# EFFECT OF HEAVY TAILS ON ULTRA HIGH DIMENSIONAL VARIABLE RANKING METHODS 

Aurore Delaigle and Peter Hall<br>University of Melbourne


#### Abstract

Contemporary problems involving sparse, high-dimensional feature selection are becoming rapidly more challenging through substantial increases in dimension. This places ever more stress on methods for analysis, since the effects of even moderately heavy-tailed feature distributions become more significant as the number of features diverges. Data transformations have a significant role to play, reducing noise and enabling an increase in dimension, and for this reason they are increasingly used. In this paper we examine the performance of a typical transformation of this type, and study the extent to which it preserves the main attributes that lead to reliable feature selection. We show both numerically and theoretically that, in the presence of heavy-tailed data, the size of the dimension for which effective variable selection is possible can be increased dramatically, from a low-degree polynomial function of sample size to one that is exponentially large.


Key words and phrases: Correlation, feature selection, heavy tail, nonparametric statistics, Studentising, variable selection.

## 1. Introduction

The future of genomic data analysis promises data vectors whose length, $p$, is not in the thousands or tens of thousands, as is common today, but in the millions or tens of millions. To appreciate the relevance of these numbers, note that there are likely between $p=20,000$ and 25,000 human protein-coding genes, but, apparently, between $p=10^{6}$ and $3 \times 10^{6} \mathrm{SNPs}$, or single-nucleotide polymorphisms, in the human genome. Data on SNPs, rather than on genes, are becoming common, and the availability of such high dimensional vectors means that methods have to be effective when $p$ is much larger than the sample size, $n$.

In this paper we consider feature ranking in such very high-dimensional problems. It has been proved in that context that if distributions are sufficiently light-tailed, standard methods work for $p$ exponentially large as a function of $n$. However, genetic data are often very noisy, and can contain unusually large observations. In other words, such data are often not sufficiently light-tailed, which is a serious issue in very high-dimensional settings because, as the number
of features diverges, the influence of even moderately heavy tailed feature distributions becomes more significant. Therefore, techniques more appropriate to heavy-tailed distributions have to be employed.

In this setting it is increasingly common to borrow tools from the more conventional low-dimensional robust literature. In particular, practitioners commonly employ simple data transformation methods. In the present paper we provide a theoretical account of the performance of such transformation methods, pointing especially to the extent of their advantages in ultra high-dimensional settings. We show that they can provide substantial improvements when some or all of the components have heavy-tailed distributions. For example, we show that standard feature ranking methods, whose performance degrades when $p$ is an exponentially large function of $n$, recover their capacity to deal with ultra high-dimensional, noisy data when used in conjunction with an appropriate data transformation. The theoretical results in Section 3 demonstrate these advantages for ranking methods based on either differences of means or on correlation. In the former setting, a standard approach consists of ranking components according to Student's $t$ scores calculated for each component. Although Studentising helps combat heavy-tailed noise, we show that transforming the variables has even more to offer. Moreover, when the distributions are light-tailed and do not contain outliers, transforming the data does not have a serious negative impact on feature ranking. In Section 4 we illustrate some of the main issues using two real datasets.

## 2. Transformation Methods

### 2.1. Feature ranking based on correlation

Suppose we observe independent and identically distributed data pairs ( $X_{i}$, $\left.Y_{i}\right)$, for $1 \leq i \leq n$, where $X_{i}=\left(X_{i 1}, \ldots, X_{i p}\right)$ is a $p$-vector, $Y_{i}$ is a scalar, and we are interested in the relation between $Y_{i}$ and $X_{i}$. For example, $Y_{i}$ might be a score variable for a disease, and $X_{i}$ a vector whose indices represent $p$ genes or parts of a chromosome, and the interest could be in uncovering indices that are relevant to the disease, and in constructing a model for the relationship between these and $Y_{i}$. Methods for model building include the lasso (Tibshiranil ([996)), non-convex penalisation (Fan and Lil (2007)) and the Dantzig selector (Candes and Tan (2007)).

Fan and Lv (20108) suggest preceding model building by a massive dimension reduction step, using "correlation learning" where the $p$ components of $X_{i}$ are ranked according to the values of $\left|\hat{\rho}_{j}\right|$, with, for $j=1, \ldots, p, \hat{\rho}_{j}=\hat{\xi}_{j} /\left(\hat{\sigma}_{X j} \hat{\sigma}_{Y}\right)$ denoting an empirical estimator of the theoretical correlation $\rho_{j}=\operatorname{corr}\left(X_{1 j}, Y_{1}\right)=$
$\xi_{j} /\left\{\sigma_{X j}^{2} \sigma_{Y}^{2}\right\}^{1 / 2}$. Here we have used the notation $\xi_{j}=\operatorname{cov}\left(X_{1 j}, Y_{1}\right), \sigma_{X j}^{2}=$ $\operatorname{var}\left(X_{1 j}\right), \sigma_{Y}^{2}=\operatorname{var}\left(Y_{1}\right)$,
$\hat{\xi}_{j}=\frac{1}{n} \sum_{i=1}^{n}\left(X_{i j} Y_{i}-\bar{X}_{j} \bar{Y}\right), \hat{\sigma}_{X j}^{2}=\frac{1}{n} \sum_{i=1}^{n}\left(X_{i j}^{2}-\bar{X}_{j}^{2}\right), \hat{\sigma}_{Y}^{2}=\frac{1}{n} \sum_{i=1}^{n}\left(Y_{i}^{2}-\bar{Y}^{2}\right)$,
where $\bar{X}_{j}=n^{-1} \sum_{i} X_{i j}$ and $\bar{Y}=n^{-1} \sum_{i} Y_{i}$. Correlation learning assigns empirical rank $\hat{\jmath}_{k}$ to the $j$ th component, where $\hat{\jmath}_{1}, \ldots, \hat{\jmath}_{p}$ denotes the permutation of $1, \ldots, p$ such that $\left|\hat{\rho}_{\hat{\jmath}_{1}}\right| \geq \cdots \geq\left|\hat{\rho}_{\hat{\jmath}_{p}}\right|$. For each $k, \hat{\jmath}_{k}$ is an estimator of the theoretical rank $j_{k}$, where $j_{1}, \ldots, j_{p}$ is a permutation of $1, \ldots, p$ such that $\left|\rho_{j_{1}}\right| \geq \cdots \geq\left|\rho_{j_{p}}\right|$. Dimension reduction is performed by removing components with low absolute empirical correlations. See also Lv and Fan (2009) and Fan and L.v (2010).

In some analyses of genomic data, practitioners report the results of robust Wilcoxon rank tests as well as, or instead of, Student's $t$ statistics; see Li and Fine (2010). This approach is relatively robust, and immune to heavy-tail problems, although the results to which it leads are perhaps a little more difficult to interpret since they address differences between the distributions of noisy approximations to (for example) gene expression levels, rather than differences between means.

### 2.2. Data transformation for correlation ranking

A difficulty with the correlation learning methodology discussed in Section 2.1 is that empirical correlations are sensitive to aberrations caused by heavytailed distributions of the explanatory vectors $X_{i}$ or experimental errors $\epsilon_{i}$. As a result, ranking based on correlations can be quite poor in the presence of heavy tails, since empirical correlations fluctuate so heavily that many of those corresponding to a theoretical correlation of zero take higher values than those for which the theoretical counterpart is nonzero. To reduce the fluctuations of empirical correlations, a common approach is to transform the data, replacing $X_{i j}$ and $Y_{i}$ by

$$
\begin{equation*}
U_{i j}=\Psi\left(X_{i j}\right), \quad Z_{i}=\Psi\left(Y_{i}\right), \tag{2.2}
\end{equation*}
$$

respectively, when computing the correlation coefficient. Here we take $\Psi$ to be a uniformly bounded, monotone function, for example a distribution function. Then we replace $\rho_{j}$ by the new correlation coefficient computed from the transformed data:

$$
\begin{equation*}
\omega_{j}=\frac{\operatorname{cov}\left(U_{1 j}, Z_{1}\right)}{\left\{\operatorname{var}\left(U_{1 j}\right) \operatorname{var}\left(Z_{1}\right)\right\}^{1 / 2}}, \tag{2.3}
\end{equation*}
$$

of which an estimator is $\widehat{\omega}_{j}=\hat{\zeta}_{j} /\left(\hat{\tau}_{U j} \hat{\tau}_{Z}\right)$ where

$$
\begin{equation*}
\hat{\zeta}_{j}=\frac{1}{n} \sum_{i=1}^{n}\left(U_{i j} Z_{i}-\bar{U}_{j} \bar{Z}\right), \hat{\tau}_{U j}^{2}=\frac{1}{n} \sum_{i=1}^{n}\left(U_{i j}^{2}-\bar{U}_{j}^{2}\right), \hat{\tau}_{Z}^{2}=\frac{1}{n} \sum_{i=1}^{n}\left(Z_{i}^{2}-\bar{Z}^{2}\right) \tag{2.4}
\end{equation*}
$$

The components are ranked in the order $\hat{\jmath}_{1}, \ldots, \hat{\jmath}_{p}$ determined by $\left|\widehat{\omega}_{\hat{\jmath}_{1}}\right| \geq \cdots \geq$ $\left|\widehat{\omega}_{\hat{\jmath}_{p}}\right|$.

The estimator $\widehat{\omega}_{j}$ is effectively the correlation estimator discussed in Section 8.3 of Huber ( 1981 ). Khan, Van Aelst, and Zamar (2007) discuss a similar approach for calculating correlation in the context of the LARS method of Efron et al. (2004), a technique of the same type as the lasso, where one of the ingredients needed is correlation. However, they treat only numerical aspects, and only in cases where $p<n$.

Our goal here is quite different. We wish to prove, both theoretically and numerically, that in ultra high-dimensional problems, where $p$ is much larger than $n$, standard correlations fail to rank components correctly, whereas transformed correlations perform well. Note that it is really the conjunction of heavy-tails and high dimension that causes poor performance of correlation ranking, since in the case of heavy-tailed distributions, when $p$ is large, with high probability many of the $p$ empirical correlations are very far from their true values. Section 3 will shed light on the impact of heavy tails on ranking a very high number of components through empirical correlations, and prove that transforming data can greatly improve ranking.

