Statistical learning - Vinciotti (2022)

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 $Github:\ https://github.com/StefanoCretti/StatisticalLearning.git$

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Part I Introduction

Chapter 1

General course information

1.1 Textbooks

- James et al (2021), Introduction to statistical learning in R, 2nd edition. (Book that is used as guidline for the course, but further concepts will be added during the lectures)
- Hastie et al (2001), Elements of statistical learning. (More advanced book for those who want to study the subject more in depth)

1.2 Assessment

- Three homework tasks during the course (Uploaded on moodle, two weeks of time for each, more practical and mainly focused on applying methods to some data. If done well they will add 2 points to the written exam score)
- Final written exam (More theoretical but still connected to the practical part, for instance by commenting on analysis output)

1.3 Topics

- Linear regression (Gauss, 1800) (Assumed to be already known from Statistical Learning 1)
- Linear discriminant analysis, LDA (Fisher, 1936) (Later extended to quadratic discriminant analysis, QDA)
- Logistic regression (1940s)
- Generalized linear models (Nelder and Wedderburn, 1972)
- Classification and regression trees (Breiman and Freidman, 1980s) (First introduction of computer intensive methods)
- Machine learning (1990s): support vector machines, neural networks/deep learning, unsupervised learning (clustering, PCA)
- \bullet Individual methods: theory, details, implementation \dots

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 \bullet General concept: model selection, inference, prediction \dots

Part II Statistical learning

Chapter 2

What is statistical learning?

2.1 Definition of statistical learning

In general, a statistical learning problem can be formalized as follows:

- ullet Y: response/dependent/outcome variable
- $\underline{X} = (X_1, \dots, X_p)$: vector of predictors/features/independent variables/covariates

We assume that there is a relationship between Y and \underline{X} , which can be written as:

$$Y = f(\underline{X}) + \varepsilon$$

Where:

- $f(\underline{X})$ is the deterministic (but unknown) function of the vector $\underline{X} = (X_1, \dots, X_p)$
- ε is the error (stochastic part), for which we assume the following properties:
 - $E[\varepsilon] = 0$ (Its expected value is zero)
 - $\varepsilon \perp \underline{X}$ (It is independent from \underline{X})

Therefore, the expression statistical learning encompasses different methods to estimate $f(\underline{X})$.

2.2 Why estimate f?

There are two main reasons to estimate f, those two being **prediction** and **inference**.

2.2.1 Prediction

Predict Y when we only have observations about \underline{X} . Since $E[\varepsilon] = 0$, we usually take:

$$\hat{Y} = \hat{f}(X)$$

With \hat{f} being our estimate of f.

If this is the only reason to estimate f, then \hat{f} can be a black-box method (deep learning). The accuracy of \hat{Y} as a predictor of Y can be described calculating the expected MSE (mean squared error):

$$\begin{split} E[(Y-\hat{f}(\underline{X}))^2 \,|\, \underline{X} &= \underline{x}\,] = & \text{where } \hat{f} \text{ is a fixed known function} \\ &= E[(f(\underline{X}) + \varepsilon - \hat{f}(\underline{X}))^2] & \text{since } Y = f(\underline{X}) + \varepsilon \\ &= E[(f(\underline{X}) - \hat{f}(\underline{X})) + \varepsilon)^2] & \text{rearranging} \\ &= E[(f(\underline{X}) - \hat{f}(\underline{X}))^2 + \varepsilon^2 + 2\varepsilon(f(\underline{X}) - \hat{f}(\underline{X}))] & \text{solving the square} \\ &= E[(f(\underline{X}) - \hat{f}(\underline{X}))^2] + E[\varepsilon^2] + 2E[\varepsilon(f(\underline{X}) - \hat{f}(\underline{X}))] & \text{separating the expectations} \end{split}$$

Furthermore, since we know that:

$$Var(\varepsilon) = E[(\varepsilon - E(\varepsilon))^2]$$
 formal definition of variance
$$= E[\varepsilon^2] - (E[\varepsilon])^2$$
 definition generally used during calculation
$$= E[\varepsilon^2]$$
 since $E[\varepsilon] = 0$

Thus we get:

$$\begin{split} E[(Y-\hat{f}(\underline{X}))^2 \,|\, \underline{X} &= \underline{x}\,] = \\ &= E[(f(\underline{X})-\hat{f}(\underline{X}))^2] + Var(\varepsilon) + 2E[\varepsilon(f(\underline{X})-\hat{f}(\underline{X}))] \qquad \text{substituting } E[\varepsilon^2] = Var(\varepsilon) \\ &= (f(\underline{X})-\hat{f}(\underline{X}))^2 + Var(\varepsilon) + 2(f(\underline{X})-\hat{f}(\underline{X}))E[\varepsilon] \qquad \text{since } f(\underline{X})-\hat{f}(\underline{X}) \text{ is a constant} \\ &= (f(\underline{X})-\hat{f}(\underline{X}))^2 + Var(\varepsilon) \qquad \qquad \text{since } E[\varepsilon] = 0 \end{split}$$

With:

- $(f(\underline{X}) \hat{f}(\underline{X}))^2$ being the **reducible error**. The model choice can increase or reduce this value, hence it is mostly controllable.
- $Var(\varepsilon)$ being the **irreducible error**. This value depends on the innate randomness present in the data, hence you can only try and minimize it by deciding which variables to use in your prediction (but it will never be zero otherwise you would have a deterministic situation).

2.2.2 Inference

Inference is used when you want to understand the relation between Y and \underline{X} (and not just be able to make predictions). Namely, inference answers questions such as:

- Which predictors/factors are most associated with the response?
- What is the relationship between Y and X_i ?

2.3 How to extimate f?

Given some **training data** (\underline{x}_i, y_i) , i = 1, ..., n, where $\underline{x}_i = (x_{i1}, ..., x_{ip})^t$ is the vector of observations of unit i while y_i is the response for unit i, broadly speaking there are two types of methods to estimate f: **parametric methods** and **non-parametric methods**.

2.3.1 Parametric methods

In order to use parametric methods, we make an assumption about the functional form of $f(\underline{X})$, that assumption being that the form of the function depends on some parameters (which we can estimate). An example of parametric method is **linear regression**, which implies that $f(\underline{X})$ is in the form:

$$f(\underline{X}) = \beta_0 + \beta_1 X_1 + \beta_2 X_2 + \dots + \beta_p X_p$$

Given this assumption, statistical learning becomes **fitting** (or training) the model on the data, which means estimating the parameters $\hat{\beta_0}, \hat{\beta_1}, \dots, \hat{\beta_p}$ such that:

$$\hat{f}(\underline{X}) = \hat{\beta}_0 + \hat{\beta}_1 X_1 + \hat{\beta}_2 X_2 + \dots + \hat{\beta}_p X_p$$

The main disadvantage of these methods is that they may be **too restrictive**. Notice that linear models are linear in the parameters, not in the predictors, hence:

- $f(X_1) = \beta_0 + \beta_1 X_1 + \beta_2 X_1^2 + \beta_3 X_1^3 + \varepsilon$ (polynomial regression) is a linear model.
- $f(X_1) = \beta_0 X_1^{\beta_1}$ is not a linear model.

