# **Introduction to Machine Learning**

2010

Prof. J. M. Buhmann

#### **Final Exam**

February 10th, 2010

First and Last name:	
ETH number:	
Signature:	

# **General Remarks**

- Please check that you have all 21 pages of this exam.
- Remove all material from your desk which is not permitted by the examination regulations.
- Fill in your first and last name and your ETH number and sign the exam. Place your student ID in front of you.
- You have 2 hours for the exam. There are five questions, where you can earn a total of 150 points. Scoring 120 points (equivalent to solving four questions) guarantees you a grade of six.
- Write your answers directly on the exam sheets. If you need more space, put your name and ETH number on top of each supplementary sheet.
- Answer the questions in English. Do not use a pencil or red color pen.
- You may provide at most one valid answer per question. Invalid solutions must be canceled out clearly.

	Topic	Max. Points	Points Achieved	Visum
1	Density Estimation	30		
2	PCA	30		
3	Classifiers	30		
4	Regression	30		
5	Boosting	30		
Total		150		

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### Question 1: Parametric Density Estimation (30 pts.)

When estimating densities from data, we assume that the data distribution is well approximated by a parametric model  $P(x; \theta)$ , where  $\theta$  denotes the parameters of the model. Estimation of the density then reduces to estimation of the parameter  $\theta$ .

a) In maximum-likelihood estimation (MLE), we usually start by assuming that the data is given as an i.i.d. sample  $x_1, x_2, ..., x_N \sim P(x; \theta)$ . Please explain how this assumption simplifies the calculations. Answer:

2 pts.

b) We model a coin flip as a random variable X, where X=1 denotes the outcome 'heads', and X=0 denotes 'tails'. We model this event using the Bernoulli distribution, parametrized by  $\mu$  that denotes P(X=1).

Bern
$$(x|\mu) = \mu^x (1-\mu)^{1-x}$$
.

Given a data set  $\mathcal{X}_N = \{x_1, ..., x_N\}$  of N i.i.d. observations of X, show us step by step how to compute the MLE for  $\mu$ . Comment your calculations where necessary. Answer:

c'	We flip	our	coin	three	times	and	obtain	the	dataset	$\chi_2 = -$	{1	1 1	1}
٠,	we mp	Oui	COIII	tillee	LIIIICS	anu	Obtain	LIIC	uataset	$\alpha_3 - \gamma$	ι Ι,	$_{1}$	١ſ٠

- 1. What is the MLE for  $\hat{\mu}$  for dataset  $\mathcal{X}_3$ ? Answer:
- 2. State one problem with this estimate. What is a possible remedy (solution)? Answer:

3 pts.

- d) Assume that there are a total of m heads in  $\mathcal{X}_N$ , i.e.,  $\sum_{i=1}^N x_i = m$ .
  - 1. In terms of  $\mu, N$  and m only, what is the probability of obtaining this dataset  $P(\mathcal{X}_N|\mu)$ ? Explain the terms of your solution. Answer:

2. What is the expected value of m? (Hint: For independent random variables, the expectation of a sum is the sum of expectations)

Answer:

$$\mathbb{E}[m] =$$

We now turn to Bayesian estimation of  $\mu$  given a data set  $\mathcal{X}$ , by introducing a prior distribution  $P(\mu)$ . Here, we use the (conjugate) prior given by the Beta distribution:

$$\mathsf{Beta}(\mu|\alpha,\beta) = \frac{\Gamma(\alpha+\beta)}{\Gamma(\alpha)\Gamma(\beta)} \mu^{\alpha-1} (1-\mu)^{\beta-1},$$

where  $\Gamma(x)$  denotes the gamma function (not necessary). Some useful properties:

- 1. Mean  $\mathbb{E}(\mu) = \frac{\alpha}{\alpha + \beta}$
- 2. Variance  $\mathbb{V}(\mu) = \frac{\alpha\beta}{(\alpha+\beta)^2(\alpha+\beta+1)}$
- e) With the mean and variance of  $\mu$ , we can understand how the parameters influence the prior. Select the best  $(\alpha_0, \beta_0)$  pair for the prior that emphasizes the following beliefs.
  - 1. My coin is biased towards 'tails'.

