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Shape-constrained estimation for modern statistical problems
by
Jake Soloff

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in
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in the
Graduate Division
of the
University of California, Berkeley

Committee in charge:
Professor Adityanand Guntuboyina, Co-chair
Professor Michael I. Jordan, Co-chair
Professor Martin J. Wainwright

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Shape-constrained estimation for modern statistical problems

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Jake Soloff


#### Abstract

Shape-constrained estimation for modern statistical problems


by
Jake Soloff

Doctor of Philosophy in Statistics
University of California, Berkeley

Professor Adityanand Guntuboyina, Co-chair
Professor Michael I. Jordan, Co-chair

Shape constraints encode a relatively weak form of prior information specifying the direction of certain relationships in an unknown signal. Classical examples include estimation of a convex function or a monotone density. Shape constraints are often strong enough to dramatically reduce statistical complexity while still yielding flexible, nonparametric estimators. This thesis brings shape constraints to bear on several recent research areas in statistics - distribution-free inference, high-dimensional covariance estimation, empirical Bayes, and multiple hypothesis testing.

Chapter 2 discusses my joint work with Professor Aditya Guntuboyina and Professor Jim Pitman on distribution-free properties of isotonic regression. In this work, we establish a distributional result for the components of the isotonic least squares estimator using its characterization as the derivative of the greatest convex minorant of a random walk. Provided the walk has exchangeable increments, we prove that the slopes of the greatest convex minorant are distributed as order statistics of the running averages. This result implies an exact formula for the squared error risk of least squares in homoscedastic isotonic regression when the true sequence is constant that holds for every exchangeable error distribution.

Chapter 3 discusses my joint work with Professor Aditya Guntuboyina and Professor Michael I. Jordan on sign-constrained precision matrix estimation. We investigate the problem of high-dimensional covariance estimation under the constraint that the partial correlations are nonnegative. The sign constraints dramatically simplify estimation: the Gaussian maximum likelihood estimator is well defined with only two observations regardless of the number of variables. We analyze its performance in the setting where the dimension may be much larger than the sample size. We establish that the estimator is both high-dimensionally consistent and minimax optimal in the symmetrized Stein loss. We also prove a negative
result which shows that the sign-constraints can introduce substantial bias for estimating the top eigenvalue of the covariance matrix.

Chapter 4 discusses my joint work with Professor Aditya Guntuboyina and Professor Bodhisattva Sen on nonparametric empirical Bayes with multivariate, heteroscedastic Gaussian errors. Multivariate, heteroscedastic errors complicate statistical inference in many largescale denoising problems. Empirical Bayes is attractive in such settings, but standard parametric approaches rest on assumptions about the form of the prior distribution which can be hard to justify and which introduce unnecessary tuning parameters. We extend the nonparametric maximum likelihood estimator (NPMLE) for Gaussian location mixture densities to allow for multivariate, heteroscedastic errors. NPMLEs estimate an arbitrary prior by solving an infinite-dimensional, convex optimization problem; we show that this convex optimization problem can be tractably approximated by a finite-dimensional version. We introduce a dual mixture density whose modes contain the atoms of every NPMLE, and we leverage the dual both to establish non-uniqueness in multivariate settings as well as to construct explicit bounds on the support of the NPMLE.

The empirical Bayes posterior means based on an NPMLE have low regret, meaning they closely target the oracle posterior means one would compute with the true prior in hand. We prove an oracle inequality implying that the empirical Bayes estimator performs at nearly the optimal level (up to logarithmic factors) for denoising without prior knowledge. We provide finite-sample bounds on the average Hellinger accuracy of an NPMLE for estimating the marginal densities of the observations. We also demonstrate the adaptive and nearlyoptimal properties of NPMLEs for deconvolution. We apply the method to two astronomy datasets, constructing a fully data-driven color-magnitude diagram of 1.4 million stars in the Milky Way and investigating the distribution of chemical abundance ratios for 27 thousand stars in the red clump.

Chapter 5 discusses my joint work with Daniel Xiang and Professor William Fithian on finite-sample control of the maximum local false discovery rate in multiple hypothesis testing. Despite the popularity of the false discovery rate (FDR) as an error control metric for largescale multiple testing, its close Bayesian counterpart the local false discovery rate (lfdr), defined as the posterior probability that a particular null hypothesis is false, is a more directly relevant standard for justifying and interpreting individual rejections. However, the lfdr is difficult to work with in small samples, as the prior distribution is typically unknown. We propose a simple multiple testing procedure and prove that it controls the expectation of the maximum lfdr across all rejections; equivalently, it controls the probability that the rejection with the largest p-value is a false discovery. Our method operates without knowledge of the prior, assuming only that the $p$-value density is uniform under the null and decreasing under the alternative. We also show that our method asymptotically implements the oracle Bayes procedure for a weighted classification risk, optimally trading off between false positives and false negatives. We derive the limiting distribution of the attained maximum lfdr over the rejections, and the limiting empirical Bayes regret relative to the oracle procedure.

To my family.

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## Chapter 1

## Introduction

The study of shape-constrained estimation and inference traces back to the mid-twentieth century, when researchers introduced isotonic regression (Ayer et al., 1955; van Eeden, 1956) and monotone density estimation (Grenander, 1956). Both problems constrain an unknown function - either a regression function or a probability density function - to be monotone. Although the target estimand need not be differentiable, it is helpful to view shape restrictions such as monotonicity or convexity as constraining the sign of a (first or second) derivative of an unknown function. By contrast, smoothness assumptions in density estimation and regression commonly constrain the magnitude of one or more derivatives of a target function. A resulting attraction of shape constraints is that one can often perform constrained maximum likelihood estimation with no explicit regularization to yield tuning-free, nonparametric estimators.

Conversely, quantifying the full benefits of such qualitative restrictions can be a delicate exercise. In recent decades, a resurgence in theoretical research has greatly expanded our mathematical toolkit for analyzing shape constraints. Some important developments include Barber and Samworth (2021), Cai and Low (2015), Cai et al. (2013), Chatterjee (2014), Dümbgen (2003), Dümbgen et al. (2011), Groeneboom et al. (2001), Guntuboyina and Sen (2013), Han et al. (2019), Kim and Samworth (2016), Meyer and Woodroofe (2000), Slawski and Hein (2013), Wei et al. (2019), and Zhang (2002); see also Guntuboyina and Sen (2018) for a recent survey and Groeneboom and Jongbloed (2014) for a general introduction.

As this now rich statistical literature routinely demonstrates, shape-constrained estimation has the potential to alleviate the challenges of nonparametric estimation while preserving the flexibility of light assumptions. Realizing this potential in a wide variety of applications, however, demands innovations on other research fronts. On one important front, designing scalable algorithms for large samples, high-dimensional data, and complex constraints presents many challenges. Koenker and Mizera (2014) approximate nonparametric maximum likelihood estimators with finite-dimensional, convex optimization problems, opening up nonparametric density estimation to a wide range of existing algorithms and off-the-shelf solvers. In multivariate convex regression, Mazumder et al. (2019) illustrate the flexibility of the alternating direction method of multipliers (ADMM), which itself is closely related to
the algorithm of Dykstra (1983) for constrained least squares.
Beyond computation and risk bounds, getting shape constraints to work with other areas of statistics requires a richer appreciation of the varied roles constraints play in framing and simplifying problems. In this dissertation, we investigate the role of shape constraints in four modern research areas, and for the rest of the introduction, we highlight the varied roles of shape constraints in these and other problems.

Problem specification. Likelihood-based estimation often fails in high-dimensional spaces, and shape constraints offer one approach to rescue nonparametric estimation from the limitations of the maximum likelihood criterion over large parameter spaces. A prototypical example comes from density estimation. Let $p_{1}, \ldots, p_{m}$ be an iid sample of continuous random variables on $[0,1]$ with density $f^{*}$ and corresponding cdf $F^{*}$. It is well known that the maximum likelihood estimator (MLE) of $F^{*}$ over the space of all cdfs is the empirical cdf $F_{m}$, defined as $F_{m}(t):=\frac{1}{m} \sum_{i=1}^{m} 1\left\{p_{i} \leq t\right\}$. The MLE of the pdf $f^{*}$, on the other hand, does not exist: essentially, since $F_{m}$ is discrete, the likelihood of a density $f$ increases without bound as we move outside of the space of all densities on $[0,1]$ and towards a discrete distribution.

Imposing the additional constraint that $f^{*}$ is monotone decreasing enables maximum likelihood estimation. To see why, first note that a density $f$ is nonincreasing precisely when the corresponding cdf $F$ is concave. The empirical cdf $F_{m}$ is not concave, so the constrained MLE $\hat{F}_{m}$ of the cdf $F^{*}$ over all concave cdfs is a different object than $F_{m}$. In fact, the likelihood under the concave MLE $\hat{F}_{m}$ is lower than the likelihood under the unconstrained MLE $F_{m}$, since the latter maximizes the likelihood over all cdfs. Roughly, in order to make the gap due to introducing the concavity constraint as small as possible, $\hat{F}_{m}$ should be as close to $F_{m}$ as possible. More precisely, it can be shown (Groeneboom \& Jongbloed, 2014, Lemma 2.2) that $\hat{F}_{m}$ is the least concave majorant of $F_{m}$, i.e. the (pointwise) smallest concave function dominating $F_{m}$. The least concave majorant $\hat{F}_{m}$ is a linear spline, so it is 'nearly' differentiable, i.e. differentiable everywhere except at the knots. Since $\hat{F}_{m}$ is concave the slope of the line to the left of a given knot is larger than the slope of the line to the right, and indeed the Grenander (1956) estimator $\hat{f}_{m}$ of a monotone density is given by the left derivative of the least concave majorant $\hat{F}_{m}$. The empirical cdf $F_{m}$ plays the role of a 'default' estimator of $F^{*}$ across statistics, e.g. in the nonparametric bootstrap, so its scope and limitations as a plug-in estimator are well understood. The Grenander estimator $\hat{f}_{m}$ plays a similar role as a 'default' density estimator of a monotone density.

The Grenander estimator $\hat{f}_{m}$ is somewhat special among shape-constrained density estimators, in that it has this exact geometric characterization, especially amenable to precise asymptotic theory. Moreover, there are many equivalent ways of thinking about the problem of monotone density estimation that connect it to other types of shape constraints. For instance, Khintchine's theorem states that the collection of monotone densities on $[0,1]$ coincides with the collection of densities of uniform scale mixtures, i.e. distributions of the form $\operatorname{Unif}(0, \Theta)$ where $\Theta \sim G^{*}$ is an arbitrary probability measure on $[0,1]$. Hence, the Grenander estimator $\hat{f}_{m}$ implicitly describes an estimator $\widehat{G}_{m}$ of the mixing measure $G^{*}$,
known as the nonparametric maximum likelihood estimator or NPMLE of $G^{*}$. Since $\hat{F}_{m}$ is a (concave) linear spline, the corresponding density $\hat{f}_{m}$ is a piecewise constant (decreasing) density with discontinuities at the knots of $\hat{F}_{m}$; equivalently, the NPMLE $\widehat{G}_{m}$ is a discrete distribution supported on the knots of $\hat{F}_{m}$. These statements are equivalent, but the latter is remarkable for its own reason: $\widehat{G}_{m}$ maximizes the likelihood of the uniform scale mixture over all probability measures supported on $[0,1]$. Hence, whereas the Grenander estimator $\hat{f}_{m}$ is a constrained MLE over the class of monotone densities, $\widehat{G}_{m}$ is an unconstrained NPMLE where the likelihood is that of a uniform scale mixture. The key takeaway is that the mixture structure can facilitate fully nonparametric modeling of unobserved heterogeneity, a point which surprisingly predates even the Grenander estimator (Robbins, 1950). The idea to use the NPMLE of a mixing distribution to model a prior distribution is foundational to nonparametric empirical Bayes, as we shall see in Chapter 4.

More recent results have begun to uncover multivariate settings with flexible 'default' estimators. Suppose we aim to estimate the density $f^{*}$ of i.i.d. observations $X_{1}, \ldots, X_{n} \in \mathbb{R}^{p}$ under the constraint that $\log f^{*}$ is concave: Cule et al. (2010) showed that the MLE is welldefined almost surely once $n>p$. Another remarkable result, due to Slawski and Hein (2015), establishes a setting where shape constraints lead to a well-defined constrained MLE even in very high-dimensional settings $p \gg n$. Instead of estimating the density of $X_{1}, \ldots, X_{n} \in \mathbb{R}^{p}$, the goal is to estimate the precision matrix $\Theta^{*}:=\left(\Sigma^{*}\right)^{-1}$ under a multivariate Gaussian likelihood $X_{1}, \ldots, X_{n} \sim \mathcal{N}\left(0, \Sigma^{*}\right)$. When $p>n$, the MLE of $\Theta^{*}$ is not defined, since the sample covariance matrix $S=\frac{1}{n} \sum_{i=1}^{n} X_{i} X_{i}^{\top}$ is rank-deficient. However, under the additional sign-constraints $\Theta_{j k}^{*} \leq 0$ for $j \neq k$, Slawski and Hein (2015) showed that the constrained MLE $\widehat{\Theta}$ is almost surely well-defined for any $p$ once $n>1$. The sign-constraints on the off-diagonal entries of $\Theta^{*}$ provide a model of positive dependence that has been widely used in economics and actuarial sciences (Karlin \& Rinott, 1983). The fact that the constrained MLE $\widehat{\Theta}$ is well-defined is nontrivial, but Lauritzen et al. (2019) give a beautiful proof, exploiting a connection to single-linkage clustering. We investigate the statistical properties of this estimator $\widehat{\Theta}$ in high-dimensional settings in Chapter 3.

It is worth noting that reasonable shape constraints do not always impose sufficient structure to produce well-defined solutions. A classical example is the estimation of a unimodal density, even with univariate data, where the likelihood is unbounded as in the unconstrained density estimation example. In Chapter 4 we show that certain NPMLEs of multivariate mixing distributions can have infinitely many solutions.

Finally, we note that shape restrictions have been leveraged to specify a problem structure beyond producing well-defined constrained MLEs. In causal inference, for example, Angrist and Imbens (1994) leveraged a monotonicity assumption on compliance to identify the local average treatment effect under two-sided noncompliance.

Structured decisions. Shape-enforcement is often viewed as a post-processing step to improve upon a base estimator (Bonakdarpour et al., 2018). Obozinski et al. (2008) followed this approach to yield logically consistent predictions in protein function annotation. In this domain, various data sources are combined to produce probabilistic predictions of protein
function. The attributes that an individual protein may satisfy fit into the gene ontology (GO), a dictionary of terms together with a topology describing their logical relations. For instance, 'nucleic acid binding' and 'protein binding' are two attributes that have 'binding' as a parent attribute. Probabilistic predictions are made on a term-by-term basis, e.g. the probability that a given protein is nucleic acid binding is estimated independently of its other attributes, so this base set of predictions need not respect the GO structure - e.g. 'binding' could be reported as less likely than one of its child attributes. These probabilistic predictions can be 'reconciled' by projecting them onto the set of probabilistic predictions respecting the GO partial order. When the initial predictions are reconciled using Euclidean projection, this is known as isotonic regression with respect to a partial order. Obozinski et al. (2008) find empirically that the isotonized predictions almost always have higher precision than the base predictions, whereas more heuristic reconciliation methods and Bayesian networks often perform worse than the base predictions.

In empirical Bayes applications, downstream decisions often satisfy desirable properties as a byproduct of the mixture constraint. Efron and Hastie (2016, Chapter 6) describe the example of an insurance company needing to estimate the number of claims a policy holder will make next year based on the number of claims they made this year. The dataset is simply a collection of nonnegative counts $Y_{1}, \ldots, Y_{n}$, and from the law of rare events it is reasonable to model each policy holder's number of claims $Y_{i}$ as a Poisson random variable with rate parameter $\lambda_{i}$, independent across $i$. The rate parameters $\lambda_{i}$ likely vary across policy holders $i=1, \ldots, n$, and we can model this unobserved heterogeneity by placing a common prior on the rate parameters

$$
\lambda_{i} \stackrel{\mathrm{iid}}{\sim} G^{*}, \text { for } i=1, \ldots, n,
$$

where $G^{*}$ is a probability measure supported on the nonnegative reals $\mathbb{R}_{+}$. Marginally, the counts $Y_{i}$ are iid with Poisson mixture $\operatorname{pmf} p_{G^{*}}$, where

$$
p_{G^{*}}(y):=\int e^{-\lambda} \lambda^{y} / y!\mathrm{d} G^{*}(\lambda) .
$$

Robbins (1956) showed that the Bayes estimator $\delta_{G^{*}}(y):=\mathbb{E}\left[\lambda_{i} \mid Y_{i}=y\right]$ can be written directly in terms of the marginal pmf

$$
\delta_{G^{*}}(y)=\frac{(y+1) p_{G^{*}}(y+1)}{p_{G^{*}}(y)} .
$$

It is not hard to show that, no matter the choice of prior $G^{*}$, the posterior mean $\delta_{G^{*}}$ is nondecreasing in the number of counts. This is a desirable property for the insurance company: any non-monotone decision rule penalizes some policy holders for making fewer claims than others. In practice, however, $G^{*}$ is unknown so the Bayes estimator is unavailable. Robbins (1956) proposed to plug-in the empirical probability mass function (pmf) $\hat{p}_{n}(y)=\frac{\#\left\{i: Y_{i}=y\right\}}{n}$ and take

$$
\hat{\delta}(y)=\frac{(y+1) \hat{p}_{n}(y+1)}{\hat{p}_{n}(y)}
$$

for any $y$ such that $\hat{p}_{n}(y)>0$, and define $\hat{\delta}(y)=0$ otherwise. This estimator remains the 'default' choice for empirical Bayes denoising of count data, despite severely violating monotonicity: the largest count $y_{\max }=\max _{i=1: n} Y_{i}$ is assigned a predicted value of $\hat{\delta}\left(y_{\max }\right)=0$. Brown et al. (2013) patched this deficiency of Robbins' estimator by isotonizing $\hat{\delta}$, among other post-processing proposals. By contrast, if we estimate $p_{G^{*}}$ by any Poisson mixture $p_{\widehat{G}_{n}}$, for instance by minimizing $D\left(\hat{p}_{n}, p_{G}\right)$ over all probability measures $G$ on $\mathbb{R}_{+}$, for some divergence $D$, the resulting plug-in procedure is automatically monotone. For instance, if $D$ is the KL divergence, then $\widehat{G}_{n}$ is the Poisson mixture NPMLE. In this example, utilizing the shape constraint in an end-to-end fashion spares us from post-processing adjustments to achieve the desired structure of a decision rule.

Sharper inferences. In other cases, the desired structure of optimal decision rules delivers a tractable shape constraint. In Chapter 5, we study multiple hypothesis testing within a Bayes two-groups model

$$
p_{i} \mid H_{i}=h \stackrel{\text { ind }}{\sim} f_{h}, \quad \text { with } \quad H_{i} \stackrel{\text { iid }}{\sim} \operatorname{Bern}\left(1-\pi_{0}\right), \quad \text { for } i=1, \ldots, m,
$$

where $H_{i}=0$ if the $i$ th hypothesis is null and $H_{i}=1$ otherwise. The $p$-values $p_{i}$ follow a density $f_{0}:=1_{[0,1]}$ under the null and $f_{1}$ under the alternative, and the null proportion is $\pi_{0} \in[0,1]$. Let $f:=\pi_{0}+\left(1-\pi_{0}\right) f_{1}$ denote the common mixture density of the $p$-values, and let $F(t):=\int_{0}^{t} f(u) d u$ denote the corresponding cumulative distribution function (cdf). Sun and Cai (2007) showed that optimally trading off false positives and false negatives is achieved by rejecting hypothesis with local false discovery rate (lfdr, Efron et al., 2001)

$$
\operatorname{lfdr}(t):=\mathbb{P}\left(H_{i}=0 \mid p_{i}=t\right)=\frac{\pi_{0}}{f(t)}
$$

falling below some level $q \in[0,1]$, depending our relative tolerance for Type I and Type II errors. In other words, the optimal procedures reject the null hypothesis for $p_{i} \in A_{q}^{*}$, where the rejection regions $A_{q}^{*}$ have the form

$$
A_{q}^{*}:=\{t: \operatorname{lfdr}(t) \leq q\}
$$

In most cases, we seek to reject all the $p$-values falling below some threshold, meaning the rejection region is an interval of the form $[0, t]$. If we posit that the optimal rejection regions also have this form $A_{q}^{*}=\left[0, \tau_{q}^{*}\right]$, note that $\tau_{q}^{*}$ is nondecreasing as a function of $q$, simply because the rejection regions $A_{q}^{*}$ are automatically nested as the tolerance $q$ increases. From this, it follows that lfdr is nondecreasing, or equivalently that the mixture density $f$ is nonincreasing. The assertion that it is optimal to reject $p$-values below some threshold is thus equivalent to the assertion that $f$ is monotone.

In other words, $f$ being monotone nonincreasing means that smaller $p$-values represent stronger evidence against the null, and if the optimal decision rule rejects small $p$ values then $f$ must be monotone. Fundamental methods in multiple hypothesis testing, such as the Benjamini and Hochberg $(1995, \mathrm{BH})$ procedure, are thresholding procedures, so such
procedures are best suited to settings where $f$ is monotone. The BH procedure builds a rejection region of the form

$$
A_{q}^{\mathrm{BH}}:=\{t: \widehat{\operatorname{Fdr}}(t) \leq q\} \quad \text { where } \quad \widehat{\operatorname{Fdr}}(t)=\frac{1}{F_{m}(t)} .
$$

This procedure is well-motivated from an empirical Bayes perspective (see, e.g., Efron et al., 2001) but it does not directly target the optimal rejection regions $A_{q}^{*}$. Moreover, enforcing the shape restriction that the true cdf $F$ is concave yields little benefit, since the concave MLE $\hat{F}_{m}$ is extremely close to $F_{m}$ at least for large $m$ (Kiefer \& Wolfowitz, 1976).

The real advantage of the monotonicity assumption is that we may leverage the Grenander estimator $\hat{f}_{m}$ as our 'default' estimator of a monotone density to directly target the optimal rejection region:

$$
A_{q}:=\{t: \widehat{\operatorname{lfdr}}(t) \leq q\} \quad \text { where } \quad \widehat{\operatorname{lfdr}}(t)=\frac{1}{\hat{f}_{m}(t)}
$$

The shape constraint thus affords us the opportunity to approach a much more ambitious inferential target in a fully nonparametric manner. We explore the properties of this testing procedure in much greater detail in Chapter 5.

We close this section with a cautionary result showing that, in some cases, shape constraints can have much worse statistical properties than the base estimators they are designed to improve upon. In particular, we show in the high-dimensional covariance estimation problem in Chapter 3 that the constrained MLE of the covariance matrix $\widehat{\Sigma}=\widehat{\Theta}^{-1}$ can be much worse than the sample covariance matrix $S$ for estimating the top eigenvalue. Note that the sample covariance $S$ is itself inconsistent in high-dimensional regimes, but the top-eigenvalue of $\widehat{\Sigma}$ diverges at an even faster rate. However, we argue that because $\widehat{\Theta}$ essentially represents a projection of $S$ under a Bregman divergence known as Stein's loss, it is more natural to study the high-dimensional behavior of $\widehat{\Theta}$ under that loss function.

## Chapter 2

## Distribution-free isotonic regression

### 2.1 Introduction

Isotonic regression with homoscedastic errors refers to the problem of estimating a monotone sequence $\theta_{1}^{*} \leq \cdots \leq \theta_{n}^{*}$ based on a noisy observation vector $Y$, assumed to be an additive perturbation of $\theta^{*}=\left(\theta_{1}^{*}, \ldots, \theta_{n}^{*}\right)$

$$
Y=\theta^{*}+\sigma Z
$$

where the components $Z_{1}, \ldots, Z_{n}$ of $Z$ are assumed to have zero mean and unit variance. It is commonly assumed that $Z_{1}, \ldots, Z_{n}$ are independent and identically distributed (i.i.d.) but we work with the more general assumption of exchangeability in this chapter. A natural estimator for $\theta^{*}$ in this setting is the isotonic Least Squares Estimator (LSE), defined as

$$
\hat{\theta}:=\Pi_{\mathcal{M}^{n}}(Y):=\underset{\theta \in \mathcal{M}^{n}}{\operatorname{argmin}}\|Y-\theta\|_{2}^{2},
$$

where $\|\cdot\|_{2}$ denotes the usual Euclidean norm on $\mathbb{R}^{n}$ and $\mathcal{M}^{n}:=\left\{\theta \in \mathbb{R}^{n}: \theta_{1} \leq \cdots \leq \theta_{n}\right\}$ is the monotone cone of length $n$ nondecreasing sequences. As $\mathcal{M}^{n}$ is a closed convex cone, $\hat{\theta}$ as defined above exists uniquely; it can also be computed in $O(n)$ time by the pool adjacent violators algorithm (Brunk et al., 1972; Grotzinger \& Witzgall, 1984).

One approach to evaluating the statistical properties of $\hat{\theta}$ is to measure the risk, or expected deviation of $\hat{\theta}$ from $\theta^{*}$. Indeed, the risk provides a convenient summary of the accuracy of $\hat{\theta}$ and many papers on isotonic regression have focused on obtaining bounds for the risk of $\hat{\theta}$ (see e.g., Bellec, 2018; Guntuboyina and Sen, 2018; Zhang, 2002). In this chapter, we primarily consider the normalized mean squared error:

$$
R\left(\hat{\theta}, \theta^{*}\right):=\frac{1}{n} \mathbb{E}_{\theta^{*}}\left\|\hat{\theta}-\theta^{*}\right\|_{2}^{2}
$$

A key quantity in understanding $R\left(\hat{\theta}, \theta^{*}\right)$ is

$$
\delta_{n}(\mu):=\mathbb{E}_{Z \sim \mu}\left\|\Pi_{\mathcal{M}^{n}}(Z)\right\|_{2}^{2}
$$

where $\mu$ denotes the law of the noise vector $Z$. Indeed, it is clear that

$$
\frac{n}{\sigma^{2}} R\left(\hat{\theta}, \theta^{*}\right)=\delta_{n}(\mu) \quad \text { when } \theta_{1}^{*}=\cdots=\theta_{n}^{*}
$$

When $\theta_{1}^{*} \leq \cdots \leq \theta_{n}^{*}$ are not all equal, let $\left(A_{1}, \ldots, A_{k}\right)$ be the coarsest partition of $\{1, \ldots, n\}$ such that $\theta^{*}$ is constant on each $A_{i}$. It has been shown (Bellec, 2018; Fang \& Guntuboyina, 2017; Oymak \& Hassibi, 2016) that

$$
\frac{n}{\sigma^{2}} R\left(\hat{\theta}, \theta^{*}\right) \begin{cases}\leq \delta_{n_{1}}\left(\mu_{A_{1}}\right)+\cdots+\delta_{n_{k}}\left(\mu_{A_{k}}\right) & \text { for every } \sigma>0  \tag{2.1}\\ \rightarrow \delta_{n_{1}}\left(\mu_{A_{1}}\right)+\cdots+\delta_{n_{k}}\left(\mu_{A_{k}}\right) & \text { as } \sigma \downarrow 0\end{cases}
$$

where $\mu_{A_{i}}$ denotes the marginal distribution of $\left(Z_{j}\right)_{j \in A_{i}}$ and $n_{i}=\left|A_{i}\right|$ is the length of the $i^{\text {th }}$ block for all $i=1, \ldots, k$. We emphasize that (2.1) holds for arbitrarily dependent $Z_{1}, \ldots, Z_{n}$ with zero mean and finite variance. It was also shown by Bellec (2018) that $\delta_{n}(\mu)$ also bounds the risk of the isotonic LSE in misspecified settings where $\theta^{*}$ does not lie in $\mathcal{M}^{n}$.

The quantity $\delta_{n}(\mu)$ therefore crucially controls the risk of the isotonic LSE. The goal of this chapter is to explicitly determine $\delta_{n}(\mu)$ for every $n \geq 1$ under the additional assumption that $Z$ is exchangeable. Specifically, under the assumption of exchangeability, we show in Corollary 2.6 that, for all $n$,

$$
\begin{equation*}
\delta_{n}(\mu)=\rho n+(1-\rho) H_{n}, \tag{2.2}
\end{equation*}
$$

where $H_{n}:=1+\frac{1}{2}+\cdots+\frac{1}{n}$ is the $n^{\text {th }}$ harmonic number, $\rho=\operatorname{Cor}\left(Z_{1}, Z_{2}\right)$ is the pairwise correlation, and $\sigma^{2}=1$. Combined with (2.1), our result provides a sharp, non-asymptotic bound on the risk of isotonic regression for any exchangeable noise vector. In the special case when $Z_{1}, \ldots, Z_{n}$ are i.i.d. with zero mean and unit variance, $\rho=0$ and thus (2.2) gives:

$$
\begin{equation*}
\delta_{n}\left(\otimes_{i=1}^{n} \eta\right)=H_{n} \quad \text { for every probability measure } \eta \text {. } \tag{2.3}
\end{equation*}
$$

Here $\eta$ is the common distribution of the independent variables $Z_{1}, \ldots, Z_{n}$.
Previously, the formula (2.3) was known when $\eta$ is the standard Gaussian probability measure on $\mathbb{R}^{n}$. This was observed by Amelunxen et al. (2014) who proved it by observing first that when $\mu=\otimes_{i=1}^{n} \eta$ and $\eta$ is the standard Gaussian measure, the formula

$$
\begin{equation*}
\mathbb{E}\left\|\Pi_{K}(Z)\right\|_{2}^{2}=\sum_{k=0}^{n} k \nu_{k}(K) \tag{2.4}
\end{equation*}
$$

holds for every closed convex cone $K \subseteq \mathbb{R}^{n}$ where $\nu_{k}(K)$ is the $k^{\text {th }}$ intrinsic volume of $K$. When $K=\mathcal{M}^{n}$ is the monotone cone, the right hand side in equation (2.4) can be shown to be equal to $H_{n}$ by using the fact that the generating function $s \mapsto \sum_{k=0}^{n} s^{k} \nu_{k}\left(\mathcal{M}^{n}\right)$ can be computed in closed form. Amelunxen et al. (2014) used the theory of finite reflection groups (Coxeter \& Moser, 2013) to obtain the exact expression for this generating function. However, the exact expression for $\sum_{k=0}^{n} s^{k} \nu_{k}\left(\mathcal{M}^{n}\right)$ can already be found in the classical
literature on isotonic regression (see Theorem 2.4.2 in Robertson et al. (1988) or Section 8 of Sparre-Andersen (1954)).

The above proof does not work for non-Gaussian $\eta$ mainly because the expression (2.4) does not hold for general $\eta$. In fact, the best available result on $\delta_{n}\left(\otimes_{i=1}^{n} \eta\right)$ for non-Gaussian $\eta$ is in equation (2.11) of Zhang (2002), who proved the asymptotic result:

$$
\delta_{n}\left(\otimes_{i=1}^{n} \eta\right)=(1+o(1))(1+\log n) \quad \text { as } n \rightarrow \infty
$$

This bound gives the right behavior as the right hand side of equation (2.3) but only as $n \rightarrow \infty$. We improve this result by proving for every $n \geq 1$ that $\delta_{n}\left(\otimes_{i=1}^{n} \eta\right)$ is always equal to the $n^{\text {th }}$ harmonic number $H_{n}$ for every probability measure $\eta$ having mean 0 and variance 1 .

We prove (2.2) by developing a precise characterization of the marginal distribution of each individual component $\left(\Pi_{\mathcal{M}^{n}}(Z)\right)_{k}$ of $\Pi_{\mathcal{M}^{n}}(Z)$. Specifically, as long as $Z$ is exchangeable, we show in Theorem 2.2 that $\left(\Pi_{\mathcal{M}^{n}}(Z)\right)_{k}$ has the same distribution as $\bar{Z}_{(k)}$, the $k^{\text {th }}$ order statistic of the running averages $\bar{Z}_{j}=\frac{Z_{1}+\cdots+Z_{j}}{j}$. We prove Theorem 2.2 in Section 2.2, using a characterization of the components of the isotonic LSE as the left-hand slopes of the greatest convex minorant of the random walk with increments $Z_{1}, \ldots, Z_{n}$. This result and its continuous-time analogue may be of independent interest outside the study of isotonic regression, so in Section 2.2 we also address consequences for the greatest convex minorant of a stochastic process with exchangeable increments. The order statistics of the running averages $\left\{\bar{Z}_{k}\right\}_{k=1}^{n}$ can be fairly complicated even when $Z$ is Gaussian; however, Theorem 2.2 easily implies results such as (2.2). In Section 2.3, we detail some risk calculations for isotonic regression and its variants which all follow from Theorem 2.2.

### 2.2 Main result

Let $S_{k}=\sum_{i=1}^{k} Z_{i}$ denote the partial sums for $k=1, \ldots, n$, started at $S_{0}=0$. Identify the random walk $\left\{S_{k}\right\}_{k=0}^{n}$ with its cumulative sum diagram $S:[0, n] \rightarrow \mathbb{R}$, where $S(k)=S_{k}$ for integers $k=0, \ldots, n$ and linearly interpolated between integers. Let $C:[0, n] \rightarrow \mathbb{R}$ denote the greatest convex minorant (GCM) of $S$, i.e. the greatest convex function that lies below $S$. See Figure 2.1 for a depiction of the GCM of $S$. With this notation, we now recall the graphical representation of the isotonic LSE as given in Theorem 1.2.1 of Robertson et al. (1988).

Lemma 2.1. For any vector $Z$, the isotonic $L S E \Pi_{\mathcal{M}^{n}}(Z)$ is given by the left-hand slopes of the greatest convex minorant of the cumulative sum diagram. For all $k=1, \ldots, n$

$$
\left(\Pi_{\mathcal{M}^{n}}(Z)\right)_{k}=C(k)-C(k-1)=\partial_{-} C(k) .
$$

For the remainder of this section let

$$
\begin{equation*}
\Delta_{k}:=\partial_{-} C(k)=\min _{k \leq v \leq n} \max _{0 \leq u<k} \frac{S_{v}-S_{u}}{v-u} \tag{2.5}
\end{equation*}
$$



Figure 2.1: Solid blue curve is the cumulative sum diagram $S$ of increments $Z_{1}, \ldots, Z_{n}$; dashed black curve is the greatest convex minorant $C$ of $S$.
denote the left-hand slope of the GCM at $k$, so $\Delta=\left(\Delta_{1}, \ldots, \Delta_{n}\right)$ is equal to $\Pi_{\mathcal{M}^{n}}(Z)$ by the lemma. In particular, when $k=1$ we have $\Delta_{1}=\min _{1 \leq v \leq n} \frac{S_{v}}{v}$. When $k=n$, we have $\Delta_{n}=\max _{0 \leq u<n} \frac{S_{n}-S_{u}}{n-u}$, and if $\left(Z_{n}, \ldots, Z_{1}\right) \stackrel{d}{=}\left(Z_{1}, \ldots, Z_{n}\right)$ then $\Delta_{n} \stackrel{d}{=} \max _{1 \leq u \leq n} \frac{S_{u}}{u}$. Our next result generalizes this observation, showing that the $k^{\text {th }}$ slope $\Delta_{k}$ is equal in distribution to the $k^{\text {th }}$ smallest running average if $Z$ is exchangeable.

Theorem 2.2. Suppose $Z=\left(Z_{1}, \ldots, Z_{n}\right)$ is exchangeable. Let $\bar{Z}_{k}:=\frac{1}{k} \sum_{i=1}^{k} Z_{i}$ denote the $k^{\text {th }}$ running average for $k=1, \ldots, n$ and let $\bar{Z}_{(1)} \leq \cdots \leq \bar{Z}_{(n)}$ denote their order statistics. Then

$$
\begin{equation*}
\Delta_{k} \stackrel{d}{=} \bar{Z}_{(k)} \tag{2.6}
\end{equation*}
$$

marginally for all $k=1, \ldots, n$.
Proof. As before, let $S_{k}$ denote the $k^{\text {th }}$ partial sum. Let $M$ be the last argmin of the sequence $\left\{S_{i}\right\}_{i=0}^{n}$, and let $N$ be the amount of time the walk is nonpositive $N:=\sum_{i=1}^{n} 1\left(S_{i} \leq 0\right)$. We will use Corollary 11.14 of Kallenberg (2006), due to Sparre-Andersen, which says $M \stackrel{d}{=} N$ as long as $Z$ is exchangeable.

Note that the slope of the GCM switches from nonpositive to positive at time $M$, since the horizontal line with intercept $S_{M}$ minorizes the GCM and touches it at time $M$. Hence, no matter the sequence of increments $Z_{i}$, there is the identity of events

$$
\begin{equation*}
\left(\Delta_{k} \leq 0\right)=(M \geq k) \tag{2.7}
\end{equation*}
$$

Also, for the time $N$ that the walk is nonpositive, since $S_{i} \leq 0$ if and only if $\bar{Z}_{i} \leq 0$, there is the identity of events

$$
\left(\bar{Z}_{(k)} \leq 0\right)=(N \geq k)
$$

The equality in distribution $M \stackrel{d}{=} N$ then implies

$$
\mathbb{P}\left(\Delta_{k} \leq 0\right)=\mathbb{P}\left(\bar{Z}_{(k)} \leq 0\right)
$$

If the sequence $\left\{Z_{i}\right\}$ is modified to $\left\{Z_{i}-z\right\}$ for some fixed $z$, the modified sequence is exchangeable, and the values of $\Delta_{k}$ and $\bar{Z}_{(k)}$ for the modified sequence are just $\Delta_{k}-z$ and $\bar{Z}_{(k)}-z$. Applying the above identity to the modified sequence gives

$$
\mathbb{P}\left(\Delta_{k} \leq z\right)=\mathbb{P}\left(\Delta_{k}-z \leq 0\right)=\mathbb{P}\left(\bar{Z}_{(k)}-z \leq 0\right)=\mathbb{P}\left(\bar{Z}_{(k)} \leq z\right)
$$

So $\Delta_{k}$ and $\bar{Z}_{(k)}$ have the same cumulative distribution function, hence the same distribution.

The proof of Theorem 2.2 generalizes to the setting where $S:[0,1] \rightarrow \mathbb{R}$ is a continuoustime stochastic process. Knight (1996) showed that the analogous distributional identity $M \stackrel{d}{=} N$ holds when $S$ has exchangeable increments and $S(0)=0$. Hence, by a similar proof, we find that the slope $\Delta(p)$ of the greatest convex minorant of $S$ at time $p \in[0,1]$ has the same distribution as the $p^{\text {th }}$ percentile point of the occupation measure for the process $\left(\frac{S(t)}{t}, 0 \leq t \leq 1\right)$. We record this result as the following corollary.

Corollary 2.3. Let $S$ denote a real-valued càdlàg stochastic process on $[0,1]$ with exchangeable increments, such that $S(0)=0$. Define $\Delta(t)$ as the slope of the greatest convex minorant of $S$ at $t$, and let $F: \mathbb{R} \rightarrow[0,1]$ denote the (random) cdf associated with the occupation measure of $\left(\frac{S(t)}{t}, 0 \leq t \leq 1\right)$,

$$
\begin{equation*}
F(x)=\lambda(\{t \in[0,1]: S(t) \leq t x\}), \tag{2.8}
\end{equation*}
$$

where $\lambda$ denotes Lebesgue measure. Then

$$
\begin{equation*}
\Delta(p)=\inf _{p \leq v \leq 1} \sup _{0 \leq u<p} \frac{S(v)-S(u)}{v-u} \stackrel{d}{=} F^{-1}(p) \tag{2.9}
\end{equation*}
$$

marginally for all $p \in[0,1]$.
See Abramson et al. (2011) for a general study of convex minorants of random walks and processes with exchangeable increments. In the special cases where $S$ is a standard Brownian motion or Brownian bridge on the unit interval, Carolan and Dykstra (2001) derive the distribution of the slope $\Delta(p)$, jointly with the process $S(p)$ and its convex minorant at $p$, for a fixed value $p \in[0,1]$. Given our corollary, their explicit formula for the slope $\Delta(p)$ provides the distribution of $F^{-1}(p)$, giving new information about the occupation measure of $\left(\frac{S(t)}{t}, 0 \leq t \leq 1\right)$ for Brownian motion and Brownian bridge. The distribution of the $p^{\text {th }}$ percentile point of the occupation measure for $(S(t), 0 \leq t \leq 1)$ has been obtained under the same generality as Corollary 2.3: see the introduction of Dassios (2005) and references therein.

### 2.3 Consequences for isotonic regression

Since the identity of Theorem 2.2 holds marginally, it allows us to simplify expectations of functions that are additive in the components of $\Pi_{\mathcal{M}^{n}}(Z)$. By Lemma 2.1, the $k^{\text {th }}$ component $\left(\Pi_{\mathcal{M}^{n}}(Z)\right)_{k}=\Delta_{k}$, which by Theorem 2.2 is equal in distribution to $\bar{Z}_{(k)}$. Hence, as long as $Z$ is exchangeable,

$$
\begin{equation*}
\sum_{k=1}^{n} \mathbb{E} h\left(\left(\Pi_{\mathcal{M}^{n}}(Z)\right)_{k}\right)=\sum_{k=1}^{n} \mathbb{E} h\left(\bar{Z}_{(k)}\right)=\sum_{k=1}^{n} \mathbb{E} h\left(\bar{Z}_{k}\right) . \tag{2.10}
\end{equation*}
$$

Taking $h(x)=|x|^{p}$, we obtain our first corollary.
Corollary 2.4. Suppose $Z=\left(Z_{1}, \ldots, Z_{n}\right)$ is exchangeable. For $p>0$,

$$
\begin{equation*}
\mathbb{E}\left\|\Pi_{\mathcal{M}^{n}}(Z)\right\|_{p}^{p}=\sum_{k=1}^{n} \mathbb{E}\left|\frac{1}{k} \sum_{i=1}^{k} Z_{i}\right|^{p} \tag{2.11}
\end{equation*}
$$

provided $\mathbb{E}\left|Z_{1}\right|^{p}<\infty$.
Remark 2.5. Viewed through its graphical representation, $\Delta_{k}=C(k)-C(k-1)$ is the leftderivative of the GCM $C$ at $k$, so when the power $p=1$, equation (2.11) yields the discrete arc-length formula

$$
\begin{equation*}
\sum_{k=1}^{n} \mathbb{E}|C(k)-C(k-1)|=\mathbb{E}| | \Pi_{\mathcal{M}^{n}}(Z) \|_{1}=\sum_{k=1}^{n} \frac{1}{k} \mathbb{E}\left|S_{k}\right| \tag{2.12}
\end{equation*}
$$

Closely related to this formula is the identity of Spitzer and Widom (1961), which takes $\tilde{Z}_{1}, \ldots, \tilde{Z}_{n}$ to be a sequence of i.i.d. random variables in $\mathbb{R}^{2}$ (or the complex plane $\mathbb{C}$ ) with finite variance. If $\tilde{S}_{k}=\sum_{i=1}^{k} \tilde{Z}_{i}$ is the partial sum and $\tilde{L}_{n}$ is the length of the perimeter of the convex hull $\operatorname{conv}\left(0, \tilde{S}_{1}, \ldots, \tilde{S}_{n}\right)$, then

$$
\begin{equation*}
\mathbb{E} \tilde{L}_{n}=2 \sum_{k=1}^{n} \frac{1}{k} \mathbb{E}\left\|\tilde{S}_{k}\right\| \tag{2.13}
\end{equation*}
$$

These formulas connect the geometry of the convex hull of a random walk to the magnitudes of the running means.

Consider the case when $p=2$. Since $Z$ is exchangeable, every pair of components has the same correlation $\rho$. If we further assume $Z_{1}$ has zero mean and unit variance, the right hand side of equation (2.11) can be computed explicitly

$$
\mathbb{E}\left(\frac{1}{k} \sum_{i=1}^{k} Z_{i}\right)^{2}=\rho+\frac{1-\rho}{k}
$$

Summing over $k$ yields our next result.

Corollary 2.6. Suppose $Z \sim \mu$ is an exchangeable random vector with zero mean, unit variance, and pairwise correlation $\rho$. Then

$$
\delta_{n}(\mu)=\rho n+(1-\rho) H_{n} .
$$

This result should be contrasted with other distribution-free identities, namely

$$
\mathbb{E}\|Z\|_{2}^{2}=n \text { and } \mathbb{E}\left\|\bar{Z}_{n} 1_{n}\right\|_{2}^{2}=1
$$

provided $Z$ has i.i.d. components with zero mean and unit variance. In particular, suppose we observe $Y=\theta^{*}+\sigma Z$ where $Z$ has i.i.d. components with zero mean and unit variance, but it turns out that $\theta^{*}=c 1_{n}$ is constant. If we know $\theta^{*}$ is constant, we can estimate it by a constant sequence $\bar{Y} 1_{n}$ and pay a price of $\frac{\sigma^{2}}{n}$ in risk (normalized mean squared error). If we know nothing about the structure of $\theta^{*}$ and use $\hat{\theta}=Y$, the risk $\sigma^{2}$ is quite large by comparison. The monotone sequence estimate resides in the middle, with a much smaller risk of $\frac{H_{n} \sigma^{2}}{n}$ and knowledge only about the relative order.

Theorem 2.2 characterizes the distribution of a component of the isotonic LSE $\hat{\theta}$ when the underlying sequence $\theta^{*}$ is constant. When $\theta^{*} \in \mathcal{M}^{n}$ is not constant, Theorem 2.2 can be applied to characterize the distribution of a component $\hat{\theta}_{i}$ in the low noise limit $\sigma \downarrow 0$. In this limit, the distribution depends only on flat regions of $\theta^{*}$ :

Corollary 2.7. Suppose $Y=\theta^{*}+\sigma Z$, for some $\theta^{*} \in \mathcal{M}^{n}$, and let $\hat{\theta}=\Pi_{\mathcal{M}^{n}}(Y)$ denote the isotonic LSE. Let $\left(A_{1}, \ldots, A_{k}\right)$ be the coarsest partition of $\{1, \ldots, n\}$ such that $\theta^{*}$ is constant on each $A_{j}$, and suppose $Z$ is exchangeable on each of these blocks. If an index $i \in\{1, \ldots, n\}$ belongs to the $j^{\text {th }}$ block, then letting $t_{i}=i+1-\min _{s \in A_{j}} s$ and $X=\left(Z_{s}\right)_{s \in A_{j}}$, we have

$$
\begin{equation*}
\frac{\hat{\theta}_{i}-\theta_{i}^{*}}{\sigma} \xrightarrow{d} \bar{X}_{\left(t_{i}\right)} \text { as } \sigma \downarrow 0 . \tag{2.14}
\end{equation*}
$$

Proof. As $\sigma \downarrow 0$, the ratio $\frac{\hat{\theta}-\theta^{*}}{\sigma}$ tends to the directional derivative $D_{Z} \Pi_{\mathcal{M}^{n}}\left(\theta^{*}\right)$. Lemma 4.6 in Zarantonello (1971) shows that this derivative exists and equals the projection of $Z$ onto the tangent cone $T_{\mathcal{M}^{n}}\left(\theta^{*}\right)$. Hence

$$
\begin{equation*}
\frac{\hat{\theta}-\theta^{*}}{\sigma} \rightarrow \Pi_{T_{\mathcal{M}^{n}}\left(\theta^{*}\right)}(Z) \text { as } \sigma \downarrow 0 . \tag{2.15}
\end{equation*}
$$

From the tangent cone computation in Bellec (2018), we have

$$
\left(\Pi_{T_{\mathcal{M}^{n}}\left(\theta^{*}\right)}(Z)\right)_{i}=\left(\Pi_{\mathcal{M}^{n_{i}}}(X)\right)_{t_{i}}
$$

where $n_{i}=\left|A_{j_{i}}\right|$. Finally, by Theorem 2.2, $\left(\Pi_{\mathcal{M}^{n_{i}}}(X)\right)_{t_{i}} \stackrel{d}{=} \bar{X}_{\left(t_{i}\right)}$.
We explained in Section 2.1 how risk calculations when $\theta^{*}=0$ generalize to MSE bounds that are sharp in the low noise limit for arbitrary $\theta^{*}$. For example, when $\theta^{*} \in \mathcal{M}^{n}$ has $k$
constant pieces, then (2.1), Corollary 2.6 and the fact that $H_{l} \leq \log (e l)$ for every $l \geq 1$ imply that

$$
\begin{equation*}
R\left(\hat{\theta}, \theta^{*}\right) \leq \frac{k \sigma^{2}}{n} \log \left(\frac{e n}{k}\right) \tag{2.16}
\end{equation*}
$$

whenever $Z_{1}, \ldots, Z_{n}$ are i.i.d. with mean zero and unit variance. The bound (2.16) should be compared with the risk of the structure-respecting estimator that averages over the constant blocks and achieves a risk of exactly $\frac{k \sigma^{2}}{n}$ when the blocks are all of size $\frac{n}{k}$. If $\theta^{*} \in \mathbb{R}^{n}$ is not necessarily in $\mathcal{M}^{n}$, then Corollary 2.6, together with the results of Bellec (2018), implies that

$$
R\left(\hat{\theta}, \theta^{*}\right) \leq \inf _{\theta \in \mathcal{M}^{n}}\left(\frac{1}{n}\left\|\theta-\theta^{*}\right\|^{2}+\sigma^{2} \frac{k(\theta)}{n} \log \left(\frac{e n}{k(\theta)}\right)\right)
$$

where $k(\theta)$ is the number of constant pieces of the vector $\theta$. These formulae (with the leading constant of 1 in front of the $\frac{k \sigma^{2}}{n} \log \frac{e n}{k}$ term on the right hand side) were previously only known when the distribution of $Z_{1}, \ldots, Z_{n}$ was standard Gaussian.

Define the $L^{p}$-risk of the isotonic LSE

$$
R^{(p)}\left(\hat{\theta}, \theta^{*}\right)=\frac{1}{n} \mathbb{E}\left\|\hat{\theta}-\theta^{*}\right\|_{p}^{p}
$$

so that $R\left(\hat{\theta}, \theta^{*}\right)=R^{(2)}\left(\hat{\theta}, \theta^{*}\right)$. We can similarly employ Theorem 2.2 to explicitly calculate the $L^{p}$-risk of the isotonic LSE $\hat{\theta}$ when $\theta^{*}$ is constant and $Z$ is Gaussian:

Corollary 2.8. Suppose $Z \sim \mathcal{N}\left(0, I_{n}\right)$. Then for any $p>0$,

$$
\mathbb{E}\left\|\Pi_{\mathcal{M}^{n}}(Z)\right\|_{p}^{p}=H_{n, p / 2} \mathbb{E}\left|Z_{1}\right|^{p}=H_{n, p / 2} \sqrt{\frac{2^{p}}{\pi}} \Gamma\left(\frac{p+1}{2}\right)
$$

where $H_{n, m}=\sum_{k=1}^{n} \frac{1}{k^{m}}$.
Proof. Note $\mathbb{E}\left|\frac{1}{k} \sum_{i=1}^{k} Z_{i}\right|^{p}=\left(\frac{2}{k}\right)^{p / 2} \frac{\Gamma\left(\frac{p+1}{2}\right)}{\sqrt{\pi}}$ and apply the theorem.
Corollary 2.8 should similarly be contrasted with the following identities when $Z \sim$ $\mathcal{N}\left(0, I_{n}\right)$ :

$$
\mathbb{E}\|Z\|_{p}^{p}=n \mathbb{E}\left|Z_{1}\right|^{p} \text { and } \mathbb{E}\left\|\bar{Z} 1_{n}\right\|_{p}^{p}=n^{1-p / 2} \mathbb{E}\left|Z_{1}\right|^{p}
$$

respectively. In particular, when $p>2$, the bound $H_{n, p / 2}<\sum_{k=1}^{\infty} \frac{1}{k^{p / 2}}<\infty$ holds for all $n$, which is to say $\mathbb{E}\left\|\Pi_{\mathcal{M}^{n}}(Z)\right\|_{p}^{p}$ is bounded when $p>2$ whereas $\mathbb{E}\|Z\|_{p}^{p}$ grows without bound as $n$ grows.

When $\theta^{*}$ is constant and $Z \sim \mathcal{N}\left(0, I_{n}\right)$, the $L^{p}$ risk of isotonic regression is

$$
\begin{equation*}
R^{(p)}\left(\hat{\theta}, \theta^{*}\right)=\frac{H_{n, p / 2}}{n} \sigma^{p} \mathbb{E}\left|Z_{1}\right|^{p} \tag{2.17}
\end{equation*}
$$

When $1 \leq p \leq 2$, Theorem 2.3 of Zhang (2002) shows an asymptotic result for the $L^{p}$ risk on constant $\theta^{*}$ that agrees with equation (2.17).

The continuous-time distributional identity in Corollary 2.3 applies to the asymptotic distribution of the isotonic least squares estimator. A standard model for studying the asymptotic behavior of isotonic regression is

$$
\theta_{k}^{*}=f^{*}\left(\frac{k}{n}\right)
$$

where $f^{*}:[0,1] \rightarrow \mathbb{R}$ is nondecreasing. We observe $Y$, a noisy version of $\theta^{*}$, and calculate $\hat{\theta}$ by projecting $Y$ onto the monotone cone. The function estimate $\hat{f}$ is defined by $\hat{f}\left(\frac{k}{n}\right)=\hat{\theta}_{k}$ and linearly interpolated between design points. Here, as before, the dependence on $n$ in $\theta^{*} \in \mathcal{M}^{n}$ is suppressed, but now we are interested in the behavior of isotonic least squares $\hat{f}(p)$ at a fixed point $p \in[0,1]$ as $n \rightarrow \infty$.

