# $\ell_{1}$-penalization for mixture regression models 

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#### Abstract

We consider a finite mixture of regressions (FMR) model for highdimensional inhomogeneous data where the number of covariates may be much larger than sample size. We propose an $\ell_{1}$-penalized maximum likelihood estimator in an appropriate parameterization. This kind of estimation belongs to a class of problems where optimization and theory for non-convex functions is needed. This distinguishes itself very clearly from high-dimensional estimation with convex loss- or objective functions as, for example, with the Lasso in linear or generalized linear models. Mixture models represent a prime and important example where non-convexity arises.

For FMR models, we develop an efficient EM algorithm for numerical optimization with provable convergence properties. Our penalized estimator is numerically better posed (e.g., boundedness of the criterion function) than unpenalized maximum likelihood estimation, and it allows for effective statistical regularization including variable selection. We also present some asymptotic theory and oracle inequalities: due to non-convexity of the negative log-likelihood function, different mathematical arguments are needed than for problems with convex losses. Finally, we apply the new method to both simulated and real data.


[^0]Keywords Adaptive Lasso • Finite mixture models • Generalized EM algorithm • High-dimensional estimation • Lasso • Oracle inequality

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## 1 Introduction

In applied statistics, tremendous number of applications deal with relating a random response variable $Y$ to a set of explanatory variables or covariates $X=$ $\left(X^{(1)}, \ldots, X^{(p)}\right)$ through a regression-type model. The homogeneity assumption that the regression coefficients are the same for different observations $\left(Y_{1}, X_{1}\right), \ldots$, $\left(Y_{n}, X_{n}\right)$ is often inadequate. Parameters may change for different subgroups of observations. Such heterogeneity can be modeled with a Finite Mixture of Regressions (FMR) model. Especially with high-dimensional data, where the number of covariates $p$ is much larger than sample size $n$, the homogeneity assumption seems rather restrictive: at least a fraction of covariates may exhibit a different influence on the response among various observations (i.e., sub-populations). Hence, addressing the issue of heterogeneity in high-dimensional data is important in many practical applications. We will empirically demonstrate with real data in Sect. 7.2 that substantial prediction improvements are possible by incorporating a heterogeneity structure to the model.

We propose here an $\ell_{1}$-penalized method, i.e., a Lasso-type estimator (Tibshirani 1996), for estimating a high-dimensional Finite Mixture of Regressions (FMR) model where $p \gg n$. Our procedure is related to the proposal in Khalili and Chen (2007). However, we argue that a different parameterization leads to more efficient computation in high-dimensional optimization for which we prove numerical convergence properties. Our algorithm can easily handle problems where $p$ is in the thousands. Furthermore, regarding statistical performance, we present an oracle inequality which includes the setting where $p \gg n$ : this is very different from Khalili and Chen (2007) who use fixed $p$ asymptotics in the low-dimensional framework. Our theory for deriving oracle inequalities in the presence of non-convex loss functions, as the negative log-likelihood in a mixture model is non-convex, is rather non-standard but sufficiently general to cover other cases than FMR models.

From a more general point of view, we show in this paper that high-dimensional estimation problems with non-convex loss functions can be addressed with high computational efficiency and good statistical accuracy. Regarding the computation, we develop a rather generic block coordinate descent generalized EM algorithm which is surprisingly fast even for large $p$. Progress in efficient gradient descent methods build on various developments by Tseng (2001) and Tseng and Yun (2008), and their use for solving Lasso-type convex problems has been worked out by, e.g., Fu (1998), Friedman et al. (2007, 2008) and Meier et al. (2008). We present in Sect. 7.3 some computation times for the more involved case with non-convex objective function using a block coordinate descent generalized EM algorithm. Regarding statistical theory, almost all results for high-dimensional Lasso-type problems have been developed for convex loss functions, e.g., the squared error in a

Gaussian regression (Greenshtein and Ritov 2004; Meinshausen and Bühlmann 2006; Zhao and Yu 2006; Bunea et al. 2007; Zhang and Huang 2008; Meinshausen and Yu 2009; Wainwright 2009; Bickel et al. 2009; Cai et al. 2009b; Candès and Plan 2009; Zhang 2009) or the negative log-likelihood in a generalized linear model (van de Geer 2008). We present a non-trivial modification of the mathematical analysis of $\ell_{1}$-penalized estimation with convex loss to non-convex but smooth likelihood problems.

When estimation is defined via optimization of a non-convex objective function, there is a major gap between the actual computation and the procedure studied in theory. The computation is typically guaranteed to find a local optimum of the objective function only, whereas the theory is usually about the estimator defined by a global optimum. Particularly in high-dimensional problems, it is difficult to compute a global optimum and it would be desirable to have some theoretical properties of estimators arising from local optima. We do not provide an answer to this difficult issue in this paper. The beauty of, e.g., the Lasso or the Dantzig selector (Candès and Tao 2007) in high-dimensional problems is the provable correctness or optimality of the estimator which is actually computed. A challenge for future research is to establish such provable correctness of estimators involving non-convex objective functions. A noticeable exception is presented in Zhang (2010) for linear models, where some theory is derived for an estimator based on a local optimum of a non-convex optimization criterion.

The rest of this article is mainly focusing on Finite Mixture of Regressions (FMR) models. Some theory for high-dimensional estimation with non-convex loss functions is presented in Sect. 5 for more general settings than FMR models. The further organization of the paper is as follows: Sect. 2 describes the FMR model with an appropriate parameterization, Sect. 3 introduces $\ell_{1}$-penalized maximum-likelihood estimation, Sects. 4 and 5 present mathematical theory for the low- and high-dimensional case, Sect. 6 develops some efficient generalized EM algorithm and describes its numerical convergence properties, and Sect. 7 reports on simulations, real data analysis and computational timings.

## 2 Finite mixture of Gaussian regressions model

Our primary focus is on the following mixture model involving Gaussian components:

$$
\begin{align*}
& Y_{i} \mid X_{i} \quad \text { independent for } i=1, \ldots, n, \\
& Y_{i} \mid X_{i}=x \sim f_{\xi}(y \mid x) d y \text { for } i=1, \ldots, n, \\
& f_{\xi}(y \mid x)=\sum_{r=1}^{k} \pi_{r} \frac{1}{\sqrt{2 \pi} \sigma_{r}} \exp \left(-\frac{\left(y-x^{T} \beta_{r}\right)^{2}}{2 \sigma_{r}^{2}}\right),  \tag{2.1}\\
& \xi=\left(\beta_{1}, \ldots, \beta_{k}, \sigma_{1}, \ldots, \sigma_{k}, \pi_{1}, \ldots, \pi_{k-1}\right) \in \mathbb{R}^{k p} \times \mathbb{R}_{>0}^{k} \times \Pi, \\
& \Pi=\left\{\pi ; \pi_{r}>0 \text { for } r=1, \ldots, k-1 \text { and } \sum_{r=1}^{k-1} \pi_{r}<1\right\} .
\end{align*}
$$

Thereby, $X_{i} \in \mathbb{R}^{p}$ are fixed or random covariates, $Y_{i} \in \mathbb{R}$ is a univariate response variable and $\xi=\left(\beta_{1}, \ldots, \beta_{k}, \sigma_{1}, \ldots, \sigma_{k}, \pi_{1}, \ldots, \pi_{k-1}\right)$ denotes the $(p+2) \cdot k-1$ free parameters and $\pi_{k}$ is given by $\pi_{k}=1-\sum_{r=1}^{k-1} \pi_{r}$. The model in (2.1) is a mixture of Gaussian regressions, where every component $r$ has its individual vector of regression coefficients $\beta_{r}$ and error variances $\sigma_{r}^{2}$. We are particularly interested in the case where $p \gg n$.

### 2.1 Reparameterized mixture of regressions model

We prefer to work with a reparameterized version of model (2.1) whose penalized maximum likelihood estimator is scale-invariant and easier to compute. The computational aspect will be discussed in greater detail in Sects. 3.1 and 6. Define new parameters

$$
\phi_{r}=\beta_{r} / \sigma_{r}, \quad \rho_{r}=\sigma_{r}^{-1}, \quad r=1, \ldots, k
$$

This yields a one-to-one mapping from $\xi$ in (2.1) to a new parameter vector $\theta=$ $\left(\phi_{1}, \ldots, \phi_{k}, \rho_{1}, \ldots, \rho_{k}, \pi_{1}, \ldots, \pi_{k-1}\right)$, and the model (2.1) in reparameterized form then equals:

$$
\begin{align*}
& Y_{i} \mid X_{i} \quad \text { independent for } i=1, \ldots, n, \\
& Y_{i} \mid X_{i}=x \sim h_{\theta}(y \mid x) d y \quad \text { for } i=1, \ldots, n, \\
& h_{\theta}(y \mid x)=\sum_{r=1}^{k} \pi_{r} \frac{\rho_{r}}{\sqrt{2 \pi}} \exp \left(-\frac{1}{2}\left(\rho_{r} y-x^{T} \phi_{r}\right)^{2}\right),  \tag{2.2}\\
& \theta=\left(\phi_{1}, \ldots, \phi_{k}, \rho_{1}, \ldots, \rho_{k}, \pi_{1}, \ldots, \pi_{k-1}\right) \in \mathbb{R}^{k p} \times \mathbb{R}_{>0}^{k} \times \Pi, \\
& \Pi=\left\{\pi ; \pi_{r}>0 \text { for } r=1, \ldots, k-1 \text { and } \sum_{r=1}^{k-1} \pi_{r}<1\right\},
\end{align*}
$$

with $\pi_{k}=1-\sum_{r=1}^{k-1} \pi_{r}$. This is the main model we are analyzing and working with.
The log-likelihood function of this model equals

$$
\begin{equation*}
\ell(\theta ; Y)=\sum_{i=1}^{n} \log \left(\sum_{r=1}^{k} \pi_{r} \frac{\rho_{r}}{\sqrt{2 \pi}} \exp \left(-\frac{1}{2}\left(\rho_{r} Y_{i}-X_{i}^{T} \phi_{r}\right)^{2}\right)\right) . \tag{2.3}
\end{equation*}
$$

Since we want to deal with the $p \gg n$ case, we have to regularize the maximum likelihood estimator (MLE) in order to obtain reasonably accurate estimates. We propose below some $\ell_{1}$-norm penalized MLE which is different from a naive $\ell_{1}$-norm penalty for the MLE in the non-transformed model (2.1). Furthermore, it is well known that the (log-) likelihood function is generally unbounded. We will see in Sect. 3.2 that our penalization will mitigate this problem.

## $3 \ell_{1}$-norm penalized maximum likelihood estimator

We argue first for the case of a (non-mixture) linear model why the reparameterization above in Sect. 2.1 is useful and quite natural.

## $3.1 \ell_{1}$-norm penalization for reparameterized linear models

Consider a Gaussian linear model

$$
\begin{align*}
& Y_{i}=\sum_{j=1}^{p} \beta_{j} X_{i}^{(j)}+\varepsilon_{i}, \quad i=1, \ldots, n, \\
& \varepsilon_{1}, \ldots, \varepsilon_{n} \text { i.i.d. } \sim \mathcal{N}\left(0, \sigma^{2}\right), \tag{3.4}
\end{align*}
$$

where $X_{i}$ are either fixed or random covariates. In short, we often write

$$
Y=\mathbf{X} \beta+\varepsilon
$$

with $n \times 1$ vectors $Y$ and $\varepsilon, p \times 1$ vector $\beta$ and $n \times p$ matrix $\mathbf{X}$. In the sequel, $\|\cdot\|$ denotes the Euclidean norm. The $\ell_{1}$-norm penalized estimator, called the Lasso (Tibshirani 1996), is defined as

$$
\begin{equation*}
\hat{\beta}_{\lambda}=\underset{\beta}{\arg \min } n^{-1}\|Y-\mathbf{X} \beta\|^{2}+\lambda \sum_{j=1}^{p}\left|\beta_{j}\right| . \tag{3.5}
\end{equation*}
$$

Here $\lambda$ is a non-negative regularization parameter. The Gaussian assumption is not crucial in model (3.4) but it is useful to make connections to the likelihood framework. The Lasso estimator in (3.5) is equivalent to minimizing the penalized negative $\log$-likelihood $n^{-1} \ell\left(\beta ; Y_{1}, \ldots, Y_{n}\right)$ as a function of the regression coefficients $\beta$ and using the $\ell_{1}$-penalty $\|\beta\|_{1}=\sum_{j=1}^{p}\left|\beta_{j}\right|$ : equivalence here means that we obtain the same estimator for a potentially different tuning parameter. But the Lasso estimator in (3.5) does not provide an estimate of the nuisance parameter $\sigma^{2}$.

In mixture models, it will be crucial to have a good estimator of $\sigma^{2}$ and the role of the scaling of the variance parameter is much more important than in homogeneous regression models. Hence, it is important to take $\sigma^{2}$ into the definition and optimization of the penalized maximum likelihood estimator: we could proceed with the following estimator,

$$
\begin{align*}
\hat{\beta}_{\lambda}, \hat{\sigma}_{\lambda}^{2} & =\underset{\beta, \sigma^{2}}{\arg \min }\left(-n^{-1} \ell\left(\beta, \sigma^{2} ; Y_{1}, \ldots, Y_{n}\right)+\lambda\|\beta\|_{1}\right) \\
& =\underset{\beta, \sigma^{2}}{\arg \min }\left(\log (\sigma)+\|Y-\mathbf{X} \beta\|^{2} /\left(2 n \sigma^{2}\right)+\lambda\|\beta\|_{1}\right) . \tag{3.6}
\end{align*}
$$

Note that we are penalizing only the $\beta$-parameter. However, the scale parameter estimate $\hat{\sigma}_{\lambda}^{2}$ is influenced indirectly by the amount of shrinkage $\lambda$.

There are two main drawbacks of the estimator in (3.6). First, it is not equivariant (Lehmann 1983) under scaling of the response. More precisely, consider the transformation

$$
\begin{equation*}
Y^{\prime}=b Y, \quad \beta^{\prime}=b \beta, \quad \sigma^{\prime}=b \sigma \quad(b>0) \tag{3.7}
\end{equation*}
$$

which leaves model (3.4) invariant. A reasonable estimator based on transformed data $Y^{\prime}$ should lead to estimators $\hat{\beta}^{\prime}, \hat{\sigma}^{\prime}$ which are related to $\hat{\beta}, \hat{\sigma}$ through $\hat{\beta}^{\prime}=b \hat{\beta}$ and $\hat{\sigma}^{\prime}=b \hat{\sigma}$. This is not the case for the estimator in (3.6). Secondly, and as important as the first issue, the optimization in (3.6) is non-convex and hence, some of the major computational advantages of Lasso for high-dimensional problems is lost. We address these drawbacks by using the penalty term $\lambda \frac{\|\beta\|_{1}}{\sigma}$ leading to the following estimator:

$$
\hat{\beta}_{\lambda}, \hat{\sigma}_{\lambda}^{2}=\underset{\beta, \sigma^{2}}{\arg \min }\left(\log (\sigma)+\|Y-\mathbf{X} \beta\|^{2} /\left(2 n \sigma^{2}\right)+\lambda \frac{\|\beta\|_{1}}{\sigma}\right) .
$$

This estimator is equivariant under the scaling transformation (3.7), i.e., the estimators $\hat{\beta}^{\prime}, \hat{\sigma}^{\prime}$ based on $Y^{\prime}$ transform as $\hat{\beta}^{\prime}=b \hat{\beta}$ and $\hat{\sigma}^{\prime}=b \hat{\sigma}$. Furthermore, it penalizes the $\ell_{1}$-norm of the coefficients and small variances $\sigma^{2}$ simultaneously which has some close relations to the Bayesian Lasso (Park and Casella 2008). For the latter, a Bayesian approach is used with a conditional Laplace prior specification of the form

$$
p\left(\beta \mid \sigma^{2}\right)=\prod_{j=1}^{p} \frac{\lambda}{2 \sqrt{\sigma^{2}}} \exp \left(-\lambda \frac{\left|\beta_{j}\right|}{\sqrt{\sigma^{2}}}\right)
$$

and a noninformative scale-invariant marginal prior $p\left(\sigma^{2}\right)=1 / \sigma^{2}$ for $\sigma^{2}$. Park and Casella (2008) argue that conditioning on $\sigma^{2}$ is important because it guarantees a unimodal full posterior.

Most importantly, we can re-parameterize to achieve convexity of the optimization problem

$$
\phi_{j}=\beta_{j} / \sigma, \quad \rho=\sigma^{-1}
$$

This then yields the following estimator which is equivariant under scaling and whose computation involves convex optimization:

$$
\begin{equation*}
\hat{\phi}_{\lambda}, \hat{\rho}_{\lambda}=\underset{\phi, \rho}{\arg \min }\left(-\log (\rho)+\frac{1}{2 n}\|\rho Y-\mathbf{X} \phi\|^{2}+\lambda\|\phi\|_{1}\right) . \tag{3.8}
\end{equation*}
$$

From an algorithmic point of view, fast algorithms are available to solve the optimization in (3.8). Shooting algorithms (Fu 1998) with coordinate-wise descent are especially suitable, as demonstrated by, e.g., Friedman et al. (2007), Meier et al. (2008) or Friedman et al. (2008). We describe in Sect. 6.1 an algorithm for estimation in a mixture of regressions model, a more complex task than the optimization for (3.8). As we will see in Sect. 6.1, we will make use of the Karush-Kuhn-Tucker (KKT) conditions in the M-step of a generalized EM algorithm. For the simpler criterion in (3.8) for a non-mixture model, the KKT conditions imply the following which we state without a proof. Denote by $\langle\cdot, \cdot\rangle$ the inner product in $n$-dimensional Euclidean space and by $\mathbf{X}_{j}$ the $j$ th column of $\mathbf{X}$.