To illustrate that it is ultra high dimension that, when associated with heavy tails, causes the difficulties, we simulated data $\left(X_{i 1}, \ldots, X_{i p}, Y_{i}\right), i=1, \ldots, n$, from the model $Y_{i}=\sum_{j=1}^{6} X_{i j}+\epsilon_{i}$, where distributions of the $\epsilon_{i}$ 's and of the $X_{i}$ 's were, respectively, $0.98 F_{1}+0.02 F_{2}$ and $0.98 F_{1}+0.02 F_{3}$, with $F_{1}, F_{2}$ and $F_{3}$ denoting the distributions of, respectively, a $U[-10,10]$, a $U[-150,-100]$, and a $U[15,25]$. Here, the $Y_{i}$ 's depend only on $X_{i 1}, \ldots, X_{i 6}$, and a small proportion of the data take unusually large values. For several values of $n$ and $p$ we generated 200 samples of size $n$ from this model, and then ranked the $p$ components $X_{i 1}, \ldots, X_{i p}$ according to their absolute empirical correlations with $Y_{i}$, as described in Section 2.1, and according to their transformed versions described above, with $\Psi$ described in Section 4 . From the 200 samples we then calculated the median (Q2) and first and third quartiles (Q1 and Q3) of ranks assigned to $X_{i 1}$ to $X_{i 6}$ (their true ranks are 1). We show the results in Table 1 for several values of $n$, when $p=20,000$.

As can be seen, it is when $n$ is too small that occasional outliers cause empirical correlation ranking to fail to identify $X_{i 1}, \ldots, X_{i 6}$ as being some of the most important components. When $n$ is larger, the method does not rank

Table 1. Median and quartile ranks of $X_{i 1}$ to $X_{i 6}$ obtained by empirical correlation ranking (first block of four rows) and by transformed empirical correlation ranking (second block of four rows), when ranking $p=20,000$ components.

|  | $X_{1}$ | $X_{2}$ |  | $X_{3}$ |  | $X_{4}$ |  | $X_{5}$ |  | $X_{6}$ |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $n$ | Q2 Q1-Q3 | Q2 | Q1-Q3 | Q2 | Q1-Q3 | Q2 | Q1-Q3 | Q2 | Q1-Q3 |  | Q1-Q3 |
| 50 | 712 75-3291 | 809 | 168-3600 | 474 | 45-2707 | 906 | 172-4786 | 672 | 120-3531 | 750 | 86-3638 |
| 75 | 236 32-1880 | 192 | 36-1394 | 282 | 30-1750 | 236 | 24-2038 | 190 | 17-1632 | 331 | 37-1951 |
| 100 | 91 7-806 | 152 | 15-854 | 110 | 10-836 | 88 | 9-795 | 145 | 14-776 | 98 | 9-575 |
| 200 | 6 3-49 | 7 | 3-56 | 7 | 2-31 | 10 | 3-59 | 6 | 3-33 | 6 | 2-56 |
| 50 | 174 28-996 | 317 | 45-1394 | 137 | 14-920 | 344 | 57-1429 | 262 | 42-1116 | 242 | 29-1570 |
| 75 | 48 5-202 | 39 | 5-199 | 40 | 5-244 | 68 | 7-388 | 25 | 5-175 | 44 | 5-288 |
| 100 | 8 3-56 | 15 | 4-76 | 12 | 3-71 | 8 | 2-37 | 10 | 3-93 | 7 | 2-68 |
| 200 | 4 2-6 | 4 | 2-5 | 4 | 2-5 | 3 | 2-6 | 3 | 2-5 | 3 | 2-5 |

them perfectly, but is able to rank them highly (a rank is high when close to 1 ). Clearly, the transformed correlation approach is able to rank highly $X_{i 1}, \ldots, X_{i 6}$ for smaller values of $n$. In this simple example the observations contain only moderate outliers, but the differences between the two methods are even more striking for higher values of $p$ and when the distributions are more heavily tailed; see Delaigle and Hall (2017) for a more extensive simulation study. Section 3 will give a theoretical account of these empirical findings.

Note that for simplicity we focus on the simple correlation ranking method but, of course, variable transformation can be used to reduce the effect of heavytailed distributions in variants of the correlation approach, such as the generalised correlation technique of Hall and Miller (2000). Conclusions similar to those drawn in this paper can be derived in such settings too.

### 2.3. Feature ranking based on mean differences

Transforming data prior to ranking can also be used in other contexts, such as the one based on mean differences. Suppose we observe a sample of size $n$ of independent $p$-vectors from populations $\Pi_{X}$ and $\Pi_{Y}$, say $X_{1}, \ldots, X_{n_{1}}$ from $\Pi_{X}$, and $Y_{1}, \ldots, Y_{n_{2}}$ from $\Pi_{Y}$, where $n_{1}+n_{2}=n, X_{i}=\left(X_{i 1}, \ldots, X_{i p}\right)$ and $Y_{i}=\left(Y_{i 1}, \ldots, Y_{i p}\right)$. For example, $X_{i j}$ and $Y_{i j}$ could represent expression levels for the $j$ th gene of the $i$ th individual coming from population $\Pi_{X}$ (a population of individuals suffering from a particular medical condition) and $\Pi_{Y}$ (a population from which the condition is absent), respectively. It is often of interest in such problems to identify the components of the $p$-vectors for which the two populations differ the most. This is often done by detecting the components that
have the largest mean differences. Let

$$
\begin{equation*}
E\left(X_{i_{1} j}\right)=\mu_{1 j} \text { and } E\left(Y_{i_{2} j}\right)=\mu_{2 j}, \quad \text { for } 1 \leq i_{k} \leq n_{k}, k=1,2 \text { and } 1 \leq j \leq p, \tag{2.5}
\end{equation*}
$$

and let $j_{1}, \ldots, j_{p}$ be a permutation of $1, \ldots, p$ such that $\mu_{2 j_{1}}-\mu_{1 j_{1}} \geq \cdots \geq$ $\mu_{2 j_{p}}-\mu_{1 j_{p}}$. When ranking is based on mean differences, the $\ell$ th component is assigned rank $j_{\ell}$.

In practice the means $\mu_{k j}$ are unknown and the ranks are estimated empirically from the data, for example by replacing the theoretical means by their empirical counterparts, and ranking the components according to the values of the raw differences,

$$
\begin{equation*}
D_{j}=\bar{Y}_{j}-\bar{X}_{j}, \tag{2.6}
\end{equation*}
$$

where $\bar{X}_{j}=n_{1}^{-1} \sum_{i} X_{i j}, \bar{Y}_{j}=n_{2}^{-1} \sum_{i} Y_{i j}$.
Like the correlation approach, this ranking method is very sensitive to random fluctuations. To some extent this sensitivity can be corrected by ranking values taken by Student's $t$ statistics $T_{j}$, corresponding to two-sample $t$-tests of the null hypothesis $H_{0 j}$ that $\mu_{2 j}=\mu_{1 j}$ against the alternative $H_{1 j}$ that $\mu_{2 j}>\mu_{1 j}$; that is, by ranking

$$
\begin{equation*}
T_{j}=\frac{\bar{Y}_{j}-\bar{X}_{j}}{\left(n_{1}^{-1} S_{1 j}^{2}+n_{2}^{-1} S_{2 j}^{2}\right)^{1 / 2}}, \tag{2.7}
\end{equation*}
$$

for $1 \leq j \leq p$, where $S_{1 j}^{2}$ and $S_{2 j}^{2}$ are conventional variance estimators computed from the data $X_{1 j}, \ldots, X_{n j}$ and $Y_{1 j}, \ldots, Y_{n j}$, respectively. Based on these $t$ statistics, the components are ranked empirically by taking $\hat{\jmath}_{1}, \ldots, \hat{\jmath}_{p}$ to be the permutation of $j_{1}, \ldots, j_{p}$ such that $T_{\hat{\jmath}_{1}} \geq \cdots \geq T_{\hat{\jmath}_{p}}$. Adopting this technique the most influential feature, for example in a genomic context, the most influential gene, is the one for which $H_{0 j}$ is rejected most resoundingly in favour of $H_{1 j}$. See also Lyons-Weiler, Patel, and Bhattacharya (20103), Xie et all (2004), Papana and Ishwaran (2006) and Yang et al. (2009) for other motivations of the Student's $t$ approach.

Student's $t$ statistics $T_{j}$ are less influenced by statistical fluctuations than the raw differences $D_{j}$ when the data $X_{i j}$ or $Y_{i j}$ are relatively heavy tailed; see for example Delaigle, Hall, and Jin (201). However they can still be affected significantly and, in consequence, $\hat{\jmath}_{1}, \ldots, \hat{\jmath}_{p}$ as above can be poor estimators of the true ranks $j_{1}, \ldots, j_{p}$. Moreover, this approach runs counter to the notion that it should emphasise location differences. To see why, note that if $j=j_{\text {small }}$ and $j=j_{\text {large }}$ denote indices for which $\sigma_{j}^{2}=n_{2}^{-1} \operatorname{var}\left(Y_{1 j}\right)+n_{1}^{-1} \operatorname{var}\left(X_{1 j}\right)$ is particularly small, or particularly large, respectively, then the feature with index $j_{\text {small }}$ is likely to be ranked well ahead of that with index $j_{\text {large }}$ even if the mean difference $\mu_{1 j}-\mu_{2 j}$ takes much the same value in both cases. Thus, a ranking of
the $T_{j} \mathrm{~s}$ generally fails to adequately reflect the sizes of the differences $\mu_{1 j}-\mu_{2 j}$; in important cases those values are distorted by standardising for scale. This focus on differences between population means reflects the fact that, in most applications of $t$ statistics, one generally has in mind not just the difference between two sample means, but rather the implications, conveyed by that difference, for the difference between the means of the respective populations.