Define the partial sum process $S^{(n)}:[0,1] \rightarrow \mathbb{R}$ by $S^{(n)}(k / n)=\frac{Y_{1}+\cdots+Y_{k}}{\sqrt{n}}$, linearly interpolated between design points. When the function $f^{*} \equiv c$ is constant, the quantity

$$
\sqrt{n}\left(\hat{f}(p)-f^{*}(p)\right)
$$

is given by the left-derivative of the greatest convex minorant of $S^{(n)}$ at $p$. By the invariance principle, this converges in distribution to the left-derivative of the greatest convex minorant of standard Brownian motion $B=(B(t), 0 \leq t \leq 1)$ at $t_{0}$. This asymptotic result is well known and a similar result was noted for the Grenander estimator by Carolan and Dykstra (1999), where Brownian motion is replaced with a Brownian bridge. Corollary 2.3 relates this asymptotic distribution to the percentile points of the occupation measure for $\left(\frac{B(t)}{t}, 0 \leq\right.$ $t \leq 1$ ).

Finally, Corollary 2.6 on the projection onto $\mathcal{M}^{n}$ extends over to that of the set of nonnegative monotone sequences $\mathcal{M}_{+}^{n}=\mathcal{M}^{n} \cap \mathbb{R}_{+}^{n}$. Theorem 1 of Németh and Németh (2012) observes that the projection of $Z$ onto $\mathcal{M}_{+}^{n}$ is given by $\Pi_{\mathcal{M}_{+}^{n}}(Z)=\Pi_{\mathcal{M}^{n}}(Z)_{+}$, the element-wise positive part of the projection onto $\mathcal{M}^{n}$. Hence the distributional identity Theorem 2.2 yields a similar set of identities for nonnegative isotonic regression.

Corollary 2.9. For any exchangeable noise vector $Z$,

$$
\begin{equation*}
\left(\Pi_{\mathcal{M}_{+}^{n}}(Z)\right)_{k} \stackrel{d}{=}\left(\bar{Z}_{(k)}\right)_{+} \tag{2.18}
\end{equation*}
$$

Provided $\mathbb{E}\left|Z_{i}\right|^{p}<\infty$,

$$
\begin{equation*}
\mathbb{E}\left\|\Pi_{\mathcal{M}_{+}^{n}}(Z)\right\|_{p}^{p}=\sum_{k=1}^{n} \mathbb{E}\left(\frac{1}{k} \sum_{i=1}^{k} Z_{i}\right)_{+}^{p} \tag{2.19}
\end{equation*}
$$

Furthermore, if $Z$ is symmetric with unit variance, the generalized statistical dimension of the monotone cone is

$$
\begin{equation*}
\mathbb{E}\left\|\Pi_{\mathcal{M}_{+}^{n}}(Z)\right\|_{2}^{2}=\frac{\rho n+(1-\rho) H_{n}}{2} \tag{2.20}
\end{equation*}
$$

where $\rho$ is the pairwise correlation.

Proof. Equation (2.19) follows from equation (2.10) by taking $h(x)=(x)_{+}^{p}$. When $Z_{i} \stackrel{d}{=}-Z_{i}$ is symmetric with unit variance,

$$
\mathbb{E}\left(\frac{1}{k} \sum_{i=1}^{k} Z_{i}\right)_{+}^{2}=\frac{1}{2} \mathbb{E}\left(\frac{1}{k} \sum_{i=1}^{k} Z_{i}\right)^{2}=\frac{1}{2}\left(\rho+\frac{1-\rho}{k}\right) .
$$

Summing over $k$ yields equation (2.20).
Equation (2.20) is also shown by Amelunxen et al. (2014) in the special case $Z \sim \mathcal{N}\left(0, I_{n}\right)$ using the theory of finite reflection groups. The identity (2.19) allows us to show equation (2.20) for a much wider variety of noise vectors, and as before also allows us to obtain relations for the expected $L^{p}$ norms of the projection of the noise vector. All of our exact formulae follow from the distributional identity in Theorem 2.2, which exploits the geometric characterization of the isotonic LSE in Lemma 2.1. An interesting open question is whether similar characterizations - such as for convex regression (Groeneboom et al., 2001) - may yield exact non-asymptotic risk calculations in other shape-constrained estimation problems.

## Chapter 3

## Sign-constrained precision matrix estimation

### 3.1 Introduction

Consider the problem of estimating a $p \times p$ covariance matrix $\Sigma^{*}$ and its inverse $\Theta^{*}:=\left(\Sigma^{*}\right)^{-1}$ from an $n \times p$ data matrix $X$ whose rows are independently distributed according to the multivariate normal distribution $\mathcal{N}\left(0, \Sigma^{*}\right)$ with mean zero and covariance matrix $\Sigma^{*}$. The maximum likelihood estimator (MLE) of $\Theta^{*}$ is given by

$$
\begin{equation*}
\tilde{\Theta}:=\underset{\Theta \in \mathcal{S}_{\succeq 0}^{p \times p}}{\operatorname{argmin}}\{\langle\Theta, S\rangle-\log \operatorname{det} \Theta\}, \tag{3.1}
\end{equation*}
$$

where $\mathcal{S}_{\succeq 0}^{p \times p}$ denotes the set of all $p \times p$ symmetric, positive semi-definite (PSD) matrices, $\langle\Theta, S\rangle:=\operatorname{tr}\left(\Theta^{\top} S\right)$ denotes the Frobenius inner product, and $S$ is the sample covariance matrix, defined as

$$
\begin{equation*}
S:=n^{-1} X^{\top} X \tag{3.2}
\end{equation*}
$$

It is well known that $\tilde{\Theta}$ exists if and only if $S$ is nonsingular, in which case $\tilde{\Theta}=S^{-1}$. In particular, in the high-dimensional setting where $p>n$, the MLE does not exist, since the minimum in (3.1) is not finite. Slawski and Hein (2015) observed, however, that if the optimizer in (3.1) is constrained to lie in the set of $p \times p$ positive semidefinite matrices with nonpositive off-diagonal entries, then, with probability one, the optimum is well-defined and attained for all $n \geq 2$ regardless of the value of $p$. Specifically, let

$$
\mathcal{M}^{p \times p}:=\left\{\Theta \in \mathcal{S}_{\succeq 0}^{p \times p}: \Theta_{j k} \leq 0 \text { for } j \neq k\right\},
$$

and observe that it is the convex cone of symmetric $M$-matrices, an important class of matrices appearing in many contexts (see, e.g., Berman \& Plemmons, 1994, Chap. 6). Slawski and Hein (2015) proved that the optimizer

$$
\begin{equation*}
\widehat{\Theta}:=\underset{\Theta \in \mathcal{M}^{p \times p}}{\operatorname{argmin}}\{\langle\Theta, S\rangle-\log \operatorname{det} \Theta\}, \tag{3.3}
\end{equation*}
$$

exists uniquely as long as, in the observed sample, no two variables are perfectly positively correlated (i.e., $S_{j k}<\sqrt{S_{j j} S_{k k}}$ for all $j \neq k$ ) and no variable is constant (i.e., $S_{j j}>0$ for all $j$ ). Both conditions hold with probability one under the assumed Gaussian model for $n \geq 2$, and thus, unlike the unconstrained MLE in (3.1), the estimator (3.3) is well-defined even in the high-dimensional regime.

The constrained MLE $\widehat{\Theta}$ presents an elegant, tuning-free method for estimating precision matrices which works for $n \geq 2$ and all values of $p$ under the assumption $\Theta^{*} \in \mathcal{M}^{p \times p}$. Efficient algorithms for computing $\widehat{\Theta}$ are given in Slawski and Hein (2015) and Lauritzen et al. (2019). Note that the precision matrix having nonpositive off-diagonal entries $\Theta_{j k}^{*}$ is equivalent to nonnegative partial correlations $-\Theta_{j k}^{*} / \sqrt{\Theta_{j j}^{*} \Theta_{k k}^{*}}$ (Bølviken, 1982). Examples of practical covariance estimation problems with nonnegative partial correlations abound (see, e.g., Agrawal et al., 2019; Lake \& Tenenbaum, 2010; Slawski \& Hein, 2015). More generally, Karlin and Rinott (1983) showed that for the normal distribution the condition that the precision matrix belongs to $\mathcal{M}^{p \times p}$ is equivalent to multivariate total positivity of order two $\left(\mathrm{MTP}_{2}\right) . \mathrm{MTP}_{2}$ is a strong form of positive dependence (Colangelo et al., 2005) that has been widely used in auction theory (Milgrom \& Weber, 1982), actuarial sciences (Denuit et al., 2006), and educational evaluation and policy analysis (Chade et al., 2014).

There is growing interest in $\widehat{\Theta}$ in the graph signal processing literature (Egilmez et al., 2017; Pavez et al., 2018; Pavez \& Ortega, 2016), where $M$-matrices are known as Generalized Graph Laplacians (GGL). Indeed, every graph Laplacian is a diagonally dominant $M$-matrix, and conversely every $M$-matrix $\Theta \in \mathcal{M}^{p \times p}$ can be viewed as a generalized graph Laplacian, in the sense that it has a sparse edge-incidence factorization $\Theta=V V^{\top}$, where $V \in \mathbb{R}^{p \times p(p+1) / 2}$ has at most two nonzero entries per column, whereas positive semidefinite matrices that have other sign patterns typically require dense factorizations (Boman et al., 2005). This connection to nonnegative weighted graphs has led to a host of other application areas in image processing and network analysis.

This chapter investigates the statistical properties of $\widehat{\Theta}$ as an estimator of the unknown precision matrix $\Theta^{*}$ in the high-dimensional regime. Even though $\widehat{\Theta}$ exists uniquely for all $n \geq 2$ regardless of the value of $p$, rigorous results have not yet been proved for the accuracy of $\widehat{\Theta}$ in the high-dimensional regime. In the classical low dimensional asymptotic regime where $p$ is fixed and $n \rightarrow \infty$, Slawski and Hein (2015) apply standard results for $M$-estimators to show consistency of $\widehat{\Theta}$. More recently, Lauritzen et al. (2019) provide an elegant perspective on $\widehat{\Theta}$ and a bound on the support graph $G(\widehat{\Theta})=\left\{(j, k): \widehat{\Theta}_{j k}<0\right\}$, and Wang et al. (2019) develop a consistent estimator of $G\left(\Theta^{*}\right)$.

The study of consistency and optimality properties of $\widehat{\Theta}$ requires fixing an appropriate loss function. Because $\widehat{\Theta}$ is defined via maximum likelihood, it is natural to work with the Stein loss:

$$
\begin{equation*}
L^{\mathrm{s}}\left(\Theta, \Theta^{*}\right):=\frac{1}{p}\left\langle\Theta, \Sigma^{*}\right\rangle-\frac{1}{p} \log \operatorname{det} \Theta \Sigma^{*}-1 \tag{3.4}
\end{equation*}
$$

which, up to scaling by $p$, is the Kullback-Leibler divergence between multivariate mean zero normal distributions with precision matrices $\Theta$ and $\Theta^{*}$ respectively. The Stein loss
has a long history of application in covariance matrix estimation (Dey \& Srinivasan, 1985; Donoho et al., 2018; James \& Stein, 1961; Ledoit \& Wolf, 2018; Stein, 1975, 1986). In this chapter, we work with the symmetrized Stein loss (alternatively known as the divergence loss), defined as

$$
\begin{equation*}
L^{\mathrm{ssym}}\left(\Theta, \Theta^{*}\right):=\frac{L^{\mathrm{s}}\left(\Theta, \Theta^{*}\right)+L^{\mathrm{s}}\left(\Theta^{*}, \Theta\right)}{2}=\frac{1}{2 p}\left\langle\Theta-\Theta^{*}, \Sigma^{*}-\Sigma\right\rangle \tag{3.5}
\end{equation*}
$$

where $\Sigma=\Theta^{-1}$. Note that $L^{\text {ssym }}\left(\Theta, \Theta^{*}\right)$ is symmetric and $2 L^{\text {ssym }}\left(\Theta, \Theta^{*}\right)$ clearly dominates both the Stein loss and the reversed Stein loss $L^{s}\left(\Theta^{*}, \Theta\right)$ (which is also known as the entropy loss). Properties of $L^{\text {ssym }}$ are further discussed in Section 3.2.

We use the $1 / p$ scaling in the loss function (3.5) because, as explained by Ledoit and Wolf (2018), this is necessary for consistency in the high-dimensional regime where the number of variables $p$ may be much larger than the sample size $n$. Indeed, in the simple case where $\Theta^{*}$ is known to be diagonal, the natural estimator is the diagonal matrix $\hat{\Theta}^{\text {DIAG }}$ with diagonal entries $1 / S_{j j}, j=1, \ldots, p$ (where $S$ is the sample covariance matrix defined in (3.2)). It is easy to see that $\left\langle\hat{\Theta}^{\text {DIAG }}-\Theta^{*}, \Sigma^{*}-\hat{\Sigma}^{\text {DIAG }}\right\rangle$ is of the order $p / n$ which will be far from zero in the high-dimensional regime where $p>n$.

We present results on the performance of $\widehat{\Theta}$ in the symmetrized Stein loss in Section 3.2. Our main result in Theorem 3.1 implies that $L^{\text {ssym }}\left(\widehat{\Theta}, \Theta^{*}\right)$ converges to zero as long as $\log p=o(n)$. This implies high-dimensional consistency of $\widehat{\Theta}$. Moreover, the rate of convergence is $\sqrt{\frac{\log p}{n}}$, which we prove in Theorem 3.2 is optimal in the minimax sense. Thus $\widehat{\Theta}$ is minimax optimal in the high-dimensional regime under the symmetrized Stein loss. Our results provide rigorous support for the assertion that the nonpositive off-diagonal constraint provides strong implicit regularization in the high-dimensional regime. In Theorem 3.4, we also lower bound the loss $L^{\text {ssym }}\left(\widehat{\Theta}, \Theta^{*}\right)$ which implies that the $\sqrt{n}$ rate is not an artifact of our analysis even when the true precision matrix $\Theta^{*}$ is diagonal.

High-dimensional consistency with the rate $\sqrt{\frac{\log p}{n}}$ has appeared previously in many papers on covariance and precision matrix estimation - see for instance Cai et al. (2011), Ravikumar et al. (2011), Rothman et al. (2008), Sun and Zhang (2013), and Yuan (2010) and Cai et al. (2016b) for a review of rates in structured covariance estimation. Most of these results are for estimators that use explicit regularizers, such as the $\ell_{1}$ penalty in the Graphical Lasso (Banerjee et al., 2008; Friedman et al., 2008; Mazumder \& Hastie, 2012), which is crucially exploited by the proof techniques and assumptions employed in these papers. By contrast, the regularization induced by the assumption $\Theta^{*} \in \mathcal{M}^{p \times p}$ is implicit and we consequently use different arguments relying on careful use of the KKT conditions underlying the optimization (3.3). Our analysis identifies a bound relating the entries of an $M$-matrix to its spectrum, providing new insight into the simplifying structure of the convex cone $\mathcal{M}^{p \times p}$.

The symmetrized Stein loss has the additional symmetry property of invariance under inversion: $L^{\text {ssym }}\left(\widehat{\Sigma}, \Sigma^{*}\right)=L^{\text {ssym }}\left(\widehat{\Theta}, \Theta^{*}\right)$ where $\widehat{\Sigma}:=\widehat{\Theta}^{-1}$. This means that $\hat{\Sigma}$ is also a high-
dimensionally-consistent estimator of $\Sigma^{*}$. The choice of the loss function is quite crucial here. In Section 3.3, using the Perron-Frobenius theorem and a careful analysis of the entrywise positive part $S_{+}$of the sample covariance, we prove a negative result which shows that, for the maximum eigenvalue, $\widehat{\Sigma}$ can be much worse as an estimator of $\Sigma^{*}$ compared to the sample covariance matrix $S$. This result indicates that enforcing the sign-constraints can exacerbate bias in the estimation of the top eigenvalue.

The chapter is organized as follows: Section 3.2 contains our main results establishing optimality of $\widehat{\Theta}$, Section 3.3 establishes suboptimality under the spectral norm, and Section 3.4 has a discussion which touches upon some related issues including misspecification (where $\Theta^{*} \notin \mathcal{M}^{p \times p}$ ), estimation of correlation matrices and connections to shape-restricted regression. Finally Section 3.5 contains proofs of all the results of the chapter.

### 3.2 Symmetrized Stein loss: consistency and optimality

This section contains our results on the high-dimensional consistency and optimality of $\widehat{\Theta}$ under the symmetrized Stein loss $L^{\text {ssym }}$ defined in (3.5). We start by describing some basic properties of $L^{\text {ssym }}$.

The expected value of the objective in (3.3), $\left\langle\Theta, \Sigma^{*}\right\rangle-\log \operatorname{det} \Theta$, agrees up to factors depending only on $\Sigma^{*}$ with the Stein loss (3.4), which is also a matrix Bregman divergence (Dhillon \& Tropp, 2008), proportional to the Kullback-Leibler (KL) divergence between centered multivariate Gaussian distributions: $\frac{2}{p} D\left(\mathcal{N}(0, \Sigma) \| \mathcal{N}\left(0, \Sigma^{*}\right)\right)$. It is well known that the KL divergence is not symmetric. When the inputs to the divergence are reversed, the resulting Bregman divergence is also known as the entropy loss, $L^{\text {ent }}\left(\Theta, \Theta^{*}\right):=L^{\mathrm{s}}\left(\Theta^{*}, \Theta\right)$. The sum of these loss functions dominates each, and conveniently does not directly involve any determinants. Following Ledoit and Wolf (2018), we define $L^{\text {ssym }}=\frac{L^{s}+L^{\text {ent }}}{2}$ to be the average of the two loss functions. Commonly known as the symmetrized Stein loss or divergence loss, $L^{\text {ssym }}$ is equal to the Jeffreys (1946) divergence between two centered multivariate Gaussian distributions, divided by $p$. Definition (3.5) entails a number of useful and important properties for the symmetrized Stein loss:
(i) (Nonnegativity) $L^{\text {ssym }}\left(\Theta, \Theta^{*}\right) \geq 0$, with equality if and only if $\Theta=\Theta^{*}$.
(ii) (Symmetry) $L^{\text {ssym }}\left(\Theta, \Theta^{*}\right)=L^{\text {ssym }}\left(\Theta^{*}, \Theta\right)$.
(iii) (Invariance under inversion) $L^{\text {ssym }}\left(\Theta, \Theta^{*}\right)=L^{\text {ssym }}\left(\Sigma, \Sigma^{*}\right)$.
(iv) (Invariance under congruent transformations) For all $p \times p$ nonsingular matrices $P$, we have the scale-invariance property:

$$
\begin{equation*}
L^{\mathrm{ssym}}\left(\Theta, \Theta^{*}\right)=L^{\mathrm{ssym}}\left(P^{\top} \Theta P, P^{\top} \Theta^{*} P\right) \tag{3.6}
\end{equation*}
$$

The symmetrized Stein loss thus induces a natural geometry on the space of PSD matricessee Moakher and Batchelor (2006) for a review and comparison to other geometries. We emphasize that triangle inequality fails to hold for both $L^{\text {ssym }}$ and $\sqrt{L^{\text {ssym }}}$. As a loss, $L^{\text {ssym }}$ treats the dual problems of estimating the covariance matrix and the precision matrix equally. It can also be shown that the symmetrized Stein loss is equivalent to the squared Frobenius norm when the input matrices $\Theta$ and $\Theta^{*}$ have bounded spectra.

In terms of the eigenvalues $\left(\lambda_{j}\right)_{j=1}^{p}$ of $\Theta \Sigma^{*}$, the symmetrized Stein loss is simply the goodness-of-fit measure

$$
\begin{equation*}
L^{\mathrm{ssym}}\left(\Theta, \Theta^{*}\right)=\frac{1}{p} \sum_{j=1}^{p} \frac{\left(\lambda_{j}-1\right)^{2}}{2 \lambda_{j}} \tag{3.7}
\end{equation*}
$$

This alternative representation provides further insight into the normalization of the loss (3.5) with a factor of $p$. The symmetrized Stein loss is the expectation of the function $\lambda \mapsto \frac{(\lambda-1)^{2}}{2 \lambda}$ with respect to the empirical spectral distribution of $\Theta \Sigma^{*}$. This expectation measures how far the spectrum of $\Theta \Sigma^{*}$ deviates from a point mass at one, which is the spectrum of the identity $I_{p}$. In asymptotic settings where $p=p(n) \rightarrow \infty$ as $n \rightarrow \infty$, a natural consistency criterion checks whether this expectation converges to zero.

Our analysis of the symmetrized Stein loss $L^{\text {ssym }}\left(\widehat{\Theta}, \Theta^{*}\right)$ involves the maximum population correlation between any two variables:

$$
\max _{j \neq k} \frac{\Sigma_{j k}^{*}}{\sqrt{\Sigma_{j j}^{*} \Sigma_{k k}^{*}}}
$$

We assume that the above quantity is strictly less than 1 which is clearly necessary for $\Sigma^{*}$ to be nonsingular i.e., for $\Theta^{*}$ to exist. Our bound on $L^{\text {ssym }}\left(\widehat{\Theta}, \Theta^{*}\right)$ will involve the quantity:

$$
\gamma\left(\Sigma^{*}\right):=\left(1-\max _{j \neq k} \frac{\Sigma_{j k}^{*}}{\sqrt{\Sigma_{j j}^{*} \Sigma_{k k}^{*}}}\right)^{-1}
$$

It is natural for $\gamma\left(\Sigma^{*}\right)$ to enter the analysis in light of the existence result of Slawski and Hein (2015) which states that the maximum sample correlation must be less than one in order for the estimator $\widehat{\Theta}$ to be well-defined. Note that $\gamma\left(\Sigma^{*}\right)$ is the smallest $\gamma \geq 1$ such that

$$
\begin{equation*}
\max _{j \neq k} \frac{\Sigma_{j k}^{*}}{\sqrt{\sum_{j j}^{*} \Sigma_{k k}^{*}}} \leq 1-\gamma^{-1}<1 \tag{3.8}
\end{equation*}
$$

Because $\gamma\left(\Sigma^{*}\right)$ is defined in terms of population correlations, it is scale-invariant. Note that $L^{\text {ssym }}$ also has this scale invariance property (see (3.6)).

Theorem 3.1. Let $S=n^{-1} X^{\top} X$ denote the sample covariance matrix based on data matrix $X \in \mathbb{R}^{n \times p}$ with i.i.d. $\mathcal{N}\left(0, \Sigma^{*}\right)$ rows, where $\Theta^{*}=\left(\Sigma^{*}\right)^{-1} \in \mathcal{M}^{p \times p}$. For all $n \geq c_{1} \gamma^{2}\left(\Sigma^{*}\right) \log p$,
the MLE $\widehat{\Theta}$ defined in (3.3) satisfies

$$
\begin{equation*}
L^{s s y m}\left(\widehat{\Theta}, \Theta^{*}\right) \leq c_{2} \gamma\left(\Sigma^{*}\right) \sqrt{\frac{\log p}{n}} \tag{3.9}
\end{equation*}
$$

with probability at least $1-c_{3} p^{-2}$. Here $c_{1}, c_{2}, c_{3}$ are universal positive constants.
Theorem 3.1 states that $\widehat{\Theta}$ is high-dimensionally consistent in the symmetrized Stein $\operatorname{loss} L^{\text {ssym }}$ as long as $\log p=o(n)$. We prove Theorem 3.1 in Section 3.5, deriving a basic inequality from the first order optimality conditions for (3.3) and showing that concentration of the intrinsic noise $\left\|S-\Sigma^{*}\right\|_{\infty}$ is sufficient to control the basic inequality. Crucially, we use the fact that every $M$-matrix $\Theta \in \mathcal{M}^{p \times p}$ is up to diagonal scaling equivalent to a diagonally dominant matrix (see Berman \& Plemmons, 1994, Chap. 6, Property $M_{34}$ ).

We emphasize that the result holds without additional assumptions on the underlying precision matrix such as sparsity. Consistency in the symmetrized Stein loss is a strong guarantee compared to the recent literature on optimal shrinkage of the sample covariance $S$ under high-dimensional asymptotics (Donoho et al., 2018; Ledoit \& Wolf, 2018), where the symmetrized Stein loss $L^{\text {ssym }}$ converges to a nonzero limit under the asymptotic regime $p / n \rightarrow \alpha>0$ as $n \rightarrow \infty$. By contrast, for the constrained MLE the loss $L^{\text {ssym }}\left(\widehat{\Theta}, \Theta^{*}\right)$ converges in probability to zero whenever $\log p=o(n)$.

Since the upper bound (3.9) depends only on the true precision matrix $\Theta^{*}$ through the population quantity $\gamma\left(\Sigma^{*}\right)$, Theorem 3.1 actually bounds the worst case risk obtained from the divergence loss over all $M$-matrices $\Theta^{*}$ with $\gamma\left(\Sigma^{*}\right)$ bounded. It is natural to question whether the $\sqrt{n}$ rate is improvable. Our next result shows that, in the high-dimensional setting where $p$ grows superlinearly in $n$, the minimax rate over the class of $M$-matrices with $\gamma\left(\Sigma^{*}\right) \leq \gamma$ matches the $\sqrt{\frac{\log p}{n}}$ rate from Theorem 3.1.

Theorem 3.2. Let $X \in \mathbb{R}^{n \times p}$ have i.i.d. $\mathcal{N}\left(0, \Sigma^{*}\right)$ rows, and suppose the number of variables $p$ satisfies $c_{1} n^{\beta} \leq p \leq \exp \left(c_{2} n\right)$. For every $\gamma>1$, we have

$$
\begin{equation*}
\inf _{\substack{\Theta}} \sup _{\substack{ \\(X)}} \mathbb{E} L^{s s y m}\left(\breve{\Theta}, \Theta^{*}\right) \geq c_{\gamma} \sqrt{\frac{\log p}{n}} . \tag{3.10}
\end{equation*}
$$

Here $c_{1}, c_{2}>0$ and $\beta>1$ are universal constants and $c_{\gamma}>0$ is a constant depending only on $\gamma$.

Paired with Theorem 3.1, this result implies that $\widehat{\Theta}$ is minimax optimal in the symmetrized Stein loss over $M$-matrices with correlations bounded away from one. Our proof adapts the construction of Cai et al. (2016a), Theorem 4.1, which lower bounds the minimax risk in the spectral norm over a parameter set of sparse precision matrices of the form $I+\varepsilon A$, where $\varepsilon$ depends on problem parameters $p$ and $n$, and $A$ is an adjacency matrix. A key aspect of this approach is to allow for different perturbations over the rows and columns of $A$, in order to recover the $\sqrt{n}$ rate (Kim, 2020).

The $M$-matrix constraint provides implicit regularization and is crucial for achieving the minimax rate $\sqrt{\frac{\log p}{n}}$. If this constraint is dropped, it is impossible for any estimator to achieve a rate better than $\sqrt{\frac{p}{n}}$ when $p>n$. This follows from the next result where we prove a minimax lower bound of $\sqrt{\frac{p}{n}}$ for the $L^{\text {ssym }}$ loss function over the entire class $\mathcal{S}_{\succeq 0}^{p \times p}$ of positive semidefinite matrices when $p>n$. On the other hand, $\mathcal{M}^{p \times p}$ is much larger than diagonal matrices because the minimax rate of estimation over the class $\mathcal{D}_{+}^{p \times p}$ of positive diagonal matrices in the $L^{\text {ssym }}$ loss function is $1 / n$ (this is also proved in the next result). In summary, the class of $M$-matrices acts as a strong high-dimensional regularizer while being considerably larger than the class of all positive diagonal matrices.

Proposition 3.3. Fix $p$ and $n>2$. The minimax risk in the symmetrized Stein loss over diagonal precision matrices satisfies

$$
\begin{equation*}
\inf _{\hat{\Theta}=\hat{\Theta}(X)} \sup _{\Theta^{*} \in \mathcal{D}_{+}^{p \times p}} \mathbb{E} L^{s s y m}\left(\hat{\Theta}, \Theta^{*}\right) \asymp \frac{1}{n} . \tag{3.11}
\end{equation*}
$$

The minimax risk in the symmetrized Stein loss over PSD matrices satisfies

$$
\begin{equation*}
\inf _{\hat{\Theta}=\hat{\Theta}(S)} \sup _{\Theta \in \in \mathcal{S}=00} \mathbb{E} L^{s s y m}\left(\hat{\Theta}, \Theta^{*}\right) \gtrsim \min \left\{\frac{p}{n}, \sqrt{\frac{p}{n}}\right\} . \tag{3.12}
\end{equation*}
$$

Theorem 3.2 implies that the $\sqrt{n}$ rate of Theorem 3.1 cannot be improved in worst case over the entire class $\mathcal{M}^{p \times p}$. In the next result, we prove that the $\sqrt{n}$ rate for $\widehat{\Theta}$ cannot be improved even when the truth $\Theta^{*}$ lies in the class $\mathcal{D}_{+}^{p \times p}$ of positive diagonal matrices. In other words, this shows that $\widehat{\Theta}$ does not adapt to the minimax rate over $\mathcal{D}_{+}^{p \times p}$.

Theorem 3.4. Suppose $\Theta^{*} \in \mathcal{D}_{+}^{p \times p}$ is a positive diagonal matrix and $c_{1} p \geq \sqrt{n}$. Then

$$
\begin{equation*}
L^{s s y m}\left(\widehat{\Theta}, \Theta^{*}\right) \geq \frac{c_{1}}{2 \sqrt{n}} \tag{3.13}
\end{equation*}
$$

with probability at least $1-3 p \exp \left(-c_{2}(n \wedge p)\right)$, where $c_{1}$ and $c_{2}$ are universal positive constants.

### 3.3 Spectral norm: suboptimality

In this section, we prove a negative result which implies that $\widehat{\Theta}$ and $\widehat{\Sigma}$ can be suboptimal for estimating spectral quantities of $\Theta^{*}$ and $\Sigma^{*}$ respectively. Consider the case when $\Sigma^{*}=I_{p}$ and consider estimation of the top eigenvalue $\lambda_{\max }\left(\Sigma^{*}\right)=1$. The performance of the sample covariance matrix $S$ is well understood. Indeed, in the asymptotic setting $p / n \rightarrow \alpha>0$, Geman (1980) proved that

$$
\lambda_{\max }(S) \rightarrow(1+\sqrt{\alpha})^{2}
$$

in probability as $n \rightarrow \infty$. This implies that $S$ is inconsistent for the estimation of $\lambda_{\max }\left(\Sigma^{*}\right)$ when $p / n$ converges to a positive constant. Our next result proves that $\widehat{\Sigma}=\widehat{\Theta}^{-1}$ is also inconsistent for estimating $\lambda_{\max }\left(\Sigma^{*}\right)$ and, more interestingly, its performance is substantially worse compared to $S$. Specifically, in the same asymptotic setting where $p / n \rightarrow \alpha>0$, we have

$$
\begin{equation*}
\lambda_{\max }(\widehat{\Sigma}) \rightarrow \infty \tag{3.14}
\end{equation*}
$$

in probability as $n \rightarrow \infty$. Thus the introduction of the sign constraints make the resulting covariance matrix estimator $\widehat{\Sigma}$ much worse compared to $S$ for estimating the principal eigenvalue. This should be contrasted with the high-dimensional minimax optimality results from the previous section in the symmetrized Stein loss.

Theorem 3.5. Suppose $\Sigma^{*}=I_{p}$ and $p \geq 17$. Then

$$
\begin{equation*}
\lambda_{\max }(\widehat{\Sigma}) \geq 1+c_{1} \frac{p}{\sqrt{n}} \tag{3.15}
\end{equation*}
$$

with probability at least $1-3 p \exp \left(-c_{2}(n \wedge p)\right)$, for some universal positive constants $c_{1}, c_{2}$.
Note that when $p / n \rightarrow \alpha>0$, the right hand side of (3.15) diverges to $\infty$ which proves (3.14).

The proof of Theorem 3.5 is crucially based on following dual formulation to the constrained MLE (3.3) (see, e.g., Slawski \& Hein, 2015):

$$
\begin{equation*}
\widehat{\Sigma}=\underset{\substack{\Sigma \in \mathcal{S}_{D=0}^{p \times 0} \\ \Sigma \geq S, D_{\Sigma}=D_{S}}}{\operatorname{argmax}} \operatorname{det} \Sigma, \tag{3.16}
\end{equation*}
$$

where the second constraint $\Sigma \geq S$ is an entry-wise inequality. This fact and the wellknown observation that the inverse of an $M$-matrix is entry-wise nonnegative (see Berman \& Plemmons, 1994, Chap. 6, Property $N_{38}$ ) together imply that $\widehat{\Sigma}_{j k} \geq S_{j k} \vee 0$ for all $j, k$. This allows us to prove Theorem 3.5 by a careful analysis of the entry-wise positive part matrix $S_{+}$of $S$.

Theorem 3.5 implies minimax suboptimality of $\widehat{\Sigma}$ in the spectral norm $\|\|\cdot\|\|_{2}$. To see this, note that, for every $K>0$, the sample covariance $S$ satisfies the worst case risk bound

$$
\sup _{\substack{\Sigma^{*} \in \mathcal{S}^{p \times p} \\ \lambda_{\max }\left(\Sigma^{*}\right) \leq K}} \mathbb{E}\left\|S-\Sigma^{*}\right\| \|_{2} \leq C K\left(\sqrt{\frac{p}{n}}+\frac{p}{n}\right),
$$

where $C>0$ is a universal constant (see, e.g., Wainwright, 2019, Example 6.3). By contrast, Theorem 3.5 implies

$$
\sup _{\substack{\Theta^{*} \in \mathcal{M}^{p \times p} \\ \lambda_{\max }\left(\Sigma^{*}\right) \leq K}} \mathbb{E}\left\|\widehat{\Sigma}-\Sigma^{*}\right\|_{2} \geq \mathbb{E}_{\Sigma^{*}=K I_{p}}\left[\lambda_{\max }(\widehat{\Sigma})-K\right] \geq c K \frac{p}{\sqrt{n}}
$$

for $n \gtrsim \log p$. Hence $\widehat{\Sigma}$ is minimax suboptimal in the spectral norm for most choices of $p$ and $n$.

Theorem 3.5 also implies inconsistency in spectral norm for the precision matrix. Since $\lambda_{\max }(\widehat{\Sigma})=\frac{1}{\lambda_{\min }(\widehat{\Theta})}$, we have

$$
\lambda_{\min }(\widehat{\Theta}) \leq \frac{1}{1+c_{1} \alpha \sqrt{n}}
$$

with probability at least $1-3 p \exp \left(-c_{2}(\alpha \wedge 1) n\right)$, where $\alpha=p / n$. As $n \rightarrow \infty$, the upper bound approaches zero: the minimum eigenvalue of $\widehat{\Theta}$ poorly estimates that of $\Theta^{*}$. We record this as a separate corollary.

Corollary 3.6. Suppose $\Sigma^{*}=I_{p}$ and $p=\alpha n \geq 17$. Then

$$
\begin{equation*}
\left\|\widehat{\Theta}-\Theta^{*}\right\|_{2} \geq 1-\lambda_{\min }(\widehat{\Theta}) \geq \frac{1}{1+1 /\left(c_{1} \alpha \sqrt{n}\right)} \tag{3.17}
\end{equation*}
$$

with probability at least $1-3 \alpha n e^{-c_{2} n(\alpha \wedge 1)}$. Hence, $\widehat{\Theta}$ is inconsistent in the spectral norm as $n \rightarrow \infty$ and $p / n \rightarrow \alpha$.

### 3.4 Discussion

In this chapter, we establish the possibility of tuning-free estimation of a large precision matrix $\Theta^{*}$ based only on the knowledge that it is an $M$-matrix i.e., it has nonpositive offdiagonal entries. Our main contribution is to identify a loss-namely, the symmetrized Stein loss-in which $\widehat{\Theta}$ is both high-dimensionally consistent and minimax optimal. As the form (3.7) for the symmetrized Stein loss suggests, the quantity $L^{\text {ssym }}\left(\Theta, \Theta^{*}\right)$ is an average measure of closeness across all of the eigenvalues. The estimator $\widehat{\Theta}$ is inadequate, however, for estimating the extreme eigenvalues when $p$ is large relative to $n$, and our other main result establishes that $\widehat{\Sigma}$ is minimax suboptimal in the spectral norm, even relative to the usual sample covariance matrix $S$. For the remainder of this section, we discuss some aspects that are naturally connected to our main results.

Misspecification. In practice, the assumption that all partial correlations are nonnegative may not hold exactly. Slawski and Hein (2015) empirically evaluate the impact of misspecification on the estimator $\widehat{\Theta}$, defining the attractive part $\Theta^{\bullet} \in \mathcal{M}^{p \times p}$ of the population precision $\Theta^{*} \notin \mathcal{M}^{p \times p}$ as the population analogue of the Bregman projection (3.3) with $S$ replaced by $\Sigma^{*}$. Under the symmetrized Stein loss, a straightforward extension of Theorem 3.1 shows that $\widehat{\Theta}$ targets the attractive part $\Theta^{\bullet}$ even under misspecification.

Theorem 3.7. Let $S=n^{-1} X^{\top} X$ denote the sample covariance based on $X \in \mathbb{R}^{n \times p}$ with i.i.d. $\mathcal{N}\left(0, \Sigma^{*}\right)$ rows. Define the attractive part $\Theta^{\bullet} \in \mathcal{M}^{p \times p}$ of the model as

$$
\Theta^{\bullet}:=\underset{\Theta \in \mathcal{M}^{p \times p}}{\operatorname{argmin}}\left\{\left\langle\Theta, \Sigma^{*}\right\rangle-\log \operatorname{det} \Theta\right\} .
$$

For all $n \geq c_{1} \gamma^{2}\left(\Sigma^{\bullet}\right) \log p$, the MLE $\widehat{\Theta}$ defined in (3.3) satisfies

$$
L^{s s y m}\left(\widehat{\Theta}, \Theta^{\bullet}\right) \leq c_{2} \gamma\left(\Sigma^{\bullet}\right) \sqrt{\frac{\log p}{n}}
$$

with probability at least $1-c_{3} p^{-2}$. Here $c_{1}, c_{2}, c_{3}$ are universal positive constants.
Estimating the correlation matrix. One may also be interested, under the same nonnegative partial correlations assumption, in estimating the population correlation matrix $\Gamma^{*}:=D_{\Sigma^{*}}^{-1 / 2} \Sigma^{*} D_{\Sigma^{*}}^{-1 / 2}$ and its inverse $\Omega^{*}=\left(\Gamma^{*}\right)^{-1}=D_{\Sigma^{*}}^{1 / 2} \Theta^{*} D_{\Sigma^{*}}^{1 / 2}$ (here $D_{\Sigma^{*}}$ denotes the diagonal matrix whose diagonal is equal to that of $\Sigma^{*}$. It is natural to use $\widehat{\Omega}:=D_{S}^{1 / 2} \widehat{\Theta} D_{S}^{1 / 2}$ to estimate $\Omega^{*}$. One can check that $\widehat{\Omega}$ satisfies

$$
\widehat{\Omega}=\underset{\Omega \in \mathcal{M}^{p \times p}}{\operatorname{argmin}}\{\langle\Omega, R\rangle-\log \operatorname{det} \Omega\} .
$$

because the optimization problem is equivariant with respect to diagonal scaling (see Lauritzen et al., 2019, Lemma 2.5). The high-dimensional consistency result of Theorem 3.1 also holds for $\widehat{\Omega}$ as an estimator of the inverse correlation matrix $\Omega^{*}$. This follows from an argument analogous to the proof of Theorem 3.1, with the tail bound for $\left\|S-\Sigma^{*}\right\|_{\infty}$ replaced by the corresponding tail bound on $\left\|R-\Gamma^{*}\right\|_{\infty}$ (see, e.g., Sun \& Zhang, 2013, Lemma 19).

Non-Gaussian observations. We state Theorems 3.1 and 3.7 under the Gaussian assumption for simplicity and to remain consistent with other results in this chapter. In general, the upper bound depends on the tail behavior of $\left\|S-\Sigma^{*}\right\|_{\infty}$-see Lemma 3.8. A similar result holds when the rows of $X$ are i.i.d. with $\sigma$-sub-Gaussian components. As Ravikumar et al. (2011) note, estimators of the form (3.3) are motivated via maximum likelihood yet remain sensible for non-Gaussian $X$. For general $X$, the estimator $\widehat{\Theta}$ is motivated as a Bregman projection of $S$ with respect to the Stein loss.

Modifying $\widehat{\Theta}$. Although we focus on properties of the tuning-free estimator $\widehat{\Theta}$, additional processing such as thresholding $\widehat{\Theta}$ or pre-processing the sample covariance $S$ may produce an estimator that is high-dimensionally consistent in the spectral norm. The tuningfree covariance estimate $\widehat{\Sigma}$ may also prove more useful for spectral analysis when the true covariance is a dense matrix. For instance, in the equicorrelation model where $\Sigma^{*}$ has unit diagonal and every off-diagonal entry equal to $r \in(0,1)$, the entry-wise inequalities in (3.16) may introduce less bias.

Related problems. Karlin and Rinott (1983), who pioneered the connection between $M$-matrices and $\mathrm{MTP}_{2}$, also considered repulsive models where the covariance matrix $\Sigma^{*} \in \mathcal{M}^{p \times p}$ has nonpositive off-diagonal, in which case all marginal and partial correlations are nonpositive. This also defines an interesting model class which may similarly simplify estimation in high-dimensional problems. Note, however, that the constraint set $\left\{\Theta: \Theta^{-1} \in \mathcal{M}^{p \times p}\right\}$ of symmetric inverse- $M$ matrices is nonconvex, presenting potential difficulties for maximum likelihood estimation.

Connection to shape-restricted regression. As a subset of the $p \times p$ symmetric positive-semidefinite matrices, the $M$-matrices $\mathcal{M}^{p \times p}$ form a closed, convex cone determined
only by sign constraints. The sign constraints on the precision matrix are analogous to a shape constraint in shape-restricted regression, enabling the use of likelihood techniques without explicit regularization. In particular, one can define the Bregman projection $\widehat{\Theta}$ of $S$ onto $\mathcal{M}^{p \times p}$ (Lauritzen et al., 2019; Slawski \& Hein, 2015). This work thus represents a first foray into the study of shape constraints for high-dimensional precision matrix estimation, inspired by results on regularization-free prediction in high-dimensional linear models via nonnegative least squares (Slawski \& Hein, 2013).

### 3.5 Proofs

### 3.5.1 Proofs of Theorems 3.1 and 3.7

We first introduce two lemmas needed in the proof of Theorem 3.1. Following previous results on sparse precision matrix estimation (see, e.g., Cai et al., 2011; Ravikumar et al., 2011; Sun \& Zhang, 2013), we rely on concentration of the entry-wise maximum deviation $\left\|S-\Sigma^{*}\right\|_{\infty}=\max _{j, k}\left|S_{j k}-\Sigma_{j k}^{*}\right|$ in the high-dimensional regime. A key technical tool in our analysis is the following lemma, which follows from an application of Bernstein's inequality.

Lemma 3.8. Jankova and Van De Geer, 2015, Lemma 6 Suppose $X \in \mathbb{R}^{n \times p}$ has i.i.d. $\mathcal{N}\left(0, \Sigma^{*}\right)$ rows and let $S=n^{-1} X^{\top} X$. For any $t>2$,

$$
\mathbb{P}\left(\left\|S-\Sigma^{*}\right\|_{\infty} \geq 2\left\|\Sigma^{*}\right\|_{\infty}\left[\sqrt{\frac{2 t \log p}{n}}+\frac{t \log p}{n}\right]\right) \leq \frac{2}{p^{t-2}}
$$

Proof. Let $\alpha=e_{j}$ and $\beta=e_{k}$ denote the standard basis vectors. Lemma 6 of Jankova and Van De Geer (2015) provides

$$
\mathbb{P}\left(\alpha^{\top}\left(S-\Sigma^{*}\right) \beta \geq 2\left\|\Sigma^{*}\right\|_{\infty}\left[\sqrt{\frac{2 x}{n}}+\frac{x}{n}\right]\right) \leq 2 e^{-x}
$$

Taking a union bound over $j \leq k$ and setting $x=\log p^{t}$ yields the claim.
The next lemma records a distinctive property of $M$-matrices, corresponding to the fact that $M$-matrices are generalized diagonally dominant (Plemmons, 1977).

Lemma 3.9. Every $M$-matrix $\Theta \in \mathcal{M}^{p \times p}$ satisfies $\|\Theta\|_{1}:=\sum_{i, j}\left|\Theta_{i j}\right| \leq 2 \operatorname{tr}(\Theta)$.
Proof. Since $\Theta$ is symmetric PSD, there are vectors $\theta_{1}, \ldots, \theta_{p}$ such that $\Theta_{i j}=\left\langle\theta_{i}, \theta_{j}\right\rangle$. Moreover, since $\Theta$ has nonpositive off-diagonal entries, $\left\langle\theta_{i}, \theta_{j}\right\rangle \leq 0$ for $i \neq j$. Hence

$$
\|\Theta\|_{1}=\sum_{i}\left\|\theta_{i}\right\|_{2}^{2}-\sum_{i \neq j}\left\langle\theta_{i}, \theta_{j}\right\rangle=2 \sum_{i}\left\|\theta_{i}\right\|_{2}^{2}-\left\|\sum_{i} \theta_{i}\right\|_{2}^{2} \leq 2 \sum_{i}\left\|\theta_{i}\right\|_{2}^{2}=2 \operatorname{tr}(\Theta) .
$$

An illustrative example is the one-parameter family of $p \times p$ symmetric matrices $A_{x}=$ $(1-x) I_{p}+x 1_{p} 1_{p}^{\prime}$ (where $1_{p}=\sum_{j=1}^{p} e_{j}$ is the all ones vector) with unit diagonal and every off-diagonal equal to $x$. Its eigenvalues are $1-x$ (with multiplicity $p-1$ ) and $1+(p-1) x$. Thus $A_{x}$ is PSD if and only if $x \in\left[-\frac{1}{p-1}, 1\right]$, whereas $A_{x}$ is an $M$-matrix if and only if $x \in\left[-\frac{1}{p-1}, 0\right]$. Finally, note $\left\|A_{x}\right\|_{1}=p+p(p-1)|x|$ and $\operatorname{tr}\left(A_{x}\right)=p$. This example shows Lemma 3.9 is tight. For general PSD matrices, the element-wise $\ell_{1}$-norm can be as large as $p$ times the trace, but for $M$-matrices it can be at most twice as large.

We are now ready to prove Theorem 3.1.
Proof of Theorem 3.1. For any positive diagonal matrix $D \in \mathcal{D}_{+}^{p \times p}$,

$$
\begin{aligned}
L^{\mathrm{ssym}}\left(\widehat{\Theta}(S), \Theta^{*}\right) & =L^{\mathrm{ssym}}\left(D \widehat{\Theta}(S) D, D \Theta^{*} D\right) \\
& =L^{\mathrm{ssym}}\left(\widehat{\Theta}\left(D^{-1} S D^{-1}\right), D^{-1} \Sigma^{*} D^{-1}\right)
\end{aligned}
$$

where the first step uses the fact that $L^{\text {ssym }}$ is invariant under congruent transformations, and the second step uses the scale-invariance of the program (3.3). With a sample covariance $S$ based on Gaussian observations, the loss $L^{\text {ssym }}\left(\widehat{\Theta}, \Theta^{*}\right)$ has the same distribution for covariance matrices of the form $\left\{D^{-1} \Sigma^{*} D^{-1}\right\}_{D \in \mathcal{D}_{+}^{p \times p}}$. In particular, taking $D=D_{\Sigma^{*}}^{1 / 2}$, we may assume without loss of generality that $\Sigma^{*}$ is normalized; i.e., $\Sigma^{*}$ has unit diagonal or equivalently $\Sigma^{*}$ equals the population correlation matrix $\Gamma^{*}$.

Let $f(\Theta)=\langle\Theta, S\rangle-\log |\Theta|$. Since the estimator solves the constrained convex optimization problem $\widehat{\Theta}=\arg \min _{\Theta \in \mathcal{M}^{p \times p}} f(\Theta)$, it is characterized by $\langle\nabla f(\widehat{\Theta}), \Theta-\widehat{\Theta}\rangle \geq 0$, for all $\Theta \in \mathcal{M}^{p \times p}$, where $\nabla f(\Theta)=S-\Theta^{-1}$. Hence

$$
\left\langle S-\widehat{\Sigma}, \Theta^{*}-\widehat{\Theta}\right\rangle \geq 0
$$

Rearranging yields the basic inequality

$$
L^{\text {ssym }}\left(\widehat{\Theta}, \Theta^{*}\right) \leq \frac{1}{2 p}\left\langle S-\Sigma^{*}, \Theta^{*}-\widehat{\Theta}\right\rangle
$$

Let $A:=\left\|S-\Sigma^{*}\right\|_{\infty}$. Using Hölder's inequality, we have:

$$
L^{\text {ssym }}\left(\widehat{\Theta}, \Theta^{*}\right) \leq \frac{A}{2 p}\left\|\Theta^{*}-\widehat{\Theta}\right\|_{1}
$$

Now applying the triangle inequality and Lemma 3.9 to the element-wise $\ell_{1}$-norm,

$$
L^{\mathrm{ssym}}\left(\widehat{\Theta}, \Theta^{*}\right) \leq \frac{A}{p}\left(\operatorname{tr}\left(\Theta^{*}\right)+\operatorname{tr}(\widehat{\Theta})\right) .
$$

Since we have assumed without loss of generality that $\Sigma^{*}=\Gamma^{*}$,

$$
\begin{aligned}
& \operatorname{tr}\left(\widehat{\Theta} \Sigma^{*}\right)=\operatorname{tr}(\widehat{\Theta})+\sum_{j \neq k} \widehat{\Theta}_{j k} \Gamma_{j k}^{*} \geq\left(1-\max _{j \neq k} \Gamma_{j k}^{*}\right) \operatorname{tr}(\widehat{\Theta}) \\
& p=\operatorname{tr}\left(\Theta^{*} \Sigma^{*}\right)=\operatorname{tr}\left(\Theta^{*}\right)+\sum_{j \neq k} \Theta_{j k}^{*} \Gamma_{j k}^{*} \geq\left(1-\max _{j \neq k} \Gamma_{j k}^{*}\right) \operatorname{tr}\left(\Theta^{*}\right),
\end{aligned}
$$

where we have again used Lemma 3.9, along with the facts that $\widehat{\Theta}_{j k}$ and $\Theta_{j k}^{*}$ are nonpositive for $j \neq k$ and $\Sigma^{*} \geq 0$ entry-wise (see Berman \& Plemmons, 1994, Chap. 6, Property $N_{38}$ ). Combining the last three displays and using the characterization of $\gamma\left(\Sigma^{*}\right)$ in (3.8), we get

$$
\begin{aligned}
L^{\text {ssym }}\left(\widehat{\Theta}, \Theta^{*}\right) & \leq \frac{\gamma\left(\Sigma^{*}\right) A}{p}\left(p+\operatorname{tr}\left(\widehat{\Theta} \Sigma^{*}\right)\right) \\
& \leq \gamma\left(\Sigma^{*}\right) A\left(3+2 L^{\mathrm{ssym}}\left(\widehat{\Theta}, \Theta^{*}\right)\right)
\end{aligned}
$$

On the event $E=\left\{2 \gamma\left(\Sigma^{*}\right) A \leq \frac{1}{2}\right\}$, we have $L^{\text {ssym }}\left(\widehat{\Theta}, \Theta^{*}\right) \leq 6 \gamma\left(\Sigma^{*}\right) A$. Applying Lemma 3.8 with $t=4$, the event $E^{\prime}=\left\{A \leq 2 \sqrt{\frac{8 \log p}{n}}+\frac{8 \log p}{n}\right\}$ occurs with probability at least $1-2 / p^{2}$.

To guarantee $E^{\prime} \subset E$, we require

$$
2 \sqrt{\frac{8 \log p}{n}}+\frac{8 \log p}{n} \leq \frac{1}{4 \gamma\left(\Sigma^{*}\right)}
$$

which is equivalent to

$$
\frac{\log p}{n} \leq 2+\frac{1}{4 \gamma\left(\Sigma^{*}\right)}-\sqrt{4+\gamma^{-1}\left(\Sigma^{*}\right)}
$$

Using $\gamma\left(\Sigma^{*}\right) \geq 1$, it is straightforward to check that the right hand side above is at least $\frac{1}{72 \gamma^{2}\left(\Sigma^{*}\right)}$. Hence, as long as $n \geq 72 \gamma^{2}\left(\Sigma^{*}\right) \log p$,

$$
L^{\text {ssym }}\left(\widehat{\Theta}, \Theta^{*}\right) \leq 6 \gamma\left(\Sigma^{*}\right)\left(2 \sqrt{\frac{8 \log p}{n}}+\frac{8 \log p}{n}\right)
$$

with probability at least $1-2 / p^{2}$. Since $\gamma\left(\Sigma^{*}\right) \geq 1$, the $\sqrt{\frac{8 \log p}{n}}$ dominates the $\frac{8 \log p}{n}$ term. In particular, we have $\sqrt{\frac{8 \log p}{n}} \leq \frac{1}{3}$, so $L^{\text {ssym }}\left(\widehat{\Theta}, \Theta^{*}\right) \leq 28 \gamma\left(\Sigma^{*}\right) \sqrt{\frac{2 \log p}{n}}$ with probability at least $1-2 / p^{2}$.

