From Gauging Accuracy of Quantity Estimates to Gauging Accuracy and Resolution of Measuring Physical Fields

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Abstract. For a numerical physical quantity v, because of the measurement imprecision, the measurement result \widetilde{v} is, in general, different from the actual value v of this quantity. Depending on what we know about the measurement uncertainty $\Delta v \stackrel{\text{def}}{=} \widetilde{v} - v$, we can use different techniques for dealing with this imprecision: probabilistic, interval, etc.

When we measure the values v(x) of physical fields at different locations x (and/or different moments of time), then, in addition to the same measurement uncertainty, we also encounter another type of *localization* uncertainty: that the measured value may come not only from the desired location x, but also from the nearby locations.

In this paper, we discuss how to handle this additional uncertainty.

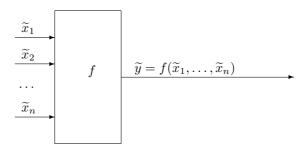
1 Need for Data Processing

Idea. In many real-life situations, we are interested in the value of a quantity which is difficult (or even impossible) to measure directly. For example, we may be interested:

- in the distance to a faraway star, or
- in the amount of water in an underground water layer.

Since we cannot measure the corresponding quantity y directly, we measure it indirectly. Specifically,

- we find easier-to-measure quantities x_1, \ldots, x_n which are related to the desired quantity y by a known dependence $y = f(x_1, \ldots, x_n)$;
- we measure the values of the auxiliary quantities x_1, \ldots, x_n ; and
- we use the results $\widetilde{x}_1, \ldots, \widetilde{x}_n$ of measuring the auxiliary quantity to compute the estimate $\widetilde{y} = f(\widetilde{x}_1, \ldots, \widetilde{x}_n)$ for the desired quantity y.



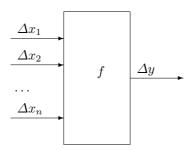
Example. To find the distance y to a faraway star, we can use the following parallax method:

- we measure the orientations x_1 and x_2 to this star at two different seasons,
- we measure the the distance x_3 between the spatial locations of the corresponding telescopes at these two seasons (i.e., in effect, the diameter of the earth orbit);
- then, reasonably simply trigonometric computations enable us to describe the desired distance y as a function of the easier-to-measure quantities x_1 , x_2 , and x_3 .

General case. In general, computations related to such indirect measurements form an important particular case of data processing.

2 Need to Take Uncertainty Into Account

Measurements are never absolutely accurate. As a result, the measurement results \widetilde{x}_i are, in general, different from the actual (unknown) values x_i of the measured quantities: $\Delta x_i \stackrel{\text{def}}{=} \widetilde{x}_i - x_i \neq 0$. Because of this, the result $\widetilde{y} = f(\widetilde{x}_1, \dots, \widetilde{x}_n)$ of data processing is, in general, different from the actual (unknown) value $y = f(x_1, \dots, x_n)$: $\Delta y \stackrel{\text{def}}{=} \widetilde{y} - y \neq 0$.



Thus, in practical applications, we need to take this uncertainty into account.

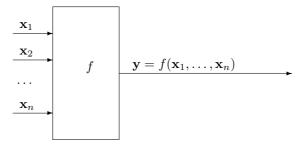
3 Interval Uncertainty

In practice, we often only know the upper bound Δ_i on the measurement errors $\Delta x_i \stackrel{\text{def}}{=} \widetilde{x}_i - x_i$: $|\Delta x_i| \leq \Delta_i$. In this case, the only information that we have about the actual values x_i is that x_i belongs to the interval $\mathbf{x}_i \stackrel{\text{def}}{=} [\widetilde{x}_i - \Delta_i, \widetilde{x}_i + \Delta_i]$.

Under such interval uncertainty, we need to find the range of possible values of y:

$$\mathbf{y} = \{ f(x_1, \dots, x_n) : x_i \in \mathbf{x}_i \}.$$

The problem of computing this range is known as *interval computations*; see, e.g., [1].



