

Introduction to Big Data

Statistical & Machine Learning

Introduction

I have started writing this handout from the notes I have taken from Olivier Rivoire's course on Big Data and Statistical learning at ESPCI Paris from March to April 2018. The LATEX sources of this document can be found on Github. Do not hesitate to update it if you feel it necessary.

Please be considerate if some mistakes crop up in this work.

Iulien

Some book reading is advised during the course, partialicularly:

- *The Elements of Statistical Learning*, T. Hastie, R. Tibshirani and J. Friedman, Springer Series in Statistics, 2008;
- *Information Theory, Inference, and Learning Algorithms*, D.J.C. MacKay, Cambridge University Press, 2003.

Dr Olivier Rivoire

Center for Interdisciplinary Research in Biology (CIRB) Collège de France olivier.rivoire@college-de-france.fr

Applications

There are plenty of applications for Big Data problems. A few examples may be given:

Post learn + identify digits on enveloppes

Biology DNA sequencing

IT Face recognition

etc.

Big Data is an issue of growing importance. As engineers, we may be familiar with such concepts.

Idea of marchine learning

The main idea of machine learning is to find models to give prediction of input data. In facts, Big Data models are deduced from a training batch of N input-output data, on which programs train to generalise models. The deduced model $input\ i \to output\ i$ can then be generalised to give prediction from a random input, as long as it relates to the training batch.

Analytically, let's start with a collection of x and y data, where x stands for the input data and y is the vector of the output data. Each sample is going to have multiple dimensions, therefore we may use an algebraic model. Let N be the number of samples used and p the dimension of each x data. We may write x as an N, p matrix and y as a vector of p dimensions.

We now N samples of p dimensions x_{ij} associated with the N output data y_i .

From now on there are two possible cases: y_i can be known or unknown. In the fist case (y_i known), the problem is said to be *supervised*. Hence we may work with a finite discrete set of data: $y_i = 1, \dots, K$. This problem is called categorical, and we can solve it with *classification*. We may also work with an infinite set of numbers: $y_i \in \mathbb{R}$. This problem is called quantitative, and we can solve it with *regression*.

The second case (y_i unknown) is said to be *unsupervised* and can be solved via *clustering* or *dimension reduction* methods.

Deep learning

In the past few years, there have been huge progress in the *deep learning* approach. It is based on so-called neural networks, that are models inspired by the brain operation.

People are trying to understand how to train these networks. It has had remarkable outcomes in image recognition, social network filtering, medical diagnoses, etc.

Deep learning is based on hidden layers, placed inbetween input and output layers, that are trained to find correlations and mathematical models.

The goal of this course is to explain whate these objects are, how do they work, and put it in relation with state of the art research.

What kind of open problems are there? How do neural networks operate? What are their unsuperised learning behaviour?

Contents

Least square regression, from small to big data 7
General principles 13
Supervised learning: classification, regression, nearest neighbours 27
Unsupervised learning: dimensionality reduction, PCA, SVD 35
Unsupervised learning: clustering, K-means, hierarchical 41
Neural networks, from single neuron to multilayer networks 43
Physics of machine learning, statistical mechanics of machine learning, applications 47

Least square regression, from small to big data

Linear Regression at One Dimension

Let us work at one dimension: p = 1. If we work with N points, then $i = 1, \dots, N$, and we work with a set of data (x_i, y_i) .

The goal here is to make a prediction of what the y data should be when x is given.

The simplest possible model is the linear regression given by the equation 1.

$$y = \alpha + \beta x. \tag{1}$$

Here, the main issue is to get the best α and β for a partialicular set of data. To know what the best choice is, we may define a cost function, that returns a number representing how well the regression permorms. In neural network problems, the cost fuction return number is associated with how well the neural network performs in mapping training examples to the correct output.

There are several choices that can be made to define the cost function. At one dimension, the simplest choice is the sum of squared residuals, defined in equation 2, usually shortened as SSR.

$$l(\alpha, \beta) = \frac{1}{N} \sum_{i=1}^{N} (y_i - \alpha - \beta x_i)^2$$
 (2)

Figure 1 illustrate a simple geometrical interpretation of what the SSR is. Actually, the lower ϵ_i , the better the fit.

For all we have done up to now, we never have never worked with big data. We need p large enough to consider this as a real big data issue.

If we're looking at a hundreds or thousands pixels picture, p will be large in comparison with N. That is a full statistics problem.

Currently, p = 1 is small data, but all we did there has been a correct introduction to clearly understand big data problems.

Striking a good fit necessitates finding the best α and the best β . For this, we may look at the optimum, defined as the points where the derivative of l versus α and β vanishes. This is given by equations 3 and 4.

TODO!

Figure 1: geometrical interpretation of the SSR, where ϵ_i is given by the relation: $\epsilon_i^2 = (y_i - \alpha - \beta x_i)^2$

$$\frac{\partial l}{\partial \alpha} = -\frac{1}{N} \sum_{i=1}^{N} (y_i - \alpha - \beta x_i) = 0$$
 (3)

$$\frac{\partial l}{\partial \beta} = -\frac{1}{N} \sum_{i=1}^{N} x_i (y_i - \alpha - \beta x_i) = 0$$
 (4)

To solve this set of equations, we require a substitution for x and у.

Let's define the mean values:

$$\bar{x} = \frac{1}{N} \sum_{i=1}^{N} x_i \tag{5}$$

$$\bar{y} = \frac{1}{N} \sum_{i=1}^{N} y_i \tag{6}$$

$$\overline{xy} = \frac{1}{N} \sum_{i=1}^{N} x_i y_i \tag{7}$$

$$\overline{x^2} = \frac{1}{N} \sum_{i=1}^{N} x_i^2 \tag{8}$$

Thus, equations 3 and 4 can be reduced as:

$$\frac{\partial l}{\partial \alpha} = -(\bar{y} - \alpha - \beta \bar{x}) \tag{9}$$

$$\frac{\partial l}{\partial \alpha} = -(\bar{y} - \alpha - \beta \bar{x}) \tag{9}$$

$$\frac{\partial l}{\partial \beta} = -(\bar{x}\bar{y} - \alpha \bar{x} - \beta \bar{x}^2) \tag{10}$$

This yields to:

$$\alpha = \bar{y} + \beta \bar{x} \tag{11}$$

$$\overline{xy} = -\bar{y}\bar{x} - \beta\bar{x}^2 - \beta\overline{x^2} \tag{12}$$

There we may substitute α and β :

$$\hat{\alpha} = \bar{y} - \hat{\beta}\bar{x}$$

$$\hat{\beta} = \frac{\overline{xy} - \bar{x}\bar{y}}{\overline{x^2} - \bar{x}^2}$$

$$= \frac{\text{Cov}(x, y)}{\text{Cov}(x, x)}$$

$$= \frac{\text{Cov}(x, y)}{\text{Var}(x)}$$

Let us define the Pearson coefficient \mathcal{R} by the relation 13.

$$\mathcal{R} = \frac{\text{Cov}(x, y)}{\sigma(x)\sigma(y)} \tag{13}$$

The Pearson coefficient R is always comprised between 0 and 1. Thus, we can define the quantity \mathbb{R}^2 , that relates to the quality of the fit:

$$\mathcal{R}^2 = 1 - \frac{\hat{l}}{\text{Var}(y)} \tag{14}$$

We must keep in mind that $\overline{x^2} \neq \overline{x}^2$.

We use the given notations:

$$Cov(x,y) = x\bar{y} - x\bar{y}$$

$$= (x - \bar{x})(y - \bar{y})$$

$$Var(x) = Cov(x,x)$$

$$\sigma(x) = \sqrt{Var(x)}$$

In the general case, we look at models where $\beta = 0$. In this case, the cost function *l* would be the sum of the square distance to the line. Figure 2 depicts two linear regressions with different parameters. The right figure shows a much better linear regression, with much lower square distances between the points and the line.

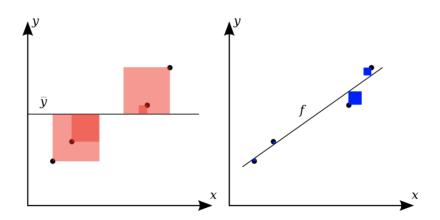


Figure 2: Two linear regression taken from two datasets. The left one shows higher square distances than the right

We may note that in some cases, it would be better to rescale the axis to fit the data with a linear regression. Logarithm axis is the most popular way of rescaling an axis to have a correct assumption. It is usually more valuable to rescale the axis and perform a linear regression, than trying to find an higher order fit.

Linear Regression at Higher Dimensions

We are now considering higher dimensions data (p > 1), that are full, meaning that $p \ll N$. It means that, for an example of data, we might add different parameters. If we take the example (given in class) of correlations between the velocity of people versus the size of towns, we might add other relevant parameters, like the average heigh of people, their ages, etc. We might then examine many potential predictors. Thus we need to generalise the same things, where each input now becomes a vector of p dimensions, as presented in equation 15

$$(x_{i,1},\cdots,x_{i,p}), \quad \forall i$$
 (15)

Let now $x_{i,j}$ be the matrix of the input data, where i is the number of samples, varying from 1 to N, and j is the dimension, ranged between 1 to p.

We may generalise the relation $\hat{y}_i = \hat{\alpha} + \hat{\beta}x_i$ in the new 16 equation:

$$\hat{y}_i = \hat{\alpha} + \sum_{j=1}^p \hat{\beta}_j x_{ij} \tag{16}$$

If we have this key figure, we can always add 1 in the x vector, as the p+1 coordinate. We can thus assume that α vanishes. In fact, we can always redefine the data so that α vanishes. We can also rescale the variable, by removing the mean:

$$x' = x - \bar{x} \tag{17}$$

$$y' = y - \bar{y} \tag{18}$$

Therefore, the output coordinate \hat{y}_i can be written as the product $\hat{y}_i = X\hat{\beta}$, that is a much more convenient way to write it.

That are just restrictions of the problems that help us to compute it

Let $l(\beta)$ be the cross-function, define with equation 19.

$$l(\beta) = \frac{1}{N} \sum_{i=1}^{N} \left(y_i - \sum_{j=1}^{p} \beta_j x_{ij} \right)^2$$
 (19)

Let Z be a vector whose components z_i are defined as follows:

$$\sum_{i=1}^{N} z_i^2 = ||Z||^2 = Z^T Z \tag{20}$$

Therefore we can write the cross-function as:

$$l(\beta) = \frac{1}{N} (Y - X\beta)^T (Y - X\beta) \tag{21}$$

This form can easily be differentiated with β , and the retrieved derivative vanishes at the extremum (eq. 22).

$$\frac{\partial l}{\partial \beta} = -\frac{Z}{N} X^{T} (Y - X\beta) = 0 \tag{22}$$

The equation 22 can be reduced as $X^TY = X^TX\beta$, which can be solved by introducing the matrix $C = X^TX$ (eq. 23)¹

$$\hat{\beta} = (X^T X)^{-1} X^T Y = C^{-1} X^T Y$$
 (23)

At higher dimension, the geometry consists in fitting with an hyperplane, as shown on figure 3.

Here, we are essentially solving a system of equations. We must consider the number of variables adapted to the number of equations that we get. When there is not enough equations (when p is too small for example), the system is undetermined. We cannot reduce it and do not have a single solution.

Actually, when p > N, we can solve this problem with the condition $\hat{l} = 0$. This is a situation where there are more parameters than there are equaltions. It is easy to solve. The solutions consists in overfitting.

For instance if we have 100 parameters and 10 equations, we can never manage to get any result. However, we can find easy solutions, but this will overfit. At this stage, the system cannot be inverted.