Proof of Theorem 3.7. Since the attractive part $\Theta^{\bullet}$ is an $M$-matrix, from the first order optimality conditions for $\widehat{\Theta}$,

$$
\left\langle S-\widehat{\Sigma}, \Theta^{\bullet}-\widehat{\Theta}\right\rangle \geq 0
$$

Using the first order optimality conditions for $\Theta^{\bullet}$ and the fact that $\widehat{\Theta} \in \mathcal{M}^{p \times p}$,

$$
\left\langle\Sigma^{*}-\Sigma^{\bullet}, \widehat{\Theta}-\Theta^{\bullet}\right\rangle \geq 0
$$

Adding these and rearranging yields the basic inequality

$$
L^{\text {ssym }}\left(\Theta^{\bullet}, \widehat{\Theta}\right) \leq \frac{1}{2 p}\left\langle S-\Sigma^{*}, \Theta^{\bullet}-\widehat{\Theta}\right\rangle
$$

The rest of the proof proceeds as the proof of Theorem 3.1, substituting $\Theta^{*}$ with $\Theta^{\bullet}$.

### 3.5.2 Proof of Theorem 3.2

Proof of Theorem 3.2. As in Cai et al. (2016a, Proof of Theorem 4.1), we consider precision matrices of the form

$$
\Theta=\left[\begin{array}{cc}
I_{\lceil p / 2\rceil} & \varepsilon A  \tag{3.18}\\
\varepsilon A^{\top} & I_{\lfloor p / 2\rfloor}
\end{array}\right],
$$

where $A$ is a sparse binary matrix with $k$ nonzero entries per row and at most $2 k$ nonzero entries per column, for some positive integer $k$ and some $\varepsilon$ to be chosen later. As long as $\varepsilon<0$ and $2 k|\varepsilon|<1$, the matrix $\Theta$ is a diagonally dominant $M$-matrix. Its inverse is given by the Neumann series

$$
\begin{aligned}
\Sigma=\Theta^{-1} & =\sum_{m=0}^{\infty}(-\varepsilon)^{m}\left[\begin{array}{cc}
0 & A \\
A^{\top} & 0
\end{array}\right]^{m} \\
& =\sum_{m=0}^{\infty} \varepsilon^{2 m}\left[\begin{array}{cc}
\left(A A^{\top}\right)^{m} & -\varepsilon A\left(A^{\top} A\right)^{m} \\
-\varepsilon A^{\top}\left(A A^{\top}\right)^{m} & \left(A^{\top} A\right)^{m}
\end{array}\right]
\end{aligned}
$$

From the last display, it is clear that $D_{\Sigma} \geq I_{p}$, so $\max _{j \neq k} \Gamma_{j k} \leq \max _{j \neq k} \Sigma_{j k}$ where $\Gamma$ is the correlation matrix corresponding to $\Sigma$. Furthermore, by triangle-inequality, the largest off-diagonal entry of the top left diagonal block is at most

$$
\left\|\sum_{m=0}^{\infty} \varepsilon^{2 m}\left(A A^{\boldsymbol{\top}}\right)^{m}\right\|_{\infty, \mathrm{off}} \leq \sum_{m=1}^{\infty} \varepsilon^{2 m}\left\|\left(A A^{\boldsymbol{\top}}\right)^{m}\right\|_{\infty, \text { off }}
$$

where we use that the first term $m=0$ has zero off-diagonal. This yields

$$
\left\|\left(A A^{\top}\right)^{m}\right\|_{\infty, \text { off }} \leq\left\|\left(A A^{\top}\right)^{m}\right\|_{2} \leq(2 k)^{2 m}
$$

By similar bounds on the other blocks of $\Sigma$, it can be shown that

$$
\max _{j \neq k} \Gamma_{j k} \leq \max _{j \neq k} \Sigma_{j k} \leq \frac{2 k|\varepsilon|}{1-(2 k \varepsilon)^{2}}
$$

A simple sufficient condition to guarantee $\gamma(\Sigma) \leq \gamma$ is thus $4 k|\varepsilon| \leq\left(1-\gamma^{-1}\right) \wedge \frac{1}{2}$.
By the Gers̆gorin circle theorem, the spectrum of $\Theta$ lies in the range $[0,2]$. Further constraining the supremum in (3.10) to $\lambda_{\max }(\Theta) \leq 2$, by Cai and Zhou (2012, Eq. (54)), we have:

$$
\inf _{\breve{\Theta}} \sup _{\substack{\Theta \in \mathcal{M} p \times p \\ \gamma(\Sigma) \leq \gamma}} \mathbb{E} L^{\operatorname{ssym}}(\breve{\Theta}, \Theta) \geq \frac{1}{4} \inf _{\breve{\Theta}} \sup _{\substack{\Theta \in \mathcal{M} p \times p \\ \gamma(\Sigma) \leq \gamma \\ \lambda_{\max }(\Theta) \leq 2}} \mathbb{E} \frac{\|\breve{\Theta}-\Theta\|_{F}^{2}}{p},
$$

so it suffices to lower bound the minimax rate in the Frobenius norm.
Now let $\mathcal{A}$ denote the set of all $\lceil p / 2\rceil \times\lfloor p / 2\rfloor$ binary matrices with $k$ nonzero entries per row and at most $2 k$ nonzero entries per column, and $\mathcal{B}=\{0,1\}^{[p / 2\rceil}$. Finally, let $e$ denote a vector of ones of length $\lfloor p / 2\rfloor$. Given $A \in \mathcal{A}$ and $b \in \mathcal{B}$, the matrix $(b \otimes e) \circ A$ has the same shape as $A$, where the $j^{\text {th }}$ row is nonzero if and only if $b_{j}=1$. Let

$$
\mathcal{F}=\left\{\Theta_{A, b}=\left[\begin{array}{cc}
I_{\lceil p / 2\rceil} & \varepsilon(b \otimes e) \circ A \\
\varepsilon\left(b^{\top} \otimes e^{\top}\right) \circ A^{\top} & I_{\lfloor p / 2\rfloor}
\end{array}\right]: A \in \mathcal{A}, b \in \mathcal{B}\right\} .
$$

As we have shown, $\mathcal{F} \subset\left\{\Theta \in \mathcal{M}^{p \times p}: \gamma(\Sigma) \leq \gamma, \lambda_{\max }(\Theta) \leq 2\right\}$. By Cai and Zhou, 2012, Lemma 3

$$
\inf _{\Theta} \max _{\Theta \in \mathcal{F}} \mathbb{E} \frac{\|\breve{\Theta}-\Theta\|_{F}^{2}}{p} \geq \frac{1}{32}\left[\min _{\substack{A, A^{\prime} \in \mathcal{A}, b, b^{\prime} \in \mathcal{B} \\ b \neq b^{\prime}}} \frac{\left\|\Theta_{A, b}-\Theta_{A^{\prime}, b^{\prime}}\right\|_{F}^{2}}{H\left(b, b^{\prime}\right)}\right]\left[\min _{1 \leq j \leq\lceil p / 2\rceil}\left\|\bar{P}_{j, 0} \wedge \bar{P}_{j, 1}\right\|\right],
$$

where $H$ denotes the Hamming distance and $\left\|\bar{P}_{j, 0} \wedge \bar{P}_{j, 1}\right\|$ denotes the total variation affinity between the measures $\bar{P}_{j, 0}$ and $\bar{P}_{j, 1}$, where $\bar{P}_{j, i}$ is the uniform mixture over $\mathcal{N}\left(0, \Theta_{A, b}^{-1}\right)$ over all $A \in \mathcal{A}$ and all $b \in \mathcal{B}$ such that $b_{j}=i$.

For the first term, fix $A, A^{\prime}$ and $b \neq b^{\prime}$. For $j$ such that $b_{j} \neq b_{j}^{\prime}$, if say $b_{j}=0$, the $j^{\text {th }}$ row of $(b \otimes e) \circ A$ is zero and the $j^{\text {th }}$ row of $\left(b^{\prime} \otimes e\right) \circ A^{\prime}$ has $k$ nonzero entries. Hence

$$
\min _{\substack{A, A^{\prime} \in \mathcal{A}, b, b^{\prime} \in \mathcal{B} \\ b \neq b^{\prime}}} \frac{\left\|\Theta_{A, b}-\Theta_{A^{\prime}, b^{\prime}}\right\|_{F}^{2}}{H\left(b, b^{\prime}\right)} \geq \min _{\substack{A, A^{\prime} \in \mathcal{A}, b, b^{\prime} \in \mathcal{B} \\ b \neq b^{\prime}}} \frac{2 \sum_{j: b_{j} \neq b_{j}^{\prime}} k \varepsilon^{2}}{H\left(b, b^{\prime}\right)}=2 k \varepsilon^{2}
$$

In particular, we have shown

$$
\inf _{\breve{\Theta}} \sup _{\substack{\Theta \in \mathcal{M} p \times p \\ \gamma(\Sigma) \leq \gamma}} \mathbb{E} L^{\operatorname{ssym}}(\breve{\Theta}, \Theta) \geq c k \varepsilon^{2} \min _{1 \leq j \leq\lceil p / 2\rceil}\left\|\bar{P}_{j, 0} \wedge \bar{P}_{j, 1}\right\| .
$$

Finally, the same argument of (Cai et al., 2016a, proof of Lemma 4.5) with $\varepsilon=c^{\prime} \sqrt{\frac{\log p}{n}}$ can be used to show $\min _{1 \leq j \leq\lceil p / 2\rceil}\left\|\bar{P}_{j, 0} \wedge \bar{P}_{j, 1}\right\| \geq c^{\prime \prime}>0$, yielding

### 3.5.3 Proof of Proposition 3.3

Proof of Proposition 3.3. Let $\tilde{\Sigma}^{\text {DIAG }}=c \cdot D_{S}$. Since $S_{11}=\frac{1}{n} \sum_{i=1}^{n} X_{i 1}^{2}$ for $X_{i 1} \stackrel{\mathrm{iid}}{\sim} \mathcal{N}\left(0, \Sigma_{11}^{*}\right)$,

$$
\mathbb{E}\left[L^{\text {ssym }}\left(\tilde{\Theta}^{\mathrm{DIAG}}, \Theta^{*}\right)\right]=\frac{1}{2} \mathbb{E}\left[\frac{\Sigma_{11}^{*}}{c S_{11}}+\frac{c S_{11}}{\Sigma_{11}^{*}}-2\right]=\frac{1}{2}\left[\frac{1}{c} \frac{n}{n-2}+c-2\right] .
$$

The minimum is achieved at $c=\sqrt{\frac{n}{n-2}}$, but taking $c=1$ suffices to prove the minimax rate (3.11) is upper bounded by $\frac{C}{n}$.

Now consider a prior $G$ on $\mathcal{D}_{+}^{p \times p}$ over which the components $\Theta_{j j}^{*}$ are i.i.d. Lower bound the minimax risk by the Bayes risk with respect to $G$ :

$$
\begin{aligned}
\inf _{\hat{\Sigma}} \sup _{\Sigma^{*} \in \mathcal{D}_{+}^{p \times p}} \mathbb{E} L^{\text {ssym }}\left(\hat{\Sigma}, \Sigma^{*}\right) & \geq \inf _{\hat{\Sigma}} \mathbb{E}_{G} L^{\text {ssym }}\left(\hat{\Sigma}, \Sigma^{*}\right) \\
& =\inf _{\hat{\Sigma}_{11}} \mathbb{E}_{G} L^{\text {ssym }}\left(\hat{\Sigma}_{11}, \Sigma_{11}^{*}\right) .
\end{aligned}
$$

If $G=[\operatorname{Gamma}(a, b)]^{\otimes n}$, such that $\Theta_{j j}^{*} \stackrel{\text { iid }}{\sim} \operatorname{Gamma}(a, b)$ under $G$, then combining with the likelihood we have:

$$
S_{11} \mid \Theta_{11}^{*} \sim \operatorname{Gamma}(a, b)
$$

By conjugacy, the posterior is readily seen to be

$$
\Theta_{11}^{*} \left\lvert\, S_{11}=s \sim \operatorname{Gamma}\left(a+\frac{n}{2}, b+\frac{n s}{2}\right) .\right.
$$

Thus, for $n>2$,

$$
\begin{aligned}
\mathbb{E}_{G}\left[L^{\mathrm{ssym}}\left(\mathfrak{d}, \Sigma_{11}^{*}\right) \mid S_{11}=s\right] & =\frac{1}{2} \mathbb{E}_{G}\left[\left.\mathfrak{d} \Theta_{11}^{*}+\frac{\sum_{11}^{*}}{\mathfrak{d}}-2 \right\rvert\, S_{11}=s\right] \\
& =\frac{\mathfrak{d}}{2} \frac{a+n / 2}{b+n s / 2}+\frac{1}{2 \mathfrak{d}} \frac{b+n s / 2}{a+n / 2-1}-1 .
\end{aligned}
$$

This is minimized at $\mathfrak{d}^{*}=\frac{b+n s / 2}{\sqrt{(a+n / 2)(a+n / 2-1)}}$, giving a Bayes risk of

$$
\mathbb{E}_{G}\left[L^{\mathrm{ssym}}\left(\mathfrak{d}^{*}, \Sigma_{11}^{*}\right)\right]=\sqrt{\frac{a+n / 2}{a+n / 2-1}}-1 .
$$

Letting $a \downarrow 0$, we find

$$
\begin{aligned}
\inf _{\hat{\Sigma}} \sup _{\Sigma^{*} \in \mathcal{D}_{+}^{p \times p}} \mathbb{E} L^{\operatorname{ssym}}\left(\hat{\Sigma}, \Sigma^{*}\right) & \geq \sqrt{1+\frac{2}{n-2}}-1 \\
& =\frac{1}{n-2}+o\left(n^{-1}\right),
\end{aligned}
$$

as $n \rightarrow \infty$. This proves the minimax rate (3.11) on $\mathcal{D}_{+}^{p \times p}$.
To prove the lower bound (3.12) on $\mathcal{S}_{\succeq 0}^{p \times p}$, place an inverse Wishart prior $\Sigma^{*} \sim \mathcal{W}^{-1}\left(\Sigma_{0}, \nu\right)$ on the covariance matrix. By conjugacy,

$$
\Sigma^{*} \mid S \sim \mathcal{W}^{-1}\left(\Sigma^{*}, \Sigma_{0}+n S, \nu+n\right)
$$

As long as $\nu+n>p+1$, the posterior loss can be written in closed form as

$$
\mathbb{E}\left[L^{\mathrm{ssym}}\left(\hat{\Theta}, \Theta^{*}\right) \mid S\right]=\frac{1}{2 p}\left[(\nu+n) \operatorname{tr}\left(\hat{\Sigma}\left(\Sigma_{0}+n S\right)^{-1}\right)+\frac{\operatorname{tr}\left(\hat{\Theta}\left(\Sigma_{0}+n S\right)\right)}{\nu+n-p-1}-2 p\right]
$$

which is minimized at $\hat{\Theta}=\sqrt{(\nu+n)(\nu+n-p-1)}\left(\Sigma_{0}+n S\right)^{-1}$, yielding a Bayes risk of

$$
\mathbb{E}\left[L^{\text {ssym }}\left(\hat{\Theta}, \Theta^{*}\right) \mid S\right]=\sqrt{\frac{\nu+n}{\nu+n-p-1}}-1
$$

independent of $\Sigma_{0}$. Setting $\nu=p+1$,

$$
\inf _{\hat{\Theta}=\hat{\Theta}(S)} \sup _{\Theta^{*} \succeq 0} \mathbb{E} L^{\operatorname{ssym}}\left(\hat{\Theta}, \Theta^{*}\right) \geq \sqrt{1+\frac{p+1}{n}}-1
$$

Finally, use $\sqrt{1+x}-1 \geq(\sqrt{2}-1)(x \wedge \sqrt{x})$ for any $x \geq 0$.

### 3.5.4 Proof of Theorem 3.4

Proof of Theorem 3.4. This proof uses Theorem 3.5 which is proved in the next subsection. Since $\Sigma^{*} \in \mathcal{D}_{+}^{p \times p}$, as in the proof of Theorem 3.1 we have

$$
L^{\text {ssym }}\left(\widehat{\Theta}(S), \Theta^{*}\right)=L^{\text {ssym }}\left(\widehat{\Theta}\left(D_{\Sigma^{*}}^{-1 / 2} S D_{\Sigma^{*}}^{-1 / 2}\right), I_{p}\right)
$$

In particular, due to scale invariance of both the estimator and the loss, the symmetrized Stein loss $L^{\text {ssym }}\left(\widehat{\Theta}, \Theta^{*}\right)$ has the same distribution for all diagonal matrices $\Sigma^{*} \in \mathcal{D}_{+}^{p \times p}$. We thus assume with no loss of generality that $\Sigma^{*}=I_{p}$.

Let $f(t)=t+t^{-1}-2$ for $t>0$. By (3.7) and nonnegativity of the function $f$,

$$
L^{\text {ssym }}\left(\widehat{\Theta}, I_{p}\right)=\frac{1}{p} \sum_{j=1}^{p} f\left(\lambda_{j}(\widehat{\Theta})\right) \geq \frac{f\left(\lambda_{\max }(\widehat{\Theta})\right)}{p}
$$

For $t>1, f^{\prime}(t)>0$, so by Theorem 3.5,

$$
L^{\mathrm{ssym}}\left(\widehat{\Theta}, I_{p}\right) \geq \frac{1}{p} f\left(1+c_{1} \frac{p}{\sqrt{n}}\right)=\frac{c_{1}}{\sqrt{n}}\left[1-\frac{1}{1+c_{1} \frac{p}{\sqrt{n}}}\right]
$$

with probability at least $1-3 p \exp \left(-c_{2}(n \wedge p)\right)$. If $c_{1} p \geq \sqrt{n}$, this implies $L^{\text {ssym }}\left(\widehat{\Theta}, I_{p}\right) \geq \frac{c_{1}}{2 \sqrt{n}}$, completing the proof.

### 3.5.5 Proof of Theorem 3.5

The most technically involved part of the proof is a lower bound on the row sums of the positive part $S_{+}$of the sample covariance matrix, which we include as a separate lemma.

Lemma 3.10. Under the conditions of Theorem 3.5,

$$
\sum_{j=1}^{p}\left(S_{+}\right)_{p j} \geq 1+c_{0} \frac{p}{\sqrt{n}}
$$

with probability at least $1-3 \exp \left(-c_{1}(n \wedge p)\right)$, for some universal positive constants $c_{0}, c_{1}$.
We give the proof of Theorem 3.5 assuming the above lemma and then prove the lemma subsequently.

Proof of Theorem 3.5. Since $\widehat{\Sigma}$ is an inverse $M$-matrix, it is entry-wise nonnegative; i.e., $\widehat{\Sigma} \geq 0$. Combining this with the first constraint $\widehat{\Sigma} \geq S$ in the dual formulation (3.16), we have that $\widehat{\Sigma} \geq S_{+} \geq 0$, where $S_{+}$is the entry-wise positive part of the sample covariance matrix $S$. The Perron-Frobenius theorem Berman and Plemmons, 1994, Corollary 1.5 gives

$$
\lambda_{\max }(\widehat{\Sigma}) \geq \lambda_{\max }\left(S_{+}\right)
$$

Thus, we want to show that $\lambda_{\max }\left(S_{+}\right)$is more severely biased than $\lambda_{\max }(S)$. To this end, we apply another standard result from the spectral theory of nonnegative matrices Berman and Plemmons, 1994, Theorem 2.35:

$$
\lambda_{\max }\left(S_{+}\right) \geq \min _{k} \sum_{j}\left(S_{+}\right)_{j k}
$$

By Lemma 3.10, $\sum_{j}\left(S_{+}\right)_{j k} \geq 1+c_{0} \frac{p}{\sqrt{n}}$ with probability at least $1-3 e^{-c_{1}(n \wedge p)}$ for each fixed $k$, so by a union bound,

$$
\min _{k} \sum_{j}\left(S_{+}\right)_{j k} \geq 1+c_{0} \frac{p}{\sqrt{n}}
$$

with probability at least $1-3 p e^{-c_{1}(n \wedge p)}$. Combining the last three displays gives the desired lower bound on $\lambda_{\text {max }}(\widehat{\Sigma})$.

We now prove the key lemma on the row sums of $S_{+}$.
Proof of Lemma 3.10. For $u>0$, write

$$
\mathbb{P}\left\{\sum_{j=1}^{p}\left(S_{+}\right)_{p j} \leq 1+u\right\} \leq \mathbb{P}\left\{S_{p p} \leq 1-u\right\}+\mathbb{P}\left\{\sum_{j<p}\left(S_{p j}\right)_{+} \leq 2 u\right\}
$$

To bound the first term, note that $n S_{p p} \sim \chi_{n}^{2}$ and the following standard chi-squared lower tail bound (see e.g., Laurent and Massart (2000, inequality (4.4))):

$$
\begin{equation*}
\mathbb{P}\left\{\frac{\chi_{n}^{2}}{n} \leq 1-u\right\} \leq \exp \left(\frac{-n u^{2}}{4}\right) \tag{3.19}
\end{equation*}
$$

gives

$$
\begin{equation*}
\mathbb{P}\left\{S_{p p} \leq 1-u\right\} \leq \exp \left(-\frac{n u^{2}}{4}\right) \tag{3.20}
\end{equation*}
$$

To bound the second term, notice that conditionally on $X_{i p}, i=1, \ldots, n$,

$$
S_{p j}, j=1, \ldots, p-1 \mid X_{i p}, i=1, \ldots, n \stackrel{\text { i.i.d }}{\sim} N\left(0, \frac{1}{n^{2}} \sum_{i=1}^{n} X_{i p}^{2}\right) .
$$

Thus, conditionally on $X_{i p}, i=1, \ldots, n$, we can write $S_{p j}=A Z_{j}$ for $j=1, \ldots, p-1$ where

$$
A^{2}:=\frac{1}{n^{2}} \sum_{i=1}^{n} X_{i p}^{2} \quad \text { and } \quad Z_{1}, \ldots, Z_{p-1} \stackrel{\text { i.i.d }}{\sim} N(0,1) .
$$

We can therefore write (using the notation $\mathbb{P}^{\mid}$for probability conditioned on $X_{i p}, i=1, \ldots, n$ )

$$
\begin{aligned}
\mathbb{P}^{\mid}\left\{\sum_{j<p}\left(S_{p j}\right)_{+} \leq 2 u\right\} & =\mathbb{P}^{\mathrm{\mid}}\left\{\sum_{j<p}\left(Z_{j}\right)_{+} \leq \frac{2 u}{A}\right\} \\
& =\mathbb{P}^{\mathrm{\mid}}\left\{\frac{1}{p-1} \sum_{j<p}\left(\left(Z_{j}\right)_{+}-c\right) \leq \frac{2 u}{(p-1) A}-c\right\}
\end{aligned}
$$

where $c:=\mathbb{E}\left(Z_{1}\right)_{+}=(2 \pi)^{-1 / 2}$ is a universal constant. We now note that

$$
\left(z_{1}, \ldots, z_{p-1}\right) \mapsto \frac{1}{p-1} \sum_{j<p}\left(z_{j}\right)_{+}
$$

is a Lipschitz function with Lipschitz constant $(p-1)^{-1 / 2}$. Thus by the usual concentration inequality for Lipschitz functions of Gaussian random vectors (see, e.g., Wainwright, 2019, Theorem 2.26), we obtain

$$
\mathbb{P}^{\mid}\left\{\frac{1}{p-1} \sum_{j<p}\left(\left(Z_{j}\right)_{+}-c\right) \leq \frac{2 u}{(p-1) A}-c\right\} \leq \exp \left(-\frac{(p-1)}{2}\left(c-\frac{2 u}{(p-1) A}\right)^{2}\right)
$$

assuming that $c>2 u /(A(p-1))$. In particular, for $c>4 u /(A(p-1))$, we get

$$
\mathbb{P}^{\mid}\left\{\frac{1}{p-1} \sum_{j<p}\left(\left(Z_{j}\right)_{+}-c\right) \leq \frac{2 u}{(p-1) A}-c\right\} \leq \exp \left(-\frac{(p-1) c^{2}}{8}\right)
$$

We have thus proved

$$
\mathbb{P}^{\mid}\left\{\sum_{j<p}\left(S_{p j}\right)_{+} \leq 2 u\right\} \leq \exp \left(-\frac{(p-1) c^{2}}{8}\right)+I\{c \leq 4 u /(A(p-1))\}
$$

Taking an expectations on both sides of this expression, we obtain

$$
\mathbb{P}\left\{\sum_{j<p}\left(S_{p j}\right)_{+} \leq 2 u\right\} \leq \exp \left(-\frac{(p-1) c^{2}}{8}\right)+\mathbb{P}\left\{A \leq \frac{4 u}{(p-1) c}\right\}
$$

Note now that $n^{2} A^{2} \sim \chi_{n}^{2}$ and thus

$$
\mathbb{P}\left\{A \leq \frac{4 u}{(p-1) c}\right\}=\mathbb{P}\left\{\frac{\chi_{n}^{2}}{n}-1 \leq \frac{16 u^{2} n}{(p-1)^{2} c^{2}}-1\right\}
$$

We now make the choice $u=\frac{(p-1) c}{4 \sqrt{2} \sqrt{n}}$, which gives (via (3.19))

$$
\mathbb{P}\left\{A \leq \frac{4 u}{(p-1) c}\right\}=\mathbb{P}\left\{\frac{\chi_{n}^{2}}{n}-1 \leq \frac{-1}{2}\right\} \leq \exp \left(-\frac{n}{16}\right)
$$

We have thus proved

$$
\mathbb{P}\left\{\sum_{j<p}\left(S_{p j}\right)_{+} \leq \frac{(p-1) c}{2 \sqrt{2} \sqrt{n}}\right\} \leq \exp \left(-\frac{(p-1) c^{2}}{8}\right)+\exp \left(-\frac{n}{16}\right)
$$

Combining this with (3.20) and using $c=(2 \pi)^{-1 / 2}$, we obtain

$$
\mathbb{P}\left\{\sum_{j=1}^{p}\left(S_{+}\right)_{p j} \leq 1+\frac{(p-1)}{8 \sqrt{\pi n}}\right\} \leq \exp \left(-\frac{(p-1)^{2}}{256 \pi}\right)+\exp \left(-\frac{p-1}{16 \pi}\right)+\exp \left(-\frac{n}{16}\right) .
$$

For $p \geq 17$ the first term is of lower order; i.e., $\exp \left(-\frac{(p-1)^{2}}{256 \pi}\right) \leq \exp \left(-\frac{p-1}{16 \pi}\right)$.

## Chapter 4

## Shrinkage for multivariate, heteroscedastic data

### 4.1 Introduction

Consider a $d$-dimensional $(d \geq 1)$, heteroscedastic normal observation model

$$
\begin{equation*}
X_{i} \mid \theta_{i}^{*} \stackrel{\text { ind }}{\sim} \mathcal{N}\left(\theta_{i}^{*}, \Sigma_{i}\right), \quad \text { with } \quad \theta_{i}^{*} \stackrel{\text { iid }}{\sim} G^{*}, \quad \text { for } i \in\{1, \ldots, n\} \tag{4.1}
\end{equation*}
$$

where $\left(\Sigma_{i}\right)_{i=1}^{n}$ is a known sequence of $d \times d$ positive-definite covariance matrices, and the underlying mean vectors $\left(\theta_{i}^{*}\right)_{i=1}^{n}$ are additionally assumed to be drawn from a common prior $G^{*}$, where $G^{*}$ belongs to the collection $\mathcal{P}\left(\mathbb{R}^{d}\right)$ of all probability measures on $\mathbb{R}^{d}$. In settings where $G^{*}$ is known, model (4.1) fully specifies a Bayesian model; this chapter studies the common empirical Bayes setting where $G^{*}$ must be estimated. The main goal of the chapter is to nonparametrically estimate $G^{*}$ and the sequence $\left(\theta_{i}^{*}\right)_{i=1}^{n}$ from the observed data $\left(X_{i}, \Sigma_{i}\right)_{i=1}^{n}$.

Empirical Bayes methods for the normal sequence model (4.1) have been studied extensively in the univariate, homoscedastic setting where $d=1$ and $\Sigma_{i} \equiv \sigma^{2}$ (see, e.g., Efron (2012, 2014), Efron and Morris (1972a, 1972b, 1973a, 1973b), James and Stein (1961), and Morris (1983) as well as Johnstone (2019) for a manuscript on estimation in Gaussian sequence models). Numerous methods extend empirical Bayes to the univariate, heteroscedastic case (see Banerjee et al., 2021; Jiang et al., 2011; Jiang, 2020; Tan, 2016; Weinstein et al., 2018; Xie et al., 2012, and references therein). Relatively little attention has been given to the general case of the present chapter.

Model (4.1) naturally arises in the analysis of astronomy data, where often a calibrated measurement error distribution comes attached to each observation, and typically these errors are heteroscedastic (Kelly, 2012); also see e.g. Akritas and Bershady (1996), Hogg et al. (2010), Anderson et al. (2018). The first part of model (4.1) indicates that the target sequence $\left(\theta_{i}^{*}\right)_{i=1}^{n}$ has, due to measurement error, been corrupted by additive, zero-mean

Gaussian noise, i.e.

$$
X_{i}=\theta_{i}^{*}+\epsilon_{i}, \quad \text { where } \quad \epsilon_{i} \stackrel{\text { ind }}{\sim} \mathcal{N}\left(0, \Sigma_{i}\right), \quad \text { for } i=1, \ldots, n .
$$

Interestingly, the $\Sigma_{i}$ 's above, which typically differ across $i$, are known in many applications where the measurement process is well-characterized. In many situations it is assumed that $\theta_{i}^{*}$ is itself random and independent of $\epsilon_{i}$ for all $i$. Although each observation has a different error distribution, the $n$ observations are tied together by the assumption that the $\theta_{i}^{*}$ 's are i.i.d. from some distribution $G^{*}$, yielding model (4.1). By allowing for arbitrary prior distributions $G^{*} \in \mathcal{P}\left(\mathbb{R}^{d}\right)$, model (4.1) captures a range of important structural assumptions on the underlying sequence $\left(\theta_{i}^{*}\right)_{i=1}^{n}$ : for instance, the clustering problem (where the terms of $\left(\theta_{i}^{*}\right)_{i=1}^{n}$ take on at most $k^{*}$ distinct values) corresponds to discrete $G^{*}$, and sparse modeling (where most of the $\left(\theta_{i}^{*}\right)_{i=1}^{n}$ are zero) corresponds to $G^{*}(\{0\}) \approx 1$. The model also accommodates more complex manifold-like structures (see e.g. Figure 4.1) as well as substantially more heterogeneous sequences (e.g. $G^{*}$ heavy tailed).

Our motivating example for model (4.1) involves the construction of a precise stellar colormagnitude diagram. A color-magnitude diagram (CMD) is a scatter plot of stars, displaying their absolute magnitude (luminosity) versus color (surface temperature) to provide a crosssectional view of stellar evolution. The continued expansion of available stellar measurements has made purely statistical models such as model (4.1) increasingly attractive for denoising. One common approach, known as Extreme Deconvolution (XD) (Bovy et al., 2011), assumes

$$
\begin{equation*}
G^{*}=\sum_{j=1}^{K} \alpha_{j}^{*} \mathcal{N}\left(\mu_{j}^{*}, V_{j}^{*}\right) \tag{4.2}
\end{equation*}
$$

and estimates the parameters $\left(\alpha_{j}^{*}, \mu_{j}^{*}, V_{j}^{*}\right)_{j=1}^{K}$ via the Expectation-Maximization (EM) algorithm with split-and-merge operations designed to avoid local optima. For instance, Anderson et al. (2018) applied XD to build a low-noise CMD with $n \approx 1.4$ million de-reddened stars from the Gaia TGAS catalogue. The XD assumption (4.2) that the prior $G^{*}$ is itself a mixture of $K$-Gaussians has a number of drawbacks. Although the class of Gaussian location-scale mixtures is flexible for large $K$, the choice of $K$ requires tuning; violations of assumption (4.2) for fixed $K$ induce bias in the estimation. To our knowledge, no theoretical results for the statistical properties of XD are available, making it difficult to quantify the misspecification error. Moreover, the class of all probability distributions of the form (4.2) is nonconvex for finite $K$, so even split-and-merge techniques employed within EM do not guarantee convergence to the global maximizer of the likelihood.

To avoid these difficulties, we extend the Kiefer and Wolfowitz (1956) nonparametric maximum likelihood estimator (NPMLE) to incorporate multivariate and heteroscedastic errors. An NPMLE is any $\widehat{G}_{n} \in \mathcal{P}\left(\mathbb{R}^{d}\right)$ which maximizes the marginal likelihood of the observations $\left(X_{i}\right)_{i=1}^{n}$. Marginally, the observations are independent, and the $i^{\text {th }}$ observation $X_{i}$ is distributed according to a Gaussian location mixture with density

$$
\begin{equation*}
f_{G^{*}, \Sigma_{i}}(x):=\int \varphi_{\Sigma_{i}}(x-\theta) \mathrm{d} G^{*}(\theta), \quad \text { for } x \in \mathbb{R}^{d} \tag{4.3}
\end{equation*}
$$

where $\varphi_{\Sigma_{i}}(x):=\frac{1}{\sqrt{\operatorname{det}\left(2 \pi \Sigma_{i}\right)}} \exp \left(-\frac{1}{2} x^{\top} \Sigma_{i}^{-1} x\right)$ denotes the density of $\mathcal{N}\left(0, \Sigma_{i}\right)$. Hence an NPMLE is any maximizer

$$
\begin{equation*}
\widehat{G}_{n} \in \underset{G \in \mathcal{P}\left(\mathbb{R}^{d}\right)}{\operatorname{argmax}} \frac{1}{n} \sum_{i=1}^{n} \log f_{G, \Sigma_{i}}\left(X_{i}\right) \tag{4.4}
\end{equation*}
$$

In contrast to the parametric model used in XD, the nonparametric domain $\mathcal{P}\left(\mathbb{R}^{d}\right)$ is convex, so $\widehat{G}_{n}$ solves a convex optimization problem, and tools from convex optimization may be leveraged to find principled approximations to $\widehat{G}_{n}$ (Kim et al., 2020; Koenker \& Mizera, 2014).

Given an estimate $\widehat{G}_{n}$ of the prior $G^{*}$, empirical Bayes imitates the optimal Bayes analysis, known as the oracle (Efron, 2019). If $G^{*}$ were known, optimal denoising of $\theta_{i}^{*}$ would be achieved through the posterior distribution $\theta_{i}^{*} \mid X_{i}$. It is well known, for instance, that the oracle posterior mean

$$
\begin{equation*}
\hat{\theta}_{i}^{*}:=\mathbb{E}_{G^{*}}\left[\theta_{i}^{*} \mid X_{i}\right], \text { where } \theta_{i}^{*} \sim G^{*} \text { and } X_{i} \mid \theta_{i}^{*} \sim \mathcal{N}\left(\theta_{i}^{*}, \Sigma_{i}\right) \tag{4.5}
\end{equation*}
$$

minimizes the squared error Bayes risk

$$
\mathbb{E}_{G^{*}}\left\|\mathfrak{d}_{i}\left(X_{i}\right)-\theta_{i}^{*}\right\|_{2}^{2}
$$

over all measurable functions $\mathfrak{d}_{i}: \mathbb{R}^{d} \rightarrow \mathbb{R}^{d}$. The NPMLE (4.4) yields a fully data-driven, empirical Bayes estimate of the oracle posterior mean via

$$
\begin{equation*}
\hat{\theta}_{i}:=\mathbb{E}_{\widehat{G}_{n}}\left[\theta_{i}^{*} \mid X_{i}\right], \text { where } \theta_{i}^{*} \sim \widehat{G}_{n} \text { and } X_{i} \mid \theta_{i}^{*} \sim \mathcal{N}\left(\theta_{i}^{*}, \Sigma_{i}\right) \tag{4.6}
\end{equation*}
$$

Figure 4.1 shows the $d=2$ dimensional dataset of Anderson et al. (2018), where each observation has a known error distribution and may be modeled as multivariate normal after a suitable transformation. The noise in the raw CMD of Figure 4.1 obscures many known features of stellar evolution, rendering the raw CMD unreliable for downstream parallax inference. The right panel of Figure 4.1 displays the empirical Bayes posterior means $\left(\hat{\theta}_{i}\right)_{i=1}^{n}$ based on the NPMLE. The substantial shrinkage of our method reveals many recognizable features of the CMD, such as the red clump and a narrow red giant branch in the upper-right region of the plot, as well as the binary sequence tail distinct from the main sequence tail in the bottom-center region. The NPMLE $\widehat{G}_{n}$ and corresponding posterior means $\left(\hat{\theta}_{i}\right)_{i=1}^{n}$ offer a powerful approach to shrinkage estimation under minimal assumptions.

The idea of using the NPMLE to estimate a prior distribution, due to Robbins (1950), has seen a resurgence in recent years (Deb et al., 2021; Dicker \& Zhao, 2016; Efron, 2019; Feng \& Dicker, 2018; Gu \& Koenker, 2016; Jiang, 2020; Jiang \& Zhang, 2009, 2010; Kim et al., 2020; Koenker \& Gu, 2017; Koenker \& Mizera, 2014; Polyanskiy \& Wu, 2020; Saha \& Guntuboyina, 2020a). These advancements, taken together, have begun to establish the NPMLE as a formidable approach to shrinkage estimation both in theory and in practice. All this prior work has focused on either the univariate setting $d=1$ or the homoscedastic


Figure 4.1: A noisy color-magnitude diagram (CMD) corresponding to the observations $X_{i}$ in model (4.1), with corresponding fully-nonparametric denoised estimates $\hat{\theta}_{i}$ in the right panel. To avoid overplotting, we display a subsample of $n=10^{5}$ stars.
setting $\Sigma_{i} \equiv \Sigma$, however. Our work extends the NPMLE to the practically important and more general setting of multivariate and heteroscedastic errors, uncovering a number of important differences.

Basic properties of the NPMLE that are well-understood in the univariate, homoscedastic setting (Lindsay, 1995) have not received careful attention in more complex settings. We verify in Lemma 4.1 that a solution $\widehat{G}_{n}$ exists for every instance of the optimization problem (4.4), and we record the first-order optimality conditions characterizing the solution set. Similar to the univariate, homoscedastic setting, there exists a solution $\widehat{G}_{n}$ which is discrete with at most $n$ atoms, and the sequence of fitted values $\hat{L} \equiv\left(\hat{L}_{1}, \ldots, \hat{L}_{n}\right)=\left(f_{\widehat{G}_{n}, \Sigma_{i}}\left(X_{i}\right)\right)_{i=1}^{n}$ is unique, i.e. every solution $\widehat{G}_{n}$ has the same sequence of fitted likelihood values $\hat{L}$.

An important contribution of Lemma 4.1 is our reinterpretation of the characterizing system of inequalities in terms of a natural 'dual' mixture density $\widehat{\psi}_{n}$. Specifically, $\widehat{\psi}_{n}$ is a heteroscedastic, $n$-component mixture density-a convex combination of Gaussian bumps centered at the datapoints $\mathcal{N}\left(X_{i}, \Sigma_{i}\right)$ with weights inversely proportional to $\hat{L}_{i}$ for $i=1, \ldots, n-$ such that the support of every NPMLE $\widehat{G}_{n}$ is contained in the set of the global maximizers of $\widehat{\psi}_{n}$. This observation has a number of important consequences that we explore in detail in Section 4.2; in particular, tools from algebraic statistics for studying the modes of Gaussian mixtures (Améndola et al., 2020; Ray \& Lindsay, 2005) translate directly into results on
the support set. We leverage this connection to establish that $\widehat{G}_{n}$ is not necessarily unique when $d>1$, even in the homoscedastic case. This finding is distinctive from the univariate, homoscedastic case where it is known that (4.4) has a unique solution for every problem instance (Lindsay \& Roeder, 1993). Our counterexample in Lemma 4.2 appears to be new and seems to invalidate prior claims of strict concavity of the log-likelihood (Koenker \& Gu, 2017; Marriott, 2002). Whereas the fitted values $\hat{L}$ are always unique, our counterexample also demonstrates that the empirical Bayes posterior means $\left(\hat{\theta}_{i}\right)_{i=1}^{n}$ are not necessarily unique. In light of the non-uniqueness of $\widehat{G}_{n}$, a natural question is whether there exist non-discrete solutions: we rule out this possibility in Corollary 4.3, however, showing every solution is indeed discrete with a finite number of atoms.

The problem of computing a solution $\widehat{G}_{n}$ is complicated by the presence of multivariate, heteroscedastic errors. The main difficulty in general is that the NPMLE solves an infinitedimensional optimization problem. Since $\widehat{G}_{n}$ may be taken to be discrete with at most $n$ atoms, a solution can in principle be found with a finite mixture model. In particular, defining the set of discrete distributions with at most $k \geq 1$ atoms,

$$
\mathcal{P}_{k}\left(\mathbb{R}^{d}\right)=\left\{\sum_{j=1}^{k} w_{j} \delta_{a_{j}}: \sum_{j} w_{j}=1, w \geq 0, a_{j} \in \mathbb{R}^{d}, j=1, \ldots, k\right\}
$$

maximum likelihood solutions over $\mathcal{P}_{n}\left(\mathbb{R}^{d}\right)$ are also NPMLEs. Hence, the EM algorithm can be applied to optimize $\left(w_{j}, a_{j}\right)_{j=1}^{n}$, as first observed by Laird (1978), though EM over discrete distributions is prohibitively slow for moderately large $n$ and suffers from the same nonconvexity issue as XD. Many algorithms (Böhning, 1985; Lesperance \& Kalbfleisch, 1992; Liu \& Zhu, 2007; Wang, 2007) have been proposed for finding approximate solutions to the optimization problem (4.4); Koenker and Mizera (2014) identified a convex, finite-dimensional, highly scalable approximation. Instead of maximizing the log-likelihood of the data $\frac{1}{n} \sum_{i=1}^{n} \log f_{G, \Sigma_{i}}\left(X_{i}\right)$ over $G \in \mathcal{P}_{n}\left(\mathbb{R}^{d}\right)$, the idea is to maximize the log-likelihood over $\mathcal{P}(\mathcal{A})$, the collection of all probability measures supported on a finite set $\mathcal{A} \subset \mathbb{R}^{d}$. If $\mathcal{A}$ has $m>0$ elements, then $\mathcal{P}(\mathcal{A})$ is isometric to the $m-1$ dimensional simplex $\Delta_{m-1}:=\left\{w \in \mathbb{R}_{+}^{m}: \sum_{j} w_{j}=1\right\}$, and maximizing the likelihood corresponds to optimizing over the mixing proportions $w$, which is a convex optimization problem. When $d=1$, it is straightforward to see that $\widehat{G}_{n}$ is supported on the range of the data $\left[X_{(1)}, X_{(n)}\right.$ ], so Koenker and Mizera (2014) proposed taking $\mathcal{A}$ to discretize this range. Jiang and Zhang (2009, Proposition 5) bounded the discretization error in $d=1$ dimension, establishing that optimizing the weights $w$ via EM can lead to a good approximation once $m \asymp(\log n) \sqrt{n}$. Dicker and Zhao (2016) further justified the discretization scheme in $d=1$ dimension by showing the discretized NPMLE is statistically indistinguishable from $\widehat{G}_{n}$ once the analyst uses at least $m=\lfloor\sqrt{n}\rfloor$ atoms.

The discretization approach naturally extends to multivariate, heteroscedastic settings, but to our knowledge, no principled recommendations are available for choosing $\mathcal{A} \subset \mathbb{R}^{d}$ in general. Feng and Dicker (2018) recommended taking $\mathcal{A}$ to be a grid over a compact region containing the data. We address the key questions of how to choose this compact region and
how the discretization error depends on the fineness of the grid. For choosing a compact region to discretize, a natural desideratum is that the region should contain the support of $\widehat{G}_{n}$. To this end, in Corollary 4.3 we present compact support bounds on the NPMLE in terms of the data $\left(X_{i}, \Sigma_{i}\right)_{i=1}^{n}$. When $d=1$ our support bounds reduce to the range of the data, reaffirming the original suggestion of Koenker and Mizera (2014), and when $d>1$ but the errors are homoscedastic, it suffices to discretize the convex hull of $\left(X_{i}\right)_{i=1}^{n}$. Interestingly, with multivariate and heteroscedastic errors, the support of the NPMLE can lie outside the convex hull of $\left(X_{i}\right)_{i=1}^{n}$, so a different region known as the ridgeline manifold $\mathcal{M}$ of $\left(X_{i}, \Sigma_{i}\right)_{i=1}^{n}$ is needed. Fortunately, this region $\mathcal{M} \subset \mathbb{R}^{d}$ is compact, and the NPMLE over $\mathcal{P}(\mathcal{M})$ agrees with the NPMLE over $\mathcal{P}\left(\mathbb{R}^{d}\right)$. This justifies the choice of $\mathcal{A}$ as a $\delta>0$ cover of $\mathcal{M}$, and in Proposition 4.5, we verify that as $\delta \downarrow 0$, the log-likelihood of the discretized NPMLE approaches that of the NPMLE. We prove a quantitative bound on the gap for fixed $\delta$, providing some guidance on how the discretization error depends on the fineness of the grid.

Our principled and efficient method of computation facilitates simulation studies assessing the performance of the empirical Bayes estimate $\hat{\theta}_{i}$ in a setting where we can actually compare to the oracle Bayes estimate $\hat{\theta}_{i}^{*}$. Figure 4.2 illustrates the method on simulated data. The means $\theta_{i}^{*}$ were drawn i.i.d. from a circle of radius two, and the data $X_{i} \mid \theta_{i}^{*}$ were drawn according to (4.1) using a variety of diagonal covariance matrices $\Sigma_{i}=\left[\begin{array}{cc}\sigma_{1, i}^{2} & 0 \\ 0 & \sigma_{2, i}^{2}\end{array}\right]$, taking each $\sigma_{j, i}^{2} \in(1 / 2,3 / 4)$. Visually, it is clear that the empirical Bayes estimates improve upon the observations by shrinking towards the underlying circle; the corresponding mean squared errors were $\frac{1}{n} \sum_{i=1}^{n}\left\|\hat{\theta}_{i}-\theta_{i}^{*}\right\|_{2}^{2}=0.87$ and $\frac{1}{n} \sum_{i=1}^{n}\left\|X_{i}-\theta_{i}^{*}\right\|_{2}^{2}=1.46$, respectively. The oracle, which minimizes the mean squared error in expectation, attained an error of $\frac{1}{n} \sum_{i=1}^{n}\left\|\hat{\theta}_{i}^{*}-\theta_{i}^{*}\right\|_{2}^{2}=0.84$. While the oracle cannot be computed in practice because $G^{*}$ is unknown, this value sets a benchmark in simulations to which we may compare the performance of bona fide estimators. The empirical Bayes estimates not only track well with this benchmark; the individual estimates also track remarkably well with the oracle. In our simulation, the regret - defined as the mean squared error between the estimator $\left(\hat{\theta}_{i}\right)_{i=1}^{n}$ and the oracle $\left(\hat{\theta}_{i}^{*}\right)_{i=1}^{n}$ —was $\frac{1}{n} \sum_{i=1}^{n}\left\|\hat{\theta}_{i}-\hat{\theta}_{i}^{*}\right\|_{2}^{2}=0.03$. Whereas $\hat{\theta}_{i}$ is a function of the observed data, the oracle $\hat{\theta}_{i}^{*}$ makes optimal use of the unknown prior $G^{*}$, making the similarity between the two especially striking.

This striking similarity between $\hat{\theta}_{i}$ and $\hat{\theta}_{i}^{*}$ affirms the empirical Bayes adage that "large data sets of parallel situations carry within them their own Bayesian information" (Efron \& Hastie, 2016). However, the setting of Figure 4.2 is complicated by the fact the situations are not directly parallel, in that each observation $X_{i}$ has a distinct error distribution. Even in heteroscedastic settings, the extent to which we glean prior information for the purpose of denoising is captured by the empirical Bayes regret $\frac{1}{n} \sum_{i=1}^{n}\left\|\hat{\theta}_{i}-\hat{\theta}_{i}^{*}\right\|_{2}^{2}$. Theorem 4.8 develops a detailed profile of the finite-sample regret properties of the NPMLE for denoising. We show that under certain tail conditions on $G^{*}$ the regret is bounded by a rate that is nearly parametric in $n$, i.e. $\frac{1}{n}$ up to logarithmic multiplicative factors. The regret still converges at a slower, nonparametric rate under less structured conditions, where $G^{*}$ may have heavy tails.


Figure 4.2: Toy data of size $n=1,000$ and $d=2$. Top: observations $X_{i}$ (left) were generated by adding heteroscedastic Gaussian errors to the underlying means $\theta_{i}^{*} \stackrel{\text { iid }}{\sim} G^{*}$ (right), generated IID uniformly from a circle of radius 2 . Our discrete estimate $\widehat{G}_{n}$ of the prior is shown in red over the prior $G^{*}$ in black. Bottom: a comparison of oracle Bayes $\hat{\theta}_{i}^{*}$ (left) based on knowledge of the prior distribution $G^{*}$ and empirical Bayes $\hat{\theta}_{i}$ (right), a function of the observed data.

Furthermore, when $G^{*}$ possesses finer structure, such as the clustering problem where $G^{*}$ is a discrete measure with $k^{*}$ atoms, we prove that the regret is bounded from above by $\frac{k^{*}}{n}$ up to logarithmic multiplicative factors in $n$. The clustering case is particularly remarkable, as
the NPMLE is completely tuning-free, with no knowledge of $k^{*}$, yet $\widehat{G}_{n}$ performs essentially as well as any estimator which knows the number of clusters $k^{*}$. Thus, Theorem 4.8 demonstrates that the NPMLE effectively discovers structure when available and also effectively learns when structure is unavailable. Theorem 4.8 generalizes the regret bounds of Saha and Guntuboyina (2020a) and Jiang (2020) who analyzed the homoscedastic $\Sigma_{i} \equiv \Sigma$ setting and the univariate $d=1$ setting, respectively. These papers in turn built upon Jiang and Zhang (2009) who studied the univariate, homoscedastic setting.

A key ingredient in the analysis of the regret is a more explicit representation of the estimator $\left(\hat{\theta}_{i}\right)_{i=1}^{n}$ and oracle $\left(\hat{\theta}_{i}^{*}\right)_{i=1}^{n}$. The oracle posterior mean (4.5) has the following alternative expression, known as Tweedie's formula (Banerjee et al., 2021; Dyson, 1926; Efron, 2011; Robbins, 1956):

$$
\begin{equation*}
\hat{\theta}_{i}^{*}=X_{i}+\Sigma_{i} \frac{\nabla f_{G^{*}, \Sigma_{i}}\left(X_{i}\right)}{f_{G^{*}, \Sigma_{i}}\left(X_{i}\right)} \tag{4.7}
\end{equation*}
$$

Similarly, our plug-in estimate can be written as

$$
\begin{equation*}
\hat{\theta}_{i}=X_{i}+\Sigma_{i} \frac{\nabla f_{\widehat{G}_{n}, \Sigma_{i}}\left(X_{i}\right)}{f_{\widehat{G}_{n}, \Sigma_{i}}\left(X_{i}\right)} . \tag{4.8}
\end{equation*}
$$

Tweedie's formula clarifies that under model (4.1) the posterior means only depend on the prior $G^{*}$ via the marginal likelihood $f_{G^{*}, \Sigma_{i}}\left(X_{i}\right)$ and its gradient. Jiang and Zhang (2009) first leveraged this observation to relate the empirical Bayes regret to the problem of estimating the marginal density. In heteroscedastic problems, there are $n$ different marginal densities, $\left(f_{G^{*}, \Sigma_{i}}\right)_{i=1}^{n}$, to estimate, and corresponding estimators $\left(f_{\widehat{G}_{n}, \Sigma_{i}}\right)_{i=1}^{n}$. We show in Theorem 4.6 and Corollary 4.7 that the NPMLE achieves similar adaptive rates in the density estimation problem under an appropriate average Hellinger distance across all $i=1, \ldots, n$ estimands $\left(f_{G^{*}, \Sigma_{i}}\right)_{i=1}^{n}$.