Proposition 1 Every solution $(\hat{\phi}, \hat{\rho})$ of (3.8) satisfies:

$$
\begin{aligned}
& -\hat{\rho}\left\langle\mathbf{X}_{j}, Y\right\rangle+\left\langle\mathbf{X}_{j}, \mathbf{X} \hat{\phi}\right\rangle+n \lambda \operatorname{sgn}\left(\hat{\phi}_{j}\right)=0 \quad \text { if } \hat{\phi}_{j} \neq 0, \\
& \left|-\hat{\rho}\left\langle\mathbf{X}_{j}, Y\right\rangle+\left\langle\mathbf{X}_{j}, \mathbf{X} \hat{\phi}\right\rangle\right| \leq n \lambda \quad \text { if } \hat{\phi}_{j}=0,
\end{aligned}
$$

and

$$
\hat{\rho}=\frac{\langle Y, \mathbf{X} \hat{\phi}\rangle+\sqrt{\langle Y, \mathbf{X} \hat{\phi}\rangle^{2}+4\|Y\|^{2} n}}{2\|Y\|^{2}}
$$

## $3.2 \ell_{1}$-norm penalized MLE for mixture of Gaussian regressions

Consider the mixture of Gaussian regressions model in (2.2). Assuming that $p$ is large, we want to regularize the MLE. In the spirit of the approach in (3.8), we propose for the unknown parameter $\theta=\left(\phi_{1}, \ldots, \phi_{k}, \rho_{1}, \ldots, \rho_{k}, \pi_{1}, \ldots, \pi_{k-1}\right)$ the estimator:

$$
\begin{align*}
& \hat{\theta}_{\lambda}^{(\gamma)}=\underset{\theta \in \Theta}{\arg \min }-n^{-1} \ell_{\text {pen }, \lambda}^{(\gamma)}(\theta), \quad \theta=\left(\phi_{1}, \ldots, \phi_{k}, \rho_{1}, \ldots, \rho_{k}, \pi_{1}, \ldots, \pi_{k-1}\right)  \tag{3.9}\\
& \begin{aligned}
-n^{-1} \ell_{\text {pen }, \lambda}^{(\gamma)}(\theta)= & -n^{-1} \sum_{i=1}^{n} \log \left(\sum_{r=1}^{k} \pi_{r} \frac{\rho_{r}}{\sqrt{2 \pi}} \exp \left(-\frac{1}{2}\left(\rho_{r} Y_{i}-X_{i}^{T} \phi_{r}\right)^{2}\right)\right) \\
& +\lambda \sum_{r=1}^{k} \pi_{r}^{\gamma}\left\|\phi_{r}\right\|_{1} \\
\Theta= & \mathbb{R}^{k p} \times \mathbb{R}_{>0}^{k} \times \Pi
\end{aligned}
\end{align*}
$$

where $\Pi=\left\{\pi ; \pi_{r}>0\right.$ for $r=1, \ldots, k-1$ and $\left.\sum_{r=1}^{k-1} \pi_{r}<1\right\}$ with $\pi_{k}=1-$ $\sum_{r=1}^{k-1} \pi_{r}$. The value of $\gamma \in\{0,1 / 2,1\}$ parameterizes three different penalties.

The first penalty function with $\gamma=0$ is independent of the component probabilities $\pi_{r}$. As we will see in Sects. 6.1 and 6.4, the optimization for computing $\hat{\theta}_{\lambda}^{(0)}$ is easiest, and we establish a rigorous result about numerical convergence of a generalized EM algorithm. The penalty with $\gamma=0$ works fine if the components are not very unbalanced, i.e., the true $\pi_{r}$ 's aren't too different. In case of strongly unbalanced components, the penalties with values $\gamma \in\{1 / 2,1\}$ are to be preferred, at the price of having to pursue a more difficult optimization problem. The value of $\gamma=1$ has been proposed by Khalili and Chen (2007) for the naively parameterized likelihood from model (2.1). We will report in Sect. 7.1 about empirical comparisons with the three different penalties involving $\gamma \in\{0,1 / 2,1\}$.

All three penalty functions involve the $\ell_{1}$-norm of the component specific ratios $\phi_{r}=\frac{\beta_{r}}{\sigma_{r}}$ and hence small variances are penalized. The penalized criteria therefore stay finite whenever $\sigma_{r} \rightarrow 0$ : this is in sharp contrast to the unpenalized MLE where the likelihood is unbounded if $\sigma_{r} \rightarrow 0$; see, for example, McLachlan and Peel (2000).

Proposition 2 Assume that $Y_{i} \neq 0$ for all $i=1, \ldots, n$. Then the penalized negative log-likelihood $-n^{-1} \ell_{\mathrm{pen}, \lambda}^{(0)}(\theta)$ is bounded from below for all values $\theta \in \Theta$ from (3.11).

A proof is given in Appendix C. Even though Proposition 2 is only stated and proved for the penalized negative log-likelihood with $\gamma=0$, we expect that the statement is also true for $\gamma=1 / 2$ or 1 .

Due to the $\ell_{1}$-norm penalty, the estimator is shrinking some of the components of $\phi_{1}, \ldots, \phi_{k}$ exactly to zero, depending on the magnitude of the regularization parameter $\lambda$. Thus, we can do variable selection as follows. Denote by

$$
\begin{equation*}
\widehat{S}=\left\{(r, j) ; 1 \leq r \leq k, 1 \leq j \leq p, \hat{\phi}_{r, j} \neq 0\right\} . \tag{3.12}
\end{equation*}
$$

Here, $\hat{\phi}_{r, j}$ is the $j$ th coefficient of the estimated regression parameter $\hat{\phi}_{r}$ belonging to mixture component $r$. The set $\widehat{S}$ denotes the collection of non-zero estimated, i.e., selected, regression coefficients in the $k$ mixture components. Note that no significance testing is involved, but, of course, $\widehat{S}=\widehat{S}_{\lambda}^{(\gamma)}$ depends on the specification of the regularization parameter $\lambda$ and the type of penalty described by $\gamma$.

### 3.3 Adaptive $\ell_{1}$-norm penalization

A two-stage adaptive $\ell_{1}$-norm penalization for linear models has been proposed by Zou (2006), called the adaptive Lasso. It is an effective way to address some bias problems of the (one-stage) Lasso which may employ strong shrinkage of coefficients corresponding to important variables.

The two-stage adaptive $\ell_{1}$-norm penalized estimator for a mixture of Gaussian regressions is defined as follows. Consider an initial estimate $\theta^{\text {ini }}$, for example, from the estimator in (3.9). The adaptive criterion to be minimized involves a re-weighted $\ell_{1}$-norm penalty term:

$$
\begin{align*}
&-n^{-1} \ell_{\mathrm{adapt}}^{(\gamma)}(\theta)=-n^{-1} \sum_{i=1}^{n} \log \left(\sum_{r=1}^{k} \pi_{r} \frac{\rho_{r}}{\sqrt{2 \pi}} \exp \left(-\frac{1}{2}\left(\rho_{r} Y_{i}-X_{i}^{T} \phi_{r}\right)^{2}\right)\right) \\
&+\lambda \sum_{r=1}^{k} \pi_{r}^{\gamma} \sum_{j=1}^{p} w_{r, j}\left|\phi_{r, j}\right| \\
& w_{r, j}=\frac{1}{\left|\phi_{r, j}^{\text {ini }}\right|}, \theta=\left(\rho_{1}, \ldots, \rho_{k}, \phi_{1}, \ldots, \phi_{k}, \pi_{1}, \ldots, \pi_{k-1}\right) \tag{3.13}
\end{align*}
$$

where $\gamma \in\{0,1 / 2,1\}$. The estimator is then defined as

$$
\begin{equation*}
\hat{\theta}_{\mathrm{adapt} ; \lambda}^{(\gamma)}=\underset{\theta \in \Theta}{\arg \min }-n^{-1} \ell_{\mathrm{adapt}}^{(\gamma)}(\theta), \tag{3.14}
\end{equation*}
$$

where $\Theta$ is as in (3.11).
The adaptive Lasso in linear models has better variable selection properties than the Lasso, see Zou (2006), Huang et al. (2008), van de Geer et al. (2010). We present some theory for the adaptive estimator in FMR models in Sect. 4. Furthermore, we report some empirical results in Sect. 7.1 indicating that the two-stage adaptive method often outperforms the one-stage $\ell_{1}$-penalized estimator.

### 3.4 Selection of the tuning parameters

The regularization parameters to be selected are the number of components $k$, the penalty parameter $\lambda$ and we may also want to select the type of the penalty function, i.e., selection of $\gamma$.

One possibility is to use a modified BIC criterion which minimizes

$$
\begin{equation*}
\text { BIC }=-2 \ell\left(\hat{\theta}_{\lambda, k}^{(\gamma)}\right)+\log (n) \mathrm{d}_{\mathrm{e}}, \tag{3.15}
\end{equation*}
$$

over a grid of candidate values for $k, \lambda$ and maybe also $\gamma$. Here, $\hat{\theta}_{\lambda, k}^{(\gamma)}$ denotes the estimator in (3.9) using the parameters $\lambda, k, \gamma$ in (3.10), and $-\ell(\cdot)$ is the negative log-likelihood. Furthermore, $\mathrm{d}_{\mathrm{e}}=k+(k-1)+\sum_{j=1, \ldots, p ; r=1, \ldots, k} 1_{\left\{\hat{\phi}_{r, j} \neq 0\right\}}$ is the effective number of parameters (Pan and Shen 2007).

Alternatively, we may use a cross-validation scheme for tuning parameter selection minimizing some cross-validated negative log-likelihood.

Regarding the grid of candidate values for $\lambda$, we consider $0 \leq \lambda_{1}<\cdots<\lambda_{M} \leq$ $\lambda_{\text {max }}$, where $\lambda_{\text {max }}$ is given by

$$
\begin{equation*}
\lambda_{\max }=\max _{j=1, \ldots, p}\left|\frac{\left\langle Y, \mathbf{X}_{j}\right\rangle}{\sqrt{n}\|Y\|}\right| . \tag{3.16}
\end{equation*}
$$

At $\lambda_{\text {max }}$, all coefficients $\hat{\phi}_{j}(j=1, \ldots, p)$ of the one-component model are exactly zero. Equation (3.16) easily follows from Proposition 1.

For the adaptive $\ell_{1}$-norm penalized estimator minimizing the criterion in (3.13), we proceed analogously but replacing $\hat{\theta}_{\lambda, k}^{(\gamma)}$ in (3.15) by $\hat{\theta}_{\text {adapt; }}^{(\gamma)}$ in (3.14). As initial estimator in the adaptive criterion, we propose to use the estimate in (3.9) which is optimally tuned using the modified BIC or some cross-validation scheme.

## 4 Asymptotic properties for fixed $\boldsymbol{p}$ and $\boldsymbol{k}$

Following the penalized likelihood theory of Fan and Li (2001), we establish first some asymptotic properties of the estimator in (3.10). As in Fan and Li (2001), we assume in this section that the design is random and that the number of covariates $p$ and the number of mixture components $k$ are fixed as sample size $n \rightarrow \infty$. Of course, this does not reflect a truly high-dimensional scenario, but the theory and methodology is much easier for this case. An extended theory for $p$ potentially very large in relation to $n$ is presented in Sect. 5.

Denote by $\theta_{0}$ the true parameter.
Theorem 1 (Consistency) Consider model (2.2) with random design, fixed $p$ and $k$. If $\lambda=O\left(n^{-1 / 2}\right)(n \rightarrow \infty)$ then, under the regularity conditions $(A)-(C)$ from Fan and Li (2001) on the joint density function of $(Y, X)$, there exists a local minimizer $\hat{\theta}_{\lambda}^{(\gamma)}$ of $-n^{-1} \ell_{\text {pen }, \lambda}^{(\gamma)}(\theta)$ in (3.10) $(\gamma \in\{0,1 / 2,1\})$ such that

$$
\sqrt{n}\left(\hat{\theta}_{\lambda}^{(\gamma)}-\theta_{0}\right)=O_{P}(1) .
$$

A proof is given in Appendix A. Theorem 1 can be easily misunderstood. It does not guarantee the existence of an asymptotically consistent sequence of estimates. The only claim is that a clairvoyant statistician (with pre-knowledge of $\theta_{0}$ ) can choose a consistent sequence of roots in the neighborhood of $\theta_{0}$ (van der Vaart 2007). In the case where $-n^{-1} \ell_{\text {pen }, \lambda}^{(\gamma)}(\theta)$ has a unique minimizer, which is the case for a FMR model with one component, the resulting estimator would be root-n consistent. But for a general FMR model with more than one component and typically several local minimizers, this does not hold anymore. In this sense, the preceding theorem might look better than it is.

Next, we present an asymptotic oracle result in the spirit of Fan and Li (2001) for the two-stage adaptive procedure described in Sect. 3.3. Denote by $S$ the population analogue of (3.12), i.e., the set of non-zero regression coefficients. Furthermore, let $\theta_{S}=\left(\left\{\phi_{r, j} ; \quad(r, j) \in S\right\}, \rho_{1}, \ldots, \rho_{k}, \pi_{1}, \ldots, \pi_{k-1}\right)$ be the sub-vector of parameters corresponding to the true non-zero regression coefficients (denoted by $S$ ) and analogously for $\hat{\theta}_{S}$.

Theorem 2 (Asymptotic oracle result for adaptive procedure) Consider model (2.2) with random design, fixed $p$ and $k$. If $\lambda=o\left(n^{-1 / 2}\right), n \lambda \rightarrow \infty$ and if $\theta^{\text {ini }}$ satisfies $\theta^{\text {ini }}-\theta_{0}=O_{P}\left(n^{-1 / 2}\right)$, then, under the regularity conditions $(A)-(C)$ from Fan and $\mathrm{Li}(2001)$ on the joint density function of $(Y, X)$, there exists a local minimizer $\hat{\theta}_{\text {adapt; } \lambda}^{(\gamma)}$ of $-n^{-1} \ell_{\text {adapt }}^{(\gamma)}(\theta)$ in (3.13) $(\gamma \in\{0,1 / 2,1\})$ which satisfies:

1. (Consistency in variable selection) $\mathbb{P}\left[\widehat{S}_{\text {adapt } ; \lambda}^{(\gamma)}=S\right] \rightarrow 1(n \rightarrow \infty)$.
2. (Oracle Property) $\sqrt{n}\left(\hat{\theta}_{\text {adapt } ; \lambda, S}^{(\gamma)}-\theta_{0, S}\right) \rightsquigarrow^{d} \mathcal{N}\left(0, I_{S}\left(\theta_{0}\right)^{-1}\right)$, where $I_{S}\left(\theta_{0}\right)$ is the Fisher information knowing that $\theta_{S^{c}}=0$ (i.e., the submatrix of the Fisher information at $\theta_{0}$ corresponding to the variables in $S$ ).

A proof is given in Appendix A. As in Theorem 1, the assertion of the theorem is only making a statement about some local optimum. Furthermore, this result only holds for the adaptive criterion with weights $w_{r, j}=\frac{1}{\left|\phi_{r, j}^{\text {ini }}\right|}$ coming from a root-n consistent initial estimator $\theta^{\text {ini }}$; this is a rather strong assumption given the fact that Theorem 1 only ensures existence of such an estimator. The non-adaptive estimator with the $\ell_{1}$-norm penalty as in (3.10) cannot achieve sparsity and maintain root-n consistency due to the bias problem mentioned in Sect. 3.3 (see also Khalili and Chen 2007).

## 5 General theory for high-dimensional setting with non-convex smooth loss

We present here some theory, entirely different from Theorems 1 and 2, which reflects some consistency and optimality behavior of the $\ell_{1}$-norm penalized maximum likelihood estimator for the potentially high-dimensional framework with $p \gg n$. In particular, we derive some oracle inequality which is non-asymptotic. We intentionally present this theory for $\ell_{1}$-penalized smooth likelihood problems which are generally non-convex; $\ell_{1}$-penalized likelihood estimation in FMR models is then a special case
discussed in Sect. 5.3. The following Sects. 5.1-5.2 introduce some mathematical conditions and derive auxiliary results and a general oracle inequality (Theorem 3); the interpretation of these conditions and of the oracle result is discussed for the case of FMR models at the end of Sect. 5.3.1.

### 5.1 The setting and notation

Let $\left\{f_{\psi} ; \psi \in \Psi\right\}$ be a collection of densities with respect to the Lebesgue measure $\mu$ on $\mathbb{R}$ (i.e., the range for the response variable). The parameter space $\Psi$ is assumed to be a bounded subset of some finite-dimensional space, say

$$
\Psi \subset\left\{\psi \in \mathbb{R}^{d} ;\|\psi\|_{\infty} \leq K\right\}
$$

where we have equipped (quite arbitrarily) the space $\mathbb{R}^{d}$ with the sup-norm $\|\psi\|_{\infty}=$ $\max _{1 \leq j \leq d}\left|\psi_{j}\right|$. In our setup, the dimension $d$ will be regarded as a fixed constant (which still covers high-dimensionality of the covariates, as we will see). Then, equivalent metrics are, e.g., the ones induced by the $\ell_{q}$-norm $\|\psi\|_{q}=\left(\sum_{j=1}^{d}\left|\psi_{j}\right|_{q}\right)^{1 / q}$ ( $q \geq 1$ ).

We observe a covariate $X$ in some space $\mathcal{X} \subset \mathbb{R}^{p}$ and a response variable $Y \in \mathbb{R}$. The true conditional density of $Y$ given $X=x$ is assumed to be equal to

$$
f_{\psi_{0}}(\cdot \mid x)=f_{\psi_{0}(x)},
$$

where

$$
\psi_{0}(x) \in \Psi, \quad \forall x \in \mathcal{X}
$$

That is, we assume that the true conditional density of $Y$ given $x$ is depending on $x$ only through some parameter function $\psi_{0}(x)$. Of course, the introduced notation also applies to fixed instead of random covariates.

The parameter $\left\{\psi_{0}(x) ; x \in \mathcal{X}\right\}$ is assumed to have a nonparametric part of interest $\left\{g_{0}(x) ; x \in \mathcal{X}\right\}$ and a low-dimensional nuisance part $\eta_{0}$, i.e.,

$$
\psi_{0}(\cdot)^{T}=\left(g_{0}(\cdot)^{T}, \eta_{0}^{T}\right),
$$

with

$$
g_{0}(x) \in \mathbb{R}^{k}, \quad \forall x \in \mathcal{X}, \quad \eta_{0} \in \mathbb{R}^{m}, \quad k+m=d .
$$

In case of FMR models, $g(x)^{T}=\left(\phi_{1}^{T} x, \phi_{2}^{T} x, \ldots, \phi_{k}^{T} x\right)$ and $\eta$ involves the parameters $\rho_{1}, \ldots, \rho_{k}, \pi_{1}, \ldots, \pi_{k-1}$. More details are given in Sect. 5.3.

