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Searching for a Best LAD-Solution of an Overdetermined System of Linear Equations Motivated by Searching for a Best LAD-Hyperplane on the Basis of Given Data

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# Searching for a Best LAD-Solution of an Overdetermined System of Linear Equations Motivated by Searching for a Best LAD-Hyperplane on the Basis of Given Data* 

Kristian Sabo $^{\dagger} \quad$ Rudolf Scitovski ${ }^{\ddagger}$ Ivan Vazler ${ }^{\S}$


#### Abstract

We consider the problem of searching for a best LAD-solution of an overdetermined system of linear equations $\mathbf{X a}=\mathbf{z}, \mathbf{X} \in \mathbb{R}^{m \times n}, m \geq n, \mathbf{a} \in \mathbb{R}^{n}, \mathbf{z} \in \mathbb{R}^{m}$. This problem is equivalent to the problem of determining a best LAD-hyperplane $\mathbf{x} \mapsto \mathbf{a}^{T} \mathbf{x}, \mathbf{x} \in \mathbb{R}^{n}$ on the basis of given data $\left(\mathbf{x}_{i}, z_{i}\right), \mathbf{x}_{i}=\left(x_{1}^{(i)}, \ldots, x_{n}^{(i)}\right)^{T} \in \mathbb{R}^{n}, z_{i} \in \mathbb{R}, i=1, \ldots, m$, whereby the minimizing functional is of the form $$
F(\mathbf{a})=\|\mathbf{z}-\mathbf{X} \mathbf{a}\|_{1}=\sum_{i=1}^{m}\left|z_{i}-\mathbf{a}^{T} \mathbf{x}_{i}\right| .
$$

An iterative procedure is constructed as a sequence of weighted median problems, which gives the solution in finitely many steps. A criterion of optimality follows from the fact that the minimizing functional $F$ is convex, and therefore the point $\mathbf{a}^{*} \in \mathbb{R}^{n}$ is the point of a global minimum of the functional $F$ if and only if $\mathbf{0} \in \partial F\left(\mathbf{a}^{*}\right)$.

Motivation for the construction of the algorithm was found in a geometrically visible algorithm for determining a best LAD-plane $(x, y) \mapsto \alpha x+\beta y$, passing through the origin of the coordinate system, on the basis of the data $\left(x_{i}, y_{i}, z_{i}\right), i=1, \ldots, m$.


Key words: LAD; least absolute deviations; overdetermined system of linear equations; $l_{1}$-norm approximation; weighted median problem; outliers; LAD-hyperplane

AMS Classification (2010): 65F20, 41A28, 57Q55, 68W25, 90C59

## 1 Introduction

We consider the problem of searching for a best Least Absolute Deviations (LAD) solution of an overdetermined system of linear equations (see e.g. $[1,2,3,4,5,6,7,8]$ ):

Let $\mathbf{X a}=\mathbf{z}$, where $\mathbf{X} \in \mathbb{R}^{m \times n}, m \geq n$, is a matrix of full column rank, $\mathbf{a} \in \mathbb{R}^{n}, \mathbf{z} \in \mathbb{R}^{m}$, and

$$
\begin{equation*}
F(\mathbf{a})=\|\mathbf{z}-\mathbf{X} \mathbf{a}\|_{1}=\sum_{i \in I}\left|r_{i}(\mathbf{a})\right|, \quad r_{i}(\mathbf{a})=z_{i}-\mathbf{a}^{T} \mathbf{x}_{i} \tag{1}
\end{equation*}
$$

[^0]where $I=\{1, \ldots, m\}$, and $\mathbf{x}_{i}^{T}$ is the $i$-th row of the matrix $\mathbf{X}$. Functional $F$ is convex and it attains its global minimum $\mathbf{a}^{*} \in \mathbb{R}^{n}$. This point is called a LAD-solution of an overdetermined system of linear equations $\mathbf{X a}=\mathbf{z}$.

The same problem can be considered as the problem of estimation of optimal parameters of a best LAD-hyperplane $\mathbf{x} \mapsto \mathbf{a}^{T} \mathbf{x}, \mathbf{x} \in \mathbb{R}^{n}$, on the basis of the given set of experimental data $\left(\mathbf{x}_{i}, z_{i}\right), \mathbf{x}_{i}=$ $\left(x_{1}^{(i)}, \ldots, x_{n}^{(i)}\right)^{T} \in \mathbb{R}^{n}, z_{i} \in \mathbb{R}, i=1, \ldots, m$ (see e.g. $[9,10,11,12,13]$ ). In statistical literature this problem is also considered as the problem of estimating LAD-optimal parameters of linear regression (see e. g. $[14,15,16])$.

Example 1.1. Let us consider the simplest case as an illustration: the system $\mathbf{X} \alpha=\mathbf{z}, \mathbf{X} \in \mathbb{R}^{m \times 1}$, $\mathbf{z} \in \mathbb{R}^{m}$.

Searching for a best LAD-solution of this system can be considered as a problem of searching for a best LAD-line $z=\alpha x$, whose graph passes through the origin of the coordinate system, on the basis of the given set of data points $\Lambda=\left\{T_{i}=\left(x_{i}, z_{i}\right): i \in I\right\}, \quad I=\{1, \ldots, m\}, m \geq 1$ (see Fig. 1). With the notation $I^{\prime}=\left\{i \in I: x_{i}=0\right\}$, our problem is reduced to the minimization problem

$$
\begin{equation*}
\min _{\alpha \in \mathbb{R}} \sum_{i \in I}\left|z_{i}-\alpha x_{i}\right|=\sum_{i \in I^{\prime}}\left|z_{i}\right|+\min _{\alpha \in \mathbb{R}} \sum_{i \in I \backslash I^{\prime}}\left|x_{i}\right|\left|\frac{z_{i}}{x_{i}}-\alpha\right|, \tag{2}
\end{equation*}
$$

known in literature as the Weighted Median Problem (see e.g. [5, 17, 18, 19]). The following lemma ([18]) gives properties and a solution of the weighted median problem.

Lemma 1.1. Let $\left(w_{i}, y_{i}\right), i \in I=\{1, \ldots, m\}, m \geq 2$, be the data, where $y_{1} \leq y_{2} \leq \ldots \leq y_{m}$ are real numbers and $w_{i}>0$ corresponding data weights. Denote

$$
J=\left\{\nu \in I: 2 \sum_{i=1}^{\nu} w_{i}-\sum_{i=1}^{m} w_{i} \leq 0\right\}
$$

For $J \neq \emptyset$, let us denote $\nu_{0}=\max J$. Furthermore, let $\varphi: \mathbb{R} \rightarrow \mathbb{R}$ be a function defined by the formula

$$
\varphi(\alpha)=\sum_{i=1}^{m} w_{i}\left|y_{i}-\alpha\right|
$$

Then
(i) if $J=\emptyset$ (i.e. $2 w_{1}>\sum_{i=1}^{m} w_{i}$ ), then the minimum of function $\varphi$ is attained at the point $\alpha^{\star}=y_{1}$.
(ii) if $J \neq \emptyset$ and $2 \sum_{i=1}^{\nu_{0}} w_{i}<\sum_{i=1}^{m} w_{i}$, then the minimum of function $\varphi$ is attained at the point $\alpha^{\star}=y_{\nu_{0}+1}$.
(iii) if $J \neq \emptyset$ and $2 \sum_{i=1}^{\nu_{0}} w_{i}=\sum_{i=1}^{m} w_{i}$, then the minimum of function $\varphi$ is attained at every point $\alpha^{\star}$ from the segment $\left[y_{\nu_{0}}, y_{\nu_{0}+1}\right]$.

By using Lemma 1.1 we can carry out a simple analysis of our problem (2):

- if $I^{\prime}=I$, then $\mathbf{X}=(0, \ldots, 0)^{T}$, $\operatorname{rank}(\mathbf{X})=0$, and each line $z=\alpha x, \alpha \in \mathbb{R}$ is a solution of problem (2) (see Fig. 1.a);
- if $I^{\prime}=\emptyset$, whereby $x_{1}=\cdots=x_{m} \neq 0$, then $\mathbf{X}=\left(x_{1}, \ldots, x_{1}\right)^{T}$, $\operatorname{rank}(\mathbf{X})=1$. Then (2) becomes (see Fig. 1.b)

$$
\min _{\alpha \in \mathbb{R}} \sum_{i \in I}\left|z_{i}-\alpha x_{i}\right|=\left|x_{1}\right| \min _{\alpha \in \mathbb{R}} \sum_{i \in I}\left|\frac{z_{i}}{x_{1}}-\alpha\right|,
$$

and $\alpha^{*}=\operatorname{med}_{i \in I}\left(\frac{z_{i}}{x_{1}}\right)=\frac{1}{x_{1}} \operatorname{med}_{i \in I} z_{i}$ is a solution of problem (2);

- if $I \backslash I^{\prime} \neq \emptyset$, whereby $0<x_{1}<x_{m}$, then $\mathbf{X}=\left(x_{1}, \ldots, x_{m}\right)^{T}, \operatorname{rank}(\mathbf{X})=1$. In this case the solution of problem (2) is $\alpha^{*}=\operatorname{med}_{i \in I \backslash I^{\prime}}\left(\left|x_{i}\right|, \frac{z_{i}}{x_{i}}\right)$ (see Fig. 1.c).


Figure 1: Best LAD-line passing through the origin
The best LAD-solution of an overdetermined system of linear equations is important in various fields of applied research, especially in the case if among the data a substantial amount of outliers (i.e. wild points) might appear (see e.g. [14, 20, 21, 5, 22, 23]).

This principle is considered to have been proposed by the Croatian mathematician J. R. Bošković in the mid-eighteenth century (see e.g. [14, 24]). The best LAD-solution has a property that it is less sensitive to extreme errors (outliers), it appears among the data, and points out the influence of the majority of data reflecting the real nature of the problem (see e.g. [24, 20, 17]).

Classical nondifferentiable minimization methods cannot be applied directly to searching the best LAD-solution since unreasonably long computing time would be necessary or we obtain a bad approximation of the solution. That is the reason why various specialized algorithms for solving this problem have been developed lately (see e.g. [2, 20, 15, 25, 10, 5, 26, 6, 8]).

