CLASSIFICATION ASYMPTOTICS IN THE RANDOM MATRIX REGIME

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ABSTRACT

This article discusses the asymptotic performance of classical machine learning classification methods (from discriminant analysis to neural networks) for simultaneously large and numerous Gaussian mixture modelled data. We first provide theoretical bounds on the minimally discriminable class means and covariances under an oracle setting, which are then compared to recent theoretical findings on the performance of machine learning. Non-obvious phenomena are discussed, among which surprising phase transitions in the optimal performance rates for specific hyperparameter settings.

Index Terms— Random matrix theory, classification, kernel methods, neural networks, LDA/QDA.

1. INTRODUCTION

The renewed interest for machine learning spurred by the big data movement leads statisticians to reconsider the asymptotic performance of statistical learning when both the number of data \( n \) and their dimension \( p \) grow simultaneously large, i.e., under a random matrix regime. In this regime, many conventional estimators, starting with the sample covariance matrix, are known to become inconsistent [1, 2, 3]. The field of random matrix theory has provided a deep understanding on the limitations and possible corrections of such “large \( p \), large \( n \)” statistics, yet mostly for linear estimators (such as in array processing [4, 5], wireless communications [6, 7], detection and estimation [8], etc.). Dealing with non-linear operators, as in machine learning (kernel methods, neural nets), is more challenging and only very recent works have provided first steps into understanding and improving machine learning for large dimensional data [9, 10, 11].

The article aims to extract from these recent articles conclusions in terms of optimal discriminative performance rates in the classification of Gaussian mixture models and compare them to the optimal oracle performance of the Neyman–Pearson test for known class distributions.

2. SYSTEM SETTING AND OBJECTIVES

Let \( x_1, \ldots, x_n \in \mathbb{R}^p \) be a set of vectors to be classified, either in a supervised, semi-supervised, or unsupervised manner. For simplicity, we consider a binary setting: \( x_1, \ldots, x_{n_1} \) belong to class \( C_1 \) and \( x_{n_1+1}, \ldots, x_n \) to class \( C_2 \), with \( n = |C_1| \).

We further restrict ourselves to the setting where the \( x_i \)’s are independent and arise from a Gaussian mixture model, i.e., for \( \mu_1, \mu_2 \in \mathbb{R}^p \) and nonnegative definite \( C_1, C_2 \in \mathbb{R}^{p \times p} \),

\[
    x_i \in C_a \iff x_i \sim N(\mu_a, C_a).
\]

We will discuss the asymptotic performance of various classification methods as \( n, p \to \infty \) with \( p/n \to c_0 > 0 \), our objective being twofold:

1. we first characterize the fundamental limits of Gaussian mixture classification under the oracle setting where \( \mu_1, \mu_2 \) and \( C_1, C_2 \) are known; these limits are expressed in terms of the minimal “distance scaling rate” between \( \mu_1 \) and \( \mu_2 \), and \( C_1 \) and \( C_2 \) with respect to \( n, p \) which ensures asymptotically non trivial classification (i.e., neither perfect nor impossible);

2. we compare the asymptotic oracle performance to those achieved by standard classification methods (discriminant analysis, kernel methods, random neural networks) with various degrees of available samples (from supervised classification to clustering).

The random matrix regime is interesting as it corresponds to the rate where \( C_1, C_2 \) can no longer be consistently estimated by sample covariances. As such, simple plug-in estimators (as we shall see with the popular QDA method) are bound to induce large performance losses with respect to oracle and thus possibly a loss in the minimally discriminable rates for \( \| \mu_1 - \mu_2 \| \) and \( \| C_1 - C_2 \| \). But random matrix theory has for long provided a series of consistent estimators for functionals of \( C_1, C_2 \) (i.e., mappings from \( \mathbb{R}^{p \times p} \) to \( \mathbb{R} \)), some of which, as we shall see, are at the core of most well-performing classifiers; these classifiers often maintain close-to-optimal discriminative power.

\(^1\)But many results to be introduced in the following are valid for a much larger scope of statistical distributions.
3. MAIN RESULTS

To avoid unnecessary complications, we make the following base hypotheses.