With minus the log-likelihood as loss function, the so-called excess risk

$$
\mathcal{E}\left(\psi \mid \psi_{0}\right)=-\int \log \left[\frac{f_{\psi}}{f_{\psi_{0}}}\right] f_{\psi_{0}} d \mu
$$

is the Kullback-Leibler information. For fixed covariates $x_{1}, \ldots, x_{n}$, we define the average excess risk

$$
\overline{\mathcal{E}}\left(\psi \mid \psi_{0}\right)=\frac{1}{n} \sum_{i=1}^{n} \mathcal{E}\left(\psi\left(x_{i}\right) \mid \psi_{0}\left(x_{i}\right)\right)
$$

and for random design, we take the expectation $\mathbb{E}\left[\mathcal{E}\left(\psi(X) \mid \psi_{0}(X)\right)\right]$.

### 5.1.1 The margin

Following Tsybakov (2004) and van de Geer (2008) we call the behavior of the excess risk $\mathcal{E}\left(\psi \mid \psi_{0}\right)$ near $\psi_{0}$ the margin. We will show in Lemma 1 that the margin is quadratic.

Denote by

$$
l_{\psi}=\log f_{\psi}
$$

the log-density. Assuming the derivatives exist, we define the score function

$$
s_{\psi}=\frac{\partial l_{\psi}}{\partial \psi}
$$

and the Fisher information

$$
I(\psi)=\int s_{\psi} s_{\psi}^{T} f_{\psi} d \mu=-\int \frac{\partial^{2} l_{\psi}}{\partial \psi \partial \psi^{T}} f_{\psi} d \mu
$$

Of course, we can then also look at $I(\psi(x))$ using the parameter function $\psi(x)$.
In the sequel, we introduce some conditions (Conditions 1-5). Their interpretation for the case of FMR models is given at the end of Sect. 5.3.1. First, we will assume boundedness of the third derivatives.

Condition 1 It holds that

$$
\sup _{\psi \in \Psi} \max _{\left(j_{1}, j_{2}, j_{3}\right) \in\{1, \ldots, d\}^{3}}\left|\frac{\partial^{3}}{\partial \psi_{j_{1}} \partial \psi_{j_{2}} \partial \psi_{j_{3}}} l_{\psi}(\cdot)\right| \leq G_{3}(\cdot),
$$

where

$$
\sup _{x \in \mathcal{X}} \int G_{3}(y) f_{\psi_{0}}(y \mid x) d \mu(y) \leq C_{3}<\infty
$$

For a symmetric, positive semi-definite matrix $A$, we let $\Lambda_{\text {min }}^{2}(A)$ be its smallest eigenvalue.

Condition 2 For all $x \in \mathcal{X}$, the Fisher information matrix $I\left(\psi_{0}(x)\right)$ is positive definite and, in fact,

$$
\Lambda_{\min }=\inf _{x \in \mathcal{X}} \Lambda_{\min }\left(I\left(\psi_{0}(x)\right)\right)>0
$$

Further, we will need the following identifiability condition.
Condition 3 For all $\varepsilon>0$, there exists an $\alpha_{\varepsilon}>0$ such that

$$
\inf _{x \in \mathcal{X}} \inf _{\substack{\psi \in \Psi \\\left\|\psi-\psi_{0}(x)\right\|_{2}>\varepsilon}} \mathcal{E}\left(\psi \mid \psi_{0}(x)\right) \geq \alpha_{\varepsilon}
$$

Based on these three conditions, we have the following result:

Lemma 1 Assume Conditions 1, 2, and 3. Then

$$
\inf _{x \in \mathcal{X}} \frac{\mathcal{E}\left(\psi \mid \psi_{0}(x)\right)}{\left\|\psi-\psi_{0}(x)\right\|_{2}^{2}} \geq \frac{1}{c_{0}^{2}},
$$

where

$$
c_{0}^{2}=\max \left[\frac{1}{\varepsilon_{0}}, \frac{d K^{2}}{\alpha_{\varepsilon_{0}}}\right], \quad \varepsilon_{0}=\frac{3 \Lambda_{\min }^{2}}{2 d^{3 / 2}} .
$$

A proof is given in Appendix B.

### 5.1.2 The empirical process

We now specialize to the case where

$$
\psi(x)^{T}=\left(g(x)^{T}, \eta^{T}\right)
$$

where (with some abuse of notation)

$$
\begin{aligned}
& g(x)^{T}=g_{\phi}(x)^{T}=\left(g_{1}(x), \ldots, g_{k}(x)\right), \\
& g_{r}(x)=g_{\phi_{r}}(x)=x^{T} \phi_{r}, \quad x \in \mathbb{R}^{p}, \phi_{r} \in \mathbb{R}^{p}, r=1, \ldots, k .
\end{aligned}
$$

We also write

$$
\psi_{\vartheta}(x)^{T}=\left(g_{\phi}(x)^{T}, \eta^{T}\right), \quad \vartheta^{T}=\left(\phi_{1}^{T}, \ldots, \phi_{k}^{T}, \eta^{T}\right)
$$

to make the dependence of the parameter function $\psi(x)$ on $\vartheta$ more explicit.
We will assume that

$$
\sup _{x \in \mathcal{X}}\left\|\phi^{T} x\right\|_{\infty}=\sup _{x \in \mathcal{X}} \max _{1 \leq r \leq k}\left|\phi_{r}^{T} x\right| \leq K .
$$

This can be viewed as a combined condition on $\mathcal{X}$ and $\phi$. For example, if $\mathcal{X}$ is bounded by a fixed constant this supremum (for fixed $\phi$ ) is finite.

Our parameter space is now

$$
\begin{equation*}
\tilde{\Theta} \subset\left\{\vartheta^{T}=\left(\phi_{1}^{T}, \ldots, \phi_{k}^{T}, \eta^{T}\right) ; \sup _{x \in \mathcal{X}}\left\|\phi^{T} x\right\|_{\infty} \leq K,\|\eta\|_{\infty} \leq K\right\} . \tag{5.17}
\end{equation*}
$$

Note that $\tilde{\Theta}$ is, in principle, $(p k+m)$-dimensional. The true parameter $\vartheta_{0}$ is assumed to be an element of $\tilde{\Theta}$.

Let us define

$$
\begin{aligned}
& L_{\vartheta}(x, \cdot)=\log f_{\psi(x)}(\cdot), \quad \psi(x)^{T}=\psi_{\vartheta}(x)^{T}=\left(g_{\phi}(x)^{T}, \eta^{T}\right), \\
& \vartheta^{T}=\left(\phi_{1}^{T}, \ldots, \phi_{k}^{T}, \eta^{T}\right),
\end{aligned}
$$

and the empirical process for fixed covariates $x_{1}, \ldots, x_{n}$

$$
V_{n}(\vartheta)=\frac{1}{n} \sum_{i=1}^{n}\left(L_{\vartheta}\left(x_{i}, Y_{i}\right)-\mathbb{E}\left[L_{\vartheta}\left(x_{i}, Y\right) \mid X=x_{i}\right]\right)
$$

We now fix some $T \geq 1$ and $\lambda_{0} \geq 0$ and define the event

$$
\begin{equation*}
\mathcal{T}=\left\{\sup _{\vartheta^{T}=\left(\phi^{T}, \eta^{T}\right) \in \tilde{\Theta}} \frac{\left|V_{n}(\vartheta)-V_{n}\left(\vartheta_{0}\right)\right|}{\left(\left\|\phi-\phi_{0}\right\|_{1}+\left\|\eta-\eta_{0}\right\|_{2}\right) \vee \lambda_{0}} \leq T \lambda_{0}\right\} . \tag{5.18}
\end{equation*}
$$

### 5.2 Oracle inequality for the Lasso for non-convex loss functions

For an optimality result, we need some condition on the design. Denote the active set, i.e., the set of non-zero coefficients, by

$$
S=\left\{(r, j) ; \phi_{r, j} \neq 0\right\}, \quad s=|S|,
$$

and let

$$
\phi_{J}=\left\{\phi_{(r, j)} ;(r, j) \in J\right\}, \quad J \subseteq\{1, \ldots, k\} \times\{1, \ldots, p\} .
$$

Further, let

$$
\Sigma_{n}=\frac{1}{n} \sum_{i=1}^{n} x_{i} x_{i}^{T}
$$

Condition 4 (Restricted eigenvalue condition) There exists a constant $\kappa \geq 1$ such that, for all $\phi \in \mathbb{R}^{p k}$ satisfying

$$
\left\|\phi_{S^{c}}\right\|_{1} \leq 6\left\|\phi_{S}\right\|_{1},
$$

it holds that

$$
\left\|\phi_{S}\right\|_{2}^{2} \leq \kappa^{2} \sum_{r=1}^{k} \phi_{r}^{T} \Sigma_{n} \phi_{r}
$$

For $\psi(\cdot)^{T}=\left(g(\cdot)^{T}, \eta^{T}\right)$, we use the notation

$$
\|\psi\|_{Q_{n}}^{2}=\frac{1}{n} \sum_{i=1}^{n} \sum_{r=1}^{k} g_{r}^{2}\left(x_{i}\right)+\sum_{j=1}^{m} \eta_{j}^{2} .
$$

We also write for $g(\cdot)=\left(g_{1}(\cdot), \ldots, g_{k}(\cdot)\right)^{T}$,

$$
\|g\|_{Q_{n}}^{2}=\frac{1}{n} \sum_{i=1}^{n} \sum_{r=1}^{k} g_{r}^{2}\left(x_{i}\right)
$$

Thus

$$
\left\|g_{\phi}\right\|_{Q_{n}}^{2}=\sum_{r=1}^{k} \phi_{r}^{T} \Sigma_{n} \phi_{r}
$$

and the bound in the restricted eigenvalue condition then reads

$$
\left\|\phi_{S}\right\|_{2}^{2} \leq \kappa^{2}\left\|g_{\phi}\right\|_{Q_{n}}^{2}
$$

Bounding $\left\|g_{\phi}\right\|_{Q_{n}}^{2}$ in terms of $\left\|\phi_{S}\right\|_{2}^{2}$ can be done directly using, e.g., the CauchySchwarz inequality. The restricted eigenvalue condition ensures a bound in the other direction which itself is needed for an oracle inequality. Some references about the restricted eigenvalue condition are provided at the end of Sect. 5.3.1.

We employ the Lasso-type estimator

$$
\begin{equation*}
\hat{\vartheta}^{T}=\left(\hat{\phi}^{T}, \hat{\eta}^{T}\right)=\underset{\vartheta^{T}=\left(\phi^{T}, \eta^{T}\right) \in \tilde{\Theta}}{\arg \min }\left\{-\frac{1}{n} \sum_{i=1}^{n} L_{\vartheta}\left(x_{i}, Y_{i}\right)+\lambda \sum_{r=1}^{k}\left\|\phi_{r}\right\|_{1}\right\} \tag{5.19}
\end{equation*}
$$

We omit in the sequel the dependence of $\hat{\vartheta}$ on $\lambda$. Note that we consider here a global minimizer which may be difficult to compute if the empirical risk $n^{-1} \sum_{i=1}^{n} L_{\vartheta}\left(x_{i}, Y_{i}\right)$ is non-convex in $\vartheta$. We then write $\|\phi\|_{1}=\sum_{r=1}^{k}\left\|\phi_{r}\right\|_{1}$. We let

$$
\hat{\psi}(x)^{T}=\left(g_{\hat{\phi}}(x)^{T}, \hat{\eta}^{T}\right)
$$

which depends only on the estimate $\hat{\vartheta}$, and we denote by

$$
\psi_{0}(x)^{T}=\left(g_{\phi_{0}}(x)^{T}, \eta_{0}^{T}\right)
$$

Theorem 3 (Oracle result for fixed design) Assume fixed covariates $x_{1}, \ldots, x_{n}$, Conditions $1-3$ and 4 , and that $\lambda \geq 2 T \lambda_{0}$ for the estimator in (5.19) with $T$ and $\lambda_{0}$ as in (5.18). Then on $\mathcal{T}$, defined in (5.18), for the average excess risk (average KullbackLeibler loss),

$$
\overline{\mathcal{E}}\left(\hat{\psi} \mid \psi_{0}\right)+2\left(\lambda-T \lambda_{0}\right)\left\|\hat{\phi}_{S^{c}}\right\|_{1} \leq 8\left(\lambda+T \lambda_{0}\right)^{2} c_{0}^{2} \kappa^{2} s
$$

where $c_{0}$ and $\kappa$ are defined in Lemma 1 and Condition 4, respectively.
A proof is given in Appendix B. We will give an interpretation of this result in Sect. 5.3.1, where we specialize to FMR models. In the case of FMR models, the probability of the set $\mathcal{T}$ is large as shown in detail by Lemma 3 below.

Before specializing to FMR models, we present more general results for lower bounding the probability of the set $\mathcal{T}$. We make the following assumption.

Condition 5 For the score function $s_{\vartheta}(\cdot)=s_{\psi_{\vartheta}}(\cdot)$, we have

$$
\sup _{\vartheta \in \tilde{\Theta}}\left\|s_{\vartheta}(\cdot)\right\|_{\infty} \leq G_{1}(\cdot),
$$

for some function $G_{1}(\cdot)$.

Condition 5 primarily has notational character. Later, in Lemma 2 and particularly in Lemma 3, the function $G_{1}(\cdot)$ needs to be sufficiently regular to ensure small corresponding probabilities.

Define

$$
\begin{equation*}
\lambda_{0}=M_{n} \log n \sqrt{\frac{\log (p \vee n)}{n}} . \tag{5.20}
\end{equation*}
$$

As we will see, we usually choose $M_{n} \asymp \sqrt{\log (n)}$. Let $\mathbb{P}_{\mathbf{x}}$ denote the conditional probability given $\left(X_{1}, \ldots, X_{n}\right)=\left(x_{1}, \ldots, x_{n}\right)=\mathbf{x}$, and with the expression $1\{\cdot\}$ we denote the indicator function.

Lemma 2 Assume Condition 5. We have for constants $c_{1}, c_{2}$ and $c_{3}$ depending on $k$ and $K$, and for all $T \geq c_{1}$,

$$
\sup _{\vartheta^{T}=\left(\phi^{T}, \eta^{T}\right) \in \Theta} \frac{\left|V_{n}(\vartheta)-V_{n}\left(\vartheta_{0}\right)\right|}{\left(\left\|\phi-\phi_{0}\right\|_{1}+\left\|\eta-\eta_{0}\right\|_{2}\right) \vee \lambda_{0}} \leq T \lambda_{0},
$$

with $\mathbb{P}_{\mathbf{x}}$ probability at least

$$
1-c_{2} \exp \left[-\frac{T^{2} \log ^{2} n \log (p \vee n)}{c_{3}^{2}}\right]-\mathbb{P}_{\mathbf{x}}\left(\frac{1}{n} \sum_{i=1}^{n} F\left(Y_{i}\right)>T \lambda_{0}^{2} /(d K)\right)
$$

where $($ for $i=1, \ldots, n)$

$$
F\left(Y_{i}\right)=G_{1}\left(Y_{i}\right) 1\left\{G_{1}\left(Y_{i}\right)>M_{n}\right\}+\mathbb{E}\left[G_{1}(Y) 1\left\{G_{1}(Y)>M_{n}\right\} \mid X=x_{i}\right]
$$

Regarding the constants $\lambda_{0}$ and $K$, see (5.20) and (5.17), respectively.
A proof is given in Appendix B.

### 5.3 FMR models

In the finite mixture of regressions model from (2.2) with $k$ components, the parameter is $\vartheta^{T}=\left(\phi^{T}, \eta^{T}\right)=\left(\phi_{1}^{T}, \ldots, \phi_{k}^{T}, \log \rho_{1}, \ldots, \log \rho_{k}, \log \pi_{1}, \ldots, \log \pi_{k-1}\right)$, where the $\rho_{r}=\sigma_{r}^{-1}$ are the inverse standard deviations in mixture component $r$ and the $\pi_{r}$ are the mixture coefficients. For mathematical convenience and simpler notation, we consider here the log-transformed $\rho$ and $\pi$ parameters in order to have lower and upper bounds for $\rho$ and $\pi$. Obviously, there is a one-to-one correspondence between $\vartheta$ and $\theta$ from Sect. 2.1.

Let the parameter space be

$$
\begin{gather*}
\tilde{\Theta} \subset\left\{\vartheta^{T} ; \sup _{x \in \mathcal{X}}\left\|\phi^{T} x\right\|_{\infty} \leq K,\|\log \rho\|_{\infty} \leq K,-K \leq \log \pi_{1} \leq 0, \ldots,\right. \\
\left.-K \leq \log \pi_{k-1} \leq 0, \quad \sum_{r=1}^{k-1} \pi_{r}<1\right\}, \tag{5.21}
\end{gather*}
$$

and $\pi_{k}=1-\sum_{r=1}^{k-1} \pi_{r}$.

We consider the estimator

$$
\begin{equation*}
\hat{\vartheta}_{\lambda}=\underset{\vartheta \in \tilde{\Theta}}{\arg \min }-n^{-1} \sum_{i=1}^{n} \log \left(\sum_{r=1}^{k} \pi_{r} \frac{\rho_{r}}{\sqrt{2 \pi}} \exp \left(-\frac{1}{2}\left(\rho_{r} Y_{i}-X_{i}^{T} \phi_{r}\right)^{2}\right)\right)+\lambda \sum_{r=1}^{k}\left\|\phi_{r}\right\|_{1} . \tag{5.22}
\end{equation*}
$$

This is the estimator from Sect. 3.2 with $\gamma=0$. We emphasize the boundedness of the parameter space by using the notation $\tilde{\Theta}$. In contrast to Sect. 4 , we focus here on any global minimizer of the penalized negative log-likelihood which is arguably difficult to compute.

In the following, we transform the estimator $\hat{\vartheta}_{\lambda}$ to $\hat{\theta}_{\lambda}$ in the parameterization $\theta$ from Sect. 2.1. Using some abuse of notation we denote the average excess risk by $\overline{\mathcal{E}}\left(\hat{\theta}_{\lambda} \mid \theta_{0}\right)$.

