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Need to Combine Interval and Probabilistic Uncertainty: What Needs to Be Computed, What Can Be Computed, What Can Be Feasibly Computed, and How Physics Can Help

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**Abstract** In many practical situations, the quantity of interest is difficult to measure directly. In such situations, to estimate this quantity, we measure easier-to-measure quantities which are related to the desired one by a known relation, and we use the results of these measurement to estimate the desired quantity. How accurate is this estimate?

Traditional engineering approach assumes that we know the probability distributions of measurement errors; however, in practice, we often only have partial information about these distributions. In some cases, we only know the upper bounds on the measurement errors; in such cases, the only thing we know about the actual value of each measured quantity is that it is somewhere in the corresponding interval. Interval computation estimates the range of possible values of the desired quantity under such interval uncertainty.

In other situations, in addition to the intervals, we also have partial information about the probabilities. In this paper, we describe how to solve this problem in the linearized case, what is computable and what is feasibly com-

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putable in the general case, and, somewhat surprisingly, how physics ideas – that initial conditions are not abnormal, that every theory is only approximate – can help with the corresponding computations.

**Keywords** Interval uncertainty  $\cdot$  Probabilistic uncertainty  $\cdot$  Feasible algorithms  $\cdot$  Physics helps computing

# 1 Need to Combine Interval and Probabilistic Uncertainty: Linearized Case

Need to take uncertainty into account when processing data. In practice, we are often interested in a quantity y which is difficult to measure directly. Examples are distance to a star, amount of oil in the well, tomorrow's weather.

A solution to this problem is to find easier-to-measure quantities  $x_1, \ldots, x_n$  related to y by a known dependence  $y = f(x_1, \ldots, x_n)$ . Then, we measure  $x_i$  and use measurement results  $\tilde{x}_i$  to compute an estimate  $\tilde{y} = f(\tilde{x}_1, \ldots, \tilde{x}_n)$  for the desired quantity y. Such computations are usually called *data processing*.

Measurements are never absolutely accurate, so even if the model f is exact,  $\tilde{x}_i \neq x_i$  leads to  $\Delta y \stackrel{\text{def}}{=} \tilde{y} - y \neq 0$ . It is important to use information about measurement errors  $\Delta x_i \stackrel{\text{def}}{=} \tilde{x}_i - x_i$  to estimate the accuracy  $\Delta y$ ; see, e.g., [23].

We often have imprecise probabilities. The usual assumption is that we know the probabilities of different values of measurement errors  $\Delta x_i$ . How can we find these probabilities?

To find them, we measure the same quantities:

- with our measuring instrument (MI) and
- with a much more accurate MI, with  $\widetilde{x}_i^{\text{st}} \approx x_i$ .

However, in two important cases, this does not work: in the case of state-of-the art-measurements, and in the case of measurements on the shop floor. In the first case, when we use state-of-the-art measuring instruments, so more accurate instruments are available. In the second case, it is, in principle, possible to accurately calibrate each sensor, but that would cost too much.

In both cases, we have partial information about probabilities. Often, all we know is an upper bound  $|\Delta x_i| \leq \Delta_i$ . Then, the only thing that we know about the actual (unknown) values  $x_i$  of the measured quantities is that  $x_i \in [\widetilde{x}_i - \Delta_i, \widetilde{x}_i + \Delta_i]$ . Then, the only thing that we know about  $y = f(x_1, \dots, x_n)$  is that

$$y \in [\underline{y}, \overline{y}] \stackrel{\text{def}}{=} \{ f(x_1, \dots, x_n) : x_i \in [\widetilde{x}_i - \Delta_i, \widetilde{x}_i + \Delta_i] \}.$$

Computing this interval  $[\underline{y}, \overline{y}]$  is known as *interval computation*; see, e.g., [4, 17, 19].

**Data processing: example.** Let us provide an example of data processing. Suppose that we want to measure coordinates  $X_j$  of an object. To find these

coordinates, we measure the distance  $Y_i$  between this object and objects with accurately known coordinates  $X_i^{(i)}$ :

$$Y_i = \sqrt{\sum_{j=1}^{3} (X_j - X_j^{(i)})^2}.$$

After the measurements, we know the results  $\widetilde{Y}_i$  of measuring  $Y_i$ . We want to estimate the desired quantities  $X_j$ .

Usually linearization is possible. In most practical situations, we know the approximate values  $X_j^{(0)}$  of the desired quantities  $X_j$ . These approximation are usually reasonably good, in the sense that the difference  $x_j \stackrel{\text{def}}{=} X_j - X_j^{(0)}$  are small.

In terms of  $x_j$ , we have  $Y_i = f(X_1^{(0)} + x_1, \dots, X_n^{(0)} + x_n)$ . When the differences  $x_i$  are small, we can safely ignore terms quadratic in  $x_j$ . Indeed, even if the estimation accuracy is 10% (0.1), its square is 1%  $\ll$  10%. We can thus expand the dependence of  $Y_i$  on  $x_j$  in Taylor series and keep only linear terms:

$$Y_i = Y_i^{(0)} + \sum_{j=1}^n a_{ij} \cdot x_j, \quad Y_i^{(0)} \stackrel{\text{def}}{=} f_i(X_1^{(0)}, \dots, X_n^{(0)}), \quad a_{ij} \stackrel{\text{def}}{=} \frac{\partial f_i}{\partial X_j}.$$

**Least squares.** Thus, to find the unknowns  $x_j$ , we need to solve a system of approximate linear equations  $\sum_{j=1}^{n} a_{ij} \cdot x_i \approx y_i$ , where  $y_i \stackrel{\text{def}}{=} \widetilde{Y}_i - Y_i^{(0)}$ . Usually, it is assumed that each measurement error is normally distributed with 0 mean and known standard deviation  $\sigma_i$ .

