# Robust estimation for the multivariate linear model based on a $\tau$-scale 

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#### Abstract

We introduce a class of robust estimates for multivariate linear models. The regression coefficients and the covariance matrix of the errors are estimated simultaneously by minimizing the determinant of the covariance matrix estimate, subject to a constraint on a robust scale of the Mahalanobis norms of the residuals. By choosing a $\tau$-estimate as a robust scale, the resulting estimates combine good robustness properties and asymptotic efficiency under Gaussian errors. These estimates are asymptotically normal and in the case where the errors have an elliptical distribution, their asymptotic covariance matrix differs only by a scalar factor from the one corresponding to the maximum likelihood estimate. We derive the influence curve and prove that the breakdown point is close to 0.5 . A Monte Carlo study shows that our estimates compare favorably with respect to $S$-estimates. © 2005 Elsevier Inc. All rights reserved.


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## 1. Introduction

Let $y_{i}=\left(y_{i 1}, \ldots, y_{i q}\right)^{\prime}$ and $x_{i}=\left(x_{i 1}, \ldots, x_{i p}\right)^{\prime}, 1 \leqslant i \leqslant n$, be the response and predictor vectors satisfying the multivariate linear model (MLM)

$$
\begin{equation*}
y_{i}=B_{0}^{\prime} x_{i}+u_{i}, \tag{1.1}
\end{equation*}
$$

[^0]where $B_{0}$ is a $p \times q$ matrix and $u_{1}, \ldots, u_{n}$ are i.i.d. $q$-dimensional vectors. The $x_{i}$ 's may be fixed or random, and in the latter case we assume that they are i.i.d. and independent of the $u_{i}$ 's. If the model includes an intercept, then $x_{i p}=1$. We denote the distributions of $u_{i}$ and $x_{i}$ by $F_{0}$ and by $M_{0}$, respectively.

Assuming that the errors $u_{i}$ 's have a multivariate normal distribution with mean 0 and covariance matrix $\Sigma_{0}\left(N_{q}\left(0, \Sigma_{0}\right)\right)$, the maximum likelihood estimator (MLE) of $B_{0}$ can be obtained by computing the least squares estimate (LSE) for each component of $y$ separately, and the MLE of $\Sigma_{0}$ is the sample covariance matrix of the corresponding residuals. As it is well known, the LSE is extremely sensitive to outliers. In fact, just one observation may have an unbounded effect on this estimate.

The first proposal of a robust estimate for the MLM was given by Koenker and Portnoy [10]. They proposed to apply a regression M-estimator based on a convex $\rho$-function to each coordinate of the response vector. This proposal has two disadvantages: lack of affine equivariance and a null breakdown point. The second problem may be overcome by replacing the M-estimator by a high breakdown point estimate, but we would still lack affine equivariance.

Recently, several robust equivariant estimates for the MLM have been proposed. Rousseeuw et al. [15] proposed estimates for the MLM based on a robust estimate of the covariance matrix of $z=\left(x^{\prime}, y^{\prime}\right)^{\prime}$. A different approach based on extending estimates of multivariate location and scatter was followed by Bilodeau and Duchesne [2] who extended the S-estimates introduced by Davies [4], and by Agulló et al. [1] who extended the minimum covariance determinant estimate introduced by Rousseeuw [13].

All these estimates have a high breakdown point and therefore a good robustness behavior. However, they are not highly efficient when the errors are Gaussian and $q$ is small. Agulló et al. [1] improved the efficiency of their estimates, maintaining their high breakdown point, by considering one-step reweighting and one-step Newton-Raphson GM-estimates.

In this paper we propose robust estimates for the MLM by extending the $\tau$-estimates of multivariate location and scatter proposed by Lopuhaä [11]. We show that these estimates simultaneously have a high breakdown point and a high efficiency under Gaussian errors.

In Section 2 we define $\tau$-estimates for MLM. In Section 3 we study their breakdown point and in Section 4 we derive the influence curve. In Section 5 we study the asymptotic properties (consistency and asymptotic normality) of the $\tau$-estimates assuming random regressors and errors with an elliptical unimodal distribution. In Section 6 we describe a computing algorithm based on an iterative weighted MLE. In Section 7 we present the results of a Monte Carlo study and a real example in Section 8. In the Appendix we derive some mathematical results.

## 2. Estimates based on a robust scale for the MLM

Let $\left(x_{i}, y_{i}\right), 1 \leqslant i \leqslant n$, satisfy the MLM (1.1), where the $u_{i}$ are i.i.d. random vectors. Let $X$ be the $n \times p$ matrix whose $i$ th row is $x_{i}^{\prime}$ and $Y$ the $n \times q$ matrix whose $i$ th row is $y_{i}^{\prime}$. In the rest of the paper we will assume that the rank of $X$ is $p$. For any $p \times q$ matrix $B$ define the residuals $\widehat{u}_{i}(B)$ as

$$
\begin{equation*}
\widehat{u}_{i}(B)=y_{i}-B^{\prime} x_{i}, \quad 1 \leqslant i \leqslant n . \tag{2.1}
\end{equation*}
$$

The MLE of $B_{0}$ and $\Sigma_{0}$ when the distribution of the $u_{i}$ 's is multivariate normal with mean 0 and covariance matrix $\Sigma_{0}$ are given by

$$
\widehat{B}=\left(X^{\prime} X\right)^{-1} X^{\prime} Y,
$$

and

$$
\begin{equation*}
\widehat{\Sigma}=\frac{1}{n} \sum_{i=1}^{n} \widehat{u}_{i}(\widehat{B}) \widehat{u}_{i}^{\prime}(\widehat{B}) . \tag{2.2}
\end{equation*}
$$

Let $d_{i}(B, \Sigma)$ be the Mahalanobis norms of the residuals, given by

$$
d_{i}(B, \Sigma)=\left(\widehat{u}_{i}^{\prime}(B) \Sigma^{-1} \widehat{u}_{i}(B)\right)^{1 / 2}, \quad 1 \leqslant i \leqslant n
$$

If the $n \times(p+q)$ matrix $(X, Y)$ has rank $(p+q)$ then $\operatorname{det}(\widehat{\Sigma}) \neq 0$. In this case the MLE of $B$ and $\Sigma$ also satisfy

$$
\begin{equation*}
(\widehat{B}, \widehat{\Sigma})=\arg \min _{B, \Sigma} \operatorname{det}(\Sigma) \tag{2.3}
\end{equation*}
$$

subject to

$$
\begin{equation*}
\frac{1}{n} \sum_{i=1}^{n} d_{i}^{2}(B, \Sigma)=q \tag{2.4}
\end{equation*}
$$

This follows from two facts: (i) the MLE of $\left(B_{0}, \Sigma_{0}\right)$ minimizes $\left(n \log |\Sigma|+\sum_{i=1}^{n} d_{i}^{2}(B\right.$, $\Sigma)$ ) and (ii) (2.2) implies that the MLE satisfies (2.4).

Given a sample $v_{1}, \ldots, v_{n}$, let $s\left(v_{1}, \ldots, v_{n}\right)$ be the scale estimate defined as the square root of the mean squared error (RMSE), that is

$$
\begin{equation*}
s\left(v_{1}, \ldots, v_{n}\right)=\left(\frac{1}{n} \sum_{i=1}^{n} v_{i}^{2}\right)^{1 / 2} \tag{2.5}
\end{equation*}
$$

Then constraint (2.4) can also be expressed as

$$
\begin{equation*}
s^{2}\left(d_{1}(B, \Sigma), \ldots, d_{n}(B, \Sigma)\right)=q \tag{2.6}
\end{equation*}
$$

One way of defining robust estimates of $B_{0}$ and of a scatter matrix $\Sigma_{0}$ of the $u_{i}$ 's is by (2.3) and (2.6), but replacing the scale $s$ defined in (2.5) by a robust scale.

Huber [9] introduced the M-estimates of scale. For a sample $v=\left(v_{1}, \ldots, v_{n}\right)$, an M-estimate of scale $s\left(v_{1}, \ldots, v_{n}\right)$ is defined by the value $s$ satisfying

$$
\begin{equation*}
\frac{1}{n} \sum_{i=1}^{n} \rho\left(\frac{\left|v_{i}\right|}{s}\right)=\kappa \tag{2.7}
\end{equation*}
$$

When $v_{1}, \ldots, v_{n}$ is a random sample of any arbitrary distribution $H, s\left(v_{1}, \ldots, v_{n}\right)$ converges to $s_{0}$ defined by $\kappa=E_{H}\left(\rho\left(|v| / s_{0}\right)\right)$. If we wish to calibrate the M -scale so that $s_{0}=1$, we should take $\kappa=E_{H}(\rho(|v|))$.

The function $\rho$ should satisfy the following properties:
(A1) $\rho(0)=0$.
(A2) $0 \leqslant v \leqslant v^{*}$ implies $\rho(v) \leqslant \rho\left(v^{*}\right)$.
(A3) $\rho$ is continuous.
(A4) $0<A=\sup _{v} \rho(v)<\infty$.
(A5) If $\rho(u)<A$ and $0 \leqslant u<v$, then $\rho(u)<\rho(v)$.
One measure of the degree of robustness of an estimate is the breakdown point introduced by Hampel [6]. Roughly speaking, the breakdown point of an estimate is the smallest fraction of outliers that is required to take the estimate to an extreme value. In the case of scale estimates, the two possible extreme values are zero or infinity ("implosion" and "explosion" breakdown, respectively). Huber [9] proved that the breakdown point to infinity of an M-estimate of scale is $\varepsilon_{\infty}^{*}=\kappa / A$ and the breakdown point to zero is $\varepsilon_{0}^{*}=1-\kappa / A$. Then, the asymptotic breakdown point of this scale estimate is $\varepsilon^{*}=\min (\kappa / A, 1-\kappa / A)$, and taking $\kappa / A=0.5$ we obtain the highest possible breakdown point, $\varepsilon^{*}=0.5$.

Bilodeau and Duchesne [2] proposed a class of robust estimates, called S-estimates, for the seemingly unrelated equations (SUE) model. The SUE model is more general than the MLM, since it allows different regressor vectors for each of the equations. When applied to the MLM, the S-estimates are defined by (2.3) and (2.6) where $s$ is an M-estimate of scale.

Croux [3] proved that M-estimates of scale can combine a high breakdown point with high efficiency under normality. However, Hossjer [8] showed that estimates of regression based on an M -scale cannot combine both properties. To overcome this limitation, Yohai and Zamar [17] proposed $\tau$-estimates of scale, defined as follows: consider two functions $\rho_{1}$ and $\rho_{2}$ satisfying A1-A5 and an arbitrary distribution $H$, and put

$$
\begin{equation*}
\kappa_{i}=E_{H}\left(\rho_{i}(|v|)\right), \quad i=1,2 . \tag{2.8}
\end{equation*}
$$

Let $s(v)$ be the M-estimate of scale defined in (2.7) with $\rho=\rho_{1}$ and $\kappa=\kappa_{1}$. Below we specify which $H$ will be used to define $\kappa_{1}$. The $\tau$-estimate of the scale of $v=\left(v_{1}, \ldots, v_{n}\right)$ is defined by

$$
\begin{equation*}
\tau^{2}(v)=s^{2}(v) \frac{1}{n} \sum_{i=1}^{n} \rho_{2}\left(\frac{\left|v_{i}\right|}{s(v)}\right) \tag{2.9}
\end{equation*}
$$

The estimate $\tau(v)$ converges to $\kappa_{2}$ when the $v_{i}$ are independent variables with distribution $H$.
Yohai and Zamar [17] defined $\tau$-estimates of univariate regression by the minimization of a $\tau$-scale of the residuals. Put $\psi_{i}(v)=\rho_{i}^{\prime}(v), i=1,2$. To guarantee the Fisher consistency of the $\tau$-estimates of regression, it is required that $\rho_{2}$ satisfy the following condition:
(A6) $\rho_{2}$ is continuously differentiable and $2 \rho_{2}(v)-\psi_{2}(v) v>0$ for $v>0$.
Yohai and Zamar [17] also showed that by properly choosing $\rho_{1}$ and $\rho_{2}$, the $\tau$-estimates of regression combine breakdown point 0.5 with a high Gaussian asymptotic efficiency.