### 5.3.1 Oracle result for FMR models

We specialize now our results from Sect. 5.2 to FMR models.
Proposition 3 For fixed design FMR models as in (2.2) with $\tilde{\Theta}$ in (5.21), Conditions 1,2 and 3 are met, for appropriate $C_{3}, \Lambda_{\min }$ and $\left\{\alpha_{\varepsilon}\right\}$, depending on $k$ and $K$. Also Condition 5 holds with

$$
G_{1}(y)=\mathrm{e}^{K}|y|+K
$$

Proof This follows from straightforward calculations.
In order to show that the probability for the set $\mathcal{T}$ is large, we invoke Lemma 2 and the following result.

Lemma 3 For fixed design FMR models as in (2.2) with $\tilde{\Theta}$ in (5.21), for some constants $c_{4}, c_{5}$ and $c_{6}$, depending on $k$, and $K$, and for $M_{n}=c_{4} \sqrt{\log n}$ and $n \geq c_{6}$, the following holds:

$$
\mathbb{P}_{\mathbf{x}}\left(\frac{1}{n} \sum_{i=1}^{n} F\left(Y_{i}\right)>c_{5} \frac{\log n}{n}\right) \leq \frac{1}{n},
$$

where $($ for $i=1, \ldots, n)$

$$
F\left(Y_{i}\right)=G_{1}\left(Y_{i}\right) 1\left\{G_{1}\left(Y_{i}\right)>M_{n}\right\}+\mathbb{E}\left[G_{1}(Y) 1\left\{G_{1}(Y)>M_{n}\right\} \mid X=x_{i}\right],
$$

and $G_{1}(\cdot)$ is as in Proposition 3.
A proof is given in Appendix B.
Hence, the oracle result in Theorem 3 for our $\ell_{1}$-norm penalized estimator in the FMR model holds on a set $\mathcal{T}$, summarized in Theorem 4, and this set $\mathcal{T}$ has large probability due to Lemma 2 and Lemma 3 as described in the following corollary.

Corollary 1 For fixed design FMR models as in (2.2) with $\tilde{\Theta}$ in (5.21), we have for constants $c_{2}, c_{4}, c_{7}, c_{8}$ depending on $k$ and $K$,

$$
\mathbb{P}_{\mathbf{x}}[\mathcal{T}] \geq 1-c_{2} \exp \left[-\frac{T^{2} \log ^{2} n \log (p \vee n)}{c_{7}^{2}}\right]-n^{-1} \text { for all } n \geq c_{8}
$$

where $\mathcal{T}$ is defined with $\lambda_{0}=M_{n} \log n \sqrt{\frac{\log (p \vee n)}{n}}$ and $M_{n}=c_{4} \sqrt{\log n}$.
Theorem 4 (Oracle result for FMR models) Consider a fixed design FMR model as in (2.2) with $\tilde{\Theta}$ in (5.21). Assume Condition 4 (restricted eigenvalue condition) and that $\lambda \geq 2 T \lambda_{0}$ for the estimator in (5.22). Then on $\mathcal{T}$, which has large probability as stated in Corollary 1, for the average excess risk (average Kullback-Leibler loss),

$$
\overline{\mathcal{E}}\left(\hat{\theta}_{\lambda} \mid \theta_{0}\right)+2\left(\lambda-T \lambda_{0}\right)\left\|\hat{\phi}_{S^{c}}\right\|_{1} \leq 8\left(\lambda+T \lambda_{0}\right)^{2} c_{0}^{2} \kappa^{2} s,
$$

where $c_{0}$ and $\kappa$ are defined in Lemma 1 and Condition 4, respectively.
The oracle inequality of Theorem 4 has the following well-known interpretation. First, we obtain

$$
\overline{\mathcal{E}}\left(\hat{\theta}_{\lambda} \mid \theta_{0}\right) \leq 8\left(\lambda+T \lambda_{0}\right)^{2} c_{0}^{2} \kappa^{2} s
$$

That is, the average Kullback-Leibler risk is of the order $O\left(s \lambda_{0}^{2}\right)=O\left(s \log ^{3} n \log (p \vee\right.$ $n) / n$ ) (take $\lambda=2 T \lambda_{0}$, use definition (5.20) and the assumption on $M_{n}$ in Lemma 3 above) which is up to the factor $\log ^{3} n \log (p \vee n)$ the optimal convergence rate if one would know the $s$ non-zero coefficients. As a second implication we obtain

$$
\left\|\hat{\phi}_{S^{c}}\right\|_{1} \leq 4\left(\lambda+T \lambda_{0}\right) c_{0}^{2} \kappa^{2} s,
$$

saying that the noise components in $S^{c}$ have small estimated values (e.g., its $\ell_{1}$-norm converges to zero at rate $\left.O\left(s \lambda_{0}\right)\right)$.

Note that the Conditions 1, 2, 3 and 5 hold automatically for FMR models, as described in Proposition 3. We do require a restricted eigenvalue condition on the design, here Condition 4. In fact, for the Lasso or Dantzig selector in linear models, restricted eigenvalue conditions (Koltchinskii 2009; Bickel et al. 2009) are considerably weaker than coherence conditions (Bunea et al. 2007; Cai et al. 2009a) or assuming the restricted isometry property (Candès and Tao 2005; Cai et al. 2009b); for an overview among the relations, see van de Geer and Bühlmann (2009).

### 5.3.2 High-dimensional consistency of FMR models

We finally give a consistency result for FMR models under weaker conditions than the oracle result from Sect. 5.3.1. Denote by $\theta_{0}$ the true parameter vector in a FMR model. In contrast to Sect. 4, the number of covariates $p$ can grow with the number of observations $n$. Therefore, also the true parameter $\theta_{0}$ depends on $n$. To guarantee consistency we have to assume some sparsity condition, i.e., the $\ell_{1}$-norm of the true parameter can only grow with $o\left(\sqrt{n /\left(\log ^{3} n \log (p \vee n)\right)}\right)$.

Theorem 5 (Consistency) Consider a fixed design FMR model (2.2) with $\tilde{\Theta}$ in (5.21) and fixed $k$. Moreover, assume that

$$
\left\|\phi_{0}\right\|_{1}=\sum_{r=1}^{k}\left\|\phi_{0, r}\right\|_{1}=o\left(\sqrt{n /\left(\log ^{3} n \log (p \vee n)\right.}\right) \quad(n \rightarrow \infty) .
$$

If $\lambda=C \sqrt{\log ^{3} n \log (p \vee n) / n}$ for some $C>0$ sufficiently large, then any (global) minimizer $\hat{\theta}_{\lambda}$ as in (5.22) satisfies

$$
\overline{\mathcal{E}}\left(\hat{\theta}_{\lambda} \mid \theta_{0}\right)=o_{P}(1) \quad(n \rightarrow \infty) .
$$

A proof is given in Appendix B. The (restricted eigenvalue) Condition 4 on the design is not required; this is typical for a high-dimensional consistency result, see Greenshtein and Ritov (2004) for the Lasso in linear models.

## 6 Numerical optimization

We present a generalized EM (GEM) algorithm for optimizing the criterion in (3.10) in Sect. 6.1. In Sects. 6.2 and 6.3, we give further details on speeding-up and on initializing the algorithm. Finally, we discuss numerical convergence properties in Sect. 6.4. For the convex penalty $(\gamma=0)$ function, we prove convergence to a stationary point.

### 6.1 GEM algorithm for optimization

Maximization of the log-likelihood of a mixture density is often done using the traditional EM algorithm of Dempster et al. (1977). Consider the complete log-likelihood

$$
\ell_{c}(\theta ; Y, \Delta)=\sum_{i=1}^{n} \sum_{r=1}^{k} \Delta_{i, r} \log \left(\frac{\rho_{r}}{\sqrt{2 \pi}} e^{-\frac{1}{2}\left(\rho_{r} Y_{i}-X_{i}^{T} \phi_{r}\right)^{2}}\right)+\Delta_{i, r} \log \left(\pi_{r}\right) .
$$

Here $\left(\Delta_{i, 1}, \ldots, \Delta_{i, k}\right), i=1, \ldots, n$, are i.i.d. unobserved multinomial variables showing the component-membership of the $i$ th observation in the FMR model: $\Delta_{i, r}=1$ if observation $i$ belongs to component $r$, and $\Delta_{i, r}=0$ otherwise. The expected complete (scaled) negative log-likelihood is then

$$
Q\left(\theta \mid \theta^{\prime}\right)=-n^{-1} \mathbb{E}_{\theta^{\prime}}\left[\ell_{c}(\theta ; Y, \Delta) \mid Y\right]
$$

and the expected complete (scaled) penalized negative log-likelihood is

$$
Q_{\text {pen }}\left(\theta \mid \theta^{\prime}\right)=Q\left(\theta \mid \theta^{\prime}\right)+\lambda \sum_{r=1}^{k} \pi_{r}^{\gamma}\left\|\phi_{r}\right\|_{1} .
$$

The EM algorithm works by alternating between the E- and M-step. Denote the parameter value at EM-iteration $m$ by $\theta^{(m)}(m=0,1,2, \ldots)$, where $\theta^{(0)}$ is a vector of starting values.

E-Step: Compute $Q\left(\theta \mid \theta^{(m)}\right)$, or equivalently,

$$
\begin{aligned}
\hat{\gamma}_{i, r} & =\mathbb{E}_{\theta^{(m)}}\left[\Delta_{i, r} \mid Y\right]=\frac{\pi_{r}^{(m)} \rho_{r}^{(m)} e^{-\frac{1}{2}\left(\rho_{r}^{(m)} Y_{i}-X_{i}^{T} \phi_{r}^{(m)}\right)^{2}}}{\sum_{l=1}^{k} \pi_{l}^{(m)} \rho_{l}^{(m)} e^{-\frac{1}{2}\left(\rho_{l}^{(m)} Y_{i}-X_{i}^{T} \phi_{l}^{(m)}\right)^{2}}}, \\
r & =1, \ldots, k, i=1, \ldots, n .
\end{aligned}
$$

Generalized M-Step: Improve $Q_{\text {pen }}\left(\theta \mid \theta^{(m)}\right)$ w.r.t. $\theta \in \Theta$.
(a) Improvement with respect to $\pi=\left(\pi_{1}, \ldots, \pi_{k}\right)$ :

Fix $\phi$ at the present value $\phi^{(m)}$ and improve

$$
\begin{equation*}
-n^{-1} \sum_{i=1}^{n} \sum_{r=1}^{k} \hat{\gamma}_{i, r} \log \left(\pi_{r}\right)+\lambda \sum_{r=1}^{k} \pi_{r}^{\gamma}\left\|\phi_{r}^{(m)}\right\|_{1} \tag{6.23}
\end{equation*}
$$

with respect to the probability simplex

$$
\left\{\pi ; \pi_{r}>0 \text { for } r=1, \ldots, k \text { and } \sum_{r=1}^{k} \pi_{r}=1\right\} .
$$

Denote by $\bar{\pi}^{(m+1)}=\frac{\sum_{i=1}^{n} \hat{\gamma}_{i}}{n}$ which is a feasible point of the simplex. We propose to update $\pi$ as

$$
\pi^{(m+1)}=\pi^{(m)}+t^{(m)}\left(\bar{\pi}^{(m+1)}-\pi^{(m)}\right),
$$

where $t^{(m)} \in(0,1]$. In practice, $t^{(m)}$ is chosen to be the largest value in the grid $\left\{\delta^{k} ; k=0,1,2, \ldots\right\}(0<\delta<1)$ such that (6.23) is not increased. In our examples, $\delta=0.1$ worked well.
(b) Coordinate descent improvement with respect to $\phi$ and $\rho$ :

A simple calculation shows that the M-step decouples for each component into $k$ distinct optimization problems of the form

$$
\begin{equation*}
-\log \left(\rho_{r}\right)+\frac{1}{2 n_{r}}\left\|\rho_{r} \tilde{Y}-\tilde{\mathbf{X}} \phi_{r}\right\|^{2}+\frac{n \lambda}{n_{r}}\left(\pi_{r}^{(m+1)}\right)^{\gamma}\left\|\phi_{r}\right\|_{1}, \quad r=1, \ldots, k \tag{6.24}
\end{equation*}
$$

with

$$
n_{r}=\sum_{i=1}^{n} \hat{\gamma}_{i, r}, \quad\left(\tilde{Y}_{i}, \tilde{X}_{i}\right)=\sqrt{\hat{\gamma}_{i, r}}\left(Y_{i}, X_{i}\right), \quad r=1, \ldots, k .
$$

Problem (6.24) has the same form as (3.8); in particular, it is convex in ( $\rho_{r}, \phi_{r, 1}, \ldots, \phi_{r, p}$ ). Instead of fully optimizing (6.24), we only minimize with respect to each of the coordinates, holding the other coordinates at their current value. Closed-form coordinate updates can easily be computed for each component $r(r=1, \ldots, k)$ using Proposition 1 :

$$
\rho_{r}^{(m+1)}=\frac{\left\langle\tilde{Y}, \tilde{\mathbf{X}} \phi_{r}^{(m)}\right\rangle+\sqrt{\left\langle\tilde{Y}, \tilde{\mathbf{X}} \phi_{r}^{(m)}\right\rangle^{2}+4\|\tilde{Y}\|^{2} n_{r}}}{2\|\tilde{Y}\|^{2}}
$$

$$
\phi_{r, j}^{(m+1)}= \begin{cases}0 & \text { if }\left|S_{j}\right| \leq n \lambda\left(\pi_{r}^{(m+1)}\right)^{\gamma}, \\ \left(n \lambda\left(\pi_{r}^{(m+1)}\right)^{\gamma}-S_{j}\right) /\left\|\tilde{\mathbf{X}}_{j}\right\|^{2} & \text { if } S_{j}>n \lambda\left(\pi_{r}^{(m+1)}\right)^{\gamma}, \\ -\left(n \lambda\left(\pi_{r}^{(m+1)}\right)^{\gamma}+S_{j}\right) /\left\|\tilde{\mathbf{X}}_{j}\right\|^{2} & \text { if } S_{j}<-n \lambda\left(\pi_{r}^{(m+1)}\right)^{\gamma},\end{cases}
$$

where $S_{j}$ is defined as

$$
S_{j}=-\rho_{r}^{(m+1)}\left\langle\tilde{\mathbf{X}}_{j}, \tilde{Y}\right\rangle+\sum_{s<j} \phi_{r, s}^{(m+1)}\left\langle\tilde{\mathbf{X}}_{j}, \tilde{\mathbf{X}}_{s}\right\rangle+\sum_{s>j} \phi_{r, s}^{(m)}\left\langle\tilde{\mathbf{X}}_{j}, \tilde{\mathbf{X}}_{s}\right\rangle
$$

and $j=1, \ldots, p$.
Because we only improve $Q_{\text {pen }}\left(\theta \mid \theta^{(m)}\right)$ instead of a full minimization, see M-step (a) and (b), this is a generalized EM (GEM) algorithm. We call it the block coordinate descent generalized EM algorithm (BCD-GEM); the word block refers to the fact that we are updating all components of $\pi$ at once. Its numerical properties are discussed in Sect. 6.4.

Remark 1 For the convex penalty function with $\gamma=0$, a minimization with respect to $\pi$ in M-step (a) is achieved with $\pi^{(m+1)}=\frac{\sum_{i=1}^{n} \hat{\gamma}_{i}}{n}$, i.e., using $t^{(m)}=1$. Then, our M-step corresponds to exact coordinate-wise minimization of $Q_{\text {pen }}\left(\theta \mid \theta^{(m)}\right)$.

### 6.2 Active set algorithm

There is a simple way to speed-up the algorithm described above. When updating the coordinates $\phi_{r, j}$ in the M-step (b), we restrict ourselves during every 10 EM -iterations to the current active set (the non-zero coordinates) and visit the remaining coordinates every 11th EM-iteration to update the active set. In very high-dimensional and sparse settings, this leads to a remarkable decrease in computational times. A similar active set strategy is also used in Friedman et al. (2007) and Meier et al. (2008). We illustrate in Sect. 7.3 the gain of speed when staying during every 10 EM-iterations within the active set.

### 6.3 Initialization

The algorithm of Sect. 6.1 requires the specification of starting values $\theta^{(0)}$. We found empirically that the following initialization works well. For each observation $i$, $i=1, \ldots, n$, draw randomly a class $\kappa \in\{1, \ldots, k\}$. Assign for observation $i$ and the corresponding component $\kappa$ the weight $\tilde{\gamma}_{i, \kappa}=0.9$ and weights $\tilde{\gamma}_{i, r}=0.1$ for all other components. Finally, normalize $\tilde{\gamma}_{i, r}, r=1, \ldots, k$, to achieve that summing over the indices $k$ yields the value one, to get the normalized values $\hat{\gamma}_{i, r}$. Note that this can be viewed as an initialization of the E-step. In the M -step which follows afterwards, we update all coordinates from the initial values $\phi_{r, j}^{(0)}=0, \rho_{r}^{(0)}=2, \pi_{r}^{(0)}=1 / k$, $r=1, \ldots, k, j=1, \ldots, p$.

### 6.4 Numerical convergence of the BCD-GEM algorithm

We address here the convergence properties of the BCD-GEM algorithm described in Sect. 6.1. A detailed account of the convergence properties of the EM algorithm in a
general setting has been given by Wu (1983). Under regularity conditions, including differentiability and continuity, convergence to stationary points is proved for the EM algorithm. For the GEM algorithm, similar statements are true under conditions which are often hard to verify.

As a GEM algorithm, our BCD-GEM algorithm has the descent property which means that the criterion function is reduced at each iteration:

$$
\begin{equation*}
-n^{-1} \ell_{\text {pen }, \lambda}^{(\gamma)}\left(\theta^{(m+1)}\right) \leq-n^{-1} \ell_{\text {pen }, \lambda}^{(\gamma)}\left(\theta^{(m)}\right) . \tag{6.25}
\end{equation*}
$$

Since $-n^{-1} \ell_{\mathrm{pen}, \lambda}^{(0)}(\theta)$ is bounded from below (Proposition 2), the following result holds.