Assumption 1 (Growth Rate Control). The matrices $C_1, C_2$ are invertible and, as $p \to \infty$, for $\parallel \cdot \parallel$ the operator norm,

$$\liminf_{p \to \infty} \{\|C_a\|, \|C^{-1}_a\|\} < \infty \quad \text{for } a \in \{1, 2\}.$$  

When convenient, we further use the shortcut definitions

$$\mu \equiv \mu_1 - \mu_2, \quad \Delta C \equiv C_1 - C_2.$$  

3.1. Oracle Classification

Under the setting of Section 2, let $x \in \mathbb{R}^p$ be a vector genuinely belonging to class $C_1$, i.e., $x \sim N(\mu_1, C_1)$. Further assume for simplicity that, when observed, $x$ has prior probability $\frac{1}{2}$ to belong to class $C_1$. Then, for perfectly known $\mu_1, \mu_2$ and $C_1, C_2$, the (decision optimal) Neyman–Pearson test for its belonging to $C_1$ consists in the comparison test

$$(x - \mu_2)^T C_2^{-1}(x - \mu_2) - (x - \mu_1)^T C_1^{-1}(x - \mu_1) \leq \log \frac{|C_1|}{|C_2|}. \quad (1)$$  

Writing $x = \mu_1 + C_1^{\frac{1}{2}}w$, so that $w \sim N(0, I_p)$, after elementary manipulation, the test is equivalent to the comparison

$$S(x) \equiv \frac{1}{p} \text{tr}(C_1^{-1} C_2^{\frac{1}{2}} - I_p)w + \frac{2}{p} \mu^T C_2^{-\frac{1}{2}} C_1^{-\frac{1}{2}} w$$

$$\quad + \frac{1}{p} \mu^T C_2^{-1} \mu - \frac{1}{p} \log \frac{|C_1|}{|C_2|} \geq 0. \quad (2)$$  

Since $Uw$ for $U \in \mathbb{R}^{p \times p}$ an eigenvector basis for $C_1^{\frac{1}{2}} C_2^{-\frac{1}{2}} C_1^{\frac{1}{2}} - I_p$ has the same distribution as $w$, the random variable $S(x)$ can be written as the sum of $p$ independent random variables. A careful application of Lyapunov’s central limit theorem [12], along with Assumption 1, reveals that, as $p \to \infty$,

$$V_S^{-\frac{1}{2}}(S(x) - \bar{S}) \to \mathcal{N}(0, 1)$$

in distribution, where

$$\bar{S} \equiv \frac{1}{p} \text{tr}(C_1^{-1} C_2^{\frac{1}{2}} - I_p) - \frac{1}{p} \mu^T C_2^{-1} \mu - \frac{1}{p} \log \frac{|C_1|}{|C_2|} \quad (3)$$

$$V_S \equiv \frac{2}{p^2} \text{tr}(C_1^{\frac{1}{2}} C_2^{\frac{1}{2}} C_1^{\frac{1}{2}} - I_p)^2 + \frac{2}{p^2} \mu^T C_2^{-1} C_1^{\frac{1}{2}} - I_p \mu. \quad (4)$$

The classification performance of $x$ is thus only non-trivial (e.g., converging neither to 0 nor 1) if both $\bar{S}$ and $\sqrt{V_S}$ are of the same order of magnitude (with respect to $p$). Assume first that $C_1 = C_2 \equiv C$; then from Assumption 1,

$$\bar{S} = \frac{1}{p} \mu^T C^{-1} \mu = O(\|\mu\|^2/p)$$

$$\sqrt{V_S} = \frac{2}{p} \sqrt{\mu^T C^{-1} \mu} = O(\|\mu\|/p).$$

As a consequence, for a non-trivial asymptotic classification, we demand here that in the worst case, in order of magnitude, $\|\mu\| \sim_p \|\mu\|^2$, i.e., $\|\mu\| = O(1)$. The associated asymptotic correct classification probability $P(S(x) > 0)$ thus satisfies

$$P(S(x) > 0) \sim Q \left(-\frac{1}{2} \sqrt{\mu^T C^{-1} \mu} \right) \to 0$$

with $Q(t) \equiv \frac{1}{\sqrt{2\pi}} \int_{t}^{\infty} e^{-\frac{1}{2}u^2} du$ the Gaussian tail function.