 $\Box(0.5, 0.5)$ 

 $\square(1,4)$ 

 $\square(4,1)$ 

 $\square(2,2)$ 

2. My coin is most probably biased (but I don't know heads or tails)

 $\Box(0.5, 0.5)$ 

 $\square(1,4)$ 

 $\square(4,1)$ 

 $\square(2,2)$ 

3 pts.

f) Show that if the prior is Beta distributed

$$P(\mu) = \text{Beta}(\alpha_0, \beta_0),$$

then the posterior takes again the same functional dependency on  $\mu$  as the prior, i.e.,

$$P(\mu|\mathcal{X}_N) \propto \mu^{\alpha_0 + m - 1} (1 - \mu)^{\beta_0 + N - m - 1}.$$

Answer:

- g) In fact, the posterior is again a Beta distribution, and we obtain  $P(X=1|\mathcal{X})=\frac{m+\alpha_0}{N+\alpha_0+\beta_0}.$ 
  - 1. For  $\mathcal{X}_3=\{1,1,1\}$ , use a choice of  $(\alpha_0,\beta_0)$  which assumes a fair coin, and obtain your estimate for  $P(X=1|\mathcal{X})$ . How does this estimate differ from the MLE? Answer:

2. What happens to the posterior as  $N \to \infty$ ? How does this relate to the MLE? Answer:

### Question 2: Principal Component Analysis (PCA) (30 pts.)

Principal component analysis is a widely applied method in data analysis and dimensionality reduction. Given a dataset  $\mathbf{X} \in \mathbb{R}^{D \times N}$  (observations as columns), where D is the number of dimensions and N is the number of observations, a linear transformation using an orthonormal matrix  $\mathbf{P} \in \mathbb{R}^{M \times D}$  is applied to make a *change of basis*, to obtain a (usually) lower-dimensional dataset  $\mathbf{Y} \in \mathbb{R}^{M \times N}$ .

a)	) We begin	by revie	wing the	steps of	applying	PCA	to a	dataset	X. Pl	ease (	complete	each
	step belov	ν by prov	iding the	appropri	ate formu	ıla to	comp	oute the	desired	l quar	ntity.	

1.	Define the zero-mean dataset $ ilde{\mathbf{X}}$ in terms of the original dataset $\mathbf{X}$ . For this purpose,
	introduce the data mean as a column-vector $oldsymbol{\mu} \in \mathbb{R}^{D  imes 1}$ .
	Answer:

2.	Define the	covariance	matrix	$\mathbf{C}_{\tilde{\mathbf{X}}}$	in	terms	of	the	zero-	mean	${\sf dataset}$	$\tilde{\mathbf{X}}$ .
	Answer:											

3.	Write d	lown	the	eiger	า-deco	mposi	ition	of	the	cova	riance	matr	x <b>C</b>	$\hat{\mathbf{X}}$ ,	in	terms	of	the
	eigenve	ctor r	natri	$\mathbf{E}$	(eigen	vector	's as	colı	umn	s) an	d the	diagor	ıal n	natri	хо	f eige	nva	lues
	D.																	

Answer:

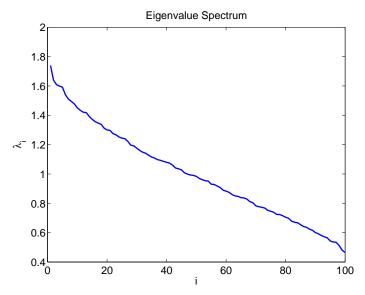
4. Define the PCA transformation matrix  ${\bf P}$  in terms of the eigenvector matrix  ${\bf E}$ . (Note: assume we don't want to reduce the dimensionality of the transformed dataset. Further assume that the eigenvectors in  ${\bf E}$  have already been sorted according to the corresponding eigenvalues, in decreasing order.)

Answer:

5. Define the transformed dataset  ${\bf Y}$  in terms of the original dataset  ${\bf X}$  and the transformation matrix  ${\bf P}$ .

