# For piecewise smooth signals, $l^{1}$-method is the best among $l^{p}$-methods: an interval-based justification of an empirical fact 

V. Kreinovich ${ }^{1, \star}$, A. Neumaier ${ }^{2}$<br>${ }^{1}$ University of Texas at El Paso, El Paso, TX 79968, USA<br>${ }^{2}$ University of Vienna, Vienna, A-1090, Austria<br>*Corresponding author: Kreinovich, Vladik, e-mail: vladik@utep.edu

Traditional engineering techniques often use the Least Squares Method (i.e., in mathematical terms, minimization of the $l^{2}$-norm) to process data. It is known that in many real-life situations, $l^{p}$-methods with $p \neq 2$ lead to better results, and different values of $p$ are optimal in different practical cases. In particular, when we need to reconstruct a piecewise smooth signal, the empirically optimal value of $p$ is close to 1 . In this paper, we provide a new theoretical explanation for this empirical fact based on ideas and results from interval analysis.

Keywords: piecewise smooth signal, $l^{1}$-method, interval uncertainty.

## 1. Formulation of the problem

$\boldsymbol{l}^{2}$-methods: brief reminder. Traditional engineering techniques frequently use the Least Squares Method - LSM (i.e., in mathematical terms, the $l^{2}$-norm) to process data. For example, if we know that measured values $b_{1}, \ldots, b_{m}$ are related to the unknowns $x_{1}, \ldots, x_{n}$ by the known dependence

$$
\sum_{j=1}^{n} A_{i j} x_{j} \approx b_{i}
$$

and we know the accuracy $\sigma_{i}$ of each measurement, then the LSM means that we find the values $x_{j}$ for which the function

$$
V=\sum_{i=1}^{m}\left(\frac{1}{\sigma_{i}} \sum_{j=1}^{n} A_{i j} x_{j}-b_{i}\right)^{2}
$$

takes the smallest possible value.
By the Gauss - Markov Theorem [1], this method is provably optimal (being the best linear unbiased estimator) under the assumption that the measurement errors

$$
\Delta b_{i} \stackrel{\text { def }}{=} \sum_{j=1}^{n} A_{i j} x_{j}-b_{i}
$$

[^0]are uncorrelated with zero mean and standard deviation $\sigma_{i}$. In addition, if the $\Delta b_{i}$ are independent and normally distributed, then the maximum likelihood method [2, 3], which requires $\rho\left(\Delta b_{1}, \ldots, \Delta b_{n}\right) \rightarrow \max$, takes the form
$$
\rho\left(\Delta b_{1}, \ldots, \Delta b_{n}\right)=\rho_{1}\left(\Delta b_{1}\right) \cdot \ldots \cdot \rho_{n}\left(\Delta b_{n}\right)
$$
where
$$
\rho_{i}\left(\Delta b_{i}\right)=\frac{1}{\sqrt{2 \pi} \sigma_{i}} \exp \left(-\frac{\Delta b_{i}^{2}}{2 \sigma_{i}^{2}}\right)
$$

Since the logarithm is a strictly increasing function, and the logarithm of a product $\rho_{1} \cdot \ldots \cdot \rho_{n}$ is equal to the sum of the logarithms, maximizing the maximal likelihood is equivalent to minimizing the sum of negative $\operatorname{logarithms}-\log \left(\rho_{i}\right)$ of $\rho_{i}$, i. e., minimizing the sum

$$
\begin{equation*}
\psi\left(\frac{\Delta b_{1}}{\sigma_{1}}\right)+\ldots+\psi\left(\frac{\Delta b_{n}}{\sigma_{n}}\right) \tag{1}
\end{equation*}
$$

with $\psi(x)=x^{2}$. We thus get the Least Squares Method.
Similarly, if we know that the next value $x_{i+1}$ is close to the previous value $x_{i}$ of the desired signal, and that the average difference between $x_{i+1}-x_{i}$ is about $\sigma_{i}$, then we can use LSM to find the values $x_{i}$ which minimize the sum

$$
\sum_{i=1}^{n-1}\left(\frac{x_{i+1}-x_{i}}{\sigma_{i}}\right)^{2}
$$

M-methods: brief reminder. In many practical situations, different measurement errors are independent, but the distribution may be different from normal; see, e.g., [4-6]. In this case, the maximum likelihood method is still equivalent to minimizing the sum (1), but with a different function $\psi(x)=-\log (\rho(x))$.