Here we extend $\tau$-estimates to the MLM model by defining

$$
\begin{equation*}
(\widetilde{B}, \widetilde{\Sigma})=\arg \min _{B, \Sigma} \operatorname{det}(\Sigma) \tag{2.10}
\end{equation*}
$$

subject to

$$
\begin{equation*}
\tau^{2}\left(d_{1}(B, \Sigma), \ldots, d_{n}(B, \Sigma)\right)=\kappa_{2}, \tag{2.11}
\end{equation*}
$$

where the scale $\tau$ is defined in (2.9). The values of $\kappa_{1}, \kappa_{2}$ are chosen using (2.8) with $H=$ $H_{00}$, the distribution of $\left(u_{i}^{\prime} \Sigma_{0}^{-1} u_{i}\right)^{1 / 2}$ under a nominal distribution. Let $H_{0}$ be the distribution of
$\left(u_{i}^{\prime} \Sigma_{0}^{-1} u_{i}\right)^{1 / 2}$. Since $H_{0}$ is unknown, $H_{00}$ may not coincide with $H_{0}$. To prevent $H_{00}$ from being dependent on $\Sigma_{0}$, the nominal model should be elliptical, for example the multivariate normal. In this case $H_{00}$ is the distribution of $\sqrt{v}$, where $v$ has a chi-square distribution with $q$ degrees of
 $\widetilde{B}$ converges to $B_{0}$ and $\widetilde{\Sigma}$ to $\lambda \Sigma_{0}$, even if $H_{00} \neq H_{0}$. If $H_{00}=H_{0}$, we have $\lambda=1$.

Let $h$ be the maximum number of observations $\left(x_{i}, y_{i}\right)$ lying in a hyperplane, i.e.,

$$
\begin{equation*}
h=\max _{\|a\|+\|b\|>0} \#\left\{i: a^{\prime} x_{i}+b^{\prime} y_{i}=0\right\} \tag{2.12}
\end{equation*}
$$

and

$$
\begin{equation*}
\eta=\frac{\kappa_{1}}{\max _{u} \rho_{1}(u)} \tag{2.13}
\end{equation*}
$$

It can be shown that if $\rho_{1}$ and $\rho_{2}$ satisfy A1-A5 and $h / n<1-\eta$, then there exists a $(\widetilde{B}, \widetilde{\Sigma})$ solution to the optimization problem defined by (2.10) and (2.11) with $|\widetilde{\Sigma}|>0$. The proof of this is similar to the proofs of Lemmas 4 and 5 of the Appendix. As we will see in Section 3, to obtain a high breakdown point $\tau$-estimate, $\eta$ should be close to 0.5 . If there are multiple solutions we take any of them.

Comparing (2.10) and (2.11) with (2.3) and (2.6), we observe that the MLE may be considered as a $\tau$-estimate corresponding to the unbounded function $\rho_{2}(v)=v^{2}$ and $\kappa_{2}=q$.

It is easy to show that when there is only one regressor variable identically equal to one, the $\tau$-estimates $\widetilde{B}$ and $\widetilde{\Sigma}$ defined here are the estimates of multivariate location and scatter introduced by Lopuhaä [11]. When $q=1$, they are the $\tau$-estimates of regression proposed by Yohai and Zamar [17].

It can be verified that $\tau$-estimates are affine and regression equivariant.
In Theorem 1 we obtain the estimating equations of $\tau$-estimates. Define

$$
d_{i}^{*}(B, \Sigma)=\frac{d_{i}(B, \Sigma)}{s\left(d_{1}(B, \Sigma), \ldots, d_{n}(B, \Sigma)\right)}
$$

and let $C_{n}=C_{n}(B, \Sigma), D_{n}=D_{n}(B, \Sigma), \psi_{n}^{*}=\psi_{n, B, \Sigma}^{*}$ and $w_{n}^{*}$ be defined by

$$
\begin{align*}
& C_{n}=\frac{1}{n} \sum_{i=1}^{n}\left(2 \rho_{2}\left(d_{i}^{*}(B, \Sigma)\right)-\psi_{2}\left(d_{i}^{*}(B, \Sigma)\right) d_{i}^{*}(B, \Sigma)\right) \\
& D_{n}=\frac{1}{n} \sum_{i=1}^{n} \psi_{1}\left(d_{i}^{*}(B, \Sigma)\right) d_{i}^{*}(B, \Sigma), \\
& \psi_{n}^{*}(v)=C_{n} \psi_{1}(v)+D_{n} \psi_{2}(v) \tag{2.14}
\end{align*}
$$

and

$$
\begin{equation*}
w_{n}^{*}(v)=\frac{\psi_{n}^{*}(v)}{v} \tag{2.15}
\end{equation*}
$$

Theorem 1. Suppose that $\rho_{1}$ and $\rho_{2}$ are differentiable. Then the $\tau$-estimates satisfy the following equations:

$$
\begin{equation*}
\sum_{i=1}^{n} w_{n}^{*}\left(d_{i}^{*}(\widetilde{B}, \widetilde{\Sigma})\right) \widehat{u}_{i}(\widetilde{B}) x_{i}^{\prime}=0 \tag{2.16}
\end{equation*}
$$

and

$$
\begin{equation*}
\widetilde{\Sigma}=\frac{q \sum_{i=1}^{n} w_{n}^{*}\left(d_{i}^{*}(\widetilde{B}, \widetilde{\Sigma})\right) \widehat{u}_{i}(\widetilde{B}) \widehat{u}_{i}^{\prime}(\widetilde{B})}{\widetilde{s}^{2} \sum_{i=1}^{n} \psi_{n}^{*}\left(d_{i}^{*}(\widetilde{B}, \widetilde{\Sigma})\right) d_{i}^{*}(\widetilde{B}, \widetilde{\Sigma})}, \tag{2.17}
\end{equation*}
$$

where $\widehat{u}_{i}(B)$ is defined in (2.1) and $\widetilde{s}=s\left(d_{1}(\widetilde{B}, \widetilde{\Sigma}), \ldots, d_{n}(\widetilde{B}, \widetilde{\Sigma})\right)$.
The proof of this theorem can be found in the Appendix.
Remark. Observe that according to (2.16), the $j$ th column of $\widetilde{B}$ is the weighted LSE corresponding to the univariate regression whose dependent variable is the $j$ th component of $y$, the vector of independent variables is $x$ and the observation $i$ receives weight $w_{n}^{*}\left(d_{i}^{*}(\widetilde{B}, \widetilde{\Sigma})\right.$ ). Besides, by (2.17), $\widetilde{\Sigma}$ is proportional to the weighted sample covariance of the residuals with the same weights. Since these weights depend on the estimates $\widetilde{B}$ and $\widetilde{\Sigma}$, we cannot use these relationships to compute the estimates, but they will be the base of the iterative algorithm described in Section 6.

## 3. Breakdown point

Donoho and Huber [5] introduced the concept of a finite sample breakdown point (FSBDP). For the case of the MLM, let $\hat{B}$ and $\hat{\Sigma}$ be estimates of $B$ and $\Sigma$. The FSBDP of $\hat{B}$ is the smallest fraction of outliers that makes this estimate unbounded, and the FSBDP of $\hat{\Sigma}$ the smallest fraction of outliers that makes this estimate unbounded or singular. To formalize this, let $Z$ be a data set of size $n$ corresponding to an MLM, $Z=\left\{z_{1}, \ldots, z_{n}\right\}, z_{i}=\left(x_{i}^{\prime}, y_{i}^{\prime}\right)^{\prime}, x_{i} \in R^{p}, y_{i} \in R^{q}$. Let $\mathcal{Z}_{m}$ be the set of all the samples $Z^{*}=\left\{z_{1}^{*}, \ldots, z_{n}^{*}\right\}$ such that $\#\left\{i: z_{i}=z_{i}^{*}\right\} \geqslant n-m$. Given estimates $\hat{B}$ and $\hat{\Sigma}$, denote by $\hat{\lambda}_{1} \leqslant \cdots \leqslant \hat{\lambda}_{q}$ the eigenvalues of $\hat{\Sigma}$,

$$
S_{m}(Z, \hat{B})=\sup \left\{\left\|\widehat{B}\left(Z^{*}\right)\right\|, Z^{*} \in \mathcal{Z}_{m}\right\}
$$

where $\left\|\|\right.$ is the $l_{2}$ norm, and let

$$
S_{m}^{-}(Z, \hat{\Sigma})=\inf \left\{\widehat{\lambda}_{1}\left(\widehat{\Sigma}\left(Z^{*}\right)\right), Z^{*} \in \mathcal{Z}_{m}\right\}
$$

and

$$
S_{m}^{+}(Z, \hat{\Sigma})=\sup \left\{\widehat{\lambda}_{q}\left(\widehat{\Sigma}\left(Z^{*}\right)\right), Z^{*} \in \mathcal{Z}_{m}\right\}
$$

Definition. The finite sample breakdown point of $\hat{B}$ is defined by $\varepsilon^{*}(Z, \hat{B})=m^{*} / n$ where $m^{*}=\min \left\{m: S_{m}(Z, \hat{B})=\infty\right\}$ and the breakdown point of $\hat{\Sigma}$ by $\varepsilon^{*}(Z, \hat{\Sigma})=m^{*} / n$ where

$$
m^{*}=\min \left\{m: \frac{1}{S_{m}^{-}(Z, \hat{\Sigma})}+S_{m}^{+}(Z, \hat{\Sigma})=\infty\right\}
$$

The following theorem, whose proof can be found in the Appendix, gives a lower bound for the breakdown point of $\tau$-estimates.

Theorem 2. Let $Z=\left\{z_{1}, z_{2}, \ldots, z_{n}\right\}$ with $z_{i}=\left(x_{i}^{\prime}, y_{i}^{\prime}\right)^{\prime}$, has defined in (2.12), and $\eta$ as defined in (2.13). Assume that $\rho_{1}$ and $\rho_{2}$ satisfy A1-A5; then a lower bound for $\varepsilon^{*}(Z, \widetilde{\Sigma})$ and $\varepsilon^{*}(Z, \widetilde{B})$ is given by $\min ((1-\eta)-(h / n), \eta)$.

The maximum breakdown point is $(n-h) /(2 n)$ and corresponds to $\eta=(n-h) /(2 n)$. The minimum value of $h$ is $p+q-1$. If $h=p+q-1$, we say that the points are in a general position. If $\eta=(n-h) /(2 n)$ and $h=p+q-1$, we obtain that $\varepsilon^{*}(Z, \widetilde{\Sigma})$ and $\varepsilon^{*}(Z, \widetilde{B})$ are larger than or equal to $0.5-((p+q-1) /(2 n))$, and then in this case the two breakdown points are close to 0.5 for large $n$.

## 4. Influence function

Consider a random variable $z \in R^{k}$ with distribution $G_{\theta}$, where $\theta \in \Theta \subset R^{m}$. Let $T$ be an estimating functional of $\theta$, i.e., $T$ is defined on a set of distributions $G$ on $R^{k}$, including the empirical distributions and the contamination neighborhoods of $G_{\theta}$, and takes values on $\Theta$. Suppose that $T$ is Fisher consistent, i.e. $T\left(G_{\theta}\right)=\theta$. Given a sample $z_{1}, z_{2}, \ldots, z_{n}$ of $G_{\theta}$, the estimate $\widehat{\theta}$ associated to $T$ is defined as $T\left(G_{n}\right)$, where $G_{n}$ is the corresponding empirical distribution. The influence function of $T$ measures the effect on the functional of a small fraction of point mass contamination ([7]). More precisely, let $\delta_{z}$ be the point mass distribution at $z$ and $G_{\theta, \varepsilon}=(1-\varepsilon) G_{\theta}+\varepsilon \delta_{z}$, then the influence function is defined by

$$
\operatorname{IF}(z, T, \theta)=\lim _{\varepsilon \rightarrow 0} \frac{T\left(G_{\theta, \varepsilon}\right)-T\left(G_{\theta}\right)}{\varepsilon}=\left.\frac{\partial}{\partial \varepsilon} T\left(G_{\theta, \varepsilon}\right)\right|_{\varepsilon=0}
$$

In our case, $z=\left(x^{\prime}, y^{\prime}\right)^{\prime}$ and $\theta=\left(B_{0}, \Sigma_{0}\right)$. The estimating functionals $T_{1}$ and $T_{2}$ corresponding to the $\tau$-estimates $\widetilde{B}$ and $\widetilde{\Sigma}$ are given in the Appendix.

For the sake of simplicity we derive the influence function of $\tau$-estimates assuming that the errors in (1.1) have an elliptical distribution with a unimodal density. Then, we need the following assumption:
(A7) The distribution $F_{0}$ of $u_{i}$ has a density of the form

$$
\begin{equation*}
f_{0}(u)=\frac{f_{0}^{*}\left(u^{\prime} \Sigma_{0}^{-1} u\right)}{\operatorname{det}\left(\Sigma_{0}\right)^{1 / 2}} \tag{4.1}
\end{equation*}
$$

where $f_{0}^{*}$ is nonincreasing and has at least one point of decrease in the interval where both functions $\rho_{1}$ and $\rho_{2}$ are strictly increasing. Besides, $\Sigma_{0}$ is a $q \times q$ positive-definite matrix.