Answer:

b) Assume we have applied PCA to some dataset (D=100). We observe the following eigenvalue spectrum of the covariance matrix of the data. ( $\lambda_i$ : eigenvalues)



1. Is the intrinsic dimensionality of this dataset low or high? Why? Answer:

2. Can this dataset be expressed in few dimensions with low approximation error? Why? Answer:

3. If yes, which dimensionality (approximately) should be chosen for the transformed dataset and why?

Answer:

c) Assume you have observed 2D data  $\mathbf{X} \in \mathbb{R}^{2 \times N}$  (observations as columns). The first row of  ${\bf X}$  corresponds to the first dimension  $x_1$ , the second row corresponds to  $x_2$ . For each of the three covariance matrices  $\mathbf{C}_{\mathbf{X}}$  below, please choose the iso-line plot (A-E) corresponding to the covariance matrix. (Note the axis labels on the figures)

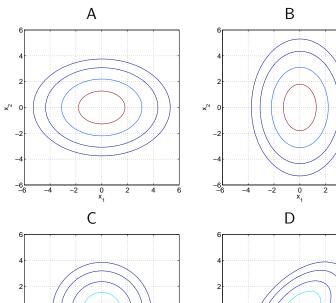
$$1. \left[ \begin{array}{cc} 1 & 0 \\ 0 & 2 \end{array} \right]$$

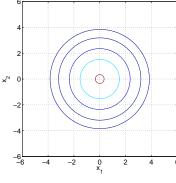
Answer: ( )

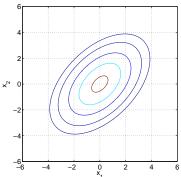
2. 
$$\begin{bmatrix} 1 & -0.5 \\ -0.5 & 1 \end{bmatrix}$$
 Answer: ( )

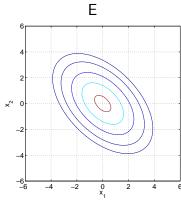
$$3. \left[ \begin{array}{cc} 1 & 0 \\ 0 & 1 \end{array} \right]$$

Answer: ( )









- d) PCA can be derived as follows: Given a dataset  $\mathbf{X} = \{\mathbf{x}_1, \mathbf{x}_2, ..., \mathbf{x}_N\} \in \mathbb{R}^{D \times N}$ , find a new set of basis vectors such that the variance of the projected dataset is maximized. We begin by finding the first projection direction  $\mathbf{e}_1 \in \mathbb{R}^D$ ,  $||\mathbf{e}_1||_2 = 1$ , which satisfies our goal.
  - 1. Please describe in words why we are interested to find a projection direction such that the projected data has large variance.

    Answer:

2. Write down formally the objective function for maximizing the variance of the projected dataset. For this purpose, define the variance of the projected dataset in terms of the original data  $\mathbf{x}_n$  and the first projection direction  $\mathbf{e}_1$ . (Hint: introduce the mean  $\boldsymbol{\mu}$  of the data  $\mathbf{X}$ .)

Answer:

e)	PCA transforms a dataset $X$ into a dataset $Y$ by defining a new basis using the eigenvectors
	of the covariance matrix $\mathbf{C}_{\mathbf{X}}$ of the dataset $\mathbf{X}$ . With this particular choice of a new basis,
	the covariance matrix $\mathbf{C}_{\mathbf{Y}}$ of the transformed dataset $\mathbf{Y}$ is diagonalized. (Note: For this
	exercise, assume a zero-mean dataset.)

1. Please explain in words, why we desire the covariance matrix of the transformed dataset to be diagonal.

Answer:

2. Show that  $C_Y = PC_XP^T$ , i.e., that the covariance matrix  $C_Y$  of the transformed dataset can be written in terms of the covariance matrix  $C_X$  of the original dataset. Answer:

3. Show that the choice  $\mathbf{P}=\mathbf{E}^T$  for the PCA transformation matrix, where  $\mathbf{E}$  is the matrix of eigenvectors of the covariance matrix  $\mathbf{C}_{\mathbf{X}}$ , actually diagonalizes the covariance matrix  $\mathbf{C}_{\mathbf{Y}}$  of the transformed dataset  $\mathbf{Y}$ . Use the eigen decomposition of  $\mathbf{C}_{\mathbf{X}}$  and the fact that  $\mathbf{P}^{-1}=\mathbf{P}^T$ .