The distribution is indeed often normal: the measurement error is a joint result of many independent factors, and the distribution of the sum of many small independent errors is close to Gaussian; this result is known as the Central Limit Theorem; see, e.g., [24].

0 mean also makes sense: we calibrate the measuring instrument by comparing it with a more accurate one, so if there was a bias (non-zero mean), we delete it by re-calibrating the scale.

It is also usually assumed that measurement errors of different measurements are independent. In this case, for each possible combination  $x = (x_1, \ldots, x_n)$ , the probability of observing  $y_1, \ldots, y_m$  is equal to the product of the corresponding probabilities:

$$\prod_{i=1}^{m} \left( \frac{1}{\sqrt{2\pi} \cdot \sigma_i} \cdot \exp\left(-\frac{\left(y_i - \sum_{j=1}^{n} a_{ij} \cdot x_j\right)^2}{2\sigma_i^2}\right) \right).$$

It is reasonable to select  $x_j$  for which this probability is the largest, i.e., equivalently, for which

$$\sum_{i=1}^{n} \frac{\left(y_i - \sum_{j=1}^{n} a_{ij} \cdot x_j\right)^2}{\sigma_i^2} \to \min.$$

For every confidence level  $\gamma$ , the confidence set  $S_{\gamma}$ , i.e., the set off all combinations x which are possible with this degree of confidence, can be determined by the formula

$$S_{\gamma} = \left\{ x : \sum_{i=1}^{n} \frac{\left( y_i - \sum_{j=1}^{n} a_{ij} \cdot x_j \right)^2}{\sigma_i^2} \le \chi_{m-n,\gamma}^2 \right\}.$$

Sometimes this set is empty; this means that some measurements are outliers.

Need to take into account systematic error. In the traditional approach, we assume that  $y_i = \sum_{j=1}^n a_{ij} \cdot x_j + e_i$ , where the measurement error  $e_i$  has 0 mean. However, sometimes, in addition to the random error  $e_i^r \stackrel{\text{def}}{=} e_i - E[e_i]$  with 0 mean, we also have a systematic error  $e_i^s \stackrel{\text{def}}{=} E[e_i]$ :

$$y_i = \sum_{j=1}^{n} a_{ij} \cdot x_j + e_i^r + e_i^s.$$

Sometimes, we know the upper bound  $\Delta_i$ :  $|e_i^s| \leq \Delta_i$ .

What can we then say about  $x_j$ ?

Comment. In other cases, we have different bounds  $\Delta_i(p)$  corresponding to different degree of confidence p; this is known as the *fuzzy case*; see, e.g., see, e.g., [1,18,5,20,21,29].

Combining probabilistic and interval uncertainty: main idea. If we knew the values  $e_i^s$ , then we would conclude that for  $e_i^r = y_i - \sum_{j=1}^n a_{ij} \cdot x_j - e_i^s$ , we have

$$\sum_{i=1}^{m} \frac{(e_i^r)^2}{\sigma_i^2} = \sum_{i=1}^{m} \frac{\left(y_i - \sum_{j=1}^{n} a_{ij} \cdot x_j - e_i^s\right)^2}{\sigma_i^2} \le \chi_{m-n,\gamma}^2.$$

In practice, we do not know the values  $e_i^s$ , we only know that these values are in the interval  $[-\Delta_i, \Delta_i]$ . Thus, we know that the above inequality holds for some  $e_i^s \in [-\Delta_i, \Delta_i]$ .

The above condition is equivalent to  $v(x) \leq \chi^2_{m-n,\gamma}$ , where

$$v(x) \stackrel{\text{def}}{=} \min_{e_i^s \in [-\Delta_i, \Delta_i]} \sum_{i=1}^m \frac{\left(y_i - \sum_{j=1}^n a_{ij} \cdot x_j - e_i^s\right)^2}{\sigma_i^2}.$$

So, the set  $S_{\gamma}$  of all combinations  $X=(x_1,\ldots,x_n)$  which are possible with confidence  $1-\gamma$  is:  $S_{\gamma}=\{x:v(x)\leq\chi^2_{m-n,\gamma}\}$ . The range of possible values of  $x_j$  can be obtained by maximizing and minimizing  $x_j$  under the constraint  $v(x)\leq\chi^2_{m-n,\gamma}$ .

Comment. In the fuzzy case, we have to repeat the computations for every p.

How to check consistency. We want to make sure that the measurements are consistent – i.e., that there are no outliers. This means that we want to check that there exists some  $x = (x_1, \ldots, x_n)$  for which  $v(x) \leq \chi^2_{m-n,\gamma}$ . This condition is equivalent to ([25]):

$$v \stackrel{\text{def}}{=} \min_{x} v(x) = \min_{x} \min_{e_{i}^{s} \in [-\Delta_{i}, \Delta_{i}]} \sum_{i=1}^{m} \frac{\left(y_{i} - \sum_{j=1}^{n} a_{ij} \cdot x_{j} - e_{i}^{s}\right)^{2}}{\sigma_{i}^{2}} \leq \chi_{m-n, \gamma}^{2}.$$

This is indeed a generalization of probabilistic and interval approaches. In the case when  $\Delta_i = 0$  for all i, i.e., when there is no interval uncertainty, we get the usual Least Squares.