Proposition 4 For the BCD-GEM algorithm, $-n^{-1} \ell_{\mathrm{pen}, \lambda}^{(0)}\left(\theta^{(m)}\right)$ decreases monotonically to some value $\bar{\ell}>-\infty$.

In Remark 1, we noted that, for the convex penalty function with $\gamma=0$, the M-step of the algorithm corresponds to exact coordinate-wise minimization of $Q_{\text {pen }}\left(\theta \mid \theta^{(m)}\right)$. In this case, convergence to a stationary point can be shown.

Theorem 6 Consider the BCD-GEM algorithm for the criterion function in (3.10) with $\gamma=0$. Then, every cluster point $\bar{\theta} \in \Theta$ of the sequence $\left\{\theta^{(m)} ; m=0,1,2, \ldots\right\}$, generated by the BCD-GEM algorithm, is a stationary point of the criterion function in (3.10).

A precise definition of a stationary point in a non-differentiable setup and a proof of the Theorem are given in Appendix C. The proof uses the crucial facts that $Q_{\text {pen }}\left(\theta \mid \theta^{\prime}\right)$ is a convex function in $\theta$ and that it is strictly convex in each coordinate of $\theta$.

## 7 Simulations, real data example and computational timings

### 7.1 Simulations

We consider four different simulation setups. Simulation scenario 1 compares the performance of the unpenalized MLE with our estimators from Sect. 3.2 (FMRLasso) and Sect. 3.3 (FMRAdapt) in a situation where the total number of noise covariates grows successively. For computing the unpenalized MLE, we used the R-package FlexMix (Leisch 2004; Grün and Leisch 2007, 2008); Simulation 2 explores sparsity; Simulation 3 compares cross-validation and BIC; and Simulation 4 compares the different penalty functions with the parameters $\gamma=0,1 / 2,1$. For every setting, the results are based on 100 independent simulation runs.

All simulations are based on Gaussian FMR models as in (2.1); the coefficients $\pi_{r}, \beta_{r}, \sigma_{r}$ and the sample size $n$ are specified below. The covariate $X$ is generated from a multivariate normal distribution with mean 0 and covariance structure as specified below.

Unless otherwise specified, the penalty with $\gamma=1$ is used in all simulations. As explored empirically in Simulation 4, in case of balanced problems (approximately equal $\pi_{r}$ ), the FMRLasso performs similarly for all three penalties. In unbalanced situations, the best results are typically achieved with $\gamma=1$. In addition, unless otherwise specified, the true number of components $k$ is assumed to be known.

For all models, training-, validation- and test data are generated of equal size $n$. The estimators are computed on the training data, with the tuning parameter (e.g., $\lambda$ ) selected by minimizing twice the negative log-likelihood (log-likelihood loss) on the validation data. As performance measure, the predictive log-likelihood loss (twice the negative log-likelihood) of the selected model is computed on the test data.

Regarding variable selection, we count a covariable $X^{(j)}$ as selected if $\hat{\beta}_{r, j} \neq 0$ for at least one $r \in\{1, \ldots, k\}$. To assess the performance of FMRLasso on recovering the sparsity structure, we report the number of truly selected covariates (True Positives) and falsely selected covariates (False Positives).

Obviously, the performances depend on the signal-to-noise ratio (SNR) which we define for an FMR model as

$$
\operatorname{SNR}=\frac{\operatorname{Var}(Y)}{\operatorname{Var}\left(Y \mid \beta_{r}=0 ; r=1, \ldots, k\right)}=\frac{\sum_{r=1}^{k} \pi_{r}\left(\beta_{r}^{T} \operatorname{Cov}(X) \beta_{r}+\sigma_{r}^{2}\right)}{\sum_{r=1}^{k} \pi_{r} \sigma_{r}^{2}}
$$

where the last equality follows since $\mathbb{E}[X]=0$.

### 7.1.1 Simulation 1

We consider five different FMR models: M1, M2, M3, M4 and M5. The parameters ( $\pi_{k}, \beta_{k}, \sigma_{k}$ ), the sample size $n$ of the training-, validation- and test-data, the correlation structure of covariates $\operatorname{corr}_{l, m}=\operatorname{corr}\left(X^{(l)}, X^{(m)}\right)$ and the signal-to-noise ratio (SNR) are specified in Table 1. Models M1, M2, M3 and M5 have two components and five active covariates, whereas model M4 has three components and six active covariates. M1, M2 and M3 differ only in their variances $\sigma_{1}^{2}, \sigma_{2}^{2}$ and hence have different signal-to-noise ratios. Model M5 has a non-diagonal covariance structure. Furthermore, in model M5, the variances $\sigma_{1}^{2}, \sigma_{2}^{2}$ are tuned to achieve the same signal-to-noise ratio as in model M1.

We compare the performances of the maximum likelihood estimator (MLE), the FMRLasso and the FMRAdapt in a situation where the number of noise covariates grows successively. For the models M1, M2, M3, M5 with two components, we start with $p_{\text {tot }}=5$ (no noise covariates) and go up to $p_{\text {tot }}=125$ (120 noise covariates). For the three component model M4, we start with $p_{\text {tot }}=6$ (no noise covariates) and go up to $p_{\text {tot }}=155$ (149 noise covariates).

The boxplots in Figs. 1-5 of the predictive log-likelihood loss, denoted by Error, the True Positives (TP) and the False Positives (FP) over 100 simulation runs summarize the results for the different models. We read off from these boxplots that the MLE performs very badly when we add noise covariates. On the other hand, our penalized estimators remain stable. For example, for M1 the MLE with $p_{\text {tot }}=20$ performs worse than the FMRLasso with $p_{\text {tot }}=125$, or for M4 the MLE with $p_{\text {tot }}=10$ performs worse than the FMRLasso with $p_{\text {tot }}=75$. Impressive is also the huge gain of the FMRAdapt method over FMRLasso in terms of log-likelihood loss and false positives.
Table 1 Models for simulation 1. $\delta_{l, m}$ denotes Kronecker's delta

|  | M1 | M2 | M3 | M4 | M5 |
| :--- | :--- | :--- | :--- | :--- | :--- |
| $n$ | 100 | 100 | 100 | 150 | 100 |
| $\beta_{1}$ | $(3,3,3,3,3)$ | $(3,3,3,3,3)$ | $(3,3,3,3,3)$ | $(3,3,0,0,0,0)$ | $(3,3,3,3,3)$ |
| $\beta_{2}$ | $(-1,-1,-1,-1,-1)$ | $(-1,-1,-1,-1,-1)$ | $(-1,-1,-1,-1,-1)$ | $(0,0,-2,-2,0,0)$ | $(-1,-1,-1,-1,-1)$ |
| $\beta_{3}$ | - | - | - | $(0,0,0,0,-3,2)$ | - |
| $\sigma$ | $0.5,0.5$ | 1,1 | $0.5,0.5,0.5$ | $0.95,0.95$ |  |
| $\pi$ | $0.5,0.5$ | $0.5,0.5$ | $0.5,0.5$ | $1 / 3,1 / 3,1 / 3$ | $0.5,0.5$ |
| corr $_{l, m}$ | $\delta_{l, m}$ | $\delta_{l, m}$ | $\delta_{l, m}$ | $0.8^{l-m \mid}$ |  |
| SNR | 101 | 26 | 12.1 | 53 | 101 |



Fig. 1 Simulation 1, Model M1. (Top) Predictive log-likelihood loss (Error) for MLE, FMRLasso, FMRAdapt; (Bottom) False Positives ( $F P$ ) and True Positives ( $T P$ ) for FMRLasso and FMRAdapt

### 7.1.2 Simulation 2

In this section, we explore the sparsity properties of the FMRLasso. The model specifications are given in Table 2. Consider the ratios $p_{\text {act }}: n: p_{\text {tot }}$. The total number of covariates $p_{\text {tot }}$ grows faster than the number of observations $n$ and the number of active covariates $p_{\text {act }}$ : when $p_{\text {tot }}$ is doubled, $p_{\text {act }}$ is raised by one and $n$ is raised by 50 from model to model. In particular, we obtain a series of models which gets "sparser" as $n$ grows (larger ratio $n / p_{\text {act }}$ ). In order to compare the performance of the FMRLasso, we report the True Positive Rate (TPR) and the False Positive Rate (FPR) defined as:

$$
\begin{aligned}
\mathrm{TPR} & =\frac{\text { \#truly selected covariates }}{\text { \#active covariates }}, \\
\mathrm{FPR} & =\frac{\text { \#falsely selected covariates }}{\text { \#inactive covariates }}
\end{aligned}
$$

These numbers are reported in Fig. 6. We see that the False Positive Rate approaches zero for sparser models, indicating that the FMRLasso recovers the true model better in sparser settings regardless of the large number of noise covariates.


Fig. 2 Simulation 1, Model M2. Same notation as in Fig. 1
Table 2 Series of models for simulation 2 which gets "sparser" as $n$ grows: when $p_{\text {tot }}$ is doubled, $p_{\text {act }}$ is raised by one and $n$ is raised by 50 from model to model

| $p_{\text {act }}$ | 3 | 4 | 5 | 6 | 7 | 8 | 9 |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- |
| $n$ | 50 | 100 | 150 | 200 | 250 | 300 | 350 |
| $p_{\text {tot }}$ | 10 | 20 | 40 | 80 | 160 | 320 | 640 |
| $\beta_{1}$ |  |  |  | $(3,3,3,0,0, \ldots)$ |  |  |  |
| $\beta_{2}$ |  |  |  | $(-1,-1,-1,0,0, \ldots)$ |  |  |  |
| $\sigma$ |  |  | $0.5,0.5$ |  |  |  |  |
| $\pi$ |  |  | $0.5,0.5$ |  |  |  |  |

### 7.1.3 Simulation 3

So far, we regarded the number $k$ of components as given, while we have chosen an optimal $\lambda_{\text {opt }}$ by minimizing the negative log-likelihood loss on validation data. In this section, we compare the performance of 10 -fold cross-validation and the BIC criterion presented in Sect. 3.4 for selecting the tuning parameters $k$ and $\lambda$. We use model M1 of Sect. 7.1.1 with $p_{\text {tot }}=25,50,75$. For each of these models, we tune the FMRLasso estimator according to the following strategies:
(1) Assume the number of components is given $(k=2)$. Choose the optimal tuning parameter $\lambda_{\text {opt }}$ using 10 -fold cross-validation.


Fig. 3 Simulation 1, Model M3. Same notation as in Fig. 1
(2) Assume the number of components is given $(k=2)$. Choose $\lambda_{\text {opt }}$ by minimizing the BIC criterion (3.15).
(3) Choose the number of components $k \in\{1,2,3\}$ and $\lambda_{\text {opt }}$ by minimizing the BIC criterion (3.15).

The results of this simulation are presented in Fig. 7, where boxplots of the loglikelihood loss (Error) are shown. All three strategies perform equally well. With $p_{\text {tot }}=25$ the BIC criterion in strategy (3) always chooses $k=2$. For the model with $p_{\text {tot }}=50$, strategy (3) chooses $k=2$ in 98 simulation runs and $k=3$ in two runs. Finally, with $p_{\text {tot }}=75$, the third strategy chooses $k=2$ in 92 runs and $k=3$ eight times.

### 7.1.4 Simulation 4

In the preceding simulations, we always used the value $\gamma=1$ in the penalty term of the FMRLasso estimator (3.10). In this section, we compare the FMRLasso for different values $\gamma=0,1 / 2$, 1 . First, we compute the FMRLasso for $\gamma=0,1 / 2,1$ on model M1 of Sect. 7.1.1 with $p_{\text {tot }}=50$. Then we do the same calculations for an "unbalanced" version of this model with $\pi_{1}=0.3$ and $\pi_{2}=0.7$.

In Fig. 8, the boxplots of the log-likelihood loss (Error), the False Positives (FP) and the True Positives (TP) over 100 simulation runs are shown. We see that the FMRLasso performs similarly for $\gamma=0,1 / 2,1$. Nevertheless, the value $\gamma=1$ is slightly preferable in the "unbalanced" setup.


Fig. 4 Simulation 1, Model M4. Same notation as in Fig. 1

### 7.2 Real data example

We now apply the FMRLasso to a dataset of riboflavin (vitamin $B_{2}$ ) production by Bacillus Subtilis. The real-valued response variable is the logarithm of the riboflavin production rate. The data has been kindly provided by DSM (Switzerland). There are $p=4088$ covariates (genes) measuring the logarithm of the expression level of 4088 genes and measurements of $n=146$ genetically engineered mutants of Bacillus Subtilis. The population seems to be rather heterogeneous as there are different strains of Bacillus Subtilis which are cultured under different fermentation conditions. We do not know the different homogeneity subgroups. For this reason, a FMR model with more than one component might be more appropriate than a simple linear regression model.

We compute the FMRLasso estimator for $k=1, \ldots, 5$ components. To keep the computational effort reasonable, we use only the 100 covariates (genes) exhibiting the highest empirical variances. We choose the optimal tuning parameter $\lambda_{\text {opt }}$ by $10-$ fold cross-validation (using the log-likelihood loss). As a result, we get five different estimators which we compare according to their cross-validated log-likelihood loss (CV Error). These numbers are plotted in Fig. 9. The estimator with three components performs clearly best, resulting in a $17 \%$ improvement in prediction over a (non-mixture) linear model, and it selects 51 genes. In Fig. 10, the coefficients of the 20 most important genes, ordered according to $\sum_{r=1}^{3}\left|\hat{\beta}_{r, j}\right|$, are shown. From the


Fig. 5 Simulation 1, Model M5. Same notation as in Fig. 1
important variables, only gene 83 shows the opposite sign of the estimated regression coefficients among the three different mixture components. However, it happens that some covariates (genes) exhibit a strong effect in one or two mixture components but none in the remaining other components. Finally, for comparison, the onecomponent (non-mixture) model selects 26 genes, of which 24 are also selected in the three-component model.

### 7.3 Computational timings

In this section, we report on the run times of the BCD-GEM algorithm on two high-dimensional examples. In particular, we focus on the substantial gain of speed achieved by using the active set version of the algorithm described in Sect. 6.2. All computations were carried out with the statistical computing language and environment R. Timings depend on the stopping criterion used in the algorithm. We stop the algorithm if the relative function improvement and the relative change of the parameter vector are small enough, i.e.,

$$
\begin{aligned}
& \left.\frac{\mid \ell_{\text {pen }, \lambda}^{(\gamma)}}{\left.1+\mid \theta^{(m+1)}\right)-\ell_{\text {pen }, \lambda}^{(\gamma)}\left(\theta^{(m+1)}\right) \mid}\left(\theta^{(m)}\right) \right\rvert\, \\
& 1, \\
& \max _{j}\left\{\frac{\left|\theta_{j}^{(m+1)}-\theta_{j}^{(m)}\right|}{1+\left|\theta_{j}^{(m+1)}\right|}\right\} \leq \sqrt{\tau}, \quad \tau=10^{-6} .
\end{aligned}
$$



Fig. 6 Simulation 2 compares the performance of the FMRLasso for a series of models which gets "sparser" as the sample size grows. (Top) True Positive Rate (TPR); (Bottom) False Positive Rate (FPR) over 100 simulation runs

We consider a high-dimensional version of the two component model M1 from Sect. 7.1.1 with $n=200, p_{\text {tot }}=1000$ and the riboflavin dataset from Sect. 7.2 with three components, $n=146$ and $p_{\text {tot }}=100$. We use the BCD-GEM algorithm with and without active set strategy to fit the FMRLasso on a small grid of eight values for $\lambda$. The corresponding BIC, CPU times (in seconds) and number of EM-iterations are reported in Tables 3 and 4. The values for the BCD-GEM without active set strategy are written in brackets. For model M1 and an appropriate $\lambda$ with minimal BIC score, the active set algorithm converges in 5.96 seconds whereas the standard BCDGEM needs 53.15 seconds. There is also a considerable gain of speed for the real data: 0.89 seconds versus 3.57 seconds for $\lambda$ with optimal BIC. Note that in Table 3,


Fig. 7 Simulation 3 compares different strategies for choosing the tuning parameters $k$ and $\lambda$. The boxplots show the predictive log-likelihood loss (Error) of the FMRLasso, tuned by strategies (1), (2) and (3), for model M1 with $p_{\text {tot }}=25,50,75$
the BIC scores sometimes differ substantially for inappropriate values of $\lambda$. For such regularization parameters, the solutions are unstable and different local optima are attained depending on the algorithm used. However, if the regularization parameter is in a reasonable range with low BIC score, the results stabilize.

## 8 Discussion

We have presented an $\ell_{1}$-penalized estimator for a finite mixture of high-dimensional Gaussian regressions where the number of covariates may greatly exceed sample size. Such a model and the corresponding Lasso-type estimator are useful to blindly account for often encountered inhomogeneity of high-dimensional data. On a highdimensional real data example, we demonstrate a $17 \%$ gain in prediction accuracy over a (non-mixture) linear model.

