MAT 5314

Assignment 1

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Exercise 1

Suppose $m < d < \infty$ then what can you conclude about the rank of the $d \times d$ matrix A and why does this imply that there are infinitely many w that minimize empirical risk?

solution:

By "Rank–nullity theorem" and the result $\operatorname{Null}(X^tX) = \operatorname{Null}(X)$ [Gram-matrix-note-course] we have: $\operatorname{Rank}(X^tX) = d - \operatorname{Nullity}(X^tX) = d - \operatorname{Nullity}(X) = d - \operatorname{Nullity}(X) = d - \operatorname{Nullity}(X) = d - \operatorname{Rank}(X) = d - (d - \operatorname{Rank}(X)) \leq d - (d - m) < d \text{ then } A = X^tX \notin \operatorname{GL}_d(\mathbb{R}) \Rightarrow \text{ we have } \infty \text{ many sol (because } A\mathcal{Y} = b \perp \operatorname{Null}(A) \text{ see "Revisiting three cases" notes)}.$

Exercise 2

a) part 1

With all the above assumptions, show that Recursive Least Squares will take time $\mathcal{O}(d^3)$ to find the optimal $w_* = w_{n+1}$ knowing the previous $w_* = w_n$ and matrices from that previous computation

solution:

Let's defined some notations that we will use on our pseudo-code :

$$x_{(n+1)} = \begin{bmatrix} x_1^t \\ \vdots \\ x_{n+1}^t \end{bmatrix} = \begin{bmatrix} x_{(n)} \\ x_{n+1}^t \end{bmatrix} \qquad \text{and} \qquad y_{(n+1)} = \begin{bmatrix} y_1 \\ \vdots \\ y_{n+1} \end{bmatrix} = \begin{bmatrix} y_{(n)} \\ y_{n+1} \end{bmatrix}$$

we sow in class that at step (n+1):

$$w_{n+1} = (x_{(n+1)}^t x_{(n+1)})^{-1} x_{(n+1)}^t y_{(n+1)}$$

On the other hand we have a rule for updating our predictor that we sow in class, it express the solution w_{n+1} in terms of things we computed before on step n, so we can show that :

$$w_{n+1} = w_n + P_{n+1}x_{n+1} \left(y_{n+1} - x_{n+1}^t w_n \right)$$

where

$$P_{n+1}^{-1} = P_n^{-1} - x_{n+1} x_{n+1}^t$$

to analyse this algorithm let sketch a pseudo-code for it:

In order to well analysing the complexity of this algorithm here is a table that summarising all the basic operations (additions and multiplications) :

Algorithm 1 Recursive Least Squares

 $w_{i+1} \leftarrow w_i + P_{i+1} x_{i+1} \left(y_{i+1} - x_{i+1}^t w_i \right)$

Require: data matrix X, labels Y1: Initialize: $P_0 \leftarrow I_d$ 2: for $(i \leftarrow 0 \text{ to } n-1)$ do \Rightarrow Complexity: $\sum_{i=0}^{n-1} x_i = 0$ 3: $P_{i+1}^{-1} \leftarrow P_i^{-1} - x_{i+1} x_{i+1}^t$ \Rightarrow Complexity: $\mathcal{O}(d^2)$ 4: inverte P_{i+1}^{-1} with Gaussien elimination algorithm to get P_{i+1} \Rightarrow Complexity: $\mathcal{O}(d^3)$ 5: compute $P_{i+1}x_{n+1}$ \Rightarrow Complexity: $\mathcal{O}(d^2)$ 6: compute $P_{i+1}x_{n+1}$ \Rightarrow Complexity: $\mathcal{O}(d^3)$

 \triangleright Complexity : $\mathcal{O}(d)$

8: **return** w_n

	number of additions	number of multiplications	Total
iteration 3	d^2	d^2	$2d^2 = \mathcal{O}(d^2)$
iteration 4	***************		$\mathcal{O}(d^3)$ (using gauss-elimination)
iteration 5	d^2	d^2	$2d^2 = \mathcal{O}(d^2)$
iteration 6	1 + d	d	$1 + 2d = \mathcal{O}(d)$
iteration 7	d	d	$2d = \mathcal{O}(d)$
			$\mathcal{O}(d^3)$