An important family of unimodal elliptical distributions is the multivariate normal. In this case, $f_{0}^{*}(v)=\left(1 /(2 \pi)^{q / 2}\right) \exp (-v / 2)$.

Given an arbitrary distribution $H$ on $R$, define $s^{*}(H)$ by

$$
\begin{equation*}
E_{H}\left(\rho_{1}\left(\frac{v}{s^{*}(H)}\right)\right)=\kappa_{1} \tag{4.2}
\end{equation*}
$$

Observe that the M-estimate of scale $s\left(v_{1}, \ldots, v_{n}\right)=s^{*}\left(H_{n}\right)$, where $H_{n}$ is the empirical distribution of $v_{1}, \ldots, v_{n}$. We denote

$$
\begin{equation*}
k_{0}=s^{*}\left(H_{0}\right), \tag{4.3}
\end{equation*}
$$

where $H_{0}$ is the true distribution of $\left(u^{\prime} \Sigma_{0}^{-1} u\right)^{1 / 2}$.

To define the limit of the sequence of the functions $\psi_{n}^{*}$ and $w_{n}^{*}$, given in (2.14) and (2.15), put

$$
\begin{aligned}
C & =E_{H_{0}}\left(2 \rho_{2}\left(\frac{v}{k_{0}}\right)-\psi_{2}\left(\frac{v}{k_{0}}\right)\left(\frac{v}{k_{0}}\right)\right), \\
D & =E_{H_{0}}\left(\psi_{1}\left(\frac{v}{k_{0}}\right)\left(\frac{v}{k_{0}}\right)\right) .
\end{aligned}
$$

Then, define $\psi^{*}(v)=C \psi_{1}(v)+D \psi_{2}(v)$ and $w^{*}(v)=\psi^{*}(v) / v$.
In the following we shall need some additional assumptions:
(A8) The $x_{i}$ 's are random, their common distribution $M_{0}$ has finite second moments and $E_{M_{0}}\left(x x^{\prime}\right)$ is nonsingular.
(A9) $\rho_{i}$ is twice differentiable $(i=1,2)$.
Theorem 3. Assume Al-A9. Then, the influence function for the estimating functional $T_{1}$ corresponding to the $\tau$-estimate $\widetilde{B}$ is

$$
\begin{align*}
& I F\left(y_{0}, x_{0}, T_{1}, B_{0}, \Sigma_{0}\right) \\
& \quad=c_{0} w^{*}\left(\frac{\left(\left(y_{0}-B_{0}^{\prime} x_{0}\right)^{\prime} \Sigma_{0}^{-1}\left(y_{0}-B_{0}^{\prime} x_{0}\right)\right)^{1 / 2}}{k_{0}}\right) E_{M_{0}}\left(x x^{\prime}\right)^{-1} x_{0}\left(y_{0}-B_{0}^{\prime} x_{0}\right)^{\prime}, \tag{4.4}
\end{align*}
$$

where

$$
\begin{equation*}
c_{0}=\frac{q}{E_{H_{0}}\left((q-1) w^{*}\left(\frac{v}{k_{0}}\right)+\psi^{* \prime}\left(\frac{v}{k_{0}}\right)\right)} . \tag{4.5}
\end{equation*}
$$

The proof of this theorem is given in the Appendix.

## 5. Consistency and asymptotic normality

In this section we state the asymptotic properties of $\tau$-estimates for the MLM assuming that the errors have an elliptical distribution. Theorems 4 and 5 establish, respectively, the consistency and the asymptotic normality of $\widetilde{B}$.

Let $\sigma_{0}=\sigma_{0}\left(\psi_{1}, \psi_{2}\right)$ be defined by

$$
\begin{equation*}
\frac{\sigma_{0}^{2}}{\kappa_{2}} E_{H_{0}}\left(\rho_{2}\left(\frac{v}{k_{0}}\right)\right)=1 \tag{5.1}
\end{equation*}
$$

where $k_{0}$ is given in (4.3). If $H_{0}=H_{00}$ we have $\sigma_{0}=k_{0}=1$.
Theorem 4. Let $\left(x_{i}, y_{i}\right), 1 \leqslant i \leqslant n$, be a random sample of an MLM with parameters $B_{0}$ and $\Sigma_{0}$, where the $x_{i}$ 's are random. Suppose also that A1-A7 hold; then the $\tau$-estimates $\widetilde{B}_{n}$ and $\widetilde{\Sigma}_{n}$ satisfy
(a) $\lim _{n \rightarrow \infty} \widetilde{\sim}_{n}=B_{0}$ a.s..
(b) $\lim _{n \rightarrow \infty} \widetilde{\Sigma}_{n}=\left(k_{0}^{2} / \sigma_{0}^{2}\right) \Sigma_{0}$ a.s.. In particular if $H_{00}=H_{0}$, then $k_{0}^{2} / \sigma_{0}^{2}=1$.

To prove the asymptotic normality of $\widetilde{B}_{n}$ we need the following additional assumption: (A10) There exist $m_{i}(i=1,2)$, such that $\rho_{i}(u)$ is constant for $|u|>m_{i}$.

Theorem 5. Let $\left(x_{i} y_{i}\right), 1 \leqslant i \leqslant n$, be a random sample of an MLM with parameters $B_{0}$ and $\Sigma_{0}$. Suppose also that A1-A10 hold; then $n^{1 / 2} \operatorname{vec}\left(\widetilde{B}_{n}-B_{0}\right) \rightarrow^{D} N(0, V)$, where $\rightarrow^{D}$ denotes convergence in distribution and

$$
\begin{equation*}
V=\frac{c_{0}^{2} k_{0}^{2}}{q} E_{H_{0}}\left(\psi^{* 2}\left(\frac{v}{k_{0}}\right)\right) \Sigma \otimes E_{M_{0}}\left(x x^{\prime}\right)^{-1} \tag{5.2}
\end{equation*}
$$

with $k_{0}$ and $c_{0}$ given in (4.3) and (4.5), respectively.
We do not have formal proofs of Theorems 4 and 5, but we think that they can be obtained using arguments similar to those in Yohai and Zamar [16,17], and Lopuhaä [1]. Heuristic proofs of both Theorems can be found in the Appendix, where we also give a rigorous proof of the Fisher consistency of these estimates (Lemma 12).

Remark. Estimates of $V$ can be obtained replacing in (5.2) $H_{0}$ and $M_{o}$ by the empirical distribution of the residuals and of the $x$ 's, respectively.

The MLE estimator $\widehat{B}_{n}$ corresponds to $\rho_{2}(u)=u^{2}$, and in this case $c_{0}=1 / 2, \psi^{*}(u)=2 u$ and the asymptotic covariance matrix is given by $\left(E_{H_{0}}\left(v^{2}\right) / q\right) \Sigma \otimes E_{M_{0}}\left(x x^{\prime}\right)^{-1}$ which differs from the one of the $\tau$-estimate by a scalar factor. Then we obtain the following result:

Corollary. Under the assumptions of Theorem 5, the asymptotic relative efficiency (ARE) of the $\tau$-estimate $\widetilde{B}$ with respect to the $M L E \widehat{B}$ is

$$
\operatorname{ARE}\left(\psi_{1}, \psi_{2}, H_{0}\right)=\frac{E_{H_{0}}\left(v^{2}\right)}{c_{0}^{2} k_{0}^{2} E_{H_{0}}\left(\psi^{* 2}\left(\frac{v}{k_{0}}\right)\right)}
$$

In order to obtain a $\tau$-estimate which is simultaneously highly robust and highly efficient under normal errors we can choose $\rho_{1}$ so that

$$
\begin{equation*}
\kappa_{1} / \max \rho_{1}=0.5, \tag{5.3}
\end{equation*}
$$

which guarantees that the initial M-estimate of scale $s$ has a breakdown point close to 0.5 , and choose as $\rho_{2}(v)$ a bounded function close enough to $v^{2}$ to obtain the desired efficiency. This can be achieved, for example, by taking $\rho_{1}$ and $\rho_{2}$ in Tukey's bisquare family defined by

$$
\rho_{B, c}(v)= \begin{cases}\frac{v^{2}}{2}\left(1-\frac{v^{2}}{c^{2}}+\frac{v^{4}}{3 c^{4}}\right) & \text { if }|v| \leqslant c \\ \frac{c^{2}}{6} & \text { if }|v|>c\end{cases}
$$

where $c$ is any positive number. Observe that when $c$ increases, $\rho_{B, c}$ approaches $v^{2}$. It is easy to verify that the functions $\rho_{B, c}(v)$ satisfy A1-A6.

Table 1 gives the values of $c_{1}$ and $\kappa_{1}$ such that $\rho_{1}=\rho_{B, c_{1}}$ satisfies (5.3) and the ARE of $S$-estimates under Gaussian errors for different values of $q$. We observe that the efficiency of the S-estimate is low for small values of $q$, but increases with $q$. It may be shown that this efficiency converges to one when $q \rightarrow \infty$. Table 2 gives the values of $c_{2}$ and $\kappa_{2}$ to achieve different levels of asymptotic efficiency (taking $\rho_{1}=\rho_{B, c_{1}}$ and $\rho_{2}=\rho_{B, c_{2}}$ ).

Table 1
Values of $c_{1}$ and $\kappa_{1}$ for the bisquare function and asymptotic relative efficiency (ARE) under Gaussian errors of the S-estimate with breakdown point 0.5

| $q$ | 1 | 2 | 3 | 4 | 5 | 10 |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- |
| $c_{1}$ | 1.55 | 2.66 | 3.45 | 4.10 | 4.65 | 6.77 |
| $\kappa_{1}$ | 0.20 | 0.59 | 0.99 | 1.40 | 1.80 | 3.82 |
| ARE | 0.29 | 0.58 | 0.72 | 0.80 | 0.85 | 0.93 |

Table 2
Values of $c_{2}$ and $\kappa_{2}$ for the bisquare function to attain given values of the asymptotic relative efficiency (ARE) under Gaussian errors

| ARE | $q$ <br>  <br>  |  |  |  |  |  |  |
| :--- | :--- | :--- | :--- | :--- | :--- | :--- | :--- |
|  | 1 | 2 | 3 | 4 | 5 | 10 |  |
| 0.80 | $c_{2}$ | 3.98 | 3.94 | 4.02 | 4.10 | 4.17 | 4.28 |
|  | $\kappa_{2}$ | 0.42 | 0.77 | 1.10 | 1.40 | 1.67 | 2.56 |
| 0.90 |  |  |  |  |  |  |  |
|  | $c_{2}$ | 4.97 | 4.97 | 5.10 | 5.25 | 5.39 | 5.98 |
|  | $\kappa_{2}$ | 0.44 | 0.85 | 1.24 | 1.61 | 1.96 | 3.54 |
| 0.95 |  |  |  |  |  |  |  |
|  | $c_{2}$ | 6.04 | 6.06 | 6.24 | 6.42 | 6.60 | 7.50 |
|  | $\kappa_{2}$ | 0.46 | 0.90 | 1.32 | 1.73 | 2.13 | 4.03 |

## 6. Computing algorithm

Based on the remark at the end of Section 2 we propose the following iterative algorithm to compute $\widetilde{B}$ and $\widetilde{\Sigma}$.