In many other practical situations, we know that the distribution is not normal, but we do not know its exact shape. In such situations of robust statistics, we can still use a similar method, with an appropriately selected function $\psi(x)$. Such methods are called $M$-methods; see, e. g., $[2,3,7]$.

In such situations, if we know that the next value $x_{i+1}$ is close to the previous value $x_{i}$ of the desired signal, and that the average difference between $x_{i+1}-x_{i}$ is about $\sigma_{i}$, then we can use LSM to find the values $x_{i}$ which minimize the sum

$$
\sum_{i=1}^{n-1} \psi\left(\frac{x_{i+1}-x_{i}}{\sigma_{i}}\right)
$$

$l^{p}$-methods: a brief reminder. Among different M-methods, empirically, $l^{p}$-methods - with $\psi(x)=|x|^{p}$ for some $p \geq 1$ - turn out to be the best for several practical applications; see, e. g., [8]. In this case, we select a signal ( $=$ tuple) $x_{i}$ for which the value

$$
V \stackrel{\text { def }}{=} \sum_{i=1}^{n-1}\left|\frac{x_{i}-x_{i+1}}{\sigma_{i}}\right|^{p}
$$

is the smallest possible. These methods have been successfully used to solve inverse problems in geophysics; see, e. g., [9, 10].

In [11], the empirical success of $l^{p}$-methods was theoretically explained: it was shown that $l^{p}$-methods are the only scale-invariant ones, and that they are the only methods optimal with respect to all reasonable scale-invariant optimality criteria. It is therefore reasonable to use $l^{p}$-methods for processing data.
$\boldsymbol{l}^{p}$-methods: how to select $\boldsymbol{p}$. The above-mentioned justification explains that with respect to each optimality criterion, one of the $l^{p}$-methods is optimal - but does not explain which one. It is known that in different practical situations, different values of $p$ lead to the best signal reconstruction.

For example, in the situation when the errors are normally distributed, $p=2$ is the best value. For other situations, we may get $p=1$ or $p \in] 1,2[$.

In each situation, we must therefore empirically select $p-$ e. g., by comparing the result of data processing with the actual (measured) values of the reconstructed quantity.

Empirical fact. In several situations, we know that the reconstructed signal is piecewise smooth. For example, in geophysics, the Earth consists of several layers with abrupt transition between layers; in image processing, an image often consists of several zones with an abrupt boundary between the zones, etc.

It turns out that in many such situations, the empirically optimal value of $p$ is close to 1 ; see, e. g., [9] for the inverse problem in geophysics, and [12-15] for image reconstruction.

How this fact is explained now (see, e. g., [12]). In the continuous approximation, the $l^{p}$-criterion leads to the minimization of $\int|\dot{x}|^{p} d t$ (in the 1D case; multidimensional case is handled similarly). For a transition of magnitude $C$ and width $\varepsilon$, the derivative $\dot{x}$ is $\approx C / \varepsilon$, so the contribution of the transition zone to the integral is of order $\varepsilon / \varepsilon^{p}=\varepsilon^{-(p-1)}$. For $p>1$, when $\varepsilon \rightarrow 0$, this contribution tends to $\infty$. Thus, for $p>1$, the minimum is never attained at the discontinuous transition ("jump") $\varepsilon=0$, but always at a smoother transition $\varepsilon>0$.

For $p=1$, the contribution is finite, so jumps are not automatically excluded - and indeed, they may be correctly reconstructed.

Limitations of our explanation. There are two limitations to this explanation:

- first, it explains why $l^{p}$-methods for $p>1$ do not reconstruct a jump, but it does not explain why $l^{1}$ methods reconstruct the jump correctly;
- second, it strongly relies on the continuous case and does not fully explain why a similar phenomenon occurs for real-life (discretized) computations.
What we do in this paper. In this paper, we provide a new interval-based theoretical explanation for the above mentioned empirical fact, an explanation that is directly applicable to real-life (discretized) computations.