4 Need to Measure Physical Fields

In practice, the situation is often more complex: the values that we measure can be:

- values v(t) of a certain dynamic quantity v at a certain moment of time t
- or, more generally, the values v(x,t) of a certain physical field v at a certain location x and at a certain moment of time t.

For example, in geophysics, we are interested in the values of the density at different locations and at different depth.

5 Need to Take Uncertainty Into Account When Measuring Physical Fields

When we measure physical fields,

- not only we get the measured value $\tilde{v} \approx v$ with some inaccuracy, but
- also the location x is not exactly known.

Moreover, the sensor picks up the "averaged" value of v at locations close to the approximately known location \widetilde{x} .

In other words,

- in addition to inaccuracy $\tilde{v} \neq v$,
- we also have a finite resolution $\tilde{x} \neq x$.

6 Estimating Uncertainty Related to Measuring Physical Fields: Challenging Problems

In general, the measured value \widetilde{v}_i differs from the averaged value v_i by the measurement imprecision $\Delta v_i = \widetilde{v}_i - v_i$. In the interval case, we know the upper bound Δ_i on this measurement error $|\Delta v_i| \leq \Delta_i$. Thus, the averaged quantity v_i can take any value from the interval $[\underline{v}_i, \overline{v}_i]$, where $\underline{v}_i \stackrel{\text{def}}{=} \widetilde{v}_i - \Delta_i$ and $\overline{v}_i \stackrel{\text{def}}{=} \widetilde{v}_i + \Delta_i$.

Based on these bounds on v_i , what can we learn about the original field v(x)? The answer to this questions depends on what we know about the averaging, i.e., on the dependence of v_i on v(x).

In principle, there are three possible situations:

- sometimes, we know exactly how the averaged values v_i are related to v(x);
- sometimes, we only know the upper bound δ on the location error $\tilde{x} x$ (this is similar to the interval case);
- sometimes, we do not even know δ .

In the following sections, we describe how to process all these types of uncertainty.

7 Possibility of Linearization

Sometimes, the signal v(x) that we are measuring is large, i.e., the values of the signal are much larger than the noise (and the measurement errors in general). In such situations, the measured values well represent the actual signal, and for many applications, the measurement errors can be safely ignored.

The need to take into account measurement errors becomes important only when the signal v(x) is relatively weak. In this case, we can expand the dependence of v_i on v(x) in Taylor series and ignore quadratic and higher order terms in this dependence. As a result, we get a linear expression for v_i in terms of v(x):

$$v_i = \int w_i(x) \cdot v(x) \, dx.$$

8 Case of Full Information about the Resolution

Description. In this section, we consider the case when we know the exact expression for this dependence, i.e., when we know the weights $w_i(x)$.

The notion of fuzzy transform. Intuitively, each "averaged" value v_i can be viewed as the value of the field v(x) at a "fuzzy" point characterized by uncertainty $w_i(x)$. Because of this interpretation, the transformation from the original function v(x) to the set of values v_1, \ldots, v_n is also known as a fuzzy transform; see, e.g., [5, 6].

What we want to predict. Based on the measurement results $\tilde{v}_1, \ldots, \tilde{v}_n$, we would like to reconstruct the field v(x). From the pragmatic viewpoint, knowing the field means being able to predict the results of all other measurements of this field.

Each such measurement can be characterized by its own averaging function w(x). Thus, predicting the result of the measurement means predicting the corresponding averaged value $y = \int w(x) \cdot v(x) dx$.

Of course, the space of functions is infinite-dimensional, which means that to uniquely reconstruct a function, we need to know infinitely many parameters. Thus, based on n numbers $\tilde{v}_1, \ldots, \tilde{v}_n$, we cannot uniquely reconstruct the function v(x) – and thus, we cannot uniquely reconstruct the desired averaged value y. So, the problem is to find the $range[y, \overline{y}]$ of this value y.