## 2 The System $m \times 2$ and the Plane $(x, y) \mapsto \alpha x+\beta y$

Searching for a best LAD-solution of the system

$$
\mathbf{X} \mathbf{a}=\mathbf{z}, \quad \mathbf{X}=\left[\begin{array}{cc}
x_{1} & y_{1}  \tag{3}\\
\vdots & \vdots \\
x_{m} & y_{m}
\end{array}\right], \quad \mathbf{z}=\left[\begin{array}{c}
z_{1} \\
\vdots \\
z_{m}
\end{array}\right], \quad \mathbf{a}=\left[\begin{array}{c}
\alpha \\
\beta
\end{array}\right]
$$

can be considered as a problem of searching for a best LAD-plane

$$
\begin{equation*}
z(x, y)=\alpha x+\beta y \tag{4}
\end{equation*}
$$

whose graph passes through the origin of the coordinate system, on the basis of the given data points

$$
\Lambda=\left\{T_{i}=\left(x_{i}, y_{i}, z_{i}\right) \in \mathbb{R}^{3}: i \in I\right\}, \quad I=\{1, \ldots, m\}, \quad m \geq 2
$$

In both cases the problem is reduced to minimizing the functional $F: \mathbb{R}^{2} \rightarrow \mathbb{R}$,

$$
\begin{equation*}
F(\alpha, \beta)=\sum_{i=1}^{m}\left|z_{i}-\alpha x_{i}-\beta y_{i}\right| \tag{5}
\end{equation*}
$$

which always attains its global minimum at $\mathbb{R}^{2}$. Solving this problem can be geometrically clearly represented creating in that way sound assumptions for solving a general problem. Searching for the aforementioned best LAD-plane is a generalization of the approach cited in [18], and a special case of the approach given in [10].

Let us introduce the following notations:

- $\mathcal{L}=\left\{P_{i}=\left(x_{i}, y_{i}\right) \in \mathbb{R}^{2}: i \in I\right\}$ - projection of the set $\Lambda$ on the $(x, y)$-plane;
- $M\left(O, T_{\mu}\right)$ - the plane containing a $z$-axis passing through the origin $O$ and the point $T_{\mu} \in \Lambda$, i.e.

$$
M\left(O, T_{\mu}\right)=\left\{(x, y, z) \in \mathbb{R}^{3}:\left|\begin{array}{cc}
x_{\mu} & y_{\mu}  \tag{6}\\
x & y
\end{array}\right|=0\right\}
$$

Note that the point $T=(x, y, z)$ lies on the plane $M\left(O, T_{\mu}\right)$ if and only if vectors $(x, y)^{T}$ and $\left(x_{\mu}, y_{\mu}\right)^{T}$ are linearly dependent or if the point $P=(x, y)$ lies on the line determined by the points $O$ and $P_{\mu} \in \mathcal{L}$;

- $I_{\mu}=\left\{i \in I: T_{i} \in M\left(O, T_{\mu}\right)\right\}$ - the set of indices of those points from $\Lambda$ lying on the plane $M\left(O, T_{\mu}\right)$. Note that $i \in I \backslash I_{\mu}$ if and only if the point $T_{i} \in \Lambda$ does not lie on the plane $M\left(O, T_{\mu}\right)$, i.e. if the corresponding point $P_{i} \in \mathcal{L}$ does not lie on the line determined by the points $O$ and $P_{\mu}$. On the other hand, in this case the $i$-th row of the matrix $X$ in (3) is linearly independent of its $\mu$-th row.

Lemma 2.1. Let $I=\{1, \ldots, m\}, m \geq 2$, be a set of indices and let
(i) $\Lambda=\left\{T_{i}=\left(x_{i}, y_{i}, z_{i}\right): i \in I\right\}$ be a set of points in space such that its projection $\mathcal{L}$ on the $(x, y)$-plane does not lie on some line passing through the origin;
(ii) $T_{\mu}=\left(x_{\mu}, y_{\mu}, z_{\mu}\right) \in \Lambda$, such that $\left(x_{\mu}, y_{\mu}\right) \neq(0,0)$ and $I_{\mu}=\left\{i \in I: T_{i} \in M\left(O, T_{\mu}\right)\right\}$.

Then there exists $\nu \in I \backslash I_{\mu}$, such that a best LAD-plane of the form (4), containing points $O$ and $T_{\mu}$, also passes through the point $T_{\nu} \in \Lambda \backslash M\left(O, T_{\mu}\right)$.

Proof. The plane of the form (4) passing through the origin $O$ and the point $T_{\mu}=\left(x_{\mu}, y_{\mu}, x_{\mu}\right)$ can be written in the form

$$
\begin{equation*}
z^{\prime}=\frac{x}{x_{\mu}}\left(z_{\mu}-\beta^{\prime} y_{\mu}\right)+\beta^{\prime} y, \quad x_{\mu} \neq 0 \quad \text { or } \quad z^{\prime \prime}=\alpha^{\prime \prime} x+\frac{y}{y_{\mu}}\left(z_{\mu}-\alpha^{\prime \prime} x_{\mu}\right), \quad y_{\mu} \neq 0 \tag{7}
\end{equation*}
$$

whereby the parameter $\beta^{\prime}$ (i.e. parameter $\alpha^{\prime \prime}$ ) can be found by minimizing the functional

$$
\begin{equation*}
\beta^{\prime} \mapsto \sum_{i \in I \backslash I_{\mu}}\left|z_{i} x_{\mu}-x_{i} z_{\mu}-\beta^{\prime}\left(y_{i} x_{\mu}-x_{i} y_{\mu}\right)\right|, \quad x_{\mu} \neq 0 \tag{8}
\end{equation*}
$$

i.e. by minimizing the functional

$$
\begin{equation*}
\alpha^{\prime \prime} \mapsto \sum_{i \in I \backslash I_{\mu}}\left|y_{i} z_{\mu}-z_{i} y_{\mu}-\alpha^{\prime \prime}\left(y_{i} x_{\mu}-x_{i} y_{\mu}\right)\right|, \quad y_{\mu} \neq 0 . \tag{9}
\end{equation*}
$$

Due to condition $(i)$, the set $I \backslash I_{\mu}=\left\{i \in I: T_{i} \notin M\left(O, T_{\mu}\right)\right\} \neq \emptyset$ and the parameter $\beta^{\prime}$, i. e. $\alpha^{\prime \prime}$, can be determined as a weighted median by minimizing functional (8), i.e. by minimizing functional (9). By the notation $w_{i}=\left|\begin{array}{cc}x_{\mu} & y_{\mu} \\ x_{i} & y_{i}\end{array}\right|$, there exists $\nu^{\prime} \in I \backslash I_{\mu}$ such that

$$
\beta^{\prime}=\operatorname{med}_{i \in I \backslash I_{\mu}}\left(\left|w_{i}\right|, \frac{1}{w_{i}}\left|\begin{array}{cc}
x_{\mu} & z_{\mu}  \tag{10}\\
x_{i} & z_{i}
\end{array}\right|\right)=\frac{1}{w_{\nu^{\prime}}}\left|\begin{array}{cc}
x_{\mu} & z_{\mu} \\
x_{\nu^{\prime}} & z_{\nu^{\prime}}
\end{array}\right|, \quad \alpha^{\prime}=\frac{1}{x_{\mu}}\left(z_{\mu}-y_{\mu} \beta^{\prime}\right), \quad x_{\mu} \neq 0
$$

i.e. there exists $\nu^{\prime \prime} \in I \backslash I_{\mu}$ such that

$$
\alpha^{\prime \prime}=\operatorname{med}_{i \in I \backslash I_{\mu}}\left(\left|w_{i}\right|, \frac{1}{w_{i}}\left|\begin{array}{cc}
z_{\mu} & y_{\mu}  \tag{11}\\
z_{i} & y_{i}
\end{array}\right|\right)=\frac{1}{w_{\nu^{\prime \prime}}}\left|\begin{array}{cc}
z_{\mu} & y_{\mu} \\
z_{\nu^{\prime \prime}} & y_{\nu^{\prime \prime}}
\end{array}\right|, \quad \beta^{\prime \prime}=\frac{1}{y_{\mu}}\left(z_{\mu}-x_{\mu} \alpha^{\prime \prime}\right), \quad y_{\mu} \neq 0
$$

Remark 2.1. Note that $\left(\alpha^{\prime}, \beta^{\prime}\right)$, i.e. $\left(\alpha^{\prime \prime}, \beta^{\prime \prime}\right)$, is a solution of the system

$$
\left[\begin{array}{cc}
x_{\mu} & y_{\mu}  \tag{12}\\
x_{\nu^{\prime}} & y_{\nu^{\prime}}
\end{array}\right]\left[\begin{array}{c}
\alpha \\
\beta
\end{array}\right]=\left[\begin{array}{c}
z_{\mu} \\
z_{\nu^{\prime}}
\end{array}\right], \quad \text { i.e. } \quad\left[\begin{array}{cc}
x_{\mu} & y_{\mu} \\
x_{\nu^{\prime \prime}} & y_{\nu^{\prime \prime}}
\end{array}\right]\left[\begin{array}{c}
\alpha \\
\beta
\end{array}\right]=\left[\begin{array}{c}
z_{\mu} \\
z_{\nu^{\prime \prime}}
\end{array}\right]
$$

whereby corresponding matrices are nonsingular.
Lemma 2.2. By assumptions as in Lemma 2.1 there holds:
(i) If $x_{\mu}=0$, i.e. $y_{\mu}=0$, then a best LAD-plane from Lemma 2.1, passing through the point $T_{\mu}$, is of the form $z^{\prime \prime}=\alpha^{\prime \prime} x+\frac{z_{\mu}}{y_{\mu}} y$, i.e of the form $z^{\prime}=\frac{z_{\mu}}{x_{\mu}} x+\beta^{\prime} y ;$
(ii) If $\left(x_{\mu}, y_{\mu}\right) \neq(0,0)$, best LAD-planes $z^{\prime}, z^{\prime \prime}$ of the form (7) correspond.

Proof. Assertion (i) is evident. For the purpose of proving assertion (ii) it suffices to show that $\alpha^{\prime}=\beta^{\prime}$ and $\alpha^{\prime \prime}=\beta^{\prime \prime}$. By using the following property of the median (see e. g. [17])

$$
\begin{equation*}
\operatorname{med}_{i}\left(p_{i}, c u_{i}+v\right)=c \operatorname{med}_{i}\left(p_{i}, u_{i}\right)+v, \quad u_{i}, c, v \in \mathbb{R}, p_{i}>0 \tag{13}
\end{equation*}
$$

we have

$$
\begin{aligned}
\alpha^{\prime}=\frac{z_{\mu}}{x_{\mu}}-\frac{y_{\mu}}{x_{\mu}} \beta^{\prime} & =\frac{z_{\mu}}{x_{\mu}}-\frac{y_{\mu}}{x_{\mu}} \operatorname{med}_{i \in I \backslash I_{\mu}}\left(\left|w_{i}\right|, \frac{1}{w_{i}}\left|\begin{array}{cc}
x_{\mu} & z_{\mu} \\
x_{i} & z_{i}
\end{array}\right|\right) \\
& =\operatorname{med}_{i \in I \backslash I_{\mu}}\left(\left|w_{i}\right|, \frac{z_{\mu}}{x_{\mu}}-\frac{y_{\mu}}{x_{\mu}} \frac{1}{w_{i}}\left|\begin{array}{cc}
x_{\mu} & z_{\mu} \\
x_{i} & z_{i}
\end{array}\right|\right) \\
& =\operatorname{med}_{i \in I \backslash I_{\mu}}\left(\left|w_{i}\right|, \frac{1}{w_{i}}\left|\begin{array}{cc}
z_{\mu} & y_{\mu} \\
z_{i} & y_{i}
\end{array}\right|\right)=\alpha^{\prime \prime} .
\end{aligned}
$$

Similarly, it can be shown that $\beta^{\prime}=\beta^{\prime \prime}$ holds.

