Convergence Rates of Active Learning for Maximum Likelihood Estimation

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Abstract

An active learner is given a class of models, a large set of unlabeled examples, and the ability to interactively query labels of a subset of these examples; the goal of the learner is to learn a model in the class that fits the data well.

Previous theoretical work has rigorously characterized label complexity of active learning, but most of this work has focused on the PAC or the agnostic PAC model. In this paper, we shift our attention to a more general setting – maximum likelihood estimation. Provided certain conditions hold on the model class, we provide a two-stage active learning algorithm for this problem. The conditions we require are fairly general, and cover the widely popular class of Generalized Linear Models, which in turn, include models for binary and multi-class classification, regression, and conditional random fields.

We provide an upper bound on the label requirement of our algorithm, and a lower bound that matches it up to lower order terms. Our analysis shows that unlike binary classification in the realizable case, just a single extra round of interaction is sufficient to achieve near-optimal performance in maximum likelihood estimation. On the empirical side, the recent work in [11] and [12] (on active linear and logistic regression) shows the promise of this approach.

1 Introduction

In active learning, we are given a sample space \mathcal{X} , a label space \mathcal{Y} , a class of models that map \mathcal{X} to \mathcal{Y} , and a large set U of unlabelled samples. The goal of the learner is to learn a model in the class with small target error while interactively querying the labels of as few of the unlabelled samples as possible.

Most theoretical work on active learning has focussed on the PAC or the agnostic PAC model, where the goal is to learn *binary classifiers* that belong to a particular hypothesis class [2, 13, 9, 6, 3, 4, 22], and there has been only a handful of exceptions [19, 8, 20]. In this paper, we shift our attention to a more general setting – maximum likelihood estimation (MLE), where $\Pr(Y|X)$ is described by a model θ belonging to a model class Θ . We show that when data is generated by a model in this class, we can do active learning provided the model class Θ has the following simple property: the Fisher information matrix for any model $\theta \in \Theta$ at any (x, y) depends only on x and θ . This condition is satisfied in a number of widely applicable model classes, such as Linear Regression and Generalized Linear Models (GLMs), which in turn includes models for Multiclass Classification and Conditional Random Fields. Consequently, we can provide active learning algorithms for maximum likelihood estimation in all these model classes.

The standard solution to active MLE estimation in the statistics literature is to select samples for label query by optimizing a class of summary statistics of the asymptotic covariance matrix of the estimator [5]. The literature, however, does not provide any guidance towards which summary statistic should be used, or any analysis of the solution quality when a finite number of labels or samples are available. There has also been some recent work in the machine learning community [11, 12, 19] on this problem; but these works focus on simple special cases (such as linear regression [19, 11] or logistic regression [12]), and only [19] involves a consistency and finite sample analysis.

In this work, we consider the problem in its full generality, with the goal of minimizing the expected log-likelihood error over the unlabelled data. We provide a two-stage active learning algorithm for this problem. In the first stage, our algorithm queries the labels of a small number of random samples from the data distribution in order to construct a crude estimate θ_1 of the optimal parameter θ^* . In the second stage, we select a set of samples for label query by optimizing a summary statistic of the covariance matrix of the estimator at θ_1 ; however, unlike the experimental design work, our choice of statistic is directly motivated by our goal of minimizing the expected log-likelihood error, which guides us towards the right objective.

We provide a finite sample analysis of our algorithm when some regularity conditions hold and when the negative log likelihood function is convex. Our analysis is still fairly general, and applies to Generalized Linear Models, for example. We match our upper bound with a corresponding lower bound, which shows that the convergence rate of our algorithm is optimal (except for lower order terms); the finite sample convergence rate of any algorithm that uses (perhaps multiple rounds of) sample selection and maximum likelihood estimation is either the same or higher than that of our algorithm. This implies that unlike what is observed in learning binary classifiers, a single round of interaction is sufficient to achieve near-optimal log likelihood error for ML estimation.

1.1 Related Work

Previous theoretical work on active learning has focussed on learning a classifier belonging to a hypothesis class \mathcal{H} in the PAC model. Both the realizable and non-realizable cases have been considered. In the realizable case, a line of work [6, 18] has looked at a generalization of binary search; while their algorithms enjoy low label complexity, this style of algorithms is inconsistent in the presence of noise. The two main styles of algorithms for the non-realizable case are disagreement-based active learning [2, 9, 4], and margin or confidence-based active learning [3, 22]. While active learning in the realizable case has been shown to achieve an exponential improvement in label complexity over passive learning [2, 6, 13], in the agnostic case, the gains are more modest (sometimes a constant factor) [13, 9, 7]. Moreover, lower bounds [14] show that the label requirement of any agnostic active learning algorithm is always at least $\Omega(\nu^2/\epsilon^2)$, where ν is the error of the best hypothesis in the class, and ϵ is the target error. In contrast, our setting is much more general than binary classification, and includes regression, multi-class classification and certain kinds of conditional random fields that are not covered by previous work.

[19] provides an active learning algorithm for linear regression problem under model mismatch. Their algorithm attempts to learn the location of the mismatch by fitting increasingly refined partitions of the domain, and then uses this information to reweight the examples. If the partition is highly refined, then the computational complexity of the resulting algorithm may be exponential in the dimension of the data domain. In contrast, our algorithm applies to a more general setting, and while we do not address model mismatch, our algorithm has polynomial time complexity. [1] provides an active learning algorithm for Generalized Linear Models in an online selective sampling setting; however, unlike ours, their input is a stream of unlabelled examples, and at each step, they need to decide whether the label of the current example should be queried.

Our work is also related to the classical statistical work on optimal experiment design, which

mostly considers maximum likelihood estimation [5]. For uni-variate estimation, they suggest selecting samples to maximize the Fisher information which corresponds to minimizing the variance of the regression coefficient. When θ is multi-variate, the Fisher information is a matrix; in this case, there are multiple notions of *optimal design* which correspond to maximizing different parameters of the Fisher information matrix. For example, D-optimality maximizes the determinant, and A-optimality maximizes the trace of the Fisher information. In contrast with this work, we directly optimize the expected log-likelihood over the unlabelled data which guides us to the appropriate objective function; moreover, we provide consistency and finite sample guarantees.