2.3.2 Non-parametric methods

Non-parametric methods do not make any explicit assumption on the function form of $f(\underline{X})$. These methods want to estimate f by getting as close as possible to the data, without being too rough or wiggly (basically overfitting).

2.3.3 Parametric vs non-parametric methods

Despite parametric methods being more restrictive than non-parametric ones, we might still choose to adopt the former for the sake of interpretability and generalizability outside of the training data.

2.4 Bias-variance trade-off

Assume you have some data pairs in the form (x_i, y_i) , i = 1, ..., n; you can then define the estimate function $\hat{f}(x)$ as a linear model with up to n-1 parameters (more parameters give the same result as n-1 parameters) and an intercept value. You could thus use, for instance, the models:

1 parameter
$$f(x) = \beta_0 + \beta_1 x + \varepsilon$$
2 parameters
$$f(x) = \beta_0 + \beta_1 x + \beta_2 x^2 + \varepsilon$$

$$\vdots$$

$$n-1 \text{ parameters}$$

$$f(x) = \beta_0 + \beta_1 x + \beta_2 x^2 + \dots + \beta_{n-1} x^{n-1} + \varepsilon$$

You could choose the model with the highest number of parameters; this model would explain all the variance of the training data $(R^2 = 1)$ yet it would be very complex and perform badly with new data points. On the other hand a model with a lower number of parameters would explain less of the variance of the training data, yet it could perform better with new data points.

Keeping in mind that Y is random and that \hat{f} is a random variable estimated from the data, when using the general formula to determine how well a model performs at generic \underline{X} , we notice:

$$\begin{split} E[(Y-\hat{f}(\underline{X}))^2 \,|\, \underline{X} &= \underline{x}\,] = \\ &= E[(f(\underline{X}) + \varepsilon - \hat{f}(\underline{X}))^2] & \text{since } Y = f(\underline{X}) + \varepsilon \\ &= E[(f(\underline{X}) + \varepsilon - \hat{f}(\underline{X}) + E[\hat{f}(\underline{X})] - E[\hat{f}(\underline{X})])^2] & \text{since } E[\hat{f}(\underline{X})] - E[\hat{f}(\underline{X})] = 0 \\ &= E[((f(\underline{X}) - E[\hat{f}(\underline{X})]) + \varepsilon + (E[\hat{f}(\underline{X})] - \hat{f}(\underline{X})))^2] & \text{grouping} \\ &= E[(f(\underline{X}) - E[\hat{f}(\underline{X})])^2] + E[\varepsilon^2] + E[(E[\hat{f}(\underline{X})] - \hat{f}(\underline{X}))^2] + \\ &\quad + 2E[(f(\underline{X}) - E[\hat{f}(\underline{X})])(E[\hat{f}(\underline{X})] - \hat{f}(\underline{X}))] + & \text{solving the square and} \\ &\quad + 2E[(E[\hat{f}(\underline{X})] - \hat{f}(\underline{X}))\varepsilon] & \text{dividing the expectations} \end{split}$$

But notice that:

$$\begin{split} E[(f(\underline{X}) - E[\hat{f}(\underline{X})])\,\varepsilon] &= E[\varepsilon](f(\underline{X}) - E[\hat{f}(\underline{X})]) & \text{since } f(\underline{X}) - E[\hat{f}(\underline{X})] \text{ is constant} \\ &= 0 & \text{since } E[\varepsilon] &= 0 \end{split}$$

$$\begin{split} E[(f(\underline{X}) - E[\hat{f}(\underline{X})])(E[\hat{f}(\underline{X})] - \hat{f}(\underline{X}))] &= \\ &= (f(\underline{X}) - E[\hat{f}(\underline{X})])E[E[\hat{f}(\underline{X})] - \hat{f}(\underline{X})] \qquad \text{since } f(\underline{X}) - E[\hat{f}(\underline{X})] \text{ is constant} \\ &= (f(\underline{X}) - E[\hat{f}(\underline{X})])E[\hat{f}(\underline{X}) - \hat{f}(\underline{X})] \qquad \text{since } \hat{f}(\underline{X}) \text{ is a constant} \\ &= (f(\underline{X}) - E[\hat{f}(\underline{X})])E[0] \\ &= 0 \end{split}$$

$$E[(E[\hat{f}(\underline{X})] - \hat{f}(\underline{X})) \varepsilon] = E[\varepsilon]E[E[\hat{f}(\underline{X})] - \hat{f}(\underline{X})] \quad \text{since } \varepsilon \perp \underline{X} \implies E[\varepsilon\underline{X}] = E[\varepsilon]E[\underline{X}]$$
$$= 0 \quad \text{since } E[\varepsilon] = 0$$

Therefore we can simplify as:

$$\begin{split} E[(Y-\hat{f}(\underline{X}))^2 \,|\, \underline{X} &= \underline{x}\,] = E[(f(\underline{X}) - E[\hat{f}(\underline{X})])^2] &\quad + E[\varepsilon^2] &\quad + E[(E[\hat{f}(\underline{X})] - \hat{f}(\underline{X}))^2] \\ &= f(\underline{X}) - E[\hat{f}(\underline{X})])^2 &\quad + Var(\varepsilon) &\quad + Var(\hat{f}(\underline{X})) \end{split}$$

Where:

- $f(\underline{X}) E[\hat{f}(\underline{X})]^2$ is the bias
- $Var(\varepsilon)$ is the irreducible error
- $Var(\hat{f}(\underline{X}))$ is the variance of the model

Bias refers to the error that is introduced by approximating a problem, which may be extremely complicated, by a much simpler model. If an estimated model performs well on the training data but it does not perform well on new data, the estimated model has high bias. If an estimated model performs well on multiple data sets, the estimated model has low bias. High bias means that an estimated model is far from the real model.

Variance refers instead to the amount by which \hat{f} would change if we estimated it using a different training data set. Ideally, the \hat{f} calculated through different datasets should be not significantly different. The more flexible models are generally more influenced by a change of dataset.

In order to minimize the expected error, we need to achieve low bias and low variance. In practice, one needs to find a good trade-off between bias and variance, since reducing one often involves increasing the other. In general:

- A **simpe model** has high bias (it is far from the real model) and low variance (when fitting using different training data you get similar estimated parameters).
- A **complex model** has low bias and high variance.

Both in parametric and non-parametric methods, you generally have at least one **tuning parameter**, which is a parameter that can be tweaked (for instance the degree of the polinomial) in order to choose the balance between bias and variance.

Chapter 3

Statistical decision theory

3.1 Definition of statistical decision theory

Statistical decision theory is a set of quantitative methods for reaching optimal decisions for well posed problems. Statistical decision theory applies to supervised learning, while it does not apply to unsupervised learning. For most of the course we will focus on supervised learning.

3.2 Supervised learning

Supervised learning is a set of methods whose objective is, given the observations (\underline{x}_i, y_i) with $i = 1, \ldots, n$, to find a rule \hat{f} that allows us to predict Y from \underline{X} , meaning that $\hat{Y} = \hat{f}(\underline{X})$. $\hat{f}(\underline{X})$ is thus an approximation of the true rule $f(\underline{X})$. Finding \hat{f} corresponds to training (or fitting) a model using labelled data pairs (both predictors and response values are known).