1. Using initial values $\widetilde{B}_{0}$ and $\tilde{\Sigma}_{0}$ satisfying (2.11), compute $\widetilde{s}_{0}=s\left(d_{1}\left(\widetilde{B}_{0}, \widetilde{\Sigma}_{0}\right), \ldots, d_{n}\left(\widetilde{B}_{0}\right.\right.$, $\left.\widetilde{\Sigma}_{0}\right)$ ) and the weights $w_{n}^{*}\left(d_{i}\left(\widetilde{B}_{0}, \widetilde{\Sigma}_{0}\right) / \widetilde{s}_{0}\right)$ for $1 \leqslant i \leqslant n$. These weights are used to compute each column of $\widetilde{B}_{1}$ separately by WLS. Now compute $\widetilde{s}_{1}=s\left(d_{1}\left(\widetilde{B}_{1}, \widetilde{\Sigma}_{0}\right), \ldots, d_{n}\left(\widetilde{B}_{1}, \widetilde{\Sigma}_{0}\right)\right)$.
2. Compute the matrix

$$
\begin{equation*}
\widetilde{\Sigma}_{1}^{*}=\frac{q \sum_{i=1}^{n} w_{n}^{*}\left(\frac{d_{i}\left(\widetilde{B}_{1}, \widetilde{\Sigma}_{0}\right)}{\widetilde{s}_{1}}\right) u_{i}\left(\widetilde{B}_{1}\right) u_{i}^{\prime}\left(\widetilde{B}_{1}\right)}{\widetilde{s}_{1}^{2} \sum_{i=1}^{n} \psi_{n}^{*}\left(\frac{d_{i}\left(\widetilde{B}_{1}, \widetilde{\Sigma}_{0}\right)}{\widetilde{s}_{1}}\right) \frac{d_{i}\left(\widetilde{B}_{1}, \widetilde{\Sigma}_{0}\right)}{\widetilde{s}_{1}}} \tag{6.1}
\end{equation*}
$$

3. Compute $\widetilde{\tau}_{1}=\tau\left(d_{1}\left(\widetilde{B}_{1}, \widetilde{\Sigma}_{1}^{*}\right), \ldots, d_{n}\left(\widetilde{B}_{1}, \widetilde{\Sigma}_{1}^{*}\right)\right)$ and $\widetilde{\Sigma}_{1}=\left(\widetilde{\tau}_{1}^{2} / \kappa_{2}\right) \widetilde{\Sigma}_{1}^{*}$. Then $\left(\widetilde{B}_{1}, \widetilde{\Sigma}_{1}\right)$ satisfy constraint (2.11).
4. Suppose now that we have already computed ( $\widetilde{B}_{h}, \widetilde{\Sigma}_{h}$ ) satisfying constraint (2.11). Then $\left(\widetilde{B}_{h+1}, \widetilde{\Sigma}_{h+1}\right)$ is computed using steps $1-3$, but starting from $\left(\widetilde{B}_{h}, \widetilde{\Sigma}_{h}\right)$ instead of $\left(\widetilde{B}_{0}, \widetilde{\Sigma}_{0}\right)$.
5. The procedure is stopped at step $h$ if the relative absolute differences of all elements of the matrices $\widetilde{B}_{h}$ and $\widetilde{B}_{h-1}$ are smaller than a given value $\delta$.
We have not proved that the reweighting step improves the value of the goal function. However in the Monte Carlo study described in Section 7 this has always occurred.

We propose to compute the initial estimates $\widetilde{B}_{0}$ and $\widetilde{\Sigma}_{0}$ by subsampling elemental sets. For this purpose we take $N$ random subsamples of size $r=p+q$ of the original sample. For the $j$ th subsample two values of $(B, \Sigma)$ are obtained. The first $\left(B_{j}^{(1)}, \Sigma_{j}^{(1)}\right)$ corresponds to the MLE of the
subsample, and the second value $\left(B_{j}^{(2)}, \Sigma_{j}^{(2)}\right)$ is the MLE of the $[n / 2]$ observations with the smallest Mahalanobis distances $d_{i}\left(B_{j}^{(1)}, \Sigma_{j}^{(1)}\right), 1 \leqslant i \leqslant n$. We now compute $d_{i}\left(B_{j}^{(2)}, \Sigma_{j}^{(2)}\right), 1 \leqslant i \leqslant n$, and $s_{j}=s_{0}\left(d_{1}\left(B_{j}^{(2)}, \Sigma_{j}^{(2)}\right), \ldots, d_{n}\left(B_{j}^{(2)}, \Sigma_{j}^{(2)}\right)\right)$, where $s_{0}\left(u_{1}, \ldots, u_{n}\right)=\operatorname{median}\left(|u|_{1}, \ldots,|u|_{n}\right)$ and standardize $\Sigma_{j}^{(2)}$ obtaining $\Sigma_{j}^{(3)}=s_{j}^{2} \Sigma_{j}^{(2)}$. Then $s_{0}^{2}\left(d_{1}\left(B_{j}^{(2)}, \Sigma_{j}^{(3)}\right), \ldots, . d_{n}\left(B_{j}^{(2)}, \Sigma_{j}^{(3)}\right)\right)=1$.

An approximation to the estimate that minimizes $\operatorname{det}(\Sigma)$ subject to $s_{0}^{2}\left(d_{1}(B, \Sigma), \ldots, d_{n}\right.$ $(B, \Sigma))=1$ is given by $\left(B_{j_{0}}^{(2)}, \Sigma_{j_{0}}^{(3)}\right)$ where $j_{0}=\arg \min _{1 \leqslant j \leqslant N} \operatorname{det}\left(\Sigma_{j}^{(3)}\right)$.
We use the scale $s_{0}$ instead of the scale $\tau$ because it is faster to compute. This scale is inefficient, but since it is used only to compute the initial estimate, it does not affect the efficiency of the final estimate.

Finally, the initial estimate is obtained by restandardizing $\Sigma_{j_{0}}^{(3)}$ so that the $\tau$-scale of the Mahalanobis distances is $\sqrt{\kappa_{2}}$. For this purpose we compute

$$
\tau_{j_{0}}=\tau\left(d_{1}\left(B_{j_{0}}^{(2)}, \Sigma_{j_{0}}^{(3)}\right), \ldots, . d_{n}\left(B_{j_{0}}^{(2)}, \Sigma_{j_{0}}^{(3)}\right)\right)
$$

and the initial estimates are $\widetilde{B}_{0}=B_{j_{0}}^{(2)}$ and $\widetilde{\Sigma}_{0}=\left(\tau_{j_{0}}^{2} / \kappa_{2}\right) \Sigma_{j_{0}}^{(3)}$.
The reason why we compute the second value $\left(B_{j}^{(2)}, \Sigma_{j}^{(2)}\right)$ is that even if the $j$ th sample does not contain outliers, it may be badly conditioned and the corresponding fit $B_{j}^{(1)}$ may be very far from the true value. However, eliminating the sample half with largest $d_{i}\left(B_{j}^{(1)}, \Sigma_{j}^{(1)}\right)$ 's increases the chance of obtaining a clean sample that produces a better value, $B_{j}^{(2)}$. This mechanism is similar to the one proposed by Rousseeuw and Van Driessen [14].

One improvement suggested by a referee would be to proceed as in Rousseeuw and Van Driessen [14] keeping the $M$ solutions $\left(B_{j}^{(2)}, \Sigma_{j}^{(3)}\right.$ ) with smallest $\left|\Sigma_{j}^{(3)}\right|$ (for example $M=10$ ) and starting the iterative process from each one of them. This gives $M$ new values ( $\left.\widetilde{B}_{h}, \widetilde{\Sigma}_{h}\right), 1 \leqslant h \leqslant M$, and the final estimate is the one with the smallest $\left|\widetilde{\Sigma}_{h}\right|$.

## 7. Monte Carlo results

In order to assess the robustness and efficiency of the proposed estimates we performed a Monte Carlo study. We consider the MLM given by (1.1) for two cases: $p=2, q=2$ and $p=2$, $q=5$. Due to the equivariance of the estimators we take, without loss of generality, $B_{0}=0$ and $\Sigma_{0}=I_{q}$, where $I_{q}$ denotes a $q \times q$ identity matrix. The errors $u_{i}$ are generated from an $\mathrm{N}_{q}\left(0, I_{q}\right)$ distribution and the regressors $x_{i}$ from an $\mathrm{N}_{p}\left(0, I_{p}\right)$ distribution. The sample size is 100 and the number of replications is 1000 . We consider uncontaminated samples and samples that contain $10 \%$ of identical outliers of the form $(x, y)$ with $x^{\prime}=\left(x_{0}, 0, \ldots, 0\right)$ and $y^{\prime}=\left(m x_{0}, 0, \ldots, 0\right)$. The values of $x_{0}$ considered are 1 (low leverage outliers) and 10 (high leverage outliers). We take a grid of values of $m$, starting at 0 . The last value of the grid was taken so that the maximum mean square error (MSE) of all the robust estimates is attained. Suppose that $\widehat{B}^{(k)}=\left(\widehat{B}_{i j}^{(k)}\right)$ is the estimate of $B_{0}$ obtained in the $k$ th replication. Then, since we are taking $B_{0}=0$, the estimate of the MSE is given by

$$
\operatorname{MSE}=\frac{1}{1000}\left(\sum_{k=1}^{1000} \sum_{i=1}^{p} \sum_{j=1}^{q}\left(\widehat{B}_{i j}^{(k)}\right)^{2}\right) .
$$

Table 3
Monte Carlo mean square error (MSE), standard error of the MSE (SE) and relative efficiency (REFF) of the estimates in the noncontaminated case for $n=100$ and $p=2$

| Estimate | $q=2$ |  |  |  | $q=5$ |  |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | MSE | SE | REFF | ARE | MSE | SE | REFF | ARE |
| MLE | 0.064 | 0.001 | 1.00 | 1.00 | 0.157 | 0.002 | 1.00 | 1.00 |
| S-estimate | 0.120 | 0.002 | 0.53 | 0.58 | 0.190 | 0.002 | 0.83 | 0.85 |
| $\tau$-estimate | 0.072 | 0.001 | 0.89 | 0.90 | 0.176 | 0.002 | 0.89 | 0.90 |

ARE is the asymptotic relative efficiency.


Fig. 1. Mean square error for $q=2$ and $x_{0}=1$.

For each case, three estimates are computed: the MLE, an S-estimate and a $\tau$-estimate. The S - and the $\tau$-estimates are based on $\rho$-functions in the bisquare family. The S -estimate is based on the M -scale defined by $\rho_{1}=\rho_{B, c_{1}}$, where $c_{1}$ and $\kappa_{1}$ were chosen so that this estimate has a breakdown point 0.5 (see Table 1). The $\tau$-estimate uses the same $\rho_{1}$ and $\kappa_{1}$ as the S-estimate and $\rho_{2}=\rho_{B, c_{2}}$, where $c_{2}$ and $\kappa_{2}$ were chosen so that the $\tau$-estimate had an ARE equal to 0.90 when the errors are Gaussian (see Table 2). We compute the initial estimate using 2000 subsamples and the value of $\delta$ in step 5 of the computing algorithm is taken equal to 0.0001 .

In Table 3 we present the MSE, its standard error (SE) and the relative efficiency (REFF) with respect to the MLE for the uncontaminated case. The efficiency of the S-estimate is low when $q=2$ and increases for $q=5$. We observe that the relative efficiencies of the S - and $\tau$-estimates are close to their asymptotic values.

In Figs. 1 and 2 we show the MSE of the different estimates under contamination. In Fig. 1, which corresponds to $q=2$ and $x_{0}=1$, we observe that the $\tau$-estimate has a smaller MSE than the S-estimate except when $m$ is (approximately) between 3 and 5 , and the maximum MSE is smaller for the $\tau$-estimate. As expected, the MSE of the MLE increases with $m$ reaching very large values. Fig. 2 shows the results for $q=2$ and $x_{0}=10$. S- and $\tau$-estimates behave similarly, with a small advantage for the $\tau$-estimate. Since for $q=5$ the behavior of S - and $\tau$-estimates is similar to the one observed for $q=2$, we do not report the results here.

As a general conclusion we can say that when there are no outliers, $\tau$-estimates are more efficient than S-estimates, and under outlier contamination, $\tau$-estimates behave better than or similar to S-estimates.


Fig. 2. Mean square error for $q=2$ and $x_{0}=10$.

## 8. Example

Observational studies have suggested that low dietary intake or low plasma concentrations of retinol, beta-carotene or other carotenoids might be associated with increased risk of developing certain types of cancer. Nieremberg et al. [12] studied the determinants of plasma concentrations of these micronutrients. In an unpublished study they have collected information on 14 variables that may be the determinants of the plasma levels of beta-carotene and retinol. The number of observations was 315 and the variables considered were: $Y_{1}$ BETAPLASMA: Plasma beta-carotene ( $\mathrm{ng} / \mathrm{ml}$ ), $Y_{2}$ RETPLASMA: Plasma Retinol ( $\mathrm{ng} / \mathrm{ml}$ ), $X_{1}$ AGE (years), $X_{2}$ SEX $(1=$ Male, $2=$ Female $), X_{3}$ SMOKSTAT: Smoking status $(1=$ Never $), X_{4}$ SMOKSTAT: Smoking status $\left(1=\right.$ Former), $X_{5}$ QUETELET (weight/(height $\left.{ }^{2}\right)$ ), $X_{6}$ VITUSE: Vitamin Use $(1=$ fairly often $), X_{7}$ VITUSE: Vitamin Use $(1=$ not often $), X_{8}$ CALORIES: Number of calories consumed per day, $X_{9}$ FAT: Grams of fat consumed per day, $X_{10}$ FIBER: Grams of fiber consumed per day, $X_{11}$ ALCOHOL: Number of alcoholic drinks consumed per week, $X_{12}$ CHOLESTEROL: Cholesterol consumed ( $\mathrm{mg} / \mathrm{day}$ ), $X_{13}$ BETADIET: Dietary beta-carotene consumed (mcg per day), $X_{14}$ RETDIET: Dietary retinol consumed (mcg/day). The data are available at http://lib.stat.cmu.edu/datasets/Plasma_Retinol.