## 2. Analysis of the problem and the main results

For simplicity, we will consider 1-D signals $x(t)$. In the interval setting, for several moments of time $t_{1}<\ldots<t_{n}$ (usually, equidistant $t_{i}=t_{1}+(i-1) \Delta t$ ), we know the intervals $\boldsymbol{x}_{i}=\left[\underline{x}_{i}, \bar{x}_{i}\right]$ that contain the actual (unknown) values $x_{i}=x\left(t_{i}\right)$. Based on this interval information, we would like to select the values $x_{i} \in \boldsymbol{x}_{i}$. According to the $l^{p}$-criterion, among all the tuples $\left(x_{1}, \ldots, x_{n}\right)$ for which $x_{1} \in \boldsymbol{x}_{1}, \ldots, x_{n} \in \boldsymbol{x}_{n}$, we select the one for which the value

$$
V=\sum_{i=1}^{n-1}\left|\frac{x_{i}-x_{i+1}}{\sigma_{i}}\right|^{p}
$$

is the smallest possible.

To select $p$, we will consider the case of a "transition zone", i. e., the case when for some values $l<u$, we know two things:

- that the value $x_{l-1}$ right before the zone cannot be equal to the value $x_{u+1}$ right after the zone - i. e., that $\boldsymbol{x}_{l-1} \cap \boldsymbol{x}_{u+1}=\emptyset$; and
- that we have (practically) no information about the values of $x_{i}$ inside the zone - i. e., at least that for all $i$ from $l$ to $u$, the interval $\boldsymbol{x}_{i}$ contains both $\boldsymbol{x}_{l-1}$ and $\boldsymbol{x}_{u+1}$.
In this case, the above criterion interpolates the values $x_{i}$ inside the zone. If we assumed that the signal is smooth, then, no matter how steep the transition, we would have had a smooth interpolation. However, since we consider the situations when the signal is only piecewise smooth, we would rather prefer to have a signal which "jumps" discontinuously from one value to another.

In this section, we will show that for $p=1$, we will indeed get such a jump, while for $p>1$, we have a smooth transition instead. Let us describe this result in precise terms.

Definition 1. By an $l^{p}$-problem, we mean the following problem:
GIVEN: $n$ intervals $\boldsymbol{x}_{i}=\left[\underline{x}_{1}, \bar{x}_{1}\right], \ldots,\left[\underline{x}_{n}, \bar{x}_{n}\right]$, $n$ real numbers $\sigma_{1}, \ldots, \sigma_{n}$, and a real number $p \geq 1$;
AMONG: tuples $x_{1} \ldots, x_{n}$ such that $x_{i} \in\left[\underline{x}_{i}, \bar{x}_{i}\right]$ for every $i$;
FIND: the tuple for which the value $V=\sum_{i=1}^{n-1}\left|\frac{x_{i}-x_{i+1}}{\sigma_{i}}\right|^{p}$ is the smallest possible.
Definition 2. An $l^{p}$-problem is called degenerate if all the values $\sigma_{i}$ are different.
Comment. Almost all combinations $\sigma_{1}, \ldots, \sigma_{n}$ are degenerate.
Definition 3. Let $l<u$ be integers. We say that an $l^{p}$-problem contains a transition zone between $l$ and $u$ if the following two conditions hold (Fig. 1):

- $\boldsymbol{x}_{l-1} \cap \boldsymbol{x}_{u+1}=\emptyset$; and
- for all $i$ from $l$ to $u$, we have $\boldsymbol{x}_{i} \supseteq \boldsymbol{x}_{l-1}$ and $\boldsymbol{x}_{i} \supseteq \boldsymbol{x}_{u+1}$.

Proposition 1. For $p=1$, for each solution $x_{i}$ to a non-degenerate $l^{p}$-problem, in each transition zone, we have $x_{l-1}=x_{l}=\ldots=x_{t}$ and $x_{t+1}=\ldots=x_{u}=x_{u+1}$ for some $t$.

In other words, for $p=1$, in each transition zone, we have a "jump" from the value $x_{l-1}$ before the transition zone to the value $x_{u+1}$ after the transition zone (Fig. 2).

Comment. In the degenerate case, when different values $\sigma_{i}$ are equal, the jump is still an optimal solution, but we may also get other solutions, with a smooth transition from


Fig. 1


Fig. 2
$x_{l-1}$ to $x_{u+1}$. For example, if all the values $\sigma_{i}$ are the same, then, as one can easily see, the minimized criterion is proportional to the sum

$$
\sum_{i=1}^{n-1}\left|x_{i}-x_{i+1}\right| .
$$

So, for each solution that monotonically changes from $x_{l-1}$ to $x_{u-1}$, the corresponding part

$$
\sum_{i=l-1}^{u}\left|x_{i}-x_{i+1}\right|
$$

of the sum is equal to $\left|x_{l-1}-x_{u+1}\right|$. Thus, the value of the minimized criterion is the same for the jump solution and for a different solution in which $x_{i}$ is the same outside $[l-1, u+1]$ but strictly monotonically changes between $l-1$ and $u+1$.