Finally, on the empirical side, [12] and [11] derive algorithms similar to ours for logistic and linear regression based on projected gradient descent. Notably, these works provide promising empirical evidence for this approach to active learning; however, no consistency guarantees or convergence rates are provided (the rates presented in these works are not stated in terms of the sample size). In contrast, our algorithm applies more generally, and we provide consistency guarantees and convergence rates. Moreover, unlike [12], our logistic regression algorithm uses a single extra round of interaction, and our results illustrate that a single round is sufficient to achieve a convergence rate that is optimal except for lower order terms.

2 The Model

We begin with some notation. We are given a pool $U = \{x_1, \ldots, x_n\}$ of n unlabelled examples drawn from some instance space \mathcal{X} , and the ability to interactively query labels belonging to a label space \mathcal{Y} of m of these examples. In addition, we are given a family of models $\mathcal{M} = \{p(y|x, \theta), \theta \in \Theta\}$ parameterized by $\theta \in \Theta \subseteq \mathbb{R}^d$. We assume that there exists an unknown parameter $\theta^* \in \Theta$ such that querying the label of an $x_i \in U$ generates a y_i drawn from the distribution $p(y|x_i, \theta^*)$. We also abuse notation and use U to denote the uniform distribution over the examples in U.

We consider the *fixed-design* (or *transductive*) setting, where our goal is to minimize the error on the fixed set of points U. For any $x \in \mathcal{X}, y \in \mathcal{Y}$ and $\theta \in \Theta$, we define the negative log-likelihood function $L(y|x, \theta)$ as:

$$L(y|x,\theta) = -\log p(y|x,\theta)$$

Our goal is to find a $\hat{\theta}$ to minimize $L_U(\hat{\theta})$, where

$$L_U(\theta) = \mathbb{E}_{X \sim U, Y \sim p(Y|X, \theta^*)} [L(Y|X, \theta)]$$

by interactively querying labels for a subset of U of size m, where we allow label queries with replacement i.e., the label of an example may be queried multiple times.

An additional quantity of interest to us is the Fisher information matrix, or the Hessian of the negative log-likelihood $L(y|x,\theta)$ function, which determines the convergence rate. For our active learning procedure to work correctly, we require the following condition.

Condition 1. For any $x \in \mathcal{X}$, $y \in \mathcal{Y}$, $\theta \in \Theta$, the Fisher information $\frac{\partial^2 L(y|x,\theta)}{\partial \theta^2}$ is a function of only x and θ (and does not depend on y.)

Condition 1 is satisfied by a number of models of practical interest; examples include linear regression and generalized linear models. Section 5.1 provides a brief derivation of Condition 1 for generalized linear models.

For any x, y and θ , we use $I(x, \theta)$ to denote the Hessian $\frac{\partial^2 L(y|x, \theta)}{\partial \theta^2}$; observe that by Assumption 1, this is just a function of x and θ . Let Γ be any distribution over the unlabelled samples in U; for any $\theta \in \Theta$, we use:

$$I_{\Gamma}(\theta) = \mathbb{E}_{X \sim \Gamma}[I(X, \theta)]$$

Algorithm 1 ActiveSetSelect

Input: Samples x_i , for $i = 1, \dots, n$

1: Draw m_1 samples u.a.r from U, and query their labels to get S_1 .

2: Use S_1 to solve the MLE problem:

$$\theta_1 = \operatorname{argmin}_{\theta \in \Theta} \sum_{(x_i, y_i) \in S_1} L(y_i | x_i, \theta)$$

3: Solve the following SDP (refer Lemma 3):

$$a^* = \operatorname{argmin}_a \operatorname{Tr} \left(S^{-1} I_U(\theta_1) \right) \quad \text{s.t.} \begin{cases} S = \sum_i a_i I(x_i, \theta_1) \\ 0 \le a_i \le 1 \\ \sum_i a_i = m_2 \end{cases}$$

4: Draw m_2 examples using probability $\overline{\Gamma} = \alpha \Gamma_1 + (1 - \alpha)U$ where the distribution $\Gamma_1 = \frac{a_i^2}{m_2}$ and $\alpha = 1 - m_2^{-1/6}$. Query their labels to get S_2 .

5: Use S_2 to solve the MLE problem:

$$\theta_2 = \operatorname{argmin}_{\theta \in \Theta} \sum_{(x_i, y_i) \in S_2} L(y_i | x_i, \theta)$$

Output: θ_2

3 Algorithm

The main idea behind our algorithm is to sample x_i from a well-designed distribution Γ over U, query the labels of these samples and perform ML estimation over them. To ensure good performance, Γ should be chosen carefully, and our choice of Γ is motivated by Lemma 1. Suppose the labels y_i are generated according to: $y_i \sim p(y|x_i, \theta^*)$. Lemma 1 states that the expected log-likelihood error of the ML estimate with respect to m samples from Γ in this case is essentially $\operatorname{Tr} \left(I_{\Gamma}(\theta^*)^{-1}I_U(\theta^*)\right)/m$.

This suggests selecting Γ as the distribution Γ^* that minimizes $\operatorname{Tr} \left(I_{\Gamma^*}(\theta^*)^{-1} I_U(\theta^*) \right)$. Unfortunately, we cannot do this as θ^* is unknown. We resolve this problem through a two stage algorithm; in the first stage, we use a small number m_1 of samples to construct a coarse estimate θ_1 of θ^* (Steps 1-2). In the second stage, we calculate a distribution Γ_1 which minimizes $\operatorname{Tr} \left(I_{\Gamma_1}(\theta_1)^{-1} I_U(\theta_1) \right)$ and draw samples from (a slight modification of) this distribution for a finer estimation of θ^* (Steps 3-5). The distribution Γ_1 is modified slightly to $\overline{\Gamma}$ (in Step 4) to ensure that $I_{\overline{\Gamma}}(\theta^*)$ is well conditioned with respect to $I_U(\theta^*)$.