Prediction problem as a particular case of linear programming (LP). The lower endpoint \underline{y} is the smallest possible value of y, the upper endpoint \overline{y} is the largest possible value of y. Thus, the problems of finding the desired endpoints \underline{y} and \overline{y} can be formulated in the following optimization form: minimize (maximize) $y - \int w(x) \cdot v(x) dx$ under the constraints

$$\underline{v}_i \le \int w_i(x) \cdot v(x) \, dx \le \overline{v}_i, \quad 1 \le i \le n.$$

In both problems, we optimize the value of a linear functional under linear constraints, so from the mathematical viewpoint, these problems are (infinite-dimensional) linear programming (LP) problems.

Without prior restrictions on the field v(x), we cannot predict anything. In general, if we do not impose any conditions on the function v(x), then both bounds are infinite – unless w(x) is a linear combination of $w_i(x)$. Indeed, it is known that every vector w which is orthogonal to all the vectors t, which are orthogonal to all the vectors w_1, \ldots, w_n , belongs to the linear space generated by the vectors w_1, \ldots, w_n – i.e., is a linear combination of w_1, \ldots, w_n . Thus, if a vector w cannot be represented as a linear combination of the vectors w_1, \ldots, w_n , then there exists a vector t which is orthogonal to all w_i but not to w. With respect to the space of all the functions, this means that if w(x) cannot be represented as a linear combination of the functions $w_i(x)$, then there exists a function t(x) which is orthogonal to all $w_i(x)$ (in the sense that $\langle w_i, t \rangle \stackrel{\text{def}}{=} \int w_i(x) \cdot t(x) \, dx = 0$) but not to w(x) ($\langle w, t \rangle \neq 0$).

For an arbitrary real number λ , instead of the actual field v(x), we can now consider a new field $v_{\lambda}(x) \stackrel{\text{def}}{=} v(x) + \lambda \cdot t(x)$. For this new field $v_{\lambda}(x)$, the values of v_i are the same as for the original field v(x) – and hence, satisfy the same inequalities. However, the new value y is equal to $y_{\lambda} = \langle w, v \rangle + \lambda \cdot \langle w, t \rangle$. Since $\langle w, t \rangle \neq 0$, for appropriate λ , we can get this value y_{λ} equal to any given real number. Thus, indeed, the smallest possible value of y is $\underline{y} = -\infty$, and the largest possible value of y is $\overline{y} = +\infty$.

Non-negative fields. In many practical problems, the field v(x) can only have non-negative values $v(x) \geq 0$. For example, in geophysics, the density v(x) cannot be negative. Under this additional restrictions, we already have non-trivial bounds y and \overline{y} .

Dual LP techniques. For solving these problems, we can use the experience of imprecise probabilities [4,8] where in similar LP problems, v(x) is the non-negative probability density function (and the weights are, e.g., functions x^2 corresponding to moments). According to this experience, many efficient algorithms come from considering dual LP problems, i.e., by computing the range $[\underline{v}, \overline{v}]$, where

$$\underline{v} = \sup \left\{ \sum y_i \cdot \underline{v}_i : \sum y_i \cdot w_i(x) \le w(x) \right\};$$

$$\overline{v} = \inf \left\{ \sum y_i \cdot \overline{v}_i : w(x) \le \sum y_i \cdot w_i(x) \right\}.$$

Indeed, if $\sum y_i \cdot w_i(x) \leq w(x)$, then, by multiplying both sides of this inequality by $v(x) \geq 0$ and integrating over x, we conclude that $\sum y_i \cdot v_i \leq y$. Since we know that $v_i \geq \underline{v}_i$, we thus get a lower bound for $y \colon y \geq \sum y_i \cdot \underline{v}_i$. Thus, y is larger than the largest of these bounds, i.e., $y \geq \underline{v}$. So, we can conclude that $\underline{y} \geq \underline{v}$. Similarly, we can conclude that $\overline{y} \leq \overline{v}$, i.e., that the dual LP interval $[\underline{v}, \overline{v}]$ is the enclosure for the desired range $[y, \overline{y}]$.