### 2.1 Searching for a Best LAD-plane

In accordance with Lemma 2.1 we construct an algorithm that will search for a best LAD-plane passing through the origin on the basis of the given data points. First, the first point $T_{\mu}$ must be selected in the algorithm such that $\left(x_{\mu}, y_{\mu}\right) \neq(0,0)$. After that, according to Lemma 2.1, we determine the next point $T_{\nu} \in \Lambda \backslash\left\{T_{\mu}\right\}$.

Next, and again according to Lemma2.1, by the point $T_{\nu}$ we determine the following point $T_{k} \in$ $\Lambda \backslash\left\{T_{\nu}\right\}$. If $T_{k}=T_{\mu}$, the procedure is finished; otherwise we repeat the procedure. Such stopping
criterion of the iterative procedure can be seen for example in $[9,10]$. However, if there exist more indices from the set $I \backslash I_{\nu}$ on which the weighted median is attained, i.e. if the plane passing through points $T_{\mu}$ and $T_{\nu}$ also contains some other points, e.g. point $T_{\nu^{\prime}}$, the algorithm should continue with points $T_{\mu}$ and $T_{\nu^{\prime}}$. Only after all points lying on the aforementioned plane have been used the algorithm can be stopped. Such situation cannot take place if the data have the property that no three points lie on the plane passing through the origin, i.e. if the augmented matrix $[\mathbf{X} ; \mathbf{z}]$ satisfies the Haar condition (see e.g. [20]). By means of geometrical analysis in Section 2.3 and by illustrative examples in Section 4 we will illustrate all cases that might occur.

## Algorithm I. ${ }^{1}$

Step 1: Input the set of points $\Lambda=\left\{T_{i}=\left(x_{i}, y_{i}, z_{i}\right): i \in I\right\}, I=\{1, \ldots, m\}$ and check condition (i) from Lemma 2.1. Choose the point $T_{\mu}\left(x_{\mu}, y_{\mu}, z_{\mu}\right) \in \Lambda$ such that $\left(x_{\mu}, y_{\mu}\right) \neq(0,0)$ and set $\nu=\mu, \gamma=0$;

Step 2: If $\gamma=\nu$, STOP;
Else set $\gamma=\mu, \mu=\nu$ and define the set $I \backslash I_{\mu}$ and the numbers $w_{i}$ in the following way:

$$
i \in I \backslash I_{\mu} \Longleftrightarrow w_{i}=\left|\begin{array}{cc}
x_{\mu} & y_{\mu} \\
x_{i} & y_{i}
\end{array}\right| \neq 0,
$$

Step 3: If $y_{\mu} \neq 0$, determine

$$
\hat{\alpha}=\operatorname{med}_{i \in I \backslash I_{\mu}}\left(\left|w_{i}\right|, \frac{1}{w_{i}}\left|\begin{array}{cc}
z_{\mu} & y_{\mu}  \tag{14}\\
z_{i} & y_{i}
\end{array}\right|\right) ; \quad \hat{\beta}=\frac{z_{\mu}}{y_{\mu}}-\frac{x_{\mu}}{y_{\mu}} \hat{\alpha} ;
$$

Else determine

$$
\hat{\beta}=\operatorname{med}_{i \in I \backslash I_{\mu}}\left(\left|w_{i}\right|, \frac{1}{w_{i}}\left|\begin{array}{cc}
x_{\mu} & z_{\mu}  \tag{15}\\
x_{i} & z_{i}
\end{array}\right|\right) ; \quad \hat{\alpha}=\frac{z_{\mu}}{x_{\mu}}-\frac{y_{\mu}}{x_{\mu}} \hat{\beta}
$$

by which $\nu \in I \backslash I_{\mu}$ is determined for which the median in (14), i.e. (15), is attained and go to Step 2.

Remark 2.2. Note that $(\hat{\alpha}, \hat{\beta})^{T}$ in some step of Algorithm I is a solution of the system (see also Remark 2.1)

$$
\mathbf{B a}=\left[\begin{array}{c}
z_{\mu}  \tag{16}\\
z_{\nu}
\end{array}\right], \quad \mathbf{B}=\left[\begin{array}{ll}
x_{\mu} & y_{\mu} \\
x_{\nu} & y_{\nu}
\end{array}\right], \quad \mathbf{a}=\left[\begin{array}{c}
\alpha \\
\beta
\end{array}\right]
$$

and the corresponding plane passes through the points $T_{\mu}, T_{\nu} \in \Lambda$. Since $\nu \in I \backslash I_{\mu}$, the matrix $\mathbf{B}$ is nonsingular, so that the solution of system (16) is unique.

The following theorem shows how a sequence of approximations $\mathbf{a}_{0}, \mathbf{a}_{1}, \mathbf{a}_{2}, \ldots$ from Algorithm I can be successively defined as an iterative process of the form

$$
\begin{equation*}
\overline{\mathbf{a}}=\mathbf{a}+\vartheta \mathbf{p} \tag{17}
\end{equation*}
$$

where $\mathbf{p} \in \mathbb{R}^{2}$ is the direction vector and $\vartheta \in \mathbb{R}$ step length. Thereby if the parameter vector a in Lemma 2.1 is defined by the points $T_{\mu}, T_{\nu}$, then a new-better parameter vector $\overline{\mathbf{a}}$ can be determined according to Lemma 2.1 only if we start from the point $T_{\nu}$ (drop the point $T_{\mu}$ ). The following theorem

[^1]shows that in iterative process (17) this corresponds to the choice of the first column of the matrix $\mathbf{B}^{-1}$ as a direction vector $\mathbf{p}$.

If we dropped the point $T_{\nu}$, i.e. if we started from the point $T_{\mu}$, then according to Lemma 2.1 we would again obtain the point $T_{\nu}$, and it means that in the iterative process the parameter vector a will not be changed. The following theorem shows that in iterative process (17) this corresponds to the choice of the second column of the matrix $\mathbf{B}^{-1}$ as a direction vector $\mathbf{p}$.

Theorem 2.1. Let $I=\{1, \ldots, m\}, m \geq 2$, be the set of indices and let
(i) $\Lambda=\left\{T_{i}=\left(x_{i}, y_{i}, z_{i}\right): i \in I\right\}$ be a set of points in space such that its projection $\mathcal{L}$ on the $(x, y)$-plane does not lie on any line passing through the origin;
(ii) $T_{\mu}$ and $T_{\nu}$ be the first and the second point obtained by AlgorithmI, respectively, and $\mathbf{B}^{-1}=$ $\frac{1}{x_{\mu} y_{\nu}-x_{\nu} y_{\mu}}\left[\begin{array}{cc}y_{\nu} & -y_{\mu} \\ -x_{\nu} & x_{\mu}\end{array}\right]=:\left[\mathbf{d}_{1}, \mathbf{d}_{2}\right]$.
Then, if $\mathbf{a}=(\alpha, \beta)^{T}$ is the solution of system (16) for which a global minimum of the functional $F$ given by (5) is not attained, then
I. Decreasing of functional values (5) can be attained by applying iterative process (17) in direction $\mathbf{d}_{1}$ of the first column of the matrix $\mathbf{B}^{-1}$, i.e. the next-better approximation obtained in Algorithm I can be written as

$$
\overline{\mathbf{a}}=\mathbf{a}+\vartheta_{1}^{*} \mathbf{d}_{1}, \quad \mathbf{d}_{1}=\frac{1}{x_{\mu} y_{\nu}-x_{\nu} y_{\mu}}\left[\begin{array}{c}
y_{\nu}  \tag{18}\\
-x_{\nu}
\end{array}\right], \quad \vartheta_{1}^{*}=\operatorname{med}_{i \in I \backslash I_{1}}\left(\left|\mathbf{d}_{1}^{T} \boldsymbol{\xi}_{i}\right|, \frac{z_{i}-\mathbf{a}^{T} \boldsymbol{\xi}_{i}}{\mathbf{d}_{1}^{T} \boldsymbol{\xi}_{i}}\right)
$$

where $\boldsymbol{\xi}_{i}=\left(x_{i}, y_{i}\right)^{T}$ and $I_{1}=\left\{i \in I: \mathbf{d}_{1}^{T} \boldsymbol{\xi}_{i}=0\right\}$.
II. By choosing the second column $\mathbf{d}_{2}=\frac{1}{x_{\mu} y_{\nu}-x_{\nu} y_{\mu}}\left(-y_{\mu}, x_{\mu}\right)^{T}$ of the matrix $\mathbf{B}^{-1}$ as a direction vector in iterative process (17) decreasing of values of the minimizing functional (5) will not be achieved, i.e. the step length in this direction is $\vartheta_{2}^{*}=\operatorname{med}_{i \in I \backslash I_{2}}\left(\left|\mathbf{d}_{2}^{T} \boldsymbol{\xi}_{i}\right|, \frac{z_{i}-\mathbf{a}^{T} \boldsymbol{\xi}_{i}}{\mathbf{d}_{2}^{T} \boldsymbol{\xi}_{i}}\right)=0$, whereby $I_{2}=\left\{i \in I: \mathbf{d}_{2}^{T} \boldsymbol{\xi}_{i}=0\right\}$.
Proof. I. By using the aforementioned property of median (13), in direction $\mathbf{d}_{1}$ we have

$$
\begin{aligned}
\alpha+\vartheta_{1}^{*} \frac{y_{\nu}}{x_{\mu} y_{\nu}-x_{\nu} y_{\mu}} & =\operatorname{med}_{i \in I \backslash I_{1}}\left(\left|\mathbf{d}_{1}^{T} \boldsymbol{\xi}_{i}\right|, y_{\nu} \frac{z_{i}-\alpha x_{i}-\beta y_{i}}{y_{\nu} x_{i}-x_{\nu} y_{i}}+\alpha\right) \\
& =\operatorname{med}_{i \in I \backslash I_{1}}\left(\left|\mathbf{d}_{1}^{T} \boldsymbol{\xi}_{i}\right|, \frac{y_{\nu} z_{i}-y_{i}\left(\alpha x_{\nu}+\beta y_{\nu}\right)}{y_{\nu} x_{i}-x_{\nu} y_{i}}\right) \\
& =\operatorname{med}_{i \in I \backslash I_{1}}\left(\left|\mathbf{d}_{1}^{T} \boldsymbol{\xi}_{i}\right|, \frac{y_{\nu} z_{i}-y_{i} z_{\nu}}{y_{\nu} x_{i}-x_{\nu} y_{i}}\right)
\end{aligned}
$$