^T denotes the transpose matrix, defined by the relation: $\left[\mathbf{A}^{\mathrm{T}}\right]_{ii} = \left[\mathbf{A}\right]_{ji}$

¹ Note that

$$C_{ij} = \sum_{k=1}^{N} x_{kj} x_{ki}$$

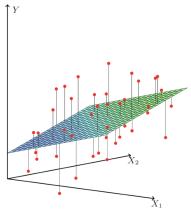


Figure 3: Linear least squares fitting with $X \in \mathbb{R}^2$. In this problem, we are looking for the linear function of X that minimises the sum of squared residuals from Y, which is an plane (hyperplane in dim 3)

When p is large, even if it is of the same order of magnitude as N, we are working with big matrixes, that can be tricky to invert with both proper mathematical accuracy and computation performances. In such cases, we should handle the data carefuly.

Large *p* are typical way of using statistical learning methods, by discriminating between datasets without having the aforesaid issues.

General principles

- Models
- · bias vs variance
- · cross-validation
- maximum likelihood
- Bayes
- etc.

Generally, big data issues are composed of many parameters, with which it is very energy-consuming. We usually want to reduce the problems to one or two interesting parameters. Even if start with a lot of data, we would need to find what are those interesting parameters, and what are the important dimensions of the problem we will be working on.

There are two advantages at reducing the number of parameters:

accuracy With less parameters, we can estimate them with much more accuracy;

easiness it is always easier to estimate things on a condensed set of parameters.

Hence the issue is to find a compromise between having enough parameters to estimate properly the result of the problem, without loosing precision when having a too large batch.

Let $\tilde{l}(\beta)$ be defined as:

$$\tilde{l}(\beta) = l(\beta) + \lambda ||\beta||_q \tag{24}$$

Now, the problem lays in finding

$$\min_{\beta} l(\beta) \text{ given } ||\beta||_0 \le C \text{ with } \beta = [0, \dots, \beta_i, 0, \dots, 0]$$
 (25)

It is not possible to provide this kind of problem with any good numerical solution. In fact, we would always try to find a compromise between what we are able to optimise efficiently and what is possible to optimise. We can write the problem in a new version, easier to solve:

$$\min_{\beta} l(\beta) \text{ given } ||\beta||_2 \le C_1 \tag{26}$$

Here $||\cdot||_0$ stands for the cardinal: $||\beta||_0 = \#(\beta_i \neq 0)$

The two models

There are two models that can be used to solve that kind of problem. The first is the *Ridge regression* that consists in using the constraints that can be solved efficiently numerically.

The second model is called the *Lasso regression* and consists in finding:

$$\min_{\beta} l(\beta) \text{ given } ||\beta||_1 \le C_2 \tag{27}$$

In the ridge problem, we assume that the problem is sparse: we only need a few parameters to capture the relationship.

In fact, the problem can be reset by writing:

$$\tilde{l}(\beta) = (Y - X\beta)^{T} (Y - X\beta) + \lambda \beta^{T} \beta$$
(28)

$$\frac{\partial \tilde{l}(\beta)}{\partial \beta} = 2(-X^T(Y - X\beta) + \lambda\beta) \tag{29}$$

$$= 2(-X^TY + (X^TX + \lambda \mathbb{1})\beta$$
 (30)

We are here adding constraints to the problem, that will allow us to solve it numerically.

$$C_{ik}p > N \tag{31}$$

$$C_{jk} = \frac{1}{N} \sum_{i=1}^{N} x_{ij} x_{ik}$$
 (32)

$$\bar{x} = 0; C = X^T X \tag{33}$$

$$X_{ij}Nxp$$
 (34)

TODO!

p > N => C is non invertible.

$$N = 1, C_{jk} = x_{1j}x_{1k}$$
$$C = XX^T$$

here, C is of rank 1.

Mathematically, $rank(C) \leq N$.

This means that there are too many parameters for too few samples. If we try to solve this with a linear regression problem, it will give too many solutions.

Here it is easy to understand that p > N cases give problems. However, because we are working with big data, the $p \sim N$ cases also gie rise to problems. We would need p << N to have correct solutions.

 $\ensuremath{\mathbb{1}}$ is the identity matrix.

remind that
$$Z^TZ = ||Z||^2$$

$$Z^Ty = < Z, y > = \vec{Z} \cdot \vec{y}$$

Unsupervised learning

Example of financial data

Let's start with stock data, as presented on figure 4.

With these data, we want to extract information. However, y data are not labelled. For this reason, we may not use the raw data, but try to find something else that better fits with the problem.

Let the price of stock i at time t be $s_i(t)$. The stock index i can vary from 1 to p = 1,000, and the data are separated with a time interval $\Delta t = 30$ min. The time variable t ranges from 1 to N = 6228(that represents 2 years).

Let us define the following variables:

$$r_{ti} = \ln \frac{s_i(t + \Delta t)}{s_i(t)}$$

$$x_{ti} = \frac{r_{ti} - \bar{r}_i}{\sigma(r_i)}$$
(35)

$$x_{ti} = \frac{r_{ti} - \bar{r}_i}{\sigma(r_i)} \tag{36}$$

$$C_{ij} = \frac{1}{N} \sum_{i=1}^{N} x_{ti} x_{tj}$$
 (37)

Here, we want to get rid of the α parameter; for this reason, we use the r_{ti} data instead of the $s_i(t)$ parameter.

Then we define x_{ti} , by substracting the mean and normalising with the standard deviation.

Therefore, the x_{ti} value has a null mean, and a standard deviation rescaled to 1. Hence we have rescaled the problem so that each data vary within the same range.

Now, let's move to the data C_{ij} that represents the correlation between two x_{ti} variables. The matrix C of C_{ij} coefficients is of great interest in our problem.

Before trying to make any calculation, we can assume that Exxon's and Chevron's data are strongly correlated, as well as JP Morgan's and Bank of America's.

Thus, we may suppose that $C_{Ex,Ch} > C_{Ex,JP}$.

In order to analyse the data, we may compute the spectrum. Or see clearly from the definition that the matrix is symmetric, and has all the properties to be diagonalised.

$$C_{ik}, C_{ii} = 1; C_{ii} = C_{ii}$$
 (38)

Thus we get the eigenvalues: $\lambda_1, \dots, \lambda_p$, and the eigenvectors v_1, \dots, v_p . The dispersion of the eigenvalues is represented on figure 5.

The bottom part of the graph presented on figure 5 is called the control part. Actually, we see that this part comprises mostly noise, and the interesting part of the data is the upper part, that describes longer range correlations.

The control part just shuffles the data, and is due to permutation of the values . If we randomly shuffle the data, to remove all the



Figure 4: Stock data used as an example

Source: Yahoo! Finance

NB: When we consider some data, it is usually very important to watch it before trying to compute anything else. Particularly, if something looks obvious to the eye, we need to set it clearly, before trying to interpret it with the math.

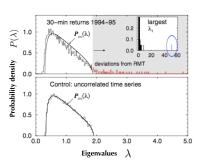


Figure 5: dispersion of the eigenvalues

interesting information (correlation between the different stocks), the control part stays the same. This is actually a way to see what kind of correlation we can get from just randomness, in a case where there is no correlation from the data. Thus, we can quantify the degree of noise. Because 98% of the eigenvalues are contained in the 1st part of the data (control), it means that the stock data are comprised of 98% of noise.

Let's write the matrix C as:

$$C = \sum_{j=1}^{p} \lambda_j v_j v_j^T \tag{39}$$

$$v_j^T v_k = \delta_j^k \tag{40}$$

With the random metrics theory (RMT - branch of statistical physics), it is possible to compute analytically the shape of the control series:

$$\min_{x} l(x) \text{ given } g(x) \le C \tag{41}$$

$$\min_{x} l(x) + \lambda g(x) \text{ with } \lambda \ge 0$$
 (42)

We assume that everything is differentiable, therefore:

$$\nabla_x l(x) = -\lambda \nabla_x g(x) \tag{43}$$

The conditions impose that the two gradients should be aligned, and directed in opposite directions. More generally, it would depend on the value C. There is a minimum value, but if the problem is not perfectly shaped, the minimum we would get can be a local minimum instead of an absolute minimum. This threatens the issue resolution.

Let's suppose we have $l(\beta)$. We may now consider the norm of β to be lower than the C value:

$$||\beta||_q \le C_q \tag{44}$$

Of course, this can also be done with different norms, like l_0 , l_2 , and sometimes also with the l_1 norm.

Example: crime data

Let's give a few examples related to such a problem. We start with the data presented in figure 7

These are not too big data: N = 50 and p = 5. With a not too big dataset, taken from a book. The goal here is to try to predict the crime rate, and to what it is correlated.

$$(x_{ij}, y_i)$$
 $i = 1...N = 50$ cities $j = 1...5$

The idea is to consider naively a simple problem.

Here we may find linear combination of all different problems. For a physicist, it seems we're not allowed to to so, because it is not homogeneous. However, this helps in finding correlations.

$$x'_{ij} = x_{ij} - \bar{x}_j \tag{45}$$

TODO!

Figure 6: scheme

city	funding	hs	not-hs	college	college4	crime rate
1	40	74	11	31	20	478
2	32	72	11	43	18	494
3	57	70	18	16	16	643
4	31	71	11	25	19	341
5	67	72	9	29	24	773
÷	:	:	÷	÷		
50	66	67	26	18	16	940

Figure 7: Example: crime data, crime rate and five predictors for N = 50 U.S.

funding: annual police funding in dollars per resident

hs: % of people 25 years and older with four years of high school

not-hs: % of 16- to 19-year olds not in high school

college: % of 18- to 24-year olds in college

college4: % of people 25 years and older with 4 years of college

In this case, we have zero mean, We may also want try to divide by the standard deviation. This is not the case here.

$$y = \sum_{j=1}^{p} \beta_j x_j \tag{46}$$

Therefore we can write the relation:

$$l(\beta) = \frac{1}{N} \sum_{i=1}^{N} (y_i - \sum_{j=1}^{p} \beta_j x_{ij})^2$$
 (47)

The result of this optimismation may be given as a function. Let's represent the results for the Ridge regression on figure 8

What is plotted is the value of the β along the x axis. This is related with the cost (c_q). We can repeat for different values of the cost. If we do it for large C, we do not put any constraints, and therefore get the same β .

If we put a very strong cost, like 0, the only solution is $\beta=0$. Hence we're looking at different solutions, constrained, and then we relax it to a state where there's no constraint anymore.

Now, let's represent the results for the Lasso regression, on figure 9.

The lasso graph is the same, but performed with l_1 norm.

There's a way to understand this, by giving an illustration of the sparsity from l_1 constraint, as shown on figure 10.