Variants of $t$ statistics have been suggested in the literature. See for example Efron et al. (2001), Smyth (2004), Wul (2005) and Opgen-Rhein and Strimmer (2007). However, the main goal of these methods is to improve estimators of the variances $\sigma_{j}^{2}$ by borrowing strengths across similar components. In particular, they do not necessarily address the problems that concern us. For an illustration on some simulated examples, see Delaigle and Hall (20II).

The discussion motivates methodology that focuses more sharply on the magnitudes of differences between means, yet also alleviates problems that arise when the data $X_{i j}$, and/or $Y_{i j}$, have heavy tails. Let $U_{i j}=\Psi\left(X_{i j}\right)$ and $V_{i j}=$ $\Psi\left(Y_{i j}\right)$, where $\Psi$ is a uniformly bounded, monotone increasing function. Here, the components are ranked by ranking values of empirical mean differences calculated from these transformed data. That is, with the transformation approach,

$$
\begin{equation*}
\hat{\jmath}_{1}, \ldots, \hat{\jmath}_{p} \text { is the permutation of } 1, \ldots, p \text { for which } \bar{V}_{\hat{\jmath}_{1}}-\bar{U}_{\hat{\jmath}_{1}} \geq \cdots \geq \bar{V}_{\hat{\jmath}_{p}}-\bar{U}_{\hat{\jmath}_{p}}, \tag{2.8}
\end{equation*}
$$

where $\bar{U}_{j}=n_{1}^{-1} \sum_{i} U_{i j}, \bar{V}_{j}=n_{2}^{-1} \sum_{i} V_{i j}$. To appreciate that this approach is well founded, let us elaborate on the model (2.5) by assuming that
$X_{i j}=\mu_{1 j}+\epsilon_{1 i j}, Y_{i j}=\mu_{2 j}+\epsilon_{2 i j}$ and, for any given $i$ and $j$, the
errors $\epsilon_{k i j}$ for $k=1,2$ are identically distributed.
(Note that the quantities $\mu_{k j}$ here need not be means and, in particular, in contrast to cases where $t$ statistics are ranked, no finite moments need be assumed.) Then $E\left\{\Psi\left(Y_{i j}\right)\right\} \geq E\left\{\Psi\left(X_{i j}\right)\right\}$ or $E\left\{\Psi\left(Y_{i j}\right)\right\} \leq E\left\{\Psi\left(X_{i j}\right)\right\}$, according as $\mu_{2 j} \geq \mu_{1 j}$, or $\mu_{2 j} \leq \mu_{1 j}$, respectively. These properties persist if the inequalities are taken to be strict, provided that $\Psi$ is strictly monotone increasing on the real line. The assumption of identical distribution of errors, imposed in (2.Y), is not a necessary condition for preserving the order of expectation in each case, for example when the means take only a finite number of distinct values, but it is perhaps the simplest sufficient condition.

Therefore, in the circumstances described by ( $\mathrm{E} . \mathrm{T}$ ),
the expected value of $\bar{V}_{j}-\bar{U}_{j}$ is of the same sign as $\mu_{2 j}-\mu_{1 j}$, and the expected value is a monotone increasing function of $\mu_{2 j}$ and a monotone decreasing function of $\mu_{1 j}$.

Moreover, since the variables $U_{i j}$ and $V_{i j}$ are uniformly bounded, their associated
large deviation probabilities are particularly small, even when the errors $\epsilon_{k i j}$ in (2.9) are heavy-tailed. This property and (2.10) underpin the attractiveness of feature ranking based on values of $\bar{V}_{j}-\bar{U}_{j}$. Note that, since the transformation
 defined even if the data are very heavy-tailed.

Transforming data before calculating a mean is not new, but our goal is to investigate the advantages of such a transformation approach when used to rank a very high number of components in the presence of heavy tails. As in the correlation context, the negative impact of heavy tails for ranking means, without transforming the data, makes itself felt mostly when $p$ is extremely large, since the probability of incorrect ranking increases with $p$. See Section 3 for a theoretical study.

To numerically illustrate the problems due to the conjunction of high dimension and heavy tails, we generated 200 samples from ( $X_{i 1}, \ldots, X_{i p}$ ) and $\left(Y_{i 1}, \ldots, Y_{i p}\right)$, where $X_{i j}-\mu_{j} \sim F$ and $Y_{i j} \sim F$, with $F$ a symmetric stable distribution with characteristic function $\phi(t)=\exp \left(-|t|^{1.5}\right)$. Here only the first six components are relevant, as we took $\mu_{j}=1\{j \leq 6\}$. We ranked all $p$ components using the method based on the raw differences at ([2.6), the one based on Student's $t$ differences at ([2.7), and the transformation approach discussed above with $\Psi$ as in Section 4. From the 200 samples we then calculated the median, first, and third quartiles of ranks assigned to $X_{i 1}$ to $X_{i 6}$ (their true ranks are 1). We show the results in Table 2 for the three methods, when $p=20,000$ and for several values of $n$.

As in the correlation context we see that it is when $n$ is very small compared to $p$ that the method based on (2.6) fails to identify these six components as being important, but the method does rank them increasingly highly as $n$ increases. Ranking from Student's $t$ differences at ([2.7) already improved the results, but the transformation method gave even better results. See Delaigle and Hall (2011) for a more extensive simulation study with more complex settings, and see Section 3 for a theoretical study of the empirical properties discussed here.

Of course, transforming the variables does not necessarily imply that all the components will be better ranked than by Student's $t$. However, the transformation approach often manages to detect components that were neglected because of heavy tailedness, and likewise can avoid false positive results caused by the same issue. Our main goal is not to rank all the components perfectly, but to identify the most important ones, i.e. those with highest ranks.

Remark 1. The expected value in (2.10) is not necessarily a monotone function of $\mu_{2 j}-\mu_{1 j}$. This need not cause difficulties, because in many settings the levels of monotonicity described in (ㄹ.IC) are adequate. However, the problem

Table 2. Median and quartile ranks of $X_{i 1}$ to $X_{i 6}$ obtained by empirical mean ranking, when ranking $p=20,000$ components. First block of four lines: using (2.6); second block: using (L2.7); third block: using $\bar{V}_{j}-\bar{U}_{j}$ at (2.8).

|  | $X_{1}$ |  | $X_{2}$ |  | $X_{3}$ |  | $X_{4}$ |  | $X_{5}$ |  | $X_{6}$ |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| $n$ | Q2 | Q1-Q3 | Q2 | Q1-Q3 | Q2 | Q1-Q3 | Q2 | Q1-Q3 | Q2 | Q1-Q3 | Q2 | Q1-Q3 |
| 50 | 1696 | 941-4085 | 1664 | 819-3140 | 1896 | 951-4672 | 1856 | 1042-3596 | 1448 | 793-3139 | 1692 | 908-3343 |
| 75 | 1442 | 698-3052 | 1378 | 786-2275 | 1382 | 796-3307 | 1344 | 717-2753 | 1318 | 663-2483 | 1492 | 736-3251 |
| 100 | 1169 | 694-2126 | 1099 | 676-1963 | 934 | 602-1760 | 1042 | 599-1739 | 1064 | 722-2514 | 1128 | 689-2377 |
| 200 | 684 | 446-1180 | 730 | 468-1254 | 694 | 471-1341 | 702 | 472-1386 | 791 | 506-1230 | 718 | 439-1132 |
| 50 | 410 | 50-4017 | 310 | 20-2420 | 627 | 53-4492 | 562 | 62-2595 | 478 | 39-2326 | 594 | 102-3389 |
| 75 | 243 | 7-3089 | 151 | 21-1249 | 380 | 18-2827 | 254 | 16-1890 | 292 | 14-2428 | 340 | 20-2818 |
| 100 | 167 | 9-1305 | 87 | 6-1168 | 81 | 6-892 | 85 | 5-788 | 108 | 9-1670 | 126 | 10-1670 |
| 200 | 11 | 2-515 | 17 | 3-639 | 11 | 2-280 | 7 | 2-304 | 25 | 3-615 | 14 | 3-315 |
| 50 | 28 | 5-133 | 36 | 7-157 | 48 | 6-240 | 39 | 6-165 | 37 | 4-188 | 35 | 7-240 |
| 75 | 6 | 2-40 | 7 | 3-36 | 6 | 2-27 | 6 | 2-38 | 6 | 2-41 | 10 | 2-43 |
| 100 | 5 | 2-13 | 4 | 2-9 | 4 | 2-6 | 4 | 2-6 | 5 | 2-12 | 4 | 2-8 |
| 200 | 3 | $2-5$ | 3 | $2-5$ | 4 | 2-5 | 3 | 2-5 | 4 | 2-5 | 3 | 2-5 |

can be eradicated by adjusting the methodology slightly. Specifically, define $W_{i_{0} i_{1} j}=\Psi\left(Y_{i_{1} j}-X_{i_{0} j}\right)$ for $1 \leq i_{0} \leq n_{1}$ and $1 \leq i_{1} \leq n_{2}$, and put

$$
\bar{W}_{j}=\frac{1}{n_{1} n_{2}} \sum_{i_{0}=1}^{n_{1}} \sum_{i_{1}=1}^{n_{2}} W_{i_{0} i_{1} j}
$$

Then, if the errors $\epsilon_{k i j}$ in ([.TM) have distributions that do not depend on $i, E\left(\bar{W}_{j}\right)$ is a monotone increasing function of $\mu_{2 j}-\mu_{1 j}$ and is strictly monotone if $\Psi$ is strictly increasing. This property motivates the following alternative method:

$$
\begin{equation*}
\text { define } \hat{\jmath}_{1}, \ldots, \hat{\jmath}_{p} \text { to be the permutation of } 1, \ldots, p \text { for which } \bar{W}_{\hat{\jmath}_{1}} \geq \cdots \geq \bar{W}_{\hat{\jmath}_{p}} \tag{2.11}
\end{equation*}
$$

We prefer the method based on (2.8), since it is considerably faster to implement in practice.