Proposition 2. For $p>1$, for each solution $x_{i}$ to an $l^{p}$-problem, in each transition zone, we have a strictly monotonic sequence $x_{l-1}<x_{l}<\ldots<x_{u}<x_{u+1}$ or $x_{l-1}>x_{l}>$ $\ldots>x_{u}>x_{u+1}$.

Proposition 3. For $p>1$, in the limit when all the values $\sigma_{i}$ tend to the same value $\sigma$, each solution $x_{i}$ to an $l^{p}$-problem, in each transition zone, is linear, i.e., has the form $x_{i}=a+b i$ for some numbers $a$ and $b$ (Fig. 3).

These results explain why $p \approx 1$ is indeed empirically best for processing piecewise smooth signals: only for $p=1, l^{p}$-interpolation leads to a piecewise smooth signal.

Comment. The fact that $l^{1}$-methods are the best among $l^{p}$-methods does not mean that they are always the best possible interpolation techniques. For example, the above results show that, with an $l^{1}$-method, we always get a jump, both

- for the steep transition from $\boldsymbol{x}_{l-1}$ to $\boldsymbol{x}_{u+1}$, where such a jump is desirable, and
- for a smoother transition from $\boldsymbol{x}_{l-1}$ to $\boldsymbol{x}_{u+1}$, where, from the physical viewpoint, we may want to prefer a smooth interpolation.
In other words,
- for small differences $x_{i}-x_{i+1}$, we would like to have smooth transitions, while
- for large differences $x_{i}-x_{i+1}$, we would like to have a jump.

Since a jump is reconstructed when $\psi(x)=|x|$ and a smooth transition, when, e. g., $\psi(x)=$ $|x|^{2}$, a natural idea is to use a Huber function $\psi(x)$ which is equal to $|x|^{2}$ when $|x|$ is below


Fig. 3
a certain threshold $x_{0}$, and which is linear $\psi(x)=C \cdot|x|$ for $|x|>x_{0}$; from the requirement that the function $\psi(x)$ be continuous, we conclude that $C=\left|x_{0}^{2}\right|=C \cdot\left|x_{0}\right|$, i. e., that $C=x_{0}$. Such technique indeed leads to a better reconstruction of piecewise smooth signals; see, e. g. [12] and references therein. Various related choices for $\psi(x)$ have been explored in the context of computer tomography by Kaufman and Neumaier [16, 17].

Huber's function $\psi(x)$, in its turn, has its own limitations; it is worth mentioning that in general, the problem of optimally reconstructing piecewise smooth 2-D signals is NP-hard; see, e. g. [18-20].

## 3. Proofs

1. First, we observe that the solution to an $l^{p}$-problem minimizes a continuous function $V$ on a bounded closed set (box) $\boldsymbol{x}_{1} \times \ldots \times \boldsymbol{x}_{n}$. Thus, this minimum is always attained, i.e., a solution always exists.
2. Let us prove that for every $p$, the solution $x_{i}$ to the $l^{p}$-problem is (non-strictly) monotonic in each transition zone, i.e., that $x_{l-1} \leq x_{l} \leq \ldots \leq x_{u} \leq x_{u+1}$ or $x_{l-1} \geq x_{l} \geq$ $\ldots \geq x_{u} \geq x_{u+1}$.

Let us prove this result by reduction to a contradiction. Namely, let us assume that the solution is attained on some non-monotonic sequence. The fact that $x_{i}$ is not monotonic on the transition zone means that not all inequalities between the neighboring values are of the same sign, i. e., that we have $x_{i-1}<x_{i}$ and $x_{j}>x_{j+1}$ for some indices $i$ and $j$ from this zone. Among such pairs $(i, j)$, let us select a one with the smallest distance $|i-j|$ between $i$ and $j$.