Broadly, there are two main types of methods of supervised learning:

- Regression methods, in which the response Y is quantitative (numerical). In this case \hat{f} is called regression function.
- Classification methods, in which the response Y is qualitative (categorical). In this case \hat{f} is called classifier.

The predictors (\underline{X}) can take any form, numerical or categorical. In general there no assumptions on the form of the predictors (but not always).

To evaluate a model you use **loss functions**, which are functions used to penalize the differences between Y and $f(\underline{X})$. Many loss functions can be used; the choice of a specific loss function determines which function is considered to be the true rule $f(\underline{X})$.

3.3 Regression setting

In a regression setting, the loss function which is generally used is the **squared error loss function**, which is defined as

$$L(Y,f(\underline{X})) = (Y - f(\underline{X}))^2$$

When using this loss function, the criterion for choosing the model becomes finding f such that it minimizes the **expected prediction error** (EPE), which is the expected value of the squared error loss function:

$$\begin{aligned} \text{EPE}(f) &= E_{Y,\underline{X}}[(Y - f(\underline{X}))^2] \\ &= \iint (y - f(\underline{x}))^2 g_{Y,\underline{X}}(y,\underline{x}) \, dy \, d\underline{x} \end{aligned} \qquad \text{if } Y \text{ is continuous}$$

This error is an expectation since you usually do not know the distribution of Y and \underline{X} . If Y is continuous you can rewrite this expectation as the double integral times the joint density function a.

Theorem 1 (Minimum of the expected prediction error). $EPE(f) = E_{Y,\underline{X}}[(Y - f(\underline{X}))^2]$ has a minimum when $f(\underline{x}) = E[Y|\underline{X} = \underline{x}]$. $f(\underline{x})$ is called regression function.

Proof. (Sketch) By definition we know that

$$EPE(f) = E_{Y,X}[(Y - f(X))^{2}]$$

Then, since the law of iterated expectations states that E[X] = E[E[X|Y]] (analogously to profile likelihood, you fix one variable and variate the other), we get

$$E_{Y,X}[(Y - f(\underline{X}))^2] = E_X[E_{Y|X}[(Y - f(\underline{X}))^2]|\underline{X}]$$

If we then consider a specific value of \underline{X} , meaning $\underline{X} = \underline{x}$, the inner expectation becomes

$$E[(Y - f(\underline{x}))^2 | \underline{X} = \underline{x}]$$

But since \underline{x} is fixed, $f(\underline{x})$ is a constant, hence this expectation can be written as a function in some parameter a

$$q(a) = E_Y[(Y - a)^2]$$

Then we want to find the (constant) value of a that minimizes g(a)

$$g(a) = E_Y[(Y - a)^2]$$

 $= E[Y^2 - 2aY + a^2]$ compute the square
 $= E[Y^2] - 2aE[Y] + E[a^2]$ separate expectations
 $= E[Y^2] - 2aE[Y] + a^2$ pull out a^2 since constant

Then, setting the derivative to zero we get

$$\frac{dg}{da} = -2E[Y] + 2a \stackrel{!}{=} 0 \implies \hat{a} = E[Y]$$

Notice that if you plug \hat{a} into $g(a) = E_Y[(Y-a)^2]$ you get

$$g(\hat{a}) = E_Y[(Y - E[Y])^2] = Var(Y)$$

Since a = E[Y], the function that minimizes the EPE is

$$f(\underline{x}) = E[Y|\underline{X} = \underline{x}]$$

In a regression setting you could use other loss functions, such as

$$L(Y|f(\underline{X})) = |Y - f(\underline{X})|$$

which is used in **least absolute deviation** (LAD) regression. In this case $f(\underline{x})$ is the median of Y given \underline{x} , therefore you have a model which is more robust to outliers.

We will choose the squared error loss and this motivates how the methods are developed. For example:

• Linear regression (parametric regression) can be expressed as

$$Y|\underline{X} = \underline{x} \sim N(\underline{x}^t \beta, \sigma^2)$$

or focussing more on the conditional mean

$$E[Y|\underline{X} = \underline{x}] = \beta_0 + \beta_1 x_1 + \dots + \beta_p x_p$$

• K nearest neighbour (non-parametric regression)

3.3.1 K nearest neighbour for regression

In the scope of regression, K nearest neighbour method is defined as

$$\hat{f}(\underline{x}) = \text{Average } (y_i | \underline{x}_i \in N_K(\underline{x}))$$

Where:

- Average is the sample mean
- $N_K(\underline{x})$ is the neighbourhood of K points closest to \underline{x}

This method makes two approximations:

- It uses sample mean to approximate population mean
- It uses a neighbourhood of x rather than only x

This method works well if you have a high amount of data and the number of parameters (p) is small. This simple method is rarely used, but it has inspired the development of more sophisticated kernel methods.

3.4 Classification setting

In a classification setting the response Y is categorical, with K categories. A classifier is deciding which class to assign to a new observation \underline{x} . In this case \hat{Y} is the predicted class.

An example of loss function that can be used in classificiation setting is the 0-1 loss function which penalizes classifying a datapoint in the wrong class. Assume for instance the binary case, meaning K=2 (but you could generalize for any number of classes); in that case the 0-1 loss function can be represented as

$$L(Y, \hat{Y}(X)) = \text{TODO Add punnet square}$$

Where:

- Y is the true class
- $\hat{Y}(X)$ is the predicted class

Notice that this function penalizes the same way all the misclassified observations (it assigns them the value 1).

Using the 0-1 loss function, the criterion for choosing the model becomes finding \hat{Y} , which is the function of X such that it minimizes the **expected 0-1 loss**, defined as:

$$E[L(Y, \hat{Y}(X))]$$

To find \hat{Y} , let us consider just one \underline{x} , therefore the problem becomes minimizing

$$E[L(Y, \hat{Y}(\underline{x}))] = \sum_{k=1}^{k} L(k, \hat{Y}(\underline{x})) \ p(k|\underline{x})$$

(Which is the weighted sum of \underline{x} belonging to each possible class). Then, considering the binary case, if the classifier predicts \underline{x} to belong to class 0 ($\hat{y} = 0$), we have

$$E[L(Y,0)] = L(0,0)p(0|\underline{x}) + L(1,0)p(1|\underline{x})$$
$$= 0 \cdot p(0|\underline{x}) + 1 \cdot p(1|\underline{x})$$
$$= p(1|x)$$

On the other hand, if the classifier predicts \underline{x} to belong to to class 1 ($\hat{y} = 1$), we have

$$E[L(Y,1)] = L(0,1)p(0|\underline{x}) + L(1,1)p(1|\underline{x})$$
$$= 1 \cdot p(0|\underline{x}) + 0 \cdot p(1|\underline{x})$$
$$= p(0|x)$$

Predicting \underline{x} to the class that minimizes the expected 1-0 loss results in classifying \underline{x} to class 1 if p(1|x) > p(0|x). Since p(0|x) = 1 - p(1|x), we can write

$$p(1|x) > 1 - p(1|x) \implies 2p(1|x) > 1 \implies p(1|x) > 0.5$$

Therefore, in the binary case, \underline{x} is predicted to belong the class with more than 0.5 probability (which is intuitive).