<u>Note:</u> Because of the iterations from 0 to n-1 (Interation 2) then the time-complexity of the Recursive Least Squares algorithm is: $\mathcal{O}(nd^3)$

Conclusion:

- ♣ The time-Complexity to find the optimal $w_* = w_{n+1}$ knowing the previous $w_* = w_n$ and matrices from that previous computation is : $\mathcal{O}(d^3)$
- \clubsuit The time-Complexity of the Recursive Least Squares algorithm is: $\mathcal{O}(nd^3)$

a) part 2

show that performing Ordinary Least Squares on the full data set (x_i, y_i) , $i = 1, \dots, n+1$ would take time $\mathcal{O}((n+1)(d+d^2)+d^3) = \mathcal{O}(nd^2+d^3)$

solution:

An ERM algorithm is one that performs empirical risk minimization, i.e.

$$w_* = \operatorname{argmin} \sum_{i=1}^{n} (y_i - wx_i)^2$$

we show in class that by differentiation we find that the gradient of the smooth function $J(w) = (Y - Xw)(Y - Xw)^t$ is zero iff the matrix equation:

$$(XX^t)w_* = X^tY$$

here is the pseudo-code for the Ordinary Least Squares :

As in the previous questions here is a table that summarising all the basic operations (additions and multiplications):

Algorithm 2 Ordinary Least Squares

	number of additions	number of multiplications	Total
iteration 1	$d^2(n+1)$	$d^2(n+1)$	$2d^{2}(n+1) = \mathcal{O}(d^{2}(n+1))$
iteration 2	d(n+1)	d(n+1)	$2d(n+1) = \mathcal{O}(d(n+1))$
iteration 3	***************		$\mathcal{O}(d^3)$
			(using gauss-elimination)
iteration 4	d^2	$\int d^2$	$2d^2 = \mathcal{O}(d^2)$
			$\mathcal{O}((n+1)(d^2+d)+d^2+d^3)$
			$= \mathcal{O}(nd^2 + nd + d^2 + d + d^2 + d^3)$
			$=\mathcal{O}(nd^2+d^3)$

Conclusion:

 \clubsuit The time-Complexity of the Ordinary Least Squares algorithm is : $\mathcal{O}(nd^2+d^3)$

b)

What situations does this suggest should be handled by Recursive Least Squares instead of Ordinary Least Squares? Justify your answer

solution:

In order to answer to this question we should compare the complexity of **Ordinary least square** and **Recursive least square** under a order condition between the size of the data "n" and the size of the features "d"

 \clubsuit if: d > n

in this case we have : $nd^2 + d^3 < 2d^3 < nd^3$ and then:

$$\mathcal{O}(nd^2 + d^3) \subset \mathcal{O}(nd^3)$$

♣ if: *d* < *n*

in this case we have : $nd^2 + d^3 > 2d^3$ and then:

$$\mathcal{O}(2d^3) \subset \mathcal{O}(nd^2 + d^3)$$

Or since n is small than d then :

$$\mathcal{O}(d^3) = \mathcal{O}(2d^3) = \mathcal{O}(nd^3) \subset \mathcal{O}(nd^2 + d^3)$$

Conclusion:

- ♣ We can see from the general perspective by analysing the time complexity the it's preferable to use Recursive Least Squares instead of Ordinary Least Squares when we have more features than observations.
- ♣ On the other hand we can think on Memory issues, in fact when we are dealing with very large datasets, matrices can become huge and memory issues may arise when using the ordinary

least square . for example in case when we are using the **data slices** the Recursive Least Squares may perform better the ordinary least square.