Without losing generality, we can assume that $i<j$ in this selected pair.
For the selected pair, for indices $k$ between $x_{i}$ and $x_{j}$, we cannot have $x_{k}<x_{k+1}$ or $x_{k}>x_{k+1}$ - otherwise we would get a pair with an even smaller difference $|i-j|$. Thus, for all intermediate indices $k$, we get $x_{k}=x_{k+1}$. Since $x_{i}=x_{i+1}=\ldots=x_{j}$, we thus have $x_{i}=x_{j}$. So, we have $x_{i-1}<x_{i}=\ldots=x_{j}>x_{j+1}$. Let $\varepsilon=\min \left(x_{i}-x_{i-1}, x_{j}-x_{j+1}\right)$. Let us now keep all the $x$-values outside $(i, j)$ intact and replace $x_{i}=\ldots=x_{j}$ with the values $x_{i}-\varepsilon=\ldots=x_{j}-\varepsilon$. The resulting value $x_{i}-\varepsilon$ is equal to either $x_{i-1} \in \boldsymbol{x}_{i-1}$ or to $x_{j+1} \in \boldsymbol{x}_{j+1}$. By the definition of a transition zone, all intermediate intervals $\boldsymbol{x}_{k}$ contain both $\boldsymbol{x}_{i-1}$ and $\boldsymbol{x}_{j+1}$. Hence, the new value of $x_{k}$ is within the corresponding interval $\boldsymbol{x}_{k}$.

By making this change, we decrease the differences $\left|x_{i}-x_{i-1}\right|$ and $\left|x_{j+1}-x_{j}\right|$ and leave all other differences intact - and hence, we decrease the value of the minimized objective function $V$ (Fig. 4).

Since the objective function $V$ attains its minimum at the original tuple $x_{i}$, the possibility to minimize even further is a contradiction. This proves that the solution is monotonic in each transition zone.
3. For the solution, we have $x_{i-1} \leq x_{i} \leq \ldots \leq x_{u} \leq x_{u+1}$ or $x_{i-1} \geq x_{i} \geq \ldots \geq x_{u} \geq x_{u+1}$ according to Part 2 of this proof. To complete the proof of Proposition 2, it is now sufficient to prove that for $p=1$ and for $k=l, \ldots, u$, we cannot have any strictly intermediate values $x_{k} \in\left(x_{l-1}, x_{u+1}\right)$ (or $\left.x_{k} \in\left(x_{u+1}, x_{l-1}\right)\right)$.

We will prove this ad absurdum. Let us assume that an intermediate value $x_{k}$ does exist. In principle, we may have values equal to $x_{k}$. Due to monotonicity, these values form an interval within $[l, u]$. Let $x_{b}$ be the first value equal to $x_{k}$, and let $x_{e}$ be the last value equal to $x_{k}$. Then, we have $\ldots \leq x_{b-1}<x_{b}=\ldots=x_{e}<x_{e+1} \leq \ldots$


Fig. 4


Fig. 5

Let us now choose a value $\varepsilon \in\left[x_{b-1}-x_{b}, x_{e+1}-x_{e}\right]$, keep all the $x$-values from outside [ $b, e]$ intact, and replace all the $x$-values from $[b, e]$ with $x_{b}+\varepsilon=\ldots=x_{e}+\varepsilon$. Similarly to Part 2 of this proof, we can show that for every $\varepsilon$, we still have $x_{b}+\varepsilon \in \boldsymbol{x}_{b}, \ldots, x_{e}+\varepsilon \in \boldsymbol{x}_{e}$.

After this replacement, we change only two differences $\left|x_{i+1}-x_{i}\right|$ :

- the difference $\left|x_{b}-x_{b-1}\right|=x_{b}-x_{b-1}$ is replaced with $x_{b}-x_{b-1}+\varepsilon$, and
- the difference $\left|x_{e+1}-x_{e}\right|=x_{e+1}-x_{e}$ is replaced with $x_{e+1}-x_{e}-\varepsilon$.

Thus, after this replacement, the original value $V$ of the minimized objective function is replaced with $V+\Delta V$, where (Fig. 5)

$$
\Delta V \stackrel{\text { def }}{=} \varepsilon\left(\frac{1}{\sigma_{b-1}}-\frac{1}{\sigma_{e}}\right) .
$$

Since the problem is non-degenerate, i. e., all the values $\sigma_{i}$ are different, the coefficient at $\varepsilon$ in $\Delta V$ is non-zero. If this coefficient is positive, we can take negative $\varepsilon$ and decrease $V$; if it is negative, we can decrease $V$ by taking $\varepsilon>0$. In both cases, we get a contradiction with the fact that the original tuple $x_{i}$ minimizes $V$. This contradiction proves that intermediate values are impossible. Proposition 2 is proven.
4. Let us now prove that the solution is strictly monotonic for $p>1$, using reduction to a contradiction once again.