TODO!: introduction à cette équation.

$$y = F(x, \theta). \tag{48}$$

Linear models yield to:

$$F(x,\theta) = \sum_{j=1}^{p} \theta_j x_j \tag{49}$$

The principle hire is to have the results of p data, and then, once we get another dataset – similar to the previous one – we will be able to fit it and to find the solution:

$$x_1, \cdots x_p, x_1^2, \cdots, x_p^2, \cdots cos x_1, \cdots$$
 (50)

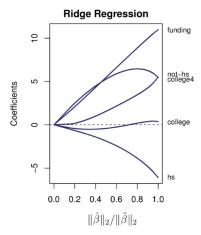


Figure 8: Coefficient path for the ridge regression, plotted versus the relative l_2 norm of the coefficient vector, relative to the norm of the unrestricted least-square estimate $\tilde{\beta}$

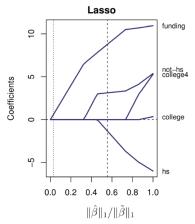


Figure 9: Coefficient path for the lasso regression, plotted against the relative l_1 norm of the coefficient vector, relative to the norm of the unrestricted least-square estimate $\tilde{\beta}$

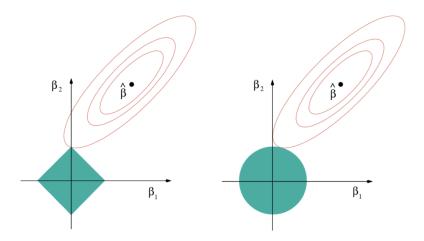


Figure 10: Estimation picture for the lasso (left) and the ridge regression (right). The solud blue areas are the constraint regions $|\beta_1|+|\beta_2| \leq t$ and $\beta_1^2+\beta_2^2 \leq t^2$, respectively, while the red ellipses are the contours of the residual- sum-of-squares function. The point $\hat{\beta}$ depicts the usual (unconstrained) least-squares estimate.

$$F(x,\theta) = \sum \theta_j h_j(x) \tag{51}$$

Later on, we will introduce neural networks. In this case, the function people are more familiar with will be:

$$h_j(x) = \frac{1}{1 + \exp(-\omega_j^T x)}$$
 (52)

Where ω_i is the weight that adjusts the coefficients.

Here, it becomes clear that there is a relation between *x* and *y*. Considering this, there is now no limit to the complexity of the problem we can take.

At this stage it becomes necessary to define a loss function L. This can be written as in equation 53.

$$L(y, F(x,\theta)) = (y - F(x,\theta))^2$$
(53)

With any given model and loss function, we need to define a training error ε_{train} , as shown in equation 54.

$$\varepsilon_{train} = \frac{1}{N} \sum_{i=1}^{N} L(y_i, F(x_i, \theta))$$
 (54)

The training error is of utter importance. However, this will not be the only quantity we will have to consider. In fact, there is also a test/generalisation error. This one would be the error we get when we are using these datapoint that have not been used in the training of the problem.

There's the training set, used to learn the parameters, and the additional data, on which we're going to apply the model, and to try to generalise the data that have been used as an input for the fit.

$$\varepsilon_{test} = \frac{1}{N'} \sum_{i=1}^{N'} L(y_i', F(x_i', \theta))$$
 (55)

Let's plot the training error ε_{train} in comparison with the test error ε_{test} on figure 11. The test errors have been traced in red while the training errors have been traced in blue.

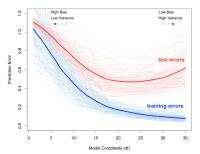


Figure 11: Training vs test errors

k-fold cross validation

Our objective here is not to get the best fit, but to generalise the given data.

The procedure we will chose is called k-fold cross validation. It consists in dividing the data in 5 datasets and then take one out of these to be used a the test, and all the others as the training sets. This is then repeated for all the other combinations. The data are divided in K=5 or 10 and then most of the data are used to train the algorithm, while the last set of data is used to test the algorithm.

In this procedure, the data set is splitted in:

- training set (for the fit) (e.g. 50% of the data)
- test set (for model selection) (e.g. 25% of the data)
- validation set (for assessment) (e.g. 25% of the data)

This is used to better understand how to find the best parameters leading to the best fit.

Let's take our previous example:

$$y = F(x) \tag{56}$$

The training set may be described as: $\hat{\theta}$

The model is $y = F(x, \theta)$

Thus we have $y = \hat{F}(x) = F(x, \theta)$.

Now, let's work with another dataset:

 x_0

$$E[(F(x_0) - \hat{F}(x_0))^2]$$
 (57)

Here, \hat{F} is the prediction.

This is considered over different training sets. It says how far we are from the value we want to predict.

$$E[(F(x_0) - \hat{F}(x_0))^2] = F(x_0)^2 - 2F(x_0)E(\hat{F}(x_0)) + E(\hat{F}(x_0)^2)$$
 (58)

Where
$$E(\hat{F}(x_0)^2)$$
 is $Var(\hat{F}(x_0) + (E[\hat{F}(x_0)])^2$
i.e. $E[(F(x_0) - \hat{F}(x_0))^2] = (F(x_0) - E[\hat{F}(x_0)]^2 + Var(\hat{F}(x_0))$.

The test error may be written as the sum of the variance and the square of the bias.

$$\varepsilon_{test} = bias^2 + Var$$
 (59)

This can be generalised with the addition of a random variable ϵ designating some noise:

$$y = F(x) + \epsilon \tag{60}$$

Here it becomes necessary to add the irreducible error Var(y) In the context of linear regression, one may apply the Gauss-Markov theorem. This one tells us that in linear regression, when

considering all the estimators that have no bias, the best one is the one that minimises the loss function presented in equation 61.

$$L(y, F(x)) = (y - F(x))^{2}$$
(61)

This theorem shows that if one is interested in minimising the bias, in the context of linear regression, he should consider $\min_{\beta} l(\beta)$

The absence of bias and the minimum variance is obtained for

$$\min_{\beta} l(\beta) \tag{62}$$

In general, the best solution is not the solution that has no bias. Hovewer, a better estimation of the parameters can be found when the set of data will be constrained.

bias VS variance

TODO!

Maximum likelihood

Let's discuss about the maximum likelihood estimation (MLE). We must keep in mind that this a general approach in statistics and not a big-data specific one. In general, we would consider a model of the form $y = f(x, \theta) + \epsilon$, where ϵ is a random variable that stands for the noise.

We may assume at a first sight that ϵ is a normal variable:

$$\epsilon(0, \sigma^2)$$
 (63)

The probability to get x with θ may be written as:

$$P(y|x,\theta) = \frac{1}{\sqrt{2\pi\sigma^2}} \exp{-\frac{(y - f(x,\theta))^2}{2\sigma^2}}$$
(64)

Now, let's imagine we are given a set of values (x_i, y_i) . We want to find the best parameter $\hat{\theta}$. For this, we can look at the parameter that make the data the most likeable.

$$\mathcal{L}(\theta|Z_{1}, \dots, Z_{n}) = \log P(Z^{N}|\theta)$$

$$= \sum_{i=1}^{N} \log P(Z_{i}|\theta)$$

$$= \sum_{i=1}^{N} \left[-\frac{1}{2} \log(2\pi\sigma^{2}) - \frac{1}{2\sigma^{2}} (y_{i} - f(x_{i}, \theta))^{2} \right]$$

$$= -\frac{N}{2} \log(2\pi\sigma^{2}) - \frac{1}{2\sigma^{2}} \sum_{i=1}^{N} N(y_{i} - f(x_{i}, \theta))^{2}$$

$$= -\frac{N}{2} \log(2\pi\sigma^{2}) - \frac{1}{2\sigma^{2}} l(\theta)$$

All we have considered until now is just a mean squared deviation (MSE) approach . The MLE approach can be written as:

$$\max_{\theta} \mathcal{L}(\theta, Z^N) \to \hat{\theta} \tag{65}$$

TODO!

Figure 12: $P(y|x,\theta)$ representation. with the standard deviation σ^2 and centered on $f(x,\theta)$

Thus we may get the following theorem:

if y can be written as the sum of f and a random variable bias ϵ : $y = f(x, \theta_0) + \epsilon$, then the expected value E tends towards θ_0 when N is sufficiently high.

$$E[\hat{\theta}] \xrightarrow[N \to \infty]{} \theta_0 \tag{66}$$

Where $\hat{\theta}(\theta_0, F(\theta_0)^{-1})$

 $F(\theta) = E\mathcal{L}(\theta)$ this is called the Fisher information.

$$I(\theta)_{ij} = -\frac{\partial^2 \mathcal{L}}{\partial \theta_i \partial \theta_k}$$

Therefore,

 $\hat{\theta} \sim (approx)N(\hat{\theta},I(\hat{\theta})^{-1})$

We want to find the best parameter, *i.e.* the one that is the more likelihood to be . . .

 ${\cal L}$ is called the log-likelihood. P is colled the likelihood. In a sense, we want to find the most likely model.

We can mention, that there are two difficulties:

- we need to find the maximum
- problem of validation: if we have a more complicated model, we would increase the log likelihood, and no way to do the validation.

Bayesian approach

Now, another approach: the Bayesian approach. Again, this is not specific to big-data. We may know that there are lot of debates in the meaning of probabilities. There are two schools: the frequentists: probabilities have a meaning only when the event is reapeating many times. Like a coin tail or head. Fundamentally, if I do it a large number of time, this is converging. If we take an event that can happen only once: the probability that there's lifeon the moon: for a frequentist, there's no meaning. the Bayesian view is different in nature: probability is not about counting, but about biliefs. This represent how I believe the event to be actually the case. Concretely, the idea of the bayesian approach is: elementary probability: when I have two random variables X, Y. The joint probability of (X,Y): P(X,Y). we can also define the marginals : P(X), P(Y), that are only the probability $P(X) = \sum_{y} P(X,Y)$. P(X|Y): conditional probability. P(X,Y) = P(X|Y)P(Y) therefore $P(X|Y) = \frac{P(X,Y)}{P(Y)}$.

 $P(Z,\theta) = \frac{P(Z,\theta)}{P(Z)} = \frac{P(Z|\theta)P(\theta)}{P(Z)}$. This is called the Bayes formula. $P(Z=1) = \theta \ P(Z=0) = 1 - \theta$, for example, for a binomial problem.

 $P(Z|\theta)$ is the likelihood I had before. In this approach, there's something new, that is $P(\theta)$ wich is called the prior. For a bayesian, we always have some *a priori* bieliefs about the distribution of the parameters. When I see the data, I have to update my beliefs. $P(\theta|Z)$ is called the posterior. P(Z) is called the evidence.

Here we have to do the inference, that is the general model. Usually, when we look at probabilities, we look at $P(Z|\theta)$. reverse

TODO! figure parabolle inversée, maximum : courbature $-\partial \mathcal{L}/\partial \theta^2$, absisse max : θ_0 . ordonnée : L.

approach. We look at the model, that also incorporates the prior.

On the slides, one partialicular example of the prior to solve a praticular problem.

In practice...

Here, there's nothing to do with big data. example taken from the book of MacKay example: Jo has a test for a disease. a = state of health. a = 1 if sick, o otherwise. the test is giving another variable b what is known about this test is that it can be positive even if there's no disease. P(b=1|a=1)=95%. same for zero. it means the test gives the right result in 95% of the cases. we need to know the case when somebody of HIS AGE has the disease. this is gonna be the prior. P(a=1)=1%. The exercise is to find what is the probability P(a=1|b=1). In the bayesian approach, we do not have a theta, but a distribution of the theta we always have a probability to have different values, especially the maximum a posteriori estimate (MAP) which is given by taking the max of this: $\max_{\theta} P(\theta|Z)$.

It will give the same results if we are assuming that this is not depending on theta. If we have a flat prior. Then, this is equivalent to the MLE.

When I'm maximising the posterior, it means I'm maximising the likelihood. Thus we can see the likelihood as a bayesian approach, with a partialicular prior value.

Let's say we have the same model as before. This time, we assume the prior is a one dimensional variable, with a gaussian distribution.___

$$f(\theta) = \sqrt{\frac{\lambda}{2\pi}} e^{-\lambda \frac{\theta^2}{2}}$$

Then, this will be very similar to the l validation, if we take the log of this. Then, this is multiplied by lambda theta square over two.

When we take a prior on these parameters, we want to give more probability to the small values of the parameters. width of the gaussian: 1/sqrt(lambda). control the probability of the parameter to have a large value. restricting the range of the parameter that we are considering, as we saw before

The goal here is to present this partialicular approach, and recover the maximum likelihood. An interesting aspect of the bayesian approach. Again, this is very general. Let's say that, in general, we have the probability of the data, given some hypotheses:

 $P(D|H_1)$ Dis the data, H_1 thehypothesis.

we want to compare with another hypothesis: $P(D|H_2)$

What the bayesian approach is telling us is that we have to consider the probability : $\frac{P(H_1|D)}{P(H_2|D)} = \frac{P(D|H_1)P(H_1)}{P(D|H_2)P(H_2)}$

it depends on the prior we put there. we do not here have no belief????