Annexe 1:(Gaussian Elimination)

(ASIDE not for credit: think about why this runtime can be achieved using Gaussian Elimination)

To find the inverse of an $n \times n$ matrix A we augment (A|I) and use Gauss-elimination. To show that here is the pseudo-code to Construct an algorithm using Gaussian elimination to find A^{-1} . But before that let's defined what's a augmented matrix:

$$B = (A|I_n) = \begin{bmatrix} a_{1,1} & a_{1,2} & \cdots & a_{1,n} & 1 & 0 & \cdots & 0 \\ a_{2,1} & \cdots & \cdots & a_{2,n} & 0 & 1 & \cdots & 0 \\ \vdots & \vdots & \vdots & \vdots & \vdots & \vdots & 0 & \ddots & 0 \\ a_{n,1} & \cdots & \cdots & a_{n,n} & 0 & 0 & \cdots & 1 \end{bmatrix}$$

Here is the pseudo-code for the Inverting-Gauss-elimination :

Algorithm 3 Inverting-Gauss-elimination

```
Input: augmented matrix B = (A|I_n)
Output: The inverse of A: A^{-1}
                                                                 \triangleright Search for maximum in this column with complexity \sum_{i=0}^{n-1} \mathcal{O}(1)
 1: for i \leftarrow 0 to n-1 do
 2:
           \max \leftarrow |A_{i,i}|
 3:
           \operatorname{argmaxRw} \leftarrow i
                                                                                                                              \triangleright Complexity : \sum_{k=i+1}^{n-1} \mathcal{O}(1)
           for (k \leftarrow i + 1 \text{ to } n - 1) do
 4:
                 if |A_{k,i}| > max then
 5:
                      \max \leftarrow A_{k,i}
 6:
                      \operatorname{argmaxRw} \leftarrow k
 7:
                                                            \triangleright Swap maximum row with current row with complexity: \sum_{k=i}^{2n} \mathcal{O}(1)
           for (k \leftarrow i \text{ to } 2n) do
 8:
                 tmp \leftarrow A_{argmaxRw,k}
 9:
10:
                 A_{argmaxRw,k} \leftarrow A_{i,k}
11:
                 A_{i,k} \leftarrow tmp
                                                         \triangleright Make all rows below this one 0 in current column with : \sum_{k=i+1}^{n} \mathcal{O}(1)
           for (k \leftarrow i + 1 \text{ to } n) do
12:
                c \leftarrow -\frac{A_{k,i}}{A_{i,i}}
13:
                                                                                                                                 \triangleright Complexity : \sum_{j=i}^{2n} \mathcal{O}(1)
                 for (k \leftarrow i \text{ to } 2n) do
14:
                      if i == j then
15:
                            A_{k,i} \leftarrow 0
16:
                       else
17:
                            A_{k,i} \leftarrow A_{k,i} + c.A_{k,i}
18:
19: return: [B_{i,j}]_{\substack{1 \le i \le n \\ 1 \le j \le 2n}}
```

By summing all this Big-O's for all the loops we have:

$$\begin{aligned} run - time &= \mathcal{O}\left(\sum_{i=0}^{n-1} \left(\sum_{k=i+1}^{n-1} 1\right)\right) + \mathcal{O}\left(\sum_{i=0}^{n-1} \left(\sum_{k=i}^{2n} 1\right)\right) + \mathcal{O}\left(\sum_{i=0}^{n-1} \left(\sum_{k=i+1}^{n-1} \sum_{j=i}^{2n} 1\right)\right) \\ &= \mathcal{O}\left(\sum_{i=0}^{n-1} \left(\sum_{k=i+1}^{n-1} \sum_{j=i}^{2n} 1\right)\right) \\ &= \mathcal{O}\left(\sum_{i=0}^{n-1} \left(\sum_{k=i+1}^{n-1} (2n-i+1)\right)\right) \\ &= \mathcal{O}\left(\sum_{i=0}^{n-1} (n-i-1)(2n-i+1)\right) \end{aligned}$$

On the other hand we know that:

$$\sum_{i=1}^{n} i = \frac{n(n+1)}{2} = \mathcal{O}(n^2)$$
$$\sum_{i=1}^{n} i^2 = \frac{n(n+1)(2n+1)}{6} = \mathcal{O}(n^3)$$

therefore:

$$\sum_{i=0}^{n-1} (n-i-1)(2n-i+1) = \underbrace{\sum_{i=0}^{n-1} (n-1)(2n+1)}_{\mathcal{O}(n^3)} + \underbrace{\sum_{i=0}^{n-1} i^2}_{\mathcal{O}(n^3)} + \underbrace{\sum_{i=0}^{n-1} -i(2n+1)}_{\mathcal{O}(n^3)} + \underbrace{\sum_{i=0}^{n-1$$

