class case using the softmax function. See pg. 42 in *Undergraduate Fundamentals of Machine Learning* for a deeper treatment of this extension.

4.1.3 Generative Model

While a **discriminative model** directly model the posterior $p(y|\mathbf{x})$ and learn a direct map from inputs x to class labels, a **generative model** incorporates the data generation process into its model via the *joint distribution* $p(\mathbf{x}, y)$ of the class y and the input data point \mathbf{x} . They make their predictions by using Bayes rule to calculate $p(y|x)$, and then pick the most likely label y . By Bayes' Rule, we know that

$$p(\mathbf{x}, y) \propto p(\mathbf{x}|y)p(y)$$

- $p(y)$ is called the **class prior** and is almost always a **categorical distribution**, and is usually just a Bernoulli distribution in the case of binary classification (classes 0 and 1, in this context).
- A **categorical distribution** (*Cat*) is a generalization of the Bernoulli distribution.
 1. If X is a categorical random variable, we write $X \sim \text{Cat}(\mathbf{x}, \boldsymbol{\pi})$ with parameters $\mathbf{x} = [x_1, \dots, x_k]$ and $\boldsymbol{\pi} = [\pi_1, \dots, \pi_k]$.
 2. \mathbf{x} is a vector of all the possible values that X can take on.
 3. $\boldsymbol{\pi}$ stores the probabilities of X taking on a particular value.
 4. Of course, all the elements in $\boldsymbol{\pi}$ must sum up to 1 and be nonnegative, because X must take on one of these values and probabilities are always nonnegative.
 5. Mathematically, we write:

$$P(X = x_i) = \pi_i$$

- The **class prior** (often also known as the “prior distribution of y ”) $p(y)$ gives an a priori probability of an observation being a certain class. Intuitively, this is our initial belief of the distribution of y before we observe any data.
- $p(\mathbf{x}|y)$ is called the **class-conditional distribution** and its form is model-specific. Intuitively, this tells us given y (our class assignment), how likely we are to see the corresponding \mathbf{x} features.

- We are interested in picking the class k that maximizes $p(y = k|\mathbf{x})$. Again, the following equation (Bayes' Rule) might be helpful:

$$p(y|\mathbf{x}) = \frac{p(\mathbf{x}|y)p(y)}{p(\mathbf{x})} \propto p(\mathbf{x}|y)p(y)$$

*note: depending on the model, \mathbf{x} can be either discrete or continuous. However, y must be discrete, by definition of "classes."

4.1.4 Naive Bayes

Naive Bayes is one type of generative model for classification. Suppose we have n data points x_n , each with D dimensions. So x_{nd} refers to the d^{th} entry of the n^{th} data point \mathbf{x}_n . For example, each data point x_n could be a person's features, with $d = 3$ for 3 dimensions of hair color, height, and age, where the correct class label y_n corresponds to their favorite color.

Each dimension can take on one of J values $\{1, \dots, J\}$. For example $J = 3$ if a student's hair color can be entered as blonde, black, or brown.

In Naive Bayes, we assume that each of x_{n1}, \dots, x_{nD} has some conditional distribution given y_n . Our naive assumption is that each dimension of data \mathbf{x}_n is independent, conditioned on the class. In other words, given the correct class label $y_n = C_k$, each dimension $d \in \{1, \dots, D\}$ of the n^{th} observed data point \mathbf{x}_n is **conditionally independent** from the other dimensions. This assumption can be mathematically formulated as follows.

$$p(\mathbf{x}_n|y_n = C_k) = \prod_{d=1}^D p(x_{nd}|y_n = C_k)$$

For our example, that would mean assuming that if we know a student's favorite color, then the distributions of their hair color, height, and age are independent.

We use Naive Bayes to limit the number of parameters needed to specify our model. If our features were dependent on each other, then we would need to explicitly model this dependence using additional parameters.

We model each feature x_d as a categorical distribution. The categorical distribution is a generalization of the Bernoulli distribution. For example, if there are 3 hair colors (blonde, black, and brown), then the vector $\pi_{dk} = [\pi_{dk1}, \pi_{dk2}, \pi_{dk3}] = [0.2, 0.5, 0.3]$ parameterizes x_{dk} , where we use notation $\pi_{kdj} \geq 0$ where π_{kdj} is the probability in class k of feature d taking on value j .

For sake of brevity, please see pg. 47 in *Undergraduate Fundamentals of Machine Learning* for a deeper treatment of Naive Bayes, and a toy example.

4.1.5 Naive Bayes Concept Questions

How many parameters does this model have? Why do we use the "naive" assumption?