Remark 2. The transformation we use is monotone, although nonlinear, and this is a critical asset when using it for the purpose of ranking. Of course, the scale and units of the data change in a nonlinear way when transformed, and that should be borne in mind if we use the transformed data for a purpose other than ranking. This issue is not as important in the case of ranking based the correlation coefficient, which is scale and unit free and is itself a somewhat arbitrary measure of association.

## 3. Theoretical Properties

In this section we study theoretical properties of the transformed feature ranking approaches discussed in Section 2. We start with a general result in

Sections 3.1 and 3.2, where we prove that overall, without any restriction on the tails of the data, the empirical rankings based on variable transformation accord with their theoretical counterpart even when $p$ is much larger than $n$. Next, in Section 3.3 we show that this nice property is not shared by standard methods based on untransformed data. Together, these results illustrate the extent to which transforming the data effectively addresses heavy-tailedness in ultra-high dimensional feature selection. We first state, and prove, our theoretical results in the case of correlation; this is more awkward than the context of means since correlation involves a random denominator. Subsequently, we consider the case of the mean.

### 3.1. Performance of the transformation method based on correlation

Suppose we observe independent and identically distributed vectors ( $X_{i}, Y_{i}$ ), for $1 \leq i \leq n$, where $X_{i}=\left(X_{i 1}, \ldots, X_{i p}\right)$ is a $p$-vector and $Y_{i}$ is a scalar. Note that we do not make any assumptions on the dependence structure of the components of the feature vectors $X_{i}=\left(X_{i 1}, \ldots, X_{i p}\right)$. Our first result, Theorem 1 below, shows that the variable transformation method is able to rank features in accordance with the values of the estimable but unknown correlation coefficients $\omega_{j}$, defined at ( 2.31 ), even if the distributions of the $X_{i j}$ 's and/or the $Y_{i}$ 's are heavy-tailed.

We choose the function $\Psi$ to define our transformed variables $U_{i j}$ and $Z_{i}$ in terms of $X_{i j}$ and $Y_{i}$ at (2.2)), so that following properties are satisfied:
$\Psi$ is uniformly bounded and monotone increasing;
$\operatorname{var}\left(U_{1 j}\right)$ and $\operatorname{var}\left(Z_{1}\right)$ are bounded away from zero uniformly in $1 \leq j \leq p$, as $n \rightarrow \infty$.

For example, we can take $\Psi$ to be the distribution function of a symmetric, unimodal, continuous distribution with unit variance and zero mean. Consider taking $\Psi=\Phi$, the standard normal distribution. For $x$ not too large we have

$$
\begin{equation*}
\Phi(x) \approx \frac{1}{2}+(2 \pi)^{-1 / 2} x \tag{3.3}
\end{equation*}
$$

and for $x$ larger than about 2 or smaller than -2 the value of $\Phi(x)$ is virtually 0 or 1, respectively. Since transformed feature ranking methods produce identical results if $\Psi$ is replaced by $\sqrt{2 \pi}(\Psi-1 / 2)$ throughout, we see that $\Psi$ is virtually the identity for $x$ not too large, although with "barriers" in both tails that prevent $x$ from taking values that are too large positive or too large negative. This explains intuitively why this approach is more resistant to problems caused by heavy tails than a method based directly on the data or, equivalently, one that takes $\Psi$ to be linear. (Note that the approximation at (3.3) is used only as an aid to intuition, and is not employed anywhere in our analysis.)

Apart from ([3.2), the only condition we assume on the distributions is the following.

The data pairs $\left(X_{1}, Y_{1}\right), \ldots,\left(X_{n}, Y_{n}\right)$ are independent and identically distributed; the common distribution is allowed to depend on $n$.

This condition clearly does not impose any restrictive tail behaviour on the distributions. Assumption ( $\mathrm{B}_{2} 2$ ) ensures that computing the correlation at ( $\mathrm{L}, 3$ ) does not involve dividing by quantities that can become arbitrarily close to zero as $n \rightarrow \infty$. Finally, we allow $p$ to be exponentially large compared to $n$. Let $\lambda_{n}$ denote a sequence diverging to infinity as $n \rightarrow \infty$, and satisfying $\lambda_{n}=o\left(n^{1 / 2}\right)$. We assume that:
$p=O\left\{\exp \left(C \lambda_{n}^{2}\right)\right\}$ for a constant $C>0$, where $C$ depends on $\sup |\Psi|$ and on the lower bounds in (3.2).

Remark 3. It is straightforward to treat cases where location and scale corrections are made empirically using, for example, the median and standardised interquartile range, respectively. Under mild additional assumptions those quantities have the properties that, for each $\epsilon>0$, they are within $\epsilon$ of the respective true values uniformly in $1 \leq i \leq p$, if $p$ satisfies (3.5). Similar remarks apply to the results in Section 3.2.

Remark 4. The normal $\mathrm{N}\left(0, \sigma^{2}\right)$ distribution function is attractive to use in practice, since it is symmetric and involves only one user-chooseable parameter. However, there may in some cases be advantages to using a distribution function where shape, as well as scale, can be adjusted. Examples include Student's $t$ distribution functions where both scale and shape, through the number of degrees of freedom, $\nu$ say, can be adjusted. The value of $\nu$ should be interpreted in the continuum.

For the statement of Theorem 1 below, recall that $\hat{\jmath}_{1}, \ldots, \hat{\jmath}_{p}$ are defined by $\left|\widehat{\omega}_{\hat{\jmath}_{p}}\right| \geq \cdots \geq\left|\widehat{\omega}_{\hat{\jmath}_{p}}\right|$. Given constants $c_{1}$ and $c_{2}$ satisfying $0<c_{1}<c_{2}<\infty$, let $\mathcal{J}_{1}$ and $\mathcal{J}_{2}$ be the respective subsets of $\{1, \ldots, p\}$ for which $\left|\omega_{j}\right| \leq c_{1} n^{-1 / 2} \lambda_{n}$ and $\left|\omega_{j}\right| \geq c_{2} n^{-1 / 2} \lambda_{n}$. Note particularly that (B.4) does not require the components $X_{i j}$ of $X_{i}$ to be independent or uncorrelated.

Theorem 1. If ([.])-(ㄹ.5) hold then, with probability converging to 1 as $n \rightarrow \infty$, all the indices from $\mathcal{J}_{2}$ are listed in the sequence $\hat{\jmath}_{1}, \ldots, \hat{\jmath}_{p}$ before any of the indices from $\mathcal{J}_{1}$.

The implications of the theorem are perhaps most readily seen by considering a case where the indices $1, \ldots, p$ can be divided into $m+2$ distinct classes $\mathcal{I}_{0}, \mathcal{I}_{1}, \ldots, \mathcal{I}_{m+1}$, corresponding to increasingly large values of $\left|\omega_{j}\right|$. Here $m$ is
fixed. In this setting the theorem implies, among other things, that for each pair ( $k_{1}, k_{2}$ ) satisfying $0 \leq k_{1}<k_{2} \leq m+1$, with probability converging to 1 as $n \rightarrow \infty$ all the indices from $\mathcal{I}_{k_{2}}$ are listed in the sequence $\hat{\jmath}_{1}, \ldots, \hat{\jmath}_{p}$ ahead of all those from $\mathcal{I}_{k_{1}}$; moreover, if the corresponding values of $\left|\omega_{j}\right|$ are at least as large as $n^{-b}$ for some $b<1 / 2$, the sizes of these classes can be exponentially large before the claimed result breaks down. The latter issue is important, since, without imposing severe conditions on the tails of the distributions of the data, only polynomially large $p$ is permitted if feature ranking using conventional correlation is employed. This point will be discussed in more detail in Theorems 3 and 4 of Section 3.3, where we shall show, respectively, that only polynomially large $p$ is permissible if feature ranking is based on conventional correlation, but that exponentially large $p$ is possible if variable transformation methods are employed.

To be more specific about the classes $\mathcal{I}_{0}, \mathcal{I}_{1}, \ldots, \mathcal{I}_{m+1}$ of indices, define these sets by asking that $j \in \mathcal{I}_{m+1}$ if and only if $\left|\omega_{j}\right|>c$, where $c>0$ is fixed; that $j \in \mathcal{I}_{k}$ (for $1 \leq k \leq m$ ) if and only if $\left|\omega_{j}\right| \in\left[c_{k 1} n^{-b_{k}}, c_{k 2} n^{-b_{k}}\right]$, where $0<c_{k 1}<c_{k 2}<\infty$ and $1 / 2>b_{1}>\ldots>b_{m}>0$; that $j \in \mathcal{I}_{0}$ if and only if $\omega_{j}=0$. Theorem 1 implies that if $p=O\left\{\exp \left(C n^{1-2 b_{1}}\right)\right\}$ then, with probability converging to 1 , all the indices from $\mathcal{J}_{k+1}$ are listed in the sequence $\hat{\jmath}_{1}, \ldots, \hat{\jmath}_{p}$ before any of the indices from $\mathcal{J}_{k}$, and that this result holds for $0 \leq k \leq m$. More general properties, in the case where $m=m(n)$ is permitted to diverge with $n$, can be derived using (5.7) in Section 5 .

### 3.2. Performance of the transformation method based on mean differences

Analogous results hold for the feature ranking methods defined at (2.8) and ([.]). To give details of those properties we recall that in the cases represented by (2.8) and (2.]T),
the data are in the form of independent $p$-vectors $X_{i_{1}}=$ $\left(X_{i_{1} 1}, \ldots, X_{i_{1} p}\right)$ and $Y_{i_{2}}=\left(Y_{i_{2} 1}, \ldots, Y_{i_{2} p}\right)$ coming from two populations, where $1 \leq i_{k} \leq n_{k}$ and $n_{k}$ denotes the size of the sample from population $k$.