Finally we have:

$$Run - time = \mathcal{O}\left(\sum_{i=0}^{n-1} (n-i-1)(2n-i+1)\right)$$
$$= \mathcal{O}(n^3)$$

Exercise 3

a)

Consider a slightly more general assumptions where the Gaussians have standard deviation that may depend on the class (value of Y). Derive the form of $\mathbb{P}(Y|X)$ under this new more general version of GNB and say whether it still corresponds to logistic regression.

solution:

In this question we will proceed as in section 3.1" Form of $\mathbb{P}(Y|X)$ for Gaussian Naive Bayes Classifier" (Chapter 3 of Mitchell). We will derive the form of $\mathbb{P}(Y|X)$ entailed by the following modeling assumptions:

- \clubsuit Y is boolean, governed by a Bernoulli distribution, with parameter $\pi = \mathbb{P}(Y=1)$
- A $X = \langle X_1, \dots, X_n \rangle$, where each X_i is a continuous random variable
- For each X_i , $\mathbb{P}(X_i|Y=k)$ is a Gaussian distribution of the form $\mathcal{N}(\mu_{i,k},\sigma_{i,k})$ where, $i=1\cdots n$ and k=0,1.

 \clubsuit For all i and $j \neq i$, X_i and X_j are conditionally independent given Y

First let's express $\pi = \mathbb{P}(Y = 1|X)$ according to $\pi = \mathbb{P}(X_i|Y = 1)$ and π :

$$\mathbb{P}(Y = 1|X) = \frac{\mathbb{P}(Y = 1)\mathbb{P}(X|Y = 1)}{\mathbb{P}(X)} \qquad (BayseRule)$$

$$= \frac{\mathbb{P}(Y = 1)\mathbb{P}(X|Y = 1)}{\mathbb{P}(Y = 1)\mathbb{P}(X|Y = 1)}$$

$$= \frac{1}{1 + \frac{\mathbb{P}(Y = 0)\mathbb{P}(X|Y = 0)}{\mathbb{P}(Y = 1)\mathbb{P}(X|Y = 1)}}$$

$$= \frac{1}{1 + \exp \ln \left(\frac{\mathbb{P}(Y = 0)\mathbb{P}(X|Y = 0)}{\mathbb{P}(Y = 1)\mathbb{P}(X|Y = 1)}\right)}$$

$$= \frac{1}{1 + \exp \left(\ln \frac{1 - \pi}{\pi} + \ln \frac{\mathbb{P}(X|Y = 0)}{\mathbb{P}(X|Y = 1)}\right)}$$

$$= \frac{1}{1 + \exp \left(\ln \frac{1 - \pi}{\pi} + \sum_{i=1}^{n} \ln \frac{\mathbb{P}(X_{i}|Y = 0)}{\mathbb{P}(X_{i}|Y = 1)}\right)}$$

On the other hand we have :