We compute two estimates of the regression coefficients: the multivariate $\tau$-estimate and the MLE. The $\tau$-estimate uses $\rho_{1}$ and $\rho_{2}$ in the bisquare family with constants equal to those used in the Monte Carlo study of Section 7. In Fig. 3 we show, for both estimates, the box-plot of the Mahalanobis norms of the residuals $d_{i}=\left(\left(\mathbf{y}_{i}-\widehat{B} \mathbf{x}_{i}\right)^{\prime} \widehat{\Sigma}^{-1}\left(\mathbf{y}_{i}-\widehat{B} \mathbf{x}_{i}\right)\right)^{1 / 2}, 1 \leqslant i \leqslant 315$. If we declare outliers those observations such that $d_{i}>\sqrt{\chi_{2,0.99}^{2}}$, the $\tau$-estimate reveals 27 outliers while the MLE reveals only 12. In Fig. 4 we present QQ-plots of the absolute values of the residuals of the $\tau$-estimate against the absolute value of the residuals of the MLE, after eliminating the 27 outliers detected by the $\tau$-estimate. Fig. 4(a) shows that for the plasma beta-carotene the $\tau$-estimate gives residuals smaller than the MLE. Instead, Fig. 4(b) shows that for the plasma retinol, the distributions of both residuals are close.

In Table 4 we show the regression coefficients and their standard errors (SE) for three estimates: the $\tau$-estimate, the MLE and the MLE after omitting the 27 outliers. We only show these values for the variables that are statistically significant at level 0.05 for at least one estimate and one


Fig. 3. Box-plot of the Mahalanobis norms of the residuals.


Fig. 4. QQ-plots of the absolute values of the residuals for both regressions (a) correspond to the plasma beta-carotene and (b) to the plasma retinol.
equation. We also show the error variances corresponding to each regression in this table. The estimated standard errors of the $\tau$-estimates were calculated as proposed in the Remark after Theorem 5.

We can observe that when the MLE and $\tau$-estimates are computed using the complete dataset, some regression coefficients and standard errors are quite different. Instead the results for the $\tau$-estimate computed with all the observations are quite close to those of the MLE after deleting the 27 outliers. This is what a robust estimate is expected to do.

Table 4
Regression coefficients and standard errors for the $\tau$-estimates, the MLE and the MLE after omitting 27 outliers (MLE_27)

| Equation <br> Estimate | 1 |  |  | 2 |  |  |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
|  | $\tau$ | MLE | MLE-27 | $\tau$ | MLE | MLE-27 |
| $X_{1}$ | 0.15 | 0.14 | 0.18 | 0.27 | 0.21 | 0.26 |
| SE | 0.07 | 0.14 | 0.07 | 0.08 | 0.08 | 0.07 |
| $X_{2}$ | -0.44 | -0.37 | -0.34 | 0.02 | 0.53 | 0.09 |
| SE | 0.19 | 0.35 | 0.18 | 0.21 | 0.20 | 0.18 |
| $X_{5}$ | -0.18 | -0.33 | -0.20 | 0.02 | 0.02 | 0.02 |
| SE | 0.04 | 0.09 | 0.04 | 0.05 | 0.05 | 0.04 |
| $X_{6}$ | 0.15 | 0.89 | 0.36 | 0.02 | 0.15 | 0.10 |
| SE | 0.13 | 0.26 | 0.13 | 0.15 | 0.14 | 0.13 |
| $X_{10}$ | 0.21 | 0.33 | 0.16 | -0.13 | -0.09 | -0.12 |
| SE | 0.07 | 0.14 | 0.07 | 0.08 | 0.08 | 0.07 |
| $X_{13}$ | 0.13 | 0.21 | 0.15 | 0.01 | -0.01 | -0.01 |
| SE | 0.05 | 0.10 | 0.05 | 0.06 | 0.05 | 0.05 |
| Error variance | 0.81 | 3.51 | 0.78 | 1.01 | 1.10 | 0.79 |

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## Appendix

## A.1. Theorem 1

Proof of Theorem 1. Put $Q(B, \Sigma)=\operatorname{det}(\Sigma) \tau^{2 q}\left(d_{1}(B, \Sigma), \ldots, d_{n}(B, \Sigma)\right)$. It is easy to show that for any real $\lambda, Q(B, \lambda \Sigma)=Q(B, \Sigma)$, then the $\tau$-estimate $(\widetilde{B}, \widetilde{\Sigma})$ also minimizes $Q(B, \Sigma)$ without restrictions (observe that $(\widetilde{B}, \lambda \widetilde{\Sigma})$ minimizes $Q(B, \Sigma)$ too), or equivalently $\log (Q(B, \Sigma))$. Therefore they should satisfy the following equations:

$$
\begin{equation*}
\frac{\partial \log (Q(B, \Sigma))}{\partial B}=0, \quad \frac{\partial \log (Q(B, \Sigma))}{\partial \Sigma}=0 . \tag{A.1}
\end{equation*}
$$

From now on, and for the sake of simplicity, we will denote $d_{i}=d_{i}(B, \Sigma), d_{i}^{*}=d_{i}(B, \Sigma) / s$ $(B, \Sigma)$ and $s=s(B, \Sigma)=s\left(d_{1}(B, \Sigma), \ldots, d_{n}(B, \Sigma)\right)$.

Differentiating $\log (Q(B, \Sigma))$ with respect to $B$, after straightforward calculations, we obtain

$$
\begin{aligned}
\frac{\partial \log (Q(B, \Sigma))}{\partial\left(\operatorname{vec}\left(B^{\prime}\right)\right)^{\prime}}= & -2 q \frac{\sum_{i=1}^{n} \psi_{1}\left(d_{i}^{*}\right) \frac{\widehat{u}_{i}^{\prime}(B)}{d_{i}} \Sigma^{-1}\left(x_{i}^{\prime} \otimes I_{q}\right)}{\sum_{i=1}^{n}\left(\psi_{1}\left(d_{i}^{*}\right) d_{i}\right)} \\
& -\frac{q \sum_{i=1}^{n} \psi_{2}\left(d_{i}^{*}\right) \frac{\widehat{u}_{i}^{\prime}(B)}{d_{i}} \Sigma^{-1}\left(x_{i}^{\prime} \otimes I_{q}\right)}{s \sum_{i=1}^{n} \rho_{2}\left(d_{i}^{*}\right)}
\end{aligned}
$$

$$
+\frac{q \sum_{i=1}^{n} d_{i}^{*} \psi_{2}\left(d_{i}^{*}\right) \sum_{j=1}^{n} \psi_{1}\left(d_{j}^{*}\right) \frac{\widehat{u}_{j}^{\prime}(B)}{d_{j}} \Sigma^{-1}\left(x_{j}^{\prime} \otimes I_{q}\right)}{\sum_{i=1}^{n} \rho_{2}\left(d_{i}^{*}\right) \sum_{j=1}^{n}\left(\psi_{1}\left(d_{j}^{*}\right) d_{j}\right)} .
$$

Then, equating this last expression to zero, we have

$$
\begin{gathered}
\frac{\sum_{i=1}^{n} \psi_{1}\left(d_{i}^{*}\right) \frac{\widehat{u}_{i}^{\prime}(B)}{d_{i}} \Sigma^{-1}\left(x_{i}^{\prime} \otimes I_{q}\right)\left[-2 \sum_{j=1}^{n} \rho_{2}\left(d_{j}^{*}\right)+\sum_{j=1}^{n} \psi_{2}\left(d_{j}^{*}\right) d_{j}^{*}\right]}{\sum_{i=1}^{n} \rho_{2}\left(d_{i}^{*}\right) \sum_{j=1}^{n}\left(\psi_{1}\left(d_{j}^{*}\right) d_{j}\right)} \\
-\frac{\frac{1}{s} \sum_{i=1}^{n} \psi_{2}\left(d_{i}^{*}\right) \frac{\widehat{u}_{i}^{\prime}(B)}{d_{i}} \Sigma^{-1}\left(x_{i}^{\prime} \otimes I_{q}\right) \sum_{j=1}^{n} \psi_{1}\left(d_{j}^{*}\right) d_{j}}{\sum_{i=1}^{n} \rho_{2}\left(d_{i}^{*}\right) \sum_{j=1}^{n}\left(\psi_{1}\left(d_{j}^{*}\right) d_{j}\right)}=0,
\end{gathered}
$$

or equivalently,

$$
\begin{equation*}
\sum_{i=1}^{n} w_{n}^{*}\left(d_{i}^{*}\right) \widehat{u}_{i}^{\prime}(B) \Sigma^{-1}\left(x_{i}^{\prime} \otimes I_{q}\right)=0 \tag{A.2}
\end{equation*}
$$

Using that $\operatorname{vec}\left(\Sigma^{-1} u_{i} x_{i}^{\prime}\right)=\left(x_{i} \otimes I_{q}\right) \Sigma^{-1} u_{i}$, we can show that expression (A.2) is equivalent to (2.16).

Differentiating $\log (Q(B, \Sigma))$ with respect to $\Sigma$, we obtain

$$
\begin{align*}
\frac{\partial \log (Q(B, \Sigma))}{\partial \Sigma}= & \frac{\partial \log (\operatorname{det}(\Sigma))}{\partial \Sigma}+\frac{2 q}{s} \frac{\partial s}{\partial \Sigma}+\frac{q}{\sum_{i=1}^{n} \rho_{2}\left(d_{i}^{*}\right)} \\
& \times \sum_{i=1}^{n} \frac{\partial}{\partial \Sigma} \rho_{2}\left(d_{i}^{*}\right) \tag{A.3}
\end{align*}
$$

It is well known that for a symmetric matrix $\Sigma$,

$$
\begin{equation*}
\frac{\partial \log (\operatorname{det}(\Sigma))}{\partial \Sigma}=\Sigma^{-1} \tag{A.4}
\end{equation*}
$$

Then, using that

$$
\begin{equation*}
\frac{\partial s}{\partial \Sigma}=-\frac{1}{2 s} \frac{\left.\Sigma^{-1}\left(\sum_{i=1}^{n} \psi_{1}\left(d_{i}^{*}\right)\right) \frac{s}{d_{i}} \widehat{u}_{i}(B) \widehat{u}_{i}^{\prime}(B)\right) \Sigma^{-1}}{\sum_{i=1}^{n} \psi_{1}\left(d_{i}^{*}\right) d_{i}^{*}} \tag{A.5}
\end{equation*}
$$

and denoting $w_{i}(v)=\psi_{i}(v) / v$ for $i=1,2$, we have

$$
\begin{equation*}
\frac{\partial d_{i}^{*}}{\partial \Sigma}=\frac{1}{2 s} \Sigma^{-1}\left[-\frac{1}{d_{i}} \widehat{u}_{i}(B) \widehat{u}_{i}^{\prime}(B)+\frac{d_{i}}{s^{2}} \frac{\sum_{j=1}^{n} w_{1}\left(d_{j}^{*}\right) \widehat{u}_{j}(B) \widehat{u}_{j}^{\prime}(B)}{\sum_{j=1}^{n} \psi_{1}\left(d_{j}^{*}\right) d_{j}^{*}}\right] \Sigma^{-1} . \tag{A.6}
\end{equation*}
$$

Replacing (A.4), (A.5) and (A.6) in (A.3), and using the fact that the $\tau$-estimate satisfies (A.1), we obtain

$$
\begin{aligned}
& \Sigma^{-1}-\frac{q}{s^{2}} \frac{\Sigma^{-1}\left(\sum_{i=1}^{n} w_{1}\left(d_{i}^{*}\right) \widehat{u}_{i}(B) \widehat{u}_{i}^{\prime}(B)\right) \Sigma^{-1}}{\sum_{i=1}^{n} \psi_{1}\left(d_{i}^{*}\right) d_{i}^{*}} \\
& \quad-\frac{q\left[\sum_{i=1}^{n} w_{2}\left(d_{i}^{*}\right) \Sigma^{-1} \widehat{u}_{i}(B) \widehat{u}_{i}^{\prime}(B) \Sigma^{-1}\right]}{2 s^{2} \sum_{i=1}^{n} \rho_{2}\left(d_{i}^{*}\right)} \\
& \quad+\frac{q \sum_{i=1}^{n} \psi_{2}\left(d_{i}^{*}\right) d_{i}^{*}\left[\Sigma^{-1}\left(\sum_{k=1}^{n} w_{1}\left(d_{k}^{*}\right) \widehat{u}_{k}(B) \widehat{u}_{k}^{\prime}(B)\right) \Sigma^{-1}\right]}{2 s^{2}\left(\sum_{i=1}^{n} \rho_{2}\left(d_{i}^{*}\right)\right)\left(\sum_{i=1}^{n} \psi_{1}\left(d_{i}^{*}\right) d_{i}^{*}\right)}=0 .
\end{aligned}
$$

Solving for $\Sigma$, we obtain

$$
\Sigma=\frac{q}{s^{2}} \frac{\sum_{i=1}^{n} w_{n}^{*}\left(d_{i}^{*}\right) \widehat{u}_{i}(B) \widehat{u}_{i}^{\prime}(B)}{\sum_{k=1}^{n} \psi_{n}^{*}\left(d_{k}^{*}\right) d_{k}^{*}},
$$

and this proves (2.17).