Condition (B. $\mathbf{H}$ ) on $\Psi$ is unchanged; (B.2) is replaced by the assumption that:
$\operatorname{var}\left\{\Psi\left(X_{1 j}\right)\right\}$ and $\operatorname{var}\left\{\Psi\left(Y_{1 j}\right)\right\}$ are uniformly bounded away from zero, if ranking features as at ([2.8); $\operatorname{var}\left\{\Psi\left(Y_{1 j}-X_{1 j}\right)\right\}$ is bounded away from zero, if ranking is according to (2.1]);
and (B.5) is unchanged. The following theorem describes properties of the ranking method. Its proof is almost identical to that of Theorem 1, hence we omit it.

Theorem 2. Under (B.]) and (B.5)-(B.7), with probability converging to 1 as $n \rightarrow \infty$, all the indices from $\mathcal{J}_{2}$ are listed in the sequence $\hat{\jmath}_{1}, \ldots, \hat{\jmath}_{p}$ defined by ([2.]1), before any of the indices from $\mathcal{J}_{1}$, where $\mathcal{J}_{1}$ and $\mathcal{J}_{2}$ are defined in Section 3.1.

### 3.3. Failure of feature ranking based on untransformed data

Next we show that feature ranking based on untransformed data is not effective when feature distributions can produce random variables taking relatively large values. For brevity, and because it is technically more difficult owing to the random denominator, we treat only the case of correlation, though similar results hold for the case of mean differences. Since our result is negative, and establishing it in generality would not confer significantly greater authority, we make simplifying assumptions that lead to a short, transparent proof. The problems that we describe here are present even more forcefully if we replace the independence and identical distribution assumption of (3.10) below by one where non-stationarity and correlation are present between components, since this reduces effective sample size.

In particular, we suppose independent and identically distributed data vectors $\left(X_{i}, Y_{i}\right)$, for $1 \leq i \leq n$, generated by a linear model,

$$
\begin{equation*}
Y_{i}=\alpha+\sum_{j=1}^{p} \beta_{j} X_{i j}+\epsilon_{i} \tag{3.8}
\end{equation*}
$$

where $\alpha, \beta_{1}, \ldots, \beta_{p}$ are scalar parameters, the experimental errors $\epsilon_{i}$ are identically distributed, and the $\epsilon_{i}$ 's are independent of the $X_{i j}$ 's. For simplicity we take the error $\epsilon_{i}$ in (5.8) to be identically zero. If feature selection methods based on conventional correlation are ineffective in this case then they will not perform well when error is present. Also for simplicity we assume that the intercept term in the model is zero, and there are just two nonzero coefficients $\beta_{j}$, these being $\beta_{1}=1$ and $\beta_{2}=n^{-b}$, where $0<b<1 / 2$. Therefore, the model at (3.8) reduces to

$$
\begin{equation*}
Y_{i}=X_{i 1}+n^{-b} X_{i 2} \tag{3.9}
\end{equation*}
$$

We take sample size, $n$, to be the index of our asymptotic theory, and interpret dimension, $p$, as a function of $n$.

We assume that the design variables $X_{i j}$ involve occasional outliers of size $n^{a}$, distributed as follows.

Simulate independent and identically distributed random variables $Q_{i j}$, where the common, nondegenerate distribution is compactly supported, has zero mean and does not depend on $i, j$ or $n$; then, for each $j$, choose a value of $i$ randomly (independently of the $Q_{i j}{ }^{\prime}$ 's) and uniformly between 1 and $n$, and replace $Q_{i j}$ by $n^{a}$.

If we were to assume that the distribution of the $X_{i j}$ 's was regularly varying at infinity with exponent $a^{-1}$, for example Student's $t$ distribution with $a-1$ degrees of freedom, then for each fixed $k$ the $k$ largest values of $X_{1 j}, \ldots, X_{n j}$ would be of size $n^{a}$ as $n \rightarrow \infty$. Since $k$ here can be arbitrarily large, the problems experienced by feature ranking methods would be more serious than those that we shall describe in Theorem 3 below. On the other hand if, rather than assume that an outlier is present in each design sequence $\mathcal{X}_{j}$, we were to suppose that it is present with probability $\pi_{n}$ where $\pi_{n} \rightarrow 0$ at rate $n^{-c}$, say, as $n \rightarrow \infty$, then the problems caused by outliers would be less conspicuous, but variable selection by feature ranking would still fail if $p$ was of larger order than $n^{1+c}$.

Under the model described by (B.IT) the variance of each $X_{i j}$ is $\left(1-n^{-1}\right)$ $\left\{\operatorname{var}(Q)+n^{2 a-1}\right\}$, and to prevent this quantity diverging to infinity we assume, in Theorem 3 below, that $a<1 / 2$. On the other hand, taking $a$ too small means that the outliers $n^{a}$ have relatively little input, and so, since Theorem 3 seeks to identify cases where the outliers cause problems, we also place a lower bound on $a$ in the theorem.

Recall that feature ranking by correlation orders the indices $1, \ldots, p$ as $\hat{\jmath}_{1}, \ldots, \hat{\jmath}_{p}$, where $\left|\hat{\rho}_{\hat{\jmath}_{p}}\right| \geq \cdots \geq\left|\hat{\rho}_{\hat{\jmath}_{p}}\right|$. Now, under the linear model at (3.8), if the vector components $X_{i 1}, \ldots, X_{i p}$ are independent, then we have $\rho_{j}=\beta_{j}\left(\sigma_{X j}^{2}\right)^{1 / 2} /$ $\left(\sigma_{Y}^{2}\right)^{1 / 2}$. Moreover, if $X_{i 1}, \ldots, X_{i p}$ are also independent of the errors $\epsilon_{i}$ and

$$
\begin{equation*}
\beta_{j}=0 \text { for } j>q, \beta_{j} \neq 0 \text { for } 1 \leq j \leq q \tag{3.11}
\end{equation*}
$$

then

$$
\begin{equation*}
\omega_{j}=0 \text { if } j>q, \text { and is nonzero for } 1 \leq j \leq q . \tag{3.12}
\end{equation*}
$$

A proof is given in Section 5.4. Hence, the actual correlations are $\rho_{1}=(1+$ $\left.n^{-2 b}\right)^{-1 / 2}, \rho_{2}=n^{-b}\left(1+n^{-2 b}\right)^{-1 / 2}$, and $\rho_{j}=0$, for $j \geq 3$, and an ideal ranking of features in terms of the estimated correlations $\hat{\rho}_{j}$ should at least preserve the order of the first two components among the $\rho_{j}$ 's. That is, it should have $\hat{\jmath}_{1}=1$ and $\hat{\jmath}_{2}=2$ with probability converging to 1 as $n \rightarrow \infty$. Theorem 3, below, shows that if $p$ is of sufficiently larger order than $n$ then this property fails in respect of the second component, although it holds for the first.

Theorem 3. Assume model (5.9) for the response variables $Y_{i}$, and model (3.10) for the design sequence $X_{i}$, where $1 / 4<a<1 / 2$ and $1-2 a<b<2(1-2 a)$ in (3.9) and (3.10); that, as $n \rightarrow \infty, p / n \rightarrow \infty$ and $p=O\left\{\exp \left(C n^{4 a-1}\right)\right\}$, for a constant $C>0$ depending on the distribution of $Q$ in (3.10). Then $P\left(\hat{\jmath}_{1}=1\right) \rightarrow 1$ and $P\left(\hat{\jmath}_{2}=2\right) \rightarrow 0$ as $n \rightarrow \infty$.

This proves that standard correlation learning is rather ineffective in the case of heavy-tailed distributions. In contrast, we have already proved in Theorem 1
that correlation learning based on variable transformation performs particularly well. To simplify comparison with the negative result of Theorem 3, in the next theorem we give a more detailed, positive, result for transformation-based correlation learning under the simple model treated in Theorem 3. In particular, we show that $P\left(\hat{\jmath}_{1}=1\right)$ and $P\left(\hat{\jmath}_{2}=2\right)$ both converge to 1 as $n \rightarrow \infty$, when we define $\hat{\jmath}_{1}, \ldots, \hat{\jmath}_{p}$ by $\left|\widehat{\omega}_{\hat{\jmath}_{p}}\right| \geq \cdots \geq\left|\widehat{\omega}_{\hat{\jmath}_{p}}\right|$, where $\widehat{\omega}_{j}=\hat{\zeta}_{j} /\left(\hat{\tau}_{X j} \hat{\tau}_{Y}\right)$ and the quantities on the right-hand side of this equation are given by (2.4).

Theorem 4. Assume the conditions of Theorem 3, except that the upper bound on $p$ there is replaced by $p=O\left\{\exp \left(C n^{1-2 b}\right)\right\}$, where again $C>0$ depends on the distribution of $Q$. Suppose that the function $\Psi$ used to transform the data is bounded and strictly monotone increasing, and has two bounded derivatives. Then both $P\left(\hat{\jmath}_{1}=1\right)$ and $P\left(\hat{\jmath}_{2}=2\right)$ converge to 1 as $n \rightarrow \infty$.

It is readily seen that Theorems 3 and 4 have a non-null intersection. That is, in the context of the models described by ([.TM) and (B.IT), the transformation approach leads to correct rankings in cases where relatively conventional methods produce incorrect results.

## 4. Data Illustrations

The simulation results reported in Section 2 illustrate the theoretical properties of Section 3 on simple examples, but simulations in more complex settings can be found in Delaigle and Hall (201). In this section we show for some data examples that our theoretical results in Section 3 translate into significant improved ranking by using the transformation approach. To apply the transformation procedure we need to choose the function $\Psi$. In Section 3.1 we gave intuitive arguments suggesting that $\Psi$ could be taken to be the standard normal distribution. Of course, for these intuitive arguments to be valid we need to rescale the data in some way before applying the transformation $\Psi$, as we need the non-aberrant data values to be located roughly between -2 and 2. Our theory is also valid when we do this; see Remark 3 in Section 3.1. We do not rescale the data in the same way for the methods based on correlation and mean differences, so below we give details for these approaches separately.