Comments.

- For discrete LP problems, the dual interval is exactly equal to the original one.
- Our problems are easier than the imprecise probability ones, since the functions $w_i(x)$ are usually localized and thus, for each x, usually at most a few functions $w_i(x)$ differ from 0. This makes checking the sums easier.
- Checking the inequalities like $\sum y_i \cdot w_i(x) \leq w(x)$ is even easier in a practically important case of piece-wise linear functions $w_i(x)$ and w(x). In this case, it is sufficient to check this inequality at endpoints of linearity intervals then, due to linearity, it will be automatically true for all internal points as well.

9 Situations in Which We Only Know Upper Bounds

General idea. In other cases – similarly to the interval setting – we do not only know the upper bounds δ on the location error $\widetilde{x} - x$. A natural question is: when is a model v(x) consistent with the given observation $(\widetilde{v}, \widetilde{x})$?

In this case, the measured value \widetilde{v} is Δ -close to a convex combination of values v(x) for x s.t. $||x - \widetilde{x}|| \leq \Delta x$. Thus, $\underline{v}_{\delta}(\widetilde{x}) - \Delta \leq \widetilde{v} \leq \overline{v}_{\delta}(\widetilde{x}) + \Delta$, where:

$$\underline{v}_{\delta}(\widetilde{x}) \stackrel{\text{def}}{=} \inf\{v(x) : \|x - \widetilde{x}\| \le \delta\}, \text{ and } \overline{v}_{\delta}(\widetilde{x}) \stackrel{\text{def}}{=} \sup\{v(x) : \|x - \widetilde{x}\| \le \delta\}.$$

Case of interval models. In real life, we rarely have an exact model v(x). Usually, we have bounds on v(x), i.e., an interval-valued model $\mathbf{v}(x) = [v^-(x), v^+(x)]$. An observation $(\widetilde{v}, \widetilde{x})$ consistent with this "interval-valued" model if these exists a model $v(x) \in \mathbf{v}(x)$ which is consistent with this

Since the values \underline{v}_{δ} and \overline{v}_{δ} monotonically depend on v(x), this consistency leads to

$$\underline{v}_{\delta}^{-}(\widetilde{x}) - \Delta \leq \widetilde{v} \leq \overline{v}_{\delta}^{+}(\widetilde{x}) + \Delta.$$

Relation to Hausdorff metric. In many practical problems, the field v(x) continuously depends on x. For continuous functions, inf and sup on a bounded closed set $\{x: \|x-\widetilde{x}\| \leq \delta\}$ are attained at some value. Thus, the above criterion for consistency between a model and observations can be simplified.

Namely, in this case, the set \widetilde{m} of all measurement results $(\widetilde{v}, \widetilde{x})$ is consistent with the model v(x) if and only if

$$\forall (\widetilde{v}, \widetilde{x}) \in \widetilde{m} \, \exists (v(x), x) \in v \, ((\widetilde{v}, \widetilde{x}) \text{ is } (\Delta, \delta) \text{-close to } (v(x), x)),$$

i.e.,
$$|\widetilde{v} - v| \leq \Delta$$
 and $||x - \widetilde{x}|| \leq \delta$.

This definition is similar to the standard definition of the Hausdorff metric d_H : $d_H(A,B) \le \varepsilon$ means that

$$\forall a \in A \,\exists b \in B \, (d(a,b) \leq \varepsilon) \text{ and } \forall b \in B \,\exists a \in A \, (d(a,b) \leq \varepsilon).$$

(This similarity was noticed in [1].)

Specifically, the above definition is an *asymmetric* version of Hausdorff metric. Let us show, on a simple example, that our "distance" is indeed asymmetric.