that according to Lemma 2.1 corresponds to the optimal value $\alpha^{\prime \prime}$ of the parameter $\alpha$ if we start from the point $T_{\nu}$. Similarly, it can be shown that $\beta+\vartheta_{1}^{*} \frac{-x_{\nu}}{x_{\mu} y_{\nu}-x_{\nu} y_{\mu}}$ corresponds to the optimal value $\beta^{\prime \prime}$ of the parameter $\beta$ if we start from the point $T_{\nu}$.
II. In direction $\mathbf{d}_{2}$ of the second column of the matrix $\mathbf{B}^{-1}$ with step length $\vartheta_{2}^{*}$ we obtain

$$
\begin{aligned}
\alpha+\vartheta_{2}^{*} \frac{-y_{\mu}}{x_{\mu} y_{\nu}-x_{\nu} y_{\mu}} & =\operatorname{med}_{i \in I \backslash I_{2}}\left(\left|\mathbf{d}_{2}^{T} \boldsymbol{\xi}_{i}\right|,-y_{\mu} \frac{z_{i}-\alpha x_{i}-\beta y_{i}}{-y_{\mu} x_{i}+x_{\mu} y_{i}}+\alpha\right) \\
& =\operatorname{med}_{i \in I \backslash I_{2}}\left(\left|\mathbf{d}_{2}^{T} \boldsymbol{\xi}_{i}\right|, \frac{-y_{\mu} z_{i}+y_{i}\left(\alpha x_{\mu}+\beta y_{\mu}\right)}{-y_{\mu} x_{i}+x_{\mu} y_{i}}\right) \\
& =\operatorname{med}_{i \in I \backslash I_{2}}\left(\left|\mathbf{d}_{2}^{T} \boldsymbol{\xi}_{i}\right|, \frac{-y_{\mu} z_{i}+y_{i} z_{\mu}}{-y_{\mu} x_{i}+x_{\mu} y_{i}}\right)=\alpha .
\end{aligned}
$$

Similarly, we obtain $\beta+\vartheta_{2}^{*} \frac{x_{\mu}}{x_{\mu} y_{\nu}-x_{\nu} y_{\mu}}=\beta$.
Remark 2.3. Note that direction vector $\mathbf{d}_{1}$ and direction vector $\mathbf{d}_{2}$ are perpendicular to the radius vector $\boldsymbol{\xi}_{\nu}=\left(x_{\nu}, y_{\nu}\right)^{T}$ of the point $T_{\nu}$ and the radius vector $\boldsymbol{\xi}_{\mu}=\left(x_{\mu}, y_{\mu}\right)^{T}$ of the point $T_{\mu}$, respectively, i.e.

$$
\mathbf{d}_{1}^{T} \boldsymbol{\xi}_{\nu}=0, \quad \mathbf{d}_{2}^{T} \boldsymbol{\xi}_{\mu}=0
$$

Note also that it is not necessary to take into account both directions $\mathbf{d}_{1}$ and $\left(-\mathbf{d}_{1}\right)$ because according to (13) the following holds: $\vartheta_{i}^{*}\left(-\mathbf{d}_{i}\right)=-\vartheta_{i}^{*}\left(\mathbf{d}_{i}\right), i=1,2$.

### 2.2 Searching for a Best LAD-solution of a System of Equations

On the basis of Algorithm I we construct a more general algorithm for searching for a best LAD-solution of system (3). Note that condition (i) from Theorem 2.1 is equivalent to the condition that the matrix $\mathbf{X}$ from (3) is of full column rank.

First, we choose two linearly independent rows of the matrix $\mathbf{X}$, by means of which we define a square nonsingular matrix $\mathbf{B}$, calculate $\mathbf{B}^{-1}=:\left[\mathbf{d}_{1}, \mathbf{d}_{2}\right]$ and define the initial approximation $\mathbf{a}_{0}$ of the solution. If the step length $\vartheta_{1}^{*}$ in direction $\mathbf{d}_{1}$ and the step length $\vartheta_{2}^{*}$ in direction $\mathbf{d}_{2}$ vanish, we suppose that we have achieved a best LAD-solution. Otherwise we search for the next approximation.

Such stopping criterion is analogous to the stopping criterion from Algorithm I and it is often mentioned in literature (see e. g. [9, 10]). A more detailed geometrical analysis in Section 2.3 and illustrative examples in Section 4 will show that it is not always implied that the point of the global minimum is attained. For further analysis of the problem the term subdifferential of the functional $F$ (see Section 3) is necessary.

## Algorithm II.

Step 1: Input matrix $\mathbf{X}=\left[\boldsymbol{\xi}_{1}, \ldots, \boldsymbol{\xi}_{m}\right]^{T}$ and vector $\mathbf{z}$ and among the rows of matrix $\mathbf{X}$ choose two linearly independent rows: $\boldsymbol{\xi}_{\mu}=\left(x_{\mu}, y_{\mu}\right)^{T}, \boldsymbol{\xi}_{\nu}=\left(x_{\nu}, y_{\nu}\right)^{T}$;
Define matrix $\mathbf{B}=\left[\begin{array}{ll}x_{\mu} & y_{\mu} \\ x_{\nu} & y_{\nu}\end{array}\right]$ and calculate $\mathbf{B}^{-1}=:\left[\mathbf{d}_{1}, \mathbf{d}_{2}\right]$ and $\mathbf{a}:=\mathbf{B}^{-1}\left[\begin{array}{l}z_{\mu} \\ z_{\nu}\end{array}\right]$;
Step 2: Define $I_{1}=\left\{i \in I: \mathbf{d}_{1}^{T} \boldsymbol{\xi}_{i}=0\right\}$ and calculate $\vartheta_{1}^{*}=\operatorname{med}_{i \in I \backslash I_{1}}\left(\left|\mathbf{d}_{1}^{T} \boldsymbol{\xi}_{i}\right|, \frac{z_{i}-\mathbf{a}^{T} \boldsymbol{\xi}_{i}}{\mathbf{d}_{1}^{\mathbf{T}} \boldsymbol{\xi}_{i}}\right)$. If $\vartheta_{1}^{*}=0$, go to Step 3 ;
Otherwise, determine $k \in I \backslash I_{1}$;
Set $\mathbf{d}:=\mathbf{d}_{1}, \vartheta^{*}=\vartheta_{1}^{*}$ and go to Step 4 ;
Step 3: Define set $I_{2}=\left\{i \in I: \mathbf{d}_{2}^{T} \boldsymbol{\xi}_{i}=0\right\}$ and calculate $\vartheta_{2}^{*}=\operatorname{med}_{i \in I \backslash I_{2}}\left(\left|\mathbf{d}_{2}^{T} \boldsymbol{\xi}_{i}\right|, \frac{z_{i}-\mathbf{a}^{T} \boldsymbol{\xi}_{i}}{\mathbf{d}_{2}^{T} \boldsymbol{\xi}_{i}}\right)$.
If $\vartheta_{2}^{*}=0, \mathrm{STOP} ;$
Otherwise, determine $k \in I \backslash I_{2}$;
Set $\mathbf{d}:=\mathbf{d}_{2}, \vartheta^{*}=\vartheta_{2}^{*}$ and go to Step 4;
Step 4: Calculate $\mathbf{a}=\mathbf{a}+\vartheta^{*} \mathbf{d}$;
Set $\mu=\nu, \nu=k$ and define matrix $\mathbf{B}=\left[\begin{array}{ll}x_{\mu} & y_{\mu} \\ x_{\nu} & y_{\nu}\end{array}\right]$;
Step 5: Calculate $\mathbf{B}^{-1}=:\left[\mathbf{d}_{1}, \mathbf{d}_{2}\right]$ by using the previously mentioned inverse matrix and go to Step 2.
Remark 2.4. Note:

1. According to Lemma 1.1, in some cases the median of the data can be attained in every point of some interval $\left[y_{\nu_{0}}, y_{\nu_{0}+1}\right]$. In order to make assertions and calculations consequent, in that case the left edge $y_{\nu_{0}}$ of the interval should be taken for the median of data.
2. The set $I \backslash I_{1}$ in Step 2 is not empty because of at least $\mu \in I \backslash I_{1}$. If that is a unique element of that set, then $\vartheta_{1}^{*}=0$. Similar also holds for the set $I \backslash I_{2}$ from Step 3.
3. Matrix $\mathbf{B}$ from Step 4 is always nonsingular because of $\operatorname{det} \mathbf{B}=c \mathbf{d}_{1}^{T} \boldsymbol{\xi}_{\nu}$, where $c$ is a determinant of the matrix $\mathbf{B}$ from the preceding step.

### 2.3 Geometrical Analysis of the Method

To every data point $\left(x_{i}, y_{i}, z_{i}\right) \in \mathbb{R}^{3}, i=1, \ldots, m$ we correspond the line

$$
p_{i}=\left\{(\alpha, \beta) \in \mathbb{R}^{2}: \alpha x_{i}+\beta y_{i}=z_{i}\right\} .
$$

Note that $\boldsymbol{\xi}_{i}=\left(x_{i}, y_{i}\right)^{T}$ is a normal vector to the line $p_{i}$.
Because of the fact that there always exists a best LAD-plane passing through at least two different data points, the global minimum of the functional $F$ is attained on the set

$$
\Omega=\left\{(\alpha, \beta):(\alpha, \beta) \in p_{i} \cap p_{j}, i \neq j, \quad i, j=1, \ldots, m\right\}
$$

Suppose we selected two data points $T_{\mu}$ and $T_{\nu}$. Lines $p_{\mu}$ and $p_{\nu}$ (see Fig. 2) are assigned to these points. Let $\left(\alpha_{0}, \beta_{0}\right)=p_{\mu} \cap p_{\nu}$ and $\mathbf{a}_{0}=\left(\alpha_{0}, \beta_{0}\right)^{T}$.


Figure 2: Geometry of the iterative process

In accordance with Algorithm II, the next approximation $\mathbf{a}_{1}=\left(\alpha_{1}, \beta_{1}\right)^{T}$ is searched for such that from the point $\mathbf{a}_{0}$ we start either in direction $\mathbf{d}_{\mu}$ perpendicular to the vector $\boldsymbol{\xi}_{\mu}=\left(x_{\mu}, y_{\mu}\right)^{T}$ or in direction $\mathbf{d}_{\nu}$ perpendicular to the vector $\boldsymbol{\xi}_{\nu}=\left(x_{\nu}, y_{\nu}\right)^{T}$, i.e. either along the line $p_{\mu}$ or along the line $p_{\nu}$. Thus, it holds $\mathbf{d}_{\mu}^{T} \boldsymbol{\xi}_{\mu}=0$ and $\mathbf{d}_{\nu}^{T} \boldsymbol{\xi}_{\nu}=0$.

The point $\left(\alpha_{0}, \beta_{0}\right)=p_{\mu} \cap p_{\nu}$, i.e. the vector $\mathbf{a}_{0}$, is obtained as a solution of system (16), and vectors $\mathbf{d}_{\mu}, \mathbf{d}_{\nu}$ should be perpendicular to the rows of the matrix $\mathbf{B}$. Columns of the matrix $\mathbf{B}^{-1}$ have such property, so that they can be used as possible direction vectors. The corresponding step length $\vartheta^{*}$ is determined by solving a weighted median problem. In such way we find the point of intersection on this line in which the functional $F$ attains the minimal value. A significant difference of the proposed method is reflected in this in relation to methods given in e.g. [2, 3, 14, 6, 8]. Namely, in those papers and books it is always the first nearest point that is searched for in which decreasing of the value of functional $F$ is attained, whereby sometimes direction $\mathbf{d}$ and $(-\mathbf{d})$ with a positive step length is especially considered.