Answer:

## Question 3: Linear Classifiers (30 pts.)

The simplest representation of a linear discriminant function is obtained by taking a linear function of the input vector x so that

$$y(\boldsymbol{x}) = \boldsymbol{w}^T \boldsymbol{x} + b$$

a) Define the term decision boundary. Explain what is the relation between the vector  ${\pmb w}$  and the decision boundary. Answer:

2 pts.

b) Consider the following set  $S \subseteq \mathbb{R}^2 \times \{+1, -1\}$ :

$$S = \{((3,3),+1),((1,4),+1),((2,1),-1),((5,1),-1),((6,2),+1),((4,3),+1),((4,1),-1)\}$$

Prove that there exist no homogeneous half-space that separates S with no errors (recall that a half-space is homogeneous if its separating hyperplane passes through the origin). Assume that the classification of a point is given by  $y(\boldsymbol{x}) = \text{sign}(\boldsymbol{w}^T\boldsymbol{x} + b)$ . Answer:

The simple form of the online perceptron algorithm is defined as follows:

Initialize by setting  ${m w}^{(0)}={m 0}.$ 

$$\mathsf{Update}\ \boldsymbol{w}^{(t+1)} = \begin{cases} \boldsymbol{w}^{(t)} & y_k(\boldsymbol{w}^{(t)^T}\boldsymbol{x}_k) > 0 \\ \boldsymbol{w}^{(t)} + y_k\boldsymbol{x}_k & \mathsf{otherwise}. \end{cases}$$

c) We would like to use the perceptron algorithm to find a separating half-space that makes no errors on the set S. Recall that you have just proven that this is impossible using a homogeneous half-space. Make the necessary adjustments and perform 7 iterations of the perceptron algorithm. For each iteration write down the training example and the current vector  $\boldsymbol{w}^{(t)}$ . Assume that the points are given in the same order as they appear in S. Answer:

5 pts.

Consider the primal form of the soft margin SVM

$$\min \frac{1}{2} \|\boldsymbol{w}\|^2 + C \sum_{i=1}^n \xi_i$$
s.t.  $y_i(\boldsymbol{w}^T \boldsymbol{x}_i + b) \ge 1 - \xi_i, \quad i = 1, \dots, n$ 

- d) Indicate which of the following statements hold as we increase the value of the parameter C (relative to any starting value, C>0 ). Use
  - D if the validity of the statement depends on the situation,
  - T if the statement is necessary true
  - N if the statement is never true
  - b will not increase

  - | | | | w | | will not decrease
  - more points will be misclassified
  - [ ] The geometric margin will not increase

The dual soft margin SVM is given below

$$\max \sum_{i=1}^{n} \alpha_i - \frac{1}{2} \sum_{i,j=1}^{n} \alpha_i \alpha_j y_i y_j K(\boldsymbol{x}_i, \boldsymbol{x}_j)$$
s.t.  $0 \le \alpha_i \le C, \quad i = 1, \dots, n$ 

We would like to solve the SVM problem using the kernel  $K(\boldsymbol{x}_i, \boldsymbol{x}_j) = 1 + (\boldsymbol{x}_i^T \boldsymbol{x}_j)^2$ .

e) What is the feature representation  $\phi(x)$  corresponding to this kernel for  $x \in \mathbb{R}^2$  ? Answer:

2 pts.

f) Using the dual variables  $\alpha$ 's and  $\phi(\boldsymbol{x})$ , give an explicit form for the primal variable  $\boldsymbol{w}$ . Answer:

2 pts.

We solved the dual SVM problem with the above kernel. We used the following training set  ${m x}_1=(1,1), y_1=+1$ ,  ${m x}_2=(3,0), y_2=-1$ ,  ${m x}_3=(-2,1), y_3=-1$ . Below are the  $\alpha$ 's and b values  $\alpha_1=0.1, \ \alpha_2=0, \ \alpha_3=0.1, \ b=1.5$ 

g) Use these values to compute the kernelized discriminant function  $y(\boldsymbol{x})$ . Answer:

### Question 4: Regression (30 pts.)

In univariate regression analysis, we study the statistical dependency of the output variable  $y \in \mathbb{R}$  on the input variable  $x \in \mathbb{R}$ .

a) Complete the equation

$$y = \qquad \qquad + \tag{1}$$

where the right hand side consists of two terms. Define each term in the equation.

Note: Be general, don't restrict yourself to linear least-squares regression.

Answer:

3 pts.

- b) Regression analysis is only meaningful if there actually is a true dependency between x and y. Prove that if x and y are statistically independent, there is no linear dependency:
  - 1. Assume that x and y are statistically independent.
  - 2. Begin with the definition of covariance,

$$cov(x, y) = \mathbb{E}\left[\left(x - \mu_x\right)\left(y - \mu_y\right)\right],\,$$

where  $\mathbb E$  is the expectation operator and  $\mu_x$  is the mean of x.