Vice versa, for very small  $\sigma_i$ , we get the case of pure interval uncertainty. In this case, the above formulas tend to the set of all the values for which  $\left|y_i - \sum\limits_{j=1}^n a_{ij} \cdot x_j\right| \leq \Delta_i$ . For example, for m repeated measurements of the same quantity, we get the intersection of the corresponding intervals.

So, the new idea is indeed a generalization of the known probabilistic and interval approaches.

From formulas to computations. The expression 
$$\left(y_i - \sum_{j=1}^n a_{ij} \cdot x_j - e_i^s\right)^2$$

is a convex function of  $x_j$ . The domain of possible values of  $e^s = (e_1^s, \ldots, e_m^s)$  is also convex: it is a box  $[-\Delta_1, \Delta_1] \times \ldots \times [-\Delta_m, \Delta_m]$ . There exist efficient algorithms for computing minima of convex functions over convex domains; these algorithms also compute locations where these minima are attained; see, e.g., [13] and references therein. Thus, for every x, we can efficiently compute v(x) and thus, efficiently check whether  $v(x) \leq \chi^2_{m-n,\gamma}$ .

Similarly, we can efficiently compute v and thus, check whether  $v \leq \chi^2_{m-n,\gamma}$  – i.e., whether we have outliers.

The set  $S_{\gamma}$  is convex. We can approximate the set  $S_{\gamma}$  by

- taking a grid G,
- checking, for each  $x \in G$ , whether  $v(x) \le \chi^2_{m-n,\gamma}$ , and
- taking the convex hull of "possible" points.

We can also efficiently find the minimum  $\underline{x}_j$  of  $x_j$  over  $x \in S_\gamma$ . By computing the min of  $-x_j$ , we can also find the maximum  $\overline{x}_j$ .

## 2 General Case: What Can Be Computed?

How do we describe imprecise probabilities? The ultimate goal of most estimates is to make decisions. It is known that a rational decision-maker maximizes expected *utility* E[u(y)].

- For smooth u(y),  $y \approx \widetilde{y}$  implies that

$$u(y) = u(\widetilde{y}) + (y - \widetilde{y}) \cdot u'(\widetilde{y}) + \frac{1}{2} \cdot (y - \widetilde{y})^2 \cdot u''(\widetilde{y}).$$

So, to find E[u(y)], we must know moments  $E[(y-\tilde{y})^k]$ .

- Often, u(y) abruptly changes: e.g., when pollution level exceeds  $y_0$ , the plant has to pay a huge fine; then E[u(y)] is proportional to the cdf:

$$E[u(y)] \sim F(y) \stackrel{\text{def}}{=} \text{Prob}(y \leq y_0).$$

So, it is enough to know moments and cdf. From the cdf F(y), we can estimate moments, so F(y) is enough.

Imprecise probabilities mean that we don't know F(y) exactly, we only know bounds  $(p\text{-}box) \underline{F}(y) \leq F(y) \leq \overline{F}(y)$ .

What is computable? Computations with p-boxes are practically important. It is thus desirable to come up with efficient algorithms which are as general as possible.

It is known that too general problems are often *not* computable. To avoid wasting time, it is therefore important to find out what *can* be computed.

At first glance, this question sounds straightforward:

- to describe a cdf, we can consider a computable function F(x);
- to describe a p-box, we consider a computable function interval  $[\underline{F}(x), \overline{F}(x)]$ .

Often, we can do that, but we will show that sometimes, we need to go beyond computable function intervals. To explain all this, let us recall what computable means in general; see, e.g., [13,22,28].

**Reminder: what is computable?** A real number x corresponds to a value of a physical quantity. We can measure x with higher and higher accuracy. So, we arrive at the following definition:

**Definition 2.1.** A real number x is called computable if there is an algorithm, that, given k, produces a rational  $r_k$  s.t.  $|x - r_k| \le 2^{-k}$ .

A computable function computes f(x) from x. We can only use approximations to x. So, an algorithm for computing a function can, given k, request a  $2^{-k}$ -approximation to x. Most usual functions are thus computable.

Not all functions are computable, an exception is a step-function f(x) = 0 for x < 0 and f(x) = 1 for  $x \ge 0$ . Indeed, no matter how accurately we know  $x \approx 0$ , from  $r_k = 0$ , we cannot tell whether x < 0 or  $x \ge 0$  [13,22,28].

Consequences for representing a cdf F(x). We would like to represent a general probability distribution by its cdf F(x). From the purely mathematical viewpoint, this is indeed the most general representation.

At first glance, it makes sense to consider computable functions F(x). For many distributions, e.g., for Gaussian, F(x) is computable.

However, when x=0 with probability 1, the cdf F(x) is exactly the step-function. And we already know that the step-function is not computable. Thus, we need to find an alternative way to represent cdf's – beyond computable functions.

Back to the drawing board. Each value F(x) is the probability that  $X \leq x$ . We cannot empirically find exact probabilities p. We can only estimate frequencies f based on a sample of size N.

For large N, the difference  $d \stackrel{\text{def}}{=} p - f$  is asymptotically normal, with  $\mu = 0$  and  $\sigma = \sqrt{\frac{p \cdot (1-p)}{N}}$ . Situations when  $|d-\mu| < 6\sigma$  are negligibly rare, so we conclude that  $|f-p| \le 6\sigma$ .