Solution

The model has $K \times D \times (J - 1)$ parameters. The $J - 1$ arises because each Categorical distributions' J parameters must sum to 1, removing one degree of freedom. We use the “naive” assumption in order to keep the number of parameters small. Naive bayes requires a number of parameters linear in the number of variables (classes, features).

4.2 Exercise: Shapes of Decision Boundaries I

Consider now a generative model with $K > 2$ classes, and output label \mathbf{y} encoded as a “one hot” vector of length K . We adopt class prior $p(\mathbf{y} = C_k; \boldsymbol{\pi}) = \pi_k$ for all $k \in \{1, \dots, K\}$ (where π_k is a parameter of the prior). Let $p(\mathbf{x} | \mathbf{y} = C_k)$ denote the class-conditional density of features \mathbf{x} (in this case for class C_k). Let the class-conditional probabilities be Gaussian distributions

$$p(\mathbf{x} | \mathbf{y} = C_k) = \mathcal{N}(\mathbf{x} | \boldsymbol{\mu}_k, \boldsymbol{\Sigma}_k), \text{ for } k \in \{1, \dots, K\}$$

We will predict the class of a new example \mathbf{x} as the class with the highest conditional probability, $p(\mathbf{y} = C_k | \mathbf{x})$. Luckily, a little bird came to the window of your dorm, and claimed that you can classify an example \mathbf{x} by finding the class that maximizes the following function:

$$f_k(\mathbf{x}) = \log(\pi_k) - \frac{1}{2} \log(|\boldsymbol{\Sigma}_k|) - \frac{1}{2}(\mathbf{x} - \boldsymbol{\mu}_k)^T \boldsymbol{\Sigma}_k^{-1}(\mathbf{x} - \boldsymbol{\mu}_k).$$

Derive this formula by comparing two different classes' conditional probabilities. What can we claim about the shape of the decision boundary given this formula?

Solution

Let's start with the conditional probability:

$$p(\mathbf{y} = C_k | \mathbf{x}) = \frac{p(\mathbf{x} | \mathbf{y} = C_k)p(\mathbf{y} = C_k)}{\sum_{\ell} p(\mathbf{x} | \mathbf{y} = C_{\ell})p(\mathbf{y} = C_{\ell})}$$

{We can take out the denominator since it will be the same for each class}

$$\begin{aligned} &\propto p(\mathbf{x} | \mathbf{y} = C_k)p(\mathbf{y} = C_k) \\ &= \frac{1}{(2\pi)^{\frac{m}{2}} |\boldsymbol{\Sigma}_k|^{\frac{1}{2}}} \exp\left(-\frac{1}{2}(\mathbf{x} - \boldsymbol{\mu}_k)^T \boldsymbol{\Sigma}_k^{-1}(\mathbf{x} - \boldsymbol{\mu}_k)\right) \pi_k \end{aligned}$$

{We can factor out constants that will be shared across classes}

$$\propto \frac{\pi_k}{|\boldsymbol{\Sigma}_k|^{\frac{1}{2}}} \exp\left(-\frac{1}{2}(\mathbf{x} - \boldsymbol{\mu}_k)^T \boldsymbol{\Sigma}_k^{-1}(\mathbf{x} - \boldsymbol{\mu}_k)\right)$$

{We take the logarithm, which is a monotonically increasing function and won't change the maximum across classes}

$$\propto \ln(\pi_k) - \frac{1}{2} \ln(|\boldsymbol{\Sigma}_k|) - \frac{1}{2}(\mathbf{x} - \boldsymbol{\mu}_k)^T \boldsymbol{\Sigma}_k^{-1}(\mathbf{x} - \boldsymbol{\mu}_k)$$

Now, what can we say about the shape of our decision boundary? $(\mathbf{x} - \boldsymbol{\mu}_k)^T \boldsymbol{\Sigma}_k^{-1} (\mathbf{x} - \boldsymbol{\mu}_k)$ gives us intuition about the shape of our decision boundary. We see it is quadratic. For further intuition, we can look at the other two terms to see what affects the probability of an example being classified to a given class. As the prior π_k increases, we see that the probability increases, which fits our intuition. As the covariance matrix $\boldsymbol{\Sigma}_k$ increases, we see that the probability decreases, which also fits our intuition.

5 Additional Exercises

5.1 Exercise: Naive Bayes Practice Problem

Let's consider just two classes. Let $p(y = 1) = \theta$ and $p(y = 0) = 1 - \theta$. For a given dataset $\{\mathbf{x}_n, y_n\}_{n=1}^N$, what are the maximum likelihood estimates of the parameters $\theta, \{\pi_{k,d,j}\}$? We do not provide any further description of the distributions – please work with a generic distribution joint distribution $p(\mathbf{x}, y)$ (given the parameters).