For the correlation-based method we take $U_{i j}=\Psi\left\{\left(X_{i j}-m_{j}\right) / s_{j}\right\}$ and $Z_{i}=\Psi\left\{\left(Y_{i}-m\right) / s\right\}$, where $m_{j}$ and $s_{j}$ are, respectively, the sample median and standardised interquartile range of $X_{i j}$, and $m$ and $s$ are defined in the same way for $Y_{i}$. Note that, since correlation is invariant under changes of scale and location, and heavy-tailedness properties are also unaffected by such changes, the theoretical results in Section 3.1 are unchanged if we replace $X_{i j}$ and $Y_{i}$ there by, respectively, $\left(X_{i j}-a_{j}\right) / b_{j}$ and $\left(Y_{i}-a\right) / b$, where $a, b, a_{j}$ and $b_{j}$ are bounded


Figure 1. Ro131 study: scatterplots of the first 16 genes selected by the variable transformation approach, that were not ranked among the top 50 by the untransformed correlation method. Left group of 16: scatterplots of the $\left(X_{i j}, Y_{i}\right)$ 's for those 16 genes; right group of 16: scatterplots of the $\left(U_{i j}, Z_{i}\right)$ 's, where $U_{i j}=\Psi\left(X_{i j}\right)$ and $Z_{i}=\Psi\left(Y_{i}\right)$. The oblique lines are the least squares lines.


Figure 2. Scatterplot of gene 1,597_at for the Affymetrix data.
nonzero constants. See Remark 3 in Section 3.1 for the case of empirical location and scale.

For the mean-based method founded on (2.8) in Section 2.3 we take $U_{i j}=$ $\Psi\left\{\left(X_{i j}-m_{j}\right) / s_{j}\right\}$ and $V_{i j}=\Psi\left\{\left(Y_{i j}-m_{j}\right) / s_{j}\right\}$, where $m_{j}$ is the minimum of the sample medians of $X_{i j}$ and $Y_{i j}$, and $s_{j}$ is the standardised interquartile range of the pooled sample of $X_{i j}$ and $Y_{i j}$ (with $j$ fixed). Note that we cannot center the components $X_{i j}$ and $Y_{i j}$ to their respective medians, since our goal is precisely
to detect differences in location. In this case, too, rescaling the data can be accommodated by our theory; see Remark 3 in Section 3.1.

### 4.1. Cardiomyopathy microarray data

Our first application concerns the ranking of genes according to their influence on overexpression of a G-protein-coupled receptor (Ro1) in mice. See Segal, Dahlquist, and Conklin (2003) and Hall and Miller (2009) for a more detailed description of the data. Here, $Y_{i}$ is the measurement of Ro1 on the $i$ th mouse, and $\left(X_{i 1}, \ldots, X_{i p}\right)$ are the expression levels of $p$ genes for the $i$ th mouse. The sample size was $n=30$, and $p=6,319$ genes were observed. We ranked the genes using both the untransformed correlations $\hat{\rho}_{j}$ and the transformed correlations $\hat{\omega}_{j}$. Out of the 50 genes ranked highest by the transformed correlation method, 40 were not ranked in the top 50 by the untransformed correlation approach. In Figure 1 we show scatterplots of the $\left(X_{i j}, Y_{i}\right)$ 's and of the $\left(U_{i j}, Z_{i}\right)$ 's for the first 16 of these 40 genes. Above each scatterplot we indicate the rank assigned by the untransformed correlation method for each gene in the left panel, and the rank assigned when we used the transformation based approach in the right panel. In each scatterplot we also show the least squares regression line.

For a number of the 16 genes the least squares lines, and thus the empirical correlations, for the untransformed data appear to be too highly influenced by outliers. As a result, some of the genes were ranked by the untransformed correlation approach much lower than they arguably should have been. This is particularly striking for the genes ranked $346,192,407,620$, and 1,143 by untransformed correlation, and ranked $12,15,17,18$, and 23 after transforming the data. For other genes, such as the genes originally ranked 55 and 126 , transforming the variables did not seem particularly useful. As we noted in Section 2.3, transforming the data does not necessarily improve the rank of each gene, but it manages to identify influential genes that were wrongly disqualified because of outliers.

### 4.2. Affymetrix spike-in data

We applied the transformed mean ranking approach to the Affymetrix spikein data described by Cope et al. (2004). This dataset contains two groups with $n=12$ observations in each group, and $p=12,626$ genes. It is available from http://strimmerlab.org/data.html. In Cope et al. (2004), a comparison with benchmark data permitted the identification of 16 genes believed to be differentially expressed (that is, for which the mean differences between the two groups were significantly different from zero). We applied the method in Section 2.3 to
these data, and found that we needed to keep the 47 highest values of $\left|\bar{V}_{j}-\bar{U}_{j}\right|$ in order to select these 16 known genes. By comparison, the method based on Student's $t$ scores needed to keep 155 components, and the method based on the untransformed means $\left|\bar{Y}_{j}-\bar{X}_{j}\right|$ had to keep 100 . We also tried the shrinkage method of Opgen-Rhein and Strimmer (2007), and had to keep 54 components in order to include all 16 identified genes. The method of Efron et all (2007) needed 61 components, and the other methods needed even more. The better performance of the transformation based procedure comes from the fact that many components have one or more outliers. For example, the scatterplot of the observations of gene 1,597_at in Figure 2 shows that the data on gene 1,597_at (which, of the 16 genes, was the most difficult to pick up for most methods) contain an outlier, and this is precisely the setting in which our method can bring considerable improvements. Note that our method did not rank all the 16 known genes higher than the other methods, but managed to pick them all more efficiently, by keeping fewer components.

## 5. Technical Arguments

### 5.1. Proof of Theorem 1

It follows from Bernstein's Inequality that, if $W, W_{1}, \ldots, W_{n}$ are independent random variables for which $P\left(|W| \leq C_{1}\right)=1, E(W)=0$ and $E\left(W^{2}\right)=1$, and if we define $S=n^{-1 / 2} \sum_{i} W_{i}$, then for each $C_{2}>0$ there exists $C_{3}>0$, depending only on $C_{1}$ and $C_{2}$, such that $P(|S|>x) \leq 2 \exp \left(-C_{3} x^{2}\right)$ for all $x \in\left(0, C_{2} n^{1 / 2}\right]$. Applying this result to $W=\left\{W^{\prime}-E\left(W^{\prime}\right)\right\} /\left(\operatorname{var} W^{\prime}\right)^{1 / 2}$, where $W^{\prime}$ is $U_{i j}, U_{i j}^{2}$, $Z_{i}, Z_{i}^{2}$ or $U_{i j} Z_{i}$, it can be deduced that

$$
P\left(n^{1 / 2}\left|\widehat{\omega}_{j}-\omega_{j}\right|>x\right) \leq 2 \exp \left(-C_{4} x^{2}\right),
$$

for all $x \in\left(0, C_{2} n^{1 / 2}\right]$, where $C_{4}>0$ depends only on sup $|\Psi|$ and $C_{2}$. Therefore, if $p=O\left\{\exp \left(C_{5} \lambda_{n}^{2}\right)\right\}$ for $C_{5}>0$ sufficiently small, then, for all $C_{6}>0$,

$$
\begin{equation*}
P\left(\max _{1 \leq j \leq p}\left|\widehat{\omega}_{j}-\omega_{j}\right|>n^{-1 / 2} C_{6} \lambda_{n}\right) \rightarrow 0 \tag{5.1}
\end{equation*}
$$

Result (5.]) implies Theorem 1.

### 5.2. Proof of Theorem 3

Under the model at (3..प) the correlation estimator $\hat{\rho}_{j}$, defined below ([.]), is given by

$$
\begin{equation*}
\hat{\rho}_{j}=\left(\hat{a}_{j 1}+n^{-b} \hat{a}_{j 2}\right)\left(\hat{b}_{j} \hat{b}_{Y}\right)^{-1} \tag{5.2}
\end{equation*}
$$

for $1 \leq j \leq p$, where $\hat{b}_{Y}=\left\{n^{-1} \sum_{i}\left(Y_{i}-\bar{Y}\right)^{2}\right\}^{1 / 2}$,

$$
\hat{a}_{j k}=\frac{1}{n} \sum_{i=1}^{n} X_{i j} \widehat{w}_{i k}, \quad \hat{b}_{j}=\left\{\frac{1}{n} \sum_{i=1}^{n}\left(X_{i j}-\bar{X}_{j}\right)^{2}\right\}^{1 / 2}, \quad \widehat{w}_{i k}=X_{i k}-\bar{X}_{k} .
$$

Let $\mathcal{E}_{j k}$ denote the event that the values of $n^{a}$ are in the same position in the sequences $\mathcal{X}_{j}$ and $\mathcal{X}_{k}$, and write $\widetilde{\mathcal{E}}_{j k}$ for the complement of $\mathcal{E}_{j k}$. For each $C_{1}>1$ there exists $C_{2}>0$ such that for all $j \neq k$ and all $d \in(0,1 / 2)$,

$$
\begin{equation*}
P\left(\left|n^{1 / 2} \hat{a}_{j k}\right|>C_{1} n^{d} \mid \widetilde{\mathcal{E}}_{j k}\right)=O\left\{\exp \left(-C_{2} n^{2 d}\right)\right\} \tag{5.3}
\end{equation*}
$$

while if $1 / 4<a<1 / 2$ then for each $C_{1}>0$ there exists $C_{2}>0$ such that, for $j \neq k$ and all $d \in(0,2 a-1 / 2]$,

$$
\begin{equation*}
P\left(\left|n^{1 / 2} \hat{a}_{j k}-n^{2 a-1 / 2}\right|>C_{1} n^{d} \mid \mathcal{E}_{j k}\right)=O\left\{\exp \left(-C_{2} n^{2 d}\right)\right\} . \tag{5.4}
\end{equation*}
$$

(If $j \neq k$ then the left-hand sides of (5.3) and (5.4) do not depend on $j$ and $k$. Each of (5.3)-(5.5) is proved using Bernstein's Inequality.)