Fig. 2 shows the flow of the iterative process for searching for a best LAD-solution of the system $\mathbf{X a}=\mathbf{z}$, where

$$
\mathbf{X}^{T}=\left[\begin{array}{cccc}
8 & 2 & 4 & -2 \\
4 & -1 & 3 & 6
\end{array}\right], \quad \mathbf{z}=(4,1,0,6)^{T}
$$

The sequence of approximations $\mathbf{a}_{0}, \mathbf{a}_{1}, \ldots$ and a direction of movement are denoted by black dots and blue arrows, respectively. In addition, beside every intersection point the value of the minimizing functional $F$ in this intersection point is given.

## 3 LAD-solution of an Overdetermined System of Linear Equations

The problem of determining a best LAD-solution of an overdetermined system of linear equations $\mathbf{X a}=\mathbf{z}$, $\mathbf{X} \in \mathbb{R}^{m \times n}, \mathbf{z} \in \mathbb{R}^{m}, m \geq n$, i.e. the problem of determining a best LAD-hyperplane $\mathbf{x} \mapsto \mathbf{a}^{T} \mathbf{x}$, on the basis of experimental data $\left(\mathbf{x}_{i}, z_{i}\right), \mathbf{x}_{i}=\left(x_{1}^{(i)}, \ldots, x_{n}^{(i)}\right)^{T} \in \mathbb{R}^{n}, z_{i} \in \mathbb{R}, i=1, \ldots, m$ is reduced to minimization of the convex functional $F: \mathbb{R}^{n} \rightarrow \mathbb{R}$,

$$
\begin{equation*}
F(\mathbf{a})=\sum_{i=1}^{m}\left|r_{i}(\mathbf{a})\right|, \quad r_{i}(\mathbf{a})=z_{i}-\mathbf{a}^{T} \mathbf{x}_{i}, \tag{19}
\end{equation*}
$$

which always attains its global minimum on $\mathbb{R}^{n}$. A subdifferential of the functional $F$ (see e.g. [27, 28, $4,29])$ is of special importance in the analysis of this problem.

Lemma 3.1. For some $\mathbf{a} \in \mathbb{R}^{n}$ a subdifferential of the functional $F$ given by (19) is

$$
\begin{equation*}
\partial F(\mathbf{a})=\sum_{i \in I_{0}}[-1,1] \mathbf{x}_{i}-\sum_{i \in I \backslash I_{0}} \sigma_{i}(\mathbf{a}) \mathbf{x}_{i}, \quad \sigma_{i}(\mathbf{a})=\operatorname{sign}\left(r_{i}(\mathbf{a})\right), \tag{20}
\end{equation*}
$$

where $[-1,1]=\{\lambda \in \mathbb{R}:-1 \leq \lambda \leq 1\}$ and $I_{0}=\left\{i \in I: r_{i}(\mathbf{a})=0\right\}$.
Proof. A subdifferential of the function $\mathbf{a} \mapsto\left|r_{i}(\mathbf{a})\right|$ is given by

$$
\partial\left(\left|r_{i}(\mathbf{a})\right|\right)=\left\{\begin{array}{rl}
-\mathbf{x}_{i}, & r_{i}(\mathbf{a})>0 \\
\mathbf{x}_{i}, & r_{i}(\mathbf{a})<0 \\
{[-1,1] \mathbf{x}_{i},} & r_{i}(\mathbf{a})=0
\end{array}=\left\{\begin{aligned}
-\operatorname{sign}\left(r_{i}(\mathbf{a})\right) \mathbf{x}_{i}, & r_{i}(\mathbf{a}) \neq 0 \\
{[-1,1] \mathbf{x}_{i}, } & r_{i}(\mathbf{a})=0
\end{aligned}\right.\right.
$$

from which there follows (20).
Definition 3.1. Let the global minimum $\hat{\mathbf{a}} \in \mathbb{R}^{n}$ of the functional $F$ given by (19) be searched by the iterative procedure of the form

$$
\overline{\mathbf{a}}=\mathbf{a}+\vartheta \mathbf{p}, \quad \mathbf{p} \in \mathbb{R}^{n}, \quad \vartheta \in \mathbb{R}
$$

Optimal step length $\vartheta^{*}$ in direction $\mathbf{p}$ implies

$$
\vartheta^{*}=\underset{\vartheta \in \mathbb{R}}{\operatorname{argmin}} \varphi(\vartheta), \quad \varphi(\vartheta)=F(\mathbf{a}+\vartheta \mathbf{p})-F(\mathbf{a}) .
$$

Theorem 3.1. Let the global minimum of the functional $F$ given by (19) be searched for by the iterative procedure of the form $\overline{\mathbf{a}}=\mathbf{a}+\vartheta \mathbf{p}$. Then the optimal step length $\vartheta^{*}$ in direction $\mathbf{p}$ is given by

$$
\begin{equation*}
\vartheta^{*}=\operatorname{med}_{i \in I \backslash I_{1}}\left(\left|\mathbf{p}^{T} \mathbf{x}_{i}\right|, \frac{z_{i}-\mathbf{a}^{T} \mathbf{x}_{i}}{\mathbf{p}^{T} \mathbf{x}_{i}}\right), \quad I_{1}=\left\{i \in I: \mathbf{p}^{T} \mathbf{x}_{i}=0\right\} \tag{21}
\end{equation*}
$$

Proof. There holds

$$
\begin{aligned}
\varphi(\vartheta)=F(\mathbf{a}+\vartheta \mathbf{p})-F(\mathbf{a}) & =\sum_{i \in I \backslash I_{1}}\left|z_{i}-\mathbf{a}^{T} \mathbf{x}_{i}-\mathbf{p}^{T} \mathbf{x}_{i} \vartheta\right|-\sum_{i \in I \backslash I_{1}}\left|z_{i}-\mathbf{a}^{T} \mathbf{x}_{i}\right| \\
& \geq \sum_{i \in I \backslash I_{1}}\left|\mathbf{p}^{T} \mathbf{x}_{i}\right|\left|\frac{z_{i}-\mathbf{a}^{T} \mathbf{x}_{i}}{\mathbf{p}^{T} \mathbf{x}_{i}}-\vartheta^{*}\right|-\sum_{i \in I \backslash I_{1}}\left|z_{i}-\mathbf{a}^{T} \mathbf{x}_{i}\right|,
\end{aligned}
$$

whereby the equality holds if and only if $\vartheta^{*}$ is given by (21).
Lemma 3.2. Let for some $\mathbf{a} \in \mathbb{R}^{n}$ and $\mathbf{p} \in \mathbb{R}^{n}$

$$
\begin{aligned}
I_{0} & :=\left\{i \in I: r_{i}(\mathbf{a})=0\right\} \\
J & :=\left\{i \in I \backslash I_{0}: \mathbf{p}^{T} \mathbf{x}_{i} \neq 0\right\}
\end{aligned}
$$

Then there exists $\varepsilon>0$ such that for every $\vartheta \in(-\varepsilon, \varepsilon)$ the following holds

$$
\begin{equation*}
F\left(\mathbf{a}+\vartheta^{*} \mathbf{p}\right)-F(\mathbf{a}) \leq-\vartheta \mathbf{p}^{T} \mathbf{h}(\mathbf{a})+|\vartheta| \sum_{i \in I_{0}}\left|\mathbf{p}^{T} \mathbf{x}_{i}\right| \tag{22}
\end{equation*}
$$

where $\mathbf{h}(\mathbf{a}):=\sum_{i \in I \backslash I_{0}} \sigma_{i}(\mathbf{a}) \mathbf{x}_{i} \in \partial F(\mathbf{a})$ and

$$
\vartheta^{*}=\operatorname{med}_{i \in I_{0} \cup J}\left(w_{i}, \rho_{i}\right), \quad w_{i}=\left|\mathbf{p}^{T} \mathbf{x}_{i}\right|, i \in I_{0} \cup J, \quad \rho_{i}= \begin{cases}0, & i \in I_{0}  \tag{23}\\ \frac{r_{i}(\mathbf{a})}{\mathbf{p}^{T} \mathbf{x}_{i}}, & i \in J .\end{cases}
$$

Proof. Because of $r_{i}(\mathbf{a})=0, \forall i \in I_{0}$, according to Theorem 3.1, for every $\vartheta \in \mathbb{R}$ the following holds

$$
\begin{aligned}
F\left(\mathbf{a}+\vartheta^{*} \mathbf{p}\right)-F(\mathbf{a}) & =\sum_{i \in I_{0}}\left|\mathbf{p}^{T} \mathbf{x}_{i}\right|\left|0-\vartheta^{*}\right|+\sum_{i \in J}\left|\mathbf{p}^{T} \mathbf{x}_{i}\right|\left|\frac{r_{i}(\mathbf{a})}{\mathbf{p}^{T} \mathbf{x}_{i}}-\vartheta^{*}\right|-\sum_{i \in J}\left|r_{i}(\mathbf{a})\right| \\
& \leq|\vartheta| \sum_{i \in I_{0}}\left|\mathbf{p}^{T} \mathbf{x}_{i}\right|+\sum_{i \in J}\left|\mathbf{p}^{T} \mathbf{x}_{i}\right|\left|\frac{r_{i}(\mathbf{a})}{\mathbf{p}^{T} \mathbf{x}_{i}}-\vartheta\right|-\sum_{i \in J}\left|r_{i}(\mathbf{a})\right| \\
& =|\vartheta| \sum_{i \in I_{0}}\left|\mathbf{p}^{T} \mathbf{x}_{i}\right|+\sum_{i \in J}\left(r_{i}(\mathbf{a})-\vartheta \mathbf{p}^{T} \mathbf{x}_{i}\right) \operatorname{sign}\left(r_{i}(\mathbf{a})-\vartheta \mathbf{p}^{T} \mathbf{x}_{i}\right)-\sum_{i \in J}\left|r_{i}(\mathbf{a})\right| .
\end{aligned}
$$

Note that $\sum_{i \in J} \mathbf{p}^{T} \mathbf{x}_{i}=\sum_{i \in I \backslash I_{0}} \mathbf{p}^{T} \mathbf{x}_{i}$ and that there always exists $\varepsilon>0$ such that for every $\vartheta \in(-\varepsilon, \varepsilon)$ there holds

$$
\begin{equation*}
\operatorname{sign}\left(r_{i}(\mathbf{a})-\vartheta \mathbf{p}^{T} \mathbf{x}_{i}\right)=\operatorname{sign}\left(r_{i}(\mathbf{a})\right)=\sigma_{i}(\mathbf{a}), \quad \forall i \in J \tag{24}
\end{equation*}
$$