One example of the maximum likelihood and the general ap-

proach.

Problem of generalization. Let's say we have the data: $x_{i=1...N}$. two approaches are possible: MLE that gives $\hat{\theta}$ bayesian that gives

what is the probability of a partialicular value?

 $P(x_{N+1}) = (MLE) P(x_{N+1}|\hat{\theta})$. if we need a best value of approximation, let's replace the parameter.

= (bayesian)
$$P(x_{N+1}|X^N) = \sum_{\theta} P(x_{N+1}|\theta P(\theta|x^N))$$

here, $P(\theta|x^N)$ should be used as the new prior. We are constantely updating our believes. this is the prior we get before knowing the value, depending on what we saw before. This has to be equal to: $\frac{P(x^N|\theta)P(\theta)}{P(\theta)}$ $P(x^N)$

There's a famous problem that has to do with that: the rule of succession. The sun is rising every morning, what is the probability it will rise tomorrow? Laplace discussed this issue.

With different approaches, we may get different results. One of the simplest example. Let's assume there's some probability θ the sun is rising in the morning. We hae a binomial model. This is the same issue as a coin that always ends up in the same edge: tail for example. Maximum likelihood estimation: the probability would be 1.

$$x_i = 0/1.$$

 $x^N = (x_1, \dots, x_N) = (1, \dots, 1)$

 $P(x_{N+1} = 1) = ?$ If we do the maximum likelihood approximation, this would be 1 everytime.

 $X^N: N_1 times1; N-N_1 times0.$ (times: frequency it happens over time)

 $P(x_i|theta) = \theta$. One parameter model, binomial model. Obviously, the probability $P(x_i = 0 | \theta) = 1 - \theta$. this example cannot be simpler than this.

theta can be anything between 0 and 1. $0 \le \theta \le 1$.

I'm giving the same probability for every θ . $P(\theta)$ ~uniform.

$$P(x_{N+1} = 1 | x^N) = \int d\theta P(x_{N+1} = 1 | \theta) P(\theta | x^N) = \int d\theta \theta \frac{P(x^N | theta) P(\theta)}{P(x^N)} where P(\theta) = 1 = \frac{N!}{N_1!(N-N_1)!} \theta^{N_1} (1-\theta)^{N-N_1}$$

The difficulty lies in the $P(x^N)$ that does not depend on theta.

$$P(\theta|x^{N}) = C(N, N_{1})\theta^{N_{1}}(1-\theta)^{N-N_{1}}$$

$$\int d\theta P(\theta|x^N) = 1$$

 $\int_0^1 d\theta \theta^a (1-\theta)^b = \frac{a!b!}{(a+b+1)!}$ this exists also for a and b that are not integer values, and this is called the gamma function. If we use this formula, we would find this to be:

 $c(N, N_1) = \frac{(N+1)!}{N_1!(N-N_1)!}$ with the proper normalisation. I need just one line to compute the stuff. If I compute this:

$$P(x_{N+1} = 1|x^N) = \int_0^1 d\theta \theta \frac{(N+1)!}{N_1!(N-N_1)!} \theta^{N_1} (1-\theta)^{N-N_1} = \frac{(N+1)!}{N_1!(N-N_1)!} \frac{(N_1+1)!(N-N_1)!}{(N+2)!} = \frac{N_1+1}{N+2} \neq \frac{N_1}{N} (MLE)$$
 this is called as the laplacian formula.

we have a non-zero probability to observe something that we have never observed before.

$$N = 3.X^N = (1,1,1)$$
.canwebetthat $x_4 = 1$? maybe it is not

very wise to say this. this rule takes this into account. here, the probability to be 0 will be 1/5, not zero.

people that carry out statistics use pseudo-counts. we are adding one zero and one one. It is a way to regularise the variation. In this case, a frequentist would say that we have too few points and that we must give up the problem. for a bayesian, the calculation would be very dependant on the prior. Prediction on the next outcomes.

commentary: $I(\theta) = -\frac{\partial^2 \mathcal{L}}{\partial \theta_i \partial \theta_i}$

 $\mathcal{F}(\theta) = \mathrm{E}[I(\theta)]$ MLE: $\hat{\theta}(\ldots)$ For each dataset, we can use another expectation. It is mathematical consideration, considering all the possibilities of my data. If we are partialificially generating data when we want to prove all mathematical results analytically, we need all the possible datasets that we can get. variance about everything we can get when generating different datasets.

second hour.

Now I want to discuss the computational issues.

 $\hat{\theta}$ that I want to maximise. log likelihood: $\mathcal{L}(\theta|x^N)$. we would have, in general, to get these data numerically, and not analytically, with optimised function. compromises to be done.

Very simple, but the problem can be complicated if the function has several minimum. Let's assume the problem is convex: the function is convex, as well as the set of data. Any global minimum is a global minimum. It can be generate, but,...

In all these problems, we can consider a gradient descent.

If we want to minimise the function,

scheme fig 2. If we look at the gradient, and spartial from a point (random). we look at the grandient, and move in the direction of it. we usually take a very small displacement. spartial iteratively.

$$\theta_{t+1} = \theta_t - \eta \left(\frac{\partial L}{\partial \theta} \right)_{\theta_t}$$

scheme: cf notes.

cf learning rate, etc.

 $L(\theta) = \sum_{i=1}^{N} (y_i - f(x_i, \theta))^2 = \sum L_i(\theta)$. for each calculation, we have to recompute the data. a version of this algorithm is used very often, and is called stochastic gradient. What we do is: compute $L_i(\theta)$ for a subset of the points. So we take a subset of samples to estimate $L(\theta) = \sum_{i \in subset}$ and we change at each iteration. The sample is called mini-batch.

There's one version of this algorithm, where everytime we take a single value as a subset: subset =1. this is called on-line learning. in this case, what we are doing is:

$$\theta_{t+1} = \theta_t - \eta \nabla_{\theta} L_i(\theta_t).$$

The issue is to get faster. In general, we have to take the sum. At the end, it is equivalent to do a move at everytime than getting the sum from the very beginning. Actually ,this algorithm is very powerful. In neural networks, backpropagation is nothing more than this algorithm put in application.

This is really a local method. If I spartial from one point, I can end somewhere different. I can be trapped in local minima. It is really difficult to find the correct minimum. That's why we usually work with convex functions, where local minima do not exist elswhere than the global minimum.

 l_0 norm, l_1 norm. $||\beta||_0 = \#nonzero\beta_j; ||\beta||_1 = \sum_j |\beta|_j$

For this reason the closest problem to the first is with l_1 , and it is convex so that it can be solved with an iterative method. Generally, these are considerations we want to take into account. We will see examples of doing this next time.

Next week: exercise as homework. practical. that will be the grade.

$$min_{\beta}L(\beta) - > \hat{\beta}$$

Lasso regression. this is one in which we are going to impose a condition: $||\beta|| \le t$ this is equivalent to the fact that we want to minimize : $min_{\beta}L(\beta) + \lambda ||\beta||$ with some parameter λ .

ridge regression: $||\beta|| = ||\beta||_2^2 = \sum \beta_i^2$

if we take the l_1 norm:

lasso regression: $||\beta|| = ||beta||_1 = \sum |\beta_i|$

At this stage, we can take it as an exercise.

Le'ts spartial with the function I want to minimize. $\mathcal{L}(\beta) =$ $\frac{1}{N}\sum_{i=1}^{N}(y_i-\beta x_i)^2+\lambda|\beta|$

differentiation: we must be careful.

$$\frac{\partial \mathcal{L}}{\partial \beta} = 0$$

the maximum, if β is positive ,let's say, we can get the value:

$$(\beta > 0)$$

$$\frac{\partial \mathcal{L}}{\partial \beta} = -\frac{2}{N} \sum_{i=1}^{N} (y_i - \beta x_i) x_i + \lambda.$$

I would take $\bar{x} = 0$, $\bar{y} = 0$ and $\bar{x}^2 = 1$. normalise all the data.

In general, everything can have completely different units. It makes sense here to normalise, so that each data has the same range of variation. So

thus
$$\hat{\beta} = \frac{X^T y}{N} - \beta + \lambda = 0$$

 $\frac{dL}{dbeta} = -2(\frac{X^T y}{N} - \beta) + \lambda = 0$

thus
$$\hat{\beta} = \frac{X^T y}{N} - \frac{\lambda}{2}$$

If
$$\frac{X^Ty}{N} > \frac{\lambda}{2}$$
, thus $\hat{\beta} = \frac{X^Ty}{N} - \frac{\lambda}{2} > 0$ if $\beta < 0$, then we can use the same approach.

$$\hat{\beta} = S_{\lambda/2} \left(\frac{X^T y}{N} \right)$$

 $S_a(u) = sign(u) max(O, |u| - a)$. soft-thresholding operator. cf figure notes.

This is a figure for p=1.

cf on the website a slide with the formulas.

for
$$p > 1$$
,

$$L(\beta) = \frac{1}{N} \sum_{i} (y_i - \sum_{k} \beta_k x_{ik})^2 + \lambda \sum_{k} |\beta_k|$$

$$= \frac{1}{N} \sum_{i} (y_i - \beta_j x_{ij} - \sum_{k \neq j} \beta_k x_{ik})^2 + \lambda |\beta_j| + \lambda \sum_{k \neq j} |\beta_k|$$

let's define $r_{ij} = y_j - \sum_{k \neq j} \beta_k x_{ik} \hat{\beta}_j < -S_{lambda/2}(\frac{1}{N}x_i^T r_j)$ cyclical coordinate descent. we do this for j, and then repeat for j + 1 until convergenc. we get an interative algorithm. β_1 , then β_2 , ...

Because the truc is convex, this is going to converge to the minimum of the function $L(\beta)$.

cf note on this algorithm. If we understand the case for p=1, then we repeat, and because the problem is convex, we're going to converge to the single minimum.

On wednesday, we will see single classification.

Supervised learning: classification, regression, nearest neighbours

regression = supervised learning for quantitative data.

$$x - > y$$

if we want to classify pictures between cats and dogs, two kinds of methods: - supervised: give information before - unsupervised: you figure out that there are two categories, understand that there are two categories, and learn from that.

fig 5

distincition on the nature of the variable that we want to predict. quantitative data: why is it a real number or a vector?

we are going to see today classification. so far we've seen just supervised quantitative : regression. there are more task to do.

today: classification. next course: clustering. after that: neural networks, that can be used for any of those tasks, that involve a bit more methodology.

Classification

Let's start with a set of training data indexed with i: (x_i, k_i) where the label k_i can vary from 1 to K: $k_i = 1, 2, ... K$.

Again, we face a challenge in using large datasets:

$$x_{N+1} \xrightarrow{?} k_{N+1}$$
 (67)

Until now, we have only used regression methods. Let's try to do the same.

Regressions methods at one dimension

Let's start with only two categories : K = 2

$$k_i \in \{0, 1\} \tag{68}$$

From these data, we can try to predict the value of K with unsupervised learning methods. It means that we would get a classification, by trying to find how close the data are from 0 or 1.

If we start with the p = 1 case, we can represent a set of data on graph 13.

TODO!

Here we need to find a value x_h somewhere to be able to split the value between two sets of data. Generally, we would work with more complicated problems, in larger dimensions.

Here, if we carry out the following linear regression:

$$y = \hat{\alpha} + \hat{\beta}x \tag{69}$$

This is actually exactly what we did in the first lecture: we had to compute the $\hat{\alpha}$ and $\hat{\beta}$ that gave us the best results.