The algorithm is formally presented in Algorithm 1.

Finally, note that Steps 1-2 are necessary because I_U and I_{Γ} are functions of θ . In certain special cases such as linear regression, I_U and I_{Γ} are independent of θ . In those cases, Steps 1-2 are unnecessary, and we may skip directly to Step 3.

4 Performance Guarantees

The following regularity conditions are essentially a quantified version of the standard Local Asymptotic Normality (LAN) conditions for studying maximum likelihood estimation (see [16, 21]).

Assumption 1. (Regularity conditions for LAN)

- 1. Smoothness: The first three derivatives of $L(y|x,\theta)$ exist in all interior points of $\Theta \subseteq \mathbb{R}^d$.
- 2. Compactness: Θ is compact and θ^* is an interior point of Θ .
- 3. Strong Convexity: $I_U(\theta^*) = \frac{1}{n} \sum_{i=1}^n I(x_i, \theta^*)$ is positive definite with smallest singular value $\sigma_{min} > 0$.
- 4. Lipschitz continuity: There exists a neighborhood B of θ^* and a constant L_3 such that for all $x \in U$, $I(x, \theta)$ is L_3 -Lipschitz in this neighborhood.

$$\left\| I_{U}(\theta^{*})^{-1/2} \left(I(x,\theta) - I(x,\theta') \right) I_{U}(\theta^{*})^{-1/2} \right\|_{2} \le L_{3} \left\| \theta - \theta' \right\|_{I_{U}(\theta^{*})},$$

for every $\theta, \theta' \in B$.

5. Concentration at θ^* : For any $x \in U$ and y, we have (with probability one),

$$\|\nabla L(y|x,\theta^*)\|_{I_U(\theta^*)^{-1}} \le L_1, \text{ and } \|I_U(\theta^*)^{-1/2}I(x,\theta^*)I_U(\theta^*)^{-1/2}\|_2 \le L_2.$$

6. Boundedness: $\max_{(x,y)} \sup_{\theta \in \Theta} |L(x,y|\theta)| \le R$.

In addition to the above, we need one extra condition which is essentially a pointwise self concordance. This condition is satisfied by a vast class of models, including the generalized linear models.

Assumption 2. Point-wise self concordance:

$$-L_4 \left\| \theta - \theta^* \right\|_2 I(x,\theta^*) \leq I(x,\theta) - I(x,\theta^*) \leq L_4 \left\| \theta - \theta^* \right\|_2 I(x,\theta^*).$$

Definition 1. [Optimal Sampling Distribution Γ^*] We define the optimal sampling distribution Γ^* over the points in U as the distribution $\Gamma^* = (\gamma_1^*, \ldots, \gamma_n^*)$ for which $\gamma_i^* \ge 0$, $\sum_i \gamma_i^* = 1$, and $Tr(I_{\Gamma^*}(\theta^*)^{-1}I_U(\theta^*))$ is as small as possible.

Definition 1 is motivated by Lemma 1, which indicates that under some mild regularity conditions, a ML estimate calculated on samples drawn from Γ^* will provide the best convergence rates (including the right constant factor) for the expected log-likelihood error.

We now present the main result of our paper. The proof of the following theorem and all the supporting lemmas will be presented in Appendix A.

Theorem 1. Suppose the regularity conditions in Assumptions 1 and 2 hold. Let $\beta \geq 10$, and the number of samples used in step (1) be $m_1 > \mathcal{O}\left(\max\left(L_2\log^2 d, L_1^2\left(L_3^2 + \frac{1}{\sigma_{\min}}\right)\log^2 d, \frac{diameter(\Theta)}{Tr(I_U(\theta^*)^{-1})}, \frac{\beta^2 L_4^2}{\delta} Tr(I_U(\theta^*)^{-1})\right)\right)$. Then with probability $\geq 1 - \delta$, the expected log likelihood error of the estimate θ_2 of Algorithm 1 is bounded as:

$$\mathbb{E}\left[L_{U}(\theta_{2})\right] - L_{U}(\theta^{*}) \leq \left(1 + \frac{2}{\beta - 1}\right)^{4} \left(1 + \widetilde{\epsilon}_{m_{2}}\right) Tr\left(I_{\Gamma^{*}}(\theta^{*})^{-1}I_{U}(\theta^{*})\right) \frac{1}{m_{2}} + \frac{R}{m_{2}^{2}}, \tag{1}$$

where Γ^* is the optimal sampling distribution in Definition 1 and $\tilde{\epsilon}_{m_2} = \mathcal{O}\left(\left(L_1L_3 + \sqrt{L_2}\right)\frac{\sqrt{\log dm_2}}{m_2^{1/6}}\right)$. Moreover, for any sampling distribution Γ satisfying $I_{\Gamma}(\theta^*) \succeq cI_U(\theta^*)$ and label constraint of m_2 , we have the following lower bound on the expected log likelihood error for ML estimate:

$$\mathbb{E}\left[L_U(\widehat{\theta}_{\Gamma})\right] - L_U(\theta^*) \ge (1 - \epsilon_{m_2}) \operatorname{Tr}\left(I_{\Gamma}(\theta^*)^{-1} I_U(\theta^*)\right) \frac{1}{m_2} - \frac{L_1^2}{cm_2^2},\tag{2}$$

where $\epsilon_{m_2} \stackrel{def}{=} \frac{\widetilde{\epsilon}_{m_2}}{c^2 m_2^{1/3}}.$

Remark 1. (Restricting to Maximum Likelihood Estimation) Our restriction to maximum likelihood estimators is minor, as this is close to minimax optimal (see [15]). Minor improvements with certain kinds of estimators, such as the James-Stein estimator, are possible.

4.1 Discussions

Several remarks about Theorem 1 are in order.

The high probability bound in Theorem 1 is with respect to the samples drawn in S_1 ; provided these samples are representative (which happens with probability $\geq 1 - \delta$), the output θ_2 of Algorithm 1 will satisfy (1). Additionally, Theorem 1 assumes that the labels are sampled with replacement; in other words, we can query the label of a point x_i multiple times. Removing this assumption is an avenue for future work.