We assume that the solution is not strictly monotonic, while usual monotonicity holds (Part 2 of the proof). Since it is monotonic, the only way for the solution to be not strictly monotonic is to have $x_{i}=x_{i+1}$ for some index $i$. We may have several indices with an $x$-value equal to this $x_{i}$; let $b$ be the first such index, and let $e$ be the last such index. Then, $x_{b}=x_{b+1}=\ldots=x_{e}$.

Since the intervals $\boldsymbol{x}_{l-1}$ and $\boldsymbol{x}_{u+1}$ have no common points, we cannot have $x_{l-1}=x_{u+1}$. Thus, either $b \neq l-1$ or $e \neq u+1$. Without losing generality, we can assume that $b \neq l-1$. Also, without losing generality, we can assume that the solution $x_{i}$ is increasing. Thus, we have $x_{b-1}<x_{b}=x_{b+1}$.

Let us now pick a small value $\varepsilon>0$ and replace $x_{b}$ with $x_{b}-\varepsilon$ - while leaving all other $x$-valued intact (Fig. 6).


Fig. 6

This replacement changed the original value $V$ of the minimized function with a new value $V+\Delta V$, where

$$
\Delta V=\frac{\left(x_{b}-x_{b-1}-\varepsilon\right)^{p}}{\sigma_{b-1}^{p}}+\frac{\varepsilon^{p}}{\sigma_{b}^{p}}-\frac{\left(x_{b}-x_{b-1}\right)^{p}}{\sigma_{b-1}^{p}}
$$

By applying the first term of Taylor expansion to the first ratio in the expression for $\Delta V$, one can conclude that

$$
\Delta V=-\frac{p\left(x_{b}-x_{b-1}\right)^{p-1}}{\sigma_{b-1}^{p}} \varepsilon+O\left(\varepsilon^{2}\right)+\frac{\varepsilon^{p}}{\sigma_{b}^{p}} .
$$

We consider the case $p>1$. Then, for sufficiently small $\varepsilon$, the first term dominates, so the difference $\Delta V$ is negative - which means that we can further decrease $V$.

This possibility contradicts to the fact that the tuple $x_{i}$ minimizes $V$. Thus, the solution is indeed strictly monotonic. Proposition 2 is proven.
5. Let us now prove Proposition 3.

By definition of the transition zone, for each index $i$ from this zone, we have $\boldsymbol{x}_{l-1} \subseteq \boldsymbol{x}_{i}$, hence $x_{l-1} \in \boldsymbol{x}_{l-1} \subseteq \boldsymbol{x}_{i}$ and $x_{l-1} \in\left[\underline{x}_{i}, \bar{x}_{i}\right]$ thence $\underline{x}_{i} \leq x_{l-1}$. Similarly, from $\boldsymbol{x}_{u+1} \subseteq \boldsymbol{x}_{i}$, we conclude that $x_{u+1} \leq \bar{x}_{i}$.

Due to strict monotonicity (Part 4 of this proof), we have $x_{l-1}<x_{i}<x_{u+1}$. Thus, $\underline{x}_{i} \leq x_{l-1}<x_{i}$ and $\underline{x}_{i}<x_{i}$ and similarly, $x_{i}<\bar{x}_{i}$.

Since the value $x_{i}$ is strictly inside the interval $\boldsymbol{x}_{i}$, the derivative of the minimized function $V$ is equal to 0 . Differentiating $V$ relative to $x_{i}$ (and taking monotonicity into account), we conclude that

$$
p\left(x_{i}-x_{i-1}\right)^{p-1} \sigma_{i-1}^{p}-p\left(x_{i+1}-x_{i}\right)^{p-1} \sigma_{i}^{p}=0 .
$$

When $\sigma_{i} \rightarrow \sigma$, we get $x_{i}-x_{i-1}=x_{i+1}-x_{i}$. So, the difference $x_{i}-x_{i-1}$ is indeed the same for all $i$ within the transition zone. Thus, we get the desired linear dependence of $x_{i}$ on $i$. The proposition is proven.

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