## A.2. Theorem 2

Lemmas A.1-A. 5 are required to prove Theorem 2
Lemma A.1. Assume that $\rho_{1}$ and $\rho_{2}$ satisfy A1-A5. Given $m<\eta n$, there exists $K$ such that for any $Z^{*}=\left\{z_{1}^{*}, \ldots, z_{n}^{*}\right\} \in \mathcal{Z}_{m}, z_{i}^{*}=\left(x_{i}^{* \prime}, y_{i}^{* \prime}\right)^{\prime}$, we have $\tau\left(\left\|y_{1}^{*}\right\|, \ldots,\left\|y_{n}^{*}\right\|\right)<K$.

This lemma is proved in [16].
Lemma A.2. Assume that $\rho_{1}$ and $\rho_{2}$ satisfy A1-A5, and let $\widetilde{\Sigma}$ be the $\tau$-estimate of $\Sigma$. Given $m<\eta n$, there exist $K^{*}$ such that for any $Z^{*}=\left\{z_{1}^{*}, \ldots, z_{n}^{*}\right\} \in \mathcal{Z}_{m}, z_{i}^{*}=\left(x_{i}^{* \prime}, y_{i}^{* \prime}\right)^{\prime}$ we have $\operatorname{det}\left(\widetilde{\Sigma}\left(Z^{*}\right)\right)<K^{*}$.

Proof. Take $B=0$ and $\Sigma=\left(\tau^{2}\left(\left\|y_{1}^{*}\right\|, \ldots,\left\|y_{n}^{*}\right\|\right) / \kappa_{2}\right) I_{q}$. Observe that $d_{i}(B, \Sigma)=\sqrt{\kappa_{2}}\left\|y_{i}^{*}\right\| / \tau\left(\left\|y_{1}^{*}\right\|, \ldots,\left\|y_{n}^{*}\right\|\right)$, and therefore $\tau^{2}\left(d_{1}(B, \Sigma), \ldots, d_{n}(B, \Sigma)\right)=\kappa_{2}$. By Lemma A. 1 there exists $K$ such that

$$
\operatorname{det}(\Sigma)=\left(\tau^{2}\left(\left\|y_{1}^{*}\right\|, \ldots,\left\|y_{n}^{*}\right\|\right) / \kappa_{2}\right)^{q}<\left(K^{2} / \kappa_{2}\right)^{q}
$$

Then using the definition of $(\widetilde{B}, \widetilde{\Sigma})$ given in (2.10) and (2.11), the lemma follows with $K^{*}=$ $\left(K^{2} / \kappa_{2}\right)^{q}$.

Lemma A.3. Assume that $\rho_{1}$ and $\rho_{2}$ satisfy A1-A5. Given $K_{1}>0$ and $r>n \eta$, there exists $K_{2}$ such that for any sample $u_{1}, \ldots, u_{n}$ such that $\#\left\{i:\left|u_{i}\right|>K_{2}\right\}>r$, we have $\tau\left(u_{1}, \ldots, u_{n}\right)>K_{1}$.

This lemma is proved in [16].
Lemma A.4. Consider the same assumptions as in Theorem 2 and $m<\min ((1-\eta) n-h, n \eta)$. Then (i) $S_{m}^{-}(Z, \widetilde{\Sigma})>0$, (ii) $S_{m}^{+}(Z, \widetilde{\Sigma})<\infty$.

Proof. Suppose that $S_{m}^{-}(Z, \widetilde{\Sigma})=0$. Then there exist $Z_{j}^{*} \in \mathcal{Z}_{m}$ such that $\lambda_{1}\left(\widetilde{\Sigma}\left(Z_{j}^{*}\right)\right) \rightarrow 0$. Put $\Sigma_{j}=\widetilde{\Sigma}\left(Z_{j}^{*}\right), B_{j}=\widetilde{B}\left(Z_{j}^{*}\right)$ and $\lambda_{1 j}=\lambda_{1}\left(\widetilde{\Sigma}\left(Z_{j}^{*}\right)\right)$. Let $U_{j}$ be an orthogonal matrix of eigenvectors of $\Sigma_{j}$ and $\Lambda_{j}$ the diagonal matrix with the corresponding eigenvalues. Let $Z_{j}^{*}=\left\{z_{j 1}^{*}, \ldots, z_{j n}^{*}\right\}$, $z_{j i}^{*}=\left(x_{j i}^{* \prime}, y_{j i}^{* \prime}\right)^{\prime}$; then

$$
\begin{align*}
d_{j i}^{* 2} & =\left(y_{j i}^{*}-B_{j}^{\prime} x_{j i}^{*}\right)^{\prime} \Sigma_{j}^{-1}\left(y_{j i}^{*}-B_{j}^{\prime} x_{j i}^{*}\right) \\
& =\left(U_{j}^{\prime} y_{j i}^{*}-U_{j}^{\prime} B_{j}^{\prime} x_{j i}^{*}\right)^{\prime} \Lambda_{j}^{-1}\left(U_{j}^{\prime} y_{j i}^{*}-U_{j}^{\prime} B_{j}^{\prime} x_{j i}^{*}\right) \geqslant\left(e_{j}^{\prime} z_{j i}^{*}\right)^{2} / \lambda_{1 j}, \tag{A.7}
\end{align*}
$$

where $e_{j}=\left(-v_{j 1}, u_{j 1}\right)$, and $u_{j 1}$ and $v_{j 1}$ are, respectively, the first rows of $U_{j}^{\prime}$ and of $V_{j}=U_{j}^{\prime} B_{j}^{\prime}$.
Put $\delta=\inf _{\|a\|=1} \sup _{1 \leqslant j_{1}<\cdots<j_{h+1} \leqslant n}\left\{\left|a^{\prime} z_{j_{1}}\right|, \ldots,\left|a^{\prime} z_{j_{h+1}}\right|\right\}$. Then by the definition of $h$, we have that $\delta>0$. Since $\left\|e_{j}\right\| \geqslant 1$, there are at least $n-m-h>\eta n$ values $d_{j i}^{*}$ larger than or equal
to $\delta /\left(\lambda_{1_{j}}\right)^{1 / 2}$. Therefore, from Lemma A. 3 we obtain $\lim _{j \rightarrow \infty} \tau\left(d_{j 1}^{*}, \ldots, d_{j n}^{*}\right)=\infty$, which contradicts the definition of the $\tau$-estimates. Then part (i) of the lemma follows.

Part (ii) follows from part (i) and Lemma A.2.
Lemma A.5. Consider the same assumptions as in Theorem 2. Then for any $m<\mathrm{min}$ $((1-\eta) n-h, n \eta)$ we have $S_{m}(Z, \widetilde{B})<\infty$.

Proof. Assume that there exists $Z_{j}^{*} \in \mathcal{Z}_{m}$, such that $\left\|\widetilde{B}\left(Z_{j}^{*}\right)\right\| \rightarrow \infty$. Let $U_{j}, \Lambda_{j}, V_{j}$ and $\delta$ be as in Lemma A.4. Since $\left\|V_{j}\right\|=\left\|U_{j}^{\prime} B_{j}^{\prime}\right\|$, we also have $\left\|V_{j}\right\| \rightarrow \infty$. Without loss of generality we can suppose that for some $i_{0},\left\|v_{j i_{0}}\right\| \rightarrow \infty$, where $v_{j i_{0}}$ is the $i_{0}$ th row of $V_{j}$. According to Lemma A. 4 we can assume that $\lambda_{q}\left(\widetilde{\Sigma}\left(Z_{j}\right)\right)<K$. Let $Z_{j}^{*}=\left\{z_{j 1}^{*}, \ldots, z_{j n}^{*}\right\}, z_{j i}^{*}=\left(x_{j i}^{* \prime}, y_{j i}^{* \prime}\right)^{\prime}$, then proceeding as in (A.7) we obtain $d_{j i}^{* 2} \geqslant(1 / K)\left(e_{j}^{\prime} z_{j i}^{*}\right)^{2}$, where $e_{j}=\left(-v_{j i_{0}}, u_{j i_{0}}\right)$. Then there are at least $n-m-h>\eta n$ values $d_{j i}^{*}$ larger than $\left(\delta / K^{1 / 2}\right)\left\|v_{j i_{0}}\right\|$. Therefore by Lemma A.3, $\tau\left(d_{j 1}, \ldots, d_{j n}\right) \rightarrow \infty$, contradicting the definition of $\tau$-estimates.

Proof of Theorem 2. It follows from Lemmas A. 4 and A.5.

## A.3. Theorem 3

To prove Theorem 3 we need to introduce some notation and Lemmas A. 6 and A. 7 below.
Define $u(B)=u(B, x, y)=y-B^{\prime} x$ and

$$
d(B, \Sigma)=d(B, \Sigma, x, y)=\left(u^{\prime}(B) \Sigma^{-1} u(B)\right)^{1 / 2}
$$

Let $\left(T_{1}, T_{2}\right)$ be the estimating functional corresponding to the $\tau$-estimates ( $\widetilde{B}, \widetilde{\Sigma}$ ). Then, according to (2.16) and (2.17), given a distribution $G$ of $(x, y),\left(T_{1}(G), T_{2}(G)\right)$ are the values ( $B, \Sigma$ ) satisfying

$$
\begin{aligned}
& E_{G}\left(w^{*}\left(d^{*}(B, \Sigma, G)\right) u(B) x^{\prime}\right)=0 \\
& \Sigma=\frac{q E_{G}\left(w^{*}\left(d^{*}(B, \Sigma, G)\right) u(B) u^{\prime}(B)\right)}{\left(k^{*}(B, \Sigma, G)\right)^{2} E_{G}\left(\psi^{*}\left(d^{*}(B, \Sigma, G)\right) d^{*}(B, \Sigma, G)\right)}
\end{aligned}
$$

where $d^{*}(B, \Sigma, G)=d(B, \Sigma) / k^{*}(B, \Sigma, G), k^{*}(B, \Sigma, G)=s^{*}(H), H$ is the distribution of $d(B, \Sigma)$ under $G$ and $s^{*}(H)$ is defined in (4.2).