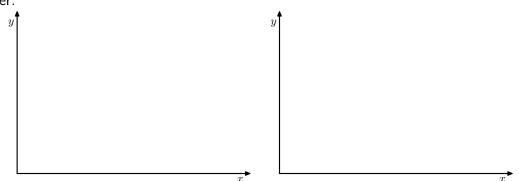
3. Show that x and y are uncorrelated.

Note: Comment your calculations where necessary.

Answer:

- c) You have shown that eq. (1) consists of two terms. Illustrate two conceptually different cases where a linear least-squares fit is inappropriate for the true dependency between x and y, one related to the first term in eq. (1), and one related to the second term in eq. (1):
  - 1. Draw a scatter plot of a data sample.
  - 2. Fit the linear least-squares solution into the sample (qualitatively).
  - 3. Explain why the solution does not capture the true dependency between x and y.

Answer:



4 pts.

So far, we considered a single input variable x. In multivariate linear regression, we assume a weighted linear functional dependency

$$y = w_0 + w_1 x_1 + \dots + w_D x_D \tag{2}$$

between D input variables  $x_1, \ldots, x_D \in \mathbb{R}$  and output variable y.

d) Define the input data matrix  $\mathbf{X}$ , weight vector  $\mathbf{w}$  and output data vector  $\mathbf{y}$  such that eq. (2) for all N data points of the sample  $(x_1^n, x_2^n, \dots, x_D^n, y^n), n = 1, \dots, N$  can be written as

$$y = Xw$$
.

Answer:

$$\mathbf{X} := \mathbf{w} := \mathbf{y} :=$$

The linear least-squares prediction  $\hat{\mathbf{y}}$  is the *orthogonal projection* of  $\mathbf{y}$  onto the space spanned by the columns of  $\mathbf{X}$ , and can be written as

$$\hat{\mathbf{y}} = \mathbf{P}\mathbf{y}$$
.

- e) Analytically derive  ${\bf P}$ , where  $\hat{{\bf y}}$  is the least-squares prediction vector. To achieve this:
  - 1. Derive the optimal weight vector  $\hat{\mathbf{w}}$  that minimizes the least-squares objective function

$$\|\mathbf{X}\mathbf{w} - \mathbf{y}\|_2^2 \tag{3}$$

Note: we have  $\frac{\partial}{\partial \mathbf{w}} \mathbf{w}^{\top} \mathbf{S} \mathbf{w} = 2 \mathbf{S} \mathbf{w}$  for symmetric matrix  $\mathbf{S}$ , and  $\frac{\partial}{\partial \mathbf{w}} \mathbf{w}^{\top} \mathbf{a} = \mathbf{a}$ .

2. Using your result for  $\hat{\mathbf{w}}$ , give the definition for  $\mathbf{P}$ .

Answer:

5 pts.

f) Prove that  ${\bf P}$  is an orthogonal projector, by showing that  ${\bf P}={\bf PP}$  and  ${\bf P}={\bf P}^\top.$  Answer:

g)	Computing $\hat{\mathbf{w}}$ involves the inversion of the symmetric matrix $\mathbf{S} := \mathbf{X}^{\top}\mathbf{X}$ , which we assume
	to be positive definite. Give a mathematical condition (in terms of the eigenvalues of S),
	when this inversion is numerically unstable.
	Answer:

2 pts.

h) Ridge regression finds the regularized weight vector  $\hat{\mathbf{w}}_{ridge}$  that minimizes eq. (3) with an additional norm penalty,

 $\|\mathbf{X}\mathbf{w} - \mathbf{y}\|_{2}^{2} + \lambda \|\mathbf{w}\|_{2}^{2}$ .

Explain why choosing  $\lambda>0$  improves stability of the inversion:

1. Derive the regularized optimal weight vector  $\hat{\mathbf{w}}_{\mathrm{ridge}}.$  Answer:

2 pts.

- 2. What is the effect of increasing  $\boldsymbol{\lambda}$  on
  - (a) the norm  $\left\|\hat{\mathbf{w}}_{ridge}\right\|_2$  and
  - (b) the numerical stability of computing  $\hat{\mathbf{w}}_{\mathrm{ridge}}?$

Provide **arguments** for your answers.