The computation and mathematical analysis in such a high-dimensional mixture model is challenging due to the non-convex behavior of the negative log-likelihood. Moreover, with high-dimensional estimation defined via optimization of a nonconvex objective function, there is a major gap between the actual computation and the procedure analyzed in theory. We do not provide an answer to this issue in this paper. Regarding the computation in FMR models, a simple reparameterization is very beneficial and the $\ell_{1}$-penalty term makes the optimization problem numerically much better behaved. We develop an efficient generalized EM algorithm and we prove its numerical convergence to a stationary point. Regarding the statistical properties, besides standard low-dimensional asymptotics, we present a non-asymptotic oracle inequality for the Lasso-type estimator in a high-dimensional setting with general, nonconvex but smooth loss functions. The mathematical arguments are different than what is typically used for convex losses.
Table 3 Model M1 with $n=200$ and $p_{\text {tot }}=1000$. Median over 10 simulation runs of BIC, CPU times and number of EM-iterations for the BCD-GEM with and without active set strategy (the latter in brackets)

| $\lambda$ | 10.0 | 15.6 | 21.1 | 26.7 | 32.2 | 37.8 | 48.9 |  |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- |
| BIC | $2033(2022)$ | $1606(1748)$ | $951(959)$ | $941(940)$ | $989(983)$ | $1236(1073)$ | $1214(1216)$ | $1206(1203)$ |
| CPU [s] | $26.78(269.91)$ | $17.05(165.78)$ | $8.63(82.78)$ | $5.96(53.15)$ | $5.08(44.23)$ | $4.23(37.27)$ | $3.35(18.99)$ | $3.30(15.62)$ |
| \# EM-iter. | $277.0(341.5)$ | $196.0(205.0)$ | $96.0(100.5)$ | $63.5(64.5)$ | $56.0(53.5)$ | $41.5(46.0)$ | $31.5(23.0)$ | $25.0(19.0)$ |

Table 4 Riboflavin data with $k=3, n=146$ and $p_{\text {tot }}=100$. BIC, CPU times and number of EM-iterations for the BCD-GEM with and without active set strategy (the latter in brackets)

| $\lambda$ | 3.0 | 13.8 | 24.6 | 35.4 | 46.2 | 57.0 | 78.6 |  |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- |
| BIC | $560(628)$ | $536(530)$ | $516(522)$ | $532(525)$ | $541(540)$ | $561(580)$ | $592(591)$ |  |
| CPU [s] | $22.40(29.98)$ | $1.35(3.28)$ | $0.89(3.57)$ | $0.86(3.34)$ | $0.78(3.87)$ | $0.69(2.42)$ | $0.37(2.56)$ | $611(613)$ |
| \# EM-iter. | $3389(2078)$ | $345(239)$ | $287(266)$ | $298(247)$ | $296(290)$ | $248(184)$ | $129(192)$ | $313(302)$ |



Fig. 8 Simulation 4 compares the FMRLasso for different values $\gamma=0,1 / 2,1$. The upper row of the panels shows the boxplots of the log-likelihood loss (Error), the False Positives (FP) and the True Positives $(T P)$ for model M1 with $p_{\mathrm{tot}}=50$ and $\pi_{1}=\pi_{2}=0.5$. The lower row of the panels shows the same boxplots for an "unbalanced" version of model M1 with $\pi_{1}=0.3$ and $\pi_{2}=0.7$

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## Appendix A: Proofs for Sect. 4

## A. 1 Proof of Theorem 1

We assume the regularity assumptions (A)-(C) of Fan and Li (2001). The theorem follows from Theorem 1 of Fan and Li (2001).

Fig. 9 Riboflavin production data. Cross-validated negative log-likelihood loss (CV Error) for the FMRLasso estimator when varying over different numbers of components


## A. 2 Proof of Theorem 2

In order to keep the notation simple, we give the proof for a two class mixture with $k=2$. All arguments in the proof can also be used for a general mixture with more than two components. Remember that $-n^{-1} \ell_{\text {adapt }}(\theta)$ is given by

$$
-n^{-1} \ell_{\text {adapt }}(\theta)=-n^{-1} \ell(\theta)+\lambda\left(\pi_{1}^{\gamma} \sum_{j=1}^{p} w_{1, j}\left|\phi_{1, j}\right|+\left(1-\pi_{1}\right)^{\gamma} \sum_{j=1}^{p} w_{2, j}\left|\phi_{2, j}\right|\right),
$$

where $\ell(\theta)$ is the log-likelihood function. The weights $w_{r, j}$ are given by $w_{r, j}=\frac{1}{\left|\phi_{r, j}^{\text {ini }}\right|}$, $r=1,2$, and $j=1, \ldots, p$.

Assertion 1. Let $\hat{\theta}$ be a root- $n$ consistent local minimizer of $-n^{-1} \ell_{\text {adapt }}(\theta)$ (construction as in Fan and Li 2001).

For all $(r, j) \in S$, we easily see from consistency of $\hat{\theta}$ that $\mathbb{P}[(r, j) \in \hat{S}] \rightarrow 1$. It then remains to show that for all $(r, j) \in S^{c}, \mathbb{P}\left[(r, j) \in \hat{S}^{c}\right] \rightarrow 1$. Assume the contrary, i.e., w.l.o.g. there is an $s \in\{1, \ldots, p\}$ with $\phi_{1, s}=0$ such that $\hat{\phi}_{1, s} \neq 0$ with nonvanishing probability.

By Taylor's theorem, applied to the function $n^{-1} \frac{\partial \ell(\theta)}{\partial \phi_{1, s}}$, there exists a (random) vector $\tilde{\theta}$ on the line segment between $\theta_{0}$ and $\hat{\theta}$ such that

$$
\begin{aligned}
& \left.\frac{1}{n} \frac{\partial \ell_{\text {adapt }}}{\partial \phi_{1, s}}\right|_{\hat{\theta}} \\
& \quad=\underbrace{\left.\frac{1}{n} \frac{\partial \ell}{\partial \phi_{1, s}}\right|_{\theta_{0}}}_{(1)}+\underbrace{\left.\frac{1}{n} \frac{\partial \ell^{\prime}}{\partial \phi_{1, s}}\right|_{\theta_{0}}}_{(2)}\left(\hat{\theta}-\theta_{0}\right)
\end{aligned}
$$



Now, using the regularity assumptions and the central limit theorem, term (1) is of order $O_{P}\left(\frac{1}{\sqrt{n}}\right)$. Similarly, term (2) is of order $O_{P}(1)$ by the law of large numbers. Term (3) is of order $O_{P}(1)$ by the law of large numbers and the regularity condition on the third derivatives (condition (C) of Fan and Li 2001 ).

Therefore, we have
$\left.\frac{1}{n} \frac{\partial \ell_{\text {adapt }}}{\partial \phi_{1, s}}\right|_{\hat{\theta}}=O_{P}\left(\frac{1}{\sqrt{n}}\right)+\left(O_{P}(1)+\left(\hat{\theta}-\theta_{0}\right)^{T} O_{P}(1)\right)\left(\hat{\theta}-\theta_{0}\right)-\lambda \hat{\pi}_{1}^{\gamma} w_{1, s} \operatorname{sgn}\left(\hat{\phi}_{1, s}\right)$.
As $\hat{\theta}$ is root- $n$ consistent we get

$$
\begin{align*}
\left.\frac{1}{n} \frac{\partial \ell_{\text {adapt }}}{\partial \phi_{1, s}}\right|_{\hat{\theta}}= & O_{P}\left(\frac{1}{\sqrt{n}}\right)+\left(O_{P}(1)+o_{P}(1) O_{P}(1)\right) \\
& \times O_{P}\left(\frac{1}{\sqrt{n}}\right)-\lambda \hat{\pi}_{1}^{\gamma} w_{1, s} \operatorname{sgn}\left(\hat{\phi}_{1, s}\right) \\
= & \frac{1}{\sqrt{n}}\left(O_{P}(1)-\frac{n \lambda}{\sqrt{n}} \hat{\pi}_{1}^{\gamma} w_{1, s} \operatorname{sgn}\left(\hat{\phi}_{1, s}\right)\right) . \tag{A.26}
\end{align*}
$$

From the assumption on the initial estimator, we have

$$
\frac{n \lambda}{\sqrt{n}} w_{1, s}=\frac{n \lambda}{\sqrt{n}\left|\phi_{1, s}^{\text {ini }}\right|}=\frac{n \lambda}{O_{P}(1)} \rightarrow \infty \quad \text { as } n \lambda \rightarrow \infty
$$

Therefore, the second term in the brackets of (A.26) dominates the first and the probability of the event

$$
\left\{\operatorname{sgn}\left(\left.\frac{1}{n} \frac{\partial \ell_{\text {adapt }}}{\partial \phi_{1, s}}\right|_{\hat{\theta}}\right)=-\operatorname{sgn}\left(\hat{\phi}_{1, s}\right) \neq 0\right\}
$$

tends to 1 . But this contradicts the assumption that $\hat{\theta}$ is a local minimizer (i.e., $\left.\left.\frac{1}{n} \frac{\partial \ell_{\text {adapt }}}{\partial \phi_{1, s}}\right|_{\hat{\theta}}=0\right)$.

Assertion 2. Write $\theta=\left(\theta_{S}, \theta_{S^{c}}\right)$. From part (1), it follows that with probability tending to one $\hat{\theta}_{S}$ is a root- $n$ local minimizer of $-n^{-1} \ell_{\text {adapt }}\left(\theta_{S}, 0\right)$.

By using a Taylor expansion,

$$
0=\left.\frac{1}{n} \ell_{\text {adapt }}^{\prime}\right|_{\hat{\theta}_{S}}=\left.\frac{1}{n} \ell^{\prime}\right|_{\theta_{0, S}}+\underbrace{\left.\frac{1}{n} \ell^{\prime \prime}\right|_{\theta_{0, S}}}_{(1)}\left(\hat{\theta}_{S}-\theta_{0, S}\right)+\frac{1}{2} \underbrace{\left(\hat{\theta}_{S}-\theta_{0, S}\right)^{T}}_{(2)} \underbrace{\left.\frac{1}{n} \ell^{\prime \prime \prime}\right|_{\tilde{\theta}_{S}}}_{(3)}\left(\hat{\theta}_{S}-\theta_{0, S}\right)
$$

$$
-\lambda\left(\begin{array}{c}
\hat{\pi}_{1}^{\gamma} w_{1, S} \operatorname{sgn}\left(\hat{\phi}_{1, S}\right) \\
\left(1-\hat{\pi}_{1}\right)^{\gamma} w_{2, S} \operatorname{sgn}\left(\hat{\phi}_{2, S}\right) \\
0 \\
0 \\
\gamma \hat{\pi}_{1}^{\gamma-1} \sum_{(1, j) \in S} w_{1, j}\left|\hat{\phi}_{1, j}\right|-\gamma\left(1-\hat{\pi}_{1}\right)^{\gamma-1} \sum_{(2, j) \in S} w_{2, j}\left|\hat{\phi}_{2, j}\right|
\end{array}\right)
$$

Now term (1) is of order $-I_{S}\left(\theta_{0}\right)+o_{P}(1)$ (law of large numbers); term (2) is of order $o_{P}(1)$ (consistency); and term (3), with some abuse of notation an $(|S|+3)$-vector of $(|S|+3) \times(|S|+3)$ matrices, is of order $O_{P}(1)$ (law of large numbers and regularity condition on the third derivatives). Therefore, we have

$$
\left.\sqrt{n} \frac{1}{n} \ell^{\prime}\right|_{\theta_{0, S}}+\left(-I_{S}\left(\theta_{0}\right)+o_{P}(1)\right) \sqrt{n}\left(\hat{\theta}_{S}-\theta_{0, S}\right)-\sqrt{n} \lambda O_{P}(1)=0,
$$

or

$$
\begin{equation*}
\left(-I_{S}\left(\theta_{0}\right)+o_{P}(1)\right) \sqrt{n}\left(\hat{\theta}_{S}-\theta_{0, S}\right)-\sqrt{n} \lambda O_{P}(1)=-\left.\frac{1}{\sqrt{n}} \ell^{\prime}\right|_{\theta_{0, S}} \tag{A.27}
\end{equation*}
$$

Notice that $\left.\frac{1}{\sqrt{n}} \ell^{\prime}\right|_{\theta_{0, S}} \rightsquigarrow{ }^{d} \mathcal{N}\left(0, I_{S}\left(\theta_{0}\right)\right)$ by the central limit theorem. Furthermore, $\sqrt{n} \lambda=o(1)$ as $\lambda=o\left(n^{-1 / 2}\right)$.

Therefore, $\sqrt{n}\left(\hat{\theta}_{S}-\theta_{0, S}\right) \rightsquigarrow^{d} \mathcal{N}\left(0, I_{S}\left(\theta_{0}\right)^{-1}\right)$ follows from (A.27).

## Appendix B: Proofs for Sect. 5

## B. 1 Proof of Lemma 1

Using a Taylor expansion,

$$
\mathcal{E}\left(\psi \mid \psi_{0}\right)=\left(\psi-\psi_{0}\right)^{T} I\left(\psi_{0}\right)\left(\psi-\psi_{0}\right) / 2+r_{\psi}
$$

where

$$
\begin{aligned}
\left|r_{\psi}\right| & \leq \frac{\left\|\psi-\psi_{0}\right\|_{1}^{3}}{6} \int \sup _{\psi \in \Psi} \max _{j_{1}, j_{2}, j_{3}}\left|\frac{\partial^{3} l_{\psi}}{\partial \psi_{j_{1}} \partial \psi_{j_{2}} \partial \psi_{j_{3}}}\right| f_{\psi_{0}} d \mu \\
& \leq \frac{d^{3 / 2} C_{3}}{6}\left\|\psi-\psi_{0}\right\|_{2}^{3} .
\end{aligned}
$$

Hence

$$
\mathcal{E}\left(\psi \mid \psi_{0}(x)\right) \geq\left\|\psi-\psi_{0}(x)\right\|_{2}^{2} \Lambda_{\min }^{2} / 2-d^{3 / 2} C_{3}\left\|\psi-\psi_{0}(x)\right\|_{2}^{3} / 6
$$

Now, apply the auxiliary lemma below, with $K_{0}^{2}=d K^{2}, \Lambda^{2}=\Lambda_{\min }^{2} / 2$, and $C=$ $d^{3 / 2} C_{3} / 6$.

Auxiliary Lemma Let $h:\left[-K_{0}, K_{0}\right] \rightarrow[0, \infty)$ have the following properties:
(i) $\forall \varepsilon>0 \exists \alpha_{\varepsilon}>0$ such that $\inf _{\varepsilon<|z| \leq K_{0}} h(z) \geq \alpha_{\varepsilon}$,
(ii) $\exists \Lambda>0, C>0$, such that $\forall|z| \leq K_{0}, h(z) \geq \Lambda^{2} z^{2}-C|z|^{3}$.

Then $\forall|z| \leq K_{0}$,

$$
h(z) \geq z^{2} / C_{0}^{2}
$$

where

$$
C_{0}^{2}=\max \left[\frac{1}{\varepsilon_{0}}, \frac{K_{0}^{2}}{\alpha_{\varepsilon_{0}}}\right], \quad \varepsilon_{0}=\frac{\Lambda^{2}}{2 C}
$$

Proof (Auxiliary Lemma) If $\varepsilon_{0}>K_{0}$, we have $h(z) \geq \Lambda^{2} z^{2} / 2$ for all $|z| \leq K_{0}$.
If $\varepsilon_{0} \leq K_{0}$ and $|z| \leq \varepsilon_{0}$, we also have $h(z) \geq\left(\Lambda^{2}-\varepsilon_{0} C\right) z^{2} \geq \Lambda^{2} z^{2} / 2$.
If $\varepsilon_{0} \leq K_{0}$ and $\varepsilon_{0}<|z| \leq K_{0}$, we have $h(z) \geq \alpha_{\varepsilon_{0}}=K_{0}^{2} \alpha_{\varepsilon_{0}} / K_{0}^{2} \geq|z|^{2} \alpha_{\varepsilon_{0}} / K_{0}^{2}$.

## B. 2 Proof of Lemma 2

In order to prove Lemma 2, we first state and prove a suitable entropy bound:
We introduce the norm

$$
\|h(\cdot, \cdot)\|_{P_{n}}=\sqrt{\frac{1}{n} \sum_{i=1}^{n} h^{2}\left(x_{i}, Y_{i}\right)}
$$

For a collection $\mathcal{H}$ of functions on $\mathcal{X} \times \mathcal{Y}$, we let $H\left(\cdot, \mathcal{H},\|\cdot\|_{P_{n}}\right)$ be the entropy of $\mathcal{H}$ equipped with the metric induced by the norm $\|\cdot\|_{P_{n}}$ (for a definition of the entropy of a metric space, see van de Geer 2000).