In general (k classes), the 0-1 loss is minimized by the **bayes classifier**, which states to assign \underline{x} to class j where

$$j = \operatorname*{argmax}_{i \, \in \, classes} p(Y = j | \underline{X} = \underline{x})$$

(Which means assigning \underline{x} to the class with the highest probability). This approach tells us how to set the problem, but it does not give any information on the probabilities, which have to be estimated by the single method.

To rewrite the binary case in another way:

assign
$$\underline{x}$$
 to class
$$\begin{cases} 1 \text{ if } p(1|\underline{x}) > 0.5 \\ 0 \text{ if } p(1|\underline{x}) < 0.5 \end{cases}$$

The line formed by the points with probability $p(1|\underline{x}) = 0.5$ is called **bayes decision boundary** or **decision surface**.

The **bayes error rate** (BER) is the probability of committing an error when classifying observations. It is defined as:

BER =
$$1 - E_{\underline{X}} \left(\max_{j} p(j|\underline{x}) \right)$$

We talk about **perfect separation** of the classes when

$$\max_{i} p(j|\underline{x}) = 1 \implies BER = 0$$

In this case it is possible to classify any \underline{x} without error. More often than not, BER > 0, therefore you have some irreducible error due to a partial overlap of the classes.

It is possible to use a more generic loss function, called **misclassification loss function**, which can be represented (if K = 2) as

$$L(Y, \hat{Y}(\underline{x})) = \text{TODO}$$
 add punnet square

This loss function is analogous to the 0-1 loss function, but it assigns different weights to different errors, therefore the misclassifications have different costs:

- C_0 is the misclassification cost of misclassifying a class 0 to 1
- \bullet C_1 is the misclassification cost of misclassifying a class 1 to 0

Just like the 0-1 loss function, misclassification loss function can be generalized for K classes. This loss function is generally better than the 0-1 loss function since in real situations it is likely for misclassifications to have different relevance (credit risk evaluation, medical context and others). Moreover, this loss function allows you to adjust for unbalanced classes. Consider for example a rare disease; using a 0-1 loss function results in a model which always classifies patients as healty, therefore it is almost always correct, but it does not perform well when trying to determine which patients are affected by the disease. When using misclassification loss function, you can instead increase the weight of the misclassification "classify a patient as healthy when they are not", which may lead to worse results in terms of identifying healthy patients, but way better results in terms of identifying diseased patients.

Repeating the same steps used for 0-1 loss function, we can determine the threshold for the misclassification function (for k = 2 but it can be generalized)

$$E[L(Y,0)] = c_1 p(1|x)$$
 and $E[L(Y,1)] = c_0 p(0|x)$

Therefore we assign \underline{x} to class 1 if

$$c_1 p(1|\underline{x}) > c_0 p(0|\underline{x})$$

But since $p(0|\underline{x}) = 1 - p(1|\underline{x})$, we can instead write

$$c_1 p(1|\underline{x}) > c_0 - c_0 p(1|\underline{x}))$$
$$p(1|\underline{x}) > \frac{c_0}{c_0 + c_1}$$

which is the classification threshold that takes into account different error weights.

3.5 Model accuracy

In practice, we would use our chosen method to estimate $f(\underline{x}) = E[Y|\underline{X} = \underline{x}]$ (regression) or $p(1|\underline{x})$ (classification) from training data $(\underline{x}_i, y_i), i = 1, \dots, n$. The **accuracy** of the model is typically measured on some test data $(\underline{x}_i^{(t)}, y_i^{(t)}), i = 1, \dots, m$. Therefore, a regression that uses MSE has for accuracy

$$MSE = \frac{1}{m} \sum_{i=1}^{m} \left(y_i^{(t)} - \hat{f}\left(\underline{x}_i^{(t)}\right) \right)^2$$

Meanwhile, for a classification, you create a confusion matrix, estimate the probabilities via some method, substitute the values into the following formulas

True positive rate (TPR) =
$$\frac{\text{TP}}{\text{TP} + \text{FN}}$$
 = sensitivity

False positive rate (FPR) = $\frac{\text{FP}}{\text{FN} + \text{FP}}$ = $1 - \frac{\text{TN}}{\text{TN} + \text{FP}} = 1 - \text{specificity}$

Error rate = $\frac{\text{FP} + \text{FN}}{m}$

Notice that this last formula only works when using a 1-0 loss function; if you are using a misclassification loss function, the formula is similar but weighted:

$$Misclassification cost = \frac{c_0 \text{ FP} + c_1 \text{ FN}}{m}$$

Since defining the exact costs of the misclassifications might prove difficult, a possible approach is using the receiver operating characteristic (ROC) curve. Basically you plot all the pairs (FPR, TPR) you obtain by varying the classification threshold from 0 to 1; you will then obtain a curve from which you can decide which value to use as a threshold (the one you see fit for your needs of specificity and sensibility).

Different models have different ROC curves. The best possible curve would touch the top left corner, meaning that you never misclassify any observation and therefore the model is deterministic. The worst possibile curve corresponds to the diagonal (from bottom left to top right corner), meaning that you have a 1/2 chance of misclassifying the observation (the same as tossing a coin to decide). The area under the curve (AUC) is the area between the curve and the main diagonal. Its value goes from 0 in the worst case to 0.5 in the best one. For these reasons, you want a curve that is as high as possible (and thus has the biggest possible AUC). When choosing between two models, if the ROC curve of one of the models is always above the ROC of the other model, the former one is the strictly better one; if the two ROC curves intersect, you choose the model whose curve in highest in the region you are most interested in (based on sensibility and specificity).

Chapter 4

Statistical decision theory II

Just like in the regression setting, there are parametric and non-parametric methods in a classification setting too, these methods being:

- logistic regression (parametric classification)
- K nearest neighbour (non-parametric classification)

4.1 K nearest neighbour for classification

K nearest neighbour is a non-parametric method for classification, meaning it allows to estimate the probability of an observation belonging to each of the classes. Formally this can be written as:

$$P(Y=j|\underline{X}=\underline{x}_0), \text{ for } j=1,\ldots,C$$

Where:

- \underline{x}_0 is the test observation
- \bullet j is a specific class
- \bullet C is the number of classes

K is a positive integer and represents the number of training data points closest to the observation data point to consider when classifying. K is called **tuning parameter**, meaning that it can be tweaked in order to change how the model works; notice that the model is still non-parametric since K does not depend on the training data points.

In order to use this method you have to:

- 1. **identify the K training observations** that are closest (according to some distance, typically we consider the Eucledian distance) to \underline{x}_0 . These observations are then indexed with \underline{N}_0
- 2. compute the probability of the observation to belong to each of the classes using the following formula:

$$P(Y = \hat{j}|\underline{X} = \underline{x}_0) = \frac{1}{k} \sum_{i \in \underline{N}_0} I(y_i = j), \text{ where } j = 1, \dots, C$$

Put into words, the equation above means: compute the probability that, given a test observation \underline{x}_0 , it belongs to a specific class j by dividing the number of neighbours belonging to that class by the number of neighbours (K). Repeat for each class.