Answer:

	uestion 5: Bagging and AdaBoost (30 pts.)	
	State two essential differences between bagging and AdaBoost.  Answer:	
b)	Explain how bagging reduces the variance of an estimator?	2 pts
٠,	Answer:	
c)	State two factors that may influence the performance of AdaBoost? Answer:	3 pts
		2 pts

2 pts.

d) How can we use AdaBoost to detect outliers (mislabeled or ambiguously labeled examples)?

Answer:

For a binary classification problem ( $\pm 1$ ), assume we have  $c_1,...,c_B$  successive base classifiers obtained via AdaBoost, and let  $\alpha_1,...,\alpha_B$  be the corresponding weights. The ensemble classifier is  $\hat{c}_B(x) = \operatorname{sgn}(\sum_{b=1}^B \alpha_b c_b(x))$ .

Reusing the same base classifiers  $c_1,...,c_B$ , we now train a *linear* classifier, by finding *new*  $\tilde{\alpha}_b$ 's that minimize the average exponential loss on the training examples, and obtain  $\tilde{\hat{c}}_B(x) = \operatorname{sgn}(\sum_{b=1}^B \tilde{\alpha}_b c_b(x))$ .

e) Is  $\tilde{\hat{c}}_B$  equivalent to  $\hat{c}_B$  ? Justify your answer. Answer:

4 pts.

Consider the following simplified ensemble learning algorithm:

We consider a dataset  $\{(x_1,y_1),(x_2,y_2),\cdots,(x_N,y_N)\}$  of N i.i.d. samples, where  $y_n\in\{-1,+1\}$ . After each iteration b of this algorithm, each training sample  $x_n$  has an associated weight  $w_{b,n}$ . At the beginning,  $w_{0,n}=1$  for n=1,...,N. Let  $0<\gamma<\frac{1}{2}$  and  $\beta=\frac{1+2\gamma}{1-2\gamma}$ . For each iteration b=1 to B:

- The algorithm finds a new weak learner  $c_b$  on the weighted training set that is guaranteed to have an error of at most  $\frac{1}{2} \gamma$  on the weighted training set.
- We then update the weights of the training samples:
  - If  $x_n$  is misclassified by  $c_b$ , i.e., if  $c_b(x_n) \neq y_n$ : Set  $w_{b,n} := \beta w_{b-1,n}$ .
  - Otherwise the weight remains unchanged:  $w_{b,n} := w_{b-1,n}$ .

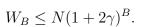
The ensemble classifier decides by combining the weak learners as follows:

$$\hat{c}_B(x) = \operatorname{sgn}(\sum_{b=1}^B c_b(x)).$$

Let  $W_b = \sum_{n=1}^N w_{b,n}$  be the sum of the weights at the end of iteration b.

f) What is the difference between AdaBoost and this algorithm in the way they combine the weak learners? Answer:

g)	Derive a bound on $W_B$ , the sum of the weights after the final iteration. Given that at each iteration $b$ , the sum of the weights increases by at most a factor of $1+2\gamma$ :
	$\frac{W_b}{W_{b-1}} \le 1 + 2\gamma,$
	show that for the final sum of weights it holds that



Answer:

4 pts.

h) Prove the following lower bound on the final weight of a training sample, if it is misclassified by the ensemble:

If 
$$\hat{c}_B(x_n) \neq y_n$$
, then  $w_{B,n} \geq \beta^{B/2}$ .

Answer:

- i) Finally, we will show a bound on the empirical error rate on the training data. Let  $M = |\{n : \hat{c}_B(x_n) \neq y_n\}|$  be the number of training samples misclassified by the ensemble classifier.
  - 1. Plugging the two previous bounds together, we have

$$M\beta^{B/2} \le \sum_{n=1}^{N} w_{n,B} = W_B \le N(1+2\gamma)^B.$$

Starting from this inequality, show that the empirical error rate on the training data is bounded by the following:

$$\frac{M}{N} \le e^{-2B\gamma^2}.$$

(Hint: note that  $(1+a)^c \le e^{ac}$  for  $c \ge 0$ , and recall that  $\beta = \frac{1+2\gamma}{1-2\gamma}$ ) Answer:

- 2. Explain what happens to this bound when we
  - (a) increase  $\boldsymbol{B}$  (the number of training iterations),
  - (b) increase  $\gamma$  (have a very good learner).

Answer:

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Supplementary Sheet

Supplementary Sheet