

# For Piecewise Smooth Signals, $l^1$ Method Is the Best Among $l^p$ : An Interval-Based Justification of an Empirical Fact

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## Abstract

Traditional engineering techniques use the Last Squares method (i.e., in mathematical terms, the  $l^2$ -norm) to process data. It is known that in many practical situations,  $l^p$ -methods with  $p \neq 2$  lead to better results. In different practical situations, different values of  $p$  are optimal. It is known that in several situations 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 interval-based theoretical explanation for this empirical fact.

**Keywords:** piecewise smooth,  $l^1$  method, interval uncertainty

## 1 Formulation of the Problem

**$l^2$ -methods: brief reminder.** Traditional engineering techniques use the Last Squares Method LSM (i.e., in mathematical terms, the  $l^2$ -norm) to process data. For example, if we know that measured values  $b_1, \dots, b_m$  are related to the unknowns  $x_1, \dots, x_n$  by the known dependence  $\sum_{i=1}^n A_{ij}x_j \approx b_i$ , 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

$$V = \sum_{i=1}^m \left( \frac{1}{\sigma_i} \cdot \sum_{j=1}^n A_{ij} x_j - b_i \right)^2$$

takes the smallest possible value.

By the Gauss-Markov Theorem [20], this method is provably optimal (the best linear unbiased estimator) under the assumption that the measurement errors  $\Delta b_i \stackrel{\text{def}}{=} \sum_{j=1}^n A_{ij} x_j - b_i$  are uncorrelated with 0 mean and standard deviation  $\sigma_i$ . In addition, if the  $\Delta b_i$  are independent and normally distributed, the maximum likelihood method [16, 18]  $\rho(\Delta b_1, \dots, \Delta b_n) \rightarrow \max$  takes the form

$$\rho(\Delta b_1, \dots, \Delta b_n) = \rho_1(\Delta b_1) \cdot \dots \cdot \rho_n(\Delta b_n)$$

$$\text{and } \rho_i(\Delta b_i) = \frac{1}{\sqrt{2\pi}\sigma_i} \exp\left(-\frac{\Delta b_i^2}{2\sigma_i^2}\right).$$

Since the logarithm is a strictly increasing function, and the logarithm of a product  $\rho_1 \cdot \dots \cdot \rho_n$  is equal to the sum of the logarithms, maximizing the maximal likelihood is equivalent to minimizing the sum of negative logarithms  $-\log(\rho_i)$  of  $\rho_i$ , i.e., minimizing the sum

$$\psi\left(\frac{\Delta b_1}{\sigma_1}\right) + \dots + \psi\left(\frac{\Delta b_n}{\sigma_n}\right) \quad (1)$$

with  $\psi(x) = x^2$  - i.e., to 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  $V = \sum_{i=1}^{n-1} \left( \frac{x_{i+1} - x_i}{\sigma_i} \right)^2$ .

**M-methods: brief reminder.** In many practical situations, different measurement errors are independent, but the distribution may be different from normal; see, e.g., [13, 14, 15]. 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., [8, 16, 18].

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

$$V = \sum_{i=1}^{n-1} \psi \left( \frac{x_{i+1} - x_i}{\sigma_i} \right).$$

**$l^p$ -methods: 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., [4]. 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., [6, 17].

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.

**$l^p$ -methods: how to select  $p$ .** 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., [6] for the inverse problem in geophysics, and [1, 7, 12, 19] for image reconstruction.

**How this fact is explained now** (see, e.g., [1]). In the continuous approximation, the  $l^p$ -criterion leads to the minimization of  $\int |\dot{x}|^p dt$  (in the 1D case; multi-D 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 this 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 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  $t_1 < \dots < t_n$  (usually, equidistant  $t_i = t_1 + (i-1)\Delta t$ ), we know the intervals  $\mathbf{x}_i = [\underline{x}_i, \bar{x}_i]$  that contain the actual (unknown) values  $x_i = x(t_i)$ . Based on this interval information, we would like to select the values  $x_i \in \mathbf{x}_i$ . According to the  $l^p$ -criterion, among all the tuples  $(x_1, \dots, x_n)$  for which  $x_1 \in \mathbf{x}_1, \dots, x_n \in \mathbf{x}_n$ , 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  $\mathbf{x}_{l-1} \cap \mathbf{x}_{u+1} = \emptyset$ ; 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  $\mathbf{x}_i$  contains both  $\mathbf{x}_{l-1}$  and  $\mathbf{x}_{u+1}$ .