Therefore

$$
\begin{aligned}
F\left(\mathbf{a}+\vartheta^{*} \mathbf{p}\right)-F(\mathbf{a}) & \leq|\vartheta| \sum_{i \in I_{0}}\left|\mathbf{p}^{T} \mathbf{x}_{i}\right|+\sum_{i \in J} r_{i}(\mathbf{a}) \sigma_{i}(\mathbf{a})-\vartheta \sum_{i \in I \backslash I_{0}}\left(\mathbf{p}^{T} \mathbf{x}_{i}\right) \sigma_{i}(\mathbf{a})-\sum_{i \in J}\left|r_{i}(\mathbf{a})\right| \\
& =|\vartheta| \sum_{i \in I_{0}}\left|\mathbf{p}^{T} \mathbf{x}_{i}\right|+\sum_{i \in J}\left|r_{i}(\mathbf{a})\right|-\vartheta \sum_{i \in I \backslash I_{0}}\left(\mathbf{p}^{T} \mathbf{x}_{i}\right) \sigma_{i}(\mathbf{a})-\sum_{i \in J}\left|r_{i}(\mathbf{a})\right| \\
& =-\vartheta \mathbf{p}^{T} \mathbf{h}(\mathbf{a})+|\vartheta| \sum_{i \in I_{0}}\left|\mathbf{p}^{T} \mathbf{x}_{i}\right| .
\end{aligned}
$$

Theorem 3.2. Let $\mathbf{X} \in \mathbb{R}^{m \times n}, m \geq n$, be the matrix of full column rank, $\mathbf{z} \in \mathbb{R}^{m}, \mathbf{B}=\left[\mathbf{x}_{i_{1}}, \ldots, \mathbf{x}_{i_{n}}\right]^{T} \in$ $\mathbb{R}^{n \times n}$, a nonsingular square submatrix of the matrix $\mathbf{X}, \mathbf{z}_{B}=\left(z_{i_{1}}, \ldots, z_{i_{n}}\right)^{T}, I_{B}=\left\{i_{1}, \ldots, i_{n}\right\}$, and let
(i) $\hat{\mathbf{a}} \in \mathbb{R}^{n}$ be the solution of the system $\mathbf{B a}=\mathbf{z}_{B}$, and $\mathbf{B}^{-1}=\left[\mathbf{d}_{1}, \ldots, \mathbf{d}_{n}\right]$,
(ii) $I_{0}=\left\{i \in I: r_{i}(\hat{\mathbf{a}})=0\right\}$,
(iii) $h(\hat{\mathbf{a}})=\sum_{i \in I \backslash I_{0}} \sigma_{i}(\hat{\mathbf{a}}) \mathbf{x}_{i} \in \partial F(\hat{\mathbf{a}})$.

If there exists $j_{0} \in I_{B}$ such that

$$
\begin{equation*}
\left|\mathbf{d}_{j_{0}}^{T} \mathbf{h}(\hat{\mathbf{a}})\right|>1+\sum_{i \in I_{0} \backslash I_{B}}\left|\mathbf{d}_{j_{0}}^{T} \mathbf{x}_{i}\right|, \tag{25}
\end{equation*}
$$

then

$$
\vartheta^{*}=\operatorname{med}_{i \in I_{0} \cup J}\left(w_{i}, \rho_{i}\right) \neq 0, \quad \text { where } \quad w_{i}=\left|\mathbf{d}_{j_{0}}^{T} \mathbf{x}_{i}\right|, \quad \rho_{i}= \begin{cases}0, & i \in I_{0}  \tag{26}\\ \frac{r_{i}(\hat{\mathbf{a}})}{\mathbf{d}_{j_{0}}^{\mathbf{x}_{i}}}, & i \in J\end{cases}
$$

where $J=\left\{i \in I \backslash I_{B}: \mathbf{d}_{j_{0}}^{T} \mathbf{x}_{i} \neq 0\right\}$, and the following holds

$$
\begin{equation*}
F\left(\hat{\mathbf{a}}+\vartheta^{*} \mathbf{d}_{j_{0}}\right)<F(\hat{\mathbf{a}}) . \tag{27}
\end{equation*}
$$

Proof. Let us first show that $\vartheta^{*} \neq 0$. Suppose contrary, i.e. that $\vartheta^{*}=0$. For that purpose define an auxiliary function $\psi: \mathbb{R} \rightarrow \mathbb{R}$,

$$
\psi(\vartheta)=\sum_{i \in I_{0} \cup I_{1}} w_{i}\left|\rho_{i}-\vartheta\right|,
$$

which attains its global minimum for $\vartheta^{*}=\operatorname{med}_{i \in I_{0} \cup J}\left(w_{i}, \rho_{i}\right)$. Thereby $\vartheta^{*}=0$ if and only if $0 \in \partial \psi(0)$, where $\partial \psi(0)$ is a subdifferential of the function $\psi$ in the point 0

$$
\partial \psi(0)=\sum_{i \in I_{0}}[-1,1] \mathbf{d}_{j_{0}}^{T} \mathbf{x}_{i}-\sum_{i \in J} \sigma_{i}(\hat{\mathbf{a}}) \mathbf{d}_{j_{0}}^{T} \mathbf{x}_{i} .
$$

Since generally $I_{B} \subseteq I_{0}$ and $\mathbf{d}_{j_{0}}^{T} \mathbf{x}_{i}=\delta_{i j_{0}}, \forall i \in I_{B}$, the condition mentioned will be fulfilled if and only if $\exists \lambda_{0}, \gamma_{i} \in[-1,1]$, such that

$$
\lambda_{0}+\sum_{i \in I_{0} \backslash I_{B}} \gamma_{i} \mathbf{d}_{j_{0}}^{T} \mathbf{x}_{i}-\sum_{i \in J} \sigma_{i}(\hat{\mathbf{a}}) \mathbf{d}_{j_{0}}^{T} \mathbf{x}_{i}=0,
$$

i.e. since $\mathbf{d}_{j_{0}}^{T} \mathbf{h}(\hat{\mathbf{a}})=\sum_{i \in J} \sigma_{i}(\hat{\mathbf{a}}) \mathbf{d}_{j_{0}}^{T} \mathbf{x}_{i}$, if and only if

$$
\begin{equation*}
\lambda_{0}+\sum_{i \in I_{0} \backslash I_{B}} \gamma_{i} \mathbf{d}_{j_{0}}^{T} \mathbf{x}_{i}=\mathbf{d}_{j_{0}}^{T} \mathbf{h}(\hat{\mathbf{a}}) \tag{28}
\end{equation*}
$$

Note that

$$
-1-\sum_{i \in I_{0} \backslash I_{B}}\left|\mathbf{d}_{j_{0}}^{T} \mathbf{x}_{i}\right| \leq \lambda_{0}+\sum_{i \in I_{0} \backslash I_{B}} \gamma_{i} \mathbf{d}_{j_{0}}^{T} \mathbf{x}_{i} \leq 1+\sum_{i \in I_{0} \backslash I_{B}}\left|\mathbf{d}_{j_{0}}^{T} \mathbf{x}_{i}\right| .
$$

According to (28), this means that

$$
\left|\mathbf{d}_{j_{0}}^{T} \mathbf{h}(\hat{\mathbf{a}})\right| \leq 1+\sum_{i \in I_{0} \backslash I_{B}}\left|\mathbf{d}_{j_{0}}^{T} \mathbf{x}_{i}\right|,
$$

which contradicts assumption (25).
For proving (27) let us first notice that

$$
\sum_{i \in I_{0}}\left|\mathbf{d}_{j_{0}}^{T} \mathbf{x}_{i}\right|=\left|\mathbf{d}_{j_{0}}^{T} \mathbf{x}_{j_{0}}\right|+\sum_{i \in I_{0} \backslash I_{B}}\left|\mathbf{d}_{j_{0}}^{T} \mathbf{x}_{i}\right|=1+\sum_{I_{0} \backslash I_{B}}\left|\mathbf{d}_{j_{0}}^{T} \mathbf{x}_{i}\right| .
$$

Due to (25), according to Lemma 3.2, there exists $\varepsilon>0$ such that for every $\vartheta \in(-\varepsilon, \varepsilon)$ the following holds

$$
\begin{align*}
F\left(\hat{\mathbf{a}}+\vartheta^{*} \mathbf{d}_{j_{0}}\right)-F(\hat{\mathbf{a}}) \leq-\vartheta \mathbf{d}_{j_{0}}^{T} \mathbf{h}(\hat{\mathbf{a}})+|\vartheta| \sum_{i \in I_{0}}\left|\mathbf{d}_{j_{0}}^{T} \mathbf{x}_{i}\right| & \leq-\vartheta \mathbf{d}_{j_{0}}^{T} \mathbf{h}(\hat{\mathbf{a}})+|\vartheta|+|\vartheta| \sum_{i \in I_{0} \backslash I_{B}}\left|\mathbf{d}_{j_{0}}^{T} \mathbf{x}_{i}\right| \\
& <-\vartheta \mathbf{d}_{j_{0}}^{T} \mathbf{h}(\hat{\mathbf{a}})+|\vartheta|\left|\mathbf{h}^{T} \mathbf{d}_{j_{0}}\right|=: A \tag{29}
\end{align*}
$$

where $\vartheta^{*}$ is given by (26). If $\mathbf{d}_{j_{0}}^{T} \mathbf{h}(\hat{\mathbf{a}})>0$, then $A=0$ for $\vartheta \in(0, \varepsilon)$, and if $\mathbf{d}_{j_{0}}^{T} \mathbf{h}(\hat{\mathbf{a}})<0$, then $A=0$ for $\vartheta \in(-\varepsilon, 0)$.

Hence, decreasing of the value of the minimizing functional $F$ is attained in the direction $\mathbf{d}_{j_{0}}$ with the step length $\vartheta^{*} \neq 0$.