Define for $\epsilon>0$,

$$
\tilde{\Theta}(\epsilon)=\left\{\vartheta^{T}=\left(\phi_{1}^{T}, \ldots, \phi_{k}^{T}, \eta^{T}\right) \in \tilde{\Theta}:\left\|\phi-\phi_{0}\right\|_{1}+\left\|\eta-\eta_{0}\right\|_{2} \leq \epsilon\right\} .
$$

Entropy Lemma For a constant $C_{0}$ depending on $k$ and $m$, we have for all $u>0$ and $M_{n}>0$,

$$
H\left(u,\left\{\left(L_{\vartheta}-L_{\vartheta^{*}}\right) 1\left\{G_{1} \leq M_{n}\right\}: \vartheta \in \tilde{\Theta}(\epsilon)\right\},\|\cdot\|_{P_{n}}\right) \leq C_{0} \frac{\epsilon^{2} M_{n}^{2}}{u^{2}} \log \left(\frac{\epsilon M_{n}}{u}\right)
$$

Proof (Entropy Lemma) We have

$$
\begin{aligned}
\left|L_{\vartheta}(x, y)-L_{\vartheta}(x, y)\right|^{2} & \leq G_{1}^{2}(y)\left[\sum_{r=1}^{k}\left|\left(\phi_{r}-\tilde{\phi}_{r}\right)^{T} x\right|+\|\eta-\tilde{\eta}\|_{1}\right]^{2} \\
& \leq d G_{1}^{2}(y)\left[\sum_{r=1}^{k}\left|\left(\phi_{r}-\tilde{\phi}_{r}\right)^{T} x\right|^{2}+\|\eta-\tilde{\eta}\|_{2}^{2}\right] .
\end{aligned}
$$

It follows that

$$
\left\|\left(L_{\vartheta}-L_{\tilde{\vartheta}}\right) 1\left\{G_{1} \leq M_{n}\right\}\right\|_{P_{n}}^{2} \leq d M_{n}^{2}\left[\sum_{r=1}^{k} \frac{1}{n} \sum_{i=1}^{n}\left|\left(\phi_{r}-\tilde{\phi}_{r}\right)^{T} x_{i}\right|^{2}+\|\eta-\tilde{\eta}\|_{2}^{2}\right]
$$

Let $N(\cdot, \Gamma, \tau)$ denote the covering number of a metric space $(\Gamma, \tau)$ with the metric (induced by the norm) $\tau$, and let $H(\cdot, \Gamma, \tau)=\log N(\cdot, \Gamma, \tau)$ be its entropy (for a definition of the covering number of a metric space, see van de Geer 2000). If $\Gamma$ is a ball with radius $\epsilon$ in the Euclidean space $\mathbb{R}^{N}$, one easily verifies that

$$
H(u, \Gamma, \tau) \leq N \log \left(\frac{5 \epsilon}{u}\right), \quad \forall u>0
$$

Thus $H\left(u,\left\{\eta \in \mathbb{R}^{m}:\left\|\eta-\eta_{0}\right\|_{2} \leq \epsilon\right\},\|\cdot\|_{2}\right) \leq m \log \left(\frac{5 \epsilon}{u}\right), \forall u>0$. Moreover, applying a bound as in Lemma 2.6.11 of van der Vaart and Wellner (1996) gives

$$
H\left(2 u,\left\{\sum_{r=1}^{k}\left(\phi_{r}-\phi_{0, r}\right)^{T} x_{r}:\left\|\phi-\phi_{0}\right\|_{1} \leq \epsilon\right\},\|\cdot\|_{P_{n}}\right) \leq\left(\frac{\epsilon^{2}}{u^{2}}+1\right) \log (1+k p)
$$

We can therefore conclude that

$$
\begin{aligned}
& H\left(3 \sqrt{d} M_{n} u,\left\{\left(L_{\vartheta}-L_{\vartheta_{0}}\right) 1\left\{G_{1} \leq M_{n}\right\}: \vartheta \in \tilde{\Theta}(\epsilon)\right\},\|\cdot\|_{P_{n}}\right) \\
& \quad \leq\left(\frac{\epsilon^{2}}{u^{2}}+m+1\right)\left(\log \left(\frac{5 \epsilon}{u}\right)+\log (1+k p)\right) .
\end{aligned}
$$

Let us now turn to the main proof of Lemma 2.
In what follows, $\left\{c_{t}\right\}$ are constants depending on $\Lambda_{\max }, k, m$ and $K$. The truncated version of the empirical process is

$$
\begin{aligned}
V_{n}^{\text {trunc }}(\vartheta)= & \frac{1}{n} \sum_{i=1}^{n}\left(L_{\vartheta}\left(x_{i}, Y_{i}\right) 1\left\{G_{1}\left(Y_{i}\right) \leq M_{n}\right\}\right. \\
& \left.-\mathbb{E}\left[L_{\vartheta}\left(x_{i}, Y\right) 1\left\{G_{1}(Y) \leq M_{n}\right\} \mid X=x_{i}\right]\right) .
\end{aligned}
$$

Let $\epsilon>0$ be arbitrary. We invoke Lemma 3.2 in van de Geer (2000), combined with a symmetrization lemma (e.g., a conditional version of Lemma 3.3 in van de Geer 2000). We apply these lemmas to the class

$$
\left\{\left(L_{\vartheta}-L_{\vartheta_{0}}\right) 1\left\{G_{1} \leq M_{n}\right\}: \vartheta \in \tilde{\Theta}(\epsilon)\right\} .
$$

In the notation used in Lemma 3.2 of van de Geer (2000), we take $\delta=c_{4} T \epsilon M_{n} \log n \sqrt{\log (p \vee n) / n}$, and $R=c_{5}(\epsilon \wedge 1) M_{n}$. This then gives

$$
\begin{aligned}
& \mathbb{P}_{\mathbf{x}}\left(\sup _{\vartheta \in \tilde{\Theta}(\epsilon)}\left|V_{n}^{\text {trunc }}(\vartheta)-V_{n}^{\text {trunc }}\left(\vartheta_{0}\right)\right| \geq c_{6} \epsilon T M_{n} \log n \sqrt{\frac{\log (p \vee n)}{n}}\right) \\
& \quad \leq c_{7} \exp \left[-\frac{T^{2} \log ^{2} n \log (p \vee n)\left(\epsilon^{2} \vee 1\right)}{c_{8}^{2}}\right] .
\end{aligned}
$$

Here, we use the bound (for $0<a \leq 1$ ),

$$
\int_{a}^{1} \frac{1}{u} \sqrt{\log \left(\frac{1}{u}\right)} d u \leq \log ^{3 / 2}\left(\frac{1}{a}\right)
$$

We then invoke the peeling device: split the set $\tilde{\Theta}$ into sets

$$
\left\{\vartheta \in \tilde{\Theta}: 2^{-(j+1)} \leq\left\|\phi-\phi_{0}\right\|_{1}+\left\|\eta-\eta_{0}\right\|_{2} \leq 2^{-j}\right\},
$$

where $j \in \mathbb{Z}$, and $2^{-j+1} \geq \lambda_{0}$. There are no more than $c_{9} \log n$ indices $j \leq 0$ with $2^{-j+1} \geq \lambda_{0}$. Hence, we get

$$
\sup _{\vartheta^{T}=\left(\phi^{T}, \eta^{T}\right) \in \tilde{\Theta}} \frac{\left|V_{n}^{\text {trunc }}(\vartheta)-V_{n}^{\text {trunc }}\left(\vartheta_{0}\right)\right|}{\left(\left\|\phi-\phi^{*}\right\|_{1}+\left\|\eta-\eta^{*}\right\|_{2}\right) \vee \lambda_{0}} \leq 2 c_{6} T M_{n} \log n \sqrt{\frac{\log (p \vee n)}{n}},
$$

with $\mathbb{P}_{\mathbf{x}}$ probability at least

$$
\begin{aligned}
1 & -c_{7}\left[c_{9} \log n\right] \exp \left[-\frac{T^{2} \log ^{2} n \log (p \vee n)}{c_{8}^{2}}\right] \\
& -\sum_{j=1}^{\infty} c_{7} \exp \left[-\frac{T^{2} 2^{2 j} \log ^{2} n \log (p \vee n)}{c_{8}^{2}}\right] \geq 1-c_{2} \exp \left[-\frac{T^{2} \log ^{2} n \log (p \vee n)}{c_{10}^{2}}\right] .
\end{aligned}
$$

Finally, to remove the truncation, we use

$$
\left|\left(L_{\vartheta}(x, y)-L_{\vartheta_{0}}(x, y)\right) 1\left\{G_{1}(y)>M_{n}\right\}\right| \leq d K G_{1}(y) 1\left\{G_{1}(y)>M_{n}\right\} .
$$

Hence

$$
\begin{aligned}
& \frac{\left|\left(V_{n}^{\text {trunc }}(\vartheta)-V_{n}^{\text {trunc }}\left(\vartheta_{0}\right)\right)-\left(V_{n}(\vartheta)-V_{n}\left(\vartheta_{0}\right)\right)\right|}{\left(\left\|\phi-\phi^{*}\right\|_{1}+\left\|\eta-\eta^{*}\right\|_{2}\right) \vee \lambda_{0}} \\
& \quad \leq \frac{d K}{n \lambda_{0}} \sum_{i=1}^{n}\left(G_{1}\left(Y_{i}\right) 1\left\{G_{1}\left(Y_{i}\right)>M_{n}\right\}+\mathbb{E}\left[G_{1}(Y) 1\left\{G_{1}(Y)>M_{n}\right\} \mid X=x_{i}\right]\right)
\end{aligned}
$$

## B. 3 Proof of Theorem 3

On $\mathcal{T}$,

$$
\overline{\mathcal{E}}\left(\hat{\psi} \mid \psi_{0}\right)+\lambda\|\hat{\phi}\|_{1} \leq T \lambda_{0}\left[\left(\left\|\hat{\phi}-\phi_{0}\right\|_{1}+\left\|\hat{\eta}-\eta_{0}\right\|_{2}\right) \vee \lambda_{0}\right]+\lambda\left\|\phi_{0}\right\|_{1}+\overline{\mathcal{E}}\left(\psi_{0} \mid \psi_{0}\right) .
$$

By Lemma 1,

$$
\overline{\mathcal{E}}\left(\hat{\psi} \mid \psi_{0}\right) \geq\left\|\hat{\psi}-\psi_{0}\right\|_{Q_{n}}^{2} / c_{0}^{2}
$$

and $\overline{\mathcal{E}}\left(\psi_{0} \mid \psi_{0}\right)=0$.
Case 1 Suppose that

$$
\left\|\hat{\phi}-\phi_{0}\right\|_{1}+\left\|\hat{\eta}-\eta_{0}\right\|_{2} \leq \lambda_{0} .
$$

Then we find

$$
\begin{aligned}
\overline{\mathcal{E}}\left(\hat{\psi} \mid \psi_{0}\right) & \leq T \lambda_{0}^{2}+\lambda\left\|\hat{\phi}-\phi_{0}\right\|_{1}+\overline{\mathcal{E}}\left(\psi_{0} \mid \psi_{0}\right) \\
& \leq\left(\lambda+T \lambda_{0}\right) \lambda_{0} .
\end{aligned}
$$

Case 2 Suppose that

$$
\left\|\hat{\phi}-\phi_{0}\right\|_{1}+\left\|\hat{\eta}-\eta_{0}\right\|_{2} \geq \lambda_{0}
$$

and that

$$
T \lambda_{0}\left\|\hat{\eta}-\eta_{0}\right\|_{2} \geq\left(\lambda+T \lambda_{0}\right)\left\|\hat{\phi}_{S}-\left(\phi_{0}\right)_{S}\right\|_{1} .
$$

Then we get

$$
\begin{aligned}
& \overline{\mathcal{E}}\left(\hat{\psi} \mid \psi_{0}\right)+\left(\lambda-T \lambda_{0}\right)\left\|\hat{\phi}_{S^{c}}\right\|_{1} \leq 2 T \lambda_{0}\left\|\hat{\eta}-\eta_{0}\right\|_{2} \\
& \quad \leq 4 T^{2} \lambda_{0}^{2} c_{0}^{2}+\left\|\hat{\eta}-\eta_{0}\right\|_{2}^{2} /\left(2 c_{0}^{2}\right) \\
& \quad \leq 4 T^{2} \lambda_{0}^{2} c_{0}^{2}+\overline{\mathcal{E}}\left(\hat{\psi} \mid \psi_{0}\right) / 2
\end{aligned}
$$

So then

$$
\overline{\mathcal{E}}\left(\hat{\psi} \mid \psi_{0}\right)+2\left(\lambda-T \lambda_{0}\right)\left\|\hat{\phi}_{S^{c}}\right\|_{1} \leq 8 T^{2} \lambda_{0}^{2} c_{0}^{2}
$$

Case 3 Suppose that

$$
\left\|\hat{\phi}-\phi_{0}\right\|_{1}+\left\|\hat{\eta}-\eta_{0}\right\|_{2} \geq \lambda_{0}
$$

and that

$$
T \lambda_{0}\left\|\hat{\eta}-\eta_{0}\right\|_{2} \leq\left(\lambda+T \lambda_{0}\right)\left\|\hat{\phi}_{S}-\left(\phi_{0}\right)_{S}\right\|_{1} .
$$

Then we have

$$
\overline{\mathcal{E}}\left(\hat{\psi} \mid \psi_{0}\right)+\left(\lambda-T \lambda_{0}\right)\left\|\hat{\phi}_{S^{c}}\right\|_{1} \leq 2\left(\lambda+T \lambda_{0}\right)\left\|\hat{\phi}_{S}-\phi_{0}\right\|_{1} .
$$

So then

$$
\left\|\hat{\phi}_{S^{c}}\right\|_{1} \leq 6\left\|\hat{\phi}_{S}-\left(\phi_{0}\right)_{S}\right\|_{1} .
$$

We can then apply the restricted eigenvalue condition to $\hat{\phi}-\phi_{0}$. This gives

$$
\begin{aligned}
\overline{\mathcal{E}}\left(\hat{\psi} \mid \psi_{0}\right)+\left(\lambda-T \lambda_{0}\right)\left\|\hat{\phi}_{S^{c}}\right\|_{1} & \leq 2\left(\lambda+T \lambda_{0}\right) \sqrt{s}\left\|\hat{\phi}_{S}-\phi_{0}\right\|_{2} \\
& \leq 2\left(\lambda+T \lambda_{0}\right) \sqrt{s} \kappa\left\|\hat{g}-g_{0}\right\|_{Q_{n}} \\
& \leq 4\left(\lambda+T \lambda_{0}\right)^{2} c_{0}^{2} \kappa^{2} s+\overline{\mathcal{E}}\left(\hat{\psi} \mid \psi_{0}\right) / 2
\end{aligned}
$$

So we arrive at

$$
\overline{\mathcal{E}}\left(\hat{\psi} \mid \psi_{0}\right)+2\left(\lambda-T \lambda_{0}\right)\left\|\hat{\phi}_{S^{c}}\right\|_{1} \leq 8\left(\lambda+T \lambda_{0}\right)^{2} c_{0}^{2} \kappa^{2} s
$$

## B. 4 Proof of Lemma 3

Let $Z$ be a standard normal random variable. Then by straightforward computations, for all $M>0$,

$$
\mathbb{E}[|Z| 1\{|Z|>M\}] \leq 2 \exp \left[-M^{2} / 2\right]
$$

and

$$
\mathbb{E}\left[|Z|^{2} 1\{|Z|>M\}\right] \leq(M+2) \exp \left[-M^{2} / 2\right]
$$

Thus, for $n$ independent copies $Z_{1}, \ldots, Z_{n}$ of $Z$, and $M=2 \sqrt{\log n}$,

$$
\begin{aligned}
& \mathbb{P}\left(\frac{1}{n} \sum_{i=1}^{n}\left|Z_{i}\right| 1\left\{\left|Z_{i}\right|>M\right\}>\frac{4 \log n}{n}\right) \\
& \quad \leq \mathbb{P}\left(\frac{1}{n} \sum_{i=1}^{n}\left|Z_{i}\right| 1\left\{\left|Z_{i}\right|>M\right\}-\mathbb{E}[|Z| 1\{|Z|>M\}]>\frac{2 \log n}{n}\right) \\
& \quad \leq \frac{n \mathbb{E}\left[|Z|^{2} 1\{|Z|>M\}\right]}{4(\log n)^{2}} \leq \frac{2}{n}
\end{aligned}
$$

The result follows from this, as

$$
G_{1}(Y)=\mathrm{e}^{K}|Y|+K,
$$

and $Y$ has a normal mixture distribution.

## B. 5 Proof of Theorem 5

On $\mathcal{T}$, defined in (5.18) with $\lambda_{0}=c_{4} \sqrt{\log ^{3} n \log (p \vee n) / n}\left(c_{4}\right.$ as in Lemma 3; i.e., $M_{n}=c_{4} \sqrt{\log (n)}$ in (5.20)), we have the basic inequality

$$
\overline{\mathcal{E}}\left(\hat{\psi} \mid \psi_{0}\right)+\lambda\|\hat{\phi}\|_{1} \leq T \lambda_{0}\left[\left(\left\|\hat{\phi}-\phi_{0}\right\|_{1}+\left\|\hat{\eta}-\eta_{0}\right\|_{2}\right) \vee \lambda_{0}\right]+\lambda\left\|\phi_{0}\right\|_{1}+\overline{\mathcal{E}}\left(\psi_{0} \mid \psi_{0}\right)
$$

Note that $\left\|\hat{\eta}-\eta_{0}\right\|_{2} \leq 2 K$ and $\overline{\mathcal{E}}\left(\psi_{0} \mid \psi_{0}\right)=0$. Hence, for $n$ sufficiently large,

$$
\begin{aligned}
\overline{\mathcal{E}}\left(\hat{\psi} \mid \psi_{0}\right)+\lambda\|\hat{\phi}\|_{1} & \leq T \lambda_{0}\left(\left\|\hat{\phi}-\phi_{0}\right\|_{1}+2 K\right)+\lambda\left\|\phi_{0}\right\|_{1}+\overline{\mathcal{E}}\left(\psi_{0} \mid \psi_{0}\right) \\
& \leq T \lambda_{0}\left(\|\hat{\phi}\|_{1}+\left\|\phi_{0}\right\|_{1}+2 K\right)+\lambda\left\|\phi_{0}\right\|_{1}+\overline{\mathcal{E}}\left(\psi_{0} \mid \psi_{0}\right)
\end{aligned}
$$

and therefore also

$$
\overline{\mathcal{E}}\left(\hat{\psi} \mid \psi_{0}\right)+\left(\lambda-T \lambda_{0}\right)\|\hat{\phi}\|_{1} \leq T \lambda_{0} 2 K+\left(\lambda+T \lambda_{0}\right)\left\|\phi_{0}\right\|_{1}+\overline{\mathcal{E}}\left(\psi_{0} \mid \psi_{0}\right)
$$

It holds that $\lambda \geq 2 T \lambda_{0}$ (since $\lambda=C \sqrt{\log ^{3} n \log (p \vee n) / n}$ for some $C>0$ sufficiently large), $\lambda_{0}=O\left(\sqrt{\log ^{3} n \log (p \vee n) / n}\right)$ and $\lambda=O\left(\sqrt{\log ^{3} n \log (p \vee n) / n}\right)$, and due to the assumption about $\left\|\phi_{0}\right\|_{1}$ we obtain on the set $\mathcal{T}$ that $\overline{\mathcal{E}}\left(\hat{\psi} \mid \psi_{0}\right) \rightarrow \overline{\mathcal{E}}\left(\psi_{0} \mid \psi_{0}\right)=$ $0(n \rightarrow \infty)$. Finally, the set $\mathcal{T}$ has large probability, as shown by Lemma 2 and using Proposition 3 and Lemma 3 for FMR models.

## Appendix C: Proofs for Sects. 3 and 6

## C. 1 Proof of Proposition 2

We restrict ourselves to a two class mixture with $k=2$. Consider the function $u(\xi)$ defined as

$$
\begin{align*}
u(\xi)= & \exp \left(\ell_{\mathrm{pen}}^{(0)}(\xi)\right) \\
= & \prod_{i=1}^{n}\left\{\left(\frac{\pi_{1}}{\sqrt{2 \pi} \sigma_{1}} \exp \left(\frac{-\left(Y_{i}-X_{i}^{T} \beta_{1}\right)^{2}}{2 \sigma_{1}^{2}}\right)+\frac{\left(1-\pi_{1}\right)}{\sqrt{2 \pi} \sigma_{2}}\right.\right. \\
& \left.\left.\times \exp \left(\frac{-\left(Y_{i}-X_{i}^{T} \beta_{2}\right)^{2}}{2 \sigma_{2}^{2}}\right)\right) \exp \left(\frac{-\lambda\left\|\beta_{1}\right\|_{1}}{\sigma_{1}}\right) \exp \left(\frac{-\lambda\left\|\beta_{2}\right\|_{1}}{\sigma_{2}}\right)\right\} . \tag{C.28}
\end{align*}
$$

We will show that $u(\xi)$ is bounded from above for $\xi=\left(\beta_{1}, \beta_{2}, \sigma_{1}, \sigma_{2}, \pi_{1}\right) \in \Xi=$ $\mathbb{R}^{2 p} \times \mathbb{R}_{>0}^{2} \times[0,1]$. Then, clearly, $-n^{-1} \ell_{\mathrm{pen}}^{(0)}(\theta)$ is bounded from below for $\theta=$ $\left(\phi_{1}, \phi_{2}, \rho_{1}, \rho_{2}, \pi_{1}\right) \in \Theta=\mathbb{R}^{2 p} \times \mathbb{R}_{>0}^{2} \times(0,1)$.