Solution

For MLE, we want to find the parameters that minimize the negated log-likelihood:

$$L(\theta, \{\pi_{k,d,j}\}) = -\ln p(\{\mathbf{x}_n, y_n\}_{n=1}^N | \theta, \{\pi_{k,d,j}\}) \quad (1)$$

$$= -\ln \prod_{n=1}^N p(\mathbf{x}_n, y_n | \theta, \{\pi_{k,d,j}\}) \quad (2)$$

$$= -\ln \left(\prod_{n:y_n=1} p(\mathbf{x}_n | \{\pi_{1,d,j}\}) p(y_n = 1 | \theta) \prod_{n:y_n=0} p(\mathbf{x}_n | \{\pi_{0,d,j}\}) p(y_n = 0 | \theta) \right) \quad (3)$$

$$= -\ln \left(\prod_{n:y_n=1} p(\mathbf{x}_n | \{\pi_{1,d,j}\}) p(y_n = 1 | \theta) \right) - \ln \left(\prod_{n:y_n=0} p(\mathbf{x}_n | \{\pi_{0,d,j}\}) p(y_n = 0 | \theta) \right) \quad (4)$$

$$= -\sum_{n:y_n=1} \ln p(\mathbf{x}_n | \{\pi_{1,d,j}\}) - \sum_{n:y_n=1} \ln p(y_n = 1 | \theta) \quad (5)$$

$$- \sum_{n:y_n=0} \ln p(\mathbf{x}_n | \{\pi_{0,d,j}\}) - \sum_{n:y_n=0} \ln p(y_n = 0 | \theta) \quad (6)$$

$$= -\sum_{n:y_n=1} \ln \prod_{d=1}^D \prod_{j=1}^J \pi_{1dj}^{\mathbb{I}[x_{nd}=j]} - \sum_{n:y_n=1} \ln(\theta) \quad (7)$$

$$- \sum_{n:y_n=0} \ln \prod_{d=1}^D \prod_{j=1}^J \pi_{0dj}^{\mathbb{I}[x_{nd}=j]} - \sum_{n:y_n=0} \ln(1 - \theta) \quad (8)$$

$$= -\sum_{n:y_n=1} \sum_{d=1}^D \sum_{j=1}^J \mathbb{I}[x_{nd}=j] \ln(\pi_{1dj}) - \sum_{n:y_n=1} \ln(\theta) \quad (9)$$

$$- \sum_{n:y_n=0} \sum_{d=1}^D \sum_{j=1}^J \mathbb{I}[x_{nd}=j] \ln(\pi_{0dj}) - \sum_{n:y_n=0} \ln(1 - \theta) \quad (10)$$

The indicator function $\mathbb{I}[\cdot]$ is useful to make sure that our likelihood only evaluates $p(\mathbf{x}, y)$ for each \mathbf{x} 's true class rather than both classes. This is used a lot, so make sure you see what is going on.

5.2 Exercise: Shapes of Decision Boundaries II

Let's say the little bird comes back and now tells you that every class has the same covariance matrix, and so $\Sigma_\ell = \Sigma'_{\ell'}$ for all classes C_ℓ and $C_{\ell'}$. Simplify this formula down further. What can we claim about the shape of the decision boundaries now?

Solution

Here was our original formula.

$$\ln(\pi_k) - \frac{1}{2} \ln(|\Sigma_k|) - \frac{1}{2}(\mathbf{x} - \boldsymbol{\mu}_k)^T \Sigma_k^{-1} (\mathbf{x} - \boldsymbol{\mu}_k) \quad (11)$$

Given that all of the classes have the same covariance matrix, which we will write as Σ , we have:

$$\begin{aligned} f_k(\mathbf{x}) &= \ln(\pi_k) - \frac{1}{2} \ln(|\Sigma|) - \frac{1}{2}(\mathbf{x} - \boldsymbol{\mu}_k)^T \Sigma^{-1} (\mathbf{x} - \boldsymbol{\mu}_k) \\ &= \ln(\pi_k) - \frac{1}{2} \ln(|\Sigma|) - \frac{1}{2}(\mathbf{x} - \boldsymbol{\mu}_k)^T (\Sigma^{-1} \mathbf{x} - \Sigma^{-1} \boldsymbol{\mu}_k) \\ &= \ln(\pi_k) - \frac{1}{2} \ln(|\Sigma|) - \frac{1}{2}(\mathbf{x}^T \Sigma^{-1} \mathbf{x} - \mathbf{x}^T \Sigma^{-1} \boldsymbol{\mu}_k - \boldsymbol{\mu}_k^T \Sigma^{-1} \mathbf{x} + \boldsymbol{\mu}_k^T \Sigma^{-1} \boldsymbol{\mu}_k) \\ &= \ln(\pi_k) - \frac{1}{2} \ln(|\Sigma|) - \frac{1}{2}(\mathbf{x} \Sigma^{-1} \mathbf{x}^T - 2\mathbf{x}^T \Sigma^{-1} \boldsymbol{\mu}_k + \boldsymbol{\mu}_k^T \Sigma^{-1} \boldsymbol{\mu}_k) \\ &= \ln(\pi_k) - \frac{1}{2} \ln(|\Sigma|) - \frac{1}{2}\mathbf{x} \Sigma^{-1} \mathbf{x}^T + \mathbf{x}^T \Sigma^{-1} \boldsymbol{\mu}_k - \frac{1}{2}\boldsymbol{\mu}_k^T \Sigma^{-1} \boldsymbol{\mu}_k \end{aligned}$$

{We can drop $-\frac{1}{2} \ln(|\Sigma|)$ and $-\frac{1}{2}\mathbf{x} \Sigma^{-1} \mathbf{x}^T$ since they are class independent}

$$\propto \ln(\pi_k) + \mathbf{x}^T \Sigma^{-1} \boldsymbol{\mu}_k - \frac{1}{2}\boldsymbol{\mu}_k^T \Sigma^{-1} \boldsymbol{\mu}_k$$

Thus, we conclude that we can adopt function

$$\tilde{f}_k(\mathbf{x}) = \mathbf{x}^T \Sigma^{-1} \boldsymbol{\mu}_k - \frac{1}{2}\boldsymbol{\mu}_k^T \Sigma^{-1} \boldsymbol{\mu}_k + \ln(\pi_k) \quad (12)$$

Looking at this formula, we can see that our decision boundaries are no longer quadratic but linear. We've proven that if the classes share the same covariance matrix, our decision boundaries will be linear!

5.3 OPTIONAL: Visualizing Decision Boundaries (This Looks Cool, But *Not* Required)

If you want to better understand what these decision boundaries look like, we can visualize them! Let's consider two classes and assume x lives in 2 dimensions. We first consider the case in which the two classes have identical covariances, here defined as

$$p(x|y=1) = \mathcal{N}\left(\mu_1 = \begin{bmatrix} 1 \\ 1 \end{bmatrix}, \Sigma_1 = \begin{bmatrix} 2 & 0 \\ 0 & 1 \end{bmatrix}\right)$$
$$p(x|y=2) = \mathcal{N}\left(\mu_2 = \begin{bmatrix} -1 \\ -1 \end{bmatrix}, \Sigma_2 = \begin{bmatrix} 2 & 0 \\ 0 & 1 \end{bmatrix}\right)$$

We use the following code to plot the contours of both Gaussians:

```
# imports
import matplotlib.pyplot as plt
import numpy as np
import plotly.graph_objects as go
import scipy.stats

# create meshgrid from -8 to 8
mesh_granularity = 100
possible_vals = np.linspace(-8., 8, mesh_granularity)
mesh_coords = np.meshgrid(possible_vals, possible_vals)
mesh_coords = np.reshape(np.stack(mesh_coords),
                          newshape=(2, mesh_granularity * mesh_granularity)).T

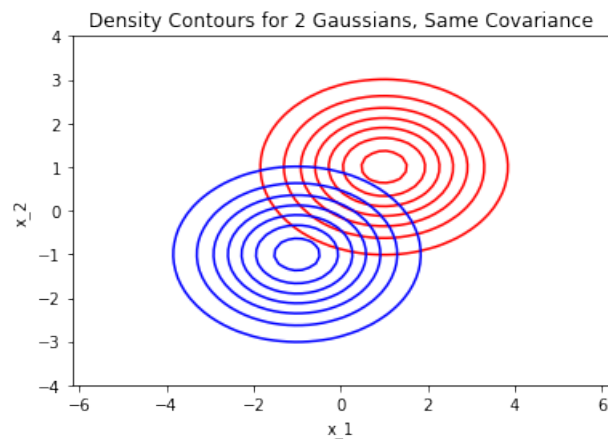
# define means of Gaussians
means = [np.array([1., 1.]),
         np.array([-1., -1.])]

# compute densities for both Gaussians assuming equal covariances
covs = [np.array([[2., 0.],
                  [0., 1.]]),
        np.array([[2., 0.],
                  [0., 1.]])]
same_cov_densities = [scipy.stats.multivariate_normal.pdf(x=mesh_coords,
                                                           mean=mean,
                                                           cov=cov)
                       for mean, cov in zip(means, covs)]

# plot
plt.contour(np.reshape(mesh_coords[:, 0],
                       newshape=(mesh_granularity, mesh_granularity)),
            np.reshape(mesh_coords[:, 1],
                       newshape=(mesh_granularity, mesh_granularity)),
            np.reshape(same_cov_densities[0],
                       newshape=(mesh_granularity, mesh_granularity)),
            colors='red')
plt.contour(np.reshape(mesh_coords[:, 0],
                       newshape=(mesh_granularity, mesh_granularity)),
            np.reshape(mesh_coords[:, 1],
                       newshape=(mesh_granularity, mesh_granularity)),
            np.reshape(same_cov_densities[1],
                       newshape=(mesh_granularity, mesh_granularity)),
            colors='blue')
```

```
plt.xlabel('x_1')
plt.ylabel('x_2')
plt.title('Density Contours for 2 Gaussians, Same Covariance')
plt.axis('equal')
plt.xlim(-4, 4)
plt.ylim(-4, 4)
plt.show()
```