4.1.1 Bias-Variance trade off in K nearest neighbour classification

The tuning parameter K relates to the complexity of the model and it is the determinant to balance the bias and the variance of the model. Usually the value of K is chosen using some training data. Intuitively (NOT formally), the model space could be considered as an area of dimension n which is divided in regions containing k training points, therefore you have $\frac{n}{k}$ regions, hence:

- if **K** is small, then there are many small regions in the space and the model is complex. For this reason, if we change \underline{x}_0 with another value, the decision surface position changes quite a lot; if **K** is too small, the model could be overfitting and thus perform poorly with new test data points (low bias, high variance).
- if **K** is large, then there are few big regions in the space and the model is simple. For this reason, if we change \underline{x}_0 with another value, the decision surface position will barely change; if **K** is too big, the model will give similar results for many different values of \underline{x}_0 (high bias, low variance).

4.2 Logistic regression

4.2.1 Why we cannot use linear regression in classification setting

Example 1: We want to predict the medical condition of a patient in the emergency room based on their symptoms. Assume that Y takes only 3 possible values:

$$Y = \begin{cases} 1 \text{ if stroke} \\ 2 \text{ if overdose} \\ 3 \text{ if seizures} \end{cases}$$

A linear regression may seem to work, but you have some problems:

- If the order of the data points changes, the model changes.
- We are arbitrarely assigning an order to not necessarely ordered classes.
- We are imposing an arbitrary distance between the categories.

Example 2: Assume that Y from the previous example only takes 2 values:

$$Y = \begin{cases} 0 \text{ if stroke} \\ 1 \text{ if overdose} \end{cases}$$

In this case Y is a random variable described by a bernoulli distribution with some probability p of taking the value 1; therefore

$$Y = \begin{cases} 1 & p \\ 0 & 1 - p \end{cases}$$

which corresponds to the following probability mass function

$$f(y) = p^y (1-p)^{1-y}, y = 0, 1$$

which has for expected value

$$E[Y] = 0 \cdot (1 - p) + 1 \cdot p = p$$

Given these premises, you can write:

$$E[Y|\underline{X} = \underline{x}] = p(1|\underline{x}) = \beta_0 + \beta_1 x_1 + \dots + \beta_p x_p$$

Or considering just the bernoulli case:

$$E[Y|\hat{X} = x] = \hat{\beta}_0 + \hat{\beta}_1 x_1$$

But notice that:

- $p(1|\underline{x}) \in (0,1)$, so the probability should go from 0 to 1
- $\beta_0 + \beta_1 x_1 + \cdots + \beta_p x_p \in (-\infty, +\infty)$, yet using linear regression the probability spans over the entirety of \mathbb{R} .

Therefore we must find a way to normalize the regression in the range (0,1).

4.2.2 Characteristics of the logistic regression

Logistic regression solves the problem of normalizing the regression in a range (0,1) by utilizing a function of the probability p(1|x), namely g, such that

$$g(p(1|\underline{x})) = \beta_0 + \beta_1 x_1 + \dots + \beta_p x_p = \beta^t \underline{x}$$

Logistic regression uses the **logistic function** (also named sigmoid function or activation function), which is defined as

$$p(1|\underline{x}) = \frac{e^{\beta_0 + \beta_1 x_1 + \dots + \beta_p x_p}}{1 + e^{\beta_0 + \beta_1 x_1 + \dots + \beta_p x_p}}$$

and it has the following properties:

- Its codomain corresponds to the interval (0,1)
- It has sigmoid shape with intercept in y = 0.5
- In the binomial case, if $\beta 1 > 0$ then the function is increasing, else it is decreasing.

The g function is the inverse of the conversion from linear regression to logistic function. In the binomial case, the logistic function can be written as:

$$\begin{split} p(1|\underline{x}) &= \frac{e^{\beta_0 + \beta_1 x_1}}{1 + e^{\beta_0 + \beta_1 x_1}} & \text{binomial case for simplicity} \\ &= \frac{e^{\beta_0 + \beta_1 x_1}}{e^0 + e^{\beta_0 + \beta_1 x_1}} & \text{reorganizing for clarity} \\ &= \frac{1}{e^{\beta_0 + \beta_1 x_1}} \cdot \frac{e^{\beta_0 + \beta_1 x_1}}{\frac{e^0}{e^{\beta_0 + \beta_1 x_1}} + 1} & \text{by collecting } e^{\beta_0 + \beta_1 x_1} & \text{from the denominator} \\ &= \frac{1}{\frac{e^0}{e^{\beta_0 + \beta_1 x_1}} + 1} & \text{simplifying} \\ &= \frac{1}{1 + e^{-\beta_0 - \beta_1 x_1}} & \text{since } \frac{a^x}{a^y} = a^{x-y} \end{split}$$

Let us define a quantity called **odds**, which is the ratio between positive and negative outcomes.

$$\begin{aligned} &\operatorname{Odds} \ = \frac{p(1|\underline{x})}{1-p(1|\underline{x})} & \text{by definition} \\ &= \frac{\frac{e^{\underline{\beta}^t \underline{x}}}{1+e^{\underline{\beta}^t \underline{x}}}}{1-\frac{e^{\underline{\beta}^t \underline{x}}}{1+e^{\underline{\beta}^t \underline{x}}}} & \operatorname{since} \ p(1|\underline{x}) = \underline{\beta}^t \underline{x} \\ &= \frac{\frac{e^{\underline{\beta}^t \underline{x}}}{1+e^{\underline{\beta}^t \underline{x}}}}{\frac{1+e^{\underline{\beta}^t \underline{x}}}{1+e^{\underline{\beta}^t \underline{x}}}} & \operatorname{minimum \ common \ denominator} \\ &= \frac{\frac{e^{\underline{\beta}^t \underline{x}}}{1+e^{\underline{\beta}^t \underline{x}}}}{\frac{1}{1+e^{\underline{\beta}^t \underline{x}}}} & \operatorname{simplifying} \\ &= \frac{e^{\underline{\beta}^t \underline{x}}}{1+e^{\underline{\beta}^t \underline{x}}} \cdot \frac{1+e^{\underline{\beta}^t \underline{x}}}{1} & \operatorname{reorganizing} \\ &= e^{\underline{\beta}^t \underline{x}} & \operatorname{simplifying} \end{aligned}$$

The predictor for this measure is called log-odds or logit (it is the g function from before) is obtained applying a logarithmic tranformation, hence:

$$\log\left(\frac{p(1|\underline{x})}{1-p(1|\underline{x})}\right) = \beta_0 + \beta_1 x_1 + \dots + \beta_p x_p = \underline{\beta}^t \underline{x}$$

The formal definition of a logistic regression model (assuming the random variable to be bernoulli distributed) then becomes:

$$Y|\underline{X} = \underline{x} \sim Bernoulli(p(1|\underline{x}))$$
$$\log\left(\frac{p(1|\underline{x})}{1 - p(1|\underline{x})}\right) = \beta_0 + \beta_1 x_1 + \dots + \beta_p x_p$$

The coefficients can be interpreted as follows: a one-unit increase in X_j (while everithing else is kept constant) induces a change in the log-odds of β_j (while it induces a multiple factor of e^{β_j} for the odds).

4.2.3 Example of logistic regression

Example: Logistic regression for a single predictor X which is binary.