- the actual field has the form v(0) = 1 and v(x) = 0 for $x \neq 0$, and
- the measurements results are all zeros, i.e., $\widetilde{v} = 0$ for all \widetilde{x} .

In this case, all the measurements are consistent with the model:

- the values $\tilde{v}=0$ for $\tilde{x}\neq 0$ are consistent with v=0 for $x=\tilde{x}$, and
- the value $\tilde{v} = 0$ for $\tilde{x} = 0$ is consistent with v(x) = 0 for $x = \delta$ s.t. $|\tilde{x} x| \leq \delta$.

In this example,

- the actual field is all zeros, i.e., v(x) = 0 for all x, and
- the measurement results are $\tilde{v}=1$ for $\tilde{x}=0$, and $\tilde{v}=0$ for all $\tilde{x}\neq 0$.

Here, when $\Delta < 1$, the measurement (1,0) is inconsistent with the model, because for all x which are δ -close to $\widetilde{x} = 0$, we have v(x) = 0 hence we should have $|\widetilde{x} - v(x)| = |\widetilde{x}| \leq \Delta$.

10 Case of Minimal Knowledge About Uncertainty

Idea. Yet another case is when we do not even know δ . It happened, e.g., when we solve the seismic inverse problem to find the velocity distribution.

In this case, a natural heuristic idea is:

- to add a perturbation of size δ_0 (e.g., sinusoidal) to the reconstructed field $\widetilde{v}(x)$,
- to simulate the new measurement results,
- to apply the same algorithm to the simulated results, and
- to reconstruct the new field $\widetilde{v}_{\text{new}}(x)$.

If the perturbations are not visible in $\tilde{v}_{\text{new}}(x) - \tilde{v}(x)$, this means that details of size δ_0 cannot be reconstructed and so, the actual resolution is $\delta > \delta_0$. This approach was partially described in [2, 7].

Linearization and its consequences. Which perturbations should we choose? To select the optimal perturbations, we will take into account the fact that since perturbations are usually small, we can safely linearize their effects. Thus, if we know the results $\Delta v_1(x), \ldots, \Delta v_k(x)$ of applying perturbations $e_1(x), \ldots, e_k(x)$, we can predict the result $\Delta v(x)$ of applying an linear combination

$$e(x) = c_1 \cdot e_1(x) + \ldots + c_k \cdot e_k(x),$$

as

$$\Delta v(x) = c_1 \cdot \Delta v_1(x) + \ldots + c_k \cdot \Delta v_k(x).$$

In other words, once we know the results of applying k different perturbations $e_1(x), \ldots, e_k(x)$, we thus also know the results of applying an arbitrary perturbation from the linear space

$$L = \{c_1 \cdot e_1(x) + \ldots + c_k \cdot e_k(x)\}.$$

From this viewpoint, it does not matter what exactly perturbations $e_i(x)$ we select as long as they are within the same space L.

Thus, the question of optimally selecting a given number k of perturbations can be formulated as the question of optimally selecting a k-dimensional linear subspace L in the space of all functions.

Shift-invariance: a natural requirement. To select the space L space, let us use the fact that in most physical problems, there is no preferred spatial location. Thus, in principle, we can choose different locations as origins (x=0) of the coordinate system.

It is reasonable to require that the optimal family of perturbations do not change if we simply change the origin x = 0. For example, if we select a point with the original coordinates x_0 as the origin of a new coordinate system, then the new coordinates will have the form $x_{\text{new}} = x - x_0$. In the original coordinates, the optimal family of perturbations has the form

$$\{c_1 \cdot e_1(x) + \ldots + c_k \cdot e_k(x)\}.$$

In the new coordinates x_{new} , we should expect the exact same family of perturbations

$$\{c_1 \cdot e_1(x_{\text{new}}) + \ldots + c_k \cdot e_k(x_{\text{new}})\}.$$

In terms of the original coordinates, this new family has the form

$$\{c_1 \cdot e_1(x-x_0) + \ldots + c_k \cdot e_k(x-x_0)\}.$$

This "shifted" family must coincide with the original one. In particular, every basis function $e_i(x-x_0)$ from the shifted basis must belong to the original family, i.e., must have the form

$$e_i(x + x_0) = \sum_{j=1}^k c_{ij}(x_0) \cdot e_j(x)$$

for some coefficients c_{ij} which are, in general, depending on the shift x_0 .