Under this minimal constraint $\|\mu\| = O(1)$, we now incorporate $C_1 \neq C_2$ in such a way that $\|\Delta C\| = o(1)$. A Taylor expansion of (3)–(4) around $C = C_1$ leads to

$$\bar{S} = \frac{1}{p} \mu^T C^{-1} \mu + \frac{1}{2p} \text{tr}(C^{-1} \Delta C)^2 + o(p^{-1})$$

$$V_S = \frac{2}{p^2} \text{tr}(C^{-1} \Delta C)^2 + \frac{4}{p^2} \mu^T C^{-1} \mu + o(p^{-2})$$

from which it unfolds that $\|\Delta C\|$ must be at least of order $O(p^{-\frac{1}{2}})$ (so that $\text{tr}(C^{-1} \Delta C)^2 = O(1)$) for the difference $\Delta C$ to have discriminative power. The associated asymptotic classification probability in this case satisfies

$$P(S(x) > 0) \sim Q \left(-\frac{\sqrt{\mu^T C^{-1} \mu} + \frac{1}{2} \text{tr}(C^{-1} \Delta C - I_p)^2}{2} \right).$$

Consequently, when $\mu_1, \mu_2$ and $C_1, C_2$ are a priori known, non-trivial classification is achieved for

$$\|\mu\| \geq O(1), \quad \|\Delta C\| \geq O(p^{-\frac{1}{2}})$$

where $X_p \geq O(p^\alpha)$ stands for $\liminf_p p^{-\alpha} X_p > 0$. This forms an optimal asymptotic baseline for Gaussian mixture classification that statistical learning methods cannot outperform. In the following, we discuss the performances of a class of such methods, retrieved from a series of recent articles based on advanced random matrix considerations.

3.2. Asymptotic Performance of Machine Learning Tools

3.2.1. Supervised learning: LDA and QDA

In [13], the authors consider the popular linear and quadratic discriminant analysis (LDA and QDA) supervised classifiers. Those classifiers assume a two-class Gaussian mixture model for the data, as introduced presently. In its most general form (QDA), the classifier consists in estimating $\mu_a$’s and $\Sigma_a$’s by sample means and covariances from two sets of data labelled in class $C_1$ or $C_2$. The estimators are then plugged in (1) in place of the genuine parameters. The respective numbers $n_1$ and $n_2$ of data from class $C_1$ and $C_2$ are such that $n_1, n_2 = O(p)$, with $n_1, n_2 > 0$ (in order for sample covariance matrices to be invertible with probability one).

In [13], it is shown that the noise induced by the estimation of $C_1$ and $C_2$ has a debilitating effect on classification. Indeed, in order to perform the comparison (2),
$S(x)$ needs to be consistently estimated. If $x = \mu_1 + C_1^\sharp w$ with $||\mu|| = O(1)$, then $\frac{1}{n}x^T C_1^{-1}x$ is a consistent estimator for $\frac{1}{n}w^T C_1^\sharp C_1^{-1}C_1^\sharp w$; but since $C_2$ is unknown and that $\hat{C}_2 = \frac{1}{n} \sum_{i=n_1+1}^{n} (x_i - \bar{x})(x_i - \bar{x})^T (\bar{x} = \frac{1}{n_2} \sum_{j=n_1+1}^{n_2} x_j)$ is not a consistent estimator for $C_2$, $\frac{1}{n}x^T C_2^{-1}x$ is a biased estimator for $\frac{1}{n}w^T C_2^\sharp C_2^{-1}C_2^\sharp w$. Nonetheless, random matrix theory provides advanced results to correct these biases, and in particular $(1 - \xi)\frac{1}{n}x^T C_2^{-1}x$ consistently estimates $\frac{1}{n}w^T C_2^\sharp C_2^{-1}C_2^\sharp w$. The major issue though, is that $(1 - \xi)\frac{1}{n}x^T C_2^{-1}x = \frac{1}{p} w^T C_2^\sharp C_2^{-1}C_2^\sharp w + O(p^{-\frac{1}{2}})$. Similarly, $(1 - \xi)\frac{1}{n}x^T C_1^{-1}x = \frac{1}{p} w^T I_p w + O(p^{-\frac{1}{2}})$. But under minimal discriminative oracle rates, $\frac{1}{p} w^T (C_2^\sharp C_2^{-1}C_2^\sharp - I_p) w = O(p^{-\frac{1}{2}})$ so that the estimator $(1 - \xi)\frac{1}{n}x^T (C_1^{-1} - C_2^{-1})x$ induces a dominating $O(p^{-\frac{1}{2}})$ noise. This noise also dominates the term $\frac{1}{p} \mu^T C_2^\sharp C_2^{-1} \mu$ in $S(x)$. As a consequence, to obtain a non-trivial classification rate, either constraint $||\Delta C|| = O(p^{-\frac{1}{2}})$ or $||\mu||^2 = O(1)$ must be relaxed to at least an order of magnitude $\sqrt{p}$. The minimal distance scaling for non-trivial classification in this case becomes