To compute the change of the log-odds for a one-unit change of X we compute

$$\begin{split} \log(\text{odds ratio}) &= \log \left(\frac{p(1|X=x_1+1)}{1-p(1|X=x_1+1)} \right) - \log \left(\frac{p(1|X=x_1)}{1-p(1|X=x_1)} \right) & \text{log odds in 1 minus log odds in 0} \\ &= \beta_0 + \beta_1(x_1+1) - \beta_0 - \beta_1 x_1 & \text{by definition of logistic regression} \\ &= \beta_1 & \text{simplifying} \end{split}$$

This quantity is called $\log(\text{odds ratio})$ since the subtraction of two logaritms in the same base is equal to the logarithm of the ratio of the arguments. Notice that, as expected, a one unit increase of the value of x_1 results in an increment of β_1 .

If we instead consider the same increment for the **odds ratio** (the exponential of the previous function), we get

odds ratio =
$$\frac{\frac{p(1|X=x_0+1)}{1-p(1|X=x_0+1)}}{\frac{p(1|X=x_0)}{1-p(1|X=x_0)}}$$
 odds in 1 divided by odds in 0
$$= \frac{e^{\beta_0+\beta_1(x_0+1)}}{e^{\beta_0+\beta_1x_0}}$$
 using the fact that $\log\left(\frac{p(1|\underline{x})}{1-p(1|\underline{x})}\right) = \underline{\beta}^t\underline{x}$
$$= \frac{e^{\beta_0+\beta_1x_0+\beta_1}}{e^{\beta_0+\beta_1x_0}}$$
 rearranging
$$= e^{\beta_1}$$
 since $\frac{x^a}{x^b} = x^{a-b}$

Therefore we have an increase of e^{β_1} as expected.

4.2.4 Parameter estimation in logistic regression

When estimating parameters for non-linear regressions you often get non-closed form expressions, therefore the maximum likelihood(s) for the parameter(s) must be maximised numerically. Notice that, contrarely to linear regression, you cannot use the least square method.

For instance, give some training data (\underline{x}_i, y_i) , i = 1, ..., n, we want to estimate $\beta_0, \beta_1, ..., \beta_p$. To do so we need to compute the maximum likelihood starting from the fact that Y is binary, thus

$$Y|\underline{X} = \underline{x} \sim Bernoulli(p(1|\underline{x})), \ f(y|\underline{x}) = p(1|\underline{x})^y (1 - p(1|\underline{x}))^{1-y}$$

The likelihood conditional to the observations hence is:

$$L(\underline{\beta}) = \prod_{i=1}^{n} f(y_i | \underline{x}_i) = \prod_{i=1}^{n} p(1 | \underline{x}_i)^{y_i} (1 - p(1 | \underline{x}_i))^{1 - y_i}$$

We then compute the log-likelihood and reorganize it in order to make evident its dependence from $\underline{\beta}$:

$$\begin{split} l(\underline{\beta}) &= \sum_{i=1}^n [y_i \log(p(1|\underline{x}_i)) + (1-y_i) \log(1-p(1|\underline{x}_i))] & \text{applying logarithm} \\ &= \sum_{i=1}^n [y_i \log(p(1|\underline{x}_i)) - y_i \log(1-p(1|\underline{x}_i)) + \log(1-p(1|\underline{x}_i))] & \text{multiplying} \\ &= \sum_{i=1}^n \left[y_i \log\left(\frac{p(1|\underline{x}_i)}{1-p(1|\underline{x}_i)}\right) + \log(1-p(1|\underline{x}_i)) \right] & \log(a) - \log(b) = \log\left(\frac{a}{b}\right) \\ & \text{but since } \log\left(\frac{p(1|\underline{x})}{1-p(1|\underline{x})}\right) = \underline{\beta}^t \underline{x} \text{ and } p(1|\underline{x}) = \frac{e^{\beta_0 + \beta_1 x_1}}{1+e^{\beta_0 + \beta_1 x_1}} \\ &= \sum_{i=1}^n \left[y_i \underline{\beta}^t \underline{x}_i + \log\left(1 - \frac{e^{\underline{\beta}^t \underline{x}_i}}{1+e^{\underline{\beta}^t \underline{x}_i}}\right) \right] & \text{using mcm and simplifying} \\ &= \sum_{i=1}^n \left[y_i \underline{\beta}^t \underline{x}_i + \log\left(1 - \log\left(1 + e^{\underline{\beta}^t \underline{x}_i}\right)\right) \right] & \text{dividing the logarithms} \\ &= \sum_{i=1}^n \left[y_i \underline{\beta}^t \underline{x}_i - \log\left(1 + e^{\underline{\beta}^t \underline{x}_i}\right) \right] & \text{since } \log 1 = 0 \end{split}$$

We then compute the first derivative in $\underline{\beta}$ (notice that since you want to derive in a vector, you need to take the partial derivatives in each of the elements of the vector):

$$\frac{\partial l(\beta)}{\partial \beta_{j}} = \frac{\partial}{\partial \beta_{j}} \left(\sum_{i=1}^{n} \left[y_{i} \underline{\beta}^{t} \underline{x}_{i} - \log \left(1 + e^{\underline{\beta}^{t} \underline{x}_{i}} \right) \right] \right) j = 0, \dots, p$$
but
$$\frac{\partial \underline{\beta}^{t} \underline{x}}{\partial \beta_{j}} = \frac{\partial}{\partial \beta_{j}} (\beta_{0} x_{i0} + \dots + \beta_{p} x_{ip}) = x_{ij} \text{ since the other terms simplify}$$

$$= \sum_{i=1}^{n} \left[y_{i} x_{ij} - \frac{e^{\underline{\beta}^{t} \underline{x}_{i}}}{1 + e^{\underline{\beta}^{t} \underline{x}_{i}}} x_{ij} \right]$$
but
$$\frac{e^{\underline{\beta}^{t} \underline{x}_{i}}}{1 + e^{\underline{\beta}^{t} \underline{x}_{i}}} = p(1|\underline{x}) \text{ therefore}$$

$$= \sum_{i=1}^{n} \left[y_{i} x_{ij} - p(1|\underline{x}_{i}) x_{ij} \right]$$

By setting this derivative to zero you obtain p+1 equations in p+1 parameters, therefore to optimize numerically the function (by solving the equations) you must use an optimization method like **Newton-Raphson**; the gml function in R uses an iterative reweighted least-squares algorithm, which is a variation of the Newton-Raphson method).

From the optimization you get $\hat{\beta} = \hat{\beta}_0, \hat{\beta}_1, \dots, \hat{\beta}_p$, which can be used to estimate $p(1|\underline{x})$, since

$$p(\hat{1}|\underline{x}) = \frac{e^{\hat{\underline{\beta}^t}\underline{x}}}{1 + e^{\hat{\underline{\beta}^t}\underline{x}}}$$

Chapter 5

Generalized linear models

Generalized linear models (GLMs) are a broad family of regression models which includes, for instance, linear, logistic and poisson regressions (and many more). The concept was created by Nelder and Wedderburn in 1972, in order to map into a single framework all regressions models. In general regression models are characterized by:

- Y: the random **response variable** you want to predict
- $\underline{X} = (X_1, \dots, X_p)$: a vector of random predictors (generally known)

Notice that, since these models describe the value of Y with respect to the vector \underline{X} $(Y|\underline{X}=\underline{x})$, they do not give information on the actual distribution of Y, but rather on the values that it takes conditionally to the values taken by the predictors, which can be written as

$$Y_i = (Y | \underline{X} = \underline{x}_i)$$

Therefore, if you have n vectors of random predictors, Y_1, \ldots, Y_n are independent random variables but not identically distributed.