Second, the highest order term in both (1) and (2) is $\operatorname{Tr}\left(I_{\Gamma^*}(\theta^*)^{-1}I_U(\theta^*)\right)/m$. The terms involving ϵ_{m_2} and $\tilde{\epsilon}_{m_2}$ are lower order as both ϵ_{m_2} and $\tilde{\epsilon}_{m_2}$ are o(1). Moreover, if $\beta = \omega(1)$, then the term involving β in (1) is of a lower order as well. Observe that β also measures the tradeoff between m_1 and m_2 , and as long as $\beta = o(\sqrt{m_2})$, m_1 is also of a lower order than m_2 . Thus, provided β is $\omega(1)$ and $o(\sqrt{m_2})$, the convergence rate of our algorithm is optimal except for lower order terms.

Finally, the lower bound (2) applies to distributions Γ for which $I_{\Gamma}(\theta^*) \geq cI_U(\theta^*)$, where c occurs in the lower order terms of the bound. This constraint is not very restrictive, and does not affect the asymptotic rate. Observe that $I_U(\theta^*)$ is full rank. If $I_{\Gamma}(\theta^*)$ is not full rank, then the expected log likelihood error of the ML estimate with respect to Γ will not be consistent, and thus such a Γ will never achieve the optimal rate. If $I_{\Gamma}(\theta^*)$ is full rank, then there always exists a c for which $I_{\Gamma}(\theta^*) \geq cI_U(\theta^*)$. Thus (2) essentially states that for distributions Γ where $I_{\Gamma}(\theta^*)$ is close to being rank-deficient, the asymptotic convergence rate of $O(\text{Tr}(I_{\Gamma}(\theta^*)^{-1}I_U(\theta^*))/m_2)$ is achieved at larger values of m_2 .

4.2 Proof Outline

Our main result relies on the following three steps.

4.2.1 Bounding the Log-likelihood Error

First, we characterize the log likelihood error (wrt U) of the empirical risk minimizer (ERM) estimate obtained using a sampling distribution Γ . Concretely, let Γ be a distribution on U. Let $\hat{\theta}_{\Gamma}$ be the ERM estimate using the distribution Γ :

$$\widehat{\theta}_{\Gamma} = \operatorname{argmin}_{\theta \in \Theta} \frac{1}{m_2} \sum_{i=1}^{m_2} L(Y_i | X_i, \theta),$$
(3)

where $X_i \sim \Gamma$ and $Y_i \sim p(y|X_i, \theta^*)$. The core of our analysis is Lemma 1, which shows a precise estimate of the log likelihood error $\mathbb{E}\left[L_U\left(\widehat{\theta}_{\Gamma}\right) - L_U\left(\theta^*\right)\right]$.

Lemma 1. Suppose L satisfies the regularity conditions in Assumptions 1 and 2. Let Γ be a distribution on U and $\hat{\theta}_{\Gamma}$ be the ERM estimate (3) using m_2 labeled examples. Suppose further that $I_{\Gamma}(\theta^*) \succeq cI_U(\theta^*)$ for some constant c < 1. Then, for any $p \ge 2$ and m_2 large enough (depending on p), we have:

$$(1-\epsilon_{m_2})\frac{\tau^2}{m_2} - \frac{L_1^2}{cm_2^{p/2}} \le \mathbb{E}\left[L_U\left(\widehat{\theta}_{\Gamma}\right) - L_U\left(\theta^*\right)\right] \le (1+\epsilon_{m_2})\frac{\tau^2}{m_2} + \frac{R}{m_2^p}$$

where $\epsilon_{m_2} = \mathcal{O}\left(\frac{1}{c^2}\left(L_1L_3 + \sqrt{L_2}\right)\sqrt{\frac{p\log dm_2}{m_2}}\right)$ and $\tau^2 \stackrel{def}{=} Tr\left(I_{\Gamma}(\theta^*)^{-1}I_U(\theta^*)\right)$.

4.2.2 Approximating θ^*

Lemma 1 motivates sampling from the optimal sampling distribution Γ^* that minimizes $\operatorname{Tr}\left(I_{\Gamma^*}(\theta^*)^{-1}I_U(\theta^*)\right)$. However, this quantity depends on θ^* , which we do not know. To resolve this issue, our algorithm first queries the labels of a small fraction of points (m_1) and solves a ML estimation problem to obtain a coarse estimate θ_1 of θ^* .

How close should θ_1 be to θ^* ? Our analysis indicates that it is sufficient for θ_1 to be close enough that for any x, $I(x, \theta_1)$ is a *constant factor spectral approximation* to $I(x, \theta^*)$; the number of samples needed to achieve this is analyzed in Lemma 2.

Lemma 2. Suppose L satisfies the regularity conditions in Assumptions 1 and 2. If the number of samples used in the first step

$$m_1 > \mathcal{O}\left(\max\left(L_2\log^2 d, L_1^2\left(L_3^2 + \frac{1}{\sigma_{min}}\right)\log^2 d, \frac{diameter(\Theta)}{Tr\left(I_U(\theta^*)^{-1}\right)}, \frac{\beta^2 L_4^2}{\delta} Tr\left(I_U(\theta^*)^{-1}\right)\right)\right),$$

then, we have:

$$-\frac{1}{\beta}I(x,\theta^*) \preceq I(x,\theta_1) - I(x,\theta^*) \preceq \frac{1}{\beta}I(x,\theta^*) \ \forall x \in X$$

with probability greater than $1 - \delta$.

4.2.3 Computing Γ_1

Third, we are left with the task of obtaining a distribution Γ_1 that minimizes the log likelihood error. We now pose this optimization problem as an SDP.