For large N, we can get  $6\sigma \leq \delta$  for any accuracy  $\delta > 0$ . We get a sample  $X_1, \ldots, X_N$ . We don't know the exact values  $X_i$ , only measured values  $\widetilde{X}_i$  such that  $|\widetilde{X}_i - X_i| \leq \varepsilon$  for some accuracy  $\varepsilon$ .

So, what we have is a frequency  $f = \text{Freq}(\widetilde{X}_i \leq x)$ .

Resulting definition. Here,  $X_i \leq x - \varepsilon$  implies that  $\widetilde{X}_i \leq x \Rightarrow X_i \leq x + \varepsilon$ , so

$$\operatorname{Freq}(X_i \le x - \varepsilon) \le f = \operatorname{Freq}(\widetilde{X}_i \le x) \le \operatorname{Freq}(X_i \le x + \varepsilon).$$

Frequencies are  $\delta$ -close to probabilities, so we arrive at the following definition [15]:

**Definition 2.2.** A cdf F(x) is called computable if there is an algorithm that, given  $x, \varepsilon > 0$ , and  $\delta > 0$ , computes a rational number f such that

$$F(x - \varepsilon) - \delta \le f \le F(x + \varepsilon) + \delta.$$

In the computer, to describe a distribution on an interval  $[\underline{T}, \overline{T}]$ : we select a grid  $x_1 = \underline{T}$ ,  $x_2 = \underline{T} + \varepsilon$ , ..., and we store the corresponding frequencies  $f_i$  with accuracy  $\delta$ . A class of possible distribution is represented, for each  $\varepsilon$  and  $\delta$ , by a finite list of such approximations.

**First equivalent definition.** It turns out that our definition is equivalent to the following one:

**Definition 2.3.** A cdf F(x) is called computable if there exists an algorithm that, given x, varepsilon > 0, and  $\delta > 0$ , computes a rational number f which is  $\delta$ -close to F(x') for some x' such that  $|x' - x| \leq \varepsilon$ .

Indeed, here is a proof of equivalence. We know that  $F(x+\varepsilon)-F(x+\varepsilon/3)\to 0$  as  $\varepsilon\to 0$ . So, for  $\varepsilon=2^{-k},\ k=1,2,\ldots$ , we take f and f' such that

$$F(x + \varepsilon/3) - \delta/4 \le f \le F(x + (2/3) \cdot \varepsilon) + \delta/4$$

$$F(x + (2/3) \cdot \varepsilon) - \delta/4 \le f' \le F(x + \varepsilon) + \delta/4.$$

We stop when f and f' are sufficiently close, i.e., when  $|f - f'| \le \delta$ . Thus, we get the desired f.

**Second equivalent definition.** We start with pairs  $(x_1, f_1)$ ,  $(x_2, f_2)$ , ... When  $f_{i+1} - f_i > \delta$ , we add intermediate pairs

$$(x_i, f_i + \delta), (x_i, f_i + 2\delta), \dots, (x_i, f_{i+1}).$$

The resulting set of pairs is  $(\varepsilon, \delta)$ -close to the graph

$$\{(x,y): F(x-0) \le y \le F(x)\}$$

in Hausdorff metric  $d_H$ . This metric can be defined as follows.

**Definition 2.4.** (x,y) and (x',y') are  $(\varepsilon,\delta)$ -close if  $|x-x'| \le \varepsilon$  and  $|y-y'| \le \delta$ .

**Definition 2.5.** The sets S and S' are  $(\varepsilon, \delta)$ -close if for every  $s \in S$ , there is a  $(\varepsilon, \delta)$ -close point  $s' \in S'$ ; for every  $s' \in S'$ , there is a  $(\varepsilon, \delta)$ -close point  $s \in S$ .

Compact sets with metric  $d_H$  form a computable compact. So, F(x) is a monotonic computable object in this compact.

What can be computed: a positive result for the 1D case. We are interested in computing the expected value  $E_{F(x)}[u(x)]$  for smooth u(x). Our result is as follows:

**Proposition 2.1.** There is an algorithm that given a computable cdf F(x), a computable function u(x), and accuracy  $\delta > 0$ , computes  $E_{F(x)}[u(x)]$  with accuracy  $\delta$ .

Comment. For computable classes  $\mathcal{F}$  of cdfs, a similar algorithm computes the range of possible values  $[\underline{u}, \overline{u}] \stackrel{\text{def}}{=} \{E_{F(x)}[u(x)] : F(x) \in \mathcal{F}\}.$ 

**Proof: main idea.** Computable functions are computably continuous: for every  $\delta > 0$ , we can compute  $\varepsilon > 0$  such that  $|x-x'| \le \varepsilon$  implies  $|f(x)-f(x')| \le \delta$ . We select  $\varepsilon$  corresponding to  $\delta/4$ , and take a grid with step  $\varepsilon/4$ .

For each  $x_i$ , the value  $f_i$  is  $(\delta/4)$ -close to  $F(x_i')$  for some  $x_i'$  which is  $(\varepsilon/4)$ -close to  $x_i$ .

The function u(x) is  $(\delta/2)$ -close to a piece-wise constant function  $u'(x) = u(x_i)$  for  $x \in [x'_i, x'_{i+1}]$ .