5.1 From linear models to generalized linear models

Linear models make some implicit assumptions that must be removed in order to generalize.

Linear models

Generalized linear models

$$Y_i \sim N(\mu_i, \sigma^2)$$

$$Y_i \sim f(y_i; \mu_i, \phi)$$

$$\mu_i = E[Y_i] = E[Y | \underline{X} = \underline{x}_i] = \eta_i$$

$$g(\mu_i) = \eta_i$$

$$\eta_i = \underline{x}_i^t \beta = \beta_0 x_{i0} + \beta_1 x_{i1} + \dots + \beta_p x_{ip}$$

$$\eta_i = \underline{x}_i^t \beta$$

First of all, linear models assume that the random variable Y_i is normal distributed, with mean μ_i , meaning that $E[Y|\underline{X}=\underline{x}_i]$, and $\sigma^2=Var(Y_i)=Var(Y|\underline{X}=\underline{x}_i)$. Notice that the variance is constant, therefore homoskedasticity is assumed. On the other hand, generalized linear models do not make any assumptions on the shape of the distribution of Y_i ; Y_i is distributed according to some

function $f(y_i)$ with the parameters μ_i and ϕ . ϕ is called **special parameter** and it may or may not be present in the model. The **linear predictor** is the function used to predict the outcome of the random variable Y_i knowing some values for the predictor \underline{X}_i ; in both cases the linear predictor is defined as $\eta_i = \underline{x}_i^t \underline{\beta} = \beta_0 x_{i0} + \beta_1 x_{i1} + \cdots + \beta_p x_{ip}$, where p is the number of parameters and the coefficients β_0, \ldots, β_p are the weights for the predictors (therefore the linear predictor is a linear combination in $\underline{\beta}$ of the predictors). The difference is in the relation between μ_i and η_i . In linear models μ_i coincides with η_i , because assuming a normal distribution $\mu_i = E[Y_i] = E[Y|\underline{X} = \underline{x}_i] = \eta_i$; in generalized linear models μ_i and η_i are related by a **link function**, which is a function g of g such that $g(\mu_i) = \eta_i$. A generalized linear model is then a choice of a distribution function f and a link function g; for instance, if f is a normal distribution (with $\sigma^2 = \phi$) and g is the identity function (therefore $\mu_i = \eta_i$), we get the linear model.

5.2 Exponential dispersion family

The function f, describing the distribution of the random variable Y_i , needs to be a member of the **exponential dispersion family** (EDF). A density belonging to the EDF can be written as

$$f(y_i; \theta_i, \phi) = \exp\left\{\frac{y_i \theta_i - b(\theta_i)}{a_i(\phi)} + c(y_i; \phi)\right\}$$

With:

- $f(y_i)$ is a generic density function
- θ_i and ϕ are some parameters. The values of θ_i are called canonical parameters.
- a_i , b and c are some functions.

5.2.1 Example: normal distribution is a member of the EDF

We want to demonstrate that the normal distribution is a member of the EDF.

$$f(y_i; \mu_i, \sigma^2) = \frac{1}{\sqrt{2\pi\sigma^2}} \exp\left\{-\frac{1}{2} \frac{(y_i - \mu_i)^2}{\sigma^2}\right\} \qquad \text{definition of normal distribution}$$

$$= \exp\left\{-\frac{1}{2} \frac{(y_i - \mu_i)^2}{\sigma^2} + \log\left(\frac{1}{\sqrt{2\pi\sigma^2}}\right)\right\} \qquad \text{pull coefficient inside the exponential}$$

$$= \exp\left\{-\frac{1}{2} \frac{(y_i - \mu_i)^2}{\sigma^2} - \frac{1}{2} \log\left(2\pi\sigma^2\right)\right\} \qquad \text{rearrange using log properties}$$

$$= \exp\left\{-\frac{1}{2} \frac{y_i^2 - 2y_i\mu_i + \mu_i^2}{\sigma^2} - \frac{1}{2} \log\left(2\pi\sigma^2\right)\right\} \qquad \text{compute the square}$$

$$= \exp\left\{\frac{y_i\mu_i + \mu_i^2/2}{\sigma^2} - \frac{1}{2\sigma^2}y_i^2 - \frac{1}{2} \log\left(2\pi\sigma^2\right)\right\} \qquad \text{group terms containing } \mu_i$$

The normal distribution therefore belongs to the EDF by setting:

$$\theta_i = \mu_i,$$

$$\phi = \sigma^2,$$

$$b(\theta_i) = \theta_i^2/2,$$

$$a_i(\phi) = \phi,$$

$$c(y_i; \phi) = -\frac{1}{2\phi}y_i^2 - \frac{1}{2}\log(2\pi\phi)$$

5.2.2 Example: binomial distribution is a member of the EDF

We want to demonstrate that the binomial distribution is a member of the EDF.

$$f(y_i; n_i, \pi_i) = \binom{\pi_i}{y_i} \pi_i^{y_i} (1 - \pi_i)^{n_i - y_i}$$
 definition of binomial distribution
$$= \exp \left\{ \log \left(\binom{\pi_i}{y_i} \pi_i^{y_i} (1 - \pi_i)^{n_i - y_i} \right) \right\}$$
 exponentiate to compare with EDF
$$= \exp \left\{ y_i \log \pi_i + (n_i - y_i) \log(1 - \pi_i) + \log \binom{n_i}{y_i} \right\}$$
 split log and simplify
$$= \exp \left\{ \frac{y_i \log \left(\frac{\pi_i}{1 - \pi_i} \right) + n_i \log(1 - \pi_i)}{1} + \log \binom{n_i}{y_i} \right\}$$
 group terms containing π_i

The binomial distribution therefore belongs to the EDF by setting:

$$\theta_i = \log\left(\frac{\pi}{1-\pi}\right) (\text{logit})$$

$$\phi = 1$$

$$b(\theta_i) = -n_i \log(1-\pi_i)$$

$$= -n_i \log(1 - \frac{e^{\theta_i}}{1+e^{\theta_i}})$$

$$= n_i \log(1 + e^{\theta_i})$$

$$a_i(\phi) = \phi$$

$$c(y_i; \phi) = \log\binom{n_i}{y_i}$$

For the binomial distribution $(Y_i \sim Bin(n_i, \pi_i))$ it is possible to write the relation between old and new parameters explicitly Link from θ_i to μ_i .

$$E[Y_i] = \mu_i = n_i \pi_i = n_i \frac{e^{\theta_i}}{1 + e^{\theta_i}}$$

Link from μ_i to θ_i .

$$\pi_i = \frac{\mu_i}{n_i} \implies \theta_i = \log\left(\frac{\pi_i}{1 - \pi_i}\right) = \log\left(\frac{\mu_i/n_i}{1 - \mu_i/n_i}\right) = \log\left(\frac{\mu_i}{n_i - \mu_i}\right)$$

CHAPTER 5. GENERALIZED LINEAR MODELS

Notice that **heteroskedasticity** in naturally embedded in the model since the variance depens on μ_i :

$$Var(Y_i) = n_i \pi_i (1 - \pi_i) = n_i \frac{\mu_i}{n_i} \left(1 - \frac{\mu_i}{n_i} \right) = \mu_i \left(\frac{n_i - \mu_i}{n_i} \right)$$

This hold true for all GLMs but linear models which is the only homoskedastic model of the group. The fact that heteroskedasticity is mandatory is not a big problem since most datasets are not homoskedastic anyway.