Whereas most recent work has focused on properties of $\widehat{G}_{n}$ for density estimation and denoising, the NPMLE is potentially much more generally applicable as a plug-in estimate of the prior. To expand our understanding of its applicability, we present the first analysis of the deconvolution error for the NPMLE. Whereas density estimation captures the problem of describing the observations $\left(X_{i}\right)_{i=1}^{n}$, deconvolution is the equally natural problem of interpreting the infinite-dimensional parameter $G^{*}$. We study the accuracy of the NPMLE under a Wasserstein distance $W_{2}\left(\widehat{G}_{n}, G^{*}\right)$. The Wasserstein distance is particularly useful for this problem since $\widehat{G}_{n}$ and $G^{*}$ are typically mutually singular; in particular, $G^{*}$ may be absolutely continuous whereas $\widehat{G}_{n}$ is always discrete. The Wasserstein distance will be discussed in detail in Section 4.3. We show in Theorem 4.10 that $\widehat{G}_{n}$ attains the minimax rate of deconvolution, which happens to be a very slow, logarithmic rate $\frac{1}{\log n}$. Inspired by the richness of the density estimation and denoising results, we hint at some of the adaptation properties of the NPMLE under the Wasserstein loss; Theorem 4.12 shows that when $G^{*}=\delta_{\mu}$ is a point mass distribution, the Wasserstein rate improves dramatically to $n^{-1 / 4}$ up to logarithmic factors.

The rest of the chapter is organized as follows: Section 4.2 systematically addresses basic properties of the NPMLE, including existence, discreteness, and non-uniqueness; Section 4.2.2 gives a full account of the approximate computation of NPMLEs. Section 4.3 establishes finite-sample risk bounds on the accuracy of $\widehat{G}_{n}$ as an estimator of $G^{*}$ for the purposes of density estimation, denoising and deconvolution. In Section 4.4, we apply the method to astronomy data to construct a fully data driven color-magnitude diagram of 1.4 million stars and compare our method to extreme deconvolution where it has previously been applied (Anderson et al., 2018). We also apply the method to chemical abundance data for a smaller subset of stars that has previously been analyzed by Ratcliffe et al. (2020). Section 4.5 concludes with some discussion of future work. The proofs are in Section 4.6.

### 4.2 Computational properties

### 4.2.1 Characterization and basic properties

In this section, we establish some basic properties of solutions to the nonparametric maximum likelihood problem (4.4), including existence, non-uniqueness, discreteness of solutions $\widehat{G}_{n}$, invariance under certain transformations, and bounds on the support. These results provide a foundation both for computing $\widehat{G}_{n}$ (Section 4.2.2) and for understanding its statistical properties (Section 4.3). Our first result extends the well-known characterization of $\widehat{G}_{n}$ for univariate, homoscedastic errors (Lindsay, 1995, Theorems 18-21) to our more general setting.

Lemma 4.1. Problem (4.4) attains its maximum: there exists a discrete solution $\widehat{G}_{n}$ with at most $n$ atoms, and the vector $\hat{L} \equiv\left(\hat{L}_{1}, \ldots, \hat{L}_{n}\right)=\left(f_{\widehat{G}_{n}, \Sigma_{i}}\left(X_{i}\right)\right)_{i=1}^{n}$ of fitted likelihood values is unique. Moreover, $\widehat{G}_{n} \in \mathcal{P}\left(\mathbb{R}^{d}\right)$ solves (4.4) if and only if

$$
D\left(\widehat{G}_{n}, \vartheta\right) \leq 0 \text { for all } \vartheta \in \mathbb{R}^{d} \text {, where } D(G, \vartheta):=\frac{1}{n} \sum_{i=1}^{n} \frac{\varphi_{\Sigma_{i}}\left(X_{i}-\vartheta\right)}{f_{G, \Sigma_{i}}\left(X_{i}\right)}-1 .
$$

The support of any $\widehat{G}_{n}$ is contained in the zero set $\mathcal{Z}:=\left\{\vartheta: D\left(\widehat{G}_{n}, \vartheta\right)=0\right\}$; the zero set $\mathcal{Z}$ is equal to the set of global maximizers of the $n$-component, heteroscedastic dual mixture density

$$
\widehat{\psi}_{n}(\vartheta):=\sum_{i=1}^{n}\left(\frac{\hat{L}_{i}^{-1}}{\sum_{\iota=1}^{n} \hat{L}_{\iota}^{-1}}\right) \varphi_{\Sigma_{i}}\left(X_{i}-\vartheta\right) .
$$

We prove Lemma 4.1, along with all results in this section, in Section 4.6.1. The first statement of the lemma guarantees the existence of a discrete solution, which we typically write as $\widehat{G}_{n}=\sum_{j=1}^{\hat{k}} \hat{w}_{j} \delta_{\hat{a}_{j}}$ (here $\hat{w}_{j} \geq 0, \sum_{j} \hat{w}_{j}=1$ and $\hat{a}_{j} \in \mathbb{R}^{d}$ ), with $\hat{k} \leq n$ providing an upper bound on the complexity of at least one solution. This implies that $\widehat{G}_{n}$ may be taken to be the maximum likelihood solution to a $\hat{k}$-component, heteroscedastic Gaussian
mixture model where $\hat{k}$ is selected in a data dependent manner. Since finite mixture models are nested by the number of components and $\hat{k} \leq n$, we may also say in general that $\widehat{G}_{n}$ is the maximum likelihood solution to an $n$-component, heteroscedastic Gaussian mixture model.

The bound $\hat{k} \leq n$ is tight: for each $n \geq 1$, there are sequences of observations $\left(X_{i}\right)_{i=1}^{n}$ and covariances $\left(\Sigma_{i}\right)_{i=1}^{n}$ such that the smallest number of components $\hat{k}$ of any solution $\widehat{G}_{n}$ to (4.4) is precisely $n$ (see, e.g., Lindsay, 1995, p. 116). However, in practice, the number of components is typically much smaller than $n$. For instance, in the univariate, homoscedastic case, Polyanskiy and $\mathrm{Wu}(2020)$ established a much stronger bound of $\hat{k}=O_{P}(\log n)$ under certain conditions on the prior distribution $G^{*}$.

The last part of Lemma 4.1 states that the atoms of $\widehat{G}_{n}$ occur at the global maximizers of the $n$-component Gaussian mixture $\widehat{\psi}_{n}$, which has component distributions of the form $\mathcal{N}\left(X_{i}, \Sigma_{i}\right)$ for $i=1, \ldots, n$ with weights inversely proportional to fitted likelihoods $\hat{L}$. Results on the modes of Gaussian mixtures (e.g. Améndola et al., 2020; Dytso et al., 2019; Ray \& Lindsay, 2005) thus provide information about the support of the NPMLE; in particular, our next two results exploit this connection to yield novel results on the NPMLE.

In the univariate $d=1$ and homoscedastic setting $\Sigma_{i} \equiv \sigma^{2}$, it is additionally known that (4.4) has a unique solution $\widehat{G}_{n}$ for all observations $X_{1}, \ldots, X_{n}$ (Lindsay \& Roeder, 1993). This means that, for every dataset $X_{1}, \ldots, X_{n}$ and every variance level $\sigma^{2}>0$, there is a unique probability measure $\widehat{G}_{n} \in \mathcal{P}(\mathbb{R})$ such that $\hat{L}_{i}=f_{\widehat{G}_{n}, \sigma^{2}}\left(X_{i}\right)$ for all $i$, where $\hat{L}$ is the unique vector of optimal likelihoods from Lemma 4.1. We observe, however, that uniqueness of the solution $\widehat{G}_{n}$ may not hold when $d>1$, even with isotropic covariances $\Sigma_{i} \equiv \sigma^{2} I_{d}$.

Lemma 4.2. Let $d=2, n=3$ and $X_{1}=(0,1), X_{2}=\left(\frac{\sqrt{3}}{2},-\frac{1}{2}\right), X_{3}=\left(-\frac{\sqrt{3}}{2},-\frac{1}{2}\right)$. Then (4.4) with data $\left(X_{i}\right)_{i=1}^{3}$, covariances $\Sigma_{i} \equiv \sigma^{2} I_{2}$ and $\sigma^{2}=3 /(\log 256)$ has infinitely many solutions of the form

$$
\widehat{G}_{n}=\alpha \delta_{0}+(1-\alpha) \frac{1}{3} \sum_{i=1}^{3} \delta_{X_{i} / 2}
$$

where $\alpha \in[0,1]$.
Figure 4.3 illustrates the counterexample given in Lemma 4.2. A key observation in the proof of Lemma 4.2 is that the dual mixture $\widehat{\psi}_{n}=f_{H, \sigma^{2} I_{2}}$ can be written explicitly as a homoscedastic mixture with uniform mixing distribution $H=\frac{1}{3} \sum_{i=1}^{3} \delta_{X_{i}}$ over the observations $\left(X_{i}\right)_{i=1}^{3}$. This set-up closely follows a construction, due to Duistermaat (see Améndola et al., 2020), exhibiting an isotropic, homoscedastic Gaussian mixture with more modes than components. Duistermaat used the same component locations $X_{i}$ but took $\sigma^{2}=0.53$ to obtain an example of a three-component mixture of isotropic, homoscedastic Gaussians such that the mixture has four modes. By specifically choosing $\sigma^{2}=\frac{3}{\log 256} \approx 0.54$, the height of the mixture $\widehat{\psi}_{n}=f_{H, \sigma^{2} I_{2}}$ is equal at all four modes, i.e. all four modes are global maximizers, and the modes are located at $\left\{X_{1} / 2, X_{2} / 2, X_{3} / 2,0\right\}$. By Lemma 4.1 any NPMLE must be
supported on these modes. Representing the fitted values $\hat{L}=\left(f_{\widehat{G}_{n}, \sigma^{2} I_{2}}\left(X_{i}\right)\right)_{i=1}^{3}$ by a probability measure $\widehat{G}_{n}=\sum_{j=1}^{3} \hat{w}_{j} \delta_{X_{j} / 2}+\hat{w}_{4} \delta_{0}$ supported on the global modes is equivalent to finding a set of weights $\hat{w} \in \mathbb{R}_{+}^{4}$ such that $\sum_{j=1}^{4} \hat{w}_{j}=1$ and $\hat{w}$ solves the under-determined linear system $\hat{L}=A \hat{w}$, where $A$ is a $3 \times 4$ matrix given by

$$
A_{i j}= \begin{cases}\varphi_{\sigma^{2} I_{2}}\left(X_{i}-X_{j} / 2\right) & j \leq 3 \\ \varphi_{\sigma^{2} I_{2}}\left(X_{i}\right) & j=4 .\end{cases}
$$

Finally, we also note that although the fitted likelihoods $f_{\widehat{G}_{n}, \sigma^{2} I_{2}}\left(X_{i}\right)$ are unique, the posterior means $\hat{\theta}_{i}$ in this example differ for the solutions $\widehat{G}_{n}$ given in Lemma 4.2.


Figure 4.3: Level sets of the dual mixture density $\widehat{\psi}_{n}=f_{H, \sigma^{2} I_{2}}$ where $n=3$ and $H=$ $\frac{1}{3} \sum_{i=1}^{3} \delta_{X_{i}}$ is uniform over the vertices of the larger equilateral triangle $\triangle X_{1} X_{2} X_{3}$. With $\sigma^{2}=\frac{3}{\log 256}$, the dual mixture density $\widehat{\psi}_{n}$ has four global modes.

Although the NPMLE searches over all probability measures $G \in \mathcal{P}\left(\mathbb{R}^{d}\right)$ supported on $\mathbb{R}^{d}$, it is useful algorithmically to reduce the search space to probability measures supported on a compact subset of $\mathbb{R}^{d}$. By Lemma 4.1, to restrict the support of the NPMLE it suffices to bound the maximizers $\mathcal{Z}$ of the $n$-component Gaussian mixture $\widehat{\psi}_{n}$. Ray and Lindsay (2005, Theorem 1) showed that all critical points of a Gaussian mixture $\widehat{\psi}_{n}(\vartheta)=$
$\sum_{i=1}^{n}\left(\frac{\hat{L}_{i}^{-1}}{\sum_{l=1}^{n} \hat{L}_{l}^{-1}}\right) \varphi_{\Sigma_{i}}\left(X_{i}-\vartheta\right)$ belong to the ridgeline manifold

$$
\begin{align*}
\mathcal{M} & :=\left\{x^{*}(\alpha): \alpha \in \mathbb{R}_{+}^{n}, \sum_{i=1}^{n} \alpha_{i}=1\right\}, \text { where } \\
x^{*}(\alpha) & :=\left(\sum_{i=1}^{n} \alpha_{i} \Sigma_{i}^{-1}\right)^{-1} \sum_{i=1}^{n} \alpha_{i} \Sigma_{i}^{-1} X_{i} . \tag{4.9}
\end{align*}
$$

In general, the ridgeline manifold $\mathcal{M}$ is a compact subset of $\mathbb{R}^{d}$ which does not depend on the weights $\left(\frac{\hat{L}_{i}^{-1}}{\sum_{\iota=1}^{n} \hat{L}_{\iota}^{-1}}\right)_{i=1}^{n}$. In the univariate case $d=1$, the rigdeline manifold $\mathcal{M}=\left[X_{(1)}, X_{(n)}\right]$ is simply the range of the data, so the univariate NPMLE is constrained to be supported on this range. In the multivariate setting, we may further simplify $\mathcal{M}$ depending on certain shape restrictions on the covariance matrices.

Corollary 4.3. Every solution to (4.4) is discrete with a finite number of atoms, supported on the ridgeline manifold $\mathcal{M}$ defined in (4.9). Depending on the values of $\left(\Sigma_{i}\right)$ we further bound the support as follows:
(i) (Homoscedastic) If $\Sigma_{i}=\Sigma$ for all $i$, or if $\Sigma_{i}=c_{i} \Sigma$ are proportional up to a sequence $\left(c_{i}\right)$ of positive scalars, the ridgeline manifold $\mathcal{M}$ is the convex hull of the data $\operatorname{conv}\left(\left\{X_{1}, \ldots, X_{n}\right\}\right)$.
(ii) (Diagonal Covariances) If $\Sigma_{i}$ is a diagonal matrix for every $i$, the ridgeline manifold $\mathcal{M}$ is contained in the axis-aligned minimum bounding box of the data

$$
\prod_{j=1}^{d}\left[\min _{i \in\{1, \ldots, n\}} X_{i j}, \max _{i \in\{1, \ldots, n\}} X_{i j}\right]
$$

where $X_{i}=\left(X_{i 1}, \ldots, X_{i d}\right)$ for all $i$.
(iii) (General Covariances) Let $\bar{k} \geq \underline{k}>0$ be chosen such that $\underline{k} I_{d} \preceq \Sigma_{i} \preceq \bar{k} I_{d}$ for all $i$, where $A \preceq B$ means $B-A$ is a symmetric positive semidefinite matrix. Choose $r>0$ and $x_{0} \in \mathbb{R}^{d}$ such that $\left\|X_{i}-x_{0}\right\|_{2} \leq r$ for all $i$. Then the ridgeline manifold $\mathcal{M}$ is contained in the ball

$$
\mathbb{B}_{\kappa r}\left(x_{0}\right):=\left\{y \in \mathbb{R}^{d}:\left\|y-x_{0}\right\|_{2} \leq \kappa r\right\}
$$

where $\kappa=\bar{k} / \underline{k}$.
The first part of Corollary 4.3 in general gives the smallest possible convex body over which the support of $\widehat{G}_{n}$ can be constrained independently of $\left\{\Sigma_{i}\right\}$. To see that the first part is tight, consider a fixed set of observations $\left(X_{i}\right)_{i=1}^{n}$ and isotropic covariance matrices
$\Sigma=\sigma^{2} I_{d}$; as $\sigma$ is made arbitrarily small, the support of $\widehat{G}_{n}$ approaches the set of observations $\left(X_{i}\right)_{i=1}^{n}$ (Lindsay, 1995). Therefore, in general the convex hull is the smallest convex body containing the support in the homoscedastic setting and more generally the setting of proportional covariance matrices. By contrast, the convex hull of the data is in general too small to capture the support of $\widehat{G}_{n}$ in the heteroscedastic setting. Figure 4.4 presents one example with diagonal covariances where the support of $\widehat{G}_{n}$ is pushed towards the corners of the minimum axis-aligned bounding box of the data. Thus, the above discussion and Figure 4.4 indicate that both parts (i) and (ii) of Corollary 4.3 give the tightest possible convex support bounds in their respective special cases.


Figure 4.4: Left: An example of observations $X_{1}=(0,1), X_{2}=(0,-1), X_{3}=(1,0)$, and $X_{4}=(-1,0)$ (blue points) with diagonal covariances $\Sigma_{1}=\Sigma_{2}=\left[\begin{array}{cc}5 & 0 \\ 0 & .05\end{array}\right]$ and $\Sigma_{3}=\Sigma_{4}=$ $\left[\begin{array}{cc}.05 & 0 \\ 0 & 5\end{array}\right]$ (dashed ellipses), where the NPMLE is supported on atoms $a_{1}, \ldots, a_{4}$ (red points) well outside the convex hull of the data, and near the corners of the minimum axis-aligned bounding box. Right: The mixture $\widehat{\psi}_{n}(\vartheta)=\frac{1}{4} \sum_{i=1}^{4} \varphi_{\Sigma_{i}}\left(X_{i}-\vartheta\right)$ only has modes at the atoms $a_{1}, \ldots, a_{4}$, so no NPMLE is supported within the convex hull of the data.

We close this section with a brief discussion on how the NPMLE behaves under certain simple transformations of the data $\left(X_{i}, \Sigma_{i}\right)_{i=1}^{n}$. Given a map $T: \mathbb{R}^{d} \rightarrow \mathbb{R}^{d}$, let $T_{\#} G \in \mathcal{P}\left(\mathbb{R}^{d}\right)$ denote the pushforward of $G \in \mathcal{P}\left(\mathbb{R}^{d}\right)$ given by $T_{\#} G(B)=G\left(T^{-1}(B)\right)$, for any Borel set $B \subseteq \mathbb{R}^{d}$. In other words, if $V \sim G$, then $T_{\#} G$ is the distribution of $T(V)$.

Lemma 4.4. Fix a dataset $\left(X_{i}, \Sigma_{i}\right)_{i=1}^{n}$, a point $x_{0} \in \mathbb{R}^{d}$ and ad $\times d$ orthogonal matrix $U_{0}$. Consider the transformed dataset $\left(X_{i}^{\prime}, \Sigma_{i}^{\prime}\right)_{i=1}^{n}$ where $\Sigma_{i}^{\prime}=U_{0} \Sigma_{i} U_{0}^{\top}$ and $X_{i}^{\prime}=T\left(X_{i}\right)$ for
$i=1, \ldots, n$, with $T(x)=U_{0} x+x_{0}$. Then

$$
f_{T_{\# G, \Sigma_{i}^{\prime}}}\left(X_{i}^{\prime}\right)=f_{G, \Sigma_{i}}\left(X_{i}\right)
$$

for all $i=1, \ldots, n$ and all $G \in \mathcal{P}\left(\mathbb{R}^{d}\right)$.
Lemma 4.4 is a straightforward consequence of the change of variables formula, but it has a number of useful corollaries. In particular, if $\widehat{G}_{n} \in \mathcal{P}\left(\mathbb{R}^{d}\right)$ is an NPMLE for the dataset $\left(X_{i}, \Sigma_{i}\right)_{i=1}^{n}$, then $T_{\#} \widehat{G}_{n}$ is an NPMLE for the modified dataset $\left(X_{i}^{\prime}, \Sigma_{i}^{\prime}\right)_{i=1}^{n}$, and the fitted likelihood values are the same, i.e.

$$
f_{T_{\#} \widehat{G}_{n}, \Sigma_{i}^{\prime}}\left(X_{i}^{\prime}\right)=f_{\widehat{G}_{n}, \Sigma_{i}}\left(X_{i}\right),
$$

for all $i=1, \ldots, n$. Thus, an NPMLE $\widehat{G}_{n}=\sum_{j=1}^{\hat{k}} \hat{w}_{j} \delta_{\hat{a}_{j}}$ is equivariant under translations $T(y)=y+x_{0}$ : if every observation is shifted by some fixed $x_{0} \in \mathbb{R}^{d}$, then the modified NPMLE $T_{\#} \widehat{G}_{n}=\sum_{j=1}^{\hat{k}} \hat{w}_{j} \delta_{\hat{a}_{j}+x_{0}}$ simply shifts every atom by $x_{0}$. Similarly, the NPMLE is equivariant under orthogonal transformations, which explains why the fitted likelihood values are all equal in the rotationally symmetric toy datasets presented in Figure 4.3 and Figure 4.4.

### 4.2.2 Grid approximation

The NPMLE solves a convex optimization problem (4.4) that is infinite-dimensional in the sense that the decision variable $G$ ranges over all probability measures on $\mathbb{R}^{d}$. Many numerical methods for approximately computing the NPMLE have been considered-including EM (Laird, 1978), vertex direction and exchange methods (Böhning, 1985), semi-infinite methods (Lesperance \& Kalbfleisch, 1992), constrained-Newton methods (Wang, 2007), and hybrid methods (Böhning, 2003; Liu \& Zhu, 2007)—typically described for the special case of univariate and homoscedastic errors. In this section, we discuss our strategy for computing the NPMLE as well as the challenges of scaling the computation to large datasets.

We follow the approach of Koenker and Mizera (2014), who approximated the infinitedimensional problem by constraining the support of $G$ to a large finite set. For a nonempty, closed set $\mathcal{A} \subseteq \mathbb{R}^{d}$, define a support-constrained NPMLE as any solution

$$
\begin{equation*}
\widehat{G}_{n}^{\mathcal{A}} \in \underset{G \in \mathcal{P}(\mathcal{A})}{\operatorname{argmax}} \frac{1}{n} \sum_{i=1}^{n} \log f_{G, \Sigma_{i}}\left(X_{i}\right), \tag{4.10}
\end{equation*}
$$

where $\mathcal{P}(\mathcal{A})$ denotes the set of probability measures supported on $\mathcal{A}$. In particular, $\widehat{G}_{n}=\widehat{G}_{n}^{\mathbb{R}^{d}}$ by definition, and by Corollary 4.3 we may write $\widehat{G}_{n}=\widehat{G}_{n}^{\mathcal{M}}$ for a compact subset $\mathcal{M}$ defined explicitly in terms of the data.

We now describe our strategy for choosing the discretization set $\mathcal{A}$. Fix $\delta>0$. Let $\mathcal{H}$ denote a covering of $\mathcal{M}$ by closed hypercubes of width $\delta$, i.e.

$$
\mathcal{H}=\left\{x_{j}+[-\delta / 2, \delta / 2]^{d}: j \in\{1, \ldots, J\}\right\}
$$

for some set of points $x_{1}, \ldots, x_{J} \in \mathbb{R}^{d}$ such that $\mathcal{M} \subseteq \bigcup_{j=1}^{J}\left(x_{j}+[-\delta / 2, \delta / 2]^{d}\right)$. Now define the discretized support $\mathcal{A}$ to be the set of corners of hypercubes in $\mathcal{H}$; specifically, for each hypercube $x_{j}+[-\delta / 2, \delta / 2]^{d}$ in $\mathcal{H}$, the point $x_{j}+\frac{\delta}{2} v \in \mathcal{A}$ for every $v \in\{-1,1\}^{d}$. Because $\mathcal{M}$ is compact, $\mathcal{A}$ is a finite set which we denote by $\left\{a_{j}\right\}_{j=1}^{m}$. Constraining the NPMLE to this finite set of atoms $a_{1}, \ldots, a_{m}$ yields a finite-dimensional convex optimization problem over the mixing proportions. That is, the solution to (4.10) can be written as $\widehat{G}_{n}^{\mathcal{A}}=\sum_{j=1}^{m} \tilde{w}_{j} \delta_{a_{j}}$, where

$$
\begin{equation*}
\tilde{w} \in \underset{w \in \Delta_{m-1}}{\operatorname{argmax}} \frac{1}{n} \sum_{i=1}^{n} \log \left(\sum_{j=1}^{m} L_{i j} w_{j}\right) \tag{4.11}
\end{equation*}
$$

and $L_{i j}=\varphi_{\Sigma_{i}}\left(X_{i}-a_{j}\right)$ encodes an $n \times m$ kernel matrix. The EM algorithm (Dempster et al., 1977) can be used to optimize directly over the mixing proportions $\tilde{w}$. While this approach was advocated by Lashkari and Golland (2008) and Jiang and Zhang (2009), EM can be prohibitively slow (Koenker \& Mizera, 2014; Redner \& Walker, 1984). A crucial observation made by Koenker and Mizera (2014) is that (4.11) is a (finite-dimensional) convex optimization problem, enabling the use of a wide array of tools from modern convex optimization; they proposed solving the dual to (4.11) using an interior point solver, and Koenker and Gu (2017) provided an R implementation to solve univariate problems. Kim et al. (2020) proposed sequential quadratic programming to solve a variant of the primal problem directly, demonstrating superior scalability with the sample size $n$. Our implementation uses the MOSEK library (MOSEK ApS, 2019) for Python.

To justify the grid approximation, some consideration of the discretization error is warranted. Our next result shows that as $\delta \downarrow 0$, the log-likelihood of the discretized NPMLE approaches that of the (unconstrained) NPMLE; moreover, the bound on the gap depends on known quantities, so it can be used to guide a suitable choice of $\delta$.

Proposition 4.5. Let $\mathcal{M} \subset \mathbb{R}^{d}$ denote any compact set such that every solution (4.4) is supported on $\mathcal{M}$. Suppose the diameter of the set $\mathcal{M}$ is at most $D$, the minimum eigenvalue of each $\Sigma_{i}$ is at least $\underline{k}$, and fix $\delta \in\left(0, \sqrt{\frac{3}{4 d}} \underline{k} D^{-1}\right)$. Let $\mathcal{H}$ denote a cover of $\mathcal{M}$ by closed hypercubes of width $\delta$, and let $\mathcal{A}$ denote the set of corners of hypercubes in $\mathcal{H}$. Every approximate $N P M L E \widehat{G}_{n}^{\mathcal{A}}$ satisfies

$$
\begin{equation*}
\sup _{G \in \mathcal{P}\left(\mathbb{R}^{d}\right)} \frac{1}{n} \sum_{i=1}^{n} \log f_{G, \Sigma_{i}}\left(X_{i}\right)-\frac{1}{n} \sum_{i=1}^{n} \log f_{\widehat{G}_{n}^{A}, \Sigma_{i}}\left(X_{i}\right) \leq d \underline{k}^{-2}\left(2 D^{2}+\frac{1}{2}\right) \delta^{2} . \tag{4.12}
\end{equation*}
$$

We prove Proposition 4.5 in Section 4.6.1. Proposition 4.5 shows that we can tractably approximate the NPMLE via a finite-dimensional, convex optimization problem. As we show in Section 4.3, our theoretical results on the statistical properties of the NPMLE hold for any approximate solution $\widehat{G}_{n}^{\mathcal{A}}$ which places nearly as high likelihood on the observations as the global optimizer, in the sense of (4.12). Hence for $\delta$ sufficiently small we can guarantee that the discretization error is negligible.

Dicker and Zhao (2016) showed in the univariate, homoscedastic case that a finely discretized NPMLE is statistically indistinguishable from the NPMLE for the purpose of density estimation. However, their analysis of the discretization error makes use of the modeling assumptions (4.1) and is statistical in nature, so their theoretical results provide little guidance on how much error is incurred due to discretization for a fixed dataset. Our result aligns more closely with and in fact essentially generalizes Jiang and Zhang (2009, Proposition 5), which bounded the optimality gap for a particular algorithm, discretization scheme and fixed dataset. The main difference between our result and Jiang and Zhang (2009, Proposition 5) is that the latter analyzed the EM algorithm for the mixing proportions (4.11), whereas by using a black-box, second-order optimization method to solve for the mixing proportions $\tilde{w}$, we can solve for the discretized NPMLE $\widehat{G}_{n}^{\mathcal{A}}$ much more accurately.

### 4.3 Statistical properties

The NPMLE $\widehat{G}_{n}$ applies as a plug-in estimator of the prior distribution $G^{*}$ for many purposes. The traditional statistical setting is density estimation, where working in a Gaussian mixture model greatly simplifies the problem of estimating the marginal density of each observation $X_{i}$. In particular, $f_{\widehat{G}_{n}, \Sigma_{i}}$ is a natural, tuning-free estimate of the true marginal density $f_{G^{*}, \Sigma_{i}}$. Another problem setting-at the heart of empirical Bayes methodology - is to imitate the Bayesian inference we would conduct if we knew $G^{*}$. Denoising, using $\left(\hat{\theta}_{i}\right)_{i=1}^{n}$ as plug-in estimators of the true posterior means $\left(\hat{\theta}_{i}^{*}\right)_{i=1}^{n}$, represents the most basic instantiation. Finally, often we wish to compare $\widehat{G}_{n}$ to the prior $G^{*}$ directly. Since we are estimating the prior given observations from a convolution model $X_{i} \stackrel{\text { ind }}{\sim} f_{G^{*}, \Sigma_{i}}$, deconvolution refers to the problem of estimating $G^{*}$.

In this section, we establish that the NPMLE is well-suited for all three disparate targets of estimation: the marginal densities $\left(f_{G^{*}, \Sigma_{i}}\right)_{i=1}^{n}$, the oracle posterior means $\left(\hat{\theta}_{i}^{*}\right)_{i=1}^{n}$ and the prior $G^{*}$. In this section, we allow for the possibility that $\widehat{G}_{n}$ is an approximate NPMLE, with the exact conditions being given in each theorem. Throughout this section, we use the standard notation $X \lesssim_{p, q} Y$ to mean $X \leq C_{p, q} Y$ for some positive constant $C_{p, q}>0$ depending only on problem parameters $p, q$.

### 4.3.1 Density estimation: average Hellinger accuracy

As the distribution of $X_{i}$ varies with $i$, we consider the density estimation quality of the NPMLE (4.4) in terms of the average squared Hellinger distance, i.e. for $G, H \in \mathcal{P}\left(\mathbb{R}^{d}\right)$,

$$
\bar{h}^{2}\left(f_{G, \bullet}, f_{H, \bullet}\right):=\frac{1}{n} \sum_{i=1}^{n} h^{2}\left(f_{G, \Sigma_{i}}, f_{H, \Sigma_{i}}\right)
$$

where $h^{2}(f, g)=\frac{1}{2} \int(\sqrt{f}-\sqrt{g})^{2}$ denotes the usual squared Hellinger distance between a pair of densities $f, g$. In the homoscedastic case where $\Sigma_{i} \equiv \Sigma$, our proposed loss func-
tion $\bar{h}^{2}\left(f_{G, \bullet}, f_{H, \bullet}\right)=h^{2}\left(f_{G, \Sigma}, f_{H, \Sigma}\right)$ agrees with the usual squared Hellinger distance. Our first result bounds the average squared Hellinger accuracy $\bar{h}^{2}\left(f_{\widehat{G}_{n}, \boldsymbol{\bullet}}, f_{G^{*}, \bullet}\right)$ of the NPMLE. In order to accommodate general heteroscedastic $\Sigma_{i}$, we state our results in terms of uniform upper and lower bounds on the spectra of all of the matrices, i.e. $\underline{k} I_{d} \preceq \Sigma_{i} \preceq \bar{k} I_{d}$ for all $i$. To state the result, some additional notation is needed. We fix a positive scalar $M \geq \sqrt{10 \bar{k} \log n}$ and a nonempty compact set $S \subset \mathbb{R}^{d}$. Define the rate function controlling the squared Hellinger distance

$$
\begin{equation*}
\varepsilon_{n}^{2}\left(M, S, G^{*}\right):=\operatorname{Vol}\left(S^{\bar{k}^{1 / 2}}\right) \frac{M^{d}}{n}(\log n)^{d / 2+1}+\inf _{q \geq(d+1) /(2 \log n)}\left(\frac{2 \mu_{q}}{M}\right)^{q} \log n \tag{4.13}
\end{equation*}
$$

where $\mu_{q}$ denotes the $q^{\text {th }}$-moment of $\mathfrak{d}_{S}(\vartheta):=\inf _{s \in S}\|\vartheta-s\|_{2}$ under $\vartheta \sim G^{*}$, and $S^{a}:=$ $\left\{y: \mathfrak{d}_{S}(y) \leq a\right\}$ denotes the $a$-enlargement of the set $S$. Note that we have suppressed the dependence of $\varepsilon_{n}^{2}$ on the upper bound $\bar{k}$.

The following result states that $\varepsilon_{n}^{2}\left(M, S, G^{*}\right)$ bounds the rate in average Hellinger accuracy both with high probability and in expectation. The scalar $M \geq \sqrt{10 \bar{k} \log n}$ and compact set $S \neq \emptyset$ are free parameters. Note that the first term on the right-hand side of (4.13) is increasing in $M$ and $S$, whereas the second is decreasing in each. In principle, then, we may tune the values of $M$ and $S$ to optimize the rate function $\varepsilon_{n}^{2}\left(M, S, G^{*}\right)$. Later in this section, we discuss a number of special cases where a more explicit rate can be obtained.

Theorem 4.6. Suppose $X_{i} \stackrel{\text { ind }}{\sim} f_{G^{*}, \Sigma_{i}}$ where $\underline{k} I_{d} \preceq \Sigma_{i} \preceq \bar{k} I_{d}$ for all $i$. Any (approximate) solution $\widehat{G}_{n} \in \mathcal{P}\left(\mathbb{R}^{d}\right)$ of (4.4) satisfying

$$
\begin{equation*}
\sup _{G \in \mathcal{P}\left(\mathbb{R}^{d}\right)} \frac{1}{n} \sum_{i=1}^{n} \log f_{G, \Sigma_{i}}\left(X_{i}\right)-\frac{1}{n} \sum_{i=1}^{n} \log f_{\widehat{G}_{n}, \Sigma_{i}}\left(X_{i}\right) \lesssim_{d, \bar{k}, \underline{k}} \varepsilon_{n}^{2}\left(M, S, G^{*}\right) \tag{4.14}
\end{equation*}
$$

satisfies

$$
\begin{equation*}
\mathbb{P}\left(\bar{h}^{2}\left(f_{\widehat{G}_{n}, \bullet}, f_{G^{*}, \bullet}\right) \gtrsim_{d, \bar{k}, \underline{,}} t^{2} \varepsilon_{n}^{2}\left(M, S, G^{*}\right)\right) \leq 2 n^{-t^{2}} \tag{4.15}
\end{equation*}
$$

for all $t \geq 1$, provided $n>\max \left(e \underline{k}^{-d / 2},(2 \pi)^{d / 2}\right)$. Moreover,

$$
\begin{equation*}
\mathbb{E}\left[\bar{h}^{2}\left(f_{\widehat{G}_{n}, \bullet}, f_{G^{*}, \bullet}\right)\right] \lesssim_{d, \bar{k}, \underline{k}} \varepsilon_{n}^{2}\left(M, S, G^{*}\right) \tag{4.16}
\end{equation*}
$$

We prove Theorem 4.6 in Section 4.6.2. Our proof extends Theorem 2.1 of Saha and Guntuboyina (2020a) on the multivariate, homoscedastic case $\Sigma_{i} \equiv I_{d}$ and Theorem 4 of Jiang (2020) on the univariate, heteroscedastic case $d=1$, which in turn build upon Theorem 1 of Zhang (2009) on the univariate, homoscedastic case. The general theory on rates of convergence for maximum likelihood estimators (van de Geer, 2000; Wong \& Shen, 1995) can in principle be used to bound $\bar{h}^{2}\left(f_{\widehat{G}_{n}, \bullet}, f_{G^{*}, \bullet}\right)$. Our proof technique deviates from the general
theory by directly bounding the likelihood ${f_{\widehat{G}_{n}, \Sigma_{i}}(x) \text { for } x \text { outside some pre-specified domain }}$ (controlled by the choice of set $S$ ), and then covering the set of densities $\left\{f_{G, \bullet}: G \in \mathcal{P}\left(\mathbb{R}^{d}\right)\right\}$ within the domain in the $L_{\infty}$ metric.

Theorem 4.6 provides a sharp bound in many special cases of $G^{*}$. For a given $G^{*}$ we need to optimize over the choices of $M \geq \sqrt{10 \bar{k} \log n}$ and the nonempty compact set $S \subset \mathbb{R}^{d}$ to obtain the smallest value of the rate function $\varepsilon_{n}^{2}\left(M, S, G^{*}\right)$. Our next result performs this calculation for various assumptions on the prior $G^{*}$.

Corollary 4.7. Suppose $X_{i} \stackrel{\text { ind }}{\sim} f_{G^{*}, \Sigma_{i}}$ where $\underline{k} I_{d} \preceq \Sigma_{i} \preceq \bar{k} I_{d}$ for all $i$. Suppose $\widehat{G}_{n} \in \mathcal{P}\left(\mathbb{R}^{d}\right)$ is any approximate NPMLE such that

$$
\begin{equation*}
\sup _{G \in \mathcal{P}\left(\mathbb{R}^{d}\right)} \frac{1}{n} \sum_{i=1}^{n} \log f_{G, \Sigma_{i}}\left(X_{i}\right)-\frac{1}{n} \sum_{i=1}^{n} \log f_{\widehat{G}_{n}, \Sigma_{i}}\left(X_{i}\right) \lesssim_{d, \bar{k}, \underline{k}} \frac{(\log n)^{d+1}}{n} \tag{4.17}
\end{equation*}
$$

(i) (Discrete support) If $G^{*}=\sum_{j=1}^{k^{*}} w_{j}^{*} \delta_{a_{j}^{*}}$, then

$$
\mathbb{E} \bar{h}^{2}\left(f_{\widehat{G}_{n}, \bullet}, f_{G^{*}, \bullet}\right) \lesssim_{d, \bar{k}, \underline{k}} \frac{k^{*}}{n}(\log n)^{d+1}
$$

(ii) (Compact support) If $G^{*}$ has compact support $S^{*}$, then

$$
\mathbb{E} \bar{h}^{2}\left(f_{\widehat{G}_{n}, \bullet}, f_{G^{*}, \bullet}\right) \lesssim_{d, \bar{k}, \underline{k}} \frac{\operatorname{Vol}\left(S^{*}+\mathbb{B}_{\bar{k}^{1 / 2}}(0)\right)}{n}(\log n)^{d+1}
$$

where $\mathbb{B}_{r}(x):=\left\{y:\|x-y\|_{2} \leq r\right\}$ denotes the d-dimensional ball of radius $r$ centered at $x$.
(iii) (Simultaneous moment control) Suppose that there is a compact $S^{*} \subset \mathbb{R}^{d}$ and $\alpha \in(0,2]$, $K \geq 1$ such that $\mu_{q}:=\mathbb{E}_{\vartheta \sim G^{*}}\left[\mathfrak{d}^{q}(\vartheta, S)\right]^{1 / q} \leq K q^{1 / \alpha}$ for all $q \geq 1$ (recall $\mathfrak{d}_{S}(\vartheta):=$ $\inf _{s \in S}\|\vartheta-s\|_{2}$ as above). Then

$$
\mathbb{E} \bar{h}^{2}\left(f_{\widehat{G}_{n}, \bullet}, f_{G^{*}, \bullet}\right) \lesssim_{\alpha, K, d, \bar{k}, \underline{k}} \frac{\operatorname{Vol}\left(S^{*}+\mathbb{B}_{\bar{k}^{1 / 2}}(0)\right)}{n}(\log n)^{\frac{2+\alpha}{2 \alpha} d+1}
$$

(iv) (Finite $q^{\text {th }}$ moment) Suppose that there is a compact $S^{*} \subset \mathbb{R}^{d}$ and $\mu, q>0$ such that $\mu_{q} \leq \mu$. Then

$$
\mathbb{E} \bar{h}^{2}\left(f_{\widehat{G}_{n}, \bullet}, f_{G^{*}, \bullet}\right) \lesssim_{\mu, q, d, \bar{k}, \underline{k}}\left(\frac{\operatorname{Vol}\left(S^{*}+\mathbb{B}_{\bar{k}^{1 / 2}}(0)\right)}{n}\right)^{\frac{q}{q+d}}(\log n)^{\frac{q}{2 q+2 d} d+1}
$$

Given the general result in Theorem 4.6, Corollary 4.7 follows directly from the calculations of Saha and Guntuboyina (2020a) in Corollary 2.2 and Theorem 2.3. Corollary 4.7 captures an important adaptation property of the NPMLE. The cases $(i)-(i v)$ described in the result are nested in the sense that (i) implies (ii), (ii) implies (iii), and (iii) implies (iv); consequently the rates get progressively worse as our assumptions weaken. This means that the NPMLE, despite searching over all probability measures $\mathcal{P}\left(\mathbb{R}^{d}\right)$, obtains better rates when structure is present in the prior $G^{*}$.

Most strikingly, when $G^{*}$ has discrete support with $k^{*}$ support points, the rate in $(i)$ is $\frac{k^{*}}{n}$ up to logarithmic factors without assuming any knowledge of $k^{*}$. This rate matches the minimax rate over all discrete distributions with at most $k^{*}$ support points (Saha \& Guntuboyina, 2020a), meaning we could not expect to do much better even if $k^{*}$ were known. In the extreme case where $k^{*}=1$, the observations actually come from a simple Gaussian, i.e. $f_{G^{*}, \Sigma_{i}}(x)=\varphi_{\Sigma_{i}}\left(x-a_{1}^{*}\right)$ with common mean $a_{1}^{*} \in \mathbb{R}^{d}$, so our result says we don't lose much in the rate when we model the density with a mixture even when it turns out to be a simple Gaussian. Similarly, in (ii), the rate adapts to the size of the support $S^{*}$ without prior knowledge of this support or even a bound on its size. Up through simultaneous moment control (iii), the dimension $d$ only affects the rate as a function of $n$ through the logarithmic factor. Hence, the NPMLE avoids the usual curse of dimensionality to some extent, while still achieving consistency in the heavier tailed setting (iv). The logarithmic factors in our bounds might be reduced slightly but cannot be eliminated as they are present in the minimax lower bounds (Kim \& Guntuboyina, 2020).

### 4.3.1.1 Implications for the Discretization Rate

Theorem 4.6 establishes that up to a multiplicative constant (depending only on the dimension $d$ and bounds $\underline{k}, \bar{k}$ on the eigenvalues of the covariance matrices) the quantity $\varepsilon_{n}^{2}\left(M, S, G^{*}\right)$ controls the average Hellinger accuracy $\mathbb{E}\left[\bar{h}^{2}\left(f_{\widehat{G}_{n}, \boldsymbol{\bullet}}, f_{G^{*}, \bullet}\right)\right]$ of the NPMLE. This also holds for approximate solutions to the optimization problem (4.4) that, in accordance with (4.14), place nearly as much likelihood on the data as does a global maximizer. It is natural to compare the requirement (4.14) with our computational guarantee on the discretization error (4.12) from Proposition 4.5. The free parameter which controls the discretization error is the resolution $\delta>0$, which represents the width of the hypercubes we use to cover the ridgeline manifold $\mathcal{M}$ or any of its outer-approximations from Corollary 4.3. Thus, in order to satisfy the main requirement of Theorem 4.6 , we need to take $\delta$ such that

$$
\varepsilon_{n}^{2}\left(M, S, G^{*}\right) \gtrsim_{d, \bar{k}, \underline{k}} d \underline{k}^{-2}\left(2 D^{2}+\frac{1}{2}\right) \delta^{2} .
$$

Observe from the definition of $\varepsilon_{n}^{2}$ that $\varepsilon_{n}^{2}\left(M, S, G^{*}\right) \gtrsim_{d, \bar{k}, \underline{k}} \frac{(\log n)^{d+1}}{n}$ for all $M \geq \sqrt{10 \bar{k} \log n}$ and all compact $S$. Absorbing additional terms depending on $d, \bar{k}$, and $\underline{k}$ and assuming for
simplicity that $D>\frac{1}{2}$, choosing $\delta$ such that

$$
\begin{equation*}
D^{2} \delta^{2} \lesssim_{d, \bar{k}, \underline{k}} \frac{(\log n)^{d+1}}{n} \tag{4.18}
\end{equation*}
$$

suffices for the discretized NPMLE to be statistically indistinguishable from a global maximizer.

The inequality (4.18) gives a preliminary bound on the rate at which the discretization level $\delta$ should decrease with $n$. Still, recall from Proposition 4.5 that $D$ denotes the diameter of the ridgeline manifold $\mathcal{M}$, so $D$ does depend on $n$. To sketch the dependence, let us consider a representative example where $G^{*}$ has sub-Gaussian tails and all of the $\Sigma_{i}$ 's are diagonal. In this case, by Corollary 4.3 part (ii), the ridgeline manifold $\mathcal{M}$ is contained in the axis-aligned minimum bounding box of the data

$$
\prod_{j=1}^{d}\left[\min _{i \in\{1, \ldots, n\}} X_{i j}, \max _{i \in\{1, \ldots, n\}} X_{i j}\right]
$$

Due to the tail condition, the length of each side of this hyper-rectangle grows like $\sqrt{\log n}$ with high probability up to multiplicative factors depending on $\bar{k}$ : hence, the diameter $D$ also scales like $\sqrt{\log n}$ with high probability up to multiplicative factors depending on $\bar{k}$ and $d$. We have thus shown that it suffices to discretize at a resolution of $\delta \asymp \sqrt{\frac{(\log n)^{d}}{n}}$. The number of points in our covering $\mathcal{A}$ is of order $m \asymp\left(\frac{n}{(\log n)^{d}}\right)^{d / 2}$. In the univariate case $d=1$, this slightly improves the finding of Theorem 2 of Dicker and Zhao (2016), who showed that an $m=\sqrt{n}$-discretization of the range of the data $\left[X_{(1)}, X_{(n)}\right]$ suffices for the same rate in Hellinger distance. Their bound on the large-deviation probability is also logarithmic, i.e. $O\left(\frac{1}{\log n}\right)$ whereas our equation (4.15) is polynomial in $n$. Our analysis also clarifies that the sense in which we need approximate NPMLE (4.14) is through the likelihood of the observations, relative to the global optimum, which could be useful for comparing alternative approaches to approximating the NPMLE.

### 4.3.2 Denoising: an oracle inequality

In this section we turn to the problem of estimating the oracle posterior means $\left(\hat{\theta}_{i}^{*}\right)_{i=1}^{n}$; see (4.5). We evaluate the performance of $\left(\hat{\theta}_{i}\right)_{i=1}^{n}$ (see (4.6)) as an estimator for $\left(\hat{\theta}_{i}^{*}\right)_{i=1}^{n}$ using the mean squared error risk measure:

$$
\frac{1}{n} \sum_{i=1}^{n} \mathbb{E}\left\|\hat{\theta}_{i}-\hat{\theta}_{i}^{*}\right\|_{2}^{2}
$$

Since $\hat{\theta}_{i}^{*}$ is the optimal estimator of $\theta_{i}^{*}$ given model (4.1), the above mean squared error quantifies the price of misspecifying $G^{*}$ with the data-driven estimator $\widehat{G}_{n}$. Hence, this loss is also known as the per-instance empirical Bayes regret.

Our next result states that the rate function $\varepsilon_{n}^{2}\left(M, S, G^{*}\right)$ governing the Hellinger accuracy (see (4.13)) also upper bounds the regret, up to additional logarithmic factors. We provide the same special cases of the rate as those stated in Corollary 4.7.

Theorem 4.8. Suppose $X_{i} \stackrel{\text { ind }}{\sim} f_{G^{*}, \Sigma_{i}}$ where $\underline{k} I_{d} \preceq \Sigma_{i} \preceq \bar{k} I_{d}$ for all $i$. Let $\widehat{G}_{n}$ denote any approximate NPMLE satisfying (4.17). Fix some $M \geq \sqrt{10 \bar{k} \log n}$ and a nonempty, compact set $S \subset \mathbb{R}^{d}$. Define $\varepsilon_{n}^{2}\left(M, S, G^{*}\right)$ as in (4.13). For all $n \geq 5 \underline{k}^{-d / 2} \vee(2 \pi)^{d / 2}$,

$$
\begin{equation*}
\frac{1}{n} \sum_{i=1}^{n} \mathbb{E}\left\|\hat{\theta}_{i}-\hat{\theta}_{i}^{*}\right\|_{2}^{2} \lesssim_{d, \bar{k}, \underline{,},} \varepsilon_{n}^{2}\left(M, S, G^{*}\right)(\log n)^{(d / 2-1) \vee 3} \tag{4.19}
\end{equation*}
$$

In particular, consider the following special cases for $G^{*}$ :
(i) (Discrete support) If $G^{*}=\sum_{j=1}^{k^{*}} w_{j}^{*} \delta_{a_{j}^{*}}$, then

$$
\frac{1}{n} \sum_{i=1}^{n} \mathbb{E}\left\|\hat{\theta}_{i}-\hat{\theta}_{i}^{*}\right\|_{2}^{2} \lesssim_{d, \bar{k}, \underline{k}} \frac{k^{*}}{n}(\log n)^{d+((d / 2) \vee 4)}
$$

(ii) (Compact support) If $G^{*}$ has compact support $S^{*}$, then

$$
\frac{1}{n} \sum_{i=1}^{n} \mathbb{E}\left\|\hat{\theta}_{i}-\hat{\theta}_{i}^{*}\right\|_{2}^{2} \lesssim_{d, \bar{k}, \underline{k}} \frac{\operatorname{Vol}\left(S^{*}+\mathbb{B}_{\bar{k}^{1 / 2}}(0)\right)}{n}(\log n)^{d+((d / 2) \vee 4)}
$$

(iii) (Simultaneous moment control) Suppose that there is a compact $S^{*} \subset \mathbb{R}^{d}$ and $\alpha \in(0,2]$, $K \geq 1$ such that $\mu_{q}:=\mathbb{E}_{\vartheta \sim G^{*}}\left[\mathfrak{d}^{q}\left(\vartheta, S^{*}\right)\right]^{1 / q} \leq K q^{1 / \alpha}$ for all $q \geq 1$. Then

$$
\frac{1}{n} \sum_{i=1}^{n} \mathbb{E}\left\|\hat{\theta}_{i}-\hat{\theta}_{i}^{*}\right\|_{2}^{2} \lesssim_{\alpha, K, d, \bar{k}, \underline{k}} \frac{\operatorname{Vol}\left(S^{*}+\mathbb{B}_{\bar{k}^{1 / 2}}(0)\right)}{n}(\log n)^{\frac{2 \alpha d}{2+\alpha}+((d / 2) \vee 4)}
$$

(iv) (Finite $q^{\text {th }}$ moment) Suppose that there exists a compact $S^{*} \subset \mathbb{R}^{d}$ and $\mu, q>0$ such that $\mu_{q} \leq \mu$. Then

$$
\frac{1}{n} \sum_{i=1}^{n} \mathbb{E}\left\|\hat{\theta}_{i}-\hat{\theta}_{i}^{*}\right\|_{2}^{2} \lesssim_{\mu, q, d, \bar{k}, \underline{k}}\left(\frac{\operatorname{Vol}\left(S^{*}+\mathbb{B}_{\bar{k}^{1 / 2}}(0)\right)}{n}\right)^{\frac{q}{q+d}}(\log n)^{\frac{q d}{2 q+2 d}}+((d / 2) \vee 4) .
$$

Theorem 4.8 shows that the denoising problem shares the adaptation features as the density estimation problem. Since we have assumed $\underline{k} I_{d} \preceq \Sigma_{i} \preceq \bar{k} I_{d}$ for all $i=1, \ldots, n$, the same set of results also hold for the scaled regret $\frac{1}{n} \sum_{i=1}^{n} \mathbb{E}\left(\hat{\theta}_{i}-\hat{\theta}_{i}^{*}\right)^{\top} \Sigma_{i}^{-1}\left(\hat{\theta}_{i}-\hat{\theta}_{i}^{*}\right)$.

Remark 4.9. (On the proof of Theorem 4.8 in Section 4.6.3) Our proof extends Theorem 3.1 of Saha and Guntuboyina (2020a) on the multivariate, homoscedastic case $\Sigma_{i} \equiv I_{d}$ and Theorem 1 of Jiang (2020) on the univariate, heteroscedastic case $d=1$, which in turn build upon Theorem 5 of Jiang and Zhang (2009) on the univariate, homoscedastic case. Jiang and Zhang (2009) and Jiang (2020) used a related notion of regret

$$
\sqrt{\frac{1}{n} \sum_{i=1}^{n} \mathbb{E}\left\|\hat{\theta}_{i}-\theta_{i}^{*}\right\|_{2}^{2}}-\sqrt{\frac{1}{n} \sum_{i=1}^{n} \mathbb{E}\left\|\hat{\theta}_{i}^{*}-\theta_{i}^{*}\right\|_{2}^{2}}
$$

Tweedie's formula relates the oracle (4.7) and empirical Bayes (4.8) posterior means to the corresponding marginal likelihoods, so the density estimation results of the previous section turn out to be useful for proving Theorem 4.8 as well. In particular, we consider Bayes rules for priors in a covering of the Hellinger ball

$$
\left\{G \in \mathcal{P}\left(\mathbb{R}^{d}\right): \bar{h}^{2}\left(f_{G, \bullet}, f_{G^{*}, \bullet}\right) \lesssim_{d, \bar{k}, \underline{k}} t^{2} \varepsilon_{n}^{2}\left(M, S, G^{*}\right)\right\}
$$

which, by Theorem 4.6, contains $\widehat{G}_{n}$ with high probability. For a fixed prior $G$, the denominator in the correction factor of Tweedie's formula

$$
X_{i}+\Sigma_{i} \frac{\nabla f_{G, \Sigma_{i}}\left(X_{i}\right)}{f_{G, \Sigma_{i}}\left(X_{i}\right)}
$$

namely $f_{G, \Sigma_{i}}\left(X_{i}\right)$, can be small. To avoid dividing by near-zero quantities, we regularize the above Bayes rule by replacing the denominator with $\max \left\{f_{G, \Sigma_{i}}\left(X_{i}\right), \rho\right\}$ for a small positive $\rho$. To handle heteroscedastic errors, we show that Tweedie's formula, even its regularized form, is equivariant under scale transformations.