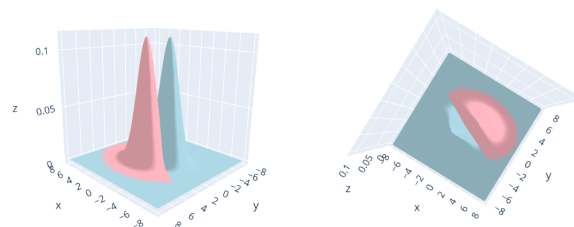
This gives us the following contours:



If you would prefer plotting in 3D, we can alternatively use Plotly:

```
fig = go.Figure(data=[
    go.Mesh3d(x=mesh_coords[:, 0],
              y=mesh_coords[:, 1],
              z=same_cov_densities[0],
              color='lightpink'),
    go.Mesh3d(x=mesh_coords[:, 0],
              y=mesh_coords[:, 1],
              z=same_cov_densities[1],
              color='lightblue')])
fig.show()
```

Rotating the plot (and ignoring peripheral floating point problems), we can see that the two Gaussians have equal density along a line:



We now plot the second case, with unequal covariances. We assume the two Gaussians are:

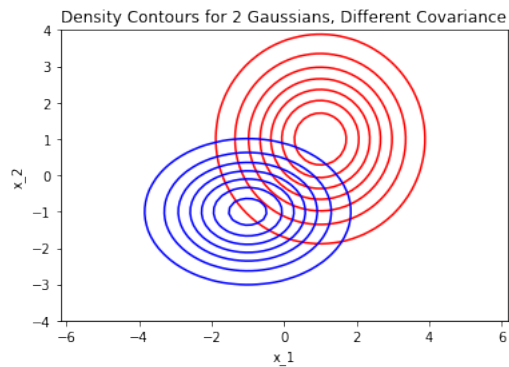
$$p(x|y=1) = \mathcal{N}\left(\mu_1 = \begin{bmatrix} 1 \\ 1 \end{bmatrix}, \Sigma_1 = \begin{bmatrix} 2 & 0 \\ 0 & 2 \end{bmatrix}\right)$$

$$p(x|y=2) = \mathcal{N}\left(\mu_2 = \begin{bmatrix} -1 \\ -1 \end{bmatrix}, \Sigma_2 = \begin{bmatrix} 2 & 0 \\ 0 & 1 \end{bmatrix}\right)$$

```
# define means of Gaussians
means = [np.array([1., 1.]),
         np.array([-1., -1.])]

# compute densities for both Gaussians assuming equal covariances
covs = [2.*np.eye(2), np.array([[2., 0.], [0., 1.]])]
diff_cov_densities = [scipy.stats.multivariate_normal.pdf(x=mesh_coords,
                                                         mean=mean,
                                                         cov=cov)
                    for mean, cov in zip(means, covs)]
```

The contour plot outputs as follows:



Rotating the plot (and ignoring the floating point problems on the periphery), we can see that the two Gaussians have equal density along a parabola:

```
fig = go.Figure(data=[
go.Mesh3d(x=mesh_coords[:, 0], y=mesh_coords[:, 1], z=diff_cov_densities[0],
          color='lightpink'),
go.Mesh3d(x=mesh_coords[:, 0], y=mesh_coords[:, 1], z=diff_cov_densities[1],
          color='lightblue')])
fig.show()
```