Thus,  $|E[u(x)] - E[u'(x)]| \le \delta/2$ . Here,

$$E[u'(x)] = \sum_{i} u(x_i) \cdot (F(x'_{i+1}) - F(x'_i)).$$

Here,  $F(x_i')$  is close to  $f_i$  and  $F(x_{i+1}')$  is close to  $f_{i+1}$ . Thus, E[u'(x)] (and hence, E[u(x)]) is computably close to a computable sum

$$\sum_{i} u(x_i) \cdot (f_{i+1} - f_i).$$

What to do in a multi-D case? For each g(x), y,  $\varepsilon > 0$ , and  $\delta > 0$ , we can find a frequency f such that:  $|P(g(x) \le y') - f| \le \varepsilon$  for some y' such that  $|y - y'| \le \delta$ . We select an  $\varepsilon$ -net  $x_1, \ldots, x_n$  for X. Then,  $X = \bigcup_i B_{\varepsilon}(x_i)$ , where

 $B_{\varepsilon}(x) \stackrel{\text{def}}{=} \{x' : d(x, x') \leq \varepsilon\}$ . We select  $f_1$  which is close to  $P(B_{\varepsilon'}(x_1))$  for all  $\varepsilon'$  from some interval  $[\underline{\varepsilon}, \overline{\varepsilon}]$  which is close to  $\varepsilon$ .

We then select  $f_2$  which is close to  $P(B_{\varepsilon'}(x_1) \cup B_{\varepsilon'}(x_2))$  for all  $\varepsilon'$  from some subinterval of  $[\varepsilon, \overline{\varepsilon}]$ , etc.

Then, we get approximations to probabilities of the sets

$$B_{\varepsilon}(x_i) - (B_{\varepsilon}(x_1) \cup \ldots \cup B_{\varepsilon}(x_{i-1})).$$

This lets us compute the desired values E[u(x)].

# 3 Taking Into Account That We Process Physical Data

Computations with real numbers: reminder. From the physical view-point, real numbers x describe values of different quantities. We get values of real numbers by measurements. Measurements are never 100% accurate, so after a measurement, we get an approximate value  $r_k$  of x. In principle, we can measure x with higher and higher accuracy.

So, from the computational viewpoint, a real number is a sequence of rational numbers  $r_k$  for which, e.g.,  $|x - r_k| \le 2^{-k}$ .

By an algorithm processing real numbers, we mean an algorithm using  $r_k$  as an "oracle" (subroutine). This is how computations with real numbers are defined in *computable analysis* [13,22,28].

**Known negative results.** The first known negative result that we will use is that no algorithm is possible that, given two numbers x and y, would check whether x = y.

Similarly, we can define a computable function f(x) from real numbers to real numbers as a mapping that,

- given an integer n, a rational number  $x_m$  and its accuracy  $2^{-m}$ ,
- produces  $y_n$  which is  $2^{-n}$ -close to all values f(x) with  $d(x, x_m) \leq 2^{-m}$  (or produces nothing)

- so that for every x and for each desired accuracy n, there is an m for which a  $y_n$  is produced.

We can similarly define a computable function f(x) on a computable compact set K.

The second negative result that we will use is that no algorithm is possible that, given f, returns x such that  $f(x) = \max_{y \in K} f(y)$ . (The maximum itself is computable.)

From the physicists' viewpoint, these negative results seem rather theoretical. In mathematics, if two numbers coincide up to 13 digits, they may still turn to be different. For example, they may be 1 and  $1 + 10^{-100}$ . In physics, if two quantities coincide up to a very high accuracy, it is a good indication that they are equal: if an experimentally value is very close to the theoretical prediction, this means that this theory is (triumphantly) true.

This is how General Relativity was confirmed. This is how physicists realized that light is formed of electromagnetic waves: their speeds are very close; see, e.g., [2,26].

**How physicists argue.** In math, if two numbers coincide up to 13 digits, they may still turn to be different: e.g., 1 and  $1+10^{-100}$ . In physics, if two quantities coincide up to a very high accuracy, it is a good indication that they are equal. A typical physicist argument is that: while numbers like  $1+10^{-100}$  (or  $c \cdot (1+10^{-100})$ ) are, in principle, possible, they are *abnormal* (not *typical*).

In physics, second order terms like  $a \cdot \Delta x^2$  of the Taylor series can be ignored if  $\Delta x$  is small, since:

- while abnormally high values of a (e.g.,  $a=10^{40}$ ) are mathematically possible,
- typical (= not abnormal) values appearing in physical equations are usually of reasonable size.

How to formalize the physicist's intuition of physically meaningful values: main idea. To some physicists, all the values of a coefficient a above 10 are abnormal. To another one, who is more cautious, all the values above 10,000 are abnormal. For every physicist, there is a value n such that all value above n are abnormal.

This argument can be generalized as a following property of the set  $\mathcal{T}$  of all physically meaningful elements. Suppose that we have a monotonically decreasing sequence of sets  $A_1 \supseteq A_2 \supseteq \ldots$  for which  $\bigcap_n A_n = \emptyset$ . In the above example,  $A_n$  is the set of all numbers  $\geq n$ . Then, there exists an integer N for which  $\mathcal{T} \cap A_N = \emptyset$ ; see, e.g., [8,3,7,9–12,14].

How to formalize the physicist's intuition: resulting definition.

**Definition 3.1.** We say that  $\mathcal{T}$  is a set of physically meaningful elements if:

- for every definable decreasing sequence  $\{A_n\}$  for which  $\bigcap_n A_n = \emptyset$ ,

- there exists an N for which  $\mathcal{T} \cap A_N = \emptyset$ .

Comment. Of course, to make this definition precise, we must restrict definability to a *subset* of properties, so that the resulting notion of definability will be defined in formal set theory (ZFC) itself.