From Lemmas 1 and 2, it is clear that we should aim to obtain a sampling distribution $\Gamma = (\frac{a_i}{m_2} : i \in [n])$ minimizing $\operatorname{Tr} \left(I_{\Gamma}(\theta_1)^{-1} I_U(\theta_1) \right)$. Let $I_U(\theta_1) = \sum_j \sigma_j v_j v_j^{\top}$ be the singular value decomposition (svd) of $I_U(\theta_1)$. Since $\operatorname{Tr} \left(I_{\Gamma}(\theta_1)^{-1} I_U(\theta_1) \right) = \sum_{j=1}^d \sigma_j v_j^{\top} I_{\Gamma}(\theta_1)^{-1} v_j$, this is equivalent to solving:

$$\min_{a,c} \sum_{j=1}^{d} \sigma_j c_j \quad \text{s.t.} \begin{cases} S = \sum_i a_i I(x_i, \theta_1) \\ v_j^{\top} S^{-1} v_j \leq c_j \\ a_i \in [0, 1] \\ \sum_i a_i = m_2. \end{cases}$$
(4)

_ d

Among the above constraints, the constraint $v_j^{\top}S^{-1}v_j \leq c_j$ seems problematic. However, Schur complement formula tells us that: $\begin{bmatrix} c_j & v_j^{\top} \\ v_j & S \end{bmatrix} \geq 0 \iff S \geq 0$ and $v_j^{\top}S^{-1}v_j \leq c_j$. In our case, we know that $S \succeq 0$, since it is a sum of positive semi definite matrices. The above argument proves the following lemma.

Lemma 3. The following two optimization programs are equivalent:

where $I_U(\theta_1) = \sum_j \sigma_j v_j v_j^\top$ denotes the svd of $I_U(\theta_1)$.

5 Illustrative Examples

We next present some examples that illustrate Theorem 1. We begin by showing that Condition 1 is satisfied by the popular class of Generalized Linear Models.

5.1 Derivations for Generalized Linear Models

A generalized linear model is specified by three parameters – a linear model, a sufficient statistic, and a member of the exponential family. Let η be a linear model: $\eta = \theta^{\top} X$. Then, in a Generalized Linear Model (GLM), Y is drawn from an exponential family distribution with parameter η . Specifically, $p(Y = y|\eta) = e^{\eta^{\top} t(y) - A(\eta)}$, where $t(\cdot)$ is the sufficient statistic and $A(\cdot)$ is the log-partition function. From properties of the exponential family, the log-likelihood is written as $\log p(y|\eta) = \eta^{\top} t(y) - A(\eta)$. If we take $\eta = \theta^{\top} x$, and take the derivative with respect to θ , we have: $\frac{\partial \log p(y|\theta,x)}{\partial \theta} = xt(y) - xA'(\theta^{\top} x)$. Taking derivatives again gives us $\frac{\partial^2 \log p(y|\theta,x)}{\partial \theta^2} = -xx^{\top}A''(\theta^{\top} x)$, which is independent of y.

5.2 Specific Examples

We next present three illustrative examples of problems that our algorithm may be applied to.

Linear Regression. Our first example is linear regression. In this case, $x \in \mathbb{R}^d$ and $Y \in \mathbb{R}$ are generated according to the distribution: $Y = \theta_*^\top X + \eta$, where η is a noise variable drawn from $\mathcal{N}(0, 1)$. In this case, the negative loglikelihood function is: $L(y|x, \theta) = (y - \theta^\top x)^2$, and the corresponding Fisher information matrix $I(x, \theta)$ is given as: $I(x, \theta) = xx^\top$. Observe that in this (very special) case, the Fisher information matrix does not depend on θ ; as a result we can eliminate the first two steps of the algorithm, and proceed directly to step 3. If $\Sigma = \frac{1}{n} \sum_i x_i x_i^\top$ is the covariance matrix of U, then Theorem 1 tells us that we need to query labels from a distribution Γ^* with covariance matrix Λ such that $\operatorname{Tr}(\Lambda^{-1}\Sigma)$ is minimized.

We illustrate the advantages of active learning through a simple example. Suppose U is the unlabelled distribution:

$$x_{i} = \begin{cases} e_{1} & \text{w.p. } 1 - \frac{d-1}{d^{2}}, \\ e_{j} & \text{w.p. } \frac{1}{d^{2}} \text{ for } j \in \{2, \cdots, d\}, \end{cases}$$

where e_j is the standard unit vector in the j^{th} direction. The covariance matrix Σ of U is a diagonal matrix with $\Sigma_{11} = 1 - \frac{d-1}{d^2}$ and $\Sigma_{jj} = \frac{1}{d^2}$ for $j \ge 2$. For passive learning over U, we query labels of examples drawn from U which gives us a convergence rate of $\frac{\text{Tr}(\Sigma^{-1}\Sigma)}{m} = \frac{d}{m}$. On the other hand, active learning chooses to sample examples from the distribution Γ^* such that

$$x_i = \begin{cases} e_1 & \text{w.p.} \sim 1 - \frac{d-1}{2d}, \\ e_j & \text{w.p.} \sim \frac{1}{2d} \text{ for } j \in \{2, \cdots, d\}, \end{cases}$$

where ~ indicates that the probabilities hold upto $\mathcal{O}\left(\frac{1}{d^2}\right)$. This has a diagonal covariance matrix Λ such that $\Lambda_{11} \sim 1 - \frac{d-1}{2d}$ and $\Lambda_{jj} \sim \frac{1}{2d}$ for $j \geq 2$, and convergence rate of $\frac{\operatorname{Tr}(\Lambda^{-1}\Sigma)}{m} \sim \frac{1}{m} \left(\frac{2d}{d+1} \cdot \left(1 - \frac{d-1}{d^2}\right) + (d-1) \cdot 2d \cdot \frac{1}{d^2}\right) \leq \frac{4}{m}$, which does not grow with d!

Logistic Regression. Our second example is logistic regression for binary classification. In this case, $x \in \mathbb{R}^d$, $Y \in \{-1, 1\}$ and the negative log-likelihood function is: $L(y|x, \theta) = \log(1 + e^{-y\theta^\top x})$, and the corresponding Fisher information $I(x, \theta)$ is given as: $I(x, \theta) = \frac{e^{\theta^\top x}}{(1 + e^{\theta^\top x})^2} \cdot xx^\top$.