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  $\mathbf{x}_i = [\underline{x}_i, \bar{x}_i]$ ,  $\dots$ ,  $[\underline{x}_n, \bar{x}_n]$ ,  $n$  real numbers  $\sigma_1, \dots, \sigma_n$ , and a real number  $p \geq 1$ ;*

AMONG: tuples  $x_1, \dots, x_n$  such that  $x_i \in [\underline{x}_i, \bar{x}_i]$  for every  $i$ ;

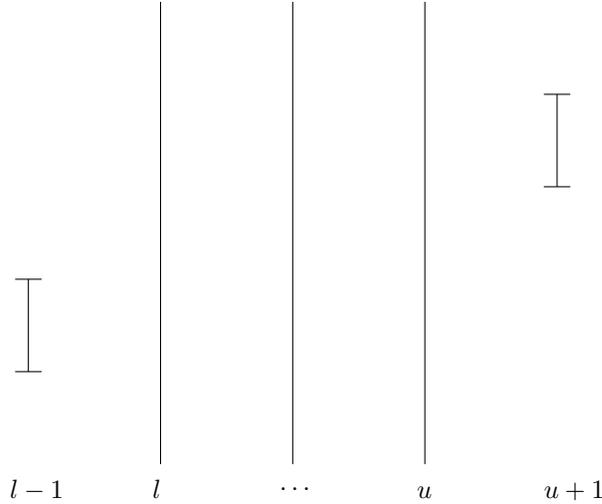
FIND: the tuple for which the value  $V = \sum_{i=1}^{n-1} \left| \frac{x_i - x_{i+1}}{\sigma_i} \right|^p$  is the smallest possible.

**Definition 2** An  $l^p$ -problem is called degenerate if all the values  $\sigma_i$  are different.

*Comment.* Almost all combinations  $\sigma_1, \dots, \sigma_n$  are degenerate.

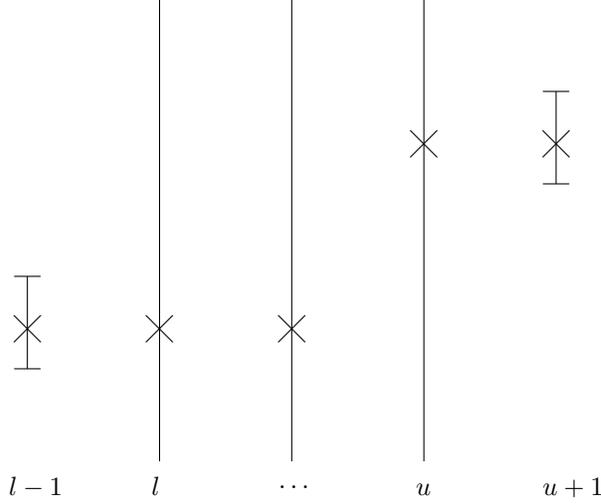
**Definition 3** Let  $l < u$  be integers. We say that an  $l^p$ -problem contains a transition zone between  $l$  and  $u$  if the following two conditions hold:

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



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

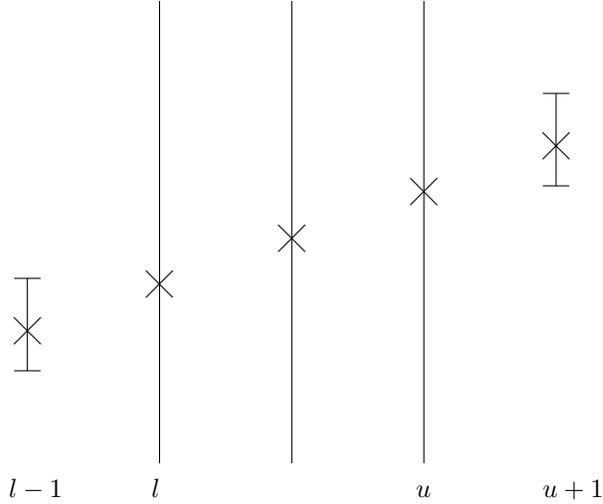
In other words, for  $p = 1$ , in each transition zone, we have a “jump” from the value  $x_{l-1}$  before the transition zone to the value  $x_{u+1}$  after the transition zone.



*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  $x_{l-1}$  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  $\sum_{i=1}^{n-1} |x_i - x_{i+1}|$ ; so, for each solution that monotonically changes from  $x_{l-1}$  to  $x_{u-1}$ , the corresponding part  $\sum_{i=l-1}^u |x_i - x_{i+1}|$  of the sum is equal to  $|x_{l-1} - x_{u+1}|$ . 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  $x_{l-1} < x_l < \dots < x_u < x_{u+1}$  or  $x_{l-1} > x_l > \dots > x_u > x_{u+1}$ .

**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  $x_i = a + bi$  for some numbers  $a$  and  $b$ .



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  $\mathbf{x}_{l-1}$  to  $\mathbf{x}_{u+1}$ , where such a jump is desirable, and
- for a smoother transition from  $\mathbf{x}_{l-1}$  to  $\mathbf{x}_{u+1}$ , 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 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  $C = |x_0^2| = C \cdot |x_0|$ , i.e., that  $C = x_0$ . Such technique indeed leads to a better reconstruction of piecewise smooth signals; see, e.g., [1] and

references therein. Various related choices for  $\psi(x)$  have been explored in the context of computer tomography by Kaufman and Neumaier [9, 10].

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., [2, 3, 5].

### 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)  $\mathbf{x}_1 \times \dots \times \mathbf{x}_n$ . Thus, this minimum is always attained, i.e., a solution always exists.

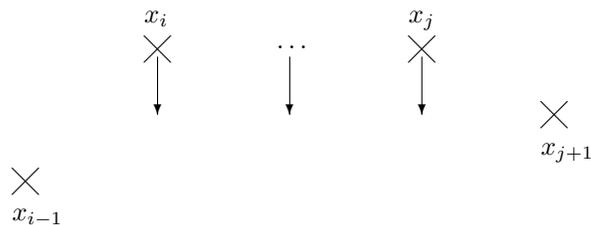
2°. Let us prove that for every  $p$ , the solution  $x_i$  to the  $l^p$ -problem is (non-strictly) monotonic in each transition zone, i.e., that  $x_{l-1} \leq x_l \leq \dots \leq x_u \leq x_{u+1}$  or  $x_{l-1} \geq x_l \geq \dots \geq x_u \geq x_{u+1}$ .

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 in this selected pair,  $i < j$ .

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  $x_i = x_{i+1} = \dots = x_j$ , we thus have  $x_i = x_j$ . So, we have  $x_{i-1} < x_i = \dots = x_j > x_{j+1}$ . Let  $\varepsilon = \min(x_i - x_{i-1}, x_j - x_{j+1})$ . Let us now keep all the  $x$ -values outside  $(i, j)$  intact and replace  $x_i = \dots = x_j$  with the values  $x_i - \varepsilon = \dots = x_j - \varepsilon$ . The resulting value  $x_i - \varepsilon$  is equal to either  $x_{i-1} \in \mathbf{x}_{i-1}$  or to  $x_{j+1} \in \mathbf{x}_{j+1}$ . By the definition of a transition zone, all intermediate intervals  $\mathbf{x}_k$  contain both  $\mathbf{x}_{i-1}$  and  $\mathbf{x}_{j+1}$ ; hence, the new value of  $x_k$  is within the corresponding interval  $\mathbf{x}_k$ .