The critical point for unboundedness is if we choose for an arbitrary sample point $i \in\{1, \ldots, n\}$ a $\beta_{1}^{*}$ such that $Y_{i}-X_{i}^{T} \beta_{1}^{*}=0$ and let $\sigma_{1} \rightarrow 0$. Without the penalty term $\exp \left(\frac{-\lambda\left\|\beta_{1}^{*}\right\|_{1}}{\sigma_{1}}\right)$ in (C.28) the function would tend to infinity as $\sigma_{1} \rightarrow 0$. But as $Y_{i} \neq 0$ for all $i \in\{1, \ldots, n\}, \beta_{1}^{*}$ cannot be zero, and therefore $\exp \left(\frac{-\lambda\left\|\beta_{1}^{*}\right\|_{1}}{\sigma_{1}}\right)$ forces $u(\xi)$ to tend to 0 as $\sigma_{1} \rightarrow 0$.

Let us give a more formal proof for the boundedness of $u(\xi)$. Choose a small $0<\varepsilon_{1}<\min _{i}\left\{Y_{i}^{2}\right\}$ and $\varepsilon_{2}>0$. As $Y_{i} \neq 0, i=1, \ldots, n$, there exists a small constant $m>0$ such that

$$
\begin{equation*}
0<\min _{i}\left\{Y_{i}^{2}\right\}-\varepsilon_{1} \leq\left(Y_{i}-X_{i}^{T} \beta_{1}\right)^{2} \tag{C.29}
\end{equation*}
$$

holds for all $i=1, \ldots, n$ as long as $\left\|\beta_{1}\right\|_{1}<m$, and

$$
\begin{equation*}
0<\min _{i}\left\{Y_{i}^{2}\right\}-\varepsilon_{1} \leq\left(Y_{i}-X_{i}^{T} \beta_{2}\right)^{2} \tag{C.30}
\end{equation*}
$$

holds for all $i=1, \ldots, n$ as long as $\left\|\beta_{2}\right\|_{1}<m$.
Furthermore, there exists a small constant $\delta>0$ such that

$$
\begin{equation*}
\frac{1}{\sqrt{2 \pi} \sigma_{1}} \exp \left(-\left(\min _{i}\left\{Y_{i}^{2}\right\}-\varepsilon_{1}\right) / 2 \sigma_{1}^{2}\right)<\varepsilon_{2} \quad \text { and } \quad \frac{1}{\sqrt{2 \pi} \sigma_{1}} \exp \left(\frac{-\lambda m}{\sigma_{1}}\right)<\varepsilon_{2} \tag{C.31}
\end{equation*}
$$

hold for all $0<\sigma_{1}<\delta$, and

$$
\begin{equation*}
\frac{1}{\sqrt{2 \pi} \sigma_{2}} \exp \left(-\left(\min _{i}\left\{Y_{i}^{2}\right\}-\varepsilon_{1}\right) / 2 \sigma_{2}^{2}\right)<\varepsilon_{2} \quad \text { and } \quad \frac{1}{\sqrt{2 \pi} \sigma_{2}} \exp \left(\frac{-\lambda m}{\sigma_{2}}\right)<\varepsilon_{2} \tag{C.32}
\end{equation*}
$$

hold for all $0<\sigma_{2}<\delta$.
Define the set $K=\left\{\left(\beta_{1}, \beta_{2}, \sigma_{1}, \sigma_{2}, \pi_{1}\right) \in \Xi ; \delta \leq \sigma_{1}, \sigma_{2}\right\}$. Now $u(\xi)$ is trivially bounded on $K$. From the construction of $K$ and (C.29)-(C.32), we easily see that $u(\xi)$ is also bounded on $K^{c}$, and therefore bounded on $\Xi$.

## C. 2 Proof of Theorem 6

The density of the complete data is given by

$$
f_{c}(Y, \Delta \mid \theta)=\prod_{i=1}^{n} \prod_{r=1}^{k} \pi_{r}^{\Delta_{i, r}}\left(\frac{\rho_{r}}{\sqrt{2 \pi}} e^{-\frac{1}{2}\left(\rho_{r} Y_{i}-X_{i}^{T} \phi_{r}\right)^{2}}\right)^{\Delta_{i, r}}
$$

whereas the density of the observed data is given by

$$
\begin{aligned}
& f_{\mathrm{obs}}(Y \mid \theta)=\prod_{i=1}^{n} \sum_{r=1}^{k} \pi_{r} \frac{\rho_{r}}{\sqrt{2 \pi}} e^{-\frac{1}{2}\left(\rho_{r} Y_{i}-X_{i}^{T} \phi_{r}\right)^{2}}, \\
& \theta=\left(\rho_{1}, \ldots, \rho_{k}, \phi_{1,1}, \phi_{1,2}, \ldots, \phi_{k, p}, \pi\right) \in \Theta=\mathbb{R}_{>0}^{k} \times \mathbb{R}^{k p} \times \Pi \subset \mathbb{R}^{k(p+2)-1}, \\
& \Pi=\left\{\pi=\left(\pi_{1}, \ldots, \pi_{k-1}\right) ; \pi_{r}>0 \text { for } r=1, \ldots, k-1 \text { and } \sum_{r=1}^{k-1} \pi_{r}<1\right\}, \\
& \quad \pi_{k}=1-\sum_{r=1}^{k-1} \pi_{r} .
\end{aligned}
$$

Furthermore, the conditional density of the complete data given the observed data is given by $k(Y, \Delta \mid Y, \theta)=f_{c}(Y, \Delta \mid \theta) / f_{\text {obs }}(Y \mid \theta)$. Then, the penalized negative loglikelihood fulfills the equation

$$
\begin{align*}
\nu_{\text {pen }}(\theta) & =-n^{-1} \ell_{\mathrm{pen}, \lambda}^{(0)}(\theta) \\
& =-n^{-1} \log f_{\mathrm{obs}}(Y \mid \theta)+\lambda \sum_{r=1}^{k}\left\|\phi_{r}\right\|_{1}=Q_{\mathrm{pen}}\left(\theta \mid \theta^{\prime}\right)-H\left(\theta \mid \theta^{\prime}\right) \tag{C.33}
\end{align*}
$$

where

$$
Q_{\text {pen }}\left(\theta \mid \theta^{\prime}\right)=-n^{-1} \mathbb{E}_{\theta^{\prime}}\left[\log f_{c}(Y, \Delta \mid \theta) \mid Y\right]+\lambda \sum_{r=1}^{k}\left\|\phi_{r}\right\|_{1} \quad \text { (compare Sect. 6.1) }
$$

and

$$
H\left(\theta \mid \theta^{\prime}\right)=-n^{-1} \mathbb{E}_{\theta^{\prime}}[\log k(Y, \Delta \mid Y, \theta) \mid Y]
$$

By Jensen's inequality, we get the following important relationship

$$
\begin{equation*}
H\left(\theta \mid \theta^{\prime}\right) \geq H\left(\theta^{\prime} \mid \theta^{\prime}\right) \quad \forall \theta \in \Theta \tag{C.34}
\end{equation*}
$$

see also Wu (1983). $Q_{\text {pen }}\left(\theta \mid \theta^{\prime}\right)$ and $H\left(\theta \mid \theta^{\prime}\right)$ are continuous functions in $\theta$ and $\theta^{\prime}$. If we think of them as functions of $\theta$ with fixed $\theta^{\prime}$ we write also $Q_{\text {pen }, \theta^{\prime}}(\theta)$ and $H_{\theta^{\prime}}(\theta)$. Furthermore, $Q_{\text {pen, } \theta^{\prime}}(\theta)$ is a convex function of $\theta$ and strictly convex in each coordinate of $\theta$. As a last preparation, we give a definition of a stationary point for non-differentiable functions (see also Tseng 2001):

Definition 1 Let $u$ be a function defined on an open set $U \subset \mathbb{R}^{k(p+2)-1}$. A point $x \in U$ is called stationary if $u^{\prime}(x ; d)=\lim _{\alpha \downarrow 0} \frac{u(x+\alpha d)-u(x)}{\alpha} \geq 0 \forall d \in \mathbb{R}^{k(p+2)-1}$.

We are now ready to start with the proof which is inspired by Bertsekas (1995). We write $\theta=\left(\theta_{1}, \ldots, \theta_{D}\right)=\left(\rho_{1}, \ldots, \rho_{k}, \phi_{1,1}, \phi_{1,2}, \ldots, \phi_{k, p}, \pi\right)$ where $D=k+k p+1$ denotes the number of coordinates. Remark that the first $D-1$ coordinates are univariate, whereas $\theta_{D}=\pi$ is a "block coordinate" of dimension $k-1$.

Proof Let $\theta^{m}=\theta^{(m)}$ be the sequence generated by the BCD-GEM algorithm. We need to prove that for a converging subsequence $\theta^{m_{j}} \rightarrow \bar{\theta} \in \Theta, \bar{\theta}$ is a stationary point of $v_{\text {pen }}(\theta)$. Taking directional derivatives of (C.33) yields

$$
v_{\mathrm{pen}}^{\prime}(\bar{\theta} ; d)=Q_{\mathrm{pen}, \bar{\theta}}^{\prime}(\bar{\theta} ; d)-\left\langle\nabla H_{\bar{\theta}}(\bar{\theta}), d\right\rangle .
$$

Note that $\nabla H_{\bar{\theta}}(\bar{\theta})=0$ as $H_{\bar{\theta}}(x)$ is minimized for $x=\bar{\theta}$ (see (C.34)). Therefore, it remains to show that $Q_{\text {pen, }}^{\prime} \bar{\theta}(\bar{\theta} ; d) \geq 0$ for all directions $d$. Let

$$
z_{i}^{m}=\left(\theta_{1}^{m+1}, \ldots, \theta_{i}^{m+1}, \theta_{i+1}^{m}, \ldots, \theta_{D}^{m}\right)
$$

Using the definition of the algorithm, we have

$$
\begin{equation*}
Q_{\text {pen }, \theta^{m}}\left(\theta^{m}\right) \geq Q_{\text {pen }, \theta^{m}}\left(z_{1}^{m}\right) \geq \cdots \geq Q_{\text {pen }, \theta^{m}}\left(z_{D-1}^{m}\right) \geq Q_{\text {pen }, \theta^{m}}\left(\theta^{m+1}\right) \tag{C.35}
\end{equation*}
$$

Additionally, from the properties of GEM ((C.33) and (C.34)), we have

$$
\begin{equation*}
v_{\text {pen }}\left(\theta^{0}\right) \geq v_{\text {pen }}\left(\theta^{1}\right) \geq \cdots \geq v_{\text {pen }}\left(\theta^{m}\right) \geq v_{\text {pen }}\left(\theta^{m+1}\right) . \tag{C.36}
\end{equation*}
$$

Equation (C.36) and the converging subsequence imply that the sequence $\left\{v_{\text {pen }}\left(\theta^{m}\right)\right.$; $m=0,1,2, \ldots\}$ converges to $\nu_{\text {pen }}(\bar{\theta})$. Further, we have

$$
\begin{align*}
0 & \leq Q_{\text {pen }, \theta^{m}}\left(\theta^{m}\right)-Q_{\text {pen }, \theta^{m}}\left(\theta^{m+1}\right) \\
& =v_{\text {pen }}\left(\theta^{m}\right)-v_{\text {pen }}\left(\theta^{m+1}\right)+\underbrace{H_{\theta^{m}}\left(\theta^{m}\right)-H_{\theta^{m}}\left(\theta^{m+1}\right)}_{\leq 0} \\
& \leq \underbrace{v_{\text {pen }}\left(\theta^{m}\right)-v_{\text {pen }}\left(\theta^{m+1}\right)}_{\rightarrow v_{\text {pen }}(\bar{\theta})-v_{\text {pen }}(\bar{\theta})=0} . \tag{C.37}
\end{align*}
$$

We conclude that the sequence $\left\{Q_{\text {pen }, \theta^{m}}\left(\theta^{m}\right)-Q_{\text {pen }, \theta^{m}}\left(\theta^{m+1}\right) ; m=0,1,2, \ldots\right\}$ converges to zero.

We now show that $\left\{\theta_{m_{j}}^{m_{j}+1}-\theta_{1}^{m_{j}}\right\}$ converges to zero $(j \rightarrow \infty)$. Assume the contrary, in particular that $\left\{z_{1}^{m_{j}}-\theta^{m_{j}}\right\}$ does not converge to 0 . Let $\gamma^{m_{j}}=\left\|z_{1}^{m_{j}}-\theta^{m_{j}}\right\|$. Without loss of generality (by restricting to a subsequence), we may assume that there exists some $\bar{\gamma}>0$ such that $\gamma^{m_{j}}>\bar{\gamma}$ for all $j$. Let $s_{1}^{m_{j}}=\frac{z_{1}^{m_{j}}-\theta^{m_{j}}}{\gamma^{m_{j}}}$. This $s_{1}^{m_{j}}$ differs from zero only along the first component. As $s_{1}^{m_{j}}$ belongs to a compact set
$\left(\left\|s_{1}^{m_{j}}\right\|=1\right)$ we may assume that $s_{1}^{m_{j}}$ converges to $\bar{s}_{1}$. Let us fix some $\varepsilon \in[0,1]$. Notice that $0 \leq \varepsilon \bar{\gamma} \leq \gamma^{m_{j}}$. Therefore, $\theta^{m_{j}}+\varepsilon \bar{\gamma} s_{1}^{m_{j}}$ lies on the segment joining $\theta^{m_{j}}$ and $z_{1}^{m_{j}}$, and belongs to $\Theta$ because $\Theta$ is convex. As $Q_{\text {pen }, \theta^{m_{j}}(\cdot)}$ is convex and $z_{1}^{m_{j}}$ minimizes this function over all values that differ from $\theta^{m_{j}}$ along the first coordinate, we obtain

$$
\begin{align*}
Q_{\text {pen }, \theta^{m_{j}}}\left(z_{1}^{m_{j}}\right)=Q_{\text {pen }, \theta^{m_{j}}}\left(\theta^{m_{j}}+\gamma^{m_{j}} s_{1}^{m_{j}}\right) & \leq Q_{\text {pen }, \theta^{m_{j}}}\left(\theta^{m_{j}}+\varepsilon \bar{\gamma} s_{1}^{m_{j}}\right) \\
& \leq Q_{\text {pen }, \theta^{m_{j}}}\left(\theta^{m_{j}}\right) \tag{C.38}
\end{align*}
$$

From (C.35) and (C.38), we conclude

$$
\begin{aligned}
& 0 \leq Q_{\text {pen }, \theta^{m_{j}}}\left(\theta^{m_{j}}\right)-Q_{\text {pen }, \theta^{m_{j}}}\left(\theta^{m_{j}}+\varepsilon \bar{\gamma} s_{1}^{m_{j}}\right) \\
& \overbrace{\leq}^{(\mathrm{C} .38)} Q_{\text {pen, } \theta^{m_{j}}}\left(\theta^{m_{j}}\right)-Q_{\text {pen }, \theta^{m_{j}}}\left(z_{1}^{m_{j}}\right) \\
& \overbrace{\leq}^{(\mathrm{C} .35)} Q_{\text {pen }, \theta^{m_{j}}}\left(\theta^{m_{j}}\right)-Q_{\text {pen }, \theta^{m_{j}}}\left(\theta^{m_{j}+1}\right) .
\end{aligned}
$$

Using (C.37) and continuity of $Q_{\text {pen }, x}(y)$ in both arguments $x$ and $y$, we conclude by taking the limit $j \rightarrow \infty$ :

$$
Q_{\mathrm{pen}, \bar{\theta}}\left(\bar{\theta}+\varepsilon \bar{\gamma} \overline{s_{1}}\right)=Q_{\mathrm{pen}, \bar{\theta}}(\bar{\theta}) \quad \forall \varepsilon \in[0,1] .
$$

Since $\bar{\gamma} \bar{s}_{1} \neq 0$ this contradicts the strict convexity of $Q_{\text {pen, } \bar{\theta}}\left(x_{1}, \bar{\theta}_{2}, \ldots, \bar{\theta}_{D}\right)$ as a function of the first block-coordinate. This contradiction establishes that $z_{1}^{m_{j}}$ converges to $\bar{\theta}$.

From the definition of the algorithm, we have

$$
Q_{\text {pen }}\left(z_{1}^{m_{j}} \mid \theta^{m_{j}}\right) \leq Q_{\text {pen }}\left(x_{1}, \theta_{2}^{m_{j}}, \ldots, \theta_{D}^{m_{j}} \mid \theta^{m_{j}}\right) \quad \forall x_{1} .
$$

By continuity and taking the limit $j \rightarrow \infty$, we obtain

$$
Q_{\text {pen }, \bar{\theta}}(\bar{\theta}) \leq Q_{\text {pen }, \bar{\theta}}\left(x_{1}, \bar{\theta}_{2}, \ldots, \bar{\theta}_{D}\right) \quad \forall x_{1} .
$$

Repeating the argument, we conclude that $\bar{\theta}$ is a coordinate-wise minimum. Therefore, following Tseng (2001), $\bar{\theta}$ is easily seen to be a stationary point of $Q_{\text {pen }, \bar{\theta}}(\cdot)$, in particular $Q_{\text {pen }, \bar{\theta}}^{\prime}(\bar{\theta} ; d) \geq 0$ for all directions $d$.

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