Lemma A.6. Suppose that we observe $z \in R^{m}$ with distribution function $G_{\boldsymbol{\theta}_{1}, \boldsymbol{\theta}_{2}}$, where $\theta_{1} \in R^{k_{1}}$ and $\theta_{2} \in R^{k_{2}}$. Consider an M-estimating functional of $\theta=\left(\theta_{1}, \theta_{2}\right), T(G)=\left(T_{1}(G), T_{2}(G)\right)$ such that

$$
E_{G}\left(h\left(z, T_{1}(G), T_{2}(G), k\left(T_{1}(G), T_{2}(G), G\right)\right)\right)=0
$$

where $h: R^{m+k_{1}+k_{2}+1} \rightarrow R^{k_{1}}$ is a differentiable function and $k: R^{k_{1}+k_{2}} \times \mathcal{F} \rightarrow R$, and $\mathcal{F}$ is the space of distributions on $R^{k_{1}+k_{2}}$. Suppose that $T$ satisfies the following strong Fisher consistency condition:

$$
\begin{equation*}
E_{G_{\theta_{1}, \theta_{2}}}\left(h\left(z, \theta_{1}, \theta_{2}, k\right)\right)=0 \quad \forall k, \tag{A.8}
\end{equation*}
$$

and

$$
\begin{equation*}
E_{G_{\theta_{1}, \theta_{2}}}\left(h_{3}\left(z, \theta_{1}, \theta_{2}, k\left(\theta_{1}, \theta_{2}, G_{\theta_{1}, \theta_{2}}\right)\right)\right)=0 \tag{A.9}
\end{equation*}
$$

where $h_{i}, 1 \leqslant i \leqslant 4$, is the derivative of $h$ with respect to the ith argument. Suppose too that $E_{G_{\theta_{1}, \theta_{2}}}\left(h\left(z, \theta_{1}, \theta_{2}, k\left(\theta_{1}, \theta_{2}, G_{\theta_{1}, \theta_{2}}\right)\right)\right)$ can be differentiated inside the expectation. Then the influence function of $T_{1}$ is given by

$$
\begin{aligned}
I C\left(z_{0}, \mathbf{T}_{1}, \theta_{1}, \theta_{2}\right)= & -\left(E_{G_{\theta_{1}, \theta_{2}}}\left(h_{2}\left(z, \theta_{1}, \theta_{2}, k\left(\theta_{1}, \theta_{2}, G_{\theta_{1}, \theta_{2}}\right)\right)\right)\right)^{-1} \\
& \times h\left(z_{0}, \theta_{1}, \theta_{2}, k\left(\theta_{1}, \theta_{2}, G_{\theta_{1}, \theta_{2}}\right)\right) .
\end{aligned}
$$

Proof. Let $G_{\varepsilon}=(1-\varepsilon) G_{\boldsymbol{\theta}_{1}, \boldsymbol{\theta}_{2}}+\varepsilon \delta_{z_{0}}$. Then $T\left(G_{\varepsilon}\right)$ satisfies

$$
\begin{aligned}
& (1-\varepsilon) E_{G_{\theta_{1}, \theta_{2}}}\left(h\left(z, T_{1}\left(G_{\varepsilon}\right), T_{2}\left(G_{\varepsilon}\right), k\left(T_{1}\left(G_{\varepsilon}\right), T_{2}\left(G_{\varepsilon}\right), G_{\varepsilon}\right)\right)\right) \\
& \quad+\varepsilon h\left(z_{0}, T_{1}\left(G_{\varepsilon}\right), T_{2}\left(G_{\varepsilon}\right), k\left(T_{1}\left(G_{\varepsilon}\right), T_{2}\left(G_{\varepsilon}\right), G_{\varepsilon}\right)\right)=0 .
\end{aligned}
$$

The proof of this lemma follows inmediately differentiating this expression with respect to $\varepsilon$ at $\varepsilon=0$ and using (A.8) and (A.9).

Lemma A.7. Consider assumptions A1-A9 and suppose $\Sigma_{0}=I_{q}$. Then, if $G$ is the distribution of $(x, y)$, we have
(a) $\quad E_{G}\left(\frac{\partial\left(\operatorname{vec}\left(w^{*}\left(d\left(B_{0}, \Sigma_{0}\right) / k\right) u\left(B_{0}\right) x^{\prime}\right)\right)}{\partial(\operatorname{vec} \Sigma)^{\prime}}\right)=0 \forall k$.
(b) $\quad E_{G}\left(\operatorname{vec}\left(w^{*}\left(d\left(B_{0}, \Sigma_{0}\right) / k\right) u\left(B_{0}\right) x^{\prime}\right)\right)=0 \forall k$.

Proof. (a) Straighforward computations lead to

$$
\begin{aligned}
& E_{G}\left(\frac{\partial\left(\operatorname{vec}\left(w^{*}\left(d\left(B_{0}, \Sigma_{0}\right) / k\right) u\left(B_{0}\right) x^{\prime}\right)\right)}{\partial(\operatorname{vec} \Sigma)^{\prime}}\right) \\
& \quad=-\frac{1}{2 k} E_{G}\left(\left(x u^{\prime} \otimes I_{q}\right)\left(\operatorname{vec} I_{q} \otimes \frac{w^{* \prime}(\|u\| / k)}{\|u\|} \operatorname{vec}\left(u u^{\prime}\right)^{\prime}\right)\right.
\end{aligned}
$$

Since the distribution of $u$ is assumed to be elliptical with $\Sigma_{0}=I_{q}$, we have that for any function $h, E_{G}\left(x_{j} u_{i} u_{k} u_{l} h(\|u\|)\right)=0$. Observe that all the elements of the right-hand side of the last equation are of this form, and therefore part (a) of the lemma follows immediately.
(b) This is derived from the fact that $E_{G}\left(x_{j} u_{i} h(\|u\|)\right)=0$.

Proof of Theorem 3. Suppose that $z=\left(x^{\prime}, y^{\prime}\right)^{\prime}$ follows MLM (1.1) and let $G_{0}$ be its distribution function. Consider first the case $\Sigma_{0}=I_{q}$; then $k^{*}\left(B_{0}, I_{q}, G_{0}\right)=k_{0}=s^{*}\left(H_{0}\right)$, where $H_{0}$ is the distribution of $\|u\|$. Using Lemma A. 6 with $\theta_{1}=\operatorname{vec}\left(B_{0}^{\prime}\right), \theta_{2}=\operatorname{vec}\left(\Sigma_{0}\right)$, and $k\left(T_{1}, T_{2}, G_{0}\right)=$ $k^{*}\left(T_{1}, T_{2}, G_{0}\right)$ and Lemma A. 7 we obtain

$$
\begin{aligned}
I F\left(y_{0}, x_{0}, T_{1}, B_{0}, I_{q}\right)= & -\left(\frac{\partial E_{G}\left(\operatorname{vec}\left(w^{*}\left(d\left(B_{0}, I_{q}\right) / k_{0}\right) u\left(B_{0}\right) x^{\prime}\right)\right)}{\partial\left(\operatorname{vec}\left(B^{\prime}\right)\right)^{\prime}}\right)^{-1} \\
& \times \operatorname{vec}\left(w^{*}\left(\frac{\left.\left(y_{0}-B_{0}^{\prime} x_{0}\right)^{\prime}\left(y_{0}-B_{0}^{\prime} x_{0}\right)\right)^{1 / 2}}{k_{0}}\right)\left(y_{0}-B_{0}^{\prime} x_{0}\right) x_{0}^{\prime}\right) .
\end{aligned}
$$

As vec $(A B C)=\left(C^{\prime} \otimes A\right) \operatorname{vec}(B)$, to prove (4.4), it will be enough to show that

$$
\begin{align*}
& \frac{\partial E_{G_{0}}\left(\operatorname{vec}\left(w^{*}\left(d\left(B_{0}, I_{q}\right) / k_{0}\right) u\left(B_{0}\right) x^{\prime}\right)\right)}{\partial\left(\operatorname{vec}\left(B^{\prime}\right)\right)^{\prime}} \\
& \quad=-\frac{1}{q} E_{G_{0}}\left((q-1) w^{*}\left(\frac{d\left(B_{0}, I_{q}\right)}{k_{0}}\right)+\psi^{* \prime}\left(\frac{d\left(B_{0}, I_{q}\right)}{k_{0}}\right)\right)\left(E_{M_{0}}\left(x x^{\prime}\right) \otimes I_{q}\right) . \tag{A.10}
\end{align*}
$$

Using that $d\left(B_{0}, I_{q}\right)=\|u\|$ and $u\left(B_{0}\right)=u$, we obtain

$$
\begin{align*}
& \frac{\partial E_{G_{0}}\left(\operatorname{vec}\left(w^{*}\left(d\left(B_{0}, I_{q}\right) / k_{0}\right) u\left(B_{0}\right) x^{\prime}\right)\right)}{\partial\left(\operatorname{vec} B^{\prime}\right)^{\prime}} \\
& \quad=E_{G_{0}}\left(\frac{\partial\left(\operatorname{vec}\left(w^{*}\left(d\left(B_{0}, I_{q}\right) / k_{0}\right) u\left(B_{0}\right) x^{\prime}\right)\right)}{\partial\left(\operatorname{vec} B^{\prime}\right)^{\prime}}\right) \\
& =E_{G_{0}}\left[-w^{*}\left(\left.\frac{\| u \mid}{k_{0}} \right\rvert\,\right)\left(x x^{\prime} \otimes I_{q}\right)+\left(x u^{\prime} \otimes I_{q}\right)\right. \\
& \quad \times\left\{\operatorname { v e c } ( I _ { q } ) \otimes \left[-\left(\frac{C \psi_{1}^{\prime}\left(\|u\| / k_{0}\right)}{\|u\| / k_{0}}+\frac{D \psi_{2}^{\prime}\left(\|u\| / k_{0}\right)}{\|u\| / k_{0}}\right.\right.\right. \\
& \left.\left.\left.\left.\quad-\frac{\psi^{*}\left(\|u\| / k_{0}\right)}{\|u\|^{2} / k_{0}^{2}}\right) \frac{u^{\prime}\left(x^{\prime} \otimes I_{q}\right)}{k_{0}\|u\|}\right]\right\}\right] . \tag{A.11}
\end{align*}
$$

Denoting

$$
g(d)=-\left[\frac{C \psi_{1}^{\prime}(d)}{d}+\frac{D \psi_{2}^{\prime}(d)}{d}-\frac{\psi^{*}(d)}{d^{2}}\right] \frac{1}{k_{0} d}
$$

and using that $E_{F_{0}}\left[u_{i} u_{j} g\left(\|u\| / k_{0}\right)\right]=0$, if $i \neq j$, we obtain

$$
\begin{align*}
& E_{G_{0}}\left(g\left(\frac{\|u\|}{k_{0}}\right)\left(x u^{\prime} \otimes I_{q}\right)\left\{\operatorname{vec}\left(I_{q}\right) \otimes\left[u^{\prime}\left(x^{\prime} \otimes I_{q}\right)\right]\right\}\right) \\
& \quad=E_{M_{0}}\left(x x^{\prime}\right) E_{F_{0}}\left[u_{i}^{2} g\left(\frac{\|u\|}{k_{0}}\right)\right] \otimes I_{q} \\
& \quad=E_{F_{0}}\left[u_{i}^{2} g\left(\frac{\|u\|}{k_{0}}\right)\right]\left(E_{M_{0}}\left(x x^{\prime}\right) \otimes I_{q}\right) . \tag{A.12}
\end{align*}
$$

We also have

$$
\begin{equation*}
E_{F_{0}}\left[u_{i}^{2} g\left(\frac{\|u\|}{k_{0}}\right)\right]=\frac{1}{q} E_{F_{0}}\left[\|u\|^{2} g\left(\frac{\|u\|}{k_{0}}\right)\right] \tag{A.13}
\end{equation*}
$$

Using (A.11), (A.12) and (A.13) we obtain (A.10), and then (4.4) holds for the case $\Sigma=I_{q}$.
Consider now the case of a general covariance matrix $\Sigma$. Take $R$ such that $\Sigma=R R^{\prime}$; then the errors $u_{i}^{*}$ of the transformed model

$$
y_{i}^{*}=R^{-1} y_{i}=R^{-1} B_{0}^{\prime} x_{i}+R^{-1} u_{i}=B_{0}^{* \prime} x_{i}+u_{i}^{*}
$$

have covariance matrix $I_{q}$. Then, (4.4) follows from the following relationship:

$$
I F\left(y_{0}, x_{0}, T_{1}, B_{0}, \Sigma\right)=R\left[I F\left(R^{-1} y_{0}, x_{0}, T_{1}, B_{0}, I_{q}\right)\right]
$$

## A.4. Theorem 4

Before proving Theorem 4 we need to introduce some notation and Lemmas A. $8-\mathrm{A} .12$ below. Define the scale-estimating functional $\tau^{*}(H)$ by

$$
\tau^{* 2}(H)=s^{* 2}(H) E_{H}\left(\rho_{2}\left(\frac{v}{s^{*}(H)}\right)\right)
$$

where the functional $s^{*}(H)$ is defined as in (4.2). Observe that $\tau\left(u_{1}, \ldots, u_{n}\right)=\tau^{*}\left(H_{n}\right)$, where $H_{n}$ is the empirical distribution of $u_{1}, \ldots, u_{n}$. Define for $a>0 r_{H}(a)=a^{2} E_{H}\left(\rho_{2}\right.$ ( $v / a)$ ).