$$\begin{split} \sum_{i=1}^{n} \ln \frac{\mathbb{P}(X_{i}|Y=0)}{\mathbb{P}(X_{i}|Y=1)} &= \sum_{i=1}^{n} \ln \frac{\frac{1}{\sqrt{2\pi\sigma_{i,0}^{2}}} \exp\left(\frac{-(X_{i}-\mu_{i,0})^{2}}{2\sigma_{i,0}^{2}}\right)}{\frac{1}{\sqrt{2\pi\sigma_{i,1}^{2}}} \exp\left(\frac{-(X_{i}-\mu_{i,1})^{2}}{2\sigma_{i,1}^{2}}\right)} \\ &= \sum_{i=1}^{n} \ln \frac{\sigma_{i,1}}{\sigma_{i,0}} + \sum_{i=1}^{n} \left(\frac{(X_{i}-\mu_{i,1})^{2}}{2\sigma_{i,1}^{2}} - \frac{(X_{i}-\mu_{i,0})^{2}}{2\sigma_{i,0}^{2}}\right) \\ &= \sum_{i=1}^{n} \ln \frac{\sigma_{i,1}}{\sigma_{i,0}} + \sum_{i=1}^{n} \frac{(\sigma_{i,0}^{2} - \sigma_{i,0}^{2})X_{i}^{2} + 2(\mu_{i,0}\sigma_{i,1}^{2} - \mu_{i,1}\sigma_{i,0}^{2})X_{i} + \mu_{i,1}^{2}\sigma_{i,0}^{2} - \mu_{i,0}^{2}\sigma_{i,1}^{2}}{2\sigma_{i,0}^{2}\sigma_{i,1}^{2}} \\ &= \sum_{i=1}^{n} \left(\ln \frac{\sigma_{i,1}}{\sigma_{i,0}} + \frac{\mu_{i,1}^{2}\sigma_{i,0}^{2} - \mu_{i,0}^{2}\sigma_{i,1}^{2}}{2\sigma_{i,0}^{2}\sigma_{i,1}^{2}}\right) + \sum_{i=1}^{n} \frac{\mu_{i,0}\sigma_{i,1}^{2} - \mu_{i,1}\sigma_{i,0}^{2}}{\sigma_{i,0}^{2}\sigma_{i,1}^{2}}X_{i} + \sum_{i=1}^{n} \frac{\sigma_{i,0}^{2} - \sigma_{i,1}^{2}}{2\sigma_{i,0}^{2}\sigma_{i,1}^{2}}X_{i}^{2} \end{split}$$

Finnally the form of $\pi = \mathbb{P}(Y = 1|X)$ is given by :

$$\mathbb{P}(Y = 1|X) = \frac{1}{1 + \exp\left(\ln\frac{1-\pi}{\pi} + \sum_{i=1}^{n} \ln\frac{\mathbb{P}(X_i|Y=0)}{\mathbb{P}(X_i|Y=1)}\right)}$$
$$= \frac{1}{1 + \exp\left(w_0 + \sum_{i=1}^{n} w_i X_i + \sum_{i=1}^{n} v_i X_i^2\right)}$$

where:

$$w_0 = \ln \frac{1 - \pi}{\pi} + \sum_{i=1}^n \left(\ln \frac{\sigma_{i,1}}{\sigma_{i,0}} + \frac{\mu_{i,1}^2 \sigma_{i,0}^2 - \mu_{i,0}^2 \sigma_{i,1}^2}{2\sigma_{i,0}^2 \sigma_{i,1}^2} \right)$$

$$w_i = \frac{\mu_{i,0} \sigma_{i,1}^2 - \mu_{i,1} \sigma_{i,0}^2}{\sigma_{i,0}^2 \sigma_{i,1}^2}$$

$$v_i = \frac{\sigma_{i,0}^2 - \sigma_{i,1}^2}{2\sigma_{i,0}^2 \sigma_{i,1}^2}$$

Conclusion:

 \clubsuit when standard deviation may depend on the class and then $\sigma_{i,1} \neq \sigma_{i,0}$ so we have a quadratic terme. Therefore the form $\mathbb{P}(Y=1|X)$ it's differente from the form of logistic regression.

b)

Finally consider a less naive algorithm where conditional independence of the random variables in $X = \langle X_1, \dots, X_n \rangle$ is not assumed - specifically suppose $\mathbb{P}(X|Y=k)$ is multivariate normal with mean depending on k, but covariance matrix now NOT depending on k. Derive $\mathbb{P}(X|Y)$ and say if it has the same form as in logistic regression.

solution:

In this question we will proceed as in the previews question. We will derive the form of $\mathbb{P}(Y|X)$ entailed by the following modeling assumptions:

- \clubsuit Y is boolean, governed by a Bernoulli distribution, with parameter $\pi = \mathbb{P}(Y=1)$
- A $X = \langle X_1, \cdots, X_n \rangle$, where each X_i is a continuous random variable
- A $X = \langle X_1, \dots, X_n \rangle$ are not conditionally independent given Y, and $\mathbb{P}(X|Y=k)$ is a multivariate normal distribution $\mathcal{N}_n(\mu_k, \Sigma)$ where, μ_k is the mean vector that depending on k and Σ is the covariance matrix not depending on k.