Theorem 3.3. By the assumption as in Theorem 3.2, let $I_{0}=I_{B}$. Then,
I. Functional $F$ attains its global minimum for $\hat{\mathbf{a}}=\mathbf{B}^{-1} \mathbf{z}_{B}$ if and only if $\left|\mathbf{d}_{j}^{T} \mathbf{h}(\hat{\mathbf{a}})\right| \leq 1 \forall j \in I_{B}$.
II. If there exists $j_{0} \in I_{B}$ such that $\left|\mathbf{d}_{j_{0}}^{T} \mathbf{h}(\hat{\mathbf{a}})\right|>1$, then

$$
F\left(\hat{\mathbf{a}}+\vartheta^{*} \mathbf{d}_{j_{0}}\right)<F(\hat{\mathbf{a}})
$$

where $\vartheta^{*}$ is given by (26).
Proof. I. According to Lemma 3.1, the subdifferential of the functional $F$ in the point $\mathbf{a} \in \mathbb{R}^{n}$ can be written as

$$
\begin{equation*}
\partial F(\mathbf{a})=\sum_{i \in I_{B}}[-1,1] \mathbf{x}_{i}-\sum_{i \in I \backslash I_{B}} \sigma_{i}(\mathbf{a}) \mathbf{x}_{i} \tag{30}
\end{equation*}
$$

The point $\hat{\mathbf{a}} \in \mathbb{R}^{n}$ is the point of the global minimum of the functional $F$ if and only if $\mathbf{0} \in \partial F(\hat{\mathbf{a}})$ (see e.g. [4]). Since the matrix $\mathbf{B}$ is nonsingular, then from (30) it follows that $\hat{\mathbf{a}}=\mathbf{B}^{-1} \mathbf{z}_{B}$ is the point of the global minimum of the functional $F$ if and only if there exists $\boldsymbol{\lambda} \in \mathbb{R}^{n},\|\boldsymbol{\lambda}\|_{\infty} \leq 1$, such that

$$
\begin{equation*}
\boldsymbol{\lambda}=\sum_{i \in I \backslash I_{B}} \sigma_{i}(\hat{\mathbf{a}})\left(\mathbf{B}^{-1}\right)^{T} \mathbf{x}_{i} \tag{31}
\end{equation*}
$$

i.e. if and only if

$$
\begin{equation*}
\left|\lambda_{j}\right|=\left|\sum_{i \in I \backslash I_{B}} \sigma_{i}(\hat{\mathbf{a}}) \mathbf{d}_{j}^{T} \mathbf{x}_{i}\right|=\left|\mathbf{d}_{j}^{T} \mathbf{h}(\hat{\mathbf{a}})\right| \leq 1, \quad \forall j \in I_{B} \tag{32}
\end{equation*}
$$

II. The proof of this assertion follows directly as a special case of Theorem 3.2.
3.1 Algorithm for Searching for a Best LAD-solution of an Overdetermined System of Linear Equations

Analogously to Algorithm II we construct an algorithm for searching for a best LAD-solution of an overdetermined system of linear equations $\mathbf{X a}=\mathbf{z}$, where $\mathbf{X} \in \mathbb{R}^{m \times n}, m \geq n$, is the matrix of full column rank, $\mathbf{a} \in \mathbb{R}^{n}, \mathbf{z} \in \mathbb{R}^{m}$, by minimizing the functional

$$
F(\mathbf{a})=\sum_{i \in I}\left|z_{i}-\mathbf{a}^{T} \mathbf{x}_{i}\right|, \quad I=\{1, \ldots, m\}
$$

where $\mathbf{x}_{i}^{T}$ is the i-th row of the matrix $\mathbf{X}$. The algorithm is constructed as an iterative process of the form

$$
\overline{\mathbf{a}}=\mathbf{a}+\vartheta \mathbf{p}
$$

where $\mathbf{p} \in \mathbb{R}^{n}$ is the direction vector, and $\vartheta$ is the step length in this direction, which are determined in accordance with Theorem 3.2, i.e. Theorem 3.3.

Remark 3.1. The approximation $\hat{\mathbf{a}}$ of the solution is considered to be optimal if and only if $\mathbf{0} \in \partial F(\hat{\mathbf{a}})$ (see e.g. [27, 28, 4, 29]). With notations $I_{0}=\left\{i \in I: z_{i}-\hat{\mathbf{a}} \mathbf{x}_{i}=0\right\}=\left\{i_{i}, \ldots, i_{l}\right\}, \mathbf{X}_{0}=\left[\mathbf{x}_{i_{1}}, \ldots, \mathbf{x}_{i_{l}}\right]$, $\mathbf{h}(\hat{\mathbf{a}})=\sum_{i \in I \backslash I_{0}} \sigma_{i}(\hat{\mathbf{a}}) \mathbf{x}_{i}$, the conditions mentioned in Theorem 3.2 and Theorem 3.3 will be fulfilled if and only if the system

$$
\begin{equation*}
\mathbf{X}_{0} \boldsymbol{\lambda}=\mathbf{h}(\hat{\mathbf{a}}) \quad \text { with condition } \quad\|\boldsymbol{\lambda}\|_{\infty} \leq 1 \tag{33}
\end{equation*}
$$

has a solution. If specially $I_{0}=I_{B}$, then the system from (33) becomes

$$
\mathbf{B}^{T} \boldsymbol{\lambda}=\mathbf{h}(\hat{\mathbf{a}})
$$

whose solution is

$$
\lambda_{j}=\mathbf{d}_{j}^{T} \mathbf{h}(\hat{\mathbf{a}}), \quad j \in I_{B}
$$

so that (33) will have a solution if and only if $\left|\mathbf{d}_{j}^{T} \mathbf{h}(\hat{\mathbf{a}})\right|<1, \forall j \in I_{B}$, which is in accordance with Theorem 3.3. For the example from Section 2.3 we obtain $\hat{\mathbf{a}}=(0,1)^{T}, I_{0}=I_{B}=\{1,4\}$, and $\boldsymbol{\lambda}=$ $\left(-\frac{5}{14},-\frac{3}{7}\right)^{T}$.

## Algorithm III.

Step $0 I=\{1, \ldots, m\}, \mathbf{X}=\left[\mathbf{x}_{1}, \ldots, \mathbf{x}_{m}\right]^{T}, \mathbf{x}_{i} \in \mathbb{R}^{n}, \mathbf{z}=\left(z_{1}, \ldots, z_{m}\right)^{T} \in \mathbb{R}^{m} ;$
Step 1: Choose $n$ linearly independent rows in $\mathbf{X}$ with ordinal numbers $I_{B}=\left\{i_{1}, \ldots, i_{n}\right\}$ and set:
$\mathbf{B}=\left[\mathbf{x}_{i_{1}}, \ldots, \mathbf{x}_{i_{n}}\right]^{T}, \quad \mathbf{z}_{B}=\left(z_{i_{1}}, \ldots, z_{i_{n}}\right)^{T}$,
$\mathbf{B}^{-1}=\left[\mathbf{d}_{1}, \ldots, \mathbf{d}_{n}\right]$,
$\mathbf{a}=\mathbf{B}^{-1} \mathbf{z}_{B}, \quad I_{0}=\left\{i \in I: z_{i}-\mathbf{a}^{T} \mathbf{x}_{i}=0\right\}$,
$\mathbf{h}=\sum_{i \in I \backslash I_{0}} \operatorname{sign}\left(z_{i}-\mathbf{a}^{T} \mathbf{x}_{i}\right) \mathbf{x}_{i} ;$
Step 2: If there exists $j_{0} \in I_{B}$ such that $\left|\mathbf{d}_{j_{0}}^{T} \mathbf{h}\right|>1+\sum_{i \in I_{0} \backslash I_{B}}\left|\mathbf{d}_{j_{0}}^{T} \mathbf{x}_{i}\right|$, go to Step 3;
Else go to Step 5.
Step 3: Define $J=\left\{i \in I \backslash I_{0}: \mathbf{d}_{j_{0}}^{T} \mathbf{x}_{i} \neq 0\right\}$ and
$\forall i \in I_{0} \cup J$ define $w_{i}=\left|\mathbf{d}_{j_{0}}^{T} \mathbf{x}_{i}\right|, \quad \rho_{i}= \begin{cases}0, & i \in I_{0} \\ \frac{z_{i}-\mathbf{a}^{T} \mathbf{x}_{i}}{\mathbf{d}_{j_{0}}^{T} \mathbf{x}_{i}}, & i \in J,\end{cases}$
and determine $\nu \in I_{0} \cup J$ on which $\vartheta^{*}=\operatorname{med}_{i \in I_{0} \cup J}^{J_{0}}\left(w_{i}, \rho_{i}\right)$ is attained;
Set $\mathbf{a}=\mathbf{a}+\vartheta^{*} \mathbf{d}_{j_{0}}$.
Step 4: Define a new matrix $B$, which is made from the old one by replacing the $j_{0}$-th row by the $\nu$-th row;

Set $I_{B}=$ ReplacePart $\left[I_{B}, j_{0} \rightarrow \nu\right], \quad I_{0}=\left\{i \in I: z_{i}-\mathbf{a}^{T} \mathbf{x}_{i}=0\right\}$,
determine a new matrix $\mathbf{B}^{-1}=\left[\mathbf{d}_{1}, \ldots, \mathbf{d}_{n}\right]$, calculate
$\mathbf{h}=\sum_{i \in I \backslash I_{0}} \operatorname{sign}\left(z_{i}-\mathbf{a}^{T} \mathbf{x}_{i}\right) \mathbf{x}_{i}$ and go to Step 2.
Step 5: Define $\mathbf{X}_{0}=\left[\mathbf{x}_{j_{0}}, \ldots, \mathbf{x}_{j_{l}}\right]^{T}$, where $\left\{j_{0}, \ldots, j_{l}\right\}=I_{0}$.
If the system $\mathbf{X}_{0} \boldsymbol{\lambda}=\mathbf{h}$ subject to $\|\boldsymbol{\lambda}\|_{\infty} \leq 1$ has a solution, STOP;
Else set Ind $=\{ \}$ and go to Step 6 .

Step 6: Choose $j_{0} \in\left(I_{0} \backslash I n d\right) \backslash I_{B}$, set $I n d=I n d \cup\left\{j_{0}\right\}$ and
define $\overline{\mathbf{B}}=\left[\mathbf{x}_{j_{0}}, \mathbf{x}_{i_{1}}, \ldots, \mathbf{x}_{i_{n}}\right]^{T}$;
Do QR factorization with pivoting $\overline{\mathbf{B}}=Q R$ such that the position of the first row is not changed;
Define a new set of indices $I_{B}=\left\{i_{1}, \ldots, i_{n}\right\}$ ( $i_{1}$ is a new row).
If there exists $j_{0} \in I_{B}$ such that $\left|\mathbf{d}_{j_{0}}^{T} \mathbf{h}\right|>1+\sum_{i \in I_{0} \backslash I_{B}}\left|\mathbf{d}_{j_{0}}^{T} \mathbf{x}_{i}\right|$, go to Step 3;
Else repeat Step 6 .
Remark 3.2. The initial approximation in the algorithm is obtained by choosing a nonsingular submatrix $\mathbf{B} \in \mathbb{R}^{n \times n}$ of the matrix $\mathbf{X}$, which can be determined by applying QR factorization with column pivoting (see e.g. [2, 30]), although other approaches can be found in literature as well (see e.g. [6]).

Since in every step of the algorithm the matrix $\mathbf{B}$ changes in only one row, calculation of the inverse matrix $\mathbf{B}^{-1}$ may be simplified significantly also by applying QR factorization (see e.g. [2, 31, 32] (that is also used in our algorithm) or by applying the Sherman-Morrison formula (see e.g. [14, 6].

As mentioned previously in Remark 2.4, the set $J$ from Step 3 is not empty since at least $j_{0} \in J$. If $j_{0}$ is a unique element of the set $J$, then $\vartheta^{*}=0$. Also, matrix $\mathbf{B}$ from Step 4 always remains nonsingular.