Lemma A.8. Suppose that $\rho_{2}$ satisfies A1-A6. Then $r_{H}(a)$ is a nondecreasing function of a.
Proof. Follows from $r_{H}^{\prime}(a)=a E_{H}\left[2 \rho_{2}(v / a)-\psi_{2}(v / a) v / a\right]$ and A6.
Lemma A.9. Suppose that $\rho$ satisfies A1-A5 and that the distribution of $u$ satisfies $A 7$ with $\Sigma=\Sigma_{0}$. Let $(\mu, \Sigma)$ be such that $\operatorname{det}(\Sigma)=\operatorname{det}\left(\Sigma_{0}\right)$ and $(\mu, \Sigma) \neq\left(0, \Sigma_{0}\right)$. Then

$$
E\left(\rho\left((u-\mu)^{\prime} \Sigma^{-1}(u-\mu)\right)^{1 / 2}\right)>E\left(\rho\left(u^{\prime} \Sigma_{0}^{-1} u\right)^{1 / 2}\right)
$$

This lemma follows immediately from Theorem 1 of [4].
Lemma A.10. Suppose that $\rho$ satisfies A1-A5 and that the distribution of $u$ satisfies A7 with $\Sigma=\Sigma_{0}$. Let $(v, \Sigma)$ be such that $v$ is a random variable independent of $u, \operatorname{det}(\Sigma)=\operatorname{det}\left(\Sigma_{0}\right)$ and either (i) $P(v \neq 0)>0$ or (ii) $\Sigma \neq \Sigma_{0}$. Then

$$
E\left(\rho\left(\left((u-v)^{\prime} \Sigma^{-1}(u-v)\right)^{1 / 2}\right)>E\left(\rho\left(u^{\prime} \Sigma_{0}^{-1} u\right)^{1 / 2}\right)\right.
$$

Proof. Suppose that (i) is true. Then, by Lemma A. 9

$$
\begin{aligned}
E\left(\rho\left(\left((u-v)^{\prime} \Sigma^{-1}(u-v)\right)^{1 / 2}\right) \mid v=\mu\right) & =E\left(\rho\left((u-\mu)^{\prime} \Sigma^{-1}(u-\mu)\right)^{1 / 2}\right) \\
& \geqslant E\left(\rho\left(u^{\prime} \Sigma_{0}^{-1} u\right)^{1 / 2}\right),
\end{aligned}
$$

and the inequality is strict with probability larger than 0 . Then the lemma follows. The proof is similar when (ii) holds.

Lemma A.11. Suppose that $\rho_{1}$ and $\rho_{2}$ satisfy A1-A6 and that the distribution of $u$ satisfies A7 with $\Sigma=\Sigma_{0}$. Let $(v, \Sigma)$ be such that $v$ is a random variable independent of $u, \operatorname{det}(\Sigma)=\operatorname{det}\left(\lambda \Sigma_{0}\right)$ and either (i) $P(v \neq 0)>0$ or (ii) $\Sigma \neq \lambda \Sigma_{0}$. Let $H^{*}$ be the distribution of $\left((u-v)^{\prime} \Sigma^{-1}(u-v)\right)^{1 / 2}$ and $H$ the distribution of $\left(u^{\prime} \Sigma_{0}^{-1} u\right)^{1 / 2} / \lambda$; then $\tau^{* 2}\left(H^{*}\right)>\tau^{* 2}(H)$.

Proof. Since $\operatorname{det}(\Sigma)=\operatorname{det}\left(\lambda \Sigma_{0}\right)$, by Lemma A.10, taking as $\rho(v)=\rho_{1}\left(v /\left(\lambda s^{*}(H)\right)\right)$, we obtain

$$
E\left(\rho_{1}\left(\frac{\left((u-v)^{\prime}(\Sigma / \lambda)^{-1}(u-v)\right)^{1 / 2}}{\lambda s^{*}(H)}\right)\right)>E\left(\rho_{1}\left(\frac{\left(u^{\prime} \Sigma_{0}^{-1} u\right)^{1 / 2}}{\lambda s^{*}(H)}\right)\right)
$$

and therefore $s^{*}\left(H^{*}\right)>s^{*}(H)$.
Using Lemmas A. 7 and A. 9 , and since $\operatorname{det}(\Sigma / \lambda)=\operatorname{det}\left(\Sigma_{0}\right)$, we obtain

$$
\begin{aligned}
\tau^{* 2}\left(H^{*}\right) & =s^{* 2}\left(H^{*}\right) E\left(\rho_{2}\left(\frac{\left((u-v)^{\prime}(\Sigma / \lambda)^{-1}(u-v)\right)^{1 / 2}}{\lambda s^{*}\left(H^{*}\right)}\right)\right) \\
& \geqslant s^{* 2}(H) E\left(\rho_{2}\left(\frac{\left((u-v)^{\prime}(\Sigma / \lambda)^{-1}(u-v)\right)^{1 / 2}}{\lambda s^{*}(H)}\right)\right) \\
& >s^{* 2}(H) E\left(\rho_{2}\left(\frac{\left(u^{\prime} \Sigma_{0}^{-1} u\right)^{1 / 2}}{\lambda s^{*}(H)}\right)\right)=\tau^{* 2}(H) .
\end{aligned}
$$

This proves the lemma.

Let $x$ and $y$ be two arbitrary random vectors of dimension $p$ and $q$, respectively, and let $G$ be any arbitrary joint distribution of $(x, y)$. Given $B$ and $\Sigma$, let $H_{G, B, \Sigma}$ be the distribution of $d(B, \Sigma)$ when $(x, y)$ has distribution $G$ and define $\tau_{G}^{*}(B, \Sigma)=\tau^{*}\left(H_{G, B, \Sigma}\right)$. Given observations $\left(x_{i}, y_{i}\right)$, $1 \leqslant i \leqslant n$, let $G_{n}$ be their empirical distribution. Then, the definition of the $\tau$-estimate $(\widetilde{B}, \widetilde{\Sigma})$ given in Section 2 is equivalent to

$$
\begin{equation*}
(\widetilde{B}, \widetilde{\Sigma})=\arg \min \operatorname{det}(\Sigma) \tag{A.14}
\end{equation*}
$$

subject to

$$
\begin{equation*}
\tau^{* 2}\left(H_{G_{n}, B, \Sigma}\right)=\kappa_{2} \tag{A.15}
\end{equation*}
$$

The following lemma shows that, if we replace in (A.15) $G_{n}$ for the true distribution, then the minimum of $\operatorname{det}(\Sigma)$ is attained at $\left(B_{0},\left(k_{0}^{2} / \sigma_{0}^{2}\right) \Sigma_{0}\right)$.

Lemma A.12. Suppose that the random vector $(x, y)$ follows an MLM with parameters $B_{0}$ and $\Sigma_{0}$ where the error $u$ has a distribution satisfying A7. Suppose also that $\rho_{k}, k=1,2$ satisfy Al$A 6$ and let $G_{0}$ be the joint distribution of $(x, y)$. Then, the problem of finding $(B, \Sigma)$ minimizing $\operatorname{det}(\Sigma)$ subject to $\tau^{* 2}\left(H_{G_{0}, B, \Sigma}\right)=\kappa_{2}$ has $\left(B_{0},\left(k_{0}^{2} / \kappa_{2}\right) \Sigma_{0}\right)$ as the unique solution.

Proof. Put $\Sigma_{0}^{*}=\left(k_{0}^{2} / \sigma_{0}^{2}\right) \Sigma_{0}$. Since $u\left(B_{0}\right)=u$, then $d\left(B_{0}, \Sigma_{0}^{*}\right)=\left(\sigma_{0} / k_{0}\right)\left(u^{\prime} \Sigma_{0}^{-1} u\right)^{1 / 2}$. Because of (4.2), we obtain $s^{*}\left(H_{G_{0} B_{0} \Sigma_{0}^{*}}\right)=\sigma_{0}$ and from (5.1), $\tau^{* 2}\left(H_{G_{0} B_{0} \Sigma_{0}^{*}}\right)=\kappa_{2}$. Now, take $(B, \Sigma) \neq$ $\left(B_{0}, \Sigma_{0}^{*}\right)$ such that $\tau^{* 2}\left(H_{G_{0} B \Sigma}\right)=\kappa_{2}$. If $B \neq B_{0}$, then $u(B)=u-\left(B_{0}^{\prime}-B^{\prime}\right) x$. Put $\lambda=$ $\left(\operatorname{det}\left(\Sigma_{0}^{*}\right) / \operatorname{det}(\Sigma)\right)^{1 / q}$, and $\Sigma^{*}=\lambda \Sigma$. Then $\operatorname{det}\left(\Sigma^{*}\right)=\operatorname{det}\left(\Sigma_{0}\right)$, and since $x$ is independent of $u$ and $P\left(\left(B_{0}^{\prime}-B^{\prime}\right) x \neq 0\right)>0$, by Lemma A. 11 we have

$$
\kappa_{2}=\tau^{* 2}\left(H_{G_{0} B_{0} \Sigma_{0}^{*}}\right)<\tau^{* 2}\left(H_{G_{0} B \Sigma^{*}}\right)=\frac{\tau^{* 2}\left(H_{G_{0} B \Sigma}\right)}{\lambda}=\frac{\kappa_{2}}{\lambda} .
$$

Then $\lambda<1$ and $\operatorname{det}\left(\Sigma_{0}^{*}\right)<\operatorname{det}(\Sigma)$, proving the lemma.
Heuristic proof of Theorem 4. The theorem follows, using standard arguments, from (A.15), (A.14), Lemma A. 12 and the fact that the empirical distribution $G_{n}$ converges a.s. to $G_{0}$ (see [11]).

## A.5. Theorem 5

Let $T(F)$ be an estimating operator with values in $R^{m}$, and let $F_{n}$ be the empirical distribution based on a random sample of size $n$ with an underlying distribution $F$. Then, under suitable differentiability conditions $n^{1 / 2}\left(T\left(F_{n}\right)-T(F)\right) \rightarrow_{D} N\left(0, E_{F}\left(I F(x, T, F) I F(x, T, F)^{\prime}\right)\right)$, where $\operatorname{IF}(x, T, F)$ is the influence function of $T$ at the point $x$ and at the distribution $F$.

Heuristic proof of Theorem 5. Let us first consider the case $\Sigma_{0}=I_{q}$ and $E_{M_{0}}\left(x x^{\prime}\right)=I_{p}$. In this case

$$
I F\left(y, x, T_{1}, B_{0}, I_{q}\right)=c_{0} w_{H_{0}}^{*}\left(\frac{\|u\|}{k_{0}}\right) x u^{\prime},
$$

where $H_{0}$ is the distribution of $\|u\|$ and we obtain

$$
\begin{aligned}
& E\left(\operatorname{vec}\left(I F\left(y, x, T_{1}, B_{0}, I_{q}\right)\right) \operatorname{vec}\left(I F\left(y, x, T_{1}, B_{0}, I_{q}\right)\right)^{\prime}\right) \\
& \quad=c_{0}^{2} E_{F_{0}}\left(\frac{k_{0}^{2} \psi_{H}^{* 2}\left(\frac{\|u\|}{k_{0}}\right)}{\|u\|^{2}} u u^{\prime}\right) \otimes E_{M_{0}}\left(x x^{\prime}\right)
\end{aligned}
$$

As the distribution of $u$ is assumed to be elliptical with $\Sigma_{0}=I_{q}$, for any function $h, E_{F_{0}}$ $\left(h(\|u\|) u_{i} u_{j}\right)=0$ if $i \neq j$ and $E_{F_{0}}\left(h(\|u\|) u_{i}^{2}\right)=E_{F_{0}}\left(h(\|u\|)\|u\|^{2}\right) / q$. Then

$$
V=\frac{c_{0}^{2} k_{0}^{2} E_{H_{0}}\left(\psi_{H_{0}}^{* 2}\left(\frac{v}{k_{0}}\right)\right)}{q} I_{q} \otimes I_{p}
$$

For the general case, let $R$ and $T$ be matrices such that $\Sigma_{0}=R R^{\prime}$ and $E_{M_{0}}\left(x x^{\prime}\right)=T T^{\prime}$, and consider the following transformation $y^{*}=R^{-1} y$ and $x^{*}=T^{-1} x$. Then if $B_{0}^{*}=T^{\prime} B_{0} R^{\prime-1}$, $y^{*}=B_{0}^{* \prime} x^{*}+u^{*}$, with $u^{*}=R^{-1} u$. Since the distribution of $u^{*}$ is given by (4.1) with $\Sigma=I_{q}$ and $E\left(x^{*} x^{* \prime}\right)=I_{q},(5.2)$ follows from the equivariance of the $\tau$-estimates and the fact that

$$
\operatorname{vec}\left(B_{0}\right)=\operatorname{vec}\left(\left(T^{-1}\right)^{\prime} B_{0}^{*} R^{\prime}\right)=\left(R \otimes\left(T^{-1}\right)^{\prime}\right) \operatorname{vec}\left(B_{0}^{*}\right)
$$

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