Checking equality of real numbers. It is known equality of real numbers is undecidable. For physically meaningful real numbers, however, a deciding algorithm *is* possible.

**Proposition 3.1.** For every set  $\mathcal{T} \subseteq \mathbb{R}^2$  which consists of physically meaningful pairs (x,y) of real numbers, there exists an algorithm deciding whether x=y.

**Proof:** We can take  $A_n = \{(x,y) : 0 < |x-y| < 2^{-n}\}$ . The intersection of all these sets is empty. Hence,  $\mathcal{T}$  has no elements from  $\bigcap_{n=1}^{N_A} A_n = A_{N_A}$ . Thus, for each  $(x,y) \in \mathcal{T}$ , x=y or  $|x-y| \geq 2^{-N_A}$ .

Indeed, we can decide which of the two alternatives is true by comparing  $2^{-(N_A+3)}$ -approximations x' and y' to x and y. Q.E.D.

**Finding roots.** In general, it is not possible, given a function f(x) attaining negative and positive values, to compute its root. This becomes possible if we restrict ourselves to physically meaningful functions.

**Proposition 3.2.** Let K be a computable compact. Let X be the set of all functions  $f: K \to \mathbb{R}$  that attain 0 value somewhere on K. Then

- for every set  $\mathcal{T} \subseteq X$  consisting of physically meaningful functions and for every  $\varepsilon > 0$ ,
- there is an algorithm that, given a f-n  $f \in \mathcal{T}$ , computes an  $\varepsilon$ -approximation to the set of roots  $R \stackrel{\text{def}}{=} \{x : f(x) = 0\}$ .

In particular, we can compute an  $\varepsilon$ -approximation to one of the roots.

**Optimization.** In general, it is not algorithmically possible to find x where f(x) attains maximum. For physically reasonable cases, it is possible:

**Proposition 3.3.** Let K be a computable compact. Let X be the set of all functions  $f: K \to \mathbb{R}$ . Then, for every set  $T \subseteq X$  consisting of physically meaningful functions and for every  $\varepsilon > 0$ , there is an algorithm that,

- given a function  $f \in \mathcal{T}$ ,
- computes an  $\varepsilon$ -approximation to  $S = \left\{ x : f(x) = \max_{y} f(y) \right\}$ .

In particular, we can compute an approximation to an individual  $x \in S$ .

**Proof:** by reduction to the roots problem, since  $f(x) = \max_{y} f(y)$  if and only if g(x) = 0, where  $g(x) \stackrel{\text{def}}{=} f(x) - \max_{y} f(y)$ .

Computing fixed points. In general, it is not possible to compute all the fixed points of a given computable function f(x). Let K be a computable compact. Let X be the set of all functions  $f: K \to K$ . Then:

**Proposition 3.4.** For every set  $\mathcal{T} \subseteq X$  consisting of physically meaningful functions and for every  $\varepsilon > 0$ , there is an algorithm that,

- given a function  $f \in \mathcal{T}$ ,
- computes an  $\varepsilon$ -approximation to the set  $\{x: f(x) = x\}$ .

In particular, we can compute an approximation to an individual fixed point.

**Proof:** reduction to roots, since f(x) = x if and only if g(x) = 0, where  $g(x) \stackrel{\text{def}}{=} d(f(x), x)$ .

Computing limits. In general, it is not algorithmically possible to find a limit  $\lim a_n$  of a convergent computable sequence.

Let K be a computable compact. Let X be the set of all convergent sequences  $a = \{a_n\}, a_n \in K$ . Then:

**Proposition 3.5.** For every set  $\mathcal{T} \subseteq X$  consisting of physically meaningful functions and for every  $\varepsilon > 0$ , there exists an algorithm that,

- given a sequence  $a \in \mathcal{T}$ ,
- computes its limit with accuracy  $\varepsilon$ .

Comment. This result enables us to compute limits of iterations and sums of Taylor series (frequent in physics).

**Proof (main idea):** for every  $\varepsilon > 0$  there exists  $\delta > 0$  such that when  $|a_n - a_{n-1}| \le \delta$ , then  $|a_n - \lim a_n| \le \varepsilon$ .

Intuitively: we stop when two consequent iterations are close to each other.

# 4 How to Take into Account that We Can Use Non-Standard Physical Phenomena to Process Data

Solving NP-complete problems is important. In practice, we often need to find a solution that satisfies a given set of constraints. At a minimum, we need to check whether such a solution is possible. Once we have a candidate, we can feasibly check whether this candidate satisfies all the constraints.

In theoretical computer science, "feasibly" is usually interpreted as computable in polynomial time.

The class of all such problems is called NP; see, e.g., [13]. A typical example of such a problem is satisfiability – checking whether a propositional formula like  $(v_1 \lor \neg v_2 \lor v_3) \& (v_4 \lor \neg v_2 \lor \neg v_5) \& \dots$  can be true.

Each problem from the class NP can be algorithmically solved by trying all possible candidates. For example, we can try all  $2^n$  possible combinations of true-or-false values  $v_1, \ldots, v_n$ .

For medium-size inputs, e.g., for  $n \approx 300$ , the resulting time  $2^n$  is larger than the lifetime of the Universe. So, these exhaustive search algorithms are not practically feasible.

It is not known whether problems from the class NP can be solved feasibly (i.e., in polynomial time). This is the famous open problem  $P \stackrel{?}{=} NP$ .

What we do know is that some problems are *NP-complete*: every problem from NP can be reduced to it. So, it is very important to be able to efficiently solve even one NP-hard problem.