For illustration, suppose $\|\theta^*\|_2$ and $\|x\|_2$ are bounded by a constant and the covariance matrix Σ is sandwiched between two multiples of identity in the PSD ordering i.e., $\frac{c}{d}I \preceq \Sigma \preceq \frac{C}{d}I$ for

some constants c and C. Then the regularity assumptions 1 and 2 are satisfied for constant values of L_1, L_2, L_3 and L_4 . In this case, Theorem 1 states that choosing m_1 to be $\omega \left(\operatorname{Tr} \left(I_U(\theta^*)^{-1} \right) \right) = \omega (d)$ gives us the optimal convergence rate of $(1 + o(1)) \frac{\operatorname{Tr} \left(I_{\Gamma^*}(\theta^*)^{-1} I_U(\theta^*) \right)}{m_2}$.

Multinomial Logistic Regression. Our third example is multinomial logistic regression for multi-class classification. In this case, $Y \in 1, ..., K$, $x \in \mathbb{R}^d$, and the parameter matrix $\theta \in \mathbb{R}^{(K-1)\times d}$. The negative log-likelihood function is written as: $L(y|x,\theta) = -\theta_y^\top x + \log(1+\sum_{k=1}^{K-1} e^{\theta_k^\top x})$, if $y \neq K$, and $L(y = k|x,\theta) = \log(1+\sum_{k=1}^{K-1} e^{\theta_k^\top x})$ otherwise. The corresponding Fisher information matrix is a $(K-1)d \times (K-1)d$ matrix, which is obtained as follows. Let F be the $(K-1) \times (K-1)$ matrix with:

$$F_{ii} = \frac{e^{\theta_i^{\top} x} (1 + \sum_{k \neq i} e^{\theta_k^{\top} x})}{(1 + \sum_k e^{\theta_k^{\top} x})^2}, \quad F_{ij} = -\frac{e^{\theta_i^{\top} x} + \theta_j^{\top} x}{(1 + \sum_k e^{\theta_k^{\top} x})^2}$$

Then, $I(x, \theta) = F \otimes xx^{\top}$.

Similar to the example in the logistic regression case, suppose $\|\theta_y^*\|_2$ and $\|x\|_2$ are bounded by a constant and the covariance matrix Σ satisfies $\frac{c}{d}I \leq \Sigma \leq \frac{C}{d}I$ for some constants c and C. Since $F^* = \text{diag}(p_i^*) - p^* p^{*\top}$, where $p_i^* = P(y = i|x,\theta^*)$, the boundedness of $\|\theta_y^*\|_2$ and $\|x\|_2$ implies that $\tilde{c}I \leq F^* \leq \tilde{C}I$ for some constants \tilde{c} and \tilde{C} (depending on K). This means that $\frac{c\tilde{c}}{d}I \leq I(x,\theta^*) \leq \frac{C\tilde{C}}{d}I$ and so the regularity assumptions 1 and 2 are satisfied with L_1, L_2, L_3 and L_4 being constants. Theorem 1 again tells us that using $\omega(d)$ samples in the first step gives us the optimal convergence rate of maximum likelihood error.

6 Conclusion

In this paper, we provide an active learning algorithm for maximum likelihood estimation which provably achieves the optimal convergence rate (upto lower order terms) and uses only two rounds of interaction. Our algorithm applies in a very general setting, which includes Generalized Linear Models.

There are several avenues of future work. Our algorithm involves solving an SDP which is computationally expensive; an open question is whether there is a more efficient, perhaps greedy, algorithm that achieves the same rate. A second open question is whether it is possible to remove the *with replacement* sampling assumption. A final question is what happens if $I_U(\theta^*)$ has a high condition number. In this case, our algorithm will require a large number of samples in the first stage; an open question is whether we can use a more sophisticated procedure in the first stage to reduce the label requirement.

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A Proofs

In order to prove Lemma 1, we use the following result which is a modification of [10]. In particular, the following lemma is a generalization of Theorem 5.1 from [10], and its proof (omitted here) follows from generalizing the proof of that theorem.

Lemma 4. Suppose $\psi_1, \dots, \psi_n : \mathbb{R}^d \to \mathbb{R}$ are random functions drawn iid from a distribution. Let $P = \mathbb{E}[\psi_i]$ and $Q : \mathbb{R}^d \to \mathbb{R}$ be another function. Let

$$\widehat{\theta} = \operatorname{argmin}_{\theta \in \mathcal{S}} \sum_{i} \psi_{i}(\theta), \text{ and } \theta^{*} = \operatorname{argmin}_{\theta \in \mathcal{S}} P(\theta).$$

Assume:

- 1. (Convexity of ψ): Assume that ψ is convex (with probability one),
- 2. (Smoothness of ψ): Assume that ψ is smooth in the following sense: the first, second and third derivatives exist at all interior points of S (with probability one),
- 3. (Regularity conditions): Suppose
 - (a) S is compact,
 - (b) θ^* is an interior point of S,
 - (c) $\nabla^2 P(\theta^*)$ is positive definite (and hence invertible),
 - (d) $\nabla Q(\theta^*) = 0$,
 - (e) There exists a neighborhood B of θ^* and a constant \widetilde{L}_3 such that (with probability one), $\nabla^2 \psi(\theta)$ and $\nabla^2 Q(\theta)$ are \widetilde{L}_3 Lipschitz, namely

$$\left\| \left(\nabla^2 P(\theta^*) \right)^{-1/2} \left(\nabla^2 \psi(\theta) - \nabla^2 \psi(\theta') \right) \left(\nabla^2 P(\theta^*) \right)^{-1/2} \right\|_2 \leq \widetilde{L_3} \left\| \theta - \theta' \right\|_{\nabla^2 P(\theta^*)}, \text{ and} \\ \left\| \left(\nabla^2 Q(\theta^*) \right)^{-1/2} \left(\nabla^2 Q(\theta) - \nabla^2 Q(\theta') \right) \left(\nabla^2 Q(\theta^*) \right)^{-1/2} \right\|_2 \leq \widetilde{L_3} \left\| \theta - \theta' \right\|_{\nabla^2 P(\theta^*)}, \text{ and} \\ \text{if } \theta \in B.$$

- for $\theta, \theta' \in B$,
- 4. (Concentration at θ^*) Suppose $\|\nabla \psi(\theta^*)\|_{\nabla^2 P(\theta^*)^{-1}} \leq \widetilde{L_1}$ and

$$\left\| \left(\nabla^2 P(\theta^*) \right)^{-1/2} \nabla^2 \psi(\theta^*) \left(\nabla^2 P(\theta^*) \right)^{-1/2} \right\|_2 \le \widetilde{L_2}$$

hold with probability one.