5.3 Connection between θ_i and $(\mu_i, Var(Y_i))$

In general, one can show that, for any function in the EDF,

$$\mu_i = E[Y_i] = b'(\theta_i)$$

Where b' is the first derivative of the function b.

$$Var(Y_i) = b''(\theta_i) \cdot a_i(\phi)$$
$$= V(\mu_i) \cdot a_i(\phi)$$

Where b'' is the second derivative of b and $V(\mu_i)$ is the variance function (basically just a way to rewrite underlining the dependence on μ).

5.3.1 Example: connection of θ_i , μ_i and $Var(Y_i)$ for the binomial model

We know, for the binomial function, that

$$b(\theta_i) = n_i \log(1 + e^{\theta_i})$$

The first derivative in θ_i then becomes

$$b'(\theta_i) = \frac{n_i e^{\theta_i}}{1 + e^{\theta_i}} = n_i \pi_i$$

But it is also known that

$$\mu_i = n_i \pi_i$$

Therefore, as expected from the general formula

$$\mu_i = b'(\theta_i)$$

We know that, for the binomial model

$$a_i(\phi) = 1$$

We then compute the second derivative of $b(\theta_i)$ in θ_i :

$$b''(\theta_i) = n_i e^{\theta_i} \left[-\left(\frac{1}{1 + e^{\theta_i}}\right)^2 e^{\theta_i} \right] + n_i \frac{e^{\theta_i}}{1 + e^{\theta_i}}$$
 by deriving $b'(\theta_i)$

$$= -n_i \left(\frac{e^{\theta_i}}{1 + e^{\theta_i}}\right)^2 + n_i \frac{e^{\theta_i}}{1 + e^{\theta_i}}$$
 pulling both e^{θ_i} inside the square
$$= -n_i \pi_i^2 + n_i \pi_i$$
 since $\frac{e^{\theta_i}}{1 + e^{\theta_i}} = \pi_i$

$$= n_i \pi_i (1 - \pi_i)$$
 collecting $n_i \pi_i$

Therefore, as expected, we get

$$Var(Y_i) = b''(\theta_i)a_i(\phi) = n_i\pi_i(1 - \pi_i)$$

5.4 Choice of link function

The choice of link function is not unique and can depend on the function f which was chosen.

5.4.1 Example: link function for a bernoulli function

Assume that Y_i is bernoulli distributed with probability π_i $(Y_i \sim Ber(\pi_i))$. Then:

- $\mu_i = \pi_i \in (0,1) \text{ (since } n_i = 1)$
- $Var(Y_i) = \pi_i(1 \pi_i) = \mu_i(1 \mu_i)$
- $\eta_i = \underline{x}_i^t \beta$ with $\eta_i \in (-\infty, +\infty)$

We must find a function g such that $g(\mu_i) = \eta_i$ which, taking into account that $\mu_i \in (0,1)$ and $\eta_i \in (-\infty, +\infty)$, means:

$$g: (0,1) \to (-\infty, +\infty)$$

 $g^{-1}: (-\infty, +\infty) \to (0,1)$

These conditions are satisfied by any cumulative density functions (CDF), as long as it is continuous (since it must be invertible). We therefore have infinitely many functions g that satisfy these conditions, but the three most common choices are:

1. Standard normal CDF: Choose $g^{-1} = \Phi$, with $\Phi \sim N(1,0)$, hence

$$g = \Phi^{-1}: (0,1) \to (-\infty, +\infty)$$

This type of regression ($f \sim \text{ bernoulli }, g = \Phi^{-1}$) is called **probit regression**.

2. Logistic CDF: Consider the logistic distribution, which has density

$$f(y;\mu,\sigma^2) = \frac{\exp\left((y-\mu)/\sigma\right)}{\sigma^2\left(1 + \exp\left((y-\mu)/\sigma\right)^2\right)} , \text{ with } -\infty < y < +\infty$$

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CHAPTER 5. GENERALIZED LINEAR MODELS

Setting $\mu = 0$, $\sigma^2 = 1$ we get the sigmoid function

$$g^{-1} = F(y; 0, 1) = \frac{e^y}{1 + e^y}$$

The link function then becomes the logit link

$$\frac{e^{\eta}}{1+e^{\eta}} = \mu \implies \eta = \log\left(\frac{\mu}{1-\mu}\right)$$

This regression ($f \sim \text{bernoulli }, g = \text{logit}$) is the **logistic regression**.

3. **Poisson CDF**: This link function is rarely used since it performs well in niche cases, such as with a binomial distributions with $\pi \sim$ Bernoulli (e.g. the number of bacteria on a plate in a dilution bioassay, in which you do not measure the number of bacteria, just whether there are or not). Assume that Z_i is a possion distributed random variable with parameter λ_i ($Z_i \sim$ Poisson (λ_i)). Then assume that we measure

$$y_i = \begin{cases} 0 & \text{if } z_i = 0 \text{ (absent)} \\ 1 & \text{if } z_i > 0 \text{ (present)} \end{cases}$$

We would like to describe how y_i changes with time (x = 0, 1, 2, ...), therefore we need a function that links μ_i and η_i . Since we know that

$$\eta_i = \log(\lambda_i)$$

and that

$$\mu_i = E[Y_i]$$
= 0 \cdot P(Y_i = 0) + 1 \cdot P(Y_i = 1)
= P(z_i > 0)
= 1 - P(z_i = 0)
= 1 - e^{-\lambda_i}

we can link μ_i and η_i passing through λ_i ($\mu_i \leftrightarrow \lambda_i \leftrightarrow \eta_i$). Firstly we rearrange the link between μ_i and λ_i :

$$\mu_i = 1 - e^{-\lambda_i} \implies e^{-\lambda_i} = 1 - \mu_i \implies \lambda_i = -\log(1 - \mu_i)$$

Then, plugging in the link between λ_i and η_i (log(λ_i) = η_i) we get

$$\eta_i = \log(-\log(1 - \mu_i))$$

This link is called complementary log-log link

5.4.2 Canonical link

The logit link is also called canonical link since

$$\eta_i = \log\left(\frac{\mu_i}{1 - \mu_i}\right) = \theta_i$$

In general a canonical link is a function g such that

$$\eta_i = g(\mu_i) = \theta_i$$

These link functions have some advantages:

5.4. CHOICE OF LINK FUNCTION

- Clear interpretation of the parameters (β)
- Some theoretical reasons (Fisher scoring and iteratively reweighted least squares algorithms are equivalent, zero sum residuals)
- It follows directly from the EDF formulation of the distribution (just derive), in fact $\mu_i = b'(\theta_i), \theta_i = (b')^{-1}(\mu_i)$

Despite these advantages the logit link is not always the correct one to use.