The linear regression equation we get on this problem allow us to find the decision boundary, defined as the \hat{x} corresponding to $\hat{y} = \frac{1}{2}$. Thus, we may claim that when \hat{y} is larger than $\frac{1}{2}$, then $\hat{k} = 1$, and when it is lower than $\frac{1}{2}$, \hat{k} vanishes.

Regressions at higher dimensions

Let's start again at two dimensions: p = 2.

We work here with the data presented on figure 14

This method works only if the data are linearly separable – i.e. if we can separate them in two sets of data by an hyperplan, defining the decision boundary. We can there define an error.

It is worth noting that among many, we cannot predict in advance which prediction method would work, without looking at the data. In practical, no method would separate perfectly the data. If we augment the dimension of the space, we can consider other functions. For example, quadratic regression would usually give better results than linear. However, the use of higher dimensions spaces may present tradeoffs: if we add more parameters, the training error would be reduced, but the variance may also rise. The issue here is to find the correct level of compexity that can fit the data, without overfitting it and that can be generalised.

Again we'll use a second class of methods. ability to solve different issues. fig 8 for example (marginfigure)

If we add to these data another class, for example y = 0,1,2. There is now no classical order between the values and the linear regression cannot be computed on the three sets in one time. We have to split the data between two classes and to find two decisions boundaries that split the three classes.

This was a very naive approach. Let's start again with p=1. Let's add more blue squares at higher abscissa. This is presented in figure 15

TODO!

Here, we added data far enough to put the blue squares in another category.

Figure 13: Representation of a dataset with blue dots indexed as 0 and red dots indexed as 1. A linear regression can be plotted on these data

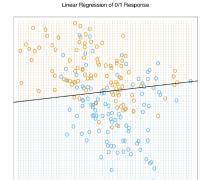


Figure 14: A classification example in two dimensions. The classes are coded as a binary variable (blue = 0, orange = 1), and then fit by linear regression. The line is the decision boundary, defined by $x^T \hat{\beta} = \frac{1}{2}$ drawn as a black line. The orange shaded region denotes that part of input space classified as orange (k = 1), while the blue region is classified as blue (k = 0).

Figure 15: Representation of a dataset with blue dots indexed as 0 and red dots indexed as 1. A linear regression can be plotted on these data

The linear regression can be written as:

$$\min_{\beta} \frac{1}{N} \sum_{i} (y_i - \beta x_i)^2 \tag{70}$$

We can see clearly that it is more convenient to work with more symmetric data. Here we can define the error ε_i for the i index.

$$\varepsilon_i = (y_i - \beta x_i)^2. \tag{71}$$

For the K = 2 problem, we can pose:

$$k_i = \pm 1bd \text{ if } \hat{\beta}x > 0 = > +1 \text{ if } \hat{\beta}x < 0 = > -1$$

What matters here is wether the quantity $\hat{y} = \hat{\beta}x$ has the correct sign or not.

$$L(\beta) = \frac{1}{N} \sum_{i} \max(0, -y_i \beta^T x_i)$$
 (72)

Here again, the idea is that when βx is of the same sign as y, the maximum is zero, and the error is zero. This is called hinge regression.

The idea of the problem is gradient descent, that as a smoother version, called logistic regression.

the idea is to replace this version fig 10.

$$L(\beta) = \frac{1}{N} \sum_{i=1}^{N} \ln(1 + e^{-y_i \beta^T x_i})$$
 (73)

he tries to minimise the error, and the split value with that logistic regression.

we can also arrive to similar results by chosing different approaches. this is connected to the kind of operations that are done in neural networks. we will see this later on.

instead of taking the reflection y to be k, it's easier to consider y to be a probability. it makes here much more sense to make a linear regression. a probabilistic model, in which I want to describe the probability, given some data. parameter ω : weight, standard notation in neural networks, whereas β is much more used in classical regression/statistical learning models.

$$P(k=1|x,\omega) \tag{74}$$

$$P(k=1|x,\omega) > P(k=-1|x,\omega) => \hat{k} = 1$$
 (75)

$$<$$
 => $\hat{k} = -1$ (76)

if we are considering the log of the proabilities:

$$\ln \frac{P(k=1|x,\omega)}{P(k=-1|x,\omega)} = \omega^T x \tag{77}$$

$$= \sum_{j=1}^{P} \omega_j x_j (+\omega_0) \tag{78}$$

I don't have to always write $+\omega_0$, for convenience, just changes in the problem that have no influence.

In general, I want to define a more relevant model.

Let's call:

$$y = P(k = 1|x, \omega) \tag{79}$$

$$\ln(\frac{y}{1-y}) = \omega^T x \tag{80}$$

$$\frac{y}{1-y} = \exp(\omega^T x) \tag{81}$$

thus

$$y = e^{\omega^T x} (1 - y) \tag{82}$$

therefore

$$y(1 + e^{\omega^T x}) = e^{\omega^T x} \tag{83}$$

$$y = \frac{e^{\omega^T x}}{1 + e^{\omega^T x}} = \frac{1}{1 + e^{-\omega^T x}}$$
(84)

$$P(k|x,\omega) = \frac{1}{1 + e^{-k\omega^T x}} \tag{85}$$

$$\mathcal{L}(\omega) = \frac{1}{N} \sum_{i} \ln P(k_i | x_i, \omega) MLE$$
 (86)

if i replace P by its form:

$$\mathcal{L}(\omega) = -\frac{1}{N} \sum_{i} \ln(1 + e^{-k_i \omega^T x_i})$$
 (87)

Essentially, the idea was to apply the maximum likelihood. I had to make hypothesis on the general form of the model. Once I do this, I end up with this function. The different with the one above: minus sign. not surprising: a likelihood we want to maximise, and not error we want to minimize.

$$\max_{\omega} \mathcal{L}(\omega) <=> \min_{\beta} L(\beta) \tag{88}$$

in neural networks,

fig11

$$u = \sum_{i=1}^{p} \omega_j x_j = \omega^T x \tag{89}$$

is called the activation.

several inputs, and one output based on this. if u > 0, then I have output y = 1. if u < 0, then y = 0.

in the context of neural networks, people prefer work with 0 and 1 rather than ± 1 .

neuron is a classifier

fig11

$$y = \phi(u) \tag{90}$$

$$\phi(u) = \frac{1}{1 + e^{-u}} \tag{91}$$

essentially, what we do, is considering the different inputs with different weights. At the end, we will end with more complex classifier models. this one is very elementary. there are very simple algorithms, like the gradient descent.

let's start again with the log likelihood:

$$\mathcal{L}(\omega) = \frac{1}{N} \sum_{i} (k_i \ln y_i + (1 - k_i) \ln(1 - y_i))$$
 (92)

$$= \frac{1}{N} \sum_{i} \ln P(k_i | x_i, \omega) \tag{93}$$

$$y_i = P(k_i = 1 | x, \omega) \tag{94}$$

$$ln(P(\ldots)) = y_i \quad \text{if } k_i = 1 \tag{95}$$

$$=1-y_i$$
 if $k_i=0$ (96)

then we make the model by assuming:

$$y_i = \phi(u_i)whereu_i = \omega^T x_i = \sum_{i=1} \omega_j x_{ij}$$
 (97)

now, we want to compute the gradient. it means to compute the derivative of the log likelihood:

$$\frac{\partial \mathcal{L}}{\partial \omega_j} = \frac{1}{N} \sum_{i} \left(k_i \frac{1}{y_i} \frac{\partial y_i}{\partial \omega_j} - (1 - k_i) \frac{1}{1 - y_i} \right) \frac{\partial y_i}{\partial \omega_j} \tag{98}$$

with

$$(\ldots) = \frac{k_i(1 - y_i) - (1 - k_i)y_i}{y_i(1 - y_i)} = \frac{k_i - y_i}{y_i(1 - y_i)}$$
(99)

and

$$\frac{\partial y_i}{\partial \omega_i} = \phi'(u_i) \frac{\partial u_i}{\partial \omega_i} = \phi'(u_i) X_{ij}$$
 (100)

thus

$$\frac{\partial \mathcal{L}}{\partial \omega_j} = \frac{1}{N} \sum_i \frac{\phi'(u_i)}{\phi(u_i)(1 - \phi(u_i))} (k_i - y_i) x_{ij}$$
 (101)

$$\phi(u) = \frac{1}{1 + e^{-u}} = \frac{\phi''}{\phi(1 - \phi)} = 1 \tag{102}$$

$$\frac{\partial \mathcal{L}}{\partial \omega_i} = \frac{1}{N} \sum_{i} (k_i - y_i) x_{ij} \tag{103}$$

for neural networks, the same idea as gradinet descent is called backpropagation , but it is a bit more complicated.

fig 2.1: representation of classification.

support vector machine: SVM problem formally attitred:

$$\min_{\omega} \sum_{i} \max(0, 1 - k_i \omega^T x_i) + \lambda ||\omega||^2$$
 (104)

the exercise is trying to find what this function is doing.

standard linear regression: take care about the distance to the line. when doing SVM: introducing some margin of errors. this approach would count as an error even if we are on the right side of the line but close to it. we're not going to explain it here, but it looks very much linke the other frameworks.

another approach: linear discriminant analysis (LDA)

we take a model (probabilistic), more gaussian like. the objective function is to fit two classes by the gaussian, in such a way we maximise the distance between the gaussians, and minimise the variance of the gaussians. that can easily be done with linear algebra.

there are many other methods on the same lines.

Now we want to be able to adress the problem of non linear separability. where the real boundary would be a circle. in and out.

fig12

we can also do this locally.

it is called k-nearest neighbours

again, let's start with a batch of values: fig13

we are given with a particular point, and say if it is red or blue. what are the neighbours doing?

if k = 3, i am looking at the three nearest points, and then take the colour of the majority.

mathematically, it can be written as:

$$\hat{y}(x) = \frac{1}{k} \sum_{i \in N_k(x)} y_i \tag{105}$$

 $N_k(x)$ is the set of k nearest neighbours to x $y_i \in \{0, 1\}$

Again, we can be in a situation described by fig 14. we have, there, a decision boundary. i'm deciding what is the coulour of the nearest points with that boundary.

fig 2.2 and 2.3. variations of K

now the question is: what should we take for K

 $\frac{N}{k} \sim p$: number of degrees of freedom.

figure slide: 'selection of K': using 10,000 datapoints and tested with 200.

minimal training error if we take k=1. very complicated thing. as we are taking larger and larger values of k, we are averaging and get more errors, even in the training set. if we have too few degrees of freedom, we are completely missing the structure of data. it is also not good if we have too many degrees of freedom.

15 is relatively well located in this area. between 20 and 10, there's not that much difference, maybe 10 is better because it uses less resources. we have to pick the right model, and there are models of different levels of complexity. different values of k that can be used.

there is nothing into this algorithm, but also nothing trivial in

finding the right datapoints.

there are many variations along this method. in this example, we are taking the same density of points, everywhere. but if I take any point, it doesn't mean anything to take just three points randomly. we may take into account the distance between the points. essentially, we would average around the point that are close to the considered point. more generally, we can do it with a hard cutoff close/far or smooth cutoffs.

$$\hat{y}(x) = \frac{\sum_{i} K_{\lambda}(x, x_{i}) y_{i}}{\sum_{i} k_{\lambda}(x, x_{i})}$$
(106)

$$K_{\lambda}(x,x_i) = \frac{1}{\lambda} e^{-\frac{||x-x'||^2}{2\lambda}} \tag{107}$$

or
$$K_{\lambda}(x, x_i) = 1$$
 if $||x - x'|| < \lambda$, o otherwise.

kernel methods also posisble. if we do this, it is also a case where we have one parameter: λ . It plays the same role. big lambda: averaging a lot. we would reduce the variance, but with a too large bias. in the opposite: reduce the bias would end in a too large variance.

again, there are variety of methods that can be used. one way to see this problem of classification is to fit a function. let's say, zero for orange points, 1 for blue points.