Before starting expressing $\mathbb{P}(Y|X)$ let's remember the density of the multivariate normal distribution $\mathcal{N}_n(\mu, \Sigma)$:

$$f(x) = f(x_1, \dots x_n) = \frac{1}{(2\pi)^{\frac{n}{2}} |\Sigma|^{1/2}} \exp^{\frac{-(x-\mu)^t \Sigma^{-1}(x-\mu)}{2}} \quad \forall x \in \mathbb{R}^n$$

Now we have all the tools to start the calculation:

$$\begin{split} \mathbb{P}(Y = 1 | X) &= \frac{\mathbb{P}(Y = 1) \mathbb{P}(X | Y = 1)}{\mathbb{P}(X)} \qquad (BayseRule) \\ &= \frac{\mathbb{P}(Y = 1) \mathbb{P}(X | Y = 1)}{\mathbb{P}(Y = 1) \mathbb{P}(X | Y = 1)} \\ &= \frac{1}{1 + \frac{\mathbb{P}(Y = 0) \mathbb{P}(X | Y = 0)}{\mathbb{P}(Y = 1) \mathbb{P}(X | Y = 1)}} \\ &= \frac{1}{1 + \exp \ln \left(\frac{\mathbb{P}(Y = 0) \mathbb{P}(X | Y = 0)}{\mathbb{P}(Y = 1) \mathbb{P}(X | Y = 1)}\right)} \\ &= \frac{1}{1 + \exp \left(\ln \frac{1 - \pi}{\pi} + \ln \frac{\mathbb{P}(X | Y = 0)}{\mathbb{P}(X | Y = 1)}\right)} \end{split}$$

On the other hand we have :

$$\ln \frac{\mathbb{P}(X|Y=0)}{\mathbb{P}(X|Y=1)} = \ln \frac{\frac{1}{(2\pi)^{\frac{1}{2}}|\Sigma|^{1/2}}}{\frac{1}{(2\pi)^{\frac{n}{2}}|\Sigma|^{1/2}}} + \ln \exp\left(\frac{1}{2}\left[(x-\mu_1)^t\Sigma^{-1}(x-\mu_1) - (x-\mu_0)^t\Sigma^{-1}(x-\mu_0)\right]\right)$$
$$= \ln \exp\left((\mu_0^t - \mu_1^t)\Sigma^{-1}x + \frac{1}{2}\mu_1^t\Sigma^{-1}\mu_1 - \frac{1}{2}\mu_0^t\Sigma^{-1}\mu_0\right)$$

Finally we got:

$$\mathbb{P}(Y = 1|X) = \frac{1}{1 + \exp\left(\ln\frac{1-\pi}{\pi} + \frac{1}{2}\mu_1^t \Sigma^{-1} \mu_1 - \frac{1}{2}\mu_0^t \Sigma^{-1} \mu_0 + (\mu_0^t - \mu_1^t) \Sigma^{-1} x\right)}$$
$$= \frac{1}{1 + \exp\left(w_0 + w^t x\right)}$$

where :

$$\begin{split} w_0 &= \ln \frac{1-\pi}{\pi} + \frac{1}{2} \mu_1^t \Sigma^{-1} \mu_1 - \frac{1}{2} \mu_0^t \Sigma^{-1} \mu_0 \\ w &= \left((\mu_0^t - \mu_1^t) \Sigma^{-1} \right)^t = \Sigma^{-1} (\mu_0 - \mu_1) \end{split} \quad \text{because the covariance matrix is symmetric}$$

Conclusion:

 \clubsuit When conditional independence of the random variables in $X = < X_1, \cdots, X_n >$ is multivariate normal with mean depending on Y, but covariance matrix now NOT depending on Y the $\mathbb{P}(X|Y)$ have the form of logistic regression