$$||\mu|| \geq O(p^{\frac{1}{2}}), \quad \text{or} \quad ||\Delta C|| \geq O(1).$$

**Remark 1** (Noise dominance when $\Delta C = 0$). Interestingly here, even when $\Delta C = 0$ (a property not known in advance), the minimal distance scaling of $||\mu||$ for non-trivial classification with QDA is still $||\mu|| = O(p^{\frac{1}{2}})$, due to the non-discriminative, yet performed, estimation of $C_1, C_2$.

When it comes to the simpler linear discriminant approach (LDA), meant to discriminate a Gaussian mixture with common covariances $C_1 = C_2 \equiv C$, a surprising phenomenon occurs. There, while $C$ is still estimated by the sample covariance, this estimator no longer degrades the minimal distance rate $||\mu|| = O(1)$. This result is due to the fact that, for LDA, only the sign of $(w - \frac{1}{n}\mu)^T C^{-1} \mu = O(1)$ must be estimated. Substituting $(1 - \frac{1}{n})C^{-1}$ (with $C$ the sample covariance of all $n$ data) for $C$ provides an estimator with fluctuations of order $O(1)$. Both estimates and fluctuations being of the same order, classification can be achieved at optimal order rate.

Even more surprisingly, the minimal discriminative distance $||\mu|| = O(1)$ is also achieved when using LDA while $C_1$ and $C_2$ are genuinely distinct (where LDA supposes equal). A much unexpected practical consequence is that, at least in scenarios where $||\Delta C||$ does not largely overtake $||\mu||$, it is often more beneficial to use LDA rather than QDA classifiers even though one knows that $C_1 \neq C_2$.2

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2This can be understood from a machine learning intuition as being related to the *overfitting* phenomenon by which, with too many unknown parameters to be estimated by a given algorithm, worse performance can be achieved than if less parameters were estimated in the first place.

### 3.2.2. Supervised learning: LS-SVM and kernel regression

Support vector machines (SVM) are probably the most popular supervised classification method. It consists in constructing a separating hyperplane between data issued from two similarity classes [14]. The asymptotic classification performance of SVM in the random matrix regime has been recently performed but offers so far little conclusive insight [15]. SVM is indeed the result of an optimization problem with no explicit solution, leading to non-trivial asymptotics. As an alternative, least square SVM (LS-SVM) [16] takes an explicit linear regression form by relaxing the SVM “hard” constraints.

In the Gaussian mixture setting, as in many classical application cases, data may not be directly linearly separated by an hyperplane and must then be “re-expressed” through a non-linear map prior to hyperplane separation. This leads to considering kernel methods [17] and in particular to use an affinity metric between $x_i$ and $y_j$ of the type, say, $K_{ij} = f(\frac{1}{\gamma} \|x_i - x_j\|^2)$, for some appropriately set function $f$. With the same notations as previously, LS-SVM (or similarly kernel regression) decides on the class of $x \in \mathbb{R}^p$ upon the test

$$S(x) = \sum_{i=1}^{n} \alpha_i f \left( \frac{1}{p} \|x_i - x\|^2 \right) + b \geq 0$$

where $\alpha = (\alpha_1, \ldots, \alpha_n)^T$ and $b \in \mathbb{R}$ are given, for some hyperparameter $\gamma > 0$ and labels $y_i = (-1)^{(i>n_1)}$, by

$$\alpha = \left( K + \frac{n}{\gamma} I_n \right)^{-1} (y - b 1_n), \quad b = \frac{1}{1_n^T (K + \frac{\gamma}{2} I_n)^{-1} 1_n}. $$