### 4.3.3 Deconvolution: estimating the prior

We turn to the fundamental question of how well $\widehat{G}_{n}$ estimates $G^{*}$. This is known as the deconvolution problem and has received much attention in the statistical literature (Meister, 2009). Indeed, the original consistency results (Kiefer \& Wolfowitz, 1956; Pfanzagl, 1988) for the NPMLE focused on weak convergence of $\widehat{G}_{n}$ to $G^{*}$ as $n \rightarrow \infty$. While most prior work on deconvolution has focused on deconvolution with homoscedastic error distributions, Delaigle and Meister (2008) allowed for heteroscedastic errors but relied on kernel estimators which contain additional smoothing parameters. By contrast, the NPMLE provides a tuning-free estimate of the mixing distribution $G^{*}$, yet to our knowledge, non-asymptotic bounds on the rate of convergence for $\widehat{G}_{n}$ in the deconvolution problem are not known.

In practice, the true prior $G^{*}$ may not be discrete even though $\widehat{G}_{n}$ always is, and even if both distributions are discrete, their supports will typically differ. Our loss function must allow for comparisons of probability measures with potentially disjoint supports. Nguyen
(2013) established that a natural loss for this problem is the Wasserstein distance from the theory of optimal transport

$$
W_{2}^{2}(G, H):=\min _{(U, V) \in \Pi_{G, H}} \mathbb{E}\|U-V\|_{2}^{2}
$$

where $G, H \in \mathcal{P}\left(\mathbb{R}^{d}\right)$ are two probability measures and $\Pi_{G, H}$ denotes the set of couplings of $G$ and $H$, i.e. joint distributions over $(U, V) \in \mathbb{R}^{2 d}$ such that $U \sim G$ and $V \sim H$. Indeed, even the likelihood criterion is intimately related to the Wasserstein distance: in the homoscedastic case $\Sigma_{i} \equiv \sigma^{2} I_{d}$, it is known that the NPMLE (4.4) equivalently solves an entropic-regularized optimal transport problem (Rigollet \& Weed, 2018).

Nguyen (2013) connected the deconvolution error $W_{2}^{2}(G, H)$ to the density estimation error between the mixtures, i.e. $h^{2}\left(f_{G, I_{d}}, f_{H, I_{d}}\right)$ in a homoscedastic Gaussian deconvolution setting. By leveraging similar techniques as well as the support bounds of Corollary 4.3, we arrive at the following upper bound on the deconvolution error.

Theorem 4.10. Suppose $X_{i} \stackrel{\text { ind }}{\sim} f_{G^{*}, \Sigma_{i}}$ where $\underline{k} I_{d} \preceq \Sigma_{i} \preceq \bar{k} I_{d}$ and $\Sigma_{i}$ is a diagonal matrix for each $i$. Suppose further that $G^{*}\left([-L, L]^{d}\right)=1$ for some $L \geq 0$. Let $\widehat{G}_{n}$ denote any approximate NPMLE supported on the minimum axis-aligned bounding box of the data satisfying (4.17). Then there is a function $n(d, \bar{k}, \underline{k}, L)$ such that, for all sample sizes $n$ with $n \geq n(d, \bar{k}, \underline{k}, L)$,

$$
W_{2}^{2}\left(G^{*}, \widehat{G}_{n}\right) \lesssim_{d, \bar{k}} \frac{1}{\log n}
$$

with probability at least $1-\frac{4 d}{n^{8}}$.
Theorem 4.10 (proved in Section 4.6.4) upper bounds the rate of convergence under the Wasserstein distance by the extremely slow logarithmic rate $\frac{1}{\log n}$. It is well known that the smoothness of the Gaussian errors makes the deconvolution more difficult; in fact, the logarithmic rate is minimax optimal (Dedecker \& Michel, 2013).

Remark 4.11. (On Theorem 4.10) To our knowledge, Theorem 4.10 is novel, and the rate of convergence for the NPMLE under a Wasserstein distance has not been studied previously. The structure of the proof follows the proof of Theorem 2 of Nguyen (2013). To deal with the fact that $\widehat{G}_{n}$ and $G^{*}$ are typically singular, we convolve each with a distribution with full support but low variance. Compared to our results on the density estimation and denoising problems, Theorem 4.10 makes additional assumptions on the problem structure, specifically that the covariance matrices are diagonal and that $G^{*}$ is compactly supported. Many practical applications satisfy the diagonal covariances restriction, including both of our applications in Section 4.4.

A common feature to our results on density estimation and denoising have been that the NPMLE adapts to the complexity of $G^{*}$. It is reasonable to conjecture, then, that
in the deconvolution problem, $\widehat{G}_{n}$ will also enjoy some adaptation properties under the Wasserstein distance. We close this section with a sharper result on the Wasserstein rate in the special case where the observations are drawn from Gaussian distributions with common mean $\mu \in \mathbb{R}^{d}$.

Theorem 4.12. Suppose $X_{i} \stackrel{\text { ind }}{\sim} \mathcal{N}\left(\mu, \Sigma_{i}\right)$, i.e. $X_{i} \stackrel{\text { ind }}{\sim} f_{G^{*}, \Sigma_{i}}$ where $G^{*}=\delta_{\mu}$ and $\underline{k} I_{d} \preceq$ $\Sigma_{i} \preceq \bar{k} I_{d}$ for all $i=1, \ldots, n$. Let $\widehat{G}_{n}$ denote any approximate NPMLE satisfying (4.17) and supported on $\mathbb{B}_{\kappa r}(\bar{X})$ where $\kappa=\bar{k} / \underline{k}, r=\max _{i}\left\|X_{i}-\bar{X}\right\|_{2}$, and $\bar{X}=\frac{1}{n} \sum_{i=1}^{n} X_{i}$. Then

$$
W_{2}\left(\widehat{G}_{n}, G^{*}\right) \lesssim_{d, \bar{k}, \underline{k}} t^{3 / 2} \frac{(\log n)^{(d+3) / 4}}{n^{1 / 4}}
$$

with probability at least $1-3 n^{-t^{2}}$ for all $t \geq 1$.
If the approximate NPMLE $\widehat{G}_{n}$ of Theorem 4.12 is selected according to the strategy described in Section 4.2.2, then by Corollary 4.3 part (iii) its support will be contained within the ball $\mathbb{B}_{\kappa r}(\bar{X})$. This additional assumption on the support of the approximate NPMLE is needed to have some control over the moments of $\widehat{G}_{n}$.

Up to logarithmic factors, the $n^{1 / 4}$-rate in Theorem 4.12 agrees with Corollary 4.1 of Ho and Nguyen (2016) for the MLE of an overfitted mixture. Specifically, their result compared the MLE of $k$-component finite Gaussian mixture to a true mixing distribution $G^{*}$ with $k^{*}<k$ components. Wu and Yang (2020) and Doss et al. (2020) also derived the $n^{1 / 4}$-rate for a different estimator under a different Wasserstein metric. All of these previous results were restricted to the homoscedastic setting. In our setting, $k^{*}=1$ and $k=\hat{k}$ is the datadependent order of the NPMLE. The best known bound on $\hat{k}$ is logarithmic in $n$ (Polyanskiy \& Wu, 2020), whereas Ho and Nguyen (2016) required $k$ to be fixed as $n \rightarrow \infty$. When $k^{*}$ is known, a faster $n^{1 / 2}$-rate is possible (Heinrich \& Kahn, 2018) and is achieved by the MLE in a well-specified finite mixture model, i.e. setting $k=k^{*}$ (Ho \& Nguyen, 2016).

While the slower $n^{1 / 4}$-rate appears to be the price of flexibility of the NPMLE, Theorem 4.12 establishes that the NPMLE indeed adapts to structure in $G^{*}$. Our analysis is greatly simplified by the assumption $G^{*}=\delta_{\mu}$, since there is only one coupling between $\widehat{G}_{n}$ and $G^{*}$. We leave for future work the important question of the extent to which $\widehat{G}_{n}$ adapts to more general distributions $G^{*}$.

### 4.4 Applications

### 4.4.1 Color-magnitude diagram

In this section, we continue our discussion of denoising the color-magnitude diagram (CMD) from Section 4.1. Our modeling strategy is closely related to the work of Anderson et al. (2018). To compare our method to extreme deconvolution (Bovy et al., 2011), we use the same stellar sample, relaxing only their assumption that the prior $G^{*}$ is a mixture of

Gaussians; by contrast, we allow $G^{*}$ to be an arbitrary probability measure. Specifically, we assume that after a suitable transformation of the color and magnitude measurements, the pair, denoted $X_{i} \in \mathbb{R}^{2}$, come from a two-dimensional Gaussian mixture $f_{G^{*}, \Sigma_{i}}$ with known covariance $\Sigma_{i}$.


Figure 4.5: Initial grid (left) of $m=10^{4}$ support points and estimated prior $\widehat{G}_{n}$ (right) where the area of each atom is proportional to its weight.

Figure 4.1 in Section 4.1 shows the plot of the observed data $X_{i}$ (left) and estimated posterior means $\hat{\theta}_{i}$ (right), the latter constituting the denoised CMD. Contrasting our CMD with theirs (Anderson et al., 2018, Figure 7), which we do not depict here, it appears that ours performs more shrinkage overall. Our CMD has rather sharp tails in the bottom of the plot (i.e. the main sequence) and the top right (i.e. the tip of the red-giant branch) as well as a definitive cluster in the center-right (i.e. the red clump).

There are also important differences between the NPMLE and extreme deconvolution in the estimated prior $\widehat{G}_{n}$. Figure 4.5 shows the initial and final iterates in the computation of the NPMLE. It is clear that we are using a discrete distribution to model the prior, and since all of the covariance matrices $\Sigma_{i}$ are diagonal, by Corollary 4.3 we have restricted the support points to lie in the minimum axis-aligned bounding box of the data. By contrast, extreme deconvolution models the prior as itself a Gaussian mixture, so the estimated prior (Anderson et al., 2018, Figure 4) actually is supported on all of $\mathbb{R}^{2}$.

### 4.4.2 Chemical abundance ratios

Our second data set is taken from the Apache Point Observatory Galactic Evolution Experiment survey (APOGEE); see Majewski et al. (2017), Abolfathi et al. (2018). We examine chemical abundance ratios for the red clump (RC) stars given in the DR14 APOGEE red clump catalog; see Ratcliffe et al. (2020) where this data set has been studied. Following the pre-processing in Ratcliffe et al. (2020) to remove the outliers with anomalous abundance measurements, the data set contains $n=27,238$ observations. We pick $d=2$ features from the 19 dimensions, namely, $[\mathrm{Si} / \mathrm{Fe}]-[\mathrm{Mg} / \mathrm{Fe}]$.

In Figure 4.6 we plot the observed data (top left) and estimated posterior means using Gaussian denoising under the estimated prior $\widehat{G}_{n}$ (top right). The initial grid (bottom left) of $m=10^{4}$ support points and estimated prior $\widehat{G}_{n}$ (bottom right), where the area of each atom is proportional to its weight, is also provided. The denoised data reveals a very interesting structure - it shows that the variables $[\mathrm{Si} / \mathrm{Fe}]$ and $[\mathrm{Mg} / \mathrm{Fe}]$ are strongly correlated, especially, the observations for the upper right cluster of stars could be lying on one dimensional manifold; something that is not at all visible when plotting the original data.

### 4.5 Concluding remarks

In this chapter we study the NPMLE $\widehat{G}_{n}$ as an estimator of a prior distribution $G^{*}$ in the presence of multivariate, heteroscedastic measurement errors. We resolve a number of basic questions on the existence, uniqueness, discreteness, and support of the NPMLE, where in several cases the answers differ significantly from the traditional univariate, homoscedastic setting. Our analysis identifies a dual mixture density $\widehat{\psi}_{n}$ with Gaussian $\mathcal{N}\left(X_{i}, \Sigma_{i}\right)$ components at each observation, whose modes contain the atoms of the NPMLE. Our characterization implies that the NPMLE is supported on the ridgeline manifold $\mathcal{M}$, which is a compact subset of $\mathbb{R}^{d}$ defined in terms of the observations $\left(X_{i}\right)_{i=1}^{n}$ and corresponding covariance matrices $\left(\Sigma_{i}\right)_{i=1}^{n}$. This support reduction allows us to approximate the NPMLE by a finite-dimensional convex optimization over the mixing proportions, and we develop a novel approach to bounding the discretization error, justifying the gridding scheme proposed by Koenker and Mizera (2014). Our real data applications show that this approach is viable for practical astronomy problems. Our theoretical results in Section 4.3 provide strong justification for using the NPMLE in a variety of contexts - estimating the prior, marginal densities, and oracle posterior means.

We conclude by outlining some possible future research directions. Computation remains an important barrier for large-scale applications. Specifically, for problems with a large number of samples, e.g. $n \gg 10^{6}$, some additional forms of approximation are warranted, such as stochastic optimization or binning via coresets (see also Ritchie and Murray (2019) on approaches for scaling Extreme Deconvolution to large datasets). Further, our result on the discretization error suggests that discretization becomes infeasible in moderate-dimensions,


Figure 4.6: Top: Observed data (left) and estimated posterior means (right); Bottom: Initial grid (left) of $m=10^{4}$ support points and estimated prior $\widehat{G}_{n}$ (right) where the area of each atom is proportional to its weight.
where the number of atoms needs to grow roughly like a polynomial in the number of dimensions, e.g. $m=O\left(\delta^{-d}\right)$. This limitation, which is common to many forms of discretization across applied mathematics, highlights the need for grid-free methods for computation of
the NPMLE in high-dimensions. The connection to entropic-regularized optimal transport established by Rigollet and Weed (2018) represents one possible direction for grid-free methods.

Next, while our framework allows the prior $G^{*}$ to be arbitrary, the underlying assumptionthat the means $\left(\theta_{i}^{*}\right)$ are identically distributed - can sometimes be difficult to justify for heteroscedastic observations. The IID assumption reflects the belief that the observation covariance $\Sigma_{i}$ is uninformative for the corresponding mean $\theta_{i}^{*}$. This assumption led to reasonable results in our applications but may be problematic in other settings. In the univariate, heteroscedastic case, Weinstein et al. (2018) proposed grouping observations with similar variances and applying a spherically symmetric estimator separately within each group. Their approach is capable of capturing dependence between $\theta_{i}$ and $\sigma_{i}^{2}$, at the expense of not sharing information across groups. Furthermore, to our knowledge, the grouping approach has not been extended to multivariate settings where binning the set of covariance matrices is more difficult. Thus, in multivariate settings there remains the important problem of how to model the relationship between $\theta_{i}$ and $\Sigma_{i}$.

Finally, there remain a number of open statistical questions for future work. Our analysis of the denoising problem focuses on estimating the posterior mean based on the unknown prior $G^{*}$, but there are numerous inferential goals one could target with an approximate prior. The analyst might summarize the empirical posteriors using a different functional, such as the posterior median or the posterior mean of some transformed parameter. This question warrants a more general analysis evaluating the quality of the empirical posterior distributions for the true, unknown posteriors.

### 4.6 Proofs

### 4.6.1 Proofs of results in sections 4.2 and 4.2.2

### 4.6.1.1 Proof of Lemma 4.1

The following uses similar techniques as Section 5.2 of Lindsay (1995), which contains a subset of our result in the homoscedastic case.
Proof of Lemma 4.1. By convexity, the first-order optimality condition for $\widehat{G}_{n}$ is

$$
D\left(\widehat{G}_{n}, G\right) \leq 0 \text { for all } G \in \mathcal{P}\left(\mathbb{R}^{d}\right)
$$

where

$$
\begin{aligned}
D\left(\widehat{G}_{n}, G\right) & :=\lim _{\alpha \downarrow 0} \frac{\frac{1}{n} \sum_{i=1}^{n}\left[\log f_{(1-\alpha) \widehat{G}_{n}+\alpha G, \Sigma_{i}}\left(X_{i}\right)-\log f_{\widehat{G}_{n}, \Sigma_{i}}\left(X_{i}\right)\right]}{\alpha} \\
& =\frac{1}{n} \sum_{i=1}^{n} \frac{1}{f_{\widehat{G}_{n}, \Sigma_{i}}\left(X_{i}\right)}\left(f_{G, \Sigma_{i}}\left(X_{i}\right)-f_{\widehat{G}_{n}, \Sigma_{i}}\left(X_{i}\right)\right)=\frac{1}{n} \sum_{i=1}^{n} \frac{f_{G, \Sigma_{i}}\left(X_{i}\right)}{f_{\widehat{G}_{n}, \Sigma_{i}}\left(X_{i}\right)}-1
\end{aligned}
$$

When $G=\delta_{\vartheta}$ is a point mass we write $D\left(\widehat{G}_{n}, \vartheta\right)$ instead of $D\left(\widehat{G}_{n}, G\right)$. It suffices to check $D\left(\widehat{G}_{n}, \vartheta\right) \leq 0$ for all $\vartheta \in \mathbb{R}^{d}$ because $D\left(\widehat{G}_{n}, G\right)=\int D\left(\widehat{G}_{n}, \vartheta\right) \mathrm{d} G[\vartheta]$.

For the first part of the Lemma, define $\mathcal{C}:=\left\{\left(f_{G, \Sigma_{i}}\left(X_{i}\right)\right)_{i=1}^{n}: G \in \mathcal{P}\left(\mathbb{R}^{d}\right)\right\} \cup\{0\}$. Observe that

$$
\mathcal{C}=\operatorname{conv}(\mathcal{L}), \text { where } \mathcal{L}:=\left\{\left(\varphi_{\Sigma_{i}}\left(X_{i}-\vartheta\right)\right)_{i=1}^{n}: \vartheta \in \mathbb{R}^{d}\right\} \cup\{0\}
$$

Since $\vartheta \mapsto\left(\varphi_{\Sigma_{i}}\left(X_{i}-\vartheta\right)\right)_{i=1}^{n}$ is continuous and $\lim _{\|\vartheta\|_{2} \rightarrow \infty}\left(\varphi_{\Sigma_{i}}\left(X_{i}-\vartheta\right)\right)_{i=1}^{n}=0$, the set $\mathcal{L}$ is closed, and by boundedness of the Gaussian likelihood, $\mathcal{L}$ is compact. Hence $\mathcal{C} \subset$ $\mathbb{R}^{n}$ is convex and compact, and $f(L)=\frac{1}{n} \sum_{i=1}^{n} \log L_{i}$ is strictly concave over $\mathcal{C}$. Thus, $f$ attains its maximum at a unique (nonzero) boundary point $\hat{L} \in \partial \mathcal{C}$. Observe $\mathcal{C}=$ $\operatorname{conv}\left\{\left(\varphi_{\Sigma_{i}}\left(X_{i}-\vartheta\right)\right)_{i=1}^{n}: \vartheta \in \mathbb{R}^{d}\right\}$ : by Carathéodory's theorem, any boundary point $\hat{L} \in \partial \mathcal{C}$ can be written as $\hat{L}_{i}=\sum_{j=1}^{\hat{k}} \hat{w}_{j} \varphi_{\Sigma_{i}}\left(X_{i}-\hat{a}_{j}\right)$ for some $\hat{k} \leq n$.

Suppose $B \subset \operatorname{supp}\left(\widehat{G}_{n}\right)$ is contained in the support of the NPMLE. Given $\widehat{G}_{n}(B)>0$, define a new probability measure $\widehat{G}_{n}^{B}$ via $\widehat{G}_{n}^{B}(A):=\frac{\widehat{G}_{n}(A \cap B)}{\widehat{G}_{n}(B)}$. Since $\widehat{G}_{n}=\alpha_{0} \widehat{G}_{n}^{B}+\left(1-\alpha_{0}\right) \widehat{G}_{n}^{B^{c}}$ for $\alpha_{0}=\widehat{G}_{n}(B)$, the mixture

$$
G_{\alpha}=(1-\alpha) \widehat{G}_{n}+\alpha \widehat{G}_{n}^{B}
$$

remains a valid probability measure for $\alpha \geq-\frac{\alpha_{0}}{1-\alpha_{0}}$. Since $\alpha=0$ maximizes the log-likelihood of $G_{\alpha}$ over a range $\alpha \in\left[-\frac{\alpha_{0}}{1-\alpha_{0}}, 1\right]$ including both negative and positive values, the derivative of the $\log$-likelihood is zero at $\alpha=0$, i.e.

$$
0=D\left(\widehat{G}_{n}, \widehat{G}_{n}^{B}\right)=\int D\left(\widehat{G}_{n}, \vartheta\right) \mathrm{d} \widehat{G}_{n}^{B}[\vartheta]
$$

so $\widehat{G}_{n}^{B}(\mathcal{Z})=1$ for all $B \subset \operatorname{supp}\left(\widehat{G}_{n}\right)$ such that $\widehat{G}_{n}(B)>0$. This implies $\widehat{G}_{n}(B \cap \mathcal{Z})=\widehat{G}_{n}(B)$ for all measurable $B$, from which we may conclude $\mathcal{Z} \supseteq \operatorname{supp}\left(\widehat{G}_{n}\right)$. Finally, observe that

$$
D\left(\widehat{G}_{n}, \vartheta\right)=\frac{1}{n} \sum_{i=1}^{n} \hat{L}_{i}^{-1} \varphi_{\Sigma_{i}}\left(X_{i}-\vartheta\right)-1=\left(\frac{1}{n} \sum_{i=1}^{n} \hat{L}_{i}^{-1}\right) \widehat{\psi}_{n}(\vartheta)-1
$$

so $D\left(\widehat{G}_{n}, \vartheta\right) \leq 0$ is equivalent to $\widehat{\psi}_{n}(\vartheta) \leq\left(\frac{1}{n} \sum_{i=1}^{n} \hat{L}_{i}^{-1}\right)^{-1}$. This proves the last statement of the Lemma, that $\mathcal{Z}$ is equal to the set of global maximizers of $\widehat{\psi}_{n}$.

### 4.6.1.2 Proof of Lemma 4.2

Proof of Lemma 4.2. By Lemma 4.4, the fitted values $\hat{L}_{1}=\hat{L}_{2}=\hat{L}_{3}$ are equal. By Lemma 4.1, the atoms of $\widehat{G}_{n}$ occur at the global modes of $\widehat{\psi}_{n}=f_{H, \sigma^{2} I_{2}}$, where $H=\frac{1}{3} \sum_{i=1}^{3} \delta_{X_{i}}$. Since $\hat{L}_{1}=\hat{L}_{2}=\hat{L}_{3}$, the fitted values are also equal to the global maximum of $\hat{\psi}_{n}$, i.e.

$$
\hat{L}_{i}=\max _{x} f_{H, \sigma^{2} I_{2}}(x)=\frac{2^{2 / 3} \log 2}{3 \pi}
$$

for each $i=1,2,3$. Note that $\hat{L}_{i}=f_{\delta_{0}, \sigma^{2} I_{2}}\left(X_{i}\right)$ for all $X_{i}$, so $\widehat{G}_{n}=\delta_{0}$ is an NPMLE. Now let $\widehat{G}_{n}^{\prime}=\frac{1}{3} \sum_{i=1}^{n} \delta_{X_{i} / 2}$. It suffices to check the fitted values of $\widehat{G}_{n}^{\prime}$ at the observations. For $i=1$,

$$
\begin{aligned}
f_{\widehat{G}_{n}^{\prime}, \sigma^{2} I_{2}}\left(X_{1}\right) & =\frac{1}{3} \sum_{i=1}^{3} \varphi_{\sigma^{2} I_{2}}\left(X_{1}-X_{i} / 2\right) \\
& =\frac{4 \log 2}{9 \pi}\left(2^{-(4 / 3)(1 / 4)}+2^{-(4 / 3)(7 / 4)}+2^{-(4 / 3)(7 / 4)}\right)=\frac{2^{2 / 3} \log 2}{3 \pi}=\hat{L}_{1} .
\end{aligned}
$$

Similarly, for $i=2$,

$$
\begin{aligned}
f_{\widehat{G}_{n}^{\prime}, \sigma^{2} I_{2}}\left(X_{2}\right) & =\frac{1}{3} \sum_{i=1}^{3} \varphi_{\sigma^{2} I_{2}}\left(X_{2}-X_{i} / 2\right) \\
& =\frac{4 \log 2}{9 \pi}\left(2^{-(4 / 3)(7 / 4)}+2^{-(4 / 3)(1 / 4)}+2^{-(4 / 3)(7 / 4)}\right)=\frac{2^{2 / 3} \log 2}{3 \pi}=\hat{L}_{2}
\end{aligned}
$$

and, for $i=3$,

$$
\begin{aligned}
f_{\widehat{G}_{n}^{\prime}, \sigma^{2} I_{2}}\left(X_{3}\right) & =\frac{1}{3} \sum_{i=1}^{3} \varphi_{\sigma^{2} I_{2}}\left(X_{3}-X_{i} / 2\right) \\
& =\frac{4 \log 2}{9 \pi}\left(2^{-(4 / 3)(7 / 4)}+2^{-(4 / 3)(7 / 4)}+2^{-(4 / 3)(1 / 4)}\right)=\frac{2^{2 / 3} \log 2}{3 \pi}=\hat{L}_{3}
\end{aligned}
$$

This verifies that $\widehat{G}_{n}^{\prime}=\frac{1}{3} \sum_{i=1}^{n} \delta_{X_{i} / 2}$ is also an NPMLE, so every convex combination $\alpha \widehat{G}_{n}+$ $(1-\alpha) \widehat{G}_{n}^{\prime}$ is an NPMLE.

### 4.6.1.3 Proof of Corollary 4.3

Proof of Corollary 4.3. We have already observed that $\mathcal{Z} \subset \mathcal{M}$ (Ray \& Lindsay, 2005). Observe that $\mathcal{M}$ is compact as it is the continuous image of the simplex, a compact set. Since any real-analytic function has a finite number of zeros, $\mathcal{Z}$ is finite. Hence any NPMLE $\widehat{G}_{n}$ is discrete with a finite number of atoms.

In the proportional covariances case $\Sigma_{i}=c_{i} \Sigma$, we have

$$
\begin{aligned}
x^{*}(\alpha) & =\left(\sum_{i=1}^{n} \alpha_{i} \Sigma_{i}^{-1}\right)^{-1} \sum_{i=1}^{n} \alpha_{i} \Sigma_{i}^{-1} X_{i} \\
& =\sum_{i=1}^{n} \frac{\alpha_{i} / c_{i}}{\sum_{\iota=1}^{n} \alpha_{\iota} / c_{\iota}} X_{i}
\end{aligned}
$$

As $\alpha$ ranges over the simplex, so does $\left(\frac{\alpha_{i} / c_{i}}{\sum_{i=1}^{n} \alpha_{\iota} / c_{l}}\right)_{i=1}^{n}$. Thus $\mathcal{M}=\operatorname{conv}\left(\left\{X_{1}, \ldots, X_{n}\right\}\right)$, proving (i). If each $\Sigma_{i}$ is diagonal, letting $x_{j}^{*}(\alpha)$ denote the $j^{\text {th }}$ coordinate of $x^{*}(\alpha) \in \mathbb{R}^{d}$,

$$
x_{j}^{*}(\alpha)=\sum_{i=1}^{n} \frac{\alpha_{i}\left(\sum_{i}\right)_{j j}}{\sum_{i^{\prime}=1}^{n} \alpha_{i^{\prime}}\left(\sum_{i^{\prime}}\right)_{j j}} X_{i j} \in\left[\min _{i \in\{1, \ldots, n\}} X_{i j}, \max _{i \in\{1, \ldots, n\}} X_{i j}\right],
$$

proving (ii). For (iii), using concavity of the minimum eigenvalue,

$$
\begin{aligned}
\left\|x^{*}(\alpha)-x\right\|_{2} & =\left\|\left(\sum_{i=1}^{n} \alpha_{i} \Sigma_{i}^{-1}\right)^{-1} \sum_{i=1}^{n} \alpha_{i} \Sigma_{i}^{-1}\left(X_{i}-x\right)\right\|_{2} \\
& \leq\left\|\left(\sum_{i=1}^{n} \alpha_{i} \Sigma_{i}^{-1}\right)^{-1}\right\|_{2}\left\|\sum_{i=1}^{n} \alpha_{i} \Sigma_{i}^{-1}\left(X_{i}-x\right)\right\|_{2} \\
& \leq\left(\sum_{i=1}^{n} \alpha_{i} \bar{k}^{-1}\right)^{-1} \sum_{i=1}^{n} \alpha_{i} \underline{k}^{-1}\left\|X_{i}-x\right\|_{2} \leq \kappa r
\end{aligned}
$$

so $\mathcal{M} \subseteq \mathbb{B}_{\kappa r}(x)$.

### 4.6.1.4 Proof of Lemma 4.4

Proof of Lemma 4.4. By the change of variables formula,

$$
\begin{aligned}
f_{T_{\#} G, \Sigma_{i}^{\prime}}\left(X_{i}^{\prime}\right) & =\int \varphi_{U_{0} \Sigma_{i} U_{0}^{\top}}\left(U_{0} X_{i}+x_{0}-\theta\right) \mathrm{d} T_{\#} G(\theta) \\
& =\int \varphi_{U_{0} \Sigma_{i} U_{0}^{\top}}\left(U_{0} X_{i}+x_{0}-T(\theta)\right) \mathrm{d} G(\theta) \\
& =\int \varphi_{\Sigma_{i}}\left(X_{i}-\theta\right) \mathrm{d} G(\theta)=f_{G, \Sigma_{i}}\left(X_{i}\right),
\end{aligned}
$$

completing the proof.

### 4.6.1.5 Proof of Proposition 4.5

Proof of Proposition 4.5. Write $\widehat{G}_{n}=\sum_{j=1}^{\hat{k}} \hat{w}_{j} \delta_{\hat{a}_{j}}$, and for each $j \in[\hat{k}]$, let $C_{j} \in \mathcal{H}$ such that $\hat{a}_{j} \in C_{j}$. Next, define a positive measure $H_{j}$ supported on the corners of $C_{j}$ such that $H_{j}\left(C_{j}\right)=\hat{w}_{j}$ and

$$
\begin{equation*}
\int_{C_{j}} u \mathrm{~d} H_{j}(u)=\hat{w}_{j} \hat{a}_{j}=\int_{C_{j}} u \mathrm{~d} \widehat{G}_{n}^{j}(u) \tag{4.20}
\end{equation*}
$$

where $\widehat{G}_{n}^{j}:=\hat{w}_{j} \delta_{\hat{a}_{j}}$. Now fix $u \in C_{j}$ and $i \in[n]$, and let $x_{j}=\Sigma_{i}^{-1 / 2}\left(X_{i}-\hat{a}_{j}\right)$ and $t=$ $\Sigma_{i}^{-1 / 2}\left(u-\hat{a}_{j}\right)$. By the moment identity (4.20) and by Jiang and Zhang (2009, A.27),

$$
\begin{aligned}
\int_{C_{j}} \varphi_{\Sigma_{i}}\left(X_{i}-u\right) \mathrm{d} \widehat{G}_{n}^{j}(u) & -\int_{C_{j}} \varphi_{\Sigma_{i}}\left(X_{i}-u\right) \mathrm{d} H_{j}(u) \\
& \leq \int_{C_{j}}\left\langle x_{j}, t\right\rangle^{2} \varphi_{\Sigma_{i}}\left(X_{i}-u\right) \mathrm{d} \widehat{G}_{n}^{j}(u)+\int_{C_{j}}\left(e^{\|t\|_{2}^{2} / 2}-1\right) \varphi_{\Sigma_{i}}\left(X_{i}-u\right) \mathrm{d} H_{j}(u) \\
& \leq \underline{k}^{-2} D^{2} d \delta^{2} \int_{C_{j}} \varphi_{\Sigma_{i}}\left(X_{i}-u\right) \mathrm{d} \widehat{G}_{n}^{j}(u)+\left(e^{\underline{k} d \delta^{2} / 2}-1\right) \int_{C_{j}} \varphi_{\Sigma_{i}}\left(X_{i}-u\right) \mathrm{d} H_{j}(u) .
\end{aligned}
$$

Let $H=\sum_{j=1}^{\hat{k}} H_{j}$. Summing the above inequality over $j$,

$$
f_{\widehat{G}_{n}, \Sigma_{i}}\left(X_{i}\right)-f_{H, \Sigma_{i}}\left(X_{i}\right) \leq \underline{k}^{-2} D^{2} d \delta^{2} f_{\widehat{G}_{n}, \Sigma_{i}}\left(X_{i}\right)+\left(e^{\underline{k} d \delta^{2} / 2}-1\right) f_{H, \Sigma_{i}}\left(X_{i}\right)
$$

Since $H$ is supported on $\mathcal{A}$, by optimality of $\widehat{G}_{n}^{\mathcal{A}}$,

$$
\prod_{i=1}^{n} f_{\widehat{G}_{n}^{A}, \Sigma_{i}}\left(X_{i}\right) \geq \prod_{i=1}^{n} f_{H, \Sigma_{i}}\left(X_{i}\right)
$$

Combining our findings,

$$
\prod_{i=1}^{n} f_{\widehat{G}_{n}^{A}, \Sigma_{i}}\left(X_{i}\right) \geq e^{-n \underline{k}^{-2} d \delta^{2} / 2}\left(1-\underline{k}^{-2} D^{2} d \delta^{2}\right)^{n} \prod_{i=1}^{n} f_{\widehat{G}_{n, \Sigma_{i}}}\left(X_{i}\right)
$$

Using the elementary inequality $1-x \geq e^{-2 x}$ for $x \leq 3 / 4$, we obtain

$$
\prod_{i=1}^{n} f_{\widehat{G}_{n}^{A}, \Sigma_{i}}\left(X_{i}\right) \geq \exp \left(-n \underline{k}^{-2} d \delta^{2} / 2-2 n \underline{k}^{-2} D^{2} d \delta^{2}\right) \prod_{i=1}^{n} f_{\widehat{G}_{n}, \Sigma_{i}}\left(X_{i}\right)
$$

for $\delta \leq \sqrt{\frac{3}{4 d}} k D^{-1}$.

### 4.6.2 Proof of Theorem 4.6

The following notation will be used throughout this section:

1. $\mathbb{B}_{r}(x)=\left\{y \in \mathbb{R}^{d}:\|x-y\|_{2} \leq r\right\}$ denotes a closed ball in $\mathbb{R}^{d}$.
2. For a positive integer $m$, let $[m]=\{1, \ldots, m\}$.
3. Given a pseudo-metric space $(M, \rho)$ and $\varepsilon>0$, let $N(\varepsilon, M, \rho)$ denote the $\varepsilon$-covering number, i.e. the smallest positive integer $N$ such that there exist $x_{1}, \ldots, x_{N} \in M$ such that

$$
M \subset \bigcup_{i=1}^{N}\left\{y: \rho\left(y, x_{i}\right) \leq \varepsilon\right\} .
$$

Any such a set $\left\{x_{i}\right\}_{i=1}^{N}$ is known as an $\varepsilon$-net or $\varepsilon$-cover of $M$ under the pseudo-metric $\rho$. When $M$ is a subset of Euclidean space we write $N(\varepsilon, M)$ instead of $N\left(\varepsilon, M,\|\cdot\|_{2}\right)$.
4. We use the shorthand $f_{G, \bullet}=\left(f_{G, \Sigma_{i}}\right)_{i=1}^{n}$, the matrices $\Sigma_{1}, \ldots, \Sigma_{n}$ being viewed as fixed. Let

$$
\mathbb{F}=\left\{f_{G, \bullet}: G \in \mathcal{P}\left(\mathbb{R}^{d}\right)\right\}
$$

5. For $S \subset \mathbb{R}^{d}$ and $M>0, S^{M}$ denotes the $M$-enlargement $S^{M}=\left\{x \in \mathbb{R}^{d}: \mathfrak{d}_{S}(x) \leq M\right\}$.
6. Define the semi-norm

$$
\left\|f_{G, \bullet}-f_{H, \bullet}\right\|_{\infty, S^{M}}:=\max _{1 \leq i \leq n} \sup _{x \in S^{M}}\left|f_{G, \Sigma_{i}}(x)-f_{H, \Sigma_{i}}(x)\right| .
$$

Similarly, define

$$
\left\|f_{G, \bullet}-f_{H, \bullet}\right\|_{\nabla, S^{M}}:=\max _{1 \leq i \leq n} \sup _{x \in S^{M}}\left|\nabla f_{G, \Sigma_{i}}(x)-\nabla f_{H, \Sigma_{i}}(x)\right|
$$

Our proof generalizes and builds upon prior techniques for analyzing the Hellinger accuracy of the NPMLE (Jiang, 2020; Saha \& Guntuboyina, 2020a; Zhang, 2009). The basic structure of our argument is to recognize, given the approximation (4.14) in the likelihood, that we may trivially rewrite the large deviation probability for the NPMLE as a joint probability
$\mathbb{P}\left(\bar{h}\left(f_{\widehat{G}_{n}, \bullet}, f_{G^{*}, \bullet}\right) \gtrsim_{d, \bar{k}, \underline{k}} t \varepsilon_{n}\right)=\mathbb{P}\left(\bar{h}\left(f_{\widehat{G}_{n}, \bullet}, f_{G^{*}, \bullet}\right) \gtrsim_{d, \bar{k}, \underline{k}} t \varepsilon_{n}, \prod_{i=1}^{n} \frac{f_{\widehat{G}_{n}, \Sigma_{i}}\left(X_{i}\right)}{f_{G^{*}, \Sigma_{i}}\left(X_{i}\right)} \geq \exp \left(-c_{d, \bar{k}, \underline{k}} n \varepsilon_{n}^{2}\right)\right)$.
If $\widehat{G}_{n}$ were a fixed probability measure $G_{0}$ such that $\bar{h}\left(f_{G_{0}, \bullet}, f_{G^{*}, \bullet \bullet}\right) \gtrsim_{d, \bar{k}, \underline{k}} t \varepsilon_{n}$, the right-hand side of the last display similarly simplifies as

$$
\begin{aligned}
& \mathbb{P}\left(\bar{h}\left(f_{G_{0}, \bullet}, f_{G^{*}, \bullet}\right) \gtrsim_{d, \bar{k}, \underline{,},} t \varepsilon_{n}, \prod_{i=1}^{n} \frac{f_{G_{0}, \Sigma_{i}}\left(X_{i}\right)}{f_{G^{*}, \Sigma_{i}}\left(X_{i}\right)} \geq \exp \left(-c_{d, \bar{k}, \underline{k}} n \varepsilon_{n}^{2}\right)\right) \\
& \quad=\mathbb{P}\left(\prod_{i=1}^{n} \frac{f_{G_{0}, \Sigma_{i}}\left(X_{i}\right)}{f_{G^{*}, \Sigma_{i}}\left(X_{i}\right)} \geq \exp \left(-c_{d, \bar{k}, \underline{k}} n \varepsilon_{n}^{2}\right)\right) .
\end{aligned}
$$

Since $\widehat{G}_{n}$ is not fixed, we first approximate it using a covering argument, and then bound the right-hand side of the previous display using Markov's inequality.

Proof of Theorem 4.6. Suppose for some $\gamma_{n}$ the NPMLE satisfies

$$
\prod_{i=1}^{n} \frac{f_{\widehat{G}_{n}, \Sigma_{i}}\left(X_{i}\right)}{f_{G^{*}, \Sigma_{i}}\left(X_{i}\right)} \geq \exp \left((\beta-\alpha) n \gamma_{n}^{2}\right) \text { for some } 0<\beta<\alpha<1
$$

We bound the probability

$$
\mathbb{P}\left(\bar{h}\left(f_{\widehat{G}_{n}, \bullet}, f_{G^{*}, \bullet}\right) \geq t \gamma_{n}\right)
$$

for $t>1$.
Take $\left\{f_{H_{j}, \bullet}\right\}_{j=1}^{N} \subset \mathbb{F}$ to be an $\eta$-net of $\mathbb{F}$ under $\|\cdot\|_{\infty, S^{M}}$. For each $j$, let $H_{0, j}$ be a distribution satisfying

$$
\left\|f_{H_{0, j}, \bullet}-f_{H_{j}, \bullet}\right\|_{\infty, S^{M}} \leq \eta \text { and } \bar{h}\left(f_{H_{0, j}, \bullet}, f_{G^{*}, \bullet}\right) \geq t \gamma_{n}
$$

and $J=\left\{j \in[N]: H_{0, j}\right.$ exists $\}$. By construction of the $\eta$-net, there is $j^{*} \in[N]$ such that

$$
\left\|f_{H_{j^{*}}, \bullet}-f_{\widehat{G}_{n}, \bullet}\right\|_{\infty, S^{M}} \leq \eta .
$$

On the event $\left\{\bar{h}\left(f_{\widehat{G}_{n}, \bullet}, f_{G^{*}, \bullet}\right) \geq t \gamma_{n}\right\}$, the NPMLE $\widehat{G}_{n}$ acts as a witness that $j^{*} \in J$, so by the triangle inequality

$$
\begin{equation*}
\left\|f_{H_{0, j^{*}}, \bullet}-f_{\widehat{G}_{n}, \bullet}\right\|_{\infty, S^{M}} \leq 2 \eta \tag{4.21}
\end{equation*}
$$

This gives

$$
f_{\widehat{G}_{n}, \Sigma_{i}}(x) \leq \begin{cases}f_{H_{0, j^{*}}, \Sigma_{i}}(x)+2 \eta, & \text { if } x \in S^{M} \\ \frac{1}{\sqrt{(2 \pi)^{d}\left|\Sigma_{i}\right|}}, & \text { otherwise }\end{cases}
$$

Defining $v(x)=\eta 1_{x \in S^{M}}+\eta\left(\frac{M}{\mathfrak{\imath}_{S}(x)}\right)^{d+1} 1_{x \notin S^{M}}$, we have

$$
\begin{equation*}
\exp \left((\beta-\alpha) n t^{2} \gamma_{n}^{2}\right) \leq \max _{j \in J}\left[\prod_{i=1}^{n} \frac{f_{H_{0, j}, \Sigma_{i}}\left(X_{i}\right)+2 v\left(X_{i}\right)}{f_{G^{*}, \Sigma_{i}}\left(X_{i}\right)}\right] \cdot\left[\prod_{i: X_{i} \notin S^{M}} \frac{1}{\sqrt{(2 \pi)^{d}\left|\Sigma_{i}\right|} \cdot 2 v\left(X_{i}\right)}\right] \tag{4.22}
\end{equation*}
$$

on the event $\left\{\bar{h}\left(f_{\widehat{G}_{n}, \bullet}, f_{G^{*}, \bullet}\right) \geq t \gamma_{n}\right\}$. Hence

$$
\begin{align*}
& \mathbb{P}\left(\bar{h}\left(f_{\widehat{G}_{n}, \bullet}, f_{G^{*}, \bullet}\right) \geq t \gamma_{n}\right)  \tag{4.23}\\
& \leq  \tag{4.24}\\
& \quad \mathbb{P}\left(\max _{j \in J} \prod_{i=1}^{n} \frac{f_{H_{0, j}, \Sigma_{i}}\left(X_{i}\right)+2 v\left(X_{i}\right)}{f_{G^{*}, \Sigma_{i}}\left(X_{i}\right)} \geq \exp \left(-\alpha n t^{2} \gamma_{n}^{2}\right)\right)  \tag{4.25}\\
& \quad+\mathbb{P}\left(\prod_{i: X_{i} \notin S^{M}} \frac{1}{\sqrt{(2 \pi)^{d}\left|\Sigma_{i}\right|} \cdot 2 v\left(X_{i}\right)} \geq \exp \left(\beta n t^{2} \gamma_{n}^{2}\right)\right)
\end{align*}
$$

By a union bound and Markov's inequality, the first term (4.24) is bounded by

$$
\begin{equation*}
e^{\alpha n t^{2} \gamma_{n}^{2} / 2} \sum_{j \in J} \mathbb{E} \prod_{i=1}^{n} \sqrt{\frac{f_{H_{0, j}, \Sigma_{i}}\left(X_{i}\right)+2 v\left(X_{i}\right)}{f_{G^{*}, \Sigma_{i}}\left(X_{i}\right)}} \tag{4.26}
\end{equation*}
$$

Writing out the expectation,

$$
\begin{aligned}
\prod_{i=1}^{n} \mathbb{E} \sqrt{\frac{f_{H_{0, j}, \Sigma_{i}}\left(X_{i}\right)+2 v\left(X_{i}\right)}{f_{G^{*}, \Sigma_{i}}\left(X_{i}\right)}} & =\exp \left(\sum_{i=1}^{n} \log \mathbb{E} \sqrt{\frac{f_{H_{0, j}, \Sigma_{i}}\left(X_{i}\right)+2 v\left(X_{i}\right)}{f_{G^{*}, \Sigma_{i}}\left(X_{i}\right)}}\right) \\
& \leq \exp \left(\sum_{i=1}^{n}\left\{\int \sqrt{f_{H_{0, j}, \Sigma_{i}}+2 v} \sqrt{f_{G^{*}, \Sigma_{i}}}-1\right\}\right) \\
& \leq \exp \left(-\frac{n t^{2} \gamma_{n}^{2}}{2}+n \sqrt{2 \int v}\right)
\end{aligned}
$$

Putting together the pieces, the first term (4.24) is bounded by

$$
\begin{align*}
& \mathbb{P}\left(\max _{j \in J} \prod_{i=1}^{n} \frac{f_{H_{0, j}, \Sigma_{i}}\left(X_{i}\right)+2 v\left(X_{i}\right)}{f_{G^{*}, \Sigma_{i}}\left(X_{i}\right)} \geq e^{-\alpha n t^{2} \gamma_{n}^{2}}\right) \\
& \quad \leq \exp \left(-(1-\alpha) \frac{n t^{2} \gamma_{n}^{2}}{2}+\log N+n \sqrt{2 \int v}\right) \tag{4.27}
\end{align*}
$$

For the second term (4.25), observe by Markov's inequality

$$
\begin{aligned}
& \mathbb{P}\left(\prod_{i: X_{i} \notin S^{M}} \frac{1}{\sqrt{(2 \pi)^{d}\left|\Sigma_{i}\right|} \cdot 2 v\left(X_{i}\right)} \geq \exp \left(\beta n t^{2} \gamma_{n}^{2}\right)\right) \\
& \quad \leq \exp \left(-\frac{\beta n t^{2} \gamma_{n}^{2}}{2 \log n}\right) \mathbb{E}\left\{\prod_{i: X_{i} \notin S^{M}}\left|\frac{1}{\sqrt{\left|2 \pi \Sigma_{i}\right|} \cdot 2 v\left(X_{i}\right)}\right|\right\}^{1 / 2 \log n} \\
& \quad=\exp \left(-\frac{\beta n t^{2} \gamma_{n}^{2}}{2 \log n}\right) \mathbb{E}\left\{\prod_{i=1}^{n}\left(\frac{\mathfrak{d}_{S}\left(X_{i}\right)}{\left|2 \pi \Sigma_{i}\right|^{1 /(2 d+2)} \cdot(2 \eta)^{1 /(d+1)} M}\right)^{1_{\left.\mathfrak{o}_{S}\left(X_{i}\right) \geq M\right)}}\right\}^{(d+1) / 2 \log n}
\end{aligned}
$$

To reduce clutter write $a=\frac{1}{\underline{k}^{d /(2 d+2)} \eta^{1 /(d+1)} M}$ and $\lambda=\frac{d+1}{2 \log n}$. The above expectation is further upper bounded by

$$
\begin{aligned}
\mathbb{E}\left\{\prod_{i=1}^{n}\left(a \mathfrak{d}_{S}\left(X_{i}\right)\right)^{1_{\mathfrak{o}_{S}\left(X_{i}\right) \geq M}}\right\}^{\lambda} & =\prod_{i=1}^{n} \mathbb{E}\left(a \mathfrak{d}_{S}\left(X_{i}\right)\right)^{\lambda 1_{\mathfrak{d}_{S}\left(X_{i}\right) \geq M}} \\
& \leq \prod_{i=1}^{n}\left(1+a^{\lambda} \mathbb{E}\left[\mathfrak{d}_{S}\left(X_{i}\right)^{\lambda} 1_{\mathfrak{d}_{S}\left(X_{i}\right) \geq M}\right]\right) \\
& \leq \exp \left(a^{\lambda} \sum_{i=1}^{n} \mathbb{E}\left[\mathfrak{d}_{S}\left(X_{i}\right)^{\lambda} 1_{\mathfrak{d}_{S}\left(X_{i}\right) \geq M}\right]\right) \\
& \leq \exp \left(n a^{\lambda}\left\{C_{d} M^{d+\lambda-2} \bar{k}^{1-d / 2} e^{-M^{2} /(8 \bar{k})}+M^{\lambda}\left(\frac{2 \mu_{q}}{M}\right)^{q}\right\}\right)
\end{aligned}
$$

The last inequality follows from Lemma 4.13. Note we need

$$
\frac{d+1}{2(1 \wedge q)} \leq \log n
$$

to ensure $\lambda \leq 1 \wedge q$. Taking $M \geq \sqrt{8 \bar{k} \log n}$, we have $e^{-M^{2} /(8 \bar{k})} \leq \frac{1}{n}$, so

$$
\mathbb{E}\left\{\prod_{i=1}^{n}\left(a \mathfrak{d}_{S}\left(X_{i}\right)\right)^{1_{\mathfrak{d}_{S}\left(X_{i}\right) \geq M}}\right\}^{\lambda} \leq \exp \left((a M)^{\lambda}\left[C_{d} M^{d-2} \bar{k}^{1-d / 2}+n\left(\frac{2 \mu_{q}}{M}\right)^{q}\right]\right)
$$

Noting $(a M)^{\lambda}=\left(\underline{k}^{d / 2} \eta\right)^{-1 /(2 \log n)}$, choose $\eta=\frac{n^{-2}}{\underline{k}^{d / 2}}$, so $(a M)^{\lambda}=e$. We directly apply Lemma A. 7 of Saha and Guntuboyina (2020b) for the integral

$$
\int v \leq C_{d} \eta \operatorname{Vol}\left(S^{M}\right)
$$

To bound the metric entropy, i.e. $\log N$ where $N$ denotes the size of our $\eta$-net $\left\{f_{H_{j}, \bullet}\right\}_{j=1}^{N} \subset \mathbb{F}$, we apply Lemma 4.16

$$
\log N=\log N\left(\eta, \mathbb{F},\|\cdot\|_{\infty, S^{M}}\right) \leq C_{d} N\left(u,\left(S^{M}\right)^{u}\right)\left(\log \frac{c_{d, \bar{k}, \underline{k}}}{\eta}\right)^{2}
$$

where the scalar $u$ in the above display corresponds to $a$ used in the lemma. Assuming $4 n \geq(2 \pi)^{d / 2}$,

$$
u=\sqrt{-2 \bar{k} \log \left(\left((2 \pi \underline{k})^{d / 2} \frac{\eta}{4}\right)\right.} \geq \sqrt{2 \bar{k} \log n}
$$

Similarly $u \leq \sqrt{6 \bar{k} \log n}$, so

$$
N\left(u,\left(S^{M}\right)^{u}\right) \leq N\left(\sqrt{2 \bar{k} \log n},\left(S^{M}\right)^{\sqrt{6 \bar{k} \log n}}\right) \leq C_{d, \bar{k}} \operatorname{Vol}\left(S^{2 M}\right)(\log n)^{-d / 2}
$$

Combining our findings,

$$
\begin{aligned}
& \mathbb{P}\left(\bar{h}\left(f_{\widehat{G}_{n}, \bullet}, f_{G^{*}, \bullet}\right) \geq t \gamma_{n}\right) \\
& \quad \leq \exp \left(-(1-\alpha) \frac{n t^{2} \gamma_{n}^{2}}{2}+C_{d, \bar{k}, \underline{k}}(\log n)^{d / 2+1} \operatorname{Vol}\left(S^{2 M}\right)+C_{d} \sqrt{\underline{k}^{-d / 2} \operatorname{Vol}\left(S^{M}\right)}\right) \\
& \quad+\exp \left(-\frac{\beta}{\log n} \frac{n t^{2} \gamma_{n}^{2}}{2}+C_{d} M^{d-2} \bar{k}^{1-d / 2}+e n \inf _{q \geq(d+1) /(2 \log n)}\left(\frac{2 \mu_{q}}{M}\right)^{q}\right)
\end{aligned}
$$

for any $t>1$. Absorbing the dependence on $d, \underline{k}$ and $\bar{k}$ into constants, take $\varepsilon_{n}^{2}=\varepsilon_{n}^{2}\left(M, S, G^{*}\right)$ such that

$$
\begin{aligned}
\max & \left\{(\log n)^{d / 2+1} \operatorname{Vol}\left(S^{2 M}\right), \sqrt{\operatorname{Vol}\left(S^{M}\right)}, M^{d-2}, e n \inf _{q \geq(d+1) /(2 \log n)}\left(\frac{2 \mu_{q}}{M}\right)^{q}\right\} \\
& \lesssim_{d, \bar{k}, \underline{,} \underline{n}} n \varepsilon_{n}^{2}\left(M, S, G^{*}\right)
\end{aligned}
$$

If we then take $\gamma_{n}^{2}=\frac{C_{d, \bar{k}, k_{n}^{2}\left(M, S, G^{*}\right)}}{4 \min (1-\alpha, \beta)}$,

$$
\mathbb{P}\left(\bar{h}\left(f_{\widehat{G}_{n}, \bullet}, f_{G^{*}, \bullet}\right) \geq t \gamma_{n}\right) \leq 2 \exp \left(-\frac{(1-\alpha) \wedge \beta}{4 \log n} n t^{2} \gamma_{n}^{2}\right)
$$

This proves (4.15). To prove (4.16), integrate the tail from (4.15),

$$
\begin{aligned}
\mathbb{E} \frac{\bar{h}^{2}\left(f_{\widehat{G}_{n}, \bullet}, f_{G^{*}, \bullet}\right)}{\gamma_{n}^{2}} & \leq 1+\int_{1}^{\infty} \mathbb{P}\left(\frac{\bar{h}^{2}\left(\widehat{G}_{n}, G^{*}\right)}{\gamma_{n}^{2}} \geq s\right) \mathrm{d} s \\
& \leq 1+\int_{1}^{\infty} 4 t n^{-t^{2}} \mathrm{~d} t=1+\frac{2}{n \log n} \leq 3
\end{aligned}
$$

for $n>1$, completing the proof.
We now state and prove the lemmas needed in the proof of Theorem 4.6.
Lemma 4.13. Let $\theta^{*} \sim G^{*}$ and $Z \sim \mathcal{N}\left(0, I_{d}\right)$ independently, and $Y=\theta^{*}+\Sigma^{1 / 2} Z$, where $\underline{k} I_{d} \preceq \Sigma \preceq \bar{k} I_{d}$. Then

$$
\mathbb{E}\left[\mathfrak{d}_{S}(Y)^{\lambda} 1_{\mathfrak{d}_{S}(Y) \geq M}\right] \leq C_{d} M^{d+\lambda-2} \bar{k}^{1-d / 2} e^{-M^{2} /(8 \bar{k})}+M^{\lambda}\left(\frac{2 \mu_{q}}{M}\right)^{q}
$$

for any $\lambda \in(0,1 \wedge q]$, where $\mu_{q}$ is the $q^{\text {th }}$-moment of $\mathfrak{d}_{S}\left(\theta^{*}\right)$ under $\theta^{*} \sim G^{*}$.
Proof. Since distance $\mathfrak{d}_{S}$ is 1-Lipschitz,

$$
\begin{equation*}
\mathbb{E}\left[\mathfrak{d}_{S}(Y)^{\lambda} 1_{\mathfrak{d}_{S}(Y) \geq M}\right] \leq \mathbb{E}\left[\left(2\left\|\Sigma^{1 / 2} Z\right\|_{2}\right)^{\lambda} 1_{2\left\|\Sigma^{1 / 2} Z\right\|_{2} \geq M}\right]+\mathbb{E}\left[\left(2 \mathfrak{d}_{S}\left(\theta^{*}\right)\right)^{\lambda} 1_{2 \mathfrak{d}_{S}\left(\theta^{*}\right) \geq M}\right] \tag{4.28}
\end{equation*}
$$

For the first term on the RHS of (4.28),

$$
\begin{aligned}
\mathbb{E}\left[\left(2\left\|\Sigma^{1 / 2} Z\right\|_{2}\right)^{\lambda} 1_{2\left\|\Sigma^{1 / 2} Z\right\|_{2} \geq M}\right] & \leq M^{\lambda} \mathbb{E}\left[\left(\frac{\left\|\Sigma^{1 / 2} Z\right\|_{2}}{M / 2}\right)^{\lambda} 1_{\left\|\Sigma^{1 / 2} Z\right\|_{2} \geq M / 2}\right] \\
& \leq 2 M^{\lambda-1} \mathbb{E}\left[\left\|\Sigma^{1 / 2} Z\right\|_{2} 1_{\left\|\Sigma^{1 / 2} Z\right\|_{2} \geq M / 2}\right] \\
& \leq 2 M^{\lambda-1} \bar{k}^{1 / 2} \mathbb{E}\left[\|Z\|_{2} 1_{\|Z\|_{2} \geq M /\left(2 \bar{k}^{1 / 2}\right)}\right] \\
& \leq 2 C_{d} M^{\lambda-1} \bar{k}^{1 / 2}\left(\frac{M}{\bar{k}^{1 / 2}}\right)^{d-1} e^{-M^{2} /(8 \bar{k})} \\
& =C_{d} M^{d+\lambda-2} \bar{k}^{1-d / 2} e^{-M^{2} /(8 \bar{k})}
\end{aligned}
$$

The penultimate inequality uses $\left\|\Sigma^{1 / 2} Z\right\|_{2} \leq \bar{k}^{1 / 2}\|Z\|_{2}$, and the last inequality directly uses Lemma A. 6 of Saha and Guntuboyina (2020b).