Choose $p \geq 2$ and define

$$\epsilon_n \stackrel{def}{=} \widetilde{c}(\widetilde{L_1}\widetilde{L_3} + \sqrt{\widetilde{L_2}})\sqrt{\frac{p\log dn}{n}},$$

where \tilde{c} is an appropriately chosen constant. Let \tilde{c}' be another appropriately chosen constant. If n is large enough so that $\sqrt{\frac{p\log dn}{n}} \leq \tilde{c}' \min\left\{\frac{1}{\sqrt{\tilde{L_2}}}, \frac{1}{\tilde{L_1}\tilde{L_3}}, \frac{diameter(B)}{\tilde{L_1}}\right\}$, then:

$$(1-\epsilon_n)\frac{\tau^2}{n} - \frac{\widetilde{L_1}^2}{n^{p/2}} \le \mathbb{E}\left[Q(\widehat{\theta}) - Q(\theta^*)\right] \le (1+\epsilon_n)\frac{\tau^2}{n} + \frac{\max_{\theta \in \mathcal{S}} Q(\theta) - Q(\theta^*)}{n^p},$$

where

$$\tau^2 \stackrel{def}{=} \frac{1}{n^2} Tr\left(\left(\sum_{i,j} \mathbb{E}\left[\nabla \psi_i(\theta^*) \nabla \psi_j(\theta^*)^\top\right]\right) P(\theta^*)^{-1} Q(\theta^*) P(\theta^*)^{-1}\right).$$

The following lemma is a fundamental result relating the variance of the gradient of the log likelihood to Fisher information matrix for a large class of probability distributions [17].

Lemma 5. Suppose L satisfies the regularity conditions in Assumptions 1 and 2. Then, for any example x, we have:

$$\mathbb{E}_{p(y|x,\theta^*)}\left[\nabla L(Y|x,\theta^*)\nabla L(Y|x,\theta^*)^\top\right] = \nabla^2 I_x(\theta^*).$$

We now prove Lemma 1.

(Proof of Lemma 1). We first define

$$\psi_i(\theta) = L\left(Y|X,\theta\right),$$

where $X \sim \Gamma$ and $Y \sim p(Y|X, \theta^*)$ for $i = 1, \dots, m_2$ and $Q(\theta) \stackrel{\text{def}}{=} L_U(\theta)$. Using the notation of Lemma 4, this means that

$$\nabla^2 P(\theta^*) = I_{\Gamma}(\theta^*) \text{ and } \nabla^2 Q(\theta^*) = I_U(\theta^*).$$

Using the regularity conditions from Section 4 and the hypothesis that $I_{\Gamma}(\theta^*) \succeq cI_U(\theta^*)$, we see that this satisfies the hypothesis of Lemma 4 with constants

$$(\widetilde{L_1}, \widetilde{L_2}, \widetilde{L_3}) = (L_1/\sqrt{c}, L_2/c, L_3/c^{3/2})$$

We now apply Lemma 4 to conclude that for large enough m_2 , we have:

$$(1 - \epsilon_{m_2})\tau^2 / m_2 - \frac{L_1^2}{cm_2^{p/2}} \le \mathbb{E}\left[L_U\left(\widehat{\theta}\right) - L_U\left(\theta^*\right)\right] \le (1 + \epsilon_{m_2})\tau^2 / m_2 + \frac{R}{m_2^p},$$

where

$$\epsilon_{m_2} = \mathcal{O}\left(\left(\widetilde{L_1}\widetilde{L_3} + \sqrt{\widetilde{L_2}}\right)\sqrt{\frac{p\log dm_2}{m_2}}\right) = \mathcal{O}\left(\frac{1}{c^2}\left(L_1L_3 + \sqrt{L_2}\right)\sqrt{\frac{p\log dm_2}{m_2}}\right) \text{ and}$$
$$\tau^2 \stackrel{\text{def}}{=} \operatorname{Tr}\left(\mathbb{E}\left[\nabla\widehat{P}(\theta^*)\nabla\widehat{P}(\theta^*)^{\top}\right]I_{\Gamma}(\theta^*)^{-1}I_U(\theta^*)I_{\Gamma}(\theta^*)^{-1}\right) = \operatorname{Tr}\left(I_{\Gamma}(\theta^*)^{-1}I_U(\theta^*)\right),$$

using Lemma 5 in the last step.

We now prove Lemma 2.

(Proof of Lemma 2). Define

$$\psi_i(\theta) \stackrel{\text{def}}{=} L(Y|X,\theta)$$

where $X \sim U$ and $Y \sim p(Y|X, \theta^*)$ for $i = 1, \dots, m_1$ and $Q(\theta) \stackrel{\text{def}}{=} \|\theta - \theta^*\|_2^2$. Using the regularity conditions from Section 4, we see that this satisfies the hypothesis of Lemma 4 with constants

$$(\widetilde{L_1}, \widetilde{L_2}, \widetilde{L_3}) = (L_1, L_2, \max\left(L_3, \frac{1}{\sqrt{\sigma_{\min}}}\right)))$$