Can non-standard physics speed up the solution of np-complete problems? NP-complete means difficult to solve on computers based on the usual physical techniques. A natural question is: can the use of non-standard physics speed up the solution of these problems?

This question has been analyzed for several specific physical theories, e.g.: for quantum field theory, for cosmological solutions with wormholes and/or casual anomalies.

No physical theory is perfect. If a speed-up is possible within a given theory, is this a satisfactory answer? In the history of physics, always new observations appear which are not fully consistent with the original theory. For example, Newton's physics was replaced by quantum and relativistic theories.

Many physicists believe that every physical theory is approximate. For each theory T, inevitably new observations will surface which require a modification of T. Let us analyze how this idea affects computations.

No physical theory is perfect: how to formalize this idea. We want to formalize a statement that for every theory, eventually there will be observations which violate this theory.

To formalize this statement, we need to formalize what are *observations* and what is a *theory*.

Most sensors already produce *observations* in the computer-readable form, as a sequence of 0s and 1s. Let  $\omega_i$  be the bit result of an experiment whose description is i. Thus, all past and future observations form a (potentially) infinite sequence  $\omega = \omega_1 \omega_2 \dots$  of 0s and 1s.

A physical theory may be very complex. All we care about is which sequences of observations  $\omega$  are consistent with this theory and which are not.

What is a physical theory? So, a physical theory T can be defined as the set of all sequences  $\omega$  which are consistent with this theory.

A physical theory must have at least one possible sequence of observations:  $T \neq \emptyset$ .

A theory must be described by a finite sequence of symbols: the set T must be definable.

How can we check that an infinite sequence  $\omega = \omega_1 \omega_2 \dots$  is consistent with the theory? The only way is check that for every n, the sequence  $\omega_1 \dots \omega_n$  is consistent with T; so:

if 
$$\forall n \,\exists \omega^{(n)} \in T \,(\omega_1^{(n)} \ldots \omega_n^{(n)} = \omega_1 \ldots \omega_n)$$
 then  $\omega \in T$ .

In mathematical terms, this means that T is closed in the Baire metric

$$d(\omega, \omega') \stackrel{\text{def}}{=} 2^{-N(\omega, \omega')},$$

where

$$N(\omega, \omega') \stackrel{\text{def}}{=} \max\{k : \omega_1 \dots \omega_k = \omega'_1 \dots \omega'_k\}.$$

A theory must predict something new. So, for every sequence  $\omega_1 \dots \omega_n$  consistent with T, there is a continuation which does not belong to T.

In mathematical terms, T is  $nowhere\ dense.$  So, we arrive at the following definition.

## What is a physical theory: definition.

**Definition 4.1.** By a physical theory, we mean a non-empty closed nowhere dense definable set T.

**Definition 4.2.** A sequence  $\omega$  is consistent with the no-perfect-theory principle if it does not belong to any physical theory.

In precise terms,  $\omega$  does not belong to the union of all definable closed nowhere dense set. There are countably many definable set, so this union is meager (=  $Baire\ first\ category$ ). Thus, due to Baire Theorem, such sequences  $\omega$  exist.

How to represent instances of an NP-complete problem. For each NP-complete problem  $\mathcal{P}$ , its instances are sequences of symbols. In the computer, each such sequence is represented as a sequence of 0s and 1s. We can append 1 in front and interpret this sequence as a binary code of a natural number i.

In principle, not all natural numbers i correspond to instances of a problem  $\mathcal{P}$ . We will denote the set of all natural numbers which correspond to such instances by  $S_{\mathcal{P}}$ . For each  $i \in S_{\mathcal{P}}$ , we denote the correct answer (true or false) to the i-th instance of the problem  $\mathcal{P}$  by  $s_{\mathcal{P},i}$ .

What we mean by using physical observations in computations. In addition to performing computations, our computational device can produce a scheme i for an experiment, and then use the result  $\omega_i$  of this experiment in future computations.

In other words, given an integer i, we can produce  $\omega_i$ .

In precise terms, the use of physical observations in computations means corresponds to using  $\omega$  as an *oracle*.

## Main result of this section.

**Definition 4.3.** A ph-algorithm  $\mathcal{A}$  is an algorithm that uses an oracle  $\omega$  consistent with the no-perfect-theory principle.

The result of applying an algorithm  $\mathcal{A}$  using  $\omega$  to an input i will be denoted by  $\mathcal{A}(\omega, i)$ .

**Definition 4.4.** We say that a feasible ph-algorithm  $\mathcal{A}$  solves almost all instances of an NP-complete problem  $\mathcal{P}$  if:

$$\forall \varepsilon_{>0} \, \forall n \, \exists N_{\geq n} \, \left( \frac{\#\{i \leq N : i \in S_{\mathcal{P}} \, \& \, \mathcal{A}(\omega, i) = s_{\mathcal{P}, i}\}}{\#\{i \leq N : i \in S_{\mathcal{P}}\}} > 1 - \varepsilon \right).$$

Restriction to sufficiently long inputs  $N \ge n$  makes sense: for short inputs, we can do exhaustive search.

**Proposition 4.1.** For every NP-complete problem  $\mathcal{P}$ , there is a feasible phalgorithm  $\mathcal{A}$  solving almost all instances of  $\mathcal{P}$ .