Note too that $\hat{a}_{j j}=\hat{b}_{j}^{2}$, and that we may assume without loss of generality that the common distribution of the variables $X_{i j}$ that do not equal $n^{a}$ has unit variance. Then, if $C_{1}>1$, there exists $C_{2}>0$ such that, for all $j$ and $d=2 a-1 / 2$,

$$
\begin{equation*}
P\left(n^{1 / 2}\left|\hat{b}_{j}-1\right|>C_{1} n^{d}\right)=O\left\{\exp \left(-C_{2} n^{2 d}\right)\right\} \tag{5.5}
\end{equation*}
$$

(The left-hand side of (5.5) does not depend on $j$.) Furthermore,

$$
\begin{equation*}
P\left(\mathcal{E}_{j k}\right)=n^{-1} . \tag{5.6}
\end{equation*}
$$

Note too that, by (5.2),
$\hat{b}_{Y} \hat{\rho}_{1}=\hat{b}_{1}+n^{-b} \hat{a}_{12} \hat{b}_{1}^{-1}, \hat{b}_{Y} \hat{\rho}_{2}=\hat{a}_{12} \hat{b}_{2}^{-1}+n^{-b} \hat{b}_{2}, \hat{b}_{Y} \hat{\rho}_{j}=\left(\hat{a}_{j 1}+n^{-b} \hat{a}_{j 2}\right) \hat{b}_{j}^{-1}$,
where the last equation will be used for $j \geq 3$.
Using (5.3), (5.4) and (5.5), each with $d=2 a-1 / 2$, it can be deduced that $\hat{a}_{j k}=O_{p}\left(n^{-1 / 2+2 a-1 / 2}\right)=O_{p}\left(n^{-(1-2 a)}\right)$ and $\hat{b}_{j}=1+O_{p}\left(n^{-(1-2 a)}\right)$, uniformly in $1 \leq j \leq p$ and in $k=1,2$. (The choice $d=2 a-1 / 2$ is responsible for the exponent $2 d=4 a-1$ in the formula $p=O\left\{\exp \left(C n^{4 a-1}\right)\right\}$ in the statement of Theorem 4.1.) Hence, using (5.7), we obtain $\hat{b}_{Y} \hat{\rho}_{1}=1+O_{p}\left(n^{-(1-2 a)}\right), \hat{b}_{Y} \hat{\rho}_{2}=$ $O_{p}\left(n^{-(1-2 a)}+n^{-b}\right)$ and $\hat{b}_{Y} \hat{\rho}_{j}=O_{p}\left(n^{-(1-2 a)}\right)$ uniformly in $3 \leq j \leq p$. Therefore

$$
\begin{equation*}
P\left[\left|\hat{\rho}_{1}\right|>\max \left\{\left|\hat{\rho}_{2}\right|, \ldots,\left|\hat{\rho}_{p}\right|\right\}\right] \rightarrow 1 . \tag{5.8}
\end{equation*}
$$

The same arguments show that

$$
\begin{equation*}
\hat{b}_{Y} \hat{\rho}_{j}=\hat{a}_{j 1}+o_{p}\left(n^{-b}\right) \quad \text { uniformly in } 3 \leq j \leq p, \tag{5.9}
\end{equation*}
$$

and that

$$
\begin{aligned}
\hat{b}_{Y} \hat{\rho}_{2} & =\hat{a}_{12} \hat{b}_{2}^{-1}+n^{-b} \hat{b}_{2} \\
& =n^{-(1-2 a)} I\left(\mathcal{E}_{12}\right)+O_{p}\left(n^{-2(1-2 a)}\right)+n^{-b}+O_{p}\left(n^{-b-(1-2 a)}\right) \\
& =n^{-(1-2 a)} I\left(\mathcal{E}_{12}\right)+n^{-b}\left\{1+o_{p}(1)\right\}=n^{-b}\left\{1+o_{p}(1)\right\},
\end{aligned}
$$

where the second last identity follows from the fact that, by assumption, $b<$ $2(1-2 a)$, and the last identity since $P\left(\mathcal{E}_{12}\right)=n^{-1} \rightarrow 0$. Therefore

$$
\begin{equation*}
\hat{b}_{Y} \hat{\rho}_{2}=n^{-b}\left\{1+o_{p}(1)\right\} . \tag{5.10}
\end{equation*}
$$

Results (5.3) and (5.4) imply that, for a constant $C_{3}>0$,

$$
P\left(\left|n^{1 / 2} \hat{a}_{j k}-n^{2 a-1 / 2} I\left(\mathcal{E}_{j k}\right)\right|>C_{1}\right)=O\left\{\exp \left(-C_{3} n^{2 d}\right)\right\} .
$$

This result, and the fact that $n^{-1 / 2}=o\left(n^{2 a-1}\right)$ (by assumption, $a>1 / 4$ ), imply that if $d=2 a-1 / 2$ and

$$
\begin{equation*}
p=o\left\{\exp \left(C n^{2 d}\right)\right\}, \tag{5.11}
\end{equation*}
$$

where the constant $C$ here and in the statement of the theorem satisfies $C<C_{3}$, then

$$
\begin{equation*}
\hat{a}_{j 1}=n^{2 a-1} I\left(\mathcal{E}_{j 1}\right)+o_{p}\left(n^{2 a-1}\right) \quad \text { uniformly in } 3 \leq j \leq p . \tag{5.12}
\end{equation*}
$$

Conditional on the dataset $\mathcal{X}_{1}=\left\{X_{11}, \ldots, X_{n 1}\right\}$, the events $\mathcal{E}_{j 1}$, for $3 \leq j \leq p$, are independent and satisfy $P\left(\mathcal{E}_{j 1} \mid \mathcal{X}_{1}\right)=n^{-1}$. Therefore if $p / n \rightarrow \infty$ then $P\left(\mathcal{E}_{31} \cup \ldots \cup \mathcal{E}_{p 1}\right) \rightarrow 1$. Combining this property with (5.工2) we deduce that if (5.لत) holds then for each $\epsilon>0$,

$$
\begin{equation*}
P\left\{\hat{a}_{j 1}>(1-\epsilon) n^{2 a-1} \quad \text { for at least one } j \text { in the range } 3 \leq j \leq p\right\} \rightarrow 1 \tag{5.13}
\end{equation*}
$$

Property (5.8) implies that $P\left(\hat{\jmath}_{1}=1\right) \rightarrow 1$. Properties (5.Y) and (5.J.3), and the inequality $1-2 a<b$ assumed in Theorem 3, imply that for each $\epsilon>0$,

$$
P\left\{\hat{b}_{Y} \hat{\rho}_{j}>(1-\epsilon) n^{2 a-1} \quad \text { for at least one } j \text { in the range } 3 \leq j \leq p\right\} \rightarrow 1
$$

which, together with (5.ll) and the property $1-2 a<b$, implies that $P\left(\hat{\jmath}_{2}=\right.$ 2) $\rightarrow 0$.

### 5.3. Proof of Theorem 4

Put $U_{i j}=\Psi\left(X_{i j}\right)$ and $Z_{i}=\Psi\left(Y_{i}\right)$, where $Y_{i}$ is as at (B.T), and define

$$
\begin{equation*}
\tilde{a}_{j}=\frac{1}{n} \sum_{i=1}^{n} Z_{i} \widetilde{w}_{i j}, \quad \tilde{b}_{j}=\left\{\frac{1}{n} \sum_{i=1}^{n}\left(U_{i j}-\bar{U}_{j}\right)^{2}\right\}^{1 / 2}, \quad \widetilde{w}_{i j}=U_{i j}-\bar{U}_{j}, \tag{5.14}
\end{equation*}
$$

$\tilde{b}_{Z}=\left\{n^{-1} \sum_{i}\left(Z_{i}-\bar{Z}\right)^{2}\right\}^{1 / 2}$ and $\widehat{\omega}_{j}=\tilde{a}_{j} /\left(\tilde{b}_{j} \tilde{b}_{Z}\right)$. Without loss of generality, the common distribution of the variables $\Psi\left(X_{i j}\right)$ when $X_{i j}$ does not equal $n^{a}$ has unit variance. (This property is achievable by simple rescaling of $\Psi$, since $\Psi$ is strictly monotone and, by ( $\mathbf{B} . \mathbf{d} \mathbf{l}), X_{i j}$ given that $X_{i j} \neq n^{a}$ has nonzero variance.) Result (5.54) below replaces both (5.3) and (5.4) in the present setting, and (5. 5 ) replaces (5.5). Both are derived using Bernstein's Inequality, treating separately the instance where $X_{i j}=n^{a}$ (for each $j$, this is true for exactly one $i$ ) and the contrary case. For each $C_{1}>1$ there exists $C_{2}>0$ such that, for all $d \in(0,1 / 2)$,

$$
\begin{align*}
\sup _{1 \leq j \leq p} P\left\{\left|n^{1 / 2}(1-E) \tilde{a}_{j}\right|>C_{1} n^{d}\right\} & =O\left\{\exp \left(-C_{2} n^{2 d}\right)\right\},  \tag{5.15}\\
\sup _{1 \leq j \leq p} P\left(n^{1 / 2}\left|\tilde{b}_{j}-1\right|>C_{1} n^{d}\right) & =O\left\{\exp \left(-C_{2} n^{2 d}\right)\right\} . \tag{5.16}
\end{align*}
$$