In [18], the asymptotic performance of LS-SVM is evaluated. To this end, the authors rely on a concentration approach allowing for a Taylor-expansion of the kernel matrix $K$ [9], as $n, p \to \infty$. Precisely, under small enough rates for $||\mu||$ and $||\Delta C||$, $\frac{1}{\gamma} \|x_i - x_j\|^2 \sim \frac{1}{\gamma} \text{tr} \left( \frac{1}{n} C_1 + \frac{2}{n} C_2 \right) \equiv \tau$ so that $K_{ij} = f'(\tau) + f'(\tau) \left( \frac{1}{\gamma} \|x_i - x_j\|^2 - \tau \right) + \ldots (i \neq j)$. Developing carefully the expansion for $K$, it is proved that the minimal distance rates follow the rules below:

- for $f'(\tau) \neq 0$, then either (i) $||\mu|| \geq O(1)$ and $\text{tr} \Delta C \geq O(\sqrt{p})$ or (ii) $\text{tr} (||\Delta C||^2) \geq O(p)$;
- for $f'(\tau) = 0$ with $\mu = o(1)$ and $\text{tr} \Delta C = o(\sqrt{p})$, then $\text{tr} (||\Delta C||^2) \geq O(p^2)$.

This result is first surprising in that, for $f'(\tau) \neq 0$, it has a better minimal rate for $||\mu||$ than QDA, even though it is not designed dedicatedly for Gaussian models as in the QDA case. This gain is here due to the fact that no full covariance matrix needs be estimated. As a trade-off though, only the traces of $\Delta C$ and $||\Delta C||^2$ are used for discriminating covariances. Also, as $||\Delta C|| = O(p^{-\frac{1}{2}}) \Rightarrow \text{tr} (||\Delta C||^2) = O(1)$, optimal covariance discriminative rates are not achieved. Besides, the best possible rate $\text{tr} (||\Delta C||^2) = O(p^2)$ is only
achieved under the stringent condition that $\|\mu\| = O(1)$ and $\text{tr} \Delta C = O(\sqrt{p})$ by selecting $f$ such that $f'(\tau) = 0$.

The latter result is nonetheless quite interesting when it comes in practice to differentiating data with mostly different covariances but almost identical means. This is the case for EEG time series as well as multi-antenna wireless communication channels as shown in [19]. A deeper investigation of the scenario $f'(\tau) = 0$ is carried out in [20] with surprising findings when compared to the $f'(\tau) \neq 0$ case. Beyond the case $f'(\tau) = 0$, improved kernel refinements [21] have recently been proposed that ensure both the best rates $\|\mu\| \geq O(1)$ and $\text{tr} (|\Delta C|^2) \geq O(p^{\frac{3}{2}})$. These rely on choosing $f'(\tau) = O(p^{-\frac{1}{2}})$ rather than $f'(\tau) = 0$.

**Remark 2** (Counter-intuitive behavior for large dimensional data). Interestingly, well-performing kernels with $f'(\tau) \simeq 0$ have a counter-intuitive behavior to the classical machine learning lore as a “good” affinity function $f$ is thus not a decreasing function of the distance between data vectors.

### 3.2.3. Supervised learning: Extreme learning machines

Extreme learning machines [22, 23] are practical single hidden-layer networks with a random (not optimized) first connection layer and a ridge-regression as the output layer. That is, the input-to-output relation for input data $x$ reads

$$S(x) = \beta^T \sigma(WX)$$

where $\beta = \frac{1}{n}(\frac{1}{n} \sigma(WX) \sigma(WX)^T + \gamma I_m)^{-1} \sigma(WX)y$, for training data $X = [x_1, \ldots, x_n]$ and $y = [y_1, \ldots, y_n] \in \{-1, 1\}^n$, $W \in \mathbb{R}^{m \times p}$ a random matrix dimensioned by the size-$m$ neuron layer, $\sigma(\cdot)$ a non-linear activation function operating entry-wise, and some regularization parameter $\gamma \geq 0$. For supervised classification, the resulting test consists again here in deciding on the sign of $S(x)$.