By making this change, we decrease the differences  $|x_i - x_{i-1}|$  and  $|x_{j+1} - x_j|$  and leave all other difference intact - and hence, we decrease the value of the minimized objective function  $V$ .



Since the objective function  $V$  attains its minimum at the original tuple  $x_i$ , the possibility to minimize even further is a contradiction. This contradiction proves that the solution is monotonic in each transition zone.

3°. According to Part 2 of this proof, for the solution, we have  $x_{i-1} \leq x_i \leq \dots \leq x_u \leq x_{u+1}$  or  $x_{i-1} \geq x_i \geq \dots \geq x_u \geq x_{u+1}$ . To complete the proof of Proposition 1, it is now sufficient to prove that for  $p = 1$  and for  $k = l, \dots, u$ , we cannot have any strictly intermediate values  $x_k \in (x_{l-1}, x_{u+1})$  (or  $x_k \in (x_{u+1}, x_{l-1})$ ).

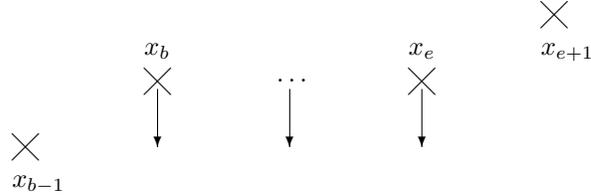
Let us prove this by reduction to a contradiction. Let us assume that an intermediate value  $x_k$  does exist. In principle, we may have values equal to  $x_k$ . Due to monotonicity, these values form an interval within  $[l, u]$ . Let  $x_b$  be the first value equal to  $x_k$ , and let  $x_e$  be the last value equal to  $x_k$ . Then, we have  $\dots \leq x_{b-1} < x_b = \dots = x_e < x_{e+1} \leq \dots$ .

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

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

- the difference  $|x_b - x_{b-1}| = x_b - x_{b-1}$  is replaced with  $x_b - x_{b-1} + \varepsilon$ , and
- the difference  $|x_{e+1} - x_e| = x_{e+1} - x_e$  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  $\Delta V \stackrel{\text{def}}{=} \varepsilon \cdot \left( \frac{1}{\sigma_{b-1}} - \frac{1}{\sigma_e} \right)$ .



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 1 is proven.

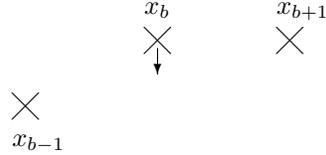
4°. Let us now prove that for  $p > 1$ , the solution is strictly monotonic.

We will prove this by reduction to a contradiction. Let us assume that it is not strictly monotonic. We know that the solution is monotonic (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  $e$  the last such index. Then,  $x_b = x_{b+1} = \dots = x_e$ .

Since the intervals  $\mathbf{x}_{l-1}$  and  $\mathbf{x}_{u+1}$  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.



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

$$\Delta V = \frac{(x_b - x_{b-1} - \varepsilon)^p}{\sigma_{b-1}^p} + \frac{\varepsilon^p}{\sigma_b^p} - \frac{(x_b - x_{b-1})^p}{\sigma_{b-1}^p}.$$

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

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

We consider the case when  $p > 1$ . In this case, 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  $\mathbf{x}_{l-1} \subseteq \mathbf{x}_i$ , hence  $x_{l-1} \in \mathbf{x}_{l-1} \subseteq \mathbf{x}_i$  and  $x_{l-1} \in [\underline{x}_i, \bar{x}_i]$  – thence  $\underline{x}_i \leq x_{l-1}$ . Similarly, from  $\mathbf{x}_{u+1} \subseteq \mathbf{x}_i$ , 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  $\mathbf{x}_i$ , 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

$$p \cdot (x_i - x_{i-1})^{p-1} \sigma_{i-1}^p - p \cdot (x_{i+1} - x_i)^{p-1} \sigma_i^p = 0.$$

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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