At one dimension, we would have fig 16.

maybe there could be errorsa

the thing is that we may smooth all those functions. we may think about these problems as fitting problems. correcting locally the errors of the functions.

kernel method smoothing, to locally have something smoother. what these things are doing is similar to the problem of smoothly inferring a function from just any perfect observation.

This is very simple, but the problem is that none of those things are working for big data. issue of classifying dogs and cats.

If we go in any of those methods, there is no way we can succeed. Several problems, the greatest is a geometrical one: the curse of dimensionality.

essentially, it has to do with all the assumption we have about neighbourhoods are false when we are working with high dimensional spaces.

example in 2D in fig 17.

volume is scaling as power p of the radius. $V = \alpha_p r^p$.

If we look at
$$\frac{V}{V_{tot}} = \left(\frac{r}{r_{tot}}\right)^p$$

$$\frac{r}{r_{tot}} = \left(\frac{V}{V_{tot}}\right)^{1/p} \tag{108}$$

If the ratio $\frac{V}{V_{tot}}=0.01$. if we take something like p=10, we'll see that $\frac{r}{r_{tot}}=0.6$.

if we want to cover even a very small fraction of the volume, we have to go very far from the point. ridiculous amount of the volume. the problem of scaling with the radius is not interesting. in high dimensional spaces, all the volume is close to the centre.

to cover a small fraction of the volume, we have to go very far. All the ideas of closest neighbourgs cannot be used in high dimensionality spaces.

we can also reduce the dimensionality of the data. that is the first approach with this problem. solutions:

- reduce dimensionality -> PCA (unsupervised)
- representation, that has to do with what kind of space we are
 woking with. we have to imagine that we have all the pictures
 of cats and dogs, all put in the same place. something invariant
 with translation, rotation, etc. so this will be seen in another
 lecture, about neural networks. essentially, one class of netwkrs:
 convolutional, is based on this. chose the right space in which it
 is easy to chose the classification.

exam In priciple, everything should be comprised in the lecture. code it ourself. not using things already written. the game here is to write ourself our own libraries. it is not very complicated, each question is just a few lines of code.

with python: using the jupyter. matlab: not sure it is doeable. livescripts. write a report in which he can see the code. in a format he can actually work the code. minimal thing: write the code and display the figures. this is a public problem, neither on the solutions presented on the internet is good, we might think about it from ourselves.

Unsupervised learning: dimensionality reduction, PCA, SVD

Dimensionality reduction

So far we've seen only supervised learning, relation between two variables (x,y). essentially, the goal was generalization of a relation between x and y

Today, we're going to see a bit of unsupervised learning. this time, we only have x, no y anymore. the goal is to find a pattern. we want to understand the data. it is a more difficult problem. we think about problems of images. we only have the picture, and it is our task to find how to classify them, and what are the common things between all the data. training set, validation set, test set: can not be done anymore. it is part of the reason it is more difficult. however, many things in common. in the supervised context, at the end, it comes to optimization of a function 1 of parameters x, y, and hyperparameters θ , λ .

= arg min $l(x, y, \theta, \lambda)$.

other approach, very similar: present the problem as a problem of optimisation. a more probabilistic one is: we assume $P(x|y,,\lambda)$. we have a statistical model. We have to decide what kind of probabilistic model we chose, and how to parametrise it. Very often, we can formulate it. Then we find the maximum likelihood $\mathcal L$ that we have to maximise, similar to the loss function l that we had to minimise.

for unsupervised learning: dimensional reduction. then we'll see clustering in the last lecture : some generative models.

Let's start with dimensional reduction. Motivation to do this dimensional reduction. problem of the curse of dimensionality.

general approach with high dimension data is to reduce the data to low dimension data. we start with the first method in dimensional reduction :

Principle component analysis, PCA

Probably one of the most useful method. the first we can do if we have large datasets. As usual, we have some data: x_{ij} where i = 1 ... N the number of data, j = 1 ... p the dimension. dimension d < p.

let's start with d = 1. how I can project the data on a line: Let's start with p = 2. I want to reduce it to d = 1

photo 1. we are looking to find a direction that explains more of the variation. along the yellow line: almost just noise. red = principal direction that explains the data.

Let's formulate it as an optimization problem.

Like in the linear regression, I am looking at the line that minimises the distances between the line.

$$l(a, m, v) = \frac{1}{N} \sum_{i=1}^{N} ||x_i - m - a_i v||^2$$
 (109)

and we assume that $||v||^2=1$. we are looking at the parameters a,m,v that minimise the function. We have to find the value of the projection as well as the data.

I am going to differentiate versus all these parameters, and check what we'll get.

First, let's write it as:

$$l(a, m, v) = \frac{1}{N} \sum_{i=1}^{N} (x_i - m - a_i v)^T (x_i - m - a_i v).$$
 (110)

with $x_i - m - a_i v$) is of p dimension.

Let's start with a:

$$\frac{\partial l}{\partial a_i} = \frac{1}{N} v^T (x_i - m - a_i v) \tag{111}$$

$$= \frac{1}{N}(v^{T}(x_{i}-m)-a_{i})=0$$
 (112)

Therefore $\hat{a}_i = v^T(x_i - m) = (x_i - m)^T v$.

$$\frac{\partial l}{\partial m} = \frac{1}{N} \sum_{i=1}^{N} (x_i - m - a_i v)$$
 (113)

$$= \frac{1}{N} \sum_{i=1}^{N} [(x_i - m) - v^T (x_i - m) v] = 0$$
 (114)

Therefore $\hat{m} = \bar{x} = \frac{1}{N} \sum_{i=1}^{N} x_i$. Then we are left with the most difficult part: optimise with respect to v.

To simplify, let's renormalise the data and consider:

$$l(\hat{a}, v) = \frac{1}{N} \sum_{i=1}^{N} (||x_i||^2 - 2x_i^T a_i v + \hat{a}_i^2 v^T v)$$
 (115)

$$= \frac{1}{N} \sum_{i=1}^{N} (||x_i||^2 - \hat{a}_i^2)$$
 (116)

The reason to do it is because $v^Tv = 1$ and $x_i^Ta_iv = \hat{a}_i^2$. Let's write:

$$\frac{1}{N} \sum_{i=1}^{N} \hat{a}_{i}^{2} = \frac{1}{N} \sum_{i=1}^{N} (v^{T} x_{i})^{2} = \frac{1}{N} \sum_{i=1}^{N} \sum_{j=1}^{p} \sum_{k=1}^{p} v_{j} x_{ij} v_{k} x_{ik} = v^{T} C v$$
 (117)

with
$$C = \frac{1}{N}X^TX$$
, or $C_{jk} = \frac{1}{N}\sum_{i=1}^{N}x_{ij}x_{ik}$.

 C_{jk} is the covariance matrix. the quantity we have to maximise is related to the covariance matrix. maximising the left term in the previous equation.

is equivalent to find the minimum of the right term.

Finally, what I have to do is to find

$$\max_{v} v^T C v$$

with the constraint of $||v||^2 = 1$.

then we have the lagrangian function : $L(v)=v^TCv-\lambda v^Tv$ $\frac{\partial L}{\partial v}=2(Cv-\lambda v)=0.$

$$C\hat{v} = \lambda \hat{v}$$

At the end, we need to solve the spectrum of the matrix. the direction v is just going to be an eigenvector of C. we need to understand which eigenvector, which one we want. we want to maximise v^TCv . If I replace Cv by λv , it is just

$$\hat{v}^T C \hat{v} = \lambda \hat{v}^T \hat{v} = \lambda \tag{118}$$

that we want to maximise.

principal component P c = eigenvector with largest eigenvalue λ Again, there are lots of calculation, but at the end, the recipee is simple. Compute C, like of correlation between every dimension. look at the eigenvector that corresppond to the largest eigenvalue. and then projection of what we want as a data.

The goal is to go from p at d=1, and then to generalise for $d \ge 1$. We would have to consider cases such as

$$l(a,v) = \frac{1}{N} \sum_{i=1}^{N} ||x_i - \sum_{k=1}^{d} a_{ik} v_k||^2$$
 (119)

with $v_k^T v_j = \delta_k^j$.

we have v_1, \ldots, v_d top eigenvectors of $C = \frac{1}{N}X^TX$. associated with the largest eigenvalues: $\lambda_1 \ge \lambda_2 \ge \cdots \ge \lambda_d \ge \lambda_{d+1} \ge \cdots \ge \lambda_p$

If I use it with generic data:

figure 14.20: what we did before.

PCA, human genomes. SNP: some bases along the genome that are different along individuals.

4 bases: ATGC, and $L \sim 10^9$. we are looking only at some SNP because there we are able to find differences then get inforamtion about polymorphism.

The data is X_{ij} : 1000 individuals, and then many other dimensions. The data that we get is: all the sequences of SNPs: N = $3e^3$ and dimension of p.

here it is only unsupervised. then we can try to find a y: associations between particular diseases and polymorphism. but not the problem here.

Here: we just have the data and try to find what are the structure of the data. all depends on what kind of data we are working with. in this case, people realise that when we compare. what are the relation ships between individuals, in relation with their genomes. compute a big correlation matrix: this is what C_{ik} is doing. if we rotate a bit the data, we find something like the map of europe. PCA and what we get as a result is a representation vs the geography. origin of the person doesn't go in the definition of the principal component. There is a link in the paper.

peple have applied such analysis in other datasets, to reconstruct the history of migrations, etc. If we compute those data, we'll find the eigenvalues.

but to get this, we need to reduce the datasets of 1000 individuals where we knew the origins (two parents coming from determined places). and then.

generally, there are a few cleaning steps before having proper data. That was a simple example.

 X_{ii} , N samples, p dim. let's compute:

 $C_{jk} = \frac{1}{N} \sum_{i=1}^{N} x_{ij} x_{ik}$ and from here, let's compute the top eigenvectors of the matrix v_1, \ldots, v_d

let's define a matrix called S: $S_{ab} = \frac{1}{p} \sum_{a=1}^{p} x_{aj} x_{bj}$ and from this one, we can get U_1, \ldots, U_d . From here, I decide: let's look at the correlations between these ones, and the correlations between SNPs.

this is where the SVD, i.e. the singular value decomposition is coming to the picture.

Singular value decomposition

This is telling us, that the fact

So, SVD is going to tell us that the two matrixes have the same spectrum: $_1 = \mu_1, \dots, \lambda_d = \mu_d$.

mathematical theorem is to say: we have a matrix X_{ij} of size $N \times p$. It is always possible to find two matrixes so that we can write:

$$X = U\Sigma V^T \tag{120}$$

with $U^TU = 1$ doublebarre $V^TV = 1$ doublebarre.

with σ is diagonal with a spectrum of dimensiond. there is always away to transform the basis in the input outputspaceinawayweareonlymappingonevectorof thebasis, withanotheroneinthe factor.

$$C = X^T X (121)$$

$$= (V\Sigma^T U^T)(U\Sigma V^T)$$
 (122)

$$= V(\Sigma^T \Sigma) V^T \tag{123}$$

And if we do p times S, we got:

$$pS = XX^{T}$$

$$= (U\Sigma V^{T})(V\Sigma^{T}U^{T})$$

$$= U(\Sigma\Sigma^{T})U^{T}$$
(124)
(125)

$$= (U\Sigma V^{T})(V\Sigma^{T}U^{T}) \tag{125}$$

$$= U(\Sigma \Sigma^{I})U^{I} \tag{126}$$

with $\Sigma \Sigma^T = \Sigma^T \Sigma = \sigma_k^2$.