In [11], the asymptotic performance of extreme learning machines is studied, as $n, m, p \to \infty$ at the same rate. It is shown that each triplet $(W, \sigma, \frac{m}{n})$ defines a kernel operator $(x,y) \mapsto f(x^T y)$ which makes the extreme learning machine behave similar to LS-SVM. The family of kernels reached in this manner is however somewhat restricted and does not allow for easily setting constraints on the derivatives of $f$ as previously. Yet the minimal discriminative rates for classification are preserved and the same as for LS-SVM.

### 3.2.4. Semi-supervised learning

The LS-SVM method has a natural extension to semi-supervised learning whereby a few samples are already labelled (their classes are known). A classical approach [24, 25], still relying on kernels, consists in computing “scores” $Y_u \in \mathbb{R}^{n_u \times 2}$ for $n_u$ unlabeled data as a function of scores $Y_l \in \mathbb{R}^{n_l \times 2}$ for $n_l$ labelled data, set as

| Datasets   | $||\mu||^2$ | $\frac{1}{\sqrt{p}} \text{tr} (\Delta C^2)$ | Gauss | Quad |
|------------|-------------|------------------------------------------|-------|------|
| MNIST (1, 7) | 613 | 1990 | 98% | 98% |
| MNIST (3, 6) | 441 | 1120 | 98% | 100% |
| MNIST (3, 8) | 212 | 658 | 84% | 95% |
| EEG (sets A, E) | 2.4 | 109 | 69% | 94% |

**Table 1.** Clustering performance for MNIST ($p = 784$) [31] and EEG ($p = 100$) [32] datasets, $n = 1024$, for $f(t) = \exp(-t^2/2)$ (Gauss) and for $f(t) = (t - \tau + \alpha p^{-\frac{1}{2}})^2$ (Quad).

$$[Y_l]_{ij} = \mathbb{1}_{\{x_i \in C_j\}}.$$ An optimization procedure leads to

$$Y_u = (I_{n_u} - K_{uu} D_{uu}^{-1})^{-1} K_{ul} D_{ul}^{-1} Y_l$$

where $D = \text{diag}(K1_n)$ and both $K$ and $D$ were divided into matrix blocks as $K = \begin{pmatrix} K_{uu} & K_{ul}^T \\ K_{ul} & K_{ll} \end{pmatrix}$, $D = \text{diag}(D_{uu}, D_{ll})$, with obvious notations. Here again, in the analysis performed in [26], under the assumption that $n_u, n_l = O(p)$, it is found that the optimal classification rate of a non-trivial number of the $n_u$ unlabeled data is the same as in Section 3.2.2.

### 3.2.5. Unsupervised spectral classification (or clustering)

Spectral clustering consists in classifying $n$ purely unlabelled data from the dominant eigenvectors of matrix $K$ (upon which $k$-means is performed) [27]. In our present scenario, similar conclusions on optimal discriminable rates as in Section 3.2.2 still hold [10, 20]. The result is surprising as almost oracle rates are achieved despite the complete absence of labelled data. This suggests that the kernel spectral clustering method is capable on its own to isolate alike data as well as if in presence of labelled information (at least in terms of rate orders).

Performance figures borrowed from [21] are provided in Table 1, which confirm the large superiority of improved versus standard kernels, particularly for datasets only discriminated through their covariance structure (EEG case). This corroborates theoretical findings on practical datasets, so beyond the Gaussian assumption.

### 4. CONCLUSION

The article showed that, under a large $p$-dimensional binary Gaussian mixture model setting ($N(\mu_0, \sigma_0^2), \alpha \in \{1, 2\}$), elementary machine learning methods achieve non-trivial classification under almost Neyman–Pearson optimal discriminable rates for $||\mu_1 - \mu_2||$ and $||C_2 - C_2||$ as a function of $p$. But the picture is not that simple and many methods need properly set hyperparameters to achieve optimal discriminative power. The article summarizes some of these “proper choices”. More importantly, the results presented here under a Gaussian setting are often shown to be applicable to real world (non Gaussian) data; as such, the article unveils a novel random matrix-improved paradigm in the understanding and optimization of machine learning methods for large dimensional data.
5. REFERENCES