Since $\lambda<q$, applying Hölder to the second term on the RHS of (4.28) yields

$$
\mathbb{E}\left[\left(2 \mathfrak{d}_{S}\left(\theta^{*}\right)\right)^{\lambda} 1_{2 \mathfrak{d}_{S}\left(\theta^{*}\right) \geq M}\right] \leq M^{\lambda}\left(\frac{2 \mu_{q}}{M}\right)^{q}
$$

Lemma 4.14. (Moment matching, part i) Let $G, H \in \mathcal{P}\left(\mathbb{R}^{d}\right)$. Suppose $A \subset \mathbb{R}^{d}$ is such that

$$
\mathbb{B}_{a}(x) \subseteq A \subseteq \mathbb{B}_{c a}(x)
$$

for some $c \geq 1$, and that

$$
\int_{A} \theta_{1}^{k_{1}} \cdots \theta_{d}^{k_{d}} d G(\theta)=\int_{A} \theta_{1}^{k_{1}} \cdots \theta_{d}^{k_{d}} d H(\theta), \text { for } k_{1}, \ldots, k_{d} \in[2 m+1]
$$

for some $m \geq 1$. Then

$$
\max _{1 \leq i \leq n}\left|f_{G, \Sigma_{i}}(x)-f_{H, \Sigma_{i}}(x)\right| \leq \frac{1}{(2 \pi \underline{k})^{d / 2}}\left(\frac{e c^{2} a^{2}}{2 \bar{k}(m+1)}\right)^{m+1}+\frac{e^{-a^{2} /(2 \bar{k})}}{(2 \pi \underline{k})^{d / 2}}
$$

Proof. For each $i \in[n]$, write

$$
f_{G, \Sigma_{i}}(x)-f_{H, \Sigma_{i}}(x)=\int_{A} \varphi_{\Sigma_{i}}(x-\theta)(d G(\theta)-d H(\theta))+\int_{A^{c}} \varphi_{\Sigma_{i}}(x-\theta)(d G(\theta)-d H(\theta))
$$

On $A^{c},\|x-\theta\|_{2} \geq a$, so

$$
\varphi_{\Sigma_{i}}(x-\theta) \leq \frac{e^{-a^{2} /(2 \bar{k})}}{(2 \pi \underline{k})^{d / 2}}
$$

Write the pdf as $\varphi_{\Sigma_{i}}(z)=P_{i}(z)+R_{i}(z)$ where $P_{i}$ is a polynomial of degree $2 m$ and the remainder $R_{i}$ satisfies

$$
\left|R_{i}(z)\right| \leq(2 \pi \underline{k})^{-d / 2}\left(\frac{e\|z\|_{2}^{2}}{2 \bar{k}(m+1)}\right)^{m+1}
$$

By hypothesis, $\int_{A} P_{i}(x-\theta)(\mathrm{d} G(\theta)-\mathrm{d} H(\theta))$, so

$$
\begin{aligned}
\left|\int_{A} \varphi_{\Sigma_{i}}(x-\theta)(\mathrm{d} G(\theta)-\mathrm{d} H(\theta))\right| & \leq\left|\int_{A} R_{i}(x-\theta)(\mathrm{d} G(\theta)-\mathrm{d} H(\theta))\right| \\
& \leq \frac{1}{(2 \pi \underline{k})^{d / 2}}\left(\frac{e c^{2} a^{2}}{2 \bar{k}(m+1)}\right)^{m+1}
\end{aligned}
$$

completing the proof.
Lemma 4.15. (Moment matching, part ii) For any $G \in \mathcal{P}\left(\mathbb{R}^{d}\right)$, there is a discrete distribution $H$ supported on $S^{a}$ with at most

$$
l:=\left(2\left\lfloor 13.5 a^{2} / \bar{k}\right\rfloor+2\right)^{d} N\left(a, S^{a}\right)+1
$$

atoms such that

$$
\left\|f_{G, \bullet}-f_{H, \bullet}\right\|_{\infty, S^{a}} \leq\left(1+\frac{1}{\sqrt{2 \pi}}\right)(2 \pi \underline{k})^{-d / 2} e^{-a^{2} /(2 \bar{k})}
$$

Proof. The idea is to choose $H$ to match moments, and then apply the previous lemma. The proof is identical to Lemma D. 3 of Saha and Guntuboyina (2020b), except that we take $m:=\left\lfloor\frac{27 a^{2}}{2 \bar{k}}\right\rfloor$.

Lemma 4.16. There exists positive constants $C_{d}$ and $c_{d, \bar{k}, \underline{k}}$ depending on $d, \bar{k}, \underline{k}$ alone such that for every compact set $S \subset \mathbb{R}^{d}, M>0$ and $\eta \in\left(0, e^{-1} \wedge 4(2 \pi \underline{k})^{-d / 2}\right)$, we have

$$
\begin{equation*}
\log N\left(\eta, \mathbb{F},\|\cdot\|_{\infty, S}\right) \leq C_{d} N\left(a, S^{a}\right)\left(\log \frac{c_{d, \bar{k}, \underline{k}}}{\eta}\right)^{d+1} \tag{4.29}
\end{equation*}
$$

Proof. The idea here is to take $f_{G, \bullet} \in \mathbb{F}$ (induced by some $G \in \mathcal{P}\left(\mathbb{R}^{d}\right)$ ), approximate $G$ by a discrete distribution $H$, and then further approximate that discrete distribution with another discrete distribution over a fixed set of atoms and weights. So let $G \in \mathcal{P}\left(\mathbb{R}^{d}\right)$, and apply the previous Lemma to obtain a discrete distribution $H$ supported on $S^{a}$ with at most $l$ atoms such that

$$
\left\|f_{G, \bullet}-f_{H, \bullet}\right\|_{\infty, S^{a}} \leq\left(1+\frac{1}{\sqrt{2 \pi}}\right)(2 \pi \underline{k})^{-d / 2} e^{-a^{2} /(2 \bar{k})}
$$

Let $\mathcal{C}$ denote a minimal $\zeta$-net of $S^{a}$, and let $H^{\prime}$ approximate each atom of $H$ with its closest element from $\mathcal{C}$. Writing $H=\sum_{j} w_{j} \delta_{a_{j}}$ and $H^{\prime}=\sum w_{j} \delta_{b_{j}}$, we have

$$
\begin{aligned}
\left\|f_{H, \bullet}-f_{H^{\prime}, \bullet}\right\|_{\infty, S^{a}} & =\max _{i \in[n]} \sup _{x \in S^{a}}\left|f_{H, \Sigma_{i}}(x)-f_{H^{\prime}, \Sigma_{i}}(x)\right| \\
& \leq \max _{i \in[n]} \sup _{x \in S^{a}} \sum_{j} w_{j}\left|\varphi_{\Sigma_{i}}\left(x-a_{j}\right)-\varphi_{\Sigma_{i}}\left(x-b_{j}\right)\right| \\
& \leq \zeta \max _{i \in[n]} \sup _{z}\left\|\nabla \varphi_{\Sigma_{i}}(z)\right\|_{2}=\zeta \max _{i \in[n]} \sup _{z} \varphi_{\Sigma_{i}}(z)\left\|\Sigma_{i}^{-1} z\right\|_{2} \\
& \leq \zeta \underline{k}^{-1}(2 \pi \underline{k})^{-d / 2} \max _{i \in[n]} \sup _{t} \exp \left(-t^{2} / 2 \bar{k}\right) t \\
& \leq \zeta \underline{k}^{-1}(2 \pi \underline{k})^{-d / 2}(\bar{k} / e)^{1 / 2}
\end{aligned}
$$

Let $\mathcal{D}$ denote a minimal $\xi$-net of $\Delta^{l-1}$ in the $\ell_{1}$ norm, and approximate the weights $w$ by their closest element $v \in \mathcal{D}$. Writing $H^{\prime \prime}=\sum_{j} v_{j} \delta_{b_{j}}$,

$$
\begin{aligned}
\left\|f_{H^{\prime}, \bullet}-f_{H^{\prime \prime}, \bullet}\right\|_{\infty, S^{a}} & =\max _{i \in[n]} \sup _{x \in S^{a}}\left|f_{H^{\prime}, \Sigma_{i}}(x)-f_{H^{\prime \prime}, \Sigma_{i}}(x)\right| \\
& \leq \max _{i \in[n]} \sup _{x \in S^{a}} \sum_{j}\left|w_{j}-v_{j}\right|\left|\varphi_{\Sigma_{i}}\left(x-b_{j}\right)\right| \leq(2 \pi \underline{k})^{-d / 2} \xi .
\end{aligned}
$$

Applying triangle inequality to the past three displays,

$$
\begin{equation*}
\left\|f_{G, \bullet}-f_{H^{\prime \prime}, \bullet}\right\|_{\infty, S^{a}} \leq(2 \pi \underline{k})^{-d / 2}\left[2 e^{-a^{2} /(2 \bar{k})}+\zeta \underline{k}^{-1}(\bar{k} / e)^{1 / 2}+\xi\right] \tag{4.30}
\end{equation*}
$$

Letting $\xi=(2 \pi \underline{k})^{d / 2} \frac{\eta}{4}, \zeta=\xi \underline{k}(\bar{k} / e)^{-1 / 2}$, and $a=\sqrt{2 \bar{k} \log \xi^{-1}}$ yields $\left\|f_{G, \bullet}-f_{H^{\prime \prime}, \bullet}\right\|_{\infty, S^{a}} \leq \eta$. In order to take $a$ as such we need $\xi<1$, or equivalently $\eta<4(2 \pi \underline{k})^{-d / 2}$.

The number of possible $H^{\prime \prime}$ is

$$
|\mathcal{C}| \cdot|\mathcal{D}|=N\left(\xi, \Delta^{l-1}\right)\binom{N\left(\zeta, S^{a}\right)}{l} \leq\left[\left(1+\frac{2}{\xi}\right) \frac{e N\left(\zeta, S^{a}\right)}{l}\right]^{l}
$$

From the previous Lemma, $l \geq N\left(a, S^{a}\right) / \bar{k}^{d}$, so

$$
\frac{N\left(\zeta, S^{a}\right)}{l} \leq \bar{k}^{d} \frac{N\left(\zeta, S^{a}\right)}{N\left(a, S^{a}\right)} \leq \bar{k}^{d}\left(1+\frac{a}{\zeta}\right)^{d}=\bar{k}^{d}\left(1+\kappa \frac{2}{\sqrt{e} \xi^{3 / 2}}\right)^{d} \leq C_{d} \frac{\bar{k}^{2 d}}{\underline{k}^{d+3 d^{2} / 4}}\left(\frac{1}{\eta}\right)^{3 p / 2}
$$

Thus,

$$
\log N\left(\eta, \mathbb{F},\|\cdot\|_{\infty, S}\right) \leq C_{d} N\left(a, S^{a}\right) \log \left(\frac{1}{\underline{k}^{d / 2}}\left(\frac{\bar{k}^{2 d}}{\underline{k}^{3(d / 2+1)}} \vee 1\right) \frac{e}{\eta}\right)^{d+1}
$$

### 4.6.3 Proof of Theorem 4.8

Throughout the proof we will group sequences of the form $\theta_{1}, \ldots, \theta_{n}$ into $n \times d$ matrices $\theta$, so that, for instance, the regret $\frac{1}{n} \sum_{i=1}^{n} \mathbb{E}\left\|\hat{\theta}_{i}-\hat{\theta}_{i}^{*}\right\|_{2}^{2}$ in the statement of the theorem may be rewritten as the expected squared Frobenius norm $\frac{1}{n} \mathbb{E}\left\|\hat{\theta}-\hat{\theta}^{*}\right\|_{F}^{2}$, where $\|\theta\|_{F}^{2}=\sum_{i=1}^{n} \sum_{j=1}^{d} \theta_{i j}^{2}$. Additionally, we use the same notation introduced at the start of Section 4.6.2.

### 4.6.3.1 Regularizing the Bayes rule

In evaluating $\hat{\theta}$, an apparent difficulty is that the denominator in Tweedie's formula can be arbitrarily small. However, since $\widehat{G}_{n}$ is an approximate NPMLE, we show that the likelihood is lower bounded at each of the observations. In accordance with (4.17), we write that

$$
\begin{equation*}
\frac{1}{n} \sum_{i=1}^{n} \log f_{\widehat{G}_{n}, \Sigma_{i}}\left(X_{i}\right)-\sup _{G \in \mathcal{P}\left(\mathbb{R}^{d}\right)} \frac{1}{n} \sum_{i=1}^{n} \log f_{G, \Sigma_{i}}\left(X_{i}\right) \geq-q, \tag{4.31}
\end{equation*}
$$

for some $q>0$. Following Jiang and Zhang, 2009, Proof of Proposition 2, for a fixed $j \in[n]$ choose $G=n^{-1} \delta_{X_{j}}+\left(1-n^{-1}\right) \widehat{G}_{n}$. Then by the previous display,

$$
\begin{aligned}
\prod_{i=1}^{n} f_{\widehat{G}_{n}, \Sigma_{i}}\left(X_{i}\right) & \geq e^{-n q} \prod_{i=1}^{n} f_{G, \Sigma_{i}}\left(X_{i}\right) \\
& \geq e^{-n q} n^{-1} \varphi_{\Sigma_{i}}(0)\left(1-n^{-1}\right)^{n-1} \prod_{i: i \neq j} f_{\widehat{G}_{n}, \Sigma_{i}}\left(X_{i}\right)
\end{aligned}
$$

Cancelling terms for $i \in[n] \backslash\{j\}$, we conclude

$$
\begin{equation*}
f_{\widehat{G}_{n}, \Sigma_{j}}\left(X_{j}\right) \geq e^{-n q-\log n} \frac{1}{e \sqrt{\left|2 \pi \Sigma_{j}\right|}} \tag{4.32}
\end{equation*}
$$

Given this, it is natural to define the regularized empirical Bayes and oracle Bayes rules

$$
\begin{align*}
& \hat{\theta}_{\rho, i}=X_{i}+\Sigma_{i} \frac{\nabla f_{\widehat{G}_{n}, \Sigma_{i}}\left(X_{i}\right)}{f_{\widehat{G}_{n}, \Sigma_{i}}\left(X_{i}\right) \vee\left(\rho / \sqrt{\left|\Sigma_{i}\right|}\right)}  \tag{4.33}\\
& \hat{\theta}_{\rho, i}^{*}=X_{i}+\Sigma_{i} \frac{\nabla f_{G^{*}, \Sigma_{i}}\left(X_{i}\right)}{f_{G^{*}, \Sigma_{i}}\left(X_{i}\right) \vee\left(\rho / \sqrt{\left|\Sigma_{i}\right|}\right)} . \tag{4.34}
\end{align*}
$$

By the lower bound (4.32) we know that $\hat{\theta}_{\rho}=\hat{\theta}$ when $\rho \leq \rho_{0}:=e^{-n q-\log n} \frac{1}{e(2 \pi)^{d / 2}}$. In particular,

$$
\begin{equation*}
\left\|\hat{\theta}-\hat{\theta}^{*}\right\|_{F}=\left\|\hat{\theta}_{\rho}-\hat{\theta}^{*}\right\|_{F} \leq\left\|\hat{\theta}_{\rho}-\hat{\theta}_{\rho}^{*}\right\|_{F}+\left\|\hat{\theta}_{\rho}^{*}-\hat{\theta}^{*}\right\|_{F} \tag{4.35}
\end{equation*}
$$

The first term $\left\|\hat{\theta}_{\rho}-\hat{\theta}_{\rho}^{*}\right\|_{F}$ represents the regret between regularized rules, which prevents the denominator in Tweedie's formula from blowing up. The second term represents the cost of introducing a small amount of regularization in the oracle Bayes rule.

### 4.6.3.2 Regularization error of oracle Bayes

Let us first consider the second term $\left\|\hat{\theta}_{\rho}^{*}-\hat{\theta}^{*}\right\|_{F}$ on the RHS of the bound (4.35). Fixing $i \in[n]$, let $G_{i}^{*}$ denote the distribution of $\xi_{i}=\Sigma_{i}^{-1 / 2} \theta_{i}^{*}$ where $\theta_{i}^{*} \sim G^{*}$. Then we may write $X_{i}=\Sigma_{i}^{1 / 2} X_{i}$ where $X_{i} \sim f_{G_{i}^{*}, I_{d}}$. Note how the scale change affects the terms in Tweedie's formula:

$$
\begin{align*}
f_{G^{*}, \Sigma_{i}}\left(X_{i}\right) & =\mathbb{E}_{\vartheta_{i} \sim G^{*}}\left[\frac{1}{\sqrt{\left|2 \pi \Sigma_{i}\right|}} \exp \left(-\frac{1}{2}\left(X_{i}-\vartheta_{i}\right)^{\prime} \Sigma_{i}^{-1}\left(X_{i}-\vartheta_{i}\right)\right)\right] \\
& =\frac{1}{\sqrt{\left|\Sigma_{i}\right|}} \mathbb{E}_{\xi_{i} \sim G_{i}^{*}}\left[\varphi_{I_{d}}\left(\Sigma_{i}^{-1 / 2} X_{i}-\xi_{i}\right)\right]=\frac{1}{\sqrt{\left|\Sigma_{i}\right|}} f_{G_{i}^{*}, I_{d}}\left(X_{i}\right) \\
\nabla f_{G^{*}, \Sigma_{i}}\left(X_{i}\right) & =\mathbb{E}_{\vartheta_{i} \sim G^{*}}\left[\Sigma_{i}^{-1}\left(\vartheta_{i}-X_{i}\right) \frac{1}{\sqrt{\left|2 \pi \Sigma_{i}\right|}} \exp \left(-\frac{1}{2}\left(X_{i}-\vartheta_{i}\right)^{\prime} \Sigma_{i}^{-1}\left(X_{i}-\vartheta_{i}\right)\right)\right]  \tag{4.36}\\
& =\frac{1}{\sqrt{\left|\Sigma_{i}\right|}} \Sigma_{i}^{-1 / 2} \nabla f_{G_{i}^{*}, I_{d}}\left(X_{i}\right)
\end{align*}
$$

In particular, Tweedie's formula, even in its regularized form, is scale equivariant:

$$
\begin{equation*}
\hat{\theta}_{\rho, i}^{*}=\Sigma_{i}^{1 / 2}\left(X_{i}+\frac{\nabla f_{G_{i}^{*}, I_{d}}\left(X_{i}\right)}{f_{G_{i}^{*}, I_{d}}\left(X_{i}\right) \vee \rho}\right) \tag{4.37}
\end{equation*}
$$

In this form, Saha and Guntuboyina (2020a, Lemma 4.3) directly applies. Specifically, defining

$$
\Delta(G, \rho):=\int\left(1-\frac{f_{G, I_{d}}}{f_{G, I_{d}} \vee \rho}\right)^{2} \frac{\left\|\nabla f_{G, I_{d}}\right\|_{2}^{2}}{f_{G, I_{d}}}
$$

for any $\rho \leq \rho_{0}$ and for all compact sets $S_{1}, \ldots, S_{n} \subset \mathbb{R}^{d}$,

$$
\begin{align*}
\mathbb{E}\left\|\hat{\theta}_{\rho}^{*}-\hat{\theta}^{*}\right\|_{F}^{2} & =\sum_{i=1}^{n} \mathbb{E}\left\|\hat{\theta}_{\rho, i}^{*}-\hat{\theta}_{i}^{*}\right\|_{2}^{2} \leq \bar{k} \sum_{i=1}^{n} \Delta\left(G_{i}^{*}, \rho\right)  \tag{4.38}\\
& \leq \bar{k} \sum_{i=1}^{n}\left\{C_{d} N\left(\frac{4}{L(\rho)}, S_{i}\right) L^{d}(\rho) \rho+d G_{i}^{*}\left(S_{i}^{\mathrm{c}}\right)\right\},
\end{align*}
$$

where $L(\rho):=\sqrt{-\log \left((2 \pi)^{d} \rho^{2}\right)}$ and $N$ denotes the usual covering number in the Euclidean norm. Choosing $\rho=(2 \pi)^{-d / 2} / n$ and $S_{i}=\Sigma_{i}^{-1 / 2} S^{M}$,

$$
\mathbb{E}\left\|\hat{\theta}_{\rho}^{*}-\hat{\theta}^{*}\right\|_{F}^{2} \leq \bar{k} n\left\{C_{d} N\left(\frac{4}{\sqrt{\log n}}, \Sigma_{i}^{-1 / 2} S^{M}\right) \frac{(\log n)^{d / 2}}{n}+d G^{*}\left(\left(S^{M}\right)^{\mathrm{c}}\right)\right\}
$$

Let $x_{1}, \ldots, x_{m}$ denote a $t$-net of $S^{M}$. Let $y \in S_{i}$ and $x=\Sigma_{i}^{1 / 2} y$. There is some $j$ s.t. $\left\|x_{j}-x\right\|_{2} \leq t$. Let $y_{j}=\Sigma_{i}^{-1 / 2} x_{j}$. Then

$$
t^{2} \geq\left(x_{j}-x\right)^{\prime}\left(x_{j}-x\right)=\left(y_{j}-y\right)^{\prime} \Sigma_{i}\left(y_{j}-y\right) \geq \underline{k}\left\|y_{j}-y\right\|_{2}^{2}
$$

so $y_{1}, \ldots, y_{m}$ is a $t / \underline{k}^{1 / 2}$-net of $S_{i}$. This shows $N\left(t / \underline{k}^{1 / 2}, S_{i}\right) \leq N\left(t, S^{M}\right)$. By Saha and Guntuboyina (2020b, Lemma F.6) and Markov's inequality,

$$
\begin{equation*}
\mathbb{E}\left\|\hat{\theta}_{\rho}^{*}-\hat{\theta}^{*}\right\|_{F}^{2} \leq C_{d} \bar{k} n\left\{\underline{k}^{-d / 2} \operatorname{Vol}\left(S^{1}\right) M^{d} \frac{(\log n)^{d}}{n}+\inf _{q \geq(d+1) / 2 \log n}\left(\frac{2 \mu_{q}}{M}\right)^{q}\right\} \tag{4.39}
\end{equation*}
$$

### 4.6.3.3 Regret of regularized rules

Now we consider the first term $\left\|\hat{\theta}_{\rho}-\hat{\theta}_{\rho}^{*}\right\|_{F}$ on the RHS of the bound (4.35). First, we will introduce some additional notation. For $\delta>0$ let $A_{\delta}=\left\{\bar{h}^{2}\left(f_{\widehat{G}_{n}, \bullet}, f_{G^{*}, \bullet}\right) \leq \delta\right\}$. Given a compact set $S \subset \mathbb{R}^{d}$, define another metric

$$
m^{S}\left(G, G^{\prime}\right):=\max _{i \in[n]} \sup _{x: \mathrm{D}_{S}(x) \leq M}\left\|\frac{\Sigma_{i} \nabla f_{G, \Sigma_{i}}(x)}{f_{G, \Sigma_{i}}(x) \vee\left(\rho / \sqrt{\left|\Sigma_{i}\right|}\right)}-\frac{\Sigma_{i} \nabla f_{G^{\prime}, \Sigma_{i}}(x)}{f_{G^{\prime}, \Sigma_{i}}(x) \vee\left(\rho / \sqrt{\left|\Sigma_{i}\right|}\right)}\right\|_{2}
$$

Let $G^{(1)}, \ldots, G^{(N)}$ denote a minimal $\eta^{*}$-covering of $\left\{G: \bar{h}^{2}\left(f_{G, \bullet}, f_{G^{*}, \bullet}\right) \leq \delta\right\}$ in the metric $m^{S}$. For $j \in[N]$ similarly define an $n \times d$ matrix $\hat{\theta}_{\rho}^{(j)}$ where the $i^{\text {th }}$ row is given by $X_{i}+$ $\Sigma_{i} \frac{\nabla f_{G}(j), \Sigma_{i}\left(X_{i}\right)}{f_{G^{(j)}, \Sigma_{i}}\left(X_{i}\right) \vee\left(\rho / \sqrt{\left|\Sigma_{i}\right|}\right)}$. We bound the regret as $\left\|\hat{\theta}_{\rho}-\hat{\theta}_{\rho}^{*}\right\|_{F} \leq \sum_{t=1}^{4} \zeta_{t}$, where

$$
\begin{align*}
& \zeta_{1}:=\left\|\hat{\theta}_{\rho}-\hat{\theta}_{\rho}^{*}\right\|_{F} 1_{A_{\delta}^{c}} \\
& \zeta_{2}:=\left(\left\|\hat{\theta}_{\rho}-\hat{\theta}_{\rho}^{*}\right\|_{F}-\max _{j \in[N]}\left\|\hat{\theta}_{\rho}^{(j)}-\hat{\theta}_{\rho}^{*}\right\|_{F}\right)_{+} 1_{A_{\delta}} \\
& \zeta_{3}:=\max _{j \in[N]}\left(\left\|\hat{\theta}_{\rho}^{(j)}-\hat{\theta}_{\rho}^{*}\right\|_{F}-\mathbb{E}\left\|\hat{\theta}_{\rho}^{(j)}-\hat{\theta}_{\rho}^{*}\right\|_{F}\right)_{+}  \tag{4.40}\\
& \zeta_{4}:=\max _{j \in[N]} \mathbb{E}\left\|\hat{\theta}_{\rho}^{(j)}-\hat{\theta}_{\rho}^{*}\right\|_{F}
\end{align*}
$$

We will control the second moment of each $\zeta_{t}$. Here's our rough overview. $\zeta_{1}$ uses Theorem 4.6 to show the NPMLE places small probability on $A_{\delta}^{\mathrm{c}} ; \zeta_{2}$ uses the fact that (on $A_{\delta}$ ) the cover $\left\{G^{(j)}\right\}$ must have some element that is close to the NPMLE in $m^{S} ; \zeta_{3}$ follows from Gaussian concentration of measure; and $\zeta_{4}$ bounds each expectation individually and uses closeness in Hellinger.

## Bounding $\mathbb{E} \zeta_{1}^{2}$

By the scaled Tweedie's formula (4.37)

$$
\left\|\hat{\theta}_{\rho, i}-\hat{\theta}_{\rho, i}^{*}\right\|_{2}^{2} \leq \bar{k}\left\|\frac{\nabla f_{\widehat{G}_{i}, I_{d}}\left(X_{i}\right)}{f_{\widehat{G}_{i}, I_{d}}\left(X_{i}\right) \vee \rho}-\frac{\nabla f_{G_{i}^{*}, I_{d}}\left(X_{i}\right)}{f_{G_{i}^{*}, I_{d}}\left(X_{i}\right) \vee \rho}\right\|_{2}^{2}
$$

Saha and Guntuboyina (2020b, Lemma F.1) provides

$$
\begin{equation*}
\mathbb{E} \zeta_{1}^{2} \leq 4 \bar{k} n \log \left(\frac{(2 \pi)^{d}}{\rho^{2}}\right) \mathbb{P}\left(A_{\delta}^{\mathrm{c}}\right) \tag{4.41}
\end{equation*}
$$

By Theorem 4.6, there is a constant $C_{d, k, k}>0$ such that $\delta=C_{d, k, \bar{k}} \varepsilon_{n}^{2}\left(M, S, G^{*}\right)$ satisfies $\mathbb{P}\left(A_{\delta}^{\mathrm{c}}\right) \leq 2 / n$. Hence

$$
\begin{equation*}
\mathbb{E} \zeta_{1}^{2} \leq 48 \bar{k} \log (n) \tag{4.42}
\end{equation*}
$$

## Bounding $\mathbb{E} \zeta_{2}^{2}$

Observe

$$
\begin{aligned}
\zeta_{2}^{2} & \leq 1_{A_{\delta}} \min _{j \in[N]}\left\|\hat{\theta}_{\rho}-\hat{\theta}_{\rho}^{(j)}\right\|_{F}^{2} \\
& =1_{A_{\delta}} \min _{j \in[N]} \sum_{i=1}^{n}\left\|\frac{\Sigma_{i} \nabla f_{\widehat{G}_{n}, \Sigma_{i}}\left(X_{i}\right)}{f_{\widehat{G}_{n}, \Sigma_{i}}\left(X_{i}\right) \vee\left(\rho / \sqrt{\left|\Sigma_{i}\right|}\right)}-\frac{\Sigma_{i} \nabla f_{G^{(j)}, \Sigma_{i}}\left(X_{i}\right)}{f_{G^{(j)}, \Sigma_{i}}\left(X_{i}\right) \vee\left(\rho / \sqrt{\left|\Sigma_{i}\right|}\right)}\right\|_{2}^{2}
\end{aligned}
$$

On $A_{\delta}$, we may take $j$ such that $m^{S}\left(\widehat{G}_{n}, G^{(j)}\right) \leq \eta^{*}$. For each $i$, consider two cases, where $X_{i} \in S^{M}$ and where $X_{i} \notin S^{M}$. When $X_{i} \in S^{M}$ bound the above $\|\cdot\|_{2}$ by the supremum over all $x \in S^{M}$. When $X_{i} \notin S^{M}$ bound the regularized rules as before. This yields

$$
\begin{equation*}
\zeta_{2}^{2} \leq 1_{A_{\delta}}\left(\#\left\{i: X_{i} \in S^{M}\right\}\left(\eta^{*}\right)^{2}+\#\left\{i: X_{i} \notin S^{M}\right\} 4 \bar{k} \log \left(\frac{(2 \pi)^{d}}{\rho^{2}}\right)\right) \tag{4.43}
\end{equation*}
$$

so in particular

$$
\begin{equation*}
\mathbb{E} \zeta_{2}^{2} \leq n\left(\eta^{*}\right)^{2}+4 \bar{k} \log \left(\frac{(2 \pi)^{d}}{\rho^{2}}\right) \sum_{i=1}^{n} \mathbb{P}\left(\mathfrak{d}_{S}\left(X_{i}\right) \geq M\right) \tag{4.44}
\end{equation*}
$$

To bound the probabilities on the RHS, write $X_{i}=\theta_{i}+\Sigma^{1 / 2} Z_{i}$. By Lemma 4.13, taking $\lambda \downarrow 0$,

$$
\begin{equation*}
\mathbb{E} \zeta_{2}^{2} / n \leq\left(\eta^{*}\right)^{2}+4 \bar{k} \log \left(\frac{(2 \pi)^{d}}{\rho^{2}}\right)\left(C_{d} \frac{M^{d-2}}{n} \bar{k}^{1-d / 2}+\inf _{q \geq(d+1) / 2 \log n}\left(\frac{2 \mu_{q}}{M}\right)^{q}\right) \tag{4.45}
\end{equation*}
$$

## Bounding $\mathbb{E} \zeta_{3}^{2}$

Fix $j \in[N]$. Let $Z_{i} \stackrel{\text { iid }}{\sim} \mathcal{N}\left(0, I_{d}\right)$ and $\xi_{i}^{*} \sim G_{i}^{*}$ and $\xi_{i}^{(j)} \sim G_{i}^{(j)}$, where $G_{i}^{(j)}$ denotes the scale change of $G^{(j)}$ by $\Sigma_{i}^{-1 / 2}$. In accordance with (4.37), we write

$$
\begin{aligned}
\| \hat{\theta}_{\rho}^{(j)} & -\hat{\theta}_{\rho}^{*} \|_{F} \\
& =\left(\sum_{i=1}^{n}\left\|\Sigma_{i}^{1 / 2}\left(\xi_{i}^{(j)}+Z_{i}+\frac{\nabla f_{G_{i}^{(j)}, I_{d}}\left(\xi_{i}^{(j)}+Z_{i}\right)}{f_{G_{i}^{(j)}, I_{d}}\left(\xi_{i}^{(j)}+Z_{i}\right) \vee \rho}\right)-\Sigma_{i}^{1 / 2}\left(\xi_{i}^{*}+Z_{i}+\frac{\nabla f_{G_{i}^{*}, I_{d}}\left(\xi_{i}^{*}+Z_{i}\right)}{f_{G_{i}^{*}, I_{d}}\left(\xi_{i}^{*}+Z_{i}\right) \vee \rho}\right)\right\|_{2}^{2}\right)^{1 / 2}
\end{aligned}
$$

Call the RHS above $F(Z)$. Then

$$
\begin{aligned}
\left|F(Z)-F\left(Z^{\prime}\right)\right| & =\left|\left\|\hat{\theta}_{\rho}^{(j)}(Z)-\hat{\theta}_{\rho}^{*}(Z)\right\|_{F}-\left\|\hat{\theta}_{\rho}^{(j)}\left(Z^{\prime}\right)-\hat{\theta}_{\rho}^{*}\left(Z^{\prime}\right)\right\|_{F}\right| \\
& \leq\left\|\hat{\theta}_{\rho}^{(j)}(Z)-\hat{\theta}_{\rho}^{(j)}\left(Z^{\prime}\right)\right\|_{F}+\left\|\hat{\theta}_{\rho}^{*}(Z)-\hat{\theta}_{\rho}^{*}\left(Z^{\prime}\right)\right\|_{F}
\end{aligned}
$$

Focusing on the second term on the RHS,

$$
\begin{aligned}
\| \hat{\theta}_{\rho}^{*}(Z) & -\hat{\theta}_{\rho}^{*}\left(Z^{\prime}\right) \|_{F} \\
& \leq \bar{k}^{1 / 2} \sqrt{\sum_{i=1}^{n}\left\|\left(\xi_{i}^{*}+Z_{i}+\frac{\nabla f_{G_{i}^{*}, I_{d}}\left(\xi_{i}^{*}+Z_{i}\right)}{f_{G_{i}^{*}, I_{d}}\left(\xi_{i}^{*}+Z_{i}\right) \vee \rho}\right)-\left(\xi_{i}^{*}+Z_{i}^{\prime}+\frac{\nabla f_{G_{i}^{*}, I_{d}}\left(\xi_{i}^{*}+Z_{i}^{\prime}\right)}{f_{G_{i}^{*}, I_{d}}\left(\xi_{i}^{*}+Z_{i}^{\prime}\right) \vee \rho}\right)\right\|_{2}^{2}} \\
& \leq \bar{k}^{1 / 2} \sqrt{\sum_{i=1}^{n}\left\|\left(Z_{i}+\frac{\nabla f_{G_{i}^{*}-\xi_{i}^{*}, I_{d}}\left(Z_{i}\right)}{f_{G_{i}^{*}-\xi_{i}^{*}, I_{d}}\left(Z_{i}\right) \vee \rho}\right)-\left(Z_{i}^{\prime}+\frac{\nabla f_{G_{i}^{*}-\xi_{i}^{*}, I_{d}}\left(Z_{i}^{\prime}\right)}{f_{G_{i}^{*}-\xi_{i}^{*}, I_{d}}\left(Z_{i}^{\prime}\right) \vee \rho}\right)\right\|_{2}^{2}}
\end{aligned}
$$

Saha and Guntuboyina (2020b, Proof of Lemma F.3) then gives

$$
\left\|\hat{\theta}_{\rho}^{*}(Z)-\hat{\theta}_{\rho}^{*}\left(Z^{\prime}\right)\right\|_{F} \leq \bar{k}^{1 / 2} L^{2}(\rho)\left\|Z-Z^{\prime}\right\|_{F}
$$

where as before $L(\rho):=\sqrt{-\log \left((2 \pi)^{d} \rho^{2}\right)}$. The same argument applies to $\left\|\hat{\theta}_{\rho}^{(j)}(Z)-\hat{\theta}_{\rho}^{(j)}\left(Z^{\prime}\right)\right\|_{F}$. Hence $F$ is $2 \bar{k} L^{2}(\rho)$-Lipschitz. By concentration of Lipschitz functions of Gaussians and a union bound

$$
\mathbb{P}\left(\zeta_{3}^{2} \geq x\right) \leq N \exp \left(-\frac{x^{2}}{8 \bar{k} L^{4}(\rho)}\right)
$$

Integrating the tail gives

$$
\begin{equation*}
\mathbb{E} \zeta_{3}^{2} \leq 8 \bar{k} L^{4}(\rho) \log (e N) \tag{4.46}
\end{equation*}
$$

## Bounding $\mathbb{E} \zeta_{4}^{2}$

Again by the scaled Tweedie's formula (4.37),

$$
\mathbb{E}\left\|\hat{\theta}_{\rho}^{(j)}-\hat{\theta}_{\rho}^{*}\right\|_{F} \leq \sqrt{\mathbb{E}\left\|\hat{\theta}_{\rho}^{(j)}-\hat{\theta}_{\rho}^{*}\right\|_{F}^{2}} \leq \sqrt{\bar{k} \sum_{i=1}^{n} \mathbb{E}_{X_{i} \sim f_{G_{i}^{*}, I_{d}}}\left\|\frac{\nabla f_{G_{i}^{(j)}, I_{d}}\left(X_{i}\right)}{f_{G_{i}^{(j)}, I_{d}}\left(X_{i}\right) \vee \rho}-\frac{\nabla f_{G_{i}^{*}, I_{d}}\left(X_{i}\right)}{f_{G_{i}^{*}, I_{d}}\left(X_{i}\right) \vee \rho}\right\|_{2}^{2}}
$$

Saha and Guntuboyina (2020a, Lemma E.1) bounds the above expectation, yielding

$$
\begin{equation*}
\left(\mathbb{E}\left\|\hat{\theta}_{\rho}^{(j)}-\hat{\theta}_{\rho}^{*}\right\|_{F}\right)^{2} \leq C_{d} \bar{k} \sum_{i=1}^{n} \max \left\{\left(\frac{L^{2}(\rho)}{2}\right)^{3},\left|\log h\left(f_{G_{i}^{*}, I_{d}}, f_{G_{i}^{(j)}, I_{d}}\right)\right|\right\} h^{2}\left(f_{G_{i}^{*}, I_{d}}, f_{G_{i}^{(j)}, I_{d}}\right) . \tag{4.47}
\end{equation*}
$$

By a change of variables,

$$
h^{2}\left(f_{G_{i}^{*}, I_{d}}, f_{G_{i}^{(j)}, I_{d}}\right)=h^{2}\left(f_{G^{*}, \Sigma_{i}}, f_{G^{(j)}, \Sigma_{i}}\right) .
$$

Using the shorthand $h_{i}^{2}=h^{2}\left(f_{G^{*}, \Sigma_{i}}, f_{G^{(j)}, \Sigma_{i}}\right)$ and using $\rho=(2 \pi)^{-d / 2} / n$,

$$
\begin{align*}
\left(\mathbb{E}\left\|\hat{\theta}_{\rho}^{(j)}-\hat{\theta}_{\rho}^{*}\right\|_{F}\right)^{2} & \leq C_{d} \bar{k} \sum_{i=1}^{n} \max \left\{(\log n)^{3},-\log h_{i}\right\} h_{i}^{2} \\
& =C_{d} \bar{k}\left(\sum_{i:(\log n)^{3} \geq-\log h_{i}}(\log n)^{3} h_{i}^{2}+\sum_{i:(\log n)^{3}<-\log h_{i}}-\left(\log h_{i}\right) h_{i}^{2}\right)  \tag{4.48}\\
& \leq C_{d} \bar{k}\left(n(\log n)^{3} \delta+\sum_{i:(\log n)^{3}<\log h_{i}^{-1}}\left(\log h_{i}^{-1}\right) h_{i}^{2}\right),
\end{align*}
$$

where in the last step we used $\frac{1}{n} \sum_{i=1}^{n} h_{i}^{2}=\bar{h}^{2}\left(f_{G^{*}, \bullet}, f_{G^{(j)}, \bullet}\right) \leq \delta$. To bound the second term, note for $n \geq 6,(\log n)^{3} \geq 3 \log n$, implying $h_{i} \leq n^{-3}$ for all $i$ such that $(\log n)^{3}<\log h_{i}^{-1}$. Since $h_{i} \log h_{i}^{-1} \leq e^{-1}$ for all $h_{i} \in[0,1]$,

$$
\sum_{i:(\log n)^{3}<\log h_{i}^{-1}}\left(\log h_{i}^{-1}\right) h_{i}^{2} \leq \sum_{i:(\log n)^{3}<\log h_{i}^{-1}} \frac{1}{e n^{3}} \leq \frac{1}{e n^{2}}
$$

The first term dominates, so

$$
\begin{equation*}
\mathbb{E} \zeta_{4}^{2} \leq C_{d} \bar{k} n(\log n)^{3} \delta \tag{4.49}
\end{equation*}
$$

## Bounding the metric entropy $\log N$

We will actually bound the larger covering number $\log N\left(\eta^{*}, \mathcal{P}\left(\mathbb{R}^{d}\right), m^{S}\right)$ of the space of all probability measures $\mathcal{P}\left(\mathbb{R}^{d}\right)$ in the metric $m^{S}$. For any measure $G$ we let $G_{i}$ denote the measure scaled by $\Sigma_{i}^{-1 / 2}$ as in the scaled Tweedie formula. For $G, H \in \mathcal{P}\left(\mathbb{R}^{d}\right)$,

$$
\begin{aligned}
m^{S}(G, H):= & \max _{i \in[n]} \sup _{x: D_{S}(x) \leq M}\left\|\frac{\Sigma_{i} \nabla f_{G, \Sigma_{i}}(x)}{f_{G, \Sigma_{i}}(x) \vee\left(\rho / \sqrt{\left|\Sigma_{i}\right|}\right)}-\frac{\Sigma_{i} \nabla f_{H, \Sigma_{i}}(x)}{f_{H, \Sigma_{i}}(x) \vee\left(\rho / \sqrt{\left|\Sigma_{i}\right|}\right)}\right\|_{2} \\
\leq & \max _{i \in[n]} \sup _{x: D_{S}(x) \leq M}\left\|\frac{\Sigma_{i} \nabla f_{G, \Sigma_{i}}(x)}{f_{G, \Sigma_{i}}(x) \vee\left(\rho / \sqrt{\left|\Sigma_{i}\right|}\right)}-\frac{\Sigma_{i} \nabla f_{G, \Sigma_{i}}(x)}{f_{H, \Sigma_{i}}(x) \vee\left(\rho / \sqrt{\left|\Sigma_{i}\right|}\right)}\right\|_{2} \\
& +\max _{i \in[n]} \sup _{x: D_{S}(x) \leq M}\left\|\frac{\Sigma_{i} \nabla f_{G, \Sigma_{i}}(x)}{f_{H, \Sigma_{i}}(x) \vee\left(\rho / \sqrt{\left|\Sigma_{i}\right|}\right)}-\frac{\Sigma_{i} \nabla f_{H, \Sigma_{i}}(x)}{f_{H, \Sigma_{i}}(x) \vee\left(\rho / \sqrt{\left|\Sigma_{i}\right|}\right)}\right\|_{2} \\
\leq & \max _{i \in[n]} \sup _{x: D_{S}(x) \leq M}\left\|\frac{\Sigma_{i} \nabla f_{G, \Sigma_{i}}(x)}{f_{G, \Sigma_{i}}(x) \vee\left(\rho / \sqrt{\left|\Sigma_{i}\right|}\right)}\right\|_{2} \\
& \times \frac{\left|f_{G, \Sigma_{i}}(x) \vee\left(\rho / \sqrt{\left|\Sigma_{i}\right|}\right)-f_{H, \Sigma_{i}}(x) \vee\left(\rho / \sqrt{\left|\Sigma_{i}\right|}\right)\right|}{f_{H, \Sigma_{i}}(x) \vee\left(\rho / \sqrt{\left|\Sigma_{i}\right|}\right)} \\
& +\max _{i \in[n]} \sup _{x: D_{S}(x) \leq M}\left\|\frac{\Sigma_{i}\left(\nabla f_{G, \Sigma_{i}}(x)-\nabla f_{H, \Sigma_{i}}(x)\right)}{f_{H, \Sigma_{i}}(x) \vee\left(\rho / \sqrt{\left|\Sigma_{i}\right|}\right)}\right\|_{2}
\end{aligned}
$$

For the first term, by Saha and Guntuboyina (2020b, Lemma F.1),

$$
\left\|\frac{\Sigma_{i} \nabla f_{G, \Sigma_{i}}(x)}{f_{G, \Sigma_{i}}(x) \vee\left(\rho / \sqrt{\left|\Sigma_{i}\right|}\right)}\right\|_{2}=\left\|\Sigma_{i}^{1 / 2} \frac{\nabla f_{G_{i}^{*}, I_{d}}\left(\Sigma_{i}^{-1 / 2} x\right)}{f_{G_{i}^{*}, I_{d}}\left(\Sigma_{i}^{-1 / 2} y\right) \vee \rho}\right\|_{2} \leq \bar{k}^{1 / 2} L(\rho) .
$$

Replacing $f \vee\left(\rho / \sqrt{\left|\Sigma_{i}\right|}\right)$ with $\rho / \sqrt{\left|\Sigma_{i}\right|}$ in the denominator can only make the denominator smaller, so

$$
\begin{align*}
m^{S}(G, H) \leq & \bar{k}^{1 / 2} \rho^{-1} L(\rho) \max _{i \in[n]} \sup _{x: \mathbb{O}_{S}(x) \leq M} \sqrt{\left|\Sigma_{i}\right|}\left|f_{G, \Sigma_{i}}(x)-f_{H, \Sigma_{i}}(x)\right|  \tag{4.50}\\
& +\rho^{-1} \max _{i \in[n]} \sup _{x: \mathcal{D}_{S}(x) \leq M} \sqrt{\left|\Sigma_{i}\right|}\left\|\Sigma_{i}\left(\nabla f_{G, \Sigma_{i}}(x)-\nabla f_{H, \Sigma_{i}}(x)\right)\right\|_{2}  \tag{4.51}\\
\leq & \bar{k}^{d / 2+1 / 2} \rho^{-1} L(\rho)\left\|f_{G, \bullet}-f_{H, \bullet}\right\|_{\infty, S^{M}}+\bar{k}^{d / 2+1} \rho^{-1}\left\|f_{G, \bullet}-f_{H, \bullet}\right\|_{\nabla, S^{M}} \tag{4.52}
\end{align*}
$$

In particular, letting

$$
\eta^{*}=\left(\bar{k}^{d / 2+1 / 2} L(\rho)+\bar{k}^{d / 2+1}\right) \frac{\eta}{\rho}
$$

we have

$$
\log N\left(\eta^{*}, \mathcal{P}\left(\mathbb{R}^{d}\right), m^{S}\right) \leq \log N\left(\eta / 2, \mathbb{F},\|\cdot\|_{\infty, S^{M}}\right)+\log N\left(\eta / 2, \mathbb{F},\|\cdot\|_{\nabla, S^{M}}\right)
$$

We already have a bound on $\log N\left(\eta / 2, \mathbb{F},\|\cdot\|_{\infty, S^{M}}\right)$ in Lemma 4.16, and we bound the other term similarly in Lemma 4.17 below. Combining these bounds,

$$
\log N\left(\eta^{*}, \mathcal{P}\left(\mathbb{R}^{d}\right), m^{S}\right) \leq C_{d} N\left(a, S^{M+a}\right)\left(\log \frac{c_{d, \underline{k}, \bar{k}}}{\eta}\right)^{d+1}
$$

where $a=\sqrt{-2 \bar{k} \log \left(\sqrt{\underline{k} \wedge 1} \frac{(2 \pi k)^{d / 2}}{5} \eta\right)}$. Take $\eta=\rho / n=(2 \pi)^{-d / 2} / n^{2}$.

$$
a=\sqrt{4 \bar{k} \log n+2 \bar{k} \log \left(\frac{5}{\sqrt{\underline{k} \wedge 1} \underline{k}^{d / 2}}\right)} \in[\sqrt{2 \bar{k} \log n}, \sqrt{6 \bar{k} \log n}]
$$

provided $1 / n \leq \frac{5}{\sqrt{\underline{k} \wedge 1} \underline{k}^{d / 2}} \leq n$. Hence by Saha and Guntuboyina (2020b, Lemma F.6) (see the argument on page 6) gives

$$
\log N \leq c_{d, k, \bar{k}}(\log n)^{1+d / 2} \operatorname{Vol}\left(S^{2 M}\right)
$$

Lemma 4.17. For all compact $S \subset \mathbb{R}^{d}, M>0$ and $\eta>0$ sufficiently small,

$$
\log N\left(\eta, \mathbb{F},\|\cdot\|_{\nabla, S^{M}}\right) \leq C_{d} N\left(a, S^{a}\right)\left(\log \frac{c_{d, k, \bar{k}}}{\eta}\right)^{d+1}
$$

where $a=\sqrt{2 \bar{k} \log \frac{c_{d, \bar{k}, \underline{b}}^{\prime}}{\eta}}$.
Proof. Fix $G \in \mathcal{P}\left(\mathbb{R}^{d}\right)$. By Lemma 4.15, there is a discrete measure $H$ supported on $S^{a}$ with at most

$$
l:=\left(2\left\lfloor 13.5 a^{2} / \bar{k}\right\rfloor+2\right)^{d} N\left(a, S^{a}\right)+1
$$

atoms such that

$$
\left\|f_{G, \bullet}-f_{H, \bullet}\right\|_{\nabla, S^{a}} \leq a\left(1+\frac{3}{\sqrt{2 \pi}}\right)(2 \pi \underline{k})^{-d / 2} e^{-a^{2} /(2 \bar{k})}
$$

Now let $\mathcal{C}$ denote a minimal $\alpha$-net of $S^{a}$. Write $H=\sum_{j} w_{j} \delta_{a_{j}}$, and define $H^{\prime}=\sum_{j} w_{j} \delta_{b_{j}}$ where $b_{j} \in \mathcal{C}$ is the closest element to $a_{j}$. Then

$$
\begin{aligned}
\left\|\nabla f_{H, \Sigma_{i}}(x)-\nabla f_{H^{\prime}, \Sigma_{i}}(x)\right\|_{2} & \leq \sum_{j} w_{j}\left\|\nabla \varphi_{\Sigma_{i}}\left(x-a_{j}\right)-\nabla \varphi_{\Sigma_{i}}\left(x-b_{j}\right)\right\|_{2} \\
& \leq \underline{k}^{-1 / 2}\left|\Sigma_{i}\right|^{-1 / 2} \sum_{j} w_{j}\left\|\nabla \varphi\left(\Sigma_{i}^{-1 / 2}\left(x-a_{j}\right)\right)-\nabla \varphi\left(\Sigma_{i}^{-1 / 2}\left(x-b_{j}\right)\right)\right\|_{2} \\
& \leq \frac{\underline{k}^{-d / 2-1 / 2} \alpha}{(2 \pi)^{d / 2}}\left[1+\frac{2}{e}+\frac{\alpha}{\sqrt{\underline{k}}}\right]
\end{aligned}
$$

Now let $\mathcal{D}_{+}^{p \times p}$ denote a minimal $\beta$-net of $\Delta_{l-1}$ under $\|\cdot\|_{1}$. Let $H^{\prime \prime}=\sum_{j} w_{j}^{\prime} \delta_{b_{j}}$ where $\left\|w^{\prime}-w\right\|_{1} \leq \beta$. Then

$$
\left\|\nabla f_{H^{\prime}, \Sigma_{i}}(x)-\nabla f_{H^{\prime \prime}, \Sigma_{i}}(x)\right\|_{2} \leq \beta \sup _{u}\left\|\nabla \varphi_{\Sigma_{i}}(u)\right\|_{2} \leq \frac{k^{-d / 2-1 / 2} \beta}{(2 \pi)^{d / 2} \sqrt{e}}
$$

By triangle inequality,

$$
\left\|f_{G, \bullet}-f_{H^{\prime \prime}, \bullet}\right\|_{\nabla, S^{M}} \leq(2 \pi \underline{k})^{-d / 2}\left[\left(1+\frac{3}{\sqrt{2 \pi}}\right) a e^{-a^{2} /(2 \bar{k})}+\frac{\alpha}{\sqrt{\underline{k}}}\left[1+\frac{2}{e}+\frac{\alpha}{\sqrt{\underline{k} e}}\right]+\frac{\beta}{\sqrt{\underline{k}} e}\right]
$$

Taking $a=\sqrt{2 \bar{k} \log \alpha^{-1}} \geq 1$ and $\alpha=\beta=\sqrt{k \wedge 1} \frac{(2 \pi k)^{d / 2}}{5} \eta$,

$$
\left\|f_{G, \bullet}-f_{H^{\prime \prime}, \bullet}\right\|_{\nabla, S^{M}} \leq \frac{5 a \alpha}{\sqrt{k \wedge 1}}(2 \pi \underline{k})^{-d / 2}=a \eta
$$

The proof is completed following same steps as the proof of Lemma 5.