For a given value of $i$ let $\mathcal{E}_{j}$ denote the event that $j$ is an index for which $X_{i j} \neq n^{a}$. Then if $\mathcal{E}_{1} \cap \mathcal{E}_{2}$ holds,

$$
\begin{equation*}
\left|\Psi\left(Y_{i}\right)-\left\{\Psi\left(X_{i 1}\right)+n^{-b} X_{i 2} \Psi^{\prime}\left(X_{i 1}\right)\right\}\right| \leq C n^{-2 b} \tag{5.17}
\end{equation*}
$$

where $C$ denotes a constant. (Here we have used the assumption that $\Psi$ is twice differentiable.) Therefore,

$$
\begin{align*}
& \operatorname{cov}\left\{\Psi\left(Y_{i}\right), \Psi\left(X_{i j}\right) \mid \mathcal{E}_{1} \cap \mathcal{E}_{2} \cap \mathcal{E}_{j}\right\} \\
& =n^{-b} \operatorname{cov}\left\{X_{i 2}, \Psi\left(X_{i j}\right) \mid \mathcal{E}_{1} \cap \mathcal{E}_{2} \cap \mathcal{E}_{j}\right\} E\left\{\Psi^{\prime}\left(X_{i 1}\right) \mid \mathcal{E}_{1} \cap \mathcal{E}_{2} \cap \mathcal{E}_{j}\right\}+O\left(n^{-2 b}\right), \tag{5.18}
\end{align*}
$$

uniformly in $i$ and $j \geq 2$.
If $j=2$ then, since $\Psi$ is strictly monotone increasing, $\operatorname{cov}\left\{X_{i 2}, \Psi\left(X_{i j}\right) \mid \mathcal{E}_{1} \cap\right.$ $\left.\mathcal{E}_{2}\right\}>0$ and $E\left\{\Psi^{\prime}\left(X_{i 1}\right) \mid \mathcal{E}_{1} \cap \mathcal{E}_{j}\right\}>0$, and so

$$
\gamma \equiv \operatorname{cov}\left\{X_{i 2}, \Psi\left(X_{i 2}\right) \mid \mathcal{E}_{1} \cap \mathcal{E}_{2}\right\} E\left\{\Psi^{\prime}\left(X_{i 1}\right) \mid \mathcal{E}_{1}\right\}>0
$$

It follows from (5.17) that

$$
\begin{equation*}
\operatorname{cov}\left\{\Psi\left(Y_{i}\right), U_{i 1} \mid \mathcal{E}_{1} \cap \mathcal{E}_{2}\right\}=\operatorname{var}\left(U_{i 1} \mid \mathcal{E}_{1}\right)+O\left(n^{-b}\right) \tag{5.19}
\end{equation*}
$$

In view of (5.18),

$$
\begin{equation*}
\operatorname{cov}\left\{\Psi\left(Y_{i}\right), U_{i 2} \mid \mathcal{E}_{1} \cap \mathcal{E}_{2}\right\}=n^{-b} \gamma+O\left(n^{-2 b}\right), \tag{5.20}
\end{equation*}
$$

and, more simply,

$$
\begin{equation*}
\operatorname{cov}\left\{\Psi\left(Y_{i}\right), U_{i j}\right\}=0 \quad \text { for } \quad 3 \leq j \leq p \tag{5.21}
\end{equation*}
$$

Results (5.15) and (5.19), and the fact that (by assumption) var $\left\{\Psi\left(X_{i j}\right)\right\}=1$ when $X_{i j} \neq n^{a}$ (i.e. $X_{i j}=Q_{i j}$, as in ( (B.I)), imply that for each $C_{1}>0$, there exists $C_{2}>0$ such that

$$
\begin{equation*}
P\left\{\left|\tilde{a}_{1}-1\right|>C_{1}\left(n^{d-1 / 2}+n^{-b}\right)\right\}=O\left\{\exp \left(-C_{2} n^{2 d}\right)\right\} \tag{5.22}
\end{equation*}
$$

whenever $d \in(0,1 / 2)$.
By (5.5), (5.20), and (5.21), for each $d \in(0,1 / 2)$ and each $C_{1}>0$, we can choose $C_{2}>0$ such that

$$
\begin{align*}
P\left\{\left|\tilde{a}_{2}-n^{-b} \gamma\right|>C_{1}\left(n^{d-1 / 2}+n^{-b}\right)\right\} & =O\left\{\exp \left(-C_{2} n^{2 d}\right)\right\},  \tag{5.23}\\
\sup _{3 \leq j \leq p} P\left\{\left|\tilde{a}_{j}\right|>C_{1}\left(n^{d-1 / 2}+n^{-b}\right)\right\} & =O\left\{\exp \left(-C_{2} n^{2 d}\right)\right\} . \tag{5.24}
\end{align*}
$$

In progressing from (5.19) $-(5.21)$ to $(5.22)-(5.24)$ we used the fact that, for each $j$, there is just one $i$ in the sequence $i=1, \ldots, n$ for which $X_{i j}=n^{a}$, and that, by the definition of $\tilde{a}_{j}$ at (5.]4), and since $P\left(\left|Z_{i} \widetilde{w}_{i j}\right| \leq 2 \sup \Psi^{2}\right)=1$, this $i$ (which is randomly chosen without regard for the values of $X_{i j}$ that do not equal $n^{a}$ ) contributes no more than $2 n^{-1} \sup \Psi^{2}$ to the value of $\tilde{a}_{j}$. That contribution is of strictly smaller order than $n^{-b}$ in (5.22) and (5.24).

Combining (5.56) and (5.22)-(5.24) we deduce that, taking $d=1 / 2-b$ and assuming that $p=O\left\{\exp \left(C n^{1-2 b}\right)\right\}$ for $C>0$ sufficiently small, we have, as $n \rightarrow \infty, \tilde{b}_{Z} \widehat{\omega}_{1}=1+o_{p}(1), \tilde{b}_{Z} \widehat{\omega}_{2}=n^{-b} \gamma+o_{p}\left(n^{-b}\right)$ where $\gamma>0$, and $\tilde{b}_{Z} \widehat{\omega}_{j}=o_{p}\left(n^{-b}\right)$ uniformly in $3 \leq j \leq p$. These three properties imply that $P\left(\hat{\jmath}_{1}=1\right)$ and $P\left(\hat{\jmath}_{2}=2\right)$ converge to 1 as $n \rightarrow \infty$, as claimed in the theorem.

### 5.4. Proof that (3.8) and (3.1]) imply (3.T2)

Note that under (3.8) and (3.1]) and the independence assumption above (3.1]), we have, trivially, $\rho_{j}=0$ if $j>q$, and is nonzero for $1 \leq j \leq q$. This result is useful in the proofs. To prove ( $\mathrm{B}_{2}[2)$ ), it suffices to show that if $\Psi$ is uniformly bounded, then the covariance between $\Psi\left(Y_{i}\right)$ and $\Psi\left(X_{i j}\right)$ vanishes when $\beta_{j}=0$, and if in addition $\Psi$ is strictly monotone then the covariance is nonzero when $\beta_{j} \neq 0$. The first of these results is trivial since, under (3.8) and the independence assumption, $\rho_{l}=0, \Psi\left(Y_{i}\right)$ and $\Psi\left(X_{i j}\right)$ are independent if $\beta_{j}=0$. To derive the second result it suffices to show that if random variables $V_{1}$ and $V_{2}$ are independent, if $V_{1}$ is essentially bounded and nondegenerate, and if functions $\Psi_{1}$ and $\Psi_{2}$ are strictly monotone increasing and bounded, then $\gamma \equiv$ $\operatorname{cov}\left[\Psi_{1}\left\{\Psi_{2}\left(V_{1}\right)+V_{2}\right\}, V_{1}\right]>0$. The function $\Psi_{3}=\Psi_{2}\left(\cdot+E V_{1}\right)$ is monotone if $\Psi_{2}$ is, and in this notation $\gamma=\operatorname{cov}\left[\Psi_{1}\left\{\Psi_{3}\left(V_{1}-E V_{1}\right)+V_{2}\right\}, V_{1}-E V_{1}\right]$, so without loss of generality $E\left(V_{1}\right)=0$, in which case, since $V_{1}$ and $V_{2}$ are independent, $\gamma=$
$E\left[\Psi_{1}\left\{\Psi_{2}\left(V_{1}\right)+V_{2}\right\} V_{1}\right]=E\left(E\left[\Psi_{1}\left\{\Psi_{2}\left(V_{1}\right)+V_{2}\right\} V_{1} \mid V_{2}\right]\right)=E\left(\operatorname{cov}\left[\Psi_{1}\left\{\Psi_{2}\left(V_{1}\right)+\right.\right.\right.$ $\left.\left.\left.V_{2}\right\}, V_{1} \mid V_{2}\right]\right)=E\left\{\Psi_{4}\left(V_{2}\right)\right\}$, where $\Psi_{4}(v)=E\left(\operatorname{cov}\left[\Psi_{1}\left\{\Psi_{2}\left(V_{1}\right)+V_{2}\right\}, V_{1} \mid V_{2}=v\right]\right)$. That is, $\gamma=E\left\{\Psi_{4}\left(V_{2}\right)\right\}$. Therefore it suffices to consider the case where $V_{2}$ is identically constant. In this case we can absorb $V_{2}$ into the definition of $\Psi_{1}$, and so it is sufficient to take $V_{2}=0$. Therefore we must show that when $\Psi_{1}$ and $\Psi_{2}$ are strictly monotone increasing and $\Psi_{1}$ is bounded, $\operatorname{cov}\left[\Psi_{1}\left\{\Psi_{2}\left(V_{1}\right)\right\}, V_{1}\right]>0$, or equivalently, if $\Psi_{3}$ is strictly monotone increasing and bounded (and $V_{1}$ is essentially bounded and not degenerate) then $\gamma_{1} \equiv \operatorname{cov}\left\{\Psi_{3}\left(V_{1}\right), V_{1}\right\}>0$. To appreciate that this inequality holds, let $V$ be a random variable distributed as $V_{1}$ and independent of it. Note that $\left\{\Psi_{3}\left(V_{1}\right)-\Psi_{3}(V)\right\}\left(V_{1}-V\right) \geq 0$, and is strictly positive whenever $V_{1} \neq V$. Therefore, $E\left[\left\{\Psi_{3}\left(V_{1}\right)-\Psi_{3}(V)\right\}\left(V_{1}-V\right)\right]>0$. Since $E\left(V_{1}\right)=E(V)=0$, the left-hand side here is $2 \operatorname{cov}\left\{\Psi_{3}\left(V_{1}\right), V_{1}\right\}$.

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Department of Mathematics and Statistics, University of Melbourne, Victoria 3010, Australia.
E-mail: A.Delaigle@ms.unimelb.edu.au
Department of Mathematics and Statistics, University of Melbourne, Victoria 3010, Australia.
E-mail: halpstat@ms.unimelb.edu.au
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