By using this kind of SVD mapping, we can relate the similarities that we see between individuals to correlations between SNPs.

still an unanswered question: what do we take for d?

here, d=2 to see it in relevant dimension, but in general? it is something like an hyperparameter. usually with supervised learning, we know what to do.

here, with unsupervised, we may use the information contained in the λ . *dataisonlyone* – *dimensional*.

let's draw $\rho(\lambda)$. first part of the spectrum = conrol data, and then we have a λ_1 . In general, we can look a the quantity:

$$\frac{\sum_{j=1}^{d} \lambda_j}{\sum_{j=1}^{p} \lambda_j} \tag{127}$$

that is the fraction of variance explained by the top dimension PCs.

the typical criterion is to use this ratio. it gives an indication of how much dimension use.

it is an empirical criterion.

let's come back to the stocks example. first lambda value: almost 50. very large variation. first mode of variation is global changes, we can assume as a first approximation that all the stocks are varying the same way. most of the variation that are present in the matrix C is just noise. this is what we see. typical spectrum we get from high dimensional problems.

then we have a few eigenvalues that are lower. Something we typically do when we analyse that kind of data. something we can do is randomise the data. for example making permutation between the data. standard approach: shuffle the data. it gives us back the control value. we will still compute some correlation matrix. still capture some spectrum. what is written by control. reproduce the lower part of the spectrum. the message here is just to say, this is a standard criterium. when we have a high dimensional data, the fact is that if we sum all the red eigenvalues, the one that are considered as signal, only 2% of the signal. we want to set appart the signal from the noise, given the fact that with high dimensional data, we have a lot of sampling noise.

we are projecting in low dimensions. companies have time to co-vary. will be put together. what would be the eigenvectors? what kind of information could we get by projecting those kind of data?

the first two are similar, the last two are similar. the reason is that they're coming from the same industry.

PCA& RMT: the sectors of the companies are taken as eigenvectors. there are peaks at some companies. here we get information of the economy. Information about different structures of the microbium, etc. compute covariance matrixes, or covariance vectors. a way to understand the data .

Unsupervised learning: clustering, K-means, hierarchical

Clustering

The idea is simple: we want to classify our x. there are two methods. one is called K-means clustering. the other hierarchical clustering.

as usual, the way is to use functions we want to maximise/minimise. cluster 1, cluster 2. we have to find the cluster by ourselves. measure of similarity/dissymilarity between points. let's define a distance $d_{ii'}$. let's take, for example, just the euclidian distance.

$$d_{ii'} = ||x_i - x_{i'}||^2 (128)$$

$$W(C) = \frac{1}{2} \sum_{k=1}^{K} \sum_{(i,i') \in C_{k}^{2}}$$
 (129)

with a partition of my points: K = clusters.

$$1, \dots, N = C_1 \cdots C_k$$

 $C_kC_l = ensemblevide$

$$\min_{C} W(C) = ? \tag{130}$$

we are not going to give all the details of the calculation.

$$T = \frac{1}{2} \sum_{i=1}^{N} \sum_{i'=1}^{N} d_{ii'} = W(C) + B(C)$$
 (131)

with B defined by:

$$B(C) = \frac{1}{2} \sum_{k=1}^{K} \sum_{i \in C_k; i' \notin C_k} d_{ii'}$$
 (132)

and finding $\min_C W(C)$ is then equivalent to find $\max_C B(C)$. warning, this only works for the euclidian distance !!!!!!!! exercise: let's work with:

$$W(C) = \sum_{k=1}^{K} N_k \sum_{i \in C_k} ||x_i - \bar{x}_k||^2$$
 (133)

with $\bar{x}_k = \frac{1}{N_k} \sum_{i \in C_k} x_i$. and $N_k = C_k$

$$\bar{x} = \arg\min_{m} \sum_{i=1}^{N} ||x_i - m_k||^2 = \frac{1}{N_k} \sum_{i=1}^{N_k} x_i$$
 (134)

$$\min_{C} \min_{m_1, \dots, m_k} \sum_{k=1}^{K} N_k \sum_{i \in C_k} ||x_i - m_k||^2$$
 (135)

$$\bar{x} = \arg\min_{m} \sum_{i=1}^{N} ||x_i - m_k||^2 = \frac{1}{N_k} \sum_{i=1}^{N_k} x_i$$
 (136)

$$\min_{C} \min_{m_1, \dots, m_k} \sum_{k=1}^{K} N_k \sum_{i \in C_k} ||x_i - m_k||^2$$
 (137)

k mean clustering. we are doing the minimisation. start initially either with m or c. we are assuming that the mean of cluster one has a particular position, mean of cluster two is elsewhere.

we have to compute it in different ways. there we redefine the clusters, and recompute the means.

figure. once we redefine the mean, we recompute the clusters, and then redoit iteratively.

the best way is to start from different points. if we end up with two different clusterings, we can say which one is better.

let's now introduce another standard algorithm of clustering a way not to solve this problem of k.

hierarchical clustering -> agglomerative -> divisive.

we start in the case we construct clustering, each point is a cluster. Then we make a clustering, and start with N clusters. make n clusters,

the first cluster that makes sense to compute, is just the minimum distances, and say the minimum distance is just replaced by a cluter. initially, we just take the minimum between all the terms. At the end, we have points, but all the points are clusters.

group average:

$$\frac{1}{|C||C'|} \sum_{i,i' \in C} d_{ii'} \tag{138}$$

single linkage:

$$\min_{i \in C, i' \in C'} d_{ii'} \tag{139}$$

complete linkage

$$\max_{i \in C, i' \in C'} d_{ii'} \tag{140}$$

in this way, we would make a series of N clusters. $N \rightarrow N-1 \rightarrow until1$.

simple way to represent the data: dendogram. raprticular representation as a tree. cf fig 14.13.

essentially ,the dedograms are different.

Neural networks, from single neuron to multilayer networks

9 apr 18

most of the reason why people are excited about machine learning is what has been made possible with neural networks. it is mostly related to what we already seen.

representation of real neuron, VS mathematical object. inspiration from biology. brain organised as a network of neurons. each of these neurons are a processor of information. many dendrites on which action potential coming in. then axon, transfer of information as action potential.

scheme.
$$y = \phi(u)$$
 $u = \sum_{j=1}^{p} \omega_j x_j$ $x_0 = 1$

first example when ϕ is linear. linear regression. example: $\phi(u) = u$. $l(\omega) = \sum_{i=1}^{N} (y_i - \phi(u_i))^2$ taking the sum of squares. $\phi(u_i) = \sum_{i=1}^{p} \omega_i x_i$. linear regression, first example.

another example that we've seen: logistic regresion (for classification).

in this context, we took: $\phi(u)=\frac{1}{1+e^{-u}}$. the cost function can be written as: $l(\omega)=-\sum_{i=1}^N k_i \ln y_i + (1-k_i) \ln(1-y_i) \ y_i=P(k_i=1|\omega,x)=\frac{e^{\omega^T x}}{1+e^{\omega^T x}}$ it is what we presented a few pages ago. $l(\omega)$ looks very much like an entropy. y_i looks like a Boltzman distribution with a two-states system.

 $P(E) = \frac{e^-}{Z}$: real boltzman distribution. two energies E1 and E2. $P(E = E_1) = \frac{e}{e^{-0+e^-\beta E_1}} = \frac{1}{1+e^{\beta\Delta E}}$. if $K \geq 2$, soft-max : $k_{im} = 1$ if i is in class m. $i \in 1, \ldots, N$

 $m \in 1, ..., K$. generalisation of the same formula when we have more than two states: $P(k_m = 1 | x, \omega) = \frac{e^{x^T \omega_m}}{\sum_{n=1}^K e^{x^T \omega_m}}$

$$l(\omega) = -\sum_{i=1}^{N} \sum_{m=1}^{K} (k_{im} \ln y_{im} + (1 - k_{im}) \ln(1 - y_{im})).$$

to find the best parameter of this model, write the function: maximum likelihood estimation (MLE). write:

$$\mathcal{L} = \ln P(\cdots)$$
 to minimise: $l = -\mathcal{L}$.

standard approach to define a model that is defined in term of probabilities. mainly to remind that we've already seen many examples on single neurons. essentially this will be discused later with classification of images.

network of neurons, the last layer will convert the likelihood into probabilities. what is the probability that it is a dog, a cat, etc.

many different non linear functions that can be used. we can see the sigmoids. nowadays in neural networks, ReLU and Leary ReLU are mostly used.

Leaky ReLU: does not saturate at high U

 $\phi(u) = 1ifu > 0$, $0ifu \le 0$. let's consider p = 2, $\omega_1 = \omega_2 = -2$, $\omega_0 = 3$. let's consider that x_1 and x_2 are binary values : 0 or 1. and then that $u = 3 - 2(x_1 + x_2)$.

Let's see what are the different options.

$$x_2$$
 x_1 0 1 0 1 1 1 1 0

it is a NAND gate it is a way whith which we can make logical combinations. this framework is one particular type of neurons. If I can do this, I can do any kind of computation. if we are doing picture recognition, input = pixels, output = what it is.

try to do some optimisation, try to organise the NAND gates, to get the circuit that is going to achieve the function that we want. The big thing is that we need an algorithm to get the perfect circuit, network. with gradient descent we can do it. we don't have to think too much about the connection. There would be two components. the weights themselves are parameters.

If I take something more smooth by playing with the weights, I can do any calculation that I want.

Multilayer feedforward network

in general, we could imagine many neural networks. we have to train the network. very simple algorithm: called multilayer feed-forward network. network organised in layers. last layer: output, classification in classes: number K of outputs. input layer: pixels for example. and inbetween, many hidden layers. feedforward, because all the connections are from one layer to the next, all oriented in the same way.

another type:

Backpropagation algorithm

smart implementation of gradient descent. we can start with a single layer. we've already seen this gradient descent algorithm.

single layer:
$$l(\omega) = \frac{1}{2} \sum_{i=1}^{N} \sum_{m=1}^{K} (y_{im} - \phi(u_{im}))^2 u_{im} = \omega_m^T x_i = \sum_{j=1}^{p} \omega_{mj} x_{ij}$$

$$\frac{\partial l}{\partial \omega_{mj}} = -\sum_{i=1}^{N} (y_{im} - \phi(u_{im}) \phi'(u_{im}) \frac{\partial u_{im}}{\partial \omega_{mj}} = -\sum_{i=1}^{N} e_{im} x_{ij}$$
gradient descent
$$\omega_{mj} \leftarrow \omega_{mj} - \eta \frac{\partial l}{\partial \omega_{mj}} = \omega_{mj} - \eta \sum_{i=1}^{N} e_{im} x_{ij}$$
on line: $\omega_{mj} \leftarrow \omega_{mj} - \eta e_{im} x_{ij}$.
we're gonna do the same in the context of neural networks.
$$l(\omega) = \frac{1}{2} \sum_{i=1}^{N} \sum_{m=1}^{K} (y_{im} - x_{im}^{(L)})^2$$
given any input, I update the data, layer by layer. that is called forward propagation: $x^{(0)} x^{(1)} \rightarrow \cdots \rightarrow x^{(L)}$

$$x_{j_k}^{(k)} = \phi(u_{j_k}^{(k)}). \ u_{j_k}^{(k)} = \sum_{j_{k-1}} \omega_{j_k j_{k-1}} x_{j_{k-1}}^{(k-1)}$$
secondly, bacward propagation: $e^{(L)} \rightarrow e^{(L-1)} \rightarrow \cdots \rightarrow e^{(1)}$

$$\begin{array}{l} e_{j_L}^{(L)} = -(y_{j_L}^{(L)} - x_{j_L}^{(L)})\phi'(u_{j_L}^{(L)} \\ e_{j_{k-1}}^{(k-1)} = \phi'(u_{j_{k-1}}^{(k-1)} \sum_{j_k} \omega_{j_k j_{k-1}}^{(k)} e_{j_k}^{(k)} \\ \text{At the end, when we've done this forward an d backward, we} \end{array}$$

have an update, that is the gradient descent. just take the product of the erors:

$$\omega_{ab}^{(n)} \leftarrow \omega_{ab}^{(n)} - \eta e_a^{(n)} x_b^{(n-1)}$$

What we'll try to do is derive??? those equations.