### 4.6.3.4 Putting together the pieces

Combining (4.39), (4.42), (4.45), (4.46), and (4.49) and pulling out any constants depending on $d, \underline{k}$, or $\bar{k}$,

$$
\begin{align*}
& \mathbb{E}\left\|\hat{\theta}-\hat{\theta}^{*}\right\|_{F}^{2} / n \leq(5 / n)[ \left.\mathbb{E}\left\|\hat{\theta}_{\rho}^{*}-\hat{\theta}^{*}\right\|_{F}^{2}+\sum_{t=1}^{4} \mathbb{E} \zeta_{t}^{2}\right] \\
& \leq c_{d, k, \bar{k}}( \varepsilon_{n}^{2}\left(M, S, G^{*}\right)(\sqrt{\log n})^{d-2} \\
&+\frac{\log n}{n}  \tag{4.53}\\
&+\left(\eta^{*}\right)^{2}+\varepsilon_{n}^{2}\left(M, S, G^{*}\right) \\
&+\log (e N) \frac{(\log n)^{2}}{n} \\
&\left.+\varepsilon_{n}^{2}\left(M, S, G^{*}\right)(\sqrt{\log n})^{6}\right) \\
& \leq c_{d, k}, \bar{k} \varepsilon_{n}^{2}\left(M, S, G^{*}\right)(\sqrt{\log n})^{(d-2) \vee 6}
\end{align*}
$$

This completes the proof of Theorem 4.8.

### 4.6.4 Proofs of Theorems 4.10 and 4.12

Proof of Theorem 4.10. We will relate the Wasserstein distance to the average Hellinger distance, so we rely on the tools of Nguyen (2013, proof of theorem 2). Fix a symmetric
density $K$ whose Fourier transform $\widetilde{K}$ is bounded with support on $[-1,1]^{d}$. For any $\delta>0$ define the scaled kernel $K_{\delta}(x)=\frac{1}{\delta^{d}} K(x / \delta)$. By the triangle inequality,

$$
W_{2}\left(G^{*}, \widehat{G}_{n}\right) \leq W_{2}\left(G^{*}, G^{*} * K_{\delta}\right)+W_{2}\left(G^{*} * K_{\delta}, \widehat{G}_{n} * K_{\delta}\right)+W_{2}\left(\widehat{G}_{n} * K_{\delta}, \widehat{G}_{n}\right)
$$

For the first and third terms, bound the minimum over all couplings by the strong coupling:

$$
W_{2}^{2}\left(G, G * K_{\delta}\right)=\min _{\theta \sim G, \theta^{\prime} \sim G, \varepsilon \sim K} \mathbb{E}\left\|\theta-\left(\theta^{\prime}+\delta \varepsilon\right)\right\|_{2}^{2} \leq \delta^{2} \mathbb{E}_{\varepsilon \sim K}\|\varepsilon\|_{2}^{2}
$$

where the inequality follows from choosing the coupling where $\theta=\theta^{\prime}$ almost surely. Letting $m_{2}(K)=\mathbb{E}_{\varepsilon \sim K}\|\varepsilon\|_{2}^{2}$ denote the second moment of the (unscaled) kernel, we have

$$
\begin{equation*}
W_{2}\left(G^{*}, \widehat{G}_{n}\right) \leq 2 \sqrt{m_{2}(K)} \delta+W_{2}\left(G^{*} * K_{\delta}, \widehat{G}_{n} * K_{\delta}\right) \tag{4.54}
\end{equation*}
$$

For the second term, Villani (2008, Theorem 6.15) yields

$$
W_{2}^{2}\left(G^{*} * K_{\delta}, \widehat{G}_{n} * K_{\delta}\right) \leq 2 \int\|x\|_{2}^{2} \mathrm{~d}\left|G^{*} * K_{\delta}-\widehat{G}_{n} * K_{\delta}\right|(x)
$$

By Nguyen (2013, Lemma 6), for any $s>2$ such that $m_{s}(K)=\mathbb{E}\|\varepsilon\|_{2}^{s}<\infty$,

$$
\begin{aligned}
W_{2}^{2}\left(G^{*} * K_{\delta}, \widehat{G}_{n} * K_{\delta}\right) \leq & 4\left\|G^{*} * K_{\delta}-\widehat{G}_{n} * K_{\delta}\right\|_{L_{1}}^{(s-2) / s} R^{2 / s} \\
\leq & 4\left[2 \operatorname{Vol}\left(B_{1}\right)^{s /(d+2 s)} R^{d /(d+2 s)}\left\|G^{*} * K_{\delta}-\widehat{G}_{n} * K_{\delta}\right\|_{L_{2}}^{2 s /(d+2 s)}\right]^{(s-2) / s} R^{2 / s} \\
= & 4 \cdot 2^{(s-2) / s} \operatorname{Vol}\left(B_{1}\right)^{(s-2) /(2 s+d)} \\
& \times R^{d(s-2) /(s(d+2 s))+2 / s}\left\|G^{*} * K_{\delta}-\widehat{G}_{n} * K_{\delta}\right\|_{L_{2}}^{2(s-2) /(d+2 s)} \\
\leq & 8 \sqrt{\operatorname{Vol}\left(B_{1}\right)} \cdot R^{d(s-2) /(s(d+2 s))+2 / s}\left\|G^{*} * K_{\delta}-\widehat{G}_{n} * K_{\delta}\right\|_{L_{2}}^{2(s-2) /(d+2 s)}
\end{aligned}
$$

where $R:=\mathbb{E}_{\theta^{*} \sim G^{*}, \varepsilon \sim K}\left\|\theta^{*}+\delta \varepsilon\right\|_{2}^{s}+\mathbb{E}_{\theta \sim \widehat{G}_{n}, \varepsilon \sim K}\|\theta+\delta \varepsilon\|_{2}^{s}$.
For moments in the term $R$, use $\mathbb{E}\|\theta+\delta \varepsilon\|_{2}^{s} \leq 2^{s}\left(\mathbb{E}\|\theta\|_{2}^{s}+\delta^{s} m_{s}(K)\right)$, so

$$
R \leq 2^{s}\left(m_{s}\left(G^{*}\right)+m_{s}\left(\widehat{G}_{n}\right)+2 \delta^{s} m_{s}(K)\right) .
$$

The quantity $m_{s}(K)$ is regarded as a constant depending only on $s>2$ and $d$. By assumption, the support of $\widehat{G}_{n}$ is contained in the minimum bounding box of the observations, which is further contained in $[-U, U]^{d}$ where $U=\max _{i, j}\left|X_{i j}\right| \leq L+\max _{i, j}\left|X_{i j}-\theta_{i j}^{*}\right|$. Since $X_{i j}-\theta_{i j}^{*} \stackrel{\text { ind }}{\sim} \mathcal{N}\left(0,\left(\Sigma_{i}\right)_{j j}\right)$, we have by a standard concentration argument that

$$
U \leq L+4 \sqrt{\bar{k} \log n}
$$

with probability at least $1-\frac{2 d}{n^{8}}$. Hence, with the same probability

$$
m_{s}\left(\widehat{G}_{n}\right)=\mathbb{E}_{\widehat{G}_{n}}\|\theta\|_{2}^{s} \leq d^{s / 2} \mathbb{E}_{\widehat{G}_{n}}\|\theta\|_{\infty}^{s} \leq d^{s / 2}(L+4 \sqrt{\bar{k} \log n})^{s}
$$

This same bound holds for $m_{s}\left(G^{*}\right)$.
For the $\|\cdot\|_{L_{2}}$ norm $\left\|G^{*} * K_{\delta}-\widehat{G}_{n} * K_{\delta}\right\|_{L_{2}}$, let $g_{\delta}^{(i)}$ denote the inverse Fourier transform of $\widetilde{K_{\delta}} / \widetilde{\varphi_{\Sigma_{i}}}$, so that $G * K_{\delta}=f_{G, \Sigma_{i}} * g_{\delta}^{(i)}$. Hence, by Proposition 8.49 of Folland (1999), we have for each $i=1, \ldots, n$,

$$
\left\|G^{*} * K_{\delta}-\widehat{G}_{n} * K_{\delta}\right\|_{L_{2}} \leq 2 d_{T V}\left(f_{G^{*}, \Sigma_{i}}, f_{\widehat{G}_{n}, \Sigma_{i}}\right)\left\|g_{\delta}^{(i)}\right\|_{L_{2}} .
$$

Using Plancherel's theorem and the fact that $\widetilde{K}$ is bounded on its support of $[-1,1]^{d}$,

$$
\begin{aligned}
\left\|g_{\delta}^{(i)}\right\|_{L_{2}}^{2} & =\int_{\mathbb{R}^{d}} \frac{\widetilde{K}(\delta \omega)^{2}}{\widetilde{\varphi_{\Sigma_{i}}}(\omega)^{2}} \mathrm{~d} \omega \leq C_{d} \int_{[-1 / \delta, 1 / \delta]^{d}} \widetilde{\varphi_{\Sigma_{i}}}(\omega)^{-2} \mathrm{~d} \omega \\
& =C_{d} \int_{[-1 / \delta, 1 / \delta]^{d}} \exp \left(\omega^{\prime} \Sigma_{i} \omega\right) \mathrm{d} \omega=C_{d} \prod_{j=1}^{d} \int_{-1 / \delta}^{1 / \delta} \exp \left(\left(\Sigma_{i}\right)_{j j} \omega_{j}^{2}\right) \mathrm{d} \omega_{j} \\
& \leq C_{d} \prod_{j=1}^{d} e^{2\left(\Sigma_{i}\right)_{j j} \delta^{-2}} \int_{-1 / \delta}^{1 / \delta} \exp \left(-\left(\Sigma_{i}\right)_{j j} \omega_{j}^{2}\right) \mathrm{d} \omega_{j} \leq C_{d}\left(\frac{\pi}{\underline{k}}\right)^{d / 2} e^{2 d \bar{k} \delta^{-2}} .
\end{aligned}
$$

Averaging over $i=1, \ldots, n$,

$$
\begin{aligned}
\left\|G^{*} * K_{\delta}-\widehat{G}_{n} * K_{\delta}\right\|_{L_{2}} & \leq C_{d} \underline{k}^{-d / 2} e^{2 d \bar{k} \delta^{-2}} \cdot \frac{1}{n} \sum_{i=1}^{n} d_{T V}\left(f_{G^{*}, \Sigma_{i}}, f_{\widehat{G}_{n}, \Sigma_{i}}\right) \\
& \leq C_{d} \underline{k}^{-d / 2} e^{2 d \bar{k} \delta^{-2}} \bar{h}\left(G^{*}, \widehat{G}_{n}\right)
\end{aligned}
$$

Combining our calculations following (4.54), we have

$$
\begin{align*}
W_{2}\left(G^{*}, \widehat{G}_{n}\right) \leq C_{d, s} \inf _{\delta \in(0,1)}\{\delta & \left(2^{s}\left(d^{s / 2}(L+4 \sqrt{\bar{k} \log n})^{s}+2 \delta^{s} m_{s}(K)\right)\right)^{3 d /(2 d+4 s)}  \tag{4.55}\\
& \left.\times\left(\underline{k}^{-d / 2} e^{2 d \bar{k} \delta^{-2}} \bar{h}\left(G^{*}, \widehat{G}_{n}\right)\right)^{(s-2) /(d+2 s)}\right\}
\end{align*}
$$

Assume $n$ is large enough that $4 \sqrt{\bar{k} \log n} \geq L$ and $2^{s} d^{s / 2}(4 \sqrt{\bar{k} \log n})^{s} \geq 2 m_{s}(K)$, so

$$
\begin{equation*}
W_{2}\left(G^{*}, \widehat{G}_{n}\right) \leq C_{d, s} \inf _{\delta \in(0,1)}\left\{\delta+(\bar{k} \log n)^{3 s d /(4(d+2 s))}\left(\underline{k}^{-d / 2} e^{2 d \bar{k} \delta-2} \bar{h}\left(G^{*}, \widehat{G}_{n}\right)\right)^{(s-2) /(d+2 s)}\right\} . \tag{4.56}
\end{equation*}
$$

Choosing $\delta^{-2}=-\frac{1}{4 \bar{k} d} \log \bar{h}\left(G^{*}, \widehat{G}_{n}\right)($ provided $\delta<1)$ and $s=d+2$ yields

$$
\begin{equation*}
W_{2}^{2}\left(G^{*}, \widehat{G}_{n}\right) \leq C_{d}\left\{\frac{\bar{k} d}{-\log \bar{h}\left(G^{*}, \widehat{G}_{n}\right)}+(\bar{k} \log n)^{d / 2}\left(\underline{k}^{-d} \bar{h}\left(G^{*}, \widehat{G}_{n}\right)\right)^{1 / 12}\right\} \tag{4.57}
\end{equation*}
$$

$\varepsilon_{n}^{2}\left(M, S, G^{*}\right)$ defined in (4.13), with $S=[-L, L]^{d}$ and $M=\sqrt{10 \bar{k} \log n}$ gives

$$
\varepsilon_{n}^{2}=\frac{\left(4 \sqrt{10}\left(L^{2} \vee \bar{k}\right)\right)^{d}}{n}(\log n)^{d+1}
$$

By Theorem 4.6,
$W_{2}^{2}\left(G^{*}, \widehat{G}_{n}\right) \leq C_{d}\left\{\frac{\bar{k} d}{\log n-\log C_{d, \bar{k}, \underline{k}, L} t^{2}(\log n)^{d+1}}+(\bar{k} \log n)^{d / 2}\left(C_{d, \bar{k}, \underline{k}, L} t^{2} \frac{(\log n)^{d+1}}{n}\right)^{1 / 24}\right\}$,
with probability at least $1-2 n^{-t^{2}}$. Take $t^{2}=8$. For $n$ sufficiently large the first term dominates, $\delta<1$, and $\log n-\log C_{d, \bar{k}, k, L} 8(\log n)^{d+1} \geq(\log n) / 2$.
Proof of Theorem 4.12. Take $\mu=0$ by location equivariance (see Lemma 4.4). Write $\widehat{G}_{n}=$ $\sum_{j=1}^{\hat{k}} \hat{w}_{j} \delta_{\hat{a}_{j}}$. Since $G^{*}=\delta_{0}$ is a point mass,

$$
W_{2}^{2}\left(\widehat{G}_{n}, G^{*}\right)=\mathbb{E}_{\vartheta \sim \sim \widehat{G}_{n}}\|\vartheta\|_{2}^{2}=\sum_{j=1}^{\hat{k}} \hat{w}_{j}\left\|\hat{a}_{j}\right\|_{2}^{2}
$$

We relate this to the marginal density $f_{\widehat{G}_{n}, \bullet}$ via

$$
\begin{aligned}
\int f_{\widehat{G}_{n}, \Sigma_{i}}(x)\|x\|_{2}^{2} \mathrm{~d} x & =\sum_{j=1}^{\hat{k}} \hat{w}_{j} \int \varphi_{\Sigma_{i}}\left(x-\hat{a}_{j}\right)\|x\|_{2}^{2} \mathrm{~d} x \\
& =\sum_{j=1}^{\hat{k}} \hat{w}_{j} \int \varphi_{\Sigma_{i}}(x)\left(\|x\|_{2}^{2}+\left\|\hat{a}_{j}\right\|_{2}^{2}+2\left\langle x, \hat{a}_{j}\right\rangle\right) \mathrm{d} x \\
& =\int f_{G^{*}, \Sigma_{i}}(x)\|x\|_{2}^{2} \mathrm{~d} x+W_{2}^{2}\left(\widehat{G}_{n}, G^{*}\right)
\end{aligned}
$$

Hence for any $i \in\{1, \ldots, n\}$,

$$
\begin{aligned}
W_{2}^{2}\left(\widehat{G}_{n}, G^{*}\right) & =\int\left(f_{\widehat{G}_{n}, \Sigma_{i}}(x)-f_{G^{*}, \Sigma_{i}}(x)\right)\|x\|_{2}^{2} \mathrm{~d} x \\
& \leq 2 h\left(f_{\widehat{G}_{n}, \Sigma_{i}}, f_{G^{*}, \Sigma_{i}}\right)\left(\int\left(f_{\widehat{G}_{n}, \Sigma_{i}}(x)+f_{G^{*}, \Sigma_{i}}(x)\right)\|x\|_{2}^{4} \mathrm{~d} x\right)^{1 / 2} .
\end{aligned}
$$

Averaging over $i \in\{1, \ldots, n\}$,

$$
W_{2}^{2}\left(\widehat{G}_{n}, G^{*}\right) \leq 2 \bar{h}\left(f_{\widehat{G}_{n}, \bullet}, f_{G^{*}, \bullet}\right) \max _{i=1: n}\left(\int\left(f_{\widehat{G}_{n}, \Sigma_{i}}(x)+f_{G^{*}, \Sigma_{i}}(x)\right)\|x\|_{2}^{4} \mathrm{~d} x\right)^{1 / 2} .
$$

Applying Theorem 4.6 with $S=\{0\}$ and $M=\sqrt{10 \bar{k} \log n}$,

$$
\bar{h}^{2}\left(f_{\widehat{G}_{n}, \bullet}, f_{G^{*}}, \bullet\right) \lesssim_{d, \bar{k}, \underline{k}} t^{2} \frac{(\log n)^{d+1}}{n}
$$

with probability at least $1-2 n^{-t^{2}}$ for all $t \geq 1$.
For the remaining terms,

$$
\int f_{G^{*}, \Sigma_{i}}(x)\|x\|_{2}^{4} \mathrm{~d} x=\mathbb{E}\left\|X_{i}\right\|_{2}^{4}=\mathbb{E}_{Z \sim \mathcal{N}\left(0, I_{d}\right)}\left\|\Sigma_{i}^{1 / 2} Z\right\|_{2}^{4} \leq \bar{k}^{2} \mathbb{E}_{A \sim \chi_{d}^{2}} A^{2}=\bar{k}^{2} d(d+2)
$$

and

$$
\begin{aligned}
\int f_{\widehat{G}_{n}, \Sigma_{i}}(x)\|x\|_{2}^{4} \mathrm{~d} x & =\sum_{j} \hat{w}_{j} \int \varphi_{\Sigma_{i}}(x)\left\|x+\hat{a}_{j}\right\|_{2}^{4} \mathrm{~d} x \\
& \leq 8 \sum_{j} \hat{w}_{j} \int \varphi_{\Sigma_{i}}(x)\left(\|x\|_{2}^{4}+\left\|\hat{a}_{j}\right\|_{2}^{4}\right) \mathrm{d} x \\
& \leq 8 \bar{k}^{2} d(d+2)+8 \sum_{j} \hat{w}_{j}\left\|\hat{a}_{j}\right\|_{2}^{4} .
\end{aligned}
$$

By our assumption on the support that each $\hat{a}_{j} \in \mathbb{B}_{\kappa r}(\bar{X})$, each $\hat{a}_{j}$ equivalently satisfies

$$
\left\|\hat{a}_{j}-\bar{X}\right\|_{2}^{4} \leq(\bar{k} / \underline{k})^{4} \max _{i}\left\|X_{i}-\bar{X}\right\|_{2}^{4} .
$$

Noting that $X_{i}-\bar{X} \sim \mathcal{N}\left(0,\left(1-2 n^{-1}\right) \Sigma_{i}+n^{-1} \bar{\Sigma}\right)$ with $\bar{\Sigma}=n^{-1} \sum_{j=1}^{n} \Sigma_{j}$, we bound

$$
\begin{aligned}
\sum_{j} \hat{w}_{j}\left\|\hat{a}_{j}\right\|_{2}^{4} & \leq 8\left(\|\bar{X}\|_{2}^{4}+\max _{j}\left\|\hat{a}_{j}-\bar{X}\right\|_{2}^{4}\right) \\
& \leq 8\left(\|\bar{X}\|_{2}^{4}+\kappa^{4} \max _{i \in[n]}\left\|X_{i}-\bar{X}\right\|_{2}^{4}\right) \\
& \leq{ }_{\text {st }} 8 \bar{k}^{2}\left(n^{-2} A_{0}^{2}+\kappa^{4} \max _{i \in[n]} A_{i}^{2}\right) \\
& \leq 16 \frac{\bar{k}^{6}}{\underline{k}^{4}} \max _{i=0: n} A_{i}^{2}
\end{aligned}
$$

where $A_{0}, A_{1}, \ldots, A_{n} \sim \chi_{d}^{2}$ are possibly dependent, and $\leq_{\text {st }}$ denotes stochastic inequality.

For $t \geq 1$, we use the following tail bound (see Laurent \& Massart, 2000, Lemma 1)

$$
\mathbb{P}\left(\max _{i=0: n} A_{i}^{2} \geq 60 t^{2}(\log n)^{2}\right) \leq n^{-t^{2}}
$$

where we have used the assumption in Theorem 4.6 that $n \geq(2 \pi)^{d / 2}$ to eliminate the dependence on $d$. We have thus shown that

$$
\begin{aligned}
& \max _{i=1: n}\left(\int\left(f_{\widehat{G}_{n}, \Sigma_{i}}(x)+f_{G^{*}, \Sigma_{i}}(x)\right)\|x\|_{2}^{4} \mathrm{~d} x\right)^{1 / 2} \\
& \quad \leq\left(9 \bar{k}^{2} d(d+2)+400 \frac{\bar{k}^{6}}{\underline{k}^{4}} t^{2}(\log n)^{2}\right)^{1 / 2} \lesssim \frac{\bar{k}^{3}}{\underline{k}^{2}} t(\log n)
\end{aligned}
$$

with probability at least $1-n^{-t^{2}}$. Combining with a union bound over our earlier estimate,

$$
W_{2}\left(\widehat{G}_{n}, G^{*}\right) \lesssim_{d, \bar{k}, \underline{,}} t^{3 / 2} \frac{(\log n)^{(d+3) / 4}}{n^{1 / 4}}
$$

with probability at least $1-3 n^{-t^{2}}$ for all $t \geq 1$.

## Chapter 5

## Local false discovery rate control

### 5.1 Introduction

A common goal in applications of multiple hypothesis testing is to identify a relatively short list of candidate "discoveries" that are sufficiently promising to undertake some costly further action. In scientific applications, for example, each discovery may be the focus of a follow-up experiment, which wastes resources if the apparent discovery was only a mirage. The false discovery rate (FDR, Benjamini \& Hochberg, 1995) has become a cornerstone of modern large-scale multiple testing because it directly measures the rate of this wastage:
[T]he proportion of errors in the pool of candidates is of great economical significance since follow-up studies are costly, and thus avoiding multiplicity control is costly. Indeed, the FDR criterion is economically interpretable; when considering a potential threshold, the adjusted FDR gives the proportion of the investment that is about to be wasted on false leads. (Reiner et al., 2003)

An analyst who controls FDR at level $q=5 \%$, then, is willing to waste resources following up on one false discovery in exchange for every nineteen real discoveries.

Carrying this reasoning further, however, we can apply the same cost-benefit analysis to each individual rejection, not only to the list of rejections taken as a whole. In economic terminology, we should consider not only the average utility of our entire rejection set, but also the marginal utility of each rejection we make, since we always have the option to exclude any rejection that is not individually promising. For example, in Section 5.4 we reproduce the simulations of Benjamini and Hochberg (1995) and find in some settings that, even while the Benjamini-Hochberg (BH) procedure controls FDR at level $q=5 \%$, the last discovery (i.e. the discovery with the largest $p$-value) is false more than $30 \%$ of the time. In such settings, unless we are willing to suffer one false discovery for every two true discoveries, we would be better served by excluding the last rejection from the BH rejection set. More generally, to decide where to set our rejection threshold, we should ask about the proportion of false leads among the incremental rejections that we would add or remove by raising or lowering it.

The likelihood that an individual discovery is a false lead is called its local false discovery rate (lfdr, Efron et al., 2001). For $i=1, \ldots, m$, let $H_{i}=0$ if the $i$ th hypothesis is null and $H_{i}=1$ otherwise, and consider the simple Bayesian two-groups model

$$
\begin{equation*}
p_{i} \mid H_{i}=h \stackrel{\text { ind }}{\sim} f_{h}, \quad \text { with } \quad H_{i} \stackrel{\text { iid }}{\sim} \operatorname{Bern}\left(1-\pi_{0}\right), \quad \text { for } i=1, \ldots, m, \tag{5.1}
\end{equation*}
$$

where $f_{0}:=1_{[0,1]}$ and $f_{1}$ are densities (null and alternative, respectively) supported on the unit interval $[0,1]$, and the null proportion is $\pi_{0} \in[0,1]$. Let $f:=\pi_{0}+\left(1-\pi_{0}\right) f_{1}$ denote the common mixture density of the $p$-values in model (5.1), and let $F(t):=\int_{0}^{t} f(u) d u$ denote the corresponding cumulative distribution function (cdf). The lfdr is then defined as the posterior probability that $H_{i}=0$, conditional on the observed $p$-value $p_{i}$ :

$$
\begin{equation*}
\operatorname{lfdr}(t):=\mathbb{P}\left(H_{i}=0 \mid p_{i}=t\right)=\frac{\pi_{0}}{f(t)} \tag{5.2}
\end{equation*}
$$

If we knew the problem parameters $\pi_{0}$ and $f_{1}$, then the definition (5.2) would neatly solve the problem posed above: we should reject only those hypotheses whose lfdr is below the breakeven threshold of our cost-benefit tradeoff. Concretely, let $\lambda>0$ define the ratio between the cost of each false discovery and the benefit of each true discovery. Then the utility of making $R$ rejections, of which $V$ are false discoveries, is proportional to $(R-V)-\lambda V$, and a simple calculation shows that we should reject the $i$ th hypothesis if and only if $\operatorname{lfdr}\left(p_{i}\right) \leq \alpha:=\frac{1}{1+\lambda}$.

We will usually work under the additional assumption that $f_{1}(t)$ is nonincreasing in $t$, or equivalently that $\operatorname{lfdr}(t)$ is nondecreasing, so that smaller $p$-values represent stronger evidence against the null. This assumption, common in multiple testing (see, e.g., Genovese \& Wasserman, 2004; Langaas et al., 2005; Strimmer, 2008), lets us restrict our attention to procedures that reject all $p$-values below a given threshold: if $f_{1}$ is nonincreasing then rejecting when $\operatorname{lfdr}\left(p_{i}\right) \leq \alpha$ is equivalent to rejecting when $p_{i}$ is sufficiently small.

In practice, $\pi_{0}$ and $f_{1}$ are typically unknown and must be estimated from the data, and many estimators have been proposed; see e.g. Aubert et al. (2004), Efron (2004, 2008), Efron et al. (2001), Liao et al. (2004), Muralidharan (2010), Patra and Sen (2016), Pounds and Cheng (2004), Pounds and Morris (2003), Robin et al. (2007), Scheid and Spang (2004), Stephens (2017), and Strimmer (2008). To the best of our knowledge, however, there are no known finite-sample lfdr control guarantees for multiple testing procedures based on these methods. By contrast, simple, robust, and well-known methods like the Benjamini-Hochberg (BH) procedure of Benjamini and Hochberg (1995) enjoy finite-sample FDR control without requiring the analyst to model the $p$-value distribution.

In this chapter, we introduce a new error control metric that measures the lfdr of a multiple testing procedure's least promising rejection. We represent a generic multiple testing method as a function $\mathcal{R}\left(p_{1}, \ldots, p_{m}\right)$ returning an index set $\mathcal{R} \subseteq\{1, \ldots, m\}$, where hypothesis $i$ is rejected if and only if $i \in \mathcal{R}$. We say the procedure's max-lfdr is

$$
\begin{equation*}
\max -\operatorname{lfdr}(\mathcal{R}):=\mathbb{E}\left[\max _{i \in \mathcal{R}} \operatorname{lfdr}\left(p_{i}\right)\right] \tag{5.3}
\end{equation*}
$$

defining the maximum as zero if no rejections are made. If $f_{1}$ is nonincreasing, then the max-lfdr of $\mathcal{R}$ coincides with the probability that the last rejection is a false discovery.

We also introduce a simple multiple testing procedure, which we call the support line (SL) procedure, that provably controls the max-lfdr under mild assumptions. Define the $p$-value order statistics $p_{(1)} \leq \cdots \leq p_{(m)}$, and let $p_{(0)}=0$ by convention. Then our procedure rejects $p$-values up to the last (and a.s. unique) minimizer

$$
\begin{equation*}
R_{q}:=\underset{k=0, \ldots, m}{\operatorname{argmin}} p_{(k)}-\frac{q k}{m} . \tag{5.4}
\end{equation*}
$$

That is, we reject $\mathcal{R}_{q}:=\left\{i: p_{i} \leq \tau_{q}\right\}$, for the threshold $\tau_{q}=p_{\left(R_{q}\right)}$. Under the two-groups model (5.1), with nonincreasing $f_{1}$, we show in Theorem 5.1 that

$$
\max -\operatorname{lfdr}\left(\mathcal{R}_{q}\right)=\pi_{0} q
$$

Our method can be implemented without knowing $\pi_{0}$ or $f_{1}$, apart from the shape constraint, and bears a close relationship to the BH procedure, which replaces $R_{q}$ in (5.4) with

$$
R_{q}^{\mathrm{BH}}:=\max \left\{k \in\{0, \ldots, m\}: p_{(k)} \leq \frac{q k}{m}\right\}
$$

rejecting $\mathcal{R}_{q}^{\mathrm{BH}}:=\left\{i: p_{i} \leq \tau_{q}^{\mathrm{BH}}\right\}$, for $\tau_{q}^{\mathrm{BH}}=q R_{q}^{\mathrm{BH}} / m \geq p_{\left(R_{q}^{\mathrm{BH})}\right.}$. Because $R_{q} \leq R_{q}^{\mathrm{BH}}$, the BH method makes at least as many rejections as the SL method, and both methods make at least one rejection if and only if $p_{(k)} \leq \frac{q k}{m}$ for some $k \geq 1$; however, as we will argue, in general, the SL method should be run with a strictly larger $q$ than we would use for BH. The left panel of Figure 5.1 illustrates the relationship between the two methods by reproducing the familiar plot of the BH procedure as an operation on the order statistics $p_{(1)}, \ldots, p_{(m)}$.

### 5.1.1 Multiple testing and the weighted classification loss

To formalize our analysis above, define the per-instance weighted classification loss:

$$
\begin{equation*}
L_{\lambda}(H, \mathcal{R}):=\frac{(1+\lambda) V-R}{m} \tag{5.5}
\end{equation*}
$$

This loss can be derived, up to additive and multiplicative constants, by viewing each of the $m$ hypotheses as a binary classification problem, where we incur a cost $c_{1}$ for each type I error or false discovery ( $i \in \mathcal{R}$, but $H_{i}=0$ ), and cost $c_{2}$ from each type II error or false non-discovery $\left(i \notin \mathcal{R}\right.$, but $\left.H_{i}=1\right)$. If the total number of non-nulls is $m_{1}=\sum_{i} H_{i}$, then there are $m_{1}-(R-V)$ false non-discoveries, so the total loss over all $m$ instances is

$$
c_{1} V+c_{2}\left(m_{1}-(R-V)\right)=c_{2} m \cdot L_{\lambda}(H, \mathcal{R})+c_{2} m_{1},
$$

where $\lambda=c_{1} / c_{2}$ is the ratio between the two misclassification costs. $L_{\lambda}$ as defined in (5.5) is normalized so that rejecting nothing incurs zero loss, and each true discovery has value $1 / \mathrm{m}$.


Figure 5.1: Left: The order statistics $p_{(k)}$ of the $p$-values as a function of the index $k$, shown in black. The BH procedure, in red, finds the largest index $R_{q}^{\mathrm{BH}}$ such that $p_{\left(R_{q}^{\mathrm{BH})}\right.}$ falls below the ray of slope $q / m$; by contrast, our procedure finds the (last and almost surely unique) boundary point $\left(R_{q}, p_{\left(R_{q}\right)}\right)$ of the supporting line of slope $q / m$. Right: The same plot with the ray through the origin of slope $q / m$ subtracted off. The black dots represent a running estimate (5.7) of the weighted classification loss (5.5), which our procedure minimizes. $\mathrm{BH}(q)$ finds the largest threshold where the estimated loss is negative.

Under the two-groups model (5.1), Sun and Cai (2007, Theorem 2) show that the corresponding Bayes risk $\mathbb{E} L_{\lambda}(H, \mathcal{R})$ is minimized by the oracle procedure

$$
\begin{equation*}
\mathcal{R}^{*}:=\left\{i: \operatorname{lfdr}\left(p_{i}\right) \leq \alpha\right\}, \quad \text { where } \quad \alpha=\frac{1}{1+\lambda} \tag{5.6}
\end{equation*}
$$

The ratio $\lambda$ specifies the "break-even exchange rate" at which we are willing to trade true discoveries for false leads; e.g., if $\lambda=19$ then we are willing to suffer a single false discovery for exactly 19 true discoveries, and we should reject a hypothesis only if its lfdr falls below the break-even tolerance $\alpha=0.05$. If $f_{1}$ is nonincreasing, then the oracle procedure reduces to thresholding $p$-values at a fixed threshold

$$
\mathcal{R}^{*}=\left\{i: p_{i} \leq \tau^{*}\right\}, \quad \text { for } \quad \tau^{*}:=\max \{t \in[0,1]: \operatorname{lfdr}(t) \leq \alpha\}
$$

with $\tau^{*}=0$ if no such threshold exists.
Our method can be directly interpreted as minimizing an empirical proxy of the weighted classification loss. For a candidate threshold $t \in[0,1]$, the expected number of null $p$-values
below the threshold is $m \pi_{0} t$. If $\pi_{0}$ is known, we obtain a running estimator of the loss:

$$
\begin{equation*}
\hat{L}_{\lambda}\left(t ; \pi_{0}\right)=\frac{(1+\lambda) m \pi_{0} t-m F_{m}(t)}{m}=(1+\lambda)\left(\pi_{0} t-\alpha F_{m}(t)\right), \tag{5.7}
\end{equation*}
$$

where $F_{m}(t)$ represents the empirical cumulative distribution function (ecdf) of the $p$-values:

$$
F_{m}(t)=\frac{1}{m} \sum_{i=1}^{m} 1\left\{p_{i} \leq t\right\}
$$

Because $\hat{L}_{\lambda}\left(t ; \pi_{0}\right)$ is increasing between successive order statistics, it is minimized at one of the order statistics, or at $p_{(0)}=0$ :

$$
\underset{k=0,1, \ldots, m}{\operatorname{argmin}} \hat{L}_{\lambda}\left(p_{(k)} ; \pi_{0}\right)=\underset{k=0,1, \ldots, m}{\operatorname{argmin}} \pi_{0} p_{(k)}-\frac{\alpha k}{m} .
$$

Comparing the last expression to the definition of our procedure in (5.4), we see that $\hat{L}_{\lambda}\left(t ; \pi_{0}\right)$ is minimized at $t=\tau_{q}$ for $q=\alpha / \pi_{0}$. By Theorem 5.1 , we then have exactly max- $\operatorname{lfdr}\left(\mathcal{R}_{q}\right)=\alpha$.

By contrast, $\tau_{q}^{\mathrm{BH}}$ for $q=\alpha / \pi_{0}$ is the largest value of $t$ that gives $\hat{L}_{\lambda}\left(t ; \pi_{0}\right)=0$, the same loss we would achieve by rejecting nothing at all. In other words, the BH procedure at level $\alpha / \pi_{0}$ only aims to break even; to do better, we should run BH at a strictly smaller level $q<\alpha / \pi_{0}$, viewing $q$ as a tuning parameter as in Neuvial and Roquain (2012).

To select $q$ for our SL procedure when $\pi_{0}$ is unknown, we can either conservatively bound $\pi_{0} \leq 1$ and run the procedure at $q=\alpha$, or estimate $\pi_{0}$ and use $q=\alpha / \hat{\pi}_{0}$. To avoid confusion, we will always use the notation $q$ to represent our method's tuning parameter, and reserve $\alpha=\frac{1}{1+\lambda}$ to represent the true target lfdr, defined in terms of the cost ratio $\lambda$.

Our procedure can alternatively be derived as a plug-in maximum likelihood estimator (MLE) of the oracle procedure $\mathcal{R}^{*}$, where we estimate $f(t)$ using Grenander's nonparametric MLE for a nonincreasing density (Grenander, 1956):

$$
\begin{equation*}
\hat{f}_{m}:=\underset{\substack{g:[0,1] \rightarrow \mathbb{R}_{+} \\ \text {nonincreasing density }}}{\operatorname{argmax}} \frac{1}{m} \sum_{i=1}^{m} \log g\left(p_{i}\right) . \tag{5.8}
\end{equation*}
$$

As we will see in Section 5.3.2, $\tau_{q}$ is also the largest value $t \in[0,1]$ for which $\hat{f}_{m}(t) \geq q^{-1}$. Thus, if we run our procedure at $q=\alpha / \pi_{0}$, we have

$$
\mathcal{R}_{\alpha / \pi_{0}}=\left\{i: \hat{f}_{m}\left(p_{i}\right) \geq\left(\alpha / \pi_{0}\right)^{-1}\right\}=\left\{i: \frac{\pi_{0}}{\hat{f}_{m}\left(p_{i}\right)} \leq \alpha\right\}
$$

As above, if $\pi_{0}$ is unknown, we can either estimate it or conservatively bound $\pi_{0} \leq 1$.
The relationship between our method and the Grenander estimator is convenient for asymptotic analysis because the latter is very well studied; see the book by Groeneboom
and Jongbloed (2014) for a thorough treatment. The Grenander estimator has previously been considered for estimating the lfdr (Strimmer, 2008) as well as for estimating the null proportion $\pi_{0}$ (Langaas et al., 2005). While $\hat{f}_{m}$ may be efficiently computed via the pool adjacent violators algorithm (Robertson et al., 1988), the definition in (5.4) is usually preferred for computational purposes.

### 5.1.2 The max-lfdr and the FDR

The max-lfdr in (5.3) and the FDR are two different error criteria that both appeal to the logic of trading off true and false discoveries. The key difference is that the FDR, defined as

$$
\operatorname{FDR}(\mathcal{R}):=\mathbb{E}\left[\frac{V}{R} \cdot 1\{R>0\}\right],
$$

measures the likelihood that a randomly selected rejection is null, whereas the max-lfdr instead measures the likelihood that the least promising rejection is null. In both cases the event in question is deemed not to have occurred if $R=0$, so that under the global null (all $H_{i}=0$, almost surely), both criteria reduce to the probability of making a single rejection.

Throughout this section, we will restrict our attention to procedures that reject the $R$ hypotheses with the smallest $p$-values. That is, we assume a procedure $\mathcal{R}$ rejects $H_{(1)}, \ldots, H_{(R)}$, where $H_{(k)}$ represents the hypothesis corresponding to $p_{(k)}$. If $f_{1}$ is nonincreasing, then the procedure's last rejection $H_{(R)}$ is the least promising, and the max-lfdr can be equivalently characterized as the probability that the last rejection is a false discovery:

$$
\begin{equation*}
\max -\operatorname{lfdr}(\mathcal{R})=\mathbb{E}\left[\operatorname{lfdr}\left(p_{(R)}\right) \cdot 1\{R>0\}\right]=\mathbb{P}\left\{H_{(R)}=0, R>0\right\} \tag{5.9}
\end{equation*}
$$

If $\max -\operatorname{lfdr}(\mathcal{R})>\alpha=\frac{1}{1+\lambda}$, then we can improve $\mathcal{R}$ by excluding its last discovery. ${ }^{1}$ Let $\mathcal{R}^{-1}$ denote the procedure that makes one fewer rejection than $\mathcal{R}$, meaning it rejects $H_{(1)}, \ldots, H_{(R-1)}$ if $R>0$, and makes no rejections if $R=0$. Then we have

$$
\begin{aligned}
\mathbb{E}\left[L_{\lambda}(H, \mathcal{R})-L_{\lambda}\left(H, \mathcal{R}^{-1}\right)\right] & =\frac{1}{m} \mathbb{E}\left[(1+\lambda) 1\left\{H_{(R)}=0, R>0\right\}-1\{R>0\}\right] \\
& =\frac{1+\lambda}{m}(\max -\operatorname{lfdr}(\mathcal{R})-\alpha \mathbb{P}\{R>0\})
\end{aligned}
$$

which is positive if $\max -\operatorname{lfdr}(\mathcal{R})>\alpha$. The converse, that dropping the last rejection does not improve the risk if $\max -\operatorname{lfdr}(\mathcal{R}) \leq \alpha$, is almost true if $\mathbb{P}\{R>0\} \approx 1$. Under the global null, however, any procedure is improved by making fewer rejections.

This thought experiment - what if we dropped the last rejection? - is at the heart of our motivation for proposing the max-lfdr as an error criterion. Even when a rejection set's

[^0]average quality is high, the rejections near the threshold may be recognizably bad bets. In that case, we are better off "trimming the fat" from our rejection set until all of the rejections that remain are individually worth following up on.

Because max-lfdr $(\mathcal{R}) \leq \operatorname{FDR}(\mathcal{R})$, controlling the max-lfdr is more conservative than controlling FDR at the same level $q$, in most cases considerably so. From this, it is tempting to conclude that max-lfdr control is an inherently more conservative goal than FDR control, but this conclusion would be mistaken. An analyst whose break-even exchange rate is $\lambda=9$ and break-even tolerance is $\alpha=0.1$, for example, would never choose a method with a $10 \%$ FDR; the resulting rejection set would be no better on average than rejecting nothing at all, so there would be no point in collecting the data in the first place. Thus, an analyst who is satisfied with a $10 \%$ FDR must have a larger break-even tolerance, say $\alpha=0.2$ or 0.3 .

By the same token, it would be unfair to evaluate the risk under $L_{\lambda}$ of the BH procedure at level $q=\alpha=\frac{1}{1+\lambda}$, since an analyst whose break-even tolerance is $\alpha$ would want to control FDR at a strictly smaller level $q$, like $\alpha / 2$ or $\alpha / 10$. However, as we show in Section 5.3.1, the performance of $\mathrm{BH}(q)$ with such a priori choices of $q$ can depend sensitively on the unknown alternative density $f_{1}$.

### 5.1.3 Outline and contributions

In Section 5.2, we state and prove our main result, that max-lfdr $\left(\mathcal{R}_{q}\right)=\pi_{0} q$ under the Bayesian two-groups model with nonincreasing $f_{1}$, applying a result of Takács (1967). Even without monotonicity of $f_{1}$, we have $\mathbb{P}\left\{H_{\left(R_{q}\right)}=0, R_{q}>0\right\}=\pi_{0} q$, but monotonicity ensures that the lfdr is not out of control for rejections in the interior of the rejection region. We also prove max-lfdr control for an adaptive method that estimates $\pi_{0}$ from the data in the same way as the procedure of Storey (2002).

In Section 5.3, we investigate our method's asymptotic performance relative to the oracle procedure $\mathcal{R}^{*}$. Extending asymptotic results for the Grenander estimator, we show that our method's attained lfdr threshold, $\operatorname{lfdr}\left(\tau_{q}\right)$, concentrates at a rate $m^{-1 / 3}$ around its expectation $\pi_{0} q$, giving an explicit formula for its asymptotic distribution. We also show that our method's asymptotic regret relative to the oracle shrinks at the rate $m^{-2 / 3}$. Section 5.4 illustrates our results with selected simulations, and Section 5.5 concludes.

### 5.2 Finite-sample max-lfdr control

### 5.2.1 Main result

Our main result is that our procedure $\mathcal{R}_{q}$ controls the max-lfdr at exactly $\pi_{0} q$.
Theorem 5.1. Suppose $p_{1}, \ldots, p_{m}$ follow the Bayesian two-groups model (5.1), with $f_{0}=$ $1_{[0,1]}$. For the procedure defined in (5.4), we have

$$
\begin{equation*}
\mathbb{E}\left[\operatorname{lfdr}\left(p_{\left(R_{q}\right)}\right) \cdot 1\left\{R_{q}>0\right\}\right]=\mathbb{P}\left\{H_{\left(R_{q}\right)}=0, R_{q}>0\right\}=\pi_{0} q . \tag{5.10}
\end{equation*}
$$

If $f_{1}$ is nonincreasing, then we have

$$
\max -\operatorname{lfdr}\left(\mathcal{R}_{q}\right)=\pi_{0} q
$$

The familiar optional-stopping arguments from the FDR control literature, introduced by Storey et al. (2004), do not seem to apply to our procedure, since the minimizer $R_{q}$ of the sequence $p_{(k)}-q k / m$ for $k=0, \ldots, m$ is not a stopping time. We instead prove Theorem 5.1 via a conditioning argument, which crucially relies on the fact that each null $p$-value has exactly a $q / m$ chance of being the last rejection $p_{\left(R_{q}\right)}$ :
Lemma 5.2. Fix $p_{1}, \ldots, p_{m-1} \in[0,1]$ and let $p_{m} \sim \operatorname{Unif}(0,1)$. Then $\mathbb{P}\left\{p_{\left(R_{q}\right)}=p_{m}\right\}=q / m$.
Given Lemma 5.2, the proof of Theorem 5.1 is straightforward:
Proof of Theorem 5.1. Because the ( $H_{i}, p_{i}$ ) pairs are independent and identically distributed, we can decompose the probability in (5.10) as

$$
\begin{aligned}
\mathbb{P}\left\{H_{\left(R_{q}\right)}=0, R_{q}>0\right\} & =\sum_{i=1}^{m} \mathbb{P}\left\{H_{i}=0, p_{\left(R_{q}\right)}=p_{i}\right\} \\
& =m \mathbb{P}\left\{H_{m}=0, p_{\left(R_{q}\right)}=p_{m}\right\} \\
& =\pi_{0} m \mathbb{P}\left\{p_{\left(R_{q}\right)}=p_{m} \mid H_{m}=0\right\} \\
& =\pi_{0} q
\end{aligned}
$$

where the last step comes from conditioning on $p_{1}, \ldots, p_{m-1}$ and applying Lemma 5.2. If $f_{1}(t)$ is nonincreasing, then $\operatorname{lfdr}(t)$ is nondecreasing, so that $\max _{i \in \mathcal{R}_{q}} \operatorname{lfdr}\left(p_{i}\right)=\operatorname{lfdr}\left(p_{\left(R_{q}\right)}\right)$ almost surely, completing the argument.

We now turn to proving Lemma 5.2. Because $p_{m}$ is uniform, the probability statement is equivalent to a showing that, for any fixed $p_{1}, \ldots, p_{m-1} \in[0,1]$, the set of "winning values" $p_{m} \in[0,1]$, for which $\tau_{q}\left(p_{1}, \ldots, p_{m}\right)=p_{m}$, has Lebesgue measure $q / m$. To prove this fact, we rely on a useful result of Takács (1967), which we state next:

Lemma 5.3. Takács, 1967, Theorem 1 Let $\varphi: \mathbb{R}_{+} \rightarrow \mathbb{R}_{+}$denote a nondecreasing step function with $\varphi(0)=0$. Assume that, for some positive $q$, we have $\varphi(u+q)=\varphi(u)+\varphi(q)$ for all $u \geq 0$, and define

$$
\delta(u)=1\{v-\varphi(v) \geq u-\varphi(u) \text { for all } v \geq u\}
$$

Then we have

$$
\int_{0}^{q} \delta(u) \mathrm{d} u=(q-\varphi(q))_{+}
$$

Lemma 5.2 is proved by designing a function $\varphi$ for which the corresponding indicator $\delta\left(p_{m}\right)$ in Lemma 5.3 checks whether $p_{m}$ is the last rejection when we run our method on $\left(p_{i}\right)_{i=1}^{m}$.

Proof of Lemma 5.2. Let $F_{m-1}(t)=\frac{1}{m-1} \sum_{i=1}^{m-1} 1\left\{p_{i} \leq t\right\}$ denote the ecdf of $p_{1}, \ldots, p_{m-1}$, and define a new function $\varphi$ on $[0, q]$

$$
\varphi(v):= \begin{cases}q F_{m-1}(v) \frac{m-1}{m} & \text { if } v<q \\ q \frac{m-1}{m} & \text { if } v=q .\end{cases}
$$

Next, extend $\varphi$ to a nondecreasing step function on all of $\mathbb{R}_{+}$by $\varphi(k q+v)=k \varphi(q)+\varphi(v)$ for all positive integers $k$ and $v \in[0, q]$.