We now apply Lemma 4 to conclude that

$$\mathbb{E}\left[\left\|\theta_1 - \theta^*\right\|_2^2\right] \le (1 + \epsilon_{m_1})\tau^2/m_1 + \frac{\operatorname{diameter}(\Theta)}{m_1^2},$$

where
$$\epsilon_{m_1} = \mathcal{O}\left((L_1 \max\left(L_3, \frac{1}{\sqrt{\sigma_{\min}}}\right) + \sqrt{L_2})\sqrt{\frac{\log dm_1}{m_1}} \right)$$
, and
 $\tau^2 \stackrel{\text{def}}{=} \operatorname{Tr}\left(\mathbb{E}\left[\nabla \widehat{L}_U(\theta^*) \nabla \widehat{L}_U(\theta^*)^\top \right] I_U(\theta^*)^{-2} \right) = \operatorname{Tr}\left(I_U(\theta^*)^{-1} \right),$

using Lemma 5 in the last step. By the choice of m_1 , we have that

$$\mathbb{E}\left[\left\|\theta_1 - \theta^*\right\|_2^2\right] \le 2\tau^2/m_1.$$

Markov's inequality then tells us that with probability at least $1 - \delta$, we have:

$$\|\theta_1 - \theta^*\|_2^2 \le \frac{2\tau^2}{\delta m_1} \le \frac{1}{\beta^2 L_4^2}$$

Using Assumption 2 on point-wise self concordancy of $I(x, \theta)$ now finishes the proof.

(Proof of Theorem 1). The proof is a careful combination of Lemmas 1, 2 and 3.

Lower Bound: For any Γ that satisfies $I_{\Gamma}(\theta^*) \succeq cI_U(\theta^*)$, we can apply Lemma 1 to write:

$$\mathbb{E}\left[L_U\left(\widehat{\theta}_{\Gamma}\right) - L_U\left(\theta^*\right)\right] \ge (1 - \epsilon_{m_2}) \frac{\operatorname{Tr}\left(I_{\Gamma}(\theta^*)^{-1}I_U(\theta^*)\right)}{m_2} - \frac{L_1^2}{cm_2^2}$$

The lower bound follows.

Upper Bound: We begin by showing that if Assumptions 1 and 2 are satisfied, then, from Lemma 2, we have that with probability $\geq 1 - \delta$, it holds that:

$$\frac{\beta - 1}{\beta} I(x, \theta^*) \preceq I(x, \theta_1) \preceq \frac{\beta + 1}{\beta} I(x, \theta^*) \ \forall \ x \in U$$

with probability $\geq 1 - \delta$. This means that the following hold for distributions Γ_1 , Γ^* and U:

$$\frac{\beta - 1}{\beta} I_{\Gamma_1}(\theta^*) \preceq I_{\Gamma_1}(\theta_1) \preceq \frac{\beta + 1}{\beta} I_{\Gamma_1}(\theta^*), \tag{5}$$

$$\frac{\beta - 1}{\beta} I_{\Gamma^*}(\theta^*) \preceq I_{\Gamma^*}(\theta_1) \preceq \frac{\beta + 1}{\beta} I_{\Gamma^*}(\theta^*), \text{ and}$$
(6)

$$\frac{\beta - 1}{\beta} I_U(\theta^*) \preceq I_U(\theta_1) \preceq \frac{\beta + 1}{\beta} I_U(\theta^*).$$
(7)

Since $\overline{\Gamma} = \alpha \Gamma_1 + (1-\alpha)U$, we have that $I_{\overline{\Gamma}}(\theta^*) \succeq \alpha I_{\Gamma_1}(\theta^*)$ which further implies that $I_{\overline{\Gamma}}(\theta^*)^{-1} \preceq \frac{1}{\alpha}I_{\Gamma_1}(\theta^*)^{-1}$. Similarly, since $I_{\overline{\Gamma}}(\theta^*) \succeq (1-\alpha)I_U(\theta^*)$, we can apply Lemma 1 on $\overline{\Gamma}$ to get:

$$\mathbb{E}\left[L_U(\theta_2) - L_U(\theta^*)\right] \le (1 + \widehat{\epsilon}_{m_2}) \frac{\operatorname{Tr}\left(I_{\overline{\Gamma}}(\theta^*)^{-1}I_U(\theta^*)\right)}{m_2} + \frac{R}{m_2^2} \le \frac{1}{\alpha} (1 + \widehat{\epsilon}_{m_2}) \frac{\operatorname{Tr}\left(I_{\Gamma_1}(\theta^*)^{-1}I_U(\theta^*)\right)}{m_2} + \frac{R}{m_2^2}$$
$$\le (1 + \widetilde{\epsilon}_{m_2}) \frac{\operatorname{Tr}\left(I_{\Gamma_1}(\theta^*)^{-1}I_U(\theta^*)\right)}{m_2} + \frac{R}{m_2^2},$$

where $\hat{\epsilon}_{m_2}$, $\tilde{\epsilon}_{m_2} = \mathcal{O}\left(\frac{1}{(1-\alpha)^2} \left(L_1 L_3 + \sqrt{L_2}\right) \sqrt{\frac{\log dm_2}{m_2}}\right) = \mathcal{O}\left(\left(L_1 L_3 + \sqrt{L_2}\right) \frac{\sqrt{\log dm_2}}{m_2^{1/6}}\right)$. From (5) and (7), the right hand side is at most:

 $(1+\widetilde{\epsilon}_{m_2})(\frac{\beta+1}{\beta-1})^2 \frac{\operatorname{Tr}\left(I_{\Gamma_1}(\theta_1)^{-1}I_U(\theta_1)\right)}{m_2} + \frac{R}{m_2^2}$

By definition of Γ_1 , this is at most:

$$(1+\tilde{\epsilon}_{m_2})(\frac{\beta+1}{\beta-1})^2 \frac{\text{Tr}\left(I_{\Gamma^*}(\theta_1)^{-1}I_U(\theta_1)\right)}{m_2} + \frac{R}{m_2^2}$$

Finally, applying (6) and (7), we get that this is at most:

$$(1+\widetilde{\epsilon}_{m_2})\left(\frac{\beta+1}{\beta-1}\right)^4 \frac{\operatorname{Tr}\left(I_{\Gamma^*}(\theta^*)^{-1}I_U(\theta^*)\right)}{m_2} + \frac{R}{m_2^2}$$

The upper bound follows.