$$\frac{\partial l}{\partial \omega_{ab}^{(n)}} = -\sum_{j_L=1}^{N} (y_{j_L} - \phi(u_{j_L})) \phi'(u_{j_L}) \frac{\partial u_{j_L}}{\partial \omega_{ab}^{(n)}}$$

let's redefine the indexes: layer k. index j for the neuron inside the layer. then we have $x_i^{(k)} = \phi(u_{i_k}^{(k)})$.

the layer. Then we have
$$x_j = \varphi(u_{j_k})$$
.
$$\frac{\partial l}{\partial \omega_{ab}^{(n)}} = \sum_{j_{L-1}} \frac{\partial x_{j_{L-1}}^{(L-1)}}{\partial \omega_{ab}^{(n)}} \sum_{j_L} \omega_{j_L j_{L-1}}^{(L)} e_{j_L}^{(L)}$$
this is true for $k = L - 1$ as shown:
$$\frac{\partial l}{\partial \omega_{ab}^{(n)}} = \sum_{j_k} \frac{x_{j_k}^{(k)}}{\partial \omega_{ab}^{(n)}} \sum_{j_{k+1}} \omega_{j_{k+1} j_k}^{(k+1)} e_{j_{k+1}}^{(k+1)}$$

$$\frac{\partial l}{\partial \omega_{ab}^{(n)}} = \sum_{j_k} \frac{x_{j_k}^{(n)}}{\partial \omega_{ab}^{(n)}} \sum_{j_{k+1}} \omega_{j_{k+1}j_k}^{(k+1)} e_{j_{k+1}}^{(k+1)}$$

with
$$\frac{x_{j_k}^{(k)}}{\partial \omega_{ab}^{(n)}} = \delta_{j_k}^a x_b^{(k)} \phi'(u_{j_k}^{(k)})$$
. for $n = k$,

frac(...) =
$$x_b^{(k)} \phi'(u_a^{(k)}) \sum_{j_{k+1}} \omega_{j_{k+1}a}^{(k+1)} e_{j_{k+1}}^{(k+1)}$$

frac(...) = $x_b^{(k)} \phi'(u_a^{(k)}) \sum_{j_{k+1}} \omega_{j_{k+1}a}^{(k+1)} e_{j_{k+1}}^{(k+1)}$ example. recognition of digits. coming from US postal services. 10 classes. as inputs, we need one neuron for each pixel. output: we used something like soft max. at the end, once we have trained these networks, set the value of the first layer based on the digits.

this example to show that only 15 neurons in the hidden layer. what is written is only the backpropagation algorithm. no hyper parameter. split into training set, and test set. training set larger. let's look at the website, start a very complicated code. relatively simple, already does 95% accuracy. state of the art:99.8%, that is what human would do. deep neural networks: much many layers.

more complicated network can perform to the level of human perfection. current application: - image recognition picture, convolution networks - process sequences that means process speech, movies or text. slightly different architecture. bit of idea of what goes in the current neural networks. essentially, the basic architecture....

in one particular application, convolution networks, ... idea before going into description. different layers, image is the first. we want a transformation of these data in such a way that at the end, we can apply a soft-max. i.e. take our points, and put together all the images that are similar. there's a problem of the distance between the images. if we take the raw images, take all the pictures of cats/dogs together. at the end, project our data. somewhat there are some special properties of images. if we make a slight rotation of the image, it doesn't change the distances. identify the edges in the picture. at the end, we convert the image in a series of activation maps, by applying many filters. question: what filter should I use. we are going to train the filters to find the good value of the weight that is going to identify the particular data in the image.

mostly done with labels???

now, come back to PCA. principal component analysis. simplest but most powerful approach that can be applied to data. that is unsupervised.

as usual ,the data are: x_{ij} , i = 1, ..., N, j = 1, ..., p. what we have to compute is the covariance matrix: C_{ij} . we normalise at first the data by substracting the mean.

$$x'_{ij} = x_{ij} - \frac{1}{N} \sum_{a=1}^{N} x_{aj}$$

$$x' = x - \bar{x} \text{ then}$$

$$x'' = \frac{x'}{\sigma(x')}$$

therefore the matrix C is p times p dimension. So we will compute: $C_{jk} = \frac{1}{N} \sum_{i=1}^{N} x_{ij} x_{ik}$. it is a correlation function. we can write it as:

 $C_{jk} = \overline{x_j x_k}$. in general, we should avoid making loops for sums, and find ways to compute it with linear algearba. in this case: $C_{jk} = \frac{1}{N} X^T X$ the eigenvalues and eigenvectors are : v_1, v_2, \ldots associated with $\lambda_1 \geq \lambda_2 \geq \cdots$. In general, we have a matrix V composed of the eigenvetors in the columns.

we can represent the data by plotting the data versus v_1 and v_2 . they are defined as a multiplicative constant. and then project the data. interpretation of that kind of diagram. each feature is contributing to the eigenvectors. we have to think of them as modes. if we look back at the examples given by financial data, we are correlating the return of our stocks. When we have two companies that are behaving in similar ways.

Physics of machine learning, statistical mechanics of machine learning, applications

application of physics ideas to this field.

Maximum entropy principle

```
equation with a thermal bath. canonical system.
```

$$p(x) = \frac{e^{-\beta E(x)}}{Z} \beta = \frac{1}{kT} Z = \sum_{x} e^{-\beta E(x)}$$

$$S[p] = -\sum_{x} p(x) \ln p(x) E* = \langle E \rangle = \sum_{x} p(x) E(x)$$
therefore $L(p) = S(p) + \lambda \langle E \rangle = -\sum_{x} p(x) \ln p(x) + \lambda \sum_{x} p(x) E(x)$

let's take the derivative $\frac{\partial L}{\partial p(x)} = -\ln p(x) - 1 + \lambda E(x) + \mu = 0$ thus $\ln p(x) = \lambda E(x) + \mu - 1$ this means that $p(x) = \exp(\lambda E(x) + \mu - 1)$ λ and μ come from the constraints. They are such that we have the equations $L(p) = \dots$ we can get it as $p(x) = \frac{e^{-\beta E(x)}}{Z}$. With β , Z such that we have $\sum p(x)E(x) = E*$ and $\sum p(x) = 1$. also $\lambda = -\beta$ and $e^{\mu-1} = \frac{1}{Z}$

grand canonical system

$$x=1,\ldots,2^N$$
 $p(x)=(p_1,\ldots,p_{2^N})$ with $p_x=p(x)$.
we get $\sum p(x)N(x)=N*$ thus $p(x)=\frac{e^{-\beta(E(x)-\mu N(x))}}{Z}=\exp(\lambda_E E(x)+\lambda_N N(x)+\lambda_1)$ with $\lambda_E=-\beta$, $\lambda_N=\beta_\mu$ and $e^{\lambda_1}=\frac{1}{Z}$.
Let's now look x_{ij} , $i=1\ldots N$ and $j=1\ldots p$.
 $\bar{E}=\frac{1}{N}\sum_{i=1}^N E(x_i) \ \phi_a(x)$.

let's take this generical formula into a particular case: $\bar{\phi}_a = \frac{1}{N} \sum_{i=1}^{N} \phi_a(x_i)$.

let's take the maximum entropy : $\max_p S[p]$ such that $<\phi_a>=\sum_x p(x)\phi_a(x)=\bar{\phi}_a$. this is the maximum entropy. i im trying to find this p. the p that we should get is something like: $p(x)=\frac{\exp\sum_a\lambda_a\phi_a(x)}{Z}$. It is a way to arrive to a particular model. We have $p(x|\lambda)$. In this model, we justify the particular form of the model we think is good to consider. There are some observation, and we get this probabilistic model, with some parameters. The parameters are uniquely defined with those constraints.

We will continue to derive the general formalism and at the end, we'll see some particular example. Last time we encountered those examples, we've seen it when we introduced logistic regression. we have the form of the probability distribution, and then we had a model where the probability of the data, given some parameters.

We find the best parameters with the maximum likelihood. We start from a particular distribution, with some data.

Let's start with some likelihood: $(\lambda_a|x_i) = \ln P(x_1 \cdots x_N|\lambda)$. $P(\dots)$ is the product $\prod_{i=1}^N P(x_i|\lambda)$. Thus $\mathcal{L}(\lambda_a|x) = \frac{1}{N} \sum_{i=1}^N \ln P(x_i|\lambda)$

Then, we want the maximum, and the gradient descent can allow us to find it: $\max_{\lambda} \mathcal{L}(\lambda|x)$. It will have a very good form.

Here we have do differentiate with respect to the parameters: $\frac{\partial \mathcal{L}}{\partial \lambda_a}$ we need to express the log likelihood as: $\mathcal{L}(\lambda|x) = \frac{1}{N} \sum_{i=1}^{N} (\sum_{a} \lambda_{a} \phi_{a}(x) - \sum_{i=1}^{N} (\sum_{a} \lambda_{a} \phi_{a}(x))$ ln Z).

$$\begin{split} &Z(\lambda) = \sum_x \exp(\sum_b \lambda_b \phi_b(x)). \ \, \tfrac{\partial Z}{\partial \lambda_a} = \sum_x \phi_a(x) \exp(\sum_b \lambda_b \phi_b(x)) \\ &\tfrac{\partial \ln Z}{\partial \lambda_a} = \tfrac{1}{Z} \tfrac{\partial Z}{\partial \lambda_a} = \sum_x \phi_a(x) p(x) = <\phi_a>. \\ &\text{gradient descent } \lambda_a \leftarrow \lambda_a + \eta(\bar{\phi}_a - <\phi_a>) \text{ then I have to} \end{split}$$

reevaluate the quantity $\langle \phi_a \rangle$ until the gradient vanishes.

This rule is called the Boltzmann machine.

In practice, there's still a difficulty. We need very simple systems to compute the analytical solution. In practice, we do this numerically. We have to sample some of the terms, to make an average. Monte Carlo methods to do it. connections with the Ising Model.

It is not really an accident that the methods people are using in statistical learning have already been invented by physicist. This is the same kind of methods. curse of dimensionality.

We need Markov Chain Monte Carlo (MCMC). It is a way to sample the chain, in such a way it is very likely to have x at a very low probability. One way to do it is to get x at random, compute p(x). There we are very likely to take x at a very low probability. Physically, we try to sample an equilibrium configuration.

reminder of statistical physics. there are many versions. one in particular is called the metropolis-hasting. we are giving a rule to say that if we start from a configuration, we are asked what is the probability p(x' | x).

in the ising model, we start with a random configuration of spins. we flip one spin, and we have to decide wether we need to accept it or to reject it. Those monte carlo are done in such a way that if we average over the trajectory, we would get a good estimate of $<\phi_a>$.

One particular rule is to say that the probability is : p(x'|x) = 1 if E(x') < E(x) and $\exp -\beta(E(x') - E(x))$ if $E(x') \ge E(x)$. $x_1 \to x_2 \to \cdots \to x_T$.