Checking whether the system $\mathbf{X}_{0} \boldsymbol{\lambda}=\mathbf{h}$ with condition $\|\boldsymbol{\lambda}\|_{\infty} \leq 1$ has a solution is carried out by a Mathematica-instruction FindInstance, which is based upon the Buchberger's algorithm and the Gröbner system (see e. g. [33]).

## 4 Illustrative Examples and Numerical Experiments

The aforementioned algorithms will first be illustrated on a $10 \times 2$ example in which different nondegenerate and degenerate situations appear, and which could be visually well observed.
Example 4.1. Given is the system $\mathbf{X a}=\mathbf{z}, \mathbf{X} \in \mathbb{R}^{10 \times 2}, \mathbf{z} \in \mathbb{R}^{10}$, where

$$
\mathbf{X}^{T}=\left[\begin{array}{cccccccccc}
-2 & -1 & 0 & 0 & 1 & 1 & 1 & 2 & 2 & 2 \\
-1 & 0 & -3 & -1 & -2 & -1 & 1 & -3 & -2 & 0
\end{array}\right], \quad \mathbf{z}=(-5,-2,-9,-1,-2,0,4,-1,2,2)^{T}
$$



Figure 3: Iterative process
In accordance with Section 2.3 in Fig. 3 left, the system is shown by the lines marked with numbers $i=1, \ldots, 10$. In the same figure, beside each intersection of the lines, the value of the minimizing
functional $F$ is denoted, whereby a yellow polygon denotes the area on which the functional $F$ attains its global minimum.

| No | Initial <br> equations | Initial <br> appr. | Iteration <br> 1 | Iteration <br> 2 | $\boldsymbol{\lambda}^{*}\left(I_{0}\right)$ |
| :---: | :---: | :---: | :---: | :---: | :---: |
| Fig.3 | $\{4,10\}$ | $\{(1,1), 14\}$ | $\{(1,1), 14\}$ | $\left\{\left(\frac{5}{3}, \frac{5}{3}\right), 10\right\}$ | $\{-1,1\}$ |
| right | $\{4,6\}$ | $\{(1,1), 14\}$ | $\left\{\left(\frac{5}{3}, \frac{5}{3}\right), 10\right\}$ | - | $\{-1,1\}$ |
|  | $\{4,8\}$ | $\{(1,1), 14\}$ | $\left\{\left(\frac{7}{4}, \frac{3}{2}\right), 10\right\}$ | - | $\{-1,-1\}$ |
| Fig.4 | $\{3,9\}$ | $\{(4,3), 20\}$ | $\left\{\left(\frac{5}{2}, \frac{3}{2}\right), 14\right\}$ | $\{(2,2), 10\}$ | $\{1,1,1,-1\}$ |
| left | $\{3,5\}$ | $\{(4,3), 20\}$ | $\{(2,2), 10\}$ | - | $\{1,1,1,-1\}$ |
|  | $\{3,8\}$ | $\{(4,3), 20\}$ | $\left\{\left(2, \frac{5}{3}\right), 10\right\}$ | - | $\{1,-1\}$ |
| Fig.4 | $\{3,10\}$ | $\{(1,3), 20\}$ | $\left\{\left(1, \frac{3}{2}\right), 14\right\}$ | $\{(2,2), 10\}$ | $\{1,1,1,-1\}$ |
| right | $\{3,7\}$ | $\{(1,3), 20\}$ | $\{(2,2), 10\}$ | - | $\{1,1,1,-1\}$ |
|  | $\{3,1\}$ | $\{(1,3), 20\}$ | $\left\{\left(\frac{5}{3}, \frac{5}{3}\right), 10\right\}$ | - | $\{-1,1\}$ |

Table 1: Iterative process

If we choose the intersection of the lines $\{4,10\}$, i.e. the point $(1,1)$, as the initial approximation, then Algorithm I and Algorithm II cannot be run since along these lines there does not exist a smaller value of the functional $F$. Algorithm III in Step 5 detects that it is not the point of the global minimum and in this point it selects a new direction along line 6 , which after that leads to a solution in one single step (a green arrow from the point $(1,1)$ to the point $\left(\frac{5}{3}, \frac{5}{3}\right)$ ). We have a similar situation if we choose the intersection of the lines $\{4,9\}$, i.e. the point $(2,1)$, as the initial approximation. Algorithm III in Step 5 detects that it is not the point of the global minimum and in this point it selects a new direction along line 1, which after that leads to a solution in one single step (a green arrow from the point $(2,1)$ to the point $\left.\left(\frac{7}{4}, \frac{3}{2}\right)\right)$. These situations are shown in Fig. 3 right. Thereby green arrows also show the direction of optimal strategy of movement from the point $(1,1)$, i.e. from the point $(2,1)$. Corresponding data can be seen in Table 1. The column denoted by $\boldsymbol{\lambda}^{*}\left(I_{0}\right)$ shows values of parameters $\lambda_{i}$ from (33) in the optimal point.

If we choose the intersection of the lines $\{3,9\}$, i.e. the point $(4,3)$, as the initial approximation, then as a solution all algorithms give the point $(2,2)$, in which the lines $\{2,5,6,7\}$ intersect. This situation is shown in Fig. 4 left by blue arrows, and the corresponding flow of the algorithm is also given in Table 1. Green arrows show directions of optimal strategy of movement from the point $(4,3)$.

If we choose the intersection of the lines $\{3,10\}$, i.e. the point $(1,3)$, as the initial approximation, then as a solution all algorithms give again the point $(2,2)$, in which the lines $\{2,5,6,7\}$ intersect. The flow of Algorithm III is shown in Fig. 4 right by blue arrows, and it is also shown in Table 1. Green arrows show directions of optimal strategy of movement from the point $(1,3)$.

Example 4.2. Similarly to [25], for the function $f:[1,2] \rightarrow \mathbb{R}, f(x)=e^{x}+\left\{\begin{array}{ll}5, & x \in(1.2,1.4) \\ 0, & x \notin(1.2,1.4)\end{array}\right.$ we will search for a best LAD-polynomial of the $(n-1)$-th degree $P_{n-1}(x)=\alpha_{1}+\alpha_{2} x+\cdots+\alpha_{n} x^{n-1}$ on the basis of given data $\left(x_{i}, z_{i}\right), i=1, \ldots, m$, where

$$
x_{i}=1+\frac{i}{m}, \quad z_{i}=f\left(x_{i}\right)+\varepsilon_{i}, \quad \varepsilon_{i} \sim \mathcal{N}\left(0, \sigma^{2}\right), \quad \sigma=0.5 .
$$

For $n=4$ and $m=18$ the problem is reduced to searching for a best LAD-solution of the system $\mathbf{X a}=\mathbf{z}$, where $\mathbf{X}_{i j}=x_{i}^{j-1}, \quad i=1, \ldots, m, \quad j=1, \ldots n$.


Figure 4: Iterative process

| $k$ | $I_{B}$ | $\mathbf{a}_{k}^{T}$ |  |  |  | $F\left(\mathbf{a}_{k}\right)$ |
| :---: | :---: | :---: | :---: | :---: | :---: | :---: |
| 0 | $\{18,1,10,5\}$ | $(-405.091$, | 852.626, | -573.812, | $125.323)$ | 35.6137 |
| 1 | $\{1,10,5,17\}$ | $(-426.959$, | 904.515, | -614.345, | $135.746)$ | 31.6462 |
| 2 | $\{1,5,17,13\}$ | $(-294.730$, | 607.759, | -398.664, | $85.3265)$ | 22.4969 |
| 3 | $\{1,17,13,2\}$ | $(-90.766$, | 191.204, | -126.184, | $27.625)$ | 21.5901 |
| 4 | $\{1,17,2,11\}$ | $(-93.870$, | 198.533, | -131.779, | $28.9859)$ | 21.4035 |
| 5 | $\{1,17,11,14\}$ | $(-133.619$, | 281.305, | -187.208, | $41.0065)$ | 21.1829 |

Table 2: Iterative process

The flow of the iterative process and the corresponding approximate polynomials are shown in Table 2 and Fig. 5, respectively. Note that the graph of each approximate polynomial passes through 4 data points, which is in accordance with the described theory and Algorithm III.


Figure 5: Polynomial LAD-approximation of the function

Example 4.3. Algorithm III will also be tested on large systems. For that purpose we consider the problem given in Example 4.2 for $n=5,10$ and $m=50,100,200,500,1000$.

The number of iterations is shown in Table 3, whereby below every number of iterations a value of the minimizing functional obtained by Algorithm III and a value of the minimizing functional obtained by the Mathematica-module NMinimize are shown. In this way it can be seen that Algorithm III is dominant.

|  | $m=50$ | $m=100$ | $m=200$ | $m=500$ | $m=1000$ |
| :--- | ---: | ---: | ---: | ---: | ---: |
| $n=5$ | 8 | 27 | 19 | 21 | 30 |
| Algorithm III | 48.9 | 107.9 | 212.9 | 536.4 | 1077.1 |
| NMinimize | 58.8 | 119.9 | 232.6 | 572.8 | 1151.1 |
| $n=10$ | 23 | 33 | 34 | 43 | 65 |
| Algorithm III | 24.7 | 48.1 | 104.6 | 261.8 | 523.6 |
| NMinimize | 90.6 | 122.5 | 236.6 | 563.6 | 1082.4 |

Table 3: Testing Algorithm III on large systems

## 5 Concluding Remarks

In this paper we consider the problem of searching for a best LAD-solution of an overdetermined system of linear equations $\mathbf{X a}=\mathbf{z}$, where $\mathbf{X} \in \mathbb{R}^{m \times n}, m \geq n$, is a matrix of full column rank, $\mathbf{a} \in \mathbb{R}^{n}$, and $\mathbf{z} \in \mathbb{R}^{m}$ (see e.g. [2, 4, 5, 6, 7, 8]). Motivated by an efficient method for solving the problem of estimation of optimal parameters of a best LAD-plane $(x, y) \mapsto \alpha x+\beta y$ on the basis of the given set of experimental data $\left(x_{i}, y_{i}, z_{i}\right), i=1, \ldots, m$ (see e.g. $[9,10,11,12,13]$ ), which can easily be geometrically visualized, we define an iterative procedure for searching for a best LAD-solution of an overdetermined system of linear equations. For the mentioned iterative procedure we construct an appropriate algorithm and give a few illustrative examples for nondegenerate and degenerate situations. The examples with large systems show efficiency of the given method, which can be easily expanded on an overdetermined system with linear constraints. The methodology used in the paper could also be used for solving a more difficult and numerically more demanding orthogonal distance linear regression problem (see e.g. [13]).

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[^1]:    ${ }^{1}$ All evaluations and illustrations were done using Mathematica 6 on a PC (CPU: 2.00 GHz Intel Core 2 Duo processor, Memory: 1.99 GB DDR2) on the basis of our own software available at http://www.mathos.hr/~scitowsk/Algorithms.nb and http://www.mathos.hr/~scitowsk/Algorithms.m