Smoothness: an additional requirement. In many physical problems, it is reasonable to consider smooth perturbations, i.e., perturbations for which the functions $e_i(x)$ are differentiable. In this case, by considering different values x, we get a system of linear equations for determining $c_{ij}(x_0)$ in terms of the smooth functions $e_i(x+x_0)$ and $e_j(x)$. The solution of a system of linear equations is – due to Cramer's rule – a smooth function of the coefficients and of the right-hand sides. Thus, the solutions $c_{ij}(x_0)$ are also smooth.

From the requirements to the description of the desired family L. Let us fix one of the spatial coordinates, e.g., the coordinate x_1 . For shifts w.r.t. this coordinate,

we have
$$e_i(x_1 + x_0, x_2, ...) = \sum_{j=1}^k c_{ij}(x_0) \cdot e_j(x_1, x_2, ...)$$

Since the functions $e_i(x_1+x_0,\ldots)$ and $c_{ij}(x_0)$ are smooth, we can differentiate both sides of the above equation with respect to x_0 and take $x_0=0$. For each components of x_0 , we get a system of linear differential equations $e_i'=\sum c_{ij}'(0)\cdot e_j$ with constant coefficients. A general solution to such a system is well known: it is a linear combination of expressions $x_1^{k_{1j}}\cdot \exp(a_{1j}\cdot x_1)$ with complex values a_{1j} (eigenvalues of the matrix $c_{ij}'(0)$) and integers $k_{1j}\geq 0$ (multiplicities of these eigenvalues).

Some of these solutions tend to infinity exponentially fast. Such solutions are not useful as perturbations, since perturbations must be uniformly small. So, it is reasonable to restrict ourselves to *bounded* perturbations.

This boundedness eliminates the terms with $\text{Re}(a_{1j}) \neq 0$. Thus, the only remaining terms correspond to imaginary values a_{1j} – i.e., to sinusoids. For these terms, boundedness also eliminates terms with $k_{1j} > 0$, so we only get pure sinusoids:

$$e_i(x_1, x_2, \ldots) = \sum_j C_j(x_2, \ldots) \cdot \sin(\omega_{1j} \cdot x_1).$$

The functions $C_j(x_2,...)$ can be computed as linear combinations of the values $e_i(x_1,x_2,...)$ corresponding to different values x_1 . On the other hand, the dependence of e_i on x_2 is also a linear combination of sinusoids. Thus, the functions $C_j(x_2,...)$ are linear combinations of sinusoids in x_2 . Substituting these linear combinations instead of $C_j(x_2,...)$ into the above formula, and taking into account that $\sin(a) \cdot \sin(b)$ is a linear combination of $\cos(a+b)$ and $\cos(a-b)$, we conclude that the dependence of e_i on x_1 and x_2 takes the form

$$e_i(x_1, x_2, x_3, \ldots) = \sum_k C_k(x_3, \ldots) \cdot \sin(\omega_{1k} \cdot x_1 + \omega_{2k} \cdot x_2).$$

Similarly, we can add x_3 , etc., and conclude that each function $e_i(x)$ is a linear combination of the sinusoids $\sin(\sum \omega_j \cdot x_j + \varphi)$.

Conclusion. We conclude that the optimal perturbations are linear combinations of sinusoids. We thus arrive at the following recommendation: use sinusoidal perturbations.

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