Now let $F_{m}(t)=\frac{1}{m} \sum_{i=1}^{m} 1\left\{p_{i} \leq t\right\}$. If $\tau_{q}=p_{\left(R_{q}\right)}=p_{m}$ then we have $p_{m}-\frac{q R_{q}}{m} \leq$ $p_{(0)}-q \frac{0}{m}=0$, so we may restrict our attention to $p_{m} \leq q$. On the range $v \in\left[p_{m}, q\right)$ we have

$$
m F_{m}(v)=1+(m-1) F_{m-1}(v), \quad \text { so } \quad \varphi(v)=q F_{m}(v)-\frac{q}{m} .
$$

On the range $v \in\left[q, q+p_{m}\right)$, we have

$$
m F_{m}(v-q)=(m-1) F_{m-1}(v-q), \quad \text { so } \quad \varphi(v)=q F_{m}(v-q)-\frac{q}{m}+q
$$

Letting $\delta\left(p_{m}\right):=1\left\{p_{m}=\tau_{q}\left(p_{1}, \ldots, p_{m}\right)\right\}$,

$$
\begin{aligned}
\delta\left(p_{m}\right) & =1\left\{v-q F_{m}(v) \geq p_{m}-q F_{m}\left(p_{m}\right) \text { for all } v \in[0,1]\right\} \\
& =1\left\{v-\varphi(v) \geq p_{m}-\varphi\left(p_{m}\right) \text { for all } v \in\left[p_{m}, q+p_{m}\right)\right\} \\
& =1\left\{v-\varphi(v) \geq p_{m}-\varphi\left(p_{m}\right) \text { for all } v \geq p_{m}\right\}
\end{aligned}
$$

where the last step follows from the fact that $\varphi(v)>\varphi(v-q)$ for all $v \geq q+p_{m}$. We have checked the conditions of Lemma 5.3 Takács, 1967, Theorem 1, from which we conclude

$$
\mathbb{P}\left(\tau_{q}=p_{m}\right)=\int_{0}^{q} \delta\left(p_{m}\right) \mathrm{d} p_{m}=(q-\varphi(q))_{+}=\frac{q}{m}
$$

To convey some intuition for our result, Figure 5.2 depicts an illustrative example, highlighting in green the "winning values" of $p_{m}$ such that $\hat{\tau}_{q}=p_{m}$.

Remark 5.4. Because the set of "winning values" in Lemma 5.2 is a subset of $[0, q]$ with Lebesgue measure $q / m$, we can trivially extend the result to conclude $\mathbb{P}\left\{p_{\left(R_{q}\right)}=p_{m}\right\} \leq q / m$, if $p_{m}$ is drawn from any density $f_{0}$ with $f_{0}(t) \leq 1$ for all $t \in[0, q]$. Likewise, we can extend Theorem 5.1 to show that $\max -\operatorname{lfdr}\left(\mathcal{R}_{q}\right) \leq \pi_{0} q$ with a more general null density $f_{0}$, as long as $\operatorname{lfdr}(t)$ is nondecreasing and $f_{0}(t) \leq 1$ for all $t \in[0, q]$.

### 5.2.2 Estimating $\pi_{0}$

Theorem 5.1 parallels the exact $\operatorname{FDR}$ guarantee $\operatorname{FDR}\left(\mathcal{R}_{q}^{\mathrm{BH}}\right)=\pi_{0} q$ for the BH procedure. If we bound $\pi_{0} \leq 1$, we can run our method at level $q=\alpha$ and ensure that we conservatively


Figure 5.2: Intuition for Lemma 5.2. Black points represent the empirical cdf (scaled by $\frac{m-1}{m}$ ) of $p_{1}, \ldots, p_{m-1}$; red points represent how the empirical cdf gets shifted after adding a point $p_{m}$ to its left. Adding a point $p_{m}$ can shift the supporting line by at most $\frac{q}{m}$, and each possible shift in $[0, q / m]$ corresponds to precisely one $p_{m}$ where $p_{m}$ becomes the new support point.
control max-lfdr at $\pi_{0} \alpha$, but our method will be overly conservative. In this section, we consider estimating $\pi_{0}$ using the Storey (2002) estimator of the null proportion, defined as

$$
\begin{equation*}
\hat{\pi}_{0}^{\zeta}:=\frac{1+\#\left\{i: p_{i}>\zeta\right\}}{(1-\zeta) m} \tag{5.11}
\end{equation*}
$$

modifying an estimator originally proposed by Schweder and Spjøtvoll (1982).
Our next result shows that plugging in $\hat{\pi}_{0}^{\zeta}$ and running a modification of our procedure at level $q=\alpha / \hat{\pi}_{0}^{\zeta}$ controls max-lfdr at level $\alpha$ in finite samples:

Theorem 5.5. Suppose $p_{1}, \ldots, p_{m}$ follow the Bayesian two-groups model (5.1), with $f_{0}=$ $1_{[0,1]}$ and $f_{1}$ nonincreasing. Fix $\zeta \in(0,1)$, and define a modified version of our SL procedure that only examines order statistics below $\zeta$ :

$$
\begin{equation*}
R_{q}^{\zeta}:=\underset{k \geq 0: p_{(k)} \leq \zeta}{\operatorname{argmin}} \hat{\pi}_{0}^{\zeta} p_{(k)}-\frac{q k}{m}, \tag{5.12}
\end{equation*}
$$

and $\mathcal{R}_{q}^{\zeta}=\left\{i: p_{i} \leq p_{\left(R_{q}^{\zeta}\right)}\right\}$. Then we have

$$
\operatorname{max-lfdr}\left(\mathcal{R}_{q}^{\zeta}\right)=q \cdot \frac{(1-\zeta) \pi_{0}}{1-F(\zeta)} \cdot\left(1-F(\zeta)^{m}\right) \leq q
$$

The proof of Theorem 5.5 is deferred to the Appendix. The method $\mathcal{R}_{\alpha}^{\zeta}$ coincides with $\mathcal{R}_{\alpha / \hat{\pi}_{0}^{\zeta}}$, our original procedure applied at the corrected level $\hat{q}=\alpha / \hat{\pi}_{0}^{\zeta}$, whenever $\tau_{\hat{q}} \leq \zeta$. Since we usually have $\tau_{\hat{q}} \ll 0.5 \leq \zeta$, the two methods are identical for all practical purposes.

In the next section, we will investigate the asymptotic regret of methods that estimate $\pi_{0}$. In particular, we will show that this estimation error is asymptotically negligible if it shrinks at a faster rate than $m^{-1 / 3}$. We can indeed achieve this with $\hat{\pi}_{0}^{\zeta}$ if $f_{1}$ has two continuous derivatives in a neighborhood of 1 , with $f_{1}^{\prime}(1)=f_{1}(1)=0$. By Taylor's theorem, we have

$$
1-F(\zeta)=(1-\zeta) \pi_{0}+\frac{\left(1-\pi_{0}\right) f_{1}^{\prime \prime}(\xi)}{6}(1-\zeta)^{3}
$$

for some $\xi \in[\zeta, 1]$. Assuming $\pi_{0} \in(0,1)$ and taking $\zeta=1-m^{-1 / 5}$, we then have

$$
\begin{equation*}
m^{2 / 5}\left(\hat{\pi}_{0}^{\zeta}-\pi_{0}\right) \sim m^{2 / 5}\left(\frac{1+\operatorname{Binom}(m, 1-F(\zeta))}{(1-\zeta) m}-\pi_{0}\right) \xrightarrow{d} \mathcal{N}\left(\frac{\left(1-\pi_{0}\right) f_{1}^{\prime \prime}(1)}{6}, \pi_{0}\right), \tag{5.13}
\end{equation*}
$$

with subgaussian errors for finite $m$, so the results in Section 5.3.3 generally apply. See Genovese and Wasserman (2004) and Patra and Sen (2016) for a discussion of estimators for $\pi_{0}$.

### 5.3 Asymptotic regret analysis

In this section, we study our procedure's empirical Bayes regret under the weighted classification risk $\mathbb{E}\left[L_{\lambda}(H, \mathcal{R})\right]$, where the expectation is taken over $H_{1}, \ldots, H_{m}$ and $p_{1}, \ldots, p_{m}$ according to (5.1), and $L_{\lambda}$ is defined as in (5.5). Throughout this section we will be considering a sequence of problems with $m \rightarrow \infty$.

A fundamental result of Sun and Cai (2007) is that the oracle (5.6) minimizes the weighted classification risk over all procedures, thus representing a benchmark against which we can compare methods that are feasible without a priori knowledge the lfdr. In the empirical Bayes literature (see, e.g., Efron, 2019), the price of our ignorance of the model parameters is measured by the regret, or average excess risk, given by the optimality gap

$$
\begin{equation*}
\operatorname{Regret}_{m}(\mathcal{R}):=\mathbb{E}\left[L_{\lambda}(H, \mathcal{R})-L_{\lambda}\left(H, \mathcal{R}^{*}\right)\right] \tag{5.14}
\end{equation*}
$$




Figure 5.3: Left: The fixed-threshold regret $\rho(t)(5.15)$ with Beta alternatives $f_{1}(t)=\theta t^{\theta}$ as a function of $\theta \in[0, .5]$. Right: a normalized version $\rho(t) / \rho(0)$, such that BH at level $\alpha / \pi_{0}$ has unit normalized regret, identical to the regret of the procedure that rejects nothing. The null proportion is $\pi_{0}=0.8$ and the cost-benefit ratio is $\lambda=4$.

### 5.3.1 Population regret

Before tackling the more delicate problem of calculating the regret for procedures with data-dependent $p$-value rejection thresholds, we first investigate the regret of fixed-threshold methods. For $t \in[0,1]$, let $\mathcal{R}_{t}^{\mathrm{Fix}}:=\left\{i: p_{i} \leq t\right\}$, and note that the oracle method is $\mathcal{R}^{*}=\mathcal{R}_{\tau^{*}}^{\mathrm{Fix}}$. We introduce the function $\rho(t)$ to represent the regret of this method, which is free of $m$ :

$$
\begin{equation*}
\rho(t):=\operatorname{Regret}_{m}\left(\mathcal{R}_{t}^{\mathrm{Fix}}\right)=F\left(\tau^{*}\right)-F(t)-\frac{\pi_{0}}{\alpha}\left(\tau^{*}-t\right) \tag{5.15}
\end{equation*}
$$

If $\operatorname{lfdr}\left(\tau^{*}\right)=\alpha$, then we also have $f\left(\tau^{*}\right)=\pi_{0} / \alpha$, and $\rho(t)$ is simply the error of the firstorder Taylor expansion of $F$ around $\tau^{*}$, also known as the Bregman divergence associated with $-F$. If $f$ is continuously differentiable between $t$ and $\tau^{*}$, then

$$
\begin{equation*}
\rho(t)=\frac{-f^{\prime}\left(\xi_{t}\right)}{2}\left(t-\tau^{*}\right)^{2}, \quad \text { for some } \xi_{t} \text { between } t \text { and } \tau^{*} . \tag{5.16}
\end{equation*}
$$

Since $F$ is concave, $\rho(t) \geq 0$. Finally, we can also rewrite (5.15) as an integral

$$
\begin{equation*}
\rho(t)=\int_{\tau}^{\tau^{*}}\left(1-\alpha^{-1} \operatorname{lfdr}(t)\right) \mathrm{d} F(t) \tag{5.17}
\end{equation*}
$$

This form for the regret underscores the relationship between the lfdr and the regret, and will prove useful for analyzing the regret with data-dependent thresholds.

We can evaluate $\rho$ to investigate the regret of population versions of our procedure and the BH procedure, i.e. versions of the procedures with rejection thresholds chosen using the true cdf $F$ in place of the empirical cdf $F_{m}$. The population BH threshold at an arbitrary level $q \in(0,1)$ is found by intersecting $F$ with the ray of slope $q^{-1}$, i.e.

$$
t_{q}^{\mathrm{BH}-\mathrm{POP}}:=\max \{t \in[0,1]: F(t)-t / q=0\}
$$

By comparison, the population version of our procedure $\tau_{q}$ is

$$
t_{q}:=\max \left\{t \in[0,1]: f(t) \leq q^{-1}\right\}
$$

which coincides with the oracle threshold $\tau^{*}$ when $q=\alpha / \pi_{0}$. Note that $t_{q}$ is equivalent to the population BH threshold $t_{q^{\prime}}^{\mathrm{BH}-\mathrm{POP}}$ at the lower level

$$
\begin{equation*}
q^{\prime}=\frac{t_{q}}{F\left(t_{q}\right)} \tag{5.18}
\end{equation*}
$$

Thus, there is always some value $q^{\prime}$ for which the BH procedure approximately reproduces the oracle, namely $t_{\alpha / \pi_{0}} / F\left(t_{\alpha / \pi_{0}}\right)$, but generally we cannot use it unless we know $f_{1}$ and $\pi_{0}$.

To illustrate the population regret in a concrete example, we consider a parametric alternative distribution

$$
f_{1}(t ; \theta):=\theta t^{\theta-1} \quad \text { for some } \theta \in(0,1)
$$

which is a $\operatorname{Beta}(\theta, 1)$ density. This form is called a Lehmann alternative in the multiple testing literature (see, e.g., Pounds \& Morris, 2003). In this case, the population procedures at level $q \in(0,1)$ use rejection thresholds

$$
t_{q}=\left(\frac{q^{-1}-\pi_{0}}{\left(1-\pi_{0}\right) \theta}\right)^{-\frac{1}{1-\theta}}, \quad \text { and } \quad t_{q}^{\mathrm{BH}-\mathrm{POP}}=\left(\frac{q^{-1}-\pi_{0}}{1-\pi_{0}}\right)^{-\frac{1}{1-\theta}}
$$

Furthermore, the threshold equivalence (5.18) gives

$$
q^{\prime}=\frac{\theta q}{1-(1-\theta) \pi_{0} q} \approx \theta q
$$

where the approximation holds for small values of $q$. Thus, the correspondence between $q$ and $q^{\prime}$ depends on the parameter $\theta$, which controls the signal strength under the alternative. For small values of $\theta$, the signal is very strong, and the "correct" choice of $q^{\prime}$ is much smaller than the desired max-lfdr level $\alpha$, but for weaker signals (larger $\theta$ ), we should choose $q^{\prime}$ closer to $\alpha$. Without knowing the signal strength in advance, it is difficult to know at what values of $q^{\prime}$ the BH method will perform well.

In Figure 5.3 we plot the population regret for various choices of the level of the procedure, $\pi_{0}=0.8$ and $\lambda=4$ and varying the parameter $\theta$. The population version of our procedure
at level $\frac{\alpha}{\pi_{0}}$ with $\alpha=\frac{1}{1+\lambda}=0.2$ is the oracle (5.6), so it achieves zero regret, while the conservative version of our procedure with $q=\alpha$ performs quite well for all values of the alternative parameter $\theta$. In this example, the asymptotic error incurred from conservatively bounding $\pi_{0}$ by one in the procedure is small compared to the error incurred by using $\mathrm{BH}\left(q^{\prime}\right)$ at an $a d$ hoc value. The BH procedure at levels $\frac{\alpha}{\pi_{0}}$ or $\alpha$ incurs substantial asymptotic regret by comparison. In particular, note that the $\mathrm{BH}\left(\alpha / \pi_{0}\right)$ procedure incurs the same asymptotic regret as the procedure that rejects nothing; i.e. $\rho\left(t_{\alpha / \pi_{0}}^{\mathrm{BH}-\mathrm{POP}}\right)=\rho(0)$. If we run BH at a lower level like $\alpha / 2, \alpha / 10$, or $\alpha / 100$, we can do well for some range of $\theta$ values, but struggle at other parts of the parameter space. No single level for BH dominates in terms of regret, so for the classification risk it is more appropriate to view the BH level as a tuning parameter (Neuvial \& Roquain, 2012).

### 5.3.2 Relationship of our method to the Grenander estimator

Since the marginal density $f$ appears in the denominator of the lfdr, bounding $\pi_{0} \leq 1$ and plugging in Grenander's estimator $\hat{f}_{m}$ (defined in (5.8)) gives the conservative estimate

$$
\widehat{\operatorname{lfdr}}(t):=\frac{1}{\hat{f}_{m}(t)}, \quad t \in[0,1]
$$

Similar to how the BH procedure chooses an interval $[0, t]$ as large as possible subject to a constraint on an estimate of the FDP, the rejection threshold of the SL procedure can be equivalently expressed as

$$
\begin{equation*}
\tau_{q}=\underset{p_{(0)}, \ldots, p_{(m)}}{\operatorname{argmax}}\left\{\frac{q k}{m}-p_{(k)}\right\}=\sup \{t \in[0,1]: \widehat{\operatorname{lfdr}}(t) \leq q\} \tag{5.19}
\end{equation*}
$$

taking the convention that $\sup \emptyset \equiv 0$. The equivalence in (5.19) is illustrated in Figure 5.4. Let $\hat{F}_{m}$ denote the least concave majorant of the empirical cdf $F_{m}$, plotted as a dotted blue line in the left panel of Figure 5.4. By definition of $\widehat{\operatorname{lfdr}}(t)$, the supremum on the right hand side is equal to the largest $t$ for which $\frac{\mathrm{d}}{\mathrm{d} t}\left(q \hat{F}_{m}(t)-t\right)=q \hat{f}_{m}(t)-1 \geq 0$, which corresponds to the maximizer of the function $q \hat{F}_{m}(t)-t$, illustrated for example in the right panel of Figure 5.4. $\hat{F}_{m} \geq F_{m}$ implies

$$
q \hat{F}_{m}(t)-t \geq q F_{m}(t)-t, \quad t \in[0,1]
$$

with equality at the knots of $\hat{F}_{m}$, and since the maximizer of the left hand side occurs at a knot of $\hat{F}_{m}$, it is also the maximizer of the right hand side, i.e. the $\operatorname{argmax}$ of $\frac{q k}{m}-p_{(k)}$.

We can again compare this result with the $\mathrm{BH}(q)$ threshold, given by

$$
\tau_{q}^{\mathrm{BH}}=\max _{k=0, \ldots, m}\left\{p_{(k)}: \frac{q k}{m}-p_{(k)} \geq 0\right\}=\sup \left\{t \in[0,1]: F_{m}(t) \geq q^{-1} t\right\}
$$



Figure 5.4: Left: empirical cdf $F_{m}$ and its least concave majorant $\hat{F}_{m}$. The support line of slope $q^{-1}$ touches both curves at the decision threshold $\tau_{q}$. Right: the same plot with the line $t / q$ subtracted off.
which is the largest $t$ for which the ray $q^{-1} t$ lies below the ecdf $F_{m}(t)$. Our procedure instead finds the last intersection of the graph of $F_{m}$ with a support line of slope $q^{-1}$, since

$$
\widehat{\operatorname{lfdr}}(t) \leq q \Longleftrightarrow \hat{f}_{m}(t) \geq q^{-1}
$$

This relationship is illustrated in the left panel of Figure 5.4.

### 5.3.3 Asymptotic behavior of our procedure

Equation (5.16) suggests that, when $f$ is sufficiently regular near $\tau^{*}$, the regret is closely related to the squared error of the rejection threshold. Our main result in this section establishes cube-root asymptotics for the behavior of our procedure $\mathcal{R}_{q}$ with $q=\alpha / \hat{\pi}_{0}$, where $\hat{\pi}_{0}$ consistently estimates $\pi_{0}$; if $\pi_{0}$ is known, then the results apply directly with $\hat{\pi}_{0}=\pi_{0}$.

We derive limiting distributions for the threshold $\tau_{q}$, the lfdr at the threshold, and the regret of $\mathcal{R}_{q}$. All three are given in terms of Chernoff's distribution (Chernoff, 1964), which is defined as the distribution of the maximizer $Z$ of a standard two-sided Brownian motion $W=$ $(W(t))_{t \in \mathbb{R}}$ with parabolic drift:

$$
\begin{equation*}
Z=\underset{t \in \mathbb{R}}{\operatorname{argmax}} W(t)-t^{2} . \tag{5.20}
\end{equation*}
$$

The random variable $Z$ has a density with respect to the Lebesgue measure on $\mathbb{R}$ that is symmetric about zero. Dykstra and Carolan (1999) suggest approximating the density and cdf of $Z$ by those of $\mathcal{N}\left(0,(.52)^{2}\right)$. This approximation can be somewhat crude but gives a rough sense for the distribution of $Z$. Groeneboom and Wellner (2001) provide much more accurate numerical methods to compute the density, cdf, quantiles and moments of $Z$.

Theorem 5.6. Suppose $p_{1}, \ldots, p_{m}$ follow the Bayesian two-groups model (5.1), with $\pi_{0} \in$ $(0,1), f_{0}=1_{[0,1]}$, and $f_{1}$ nonincreasing. For $q \in\left(0, \pi_{0}^{-1}\right)$, assume additionally that
(i) there is a unique value $t_{q} \in(0,1)$ for which $f\left(t_{q}\right)=q^{-1}$,
(ii) $f$ is continuously differentiable in a neighborhood of $t_{q}$ with $f^{\prime}\left(t_{q}\right)<0$, and
(iii) $\hat{q}$ is any random variable with $m^{1 / 3}(\hat{q}-q) \xrightarrow{p} 0$ as $m \rightarrow \infty$.

Then we have, as $m \rightarrow \infty$,

$$
\begin{align*}
& m^{1 / 3}\left(\tau_{\hat{q}}-t_{q}\right) \xrightarrow{d}\left(\frac{q}{4} \cdot f^{\prime}\left(t_{q}\right)^{2}\right)^{-1 / 3} Z, \quad \text { and }  \tag{5.21}\\
& m^{1 / 3} \cdot \frac{\operatorname{lfdr}\left(\tau_{\hat{q}}\right)-\pi_{0} q}{\pi_{0} q} \xrightarrow{d}\left(4 q^{2} \cdot\left|f^{\prime}\left(t_{q}\right)\right|\right)^{1 / 3} Z . \tag{5.22}
\end{align*}
$$

where $Z$ follows Chernoff's distribution defined in (5.20). Further, suppose that

$$
\begin{equation*}
\mathbb{P}\left\{m^{-1 / 3}(\hat{q}-q)>\varepsilon\right\}=o\left(m^{-2 / 3}\right), \quad \text { for all } \varepsilon>0 . \tag{5.23}
\end{equation*}
$$

Then we also have $m^{1 / 3} \mathbb{E}\left[\tau_{\hat{q}}\right] \rightarrow t_{q}$. In addition,

$$
\begin{align*}
m^{2 / 3} \operatorname{Var}\left(\tau_{\hat{q}}\right) & \rightarrow\left(\frac{q}{4} \cdot f^{\prime}\left(t_{q}\right)^{2}\right)^{-2 / 3} \operatorname{Var}(Z), \quad \text { and }  \tag{5.24}\\
m^{2 / 3} \operatorname{Var}\left(\frac{\operatorname{lfdr}\left(\tau_{\hat{q}}\right)-\pi_{0} q}{\pi_{0} q}\right) & \rightarrow\left(4 q^{2} \cdot\left|f^{\prime}\left(t_{q}\right)\right|\right)^{2 / 3} \operatorname{Var}(Z), \tag{5.25}
\end{align*}
$$

where $\operatorname{Var}(Z) \approx 0.26$.
The proof of Theorem 5.6 is deferred to Appendix 5.6.2. It is well-known that the Grenander estimator $\hat{f}_{m}$ estimates $f$ at a cube root rate pointwise, away from zero, but this result, due to (Rao, 1969), is too weak to describe the behavior of our procedure. We rely on a stronger version of this result due to Dümbgen et al. (2016) that approximates the local behavior of the Grenander estimator near $t_{q}$.

The distributional result (5.22) complements our result from Theorem 5.1, by showing that $\operatorname{lfdr}\left(\tau_{q}\right)=\max _{i \in \mathcal{R}_{q}} \operatorname{lfdr}\left(p_{i}\right)$ is not only controlled in expectation, but also concentrates at rate $m^{-1 / 3}$ around its expectation. In particular, because $\mathbb{P}\{Z \geq 1\} \approx 0.05$, we have

$$
\frac{\operatorname{lfdr}\left(\tau_{q}\right)-\pi_{0} q}{\pi_{0} q} \leq m^{-1 / 3}\left(4 q^{2} \cdot\left|f^{\prime}\left(t_{q}\right)\right|\right)^{1 / 3}
$$

with roughly $95 \%$ probability in large samples. For example, suppose we use $q=0.2$, so $f\left(t_{q}\right)=5$, and suppose that $f^{\prime}\left(t_{q}\right)=-50$. Then, whereas Theorem 5.1 guarantees $\mathbb{E}\left[\operatorname{fdr}\left(\tau_{q}\right)\right] \leq 0.2$ exactly, the asymptotic estimate from Theorem 5.6 bounds the 95 th percentile of $\operatorname{lfdr}\left(\tau_{q}\right)$ at 0.24 if $m=1000$, or at 0.21 if $m=64,000$.

To understand why the error is of order $m^{-1 / 3}$, consider fixed $q$ and recall that the threshold $\tau_{q}$ maximizes the stochastic process

$$
U(t):=F_{m}(t)-F_{m}\left(t_{q}\right)-\frac{t-t_{q}}{q} .
$$

Because $f\left(t_{q}\right)=q^{-1}$, we have for $t$ near $t_{q}$,

$$
F(t)-F\left(t_{q}\right) \approx \frac{t-t_{q}}{q}+\frac{f^{\prime}\left(t_{q}\right)}{2}\left(t-t_{q}\right)^{2} .
$$

Introducing the local parameterization $t=t_{q}+m^{-a} h$ for $a>0$ leads to

$$
U\left(t_{q}+m^{-a} h\right) \approx-\frac{\left|f^{\prime}\left(t_{q}\right)\right|}{2} \cdot \frac{h^{2}}{m^{2 a}}+\mathcal{N}\left(0, \frac{h}{q m^{a+1}}\right) .
$$

Setting $a=1 / 3$ balances the mean and variance, giving

$$
m^{2 / 3} U\left(t_{q}+m^{-1 / 3} h\right) \xrightarrow{d}-\frac{\left|f^{\prime}\left(t_{q}\right)\right|}{2} h^{2}+\mathcal{N}\left(0, \frac{h}{q}\right) .
$$

Under this local scaling, $U(t)$ converges to a Brownian motion with parabolic drift, and its maximizer $\tau_{q}$ converges to Chernoff's distribution. Theorem 5.6 applies a more careful version of this argument, replacing $F_{m}(t)$ with its LCM $\hat{F}_{m}(t)$ and using a result of Dümbgen et al. (2016) to characterize the process $\hat{f}_{m}(t)$ under the same local scaling. The corresponding results for $\operatorname{lfdr}\left(\tau_{q}\right)$ follow from first-order Taylor expansion of $\operatorname{lfdr}(t)=\pi_{0} / f(t)$ around $t_{q}$.

By specializing Theorem 5.6 to $q=\alpha / \pi_{0}$ and $\hat{q}=\alpha / \hat{\pi}_{0}$, we obtain the limiting regret for our procedure with a known or accurately estimated null proportion.

Theorem 5.7. Suppose $p_{1}, \ldots, p_{m}$ follow the Bayesian two-groups model (5.1), with $\pi_{0} \in$ $(0,1), f_{0}=1_{[0,1]}$, and $f_{1}$ nonincreasing. Assume additionally that
(i) there is a unique value $\tau^{*} \in(0,1)$ for which $\operatorname{lfdr}\left(\tau^{*}\right)=\frac{\pi_{0}}{f\left(\tau^{*}\right)}=\alpha$,
(ii) $f$ is continuously differentiable in a neighborhood of $\tau^{*}$ with $f^{\prime}\left(\tau^{*}\right)<0$, and
(iii) $\hat{\pi}_{0}$ is any estimator of $\pi_{0}$ with $\mathbb{P}\left\{m^{1 / 3}\left(\hat{\pi}_{0}-\pi_{0}\right)>\varepsilon\right\}=o\left(m^{-2 / 3}\right)$ for all $\varepsilon>0$.

Then we have, as $m \rightarrow \infty$,

$$
\begin{equation*}
m^{2 / 3} \operatorname{Regret}_{m}\left(\mathcal{R}_{\alpha / \hat{\pi}_{0}}\right) \rightarrow\left(\frac{\alpha^{2}}{2 \pi_{0}^{2}} \cdot\left|f^{\prime}\left(\tau^{*}\right)\right|\right)^{-1 / 3} \operatorname{Var}(Z) \tag{5.26}
\end{equation*}
$$

where $Z$ follows Chernoff's distribution defined in (5.20), and $\operatorname{Var}(Z) \approx 0.26$.

Theorems 5.6-5.7 deal with the regret for $\pi_{0} \in(0,1)$. Under the global null, represented in the Bayesian model by $\pi_{0}=1$, the behavior is different and the regret is simply $\lambda \mathbb{E} V$, which is $O\left(m^{-1}\right)$, as we see next.

Proposition 5.8. Suppose $\left(p_{i}\right)_{i=1}^{m}$ follow a two-groups model (5.1) with $f_{0}=1_{[0,1]}$ and $\pi_{0}=$ 1, i.e. $H_{i}=0$ for all $i$ and $p_{i} \stackrel{\text { iid }}{\sim} \operatorname{Unif}(0,1)$. Then as $m \rightarrow \infty$, we have

$$
m \operatorname{Regret}_{m}\left(\mathcal{R}_{q}\right) \rightarrow \lambda \sum_{k=1}^{\infty} \mathbb{P}\left\{U_{k} \leq q\right\}, \quad \text { for } U_{k} \sim \operatorname{Gamma}(k, k)
$$

which is finite for every $q \in[0,1)$.
Proposition 5.8 is closely related to results derived in Finner and Roters (2001).

### 5.4 Numerical results

This section highlights our main results on simulation experiments. We adapt a simulation setting of Benjamini and Hochberg (1995) to the two-groups model (5.1). Specifically, define the alternative density

$$
\begin{equation*}
f_{1}(t)=\frac{\frac{1}{4} \sum_{i=1}^{4} \phi\left(\bar{\Phi}^{-1}(t)-5 \frac{i}{4}\right)}{\phi\left(\bar{\Phi}^{-1}(t)\right)} \quad \text { for } 0 \leq t \leq 1 \tag{5.27}
\end{equation*}
$$

where $\phi$ and $\bar{\Phi}$ denote the density and survival function of the standard Gaussian distribution. Concretely, a non-null $p$-value $p_{i} \sim f_{1}$ can be constructed by first taking $Y_{i} \sim \mathcal{N}\left(\mu_{i}, 1\right)$ where $\mu_{i}$ is drawn uniformly at random from the set $\left\{5 \frac{i}{4}: i=1,2,3,4\right\}$; then, $p_{i}=\bar{\Phi}\left(Y_{i}\right)$ is a one sided $p$-value for the null-hypothesis that $\mu_{i}=0$. We use a null proportion of $\pi_{0}=0.75$. Figure 5.5 shows the mixture density and corresponding lfdr.

We repeatedly sampled from the above two-groups model with $m=64$ hypotheses. Figure 5.5 shows the FDR (left panel) and max-lfdr (right panel) for both our procedure and the BH procedure, at conservative level $q=\alpha$ and estimated level $\hat{q}=\alpha / \hat{\pi}_{0}^{\zeta}$. The BH procedure, shown in red, achieves FDR exactly $\pi_{0} q$, whereas the max-lfdr can be much larger. By contrast, our procedure, shown in blue, conservatively controls FDR substantially below the level $\pi_{0} q$ but has max-lfdr equal to $\pi_{0} q$.

Figure 5.6 shows a log-log plot of the regret as a function of the sample size $m$. The red curve shows the regret of our uncorrected procedure $\mathcal{R}_{\alpha}$ for $\alpha=0.05$, which asymptotically tends to $\rho\left(t_{\alpha}\right)$ and hence asymptotically incurs some non-vanishing regret described in Section 5.3.1. The blue curve shows the regret of the corrected procedure $\mathcal{R}_{\alpha / \pi_{0}}$ with known $\pi_{0}$. For larger samples, the simulated regret closely matches the asymptotic prediction from (5.26), shown in black. The green curve (which is nearly indistinguishable from the blue curve) shows the corrected procedure with an estimated null proportion $\hat{\pi}_{0}^{\zeta}$ based on (5.11) with $\zeta=1-m^{-1 / 5}$.


Figure 5.5: Above: Mixture density $f$ (left) and lfdr (right), with alternative density $f_{1}$ defined in (5.27) and null proportion $\pi_{0}=0.75$. Note $f_{1}$ diverges as $t \downarrow 0$. Below: Comparison of FDR control (left) and max-lfdr control (right) on simulated data. The estimate of the null proportion is (5.11) with $\zeta=0.5$.




Figure 5.6: A log-log plot of the regret (5.14) as a function of the sample size. The black line shows the asymptotic prediction (5.26) of Theorem 5.7. For this simulation, the alternative density $f_{1}$ is defined in (5.27), cost-benefit ratio $\lambda=19$ and null proportion $\pi_{0}=0.75$.

### 5.5 Discussion

In this chapter we introduced a new error criterion, the max-lfdr, which modifies the FDR by redirecting attention away from the average quality of the rejection set and toward the rejections that are close to the rejection boundary. Despite the seeming difficulty of measuring the quality of a single rejection, we also introduce a simple new multiple testing procedure that controls the max-lfdr at level $\pi_{0} q$ in finite samples, where $q$ is a tuning parameter and $\pi_{0}$ is the null proportion. We assume only that the data follow a Bayesian two-groups model in which smaller $p$-values reflect stronger evidence against the null. We find that our method is better able than the BH method to adapt to the unknown problem structure, and to perform well without knowledge of the true underlying distribution.

The BH procedure owes its enduring utility for FDR control in part to its versatility beyond this basic setting, however. It is known to still control FDR, for instance, when the null $p$-values are super-uniform and under certain forms of positive dependence, two of many possible extensions that we leave open for our procedure.

Another seeming advantage of the FDR criterion is that it requires no Bayesian assumptions, whereas the max-lfdr is only defined with reference to a Bayesian model. A possible avenue for generalizing the max-lfdr to frequentist settings is to work with its characterization as the probability that the last rejection is a false discovery. Indeed, our proof of

Theorem 5.1 implies that max-lfdr is controlled even conditional on $H_{1}, \ldots, H_{m}$. This is initially puzzling: if each $H_{i}$ is fixed, then how can we speak of the probability that the last rejection is a false discovery? The answer is that $H_{(R)}$ is random even if $H_{1}, \ldots, H_{m}$ are fixed, since its index is random. We leave further development of the frequentist connection to the max-lfdr to future work.

### 5.6 Proofs

### 5.6.1 Proofs of results from Section 5.2

Proof of Theorem 5.5. As in the proof of Theorem 5.1, we have

$$
\max -\operatorname{lfdr}\left(\mathcal{R}_{q}^{\zeta}\right)=\mathbb{P}\left\{H_{\left(R_{q}^{\zeta}\right)}=0, R_{q}^{\zeta}>0\right\}=m \mathbb{P}\left\{H_{m}=0, p_{\left(R_{q}^{\zeta}\right)}=p_{m}\right\} .
$$

Define the $\sigma$-field $\mathcal{F}=\sigma\left(p_{1}, \ldots, p_{m-1}, H_{m}, 1\left\{p_{m} \leq \zeta\right\}\right)$. We restrict our attention to the event $A=\left\{H_{m}=0, p_{m} \leq \zeta\right\}$, since the event $\left\{H_{m}=0, p_{\left(R_{q}^{\zeta}\right)}=p_{m}\right\}$ cannot occur except on $A$. On $A$, which is $\mathcal{F}$-measurable, we have $p_{m} / \zeta \mid \mathcal{F} \sim U[0,1]$.

Let $m^{\zeta}=\#\left\{i: p_{i} \leq \zeta\right\}$, which is also $\mathcal{F}$-measurable. If $j_{1} \leq \cdots \leq j_{m}=m$ are the indices of the $p$-values that are below $\zeta$, define the modified $p$-values $p_{i}^{\zeta}=p_{j_{i}} / \zeta$, for $i=1, \ldots, m^{\zeta}$. Because the order statistics of $\zeta p_{1}^{\zeta}, \ldots, \zeta p_{m \zeta}^{\zeta}$ are also the first $m^{\zeta}$ order statistics of $p_{1}, \ldots, p_{m}$, the quantity $R_{q}^{\zeta}$ defined in (5.12) can be rewritten as

$$
\begin{aligned}
R_{q}^{\zeta} & =\underset{k=0, \ldots, m \zeta}{\operatorname{argmin}} \zeta p_{(k)}^{\zeta}-\frac{q}{\hat{\pi}_{0}^{\zeta}} \cdot \frac{k}{m} \\
& =\underset{k=0, \ldots, m^{\zeta}}{\operatorname{argmin}} p_{(k)}^{\zeta}-\frac{q^{\zeta} k}{m^{\zeta}}, \quad \text { for } q^{\zeta}=\frac{q m^{\zeta}}{\zeta \hat{\pi}_{0}^{\zeta} m}
\end{aligned}
$$

Applying Lemma 5.2, we have

$$
\mathbb{P}\left\{H_{m}=0, p_{\left(R_{q}^{\zeta}\right)}=p_{m} \mid \mathcal{F}\right\}=\frac{q^{\zeta}}{m^{\zeta}} \cdot 1_{A}=\frac{q}{\zeta \hat{\pi}_{0}^{\zeta} m} \cdot 1_{A}
$$

Marginalizing over $\mathcal{F}$, and noting that $\mathbb{P}(A)=\pi_{0} \zeta$, we obtain

$$
\begin{aligned}
\mathbb{P}\left\{H_{m}=0, p_{\left(R_{q}^{\zeta}\right)}=p_{m}\right\} & =\frac{q}{m} \cdot \mathbb{E}\left[\left.\frac{\pi_{0}}{\hat{\pi}_{0}^{\zeta}} \right\rvert\, A\right] \\
& =\frac{q}{m} \cdot \frac{(1-\zeta) \pi_{0}}{1-F(\zeta)} \cdot \mathbb{E}\left[\frac{(1-F(\zeta)) m}{1+\#\left\{i<m: p_{i}>\zeta\right\}}\right] \\
& =\frac{q}{m} \cdot \frac{(1-\zeta) \pi_{0}}{1-F(\zeta)} \cdot\left(1-F(\zeta)^{m}\right) \\
& \leq \frac{q}{m}
\end{aligned}
$$

completing the proof. The final inequality is a standard binomial identity:

$$
\begin{aligned}
\mathbb{E}\left[\frac{\beta m}{1+\operatorname{Binom}(m-1, \beta)}\right] & =\sum_{k=0}^{m-1} \frac{\beta m}{1+k}\binom{m-1}{k} \beta^{k}(1-\beta)^{m-1-k} \\
& =\sum_{k=0}^{m-1}\binom{m}{k+1} \beta^{k+1}(1-\beta)^{m-(k+1)} \\
& =\sum_{j=1}^{m}\binom{m}{j} \beta^{j}(1-\beta)^{m-j} \\
& =\mathbb{P}\{\operatorname{Binom}(m, \beta) \geq 1\} \\
& =1-(1-\beta)^{m} .
\end{aligned}
$$

### 5.6.2 Proofs of results from Section 5.3

Proof of Theorem 5.6. Our proof will use the switching relation that states, for any $t \in(0,1)$, we have almost surely

$$
\tau_{\hat{q}} \leq t \Longleftrightarrow \hat{f}_{m}(t) \leq \hat{q}^{-1}
$$

We will work with a local expansion of $\hat{f}_{m}(t)$ around $t_{q}$ using the local parameterization $t=t_{q}+m^{-1 / 3} h$. Using $f\left(t_{q}\right)=q^{-1}$, the switching relation becomes

$$
m^{-1 / 3}\left(\tau_{\hat{q}}-t_{q}\right) \leq h \Longleftrightarrow \hat{f}_{m}\left(t_{q}+m^{-1 / 3} h\right)-f\left(t_{q}\right) \leq \hat{q}^{-1}-q^{-1} .
$$

Now let $W$ denote a standard two-sided Brownian motion, and let $\mathbb{S}_{a, b}$ denote the process of left derivatives of the least concave majorant of $X_{a, b}(t)=a W(t)-b t^{2}$, where $a=\sqrt{f\left(t_{q}\right)}$ and $b=\left|f^{\prime}\left(t_{q}\right)\right| / 2$. Under our regularity assumptions, Dümbgen et al. (2016) show

$$
m^{1 / 3}\left(\hat{f}_{m}\left(t_{q}+m^{-1 / 3} h\right)-f\left(t_{q}\right)\right) \Rightarrow \mathbb{S}_{a, b}(h)
$$

in the Skorokhod topology on $D[-K, K]$ for every finite $K>0$. Since $m^{1 / 3}\left(\hat{q}^{-1}-q^{-1}\right) \xrightarrow{p} 0$ by assumption, we have

$$
\mathbb{P}\left\{m^{1 / 3}\left(\tau_{\hat{q}}-t_{q}\right) \leq h\right\} \rightarrow \mathbb{P}\left\{\mathbb{S}_{a, b}(h) \leq 0\right\}
$$

Observe that $\mathbb{S}_{a, b}(h) \leq 0$ iff $t_{a, b}^{*} \leq h$, where $t_{a, b}^{*}$ is the (a.s. unique) maximizer of $X_{a, b}$ (note the maximizer $t_{a, b}^{*}$ is always a knot in the concave majorant since the horizontal line with intercept $X_{a, b}\left(t_{a, b}^{*}\right)$ is a supporting line intersecting $\left(t_{a, b}^{*}, X_{a, b}\left(t_{a, b}^{*}\right)\right)$ ). Combining this observation with the previous display, we have

$$
m^{1 / 3}\left(\tau_{\hat{q}}-t_{q}\right) \xrightarrow{d} t_{a, b}^{*} \stackrel{d}{=}(b / a)^{-2 / 3} Z=\left(\frac{q}{4} \cdot f^{\prime}\left(t_{q}\right)^{2}\right)^{-1 / 3} Z,
$$

proving (5.21). Next we turn to the lfdr asymptotics. By Taylor's theorem,

$$
m^{1 / 3}\left(\operatorname{lfdr}\left(\tau_{\hat{q}}\right)-\pi_{0} q\right)=\operatorname{lfdr}^{\prime}(\omega) \cdot m^{1 / 3}\left(\tau_{\hat{q}}-t_{q}\right)
$$

for some $\omega$ between $\tau_{\hat{q}}$ and $t_{q}$. Using

$$
\operatorname{lfdr}^{\prime}\left(t_{q}\right)=\frac{-\pi_{0} f^{\prime}\left(t_{q}\right)}{f\left(t_{q}\right)^{2}}=\pi_{0} q^{2} \cdot\left|f^{\prime}\left(t_{q}\right)\right|
$$

and applying the continuous mapping theorem and Slutsky's theorem, we obtain

$$
\operatorname{lfdr}^{\prime}(\omega) \cdot m^{1 / 3}\left(\tau_{\hat{q}}-t_{q}\right) \xrightarrow{d} \operatorname{lfdr}^{\prime}\left(t_{q}\right) \cdot\left(\frac{q}{4} \cdot f^{\prime}\left(t_{q}\right)^{2}\right)^{-1 / 3} Z=\pi_{0} q \cdot\left(4 q^{2} \cdot\left|f^{\prime}\left(t_{q}\right)\right|\right)^{1 / 3} Z,
$$

proving (5.22). Next, under the strengthened assumption (5.23), fix $\varepsilon>0$ and define the event

$$
\begin{equation*}
A_{\varepsilon}=\left\{|\hat{q}-q| \leq m^{-1 / 3} \varepsilon,\left|\tau_{\hat{q}}-t_{q}\right| \leq m^{-2 / 9}\right\} \tag{5.28}
\end{equation*}
$$

and the truncated random variable

$$
Y_{m}=m^{1 / 3}\left(\tau_{\hat{q}}-t_{q}\right) \cdot 1_{A_{\varepsilon}},
$$

We will show that $\mathbb{P}\left(A_{\varepsilon}^{c}\right)=o\left(m^{-2 / 3}\right)$. As a result, $Y_{m}$ has the same limit in distribution as $m^{1 / 3}\left(\tau_{\hat{q}}-t_{q}\right)$. If we can show that the sequence $Y_{m}^{2}$ is uniformly integrable, we will have convergence of its mean and variance to the mean and variance of its limiting distribution. Then, because

$$
\mathbb{E}\left[\left(m^{1 / 3}\left(\tau_{\hat{q}}-t_{q}\right)-Y_{m}\right)^{2}\right] \leq m^{2 / 3} \mathbb{P}\left(A_{\varepsilon}^{c}\right) \rightarrow 0
$$

we will have the same limiting mean and variance for $m^{1 / 3}\left(\tau_{\hat{q}}-t_{q}\right)$.
To show that $\mathbb{P}\left(A_{\varepsilon}^{c}\right)=o\left(m^{-2 / 3}\right)$, let $q_{1}=q-m^{-1 / 3} \varepsilon$ and $q_{2}=q+m^{-1 / 3} \varepsilon$ and assume that $m$ is sufficiently large that $m^{-1 / 3} \varepsilon \leq m^{-2 / 9} / 2$, and

$$
f^{\prime}(t) \leq f^{\prime}\left(t_{q}\right) / 2, \quad \text { for all } t \in\left[t_{q}-m^{-2 / 9}, t_{q}+m^{-2 / 9}\right] .
$$

As a result, for all $t \geq t_{q_{2}}+m^{-2 / 9} / 2$, we have

$$
\begin{aligned}
F(t)-F\left(t_{q_{2}}\right)-\frac{t-t_{q_{2}}}{q_{2}} & \leq F\left(t_{q_{2}}+m^{-2 / 9} / 2\right)-F\left(t_{q_{2}}\right)-\frac{m^{-2 / 9}}{2 q_{2}} \\
& \leq \frac{f^{\prime}\left(t_{q}\right)}{16} \cdot m^{-4 / 9}
\end{aligned}
$$

Then, since $\tau_{\hat{q}} \leq \tau_{q_{2}}$ a.s. on $A_{\varepsilon}$, we have

$$
\begin{aligned}
\mathbb{P}\left\{\tau_{\hat{q}}>t_{q}+m^{-2 / 9}, A_{\varepsilon}\right\} & \leq \mathbb{P}\left\{\tau_{q_{2}}>t_{q_{2}}+m^{-2 / 9} / 2\right\} \\
& \leq \mathbb{P}\left\{\sup _{t \geq t_{q_{2}}+m^{-2 / 9} / 2} F_{m}(t)-F_{m}\left(t_{q_{2}}\right)-\frac{t-t_{q_{2}}}{q_{2}} \geq 0\right\} \\
& \leq \mathbb{P}\left\{\sup _{t \geq t_{q_{2}}+m^{-2 / 9} / 2} F_{m}(t)-F(t)-\left(F_{m}\left(t_{q_{2}}\right)-F\left(t_{q_{2}}\right)\right) \geq \frac{\left|f^{\prime}\left(t_{q}\right)\right|}{16} \cdot m^{-4 / 9}\right\} \\
& \leq \mathbb{P}\left\{\sup _{t \in[0,1]}\left|F_{m}(t)-F(t)\right| \geq \frac{\left|f^{\prime}\left(t_{q}\right)\right|}{32} \cdot m^{-4 / 9}\right\} \\
& \leq C_{\mathrm{DKW}} \exp \left\{-\frac{f^{\prime}\left(t_{q}\right)^{2}}{512} \cdot m^{1 / 9}\right\},
\end{aligned}
$$

where $C_{\mathrm{DKW}}$ is the constant for the Dvoretzky-Kiefer-Wolfowitz inequality. An analogous argument yields the same bound for $\mathbb{P}\left\{\tau_{\hat{q}} \leq t_{q}-m^{-2 / 9}\right\}$.

Proof of Theorem 5.7. Define $q=\alpha / \pi_{0}$ and $\hat{q}=\alpha / \hat{\pi}_{0}$, and let $\Delta \subseteq\{1, \ldots, m\}$ denote the symmetric difference between the two rejection sets:

$$
\Delta= \begin{cases}\left\{R_{\hat{q}}+1, \ldots, R^{*}\right\} & \text { if } R_{\hat{q}}<R^{*} \\ \left\{R^{*}+1, \ldots, R_{\hat{q}}\right\} & \text { if } R_{\hat{q}}>R^{*} \\ \emptyset & \text { if } R_{\hat{q}}=R^{*}\end{cases}
$$

Then we have

$$
L_{\lambda}\left(H, \mathcal{R}_{\hat{q}}\right)-L_{\lambda}\left(H, \mathcal{R}^{*}\right)=\frac{1}{m}\left(R^{*}-R_{\hat{q}}+\frac{\operatorname{sgn}\left(R_{\hat{q}}-R^{*}\right)}{\alpha} \sum_{i \in \Delta}\left(1-H_{i}\right)\right) .
$$

Conditional on $F_{m}$, we have $H_{i} \stackrel{\text { ind }}{\sim} \operatorname{Bern}\left(1-\operatorname{lfdr}\left(p_{(i)}\right)\right)$, giving conditional expectation

$$
\begin{aligned}
\Gamma_{m} & :=\mathbb{E}\left[L_{\lambda}\left(H, \mathcal{R}_{\hat{q}}\right)-L_{\lambda}\left(H, \mathcal{R}^{*}\right) \mid F_{m}\right] \\
& =\frac{1}{m}\left(R^{*}-R_{\hat{q}}+\frac{\operatorname{sgn}\left(R_{\hat{q}}-R^{*}\right)}{\alpha} \sum_{i \in \Delta} \operatorname{lfdr}\left(p_{(i)}\right)\right) \\
& =\int_{\tau_{\hat{q}}}^{\tau^{*}}\left(1-\alpha^{-1} \operatorname{lfdr}(t)\right) d F_{m}(t) \\
& =\rho\left(\tau_{\hat{q}}\right)+\alpha^{-1} \int_{\tau_{\hat{q}}}^{\tau^{*}}(\alpha-\operatorname{lfdr}(u))\left(d F_{m}(u)-d F(u)\right)
\end{aligned}
$$

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Define the same truncation event $A_{\varepsilon}$ as in (5.28).

$$
A_{\varepsilon}=\left\{|\hat{q}-q| \leq m^{-1 / 3} \varepsilon,\left|\tau_{\hat{q}}-\tau^{*}\right| \leq m^{-2 / 9}\right\}
$$

Then, because $\left|\Gamma_{m}\right| \leq \alpha^{-1}$ we have

$$
\begin{align*}
\mid \operatorname{Regret}_{m}\left(\mathcal{R}_{\hat{q}}\right)- & \mathbb{E}\left[\rho\left(\tau_{\hat{q}}\right) 1_{A_{\varepsilon}}\right] \mid \\
& \leq \alpha^{-1} \mathbb{E}\left[\left|\int_{\tau_{\hat{q}}}^{\tau^{*}}(\alpha-\operatorname{lfdr}(u))\left(d F_{m}(u)-d F(u)\right)\right| 1_{A_{\varepsilon}}\right]+\alpha^{-1} \mathbb{P}\left(A_{\varepsilon}^{c}\right) \tag{5.29}
\end{align*}
$$

We showed in the proof of Theorem 5.6 that $\mathbb{P}\left(A_{\varepsilon}^{c}\right)=o\left(m^{-2 / 3}\right)$. Furthermore,

$$
\begin{aligned}
m^{2 / 3} \mathbb{E}\left[\rho\left(\tau_{\hat{q}}\right) 1_{A_{\varepsilon}}\right] & =\mathbb{E}\left[\frac{f^{\prime}\left(\xi_{\tau_{\hat{q}}}\right)}{2} \cdot m^{2 / 3}\left(\tau_{\hat{q}}-\tau^{*}\right)^{2} \cdot 1_{A_{\varepsilon}}\right] \\
& \rightarrow \frac{f^{\prime}\left(\tau^{*}\right)}{2}\left(\frac{\alpha}{4 \pi_{0}} \cdot f^{\prime}\left(\tau^{*}\right)^{2}\right)^{-2 / 3} \operatorname{Var}(Z) \\
& =\left(\frac{\alpha^{2}}{2 \pi_{0}^{2}} \cdot\left|f^{\prime}\left(\tau^{*}\right)\right|\right)^{-1 / 3} \operatorname{Var}(Z),
\end{aligned}
$$

where we have used the fact that $f^{\prime}\left(\xi_{\tau_{\tilde{q}}}\right)$ is uniformly close to $f^{\prime}\left(\tau^{*}\right)$ on $A_{\varepsilon}$.
Proof of Proposition 5.8. Since $H_{i}=0$ for all $i$

$$
L_{\lambda}\left(H, \mathcal{R}_{\alpha}\right)-L_{\lambda}\left(H, \mathcal{R}_{\alpha}^{\mathrm{OPT}}\right)=\frac{\lambda R_{\alpha}}{m}
$$

Recall $R_{\alpha}$ is the argmax of the random walk $k \mapsto \alpha \frac{k}{m}-p_{(k)}$, which has exchangeable increments. We will use Corollary 11.14 of Kallenberg (2006), due to Sparre-Andersen, that, by exchangeability, the number of rejections $R_{\alpha}$ is equal in distribution to the time the walk stays positive:

$$
R_{\alpha} \stackrel{d}{=} P_{\alpha}:=\sum_{k=1}^{m} 1\left\{p_{(k)} \leq \alpha \frac{k}{m}\right\}
$$

Under the global null, the regret thus has mean

$$
\begin{aligned}
m \mathbb{E}\left[L_{\lambda}\left(H, \mathcal{R}_{\alpha}\right)-L_{\lambda}\left(H, \mathcal{R}_{\alpha}^{\mathrm{OPT}}\right)\right] & =\lambda \mathbb{E} R_{\alpha}=\lambda \sum_{k=1}^{m} \mathbb{P}\left\{p_{(k)} \leq \alpha \frac{k}{m}\right\} \\
& \rightarrow \lambda \sum_{k=1}^{\infty} \mathbb{P}_{U_{k} \sim \operatorname{Gamma}(k, k)}\left\{U_{k} \leq \alpha\right\}
\end{aligned}
$$

where the last step follows from the law of rare events.

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[^0]:    ${ }^{1}$ Without the shape constraint on $f_{1}$, max-lfdr $>\alpha$ still implies that the analyst could improve the procedure by removing the least promising rejection, which may not be the same as the last rejection. However, this improvement is only feasible if the analyst can recognize which rejection is least promising.