This result is the best possible. Our result is the best possible, in the sense that the use of physical observations cannot solve *all* instances:

**Proposition 4.2.** If  $P \neq NP$ , then no feasible ph-algorithm A can solve all instances of P.

Can we prove the result for all N starting with some  $N_0$ ? We say that a feasible ph-algorithm  $\mathcal{A}$   $\delta$ -solves  $\mathcal{P}$  if

$$\exists N_0 \, \forall N \geq N_0 \, \left( \frac{\#\{i \leq N : i \in S_{\mathcal{P}} \, \& \, \mathcal{A}(\omega, i) = s_{\mathcal{P}, i}\}}{\#\{i \leq N : i \in S_{\mathcal{P}}\}} > \delta \right).$$

**Proposition 4.3.** For every NP-complete problem  $\mathcal{P}$  and for every  $\delta > 0$ ,

- if there exists a feasible ph-algorithm A that  $\delta$ -solves  $\mathcal{P}$ ,
- then there is a feasible algorithm  $\mathcal{A}'$  that also  $\delta$ -solves  $\mathcal{P}$ .

## 5 Physical and Computational Consequences

**Justification of physical induction.** What is physical induction? It means that if a property P is satisfied in the first N experiments, then it is satisfied always.

Comment: N should be sufficiently large.

**Proposition 5.1.** For every set  $\mathcal{T}$  of physically meaningful sequences  $s = s_1 s_2 \ldots$ , and for every definable property P, there exists a natural number N such that if  $P(s_i)$  holds for all  $i \leq N$ , then  $P(s_i)$  holds for all i.

Proof: Let us take

$$A_n \stackrel{\text{def}}{=} \{s : P(s_1) \& \dots \& P(s_n), \& \exists m \, \neg P(s_m)\}.$$

Then  $A_n \supseteq A_{n+1}$  and  $\cup A_n = \emptyset$  so  $\exists N (A_N \cap \mathcal{T} = \emptyset)$ .

The meaning of  $A_N \cap \mathcal{T} = \emptyset$  is that if  $P(s_i)$  holds for all  $i \leq N$ , then this property holds for all i. Q.E.D.

**Ill-posed problem: brief reminder.** The main *objectives* of science are to produce:

- quaranteed estimates for physical quantities; and
- guaranteed predictions for these quantities.

The problem is that estimation and prediction are ill-posed problems, i.e., small changes in the measurement result can lead to drastic changes in the resulting estimates.

Example: measurement devices are inertial, hence they suppress high frequencies  $\omega$ . So, the signals  $\varphi(x)$  and  $\varphi(x) + \sin(\omega \cdot t)$  are indistinguishable.

There exist many approaches to solve ill-posed problems: statistical regularization (filtering); Tikhonov regularization (e.g., assuming that  $|\dot{x}| \leq \Delta$ ); expert-based regularization, etc.; see, e.g., [27]. The main problem of all these approaches is that they provide no guaranteed bounds.

On physically meaningful solutions, problems become well-posed. Indeed, let us consider state estimation – an ill-posed problem.

A measurement process is a function f that maps state  $s \in S$  into observation  $r = f(s) \in R$ .

In principle, we can reconstruct s from r as  $s = f^{-1}(r)$ . The problem is that small changes in r can lead to huge changes in s, i.e., the inverse function  $f^{-1}$  not continuous.

**Proposition 5.2.** Let S be a definably separable metric space. Let T be a set of physically meaningful elements of S. Let  $f: S \to R$  be a continuous 1-1 function. Then, the inverse mapping  $f^{-1}: R \to S$  is continuous for every  $r \in f(T)$ .

**Everything is related: EPR paradox.** Due to *Relativity Theory*, two spatially separated simultaneous events cannot influence each other. By their paradox (see, e.g., [2,26]) Einstein, Podolsky, and Rosen (EPR) intended to show that in quantum physics, such influence is possible.

In formal terms, let x and x' be measured values at these two events. Independence means that possible values of x do not depend on x', i.e.,  $\mathcal{T} = X \times X'$  for some X and X'.

Physical induction implies that the pair (x, x') belongs to a set S of physically meaningful pairs.

**Proposition 5.3.** A set  $\mathcal{T}$  of physically meaningful pairs cannot be represented as  $X \times X'$ .

Thus, everything is related – but we probably can't use this relation to pass information (since the set  $\mathcal{T}$  isn't computable).

When to stop an iterative algorithm? The following situation is typical in numerical mathematics:

- we know an iterative process whose results  $x_k$  are known to converge to the desired solution x, but
- we do not know when to stop to guarantee that  $d_X(x_k, x) \leq \varepsilon$ .

A usual heuristic approach is to stop when  $d_X(x_k, x_{k+1}) \leq \delta$  for some  $\delta > 0$ .

For example, in physics, if 2nd order terms are small, we use the linear expression as an approximation.

When to stop an iterative algorithm: result.

**Definition 5.1.** Let  $\{x_k\} \in \mathcal{T}$ , k be an integer, and  $\varepsilon > 0$  a real number. We say that  $x_k$  is  $\varepsilon$ -accurate if  $d_X(x_k, \lim x_p) \leq \varepsilon$ .

**Definition 5.2.** Let  $d \ge 1$  be an integer. By a stopping criterion, we mean a function  $c: X^d \to R_0^+$  that satisfies the following two properties:

- If  $\{x_k\} \in \mathcal{T}$ , then  $c(x_k, \dots, x_{k+d-1}) \to 0$ . - If for some  $\{x_n\} \in \mathcal{T}$  and k,  $c(x_k, \dots, x_{k+d-1}) = 0$ , then

$$x_k = \ldots = x_{k+d-1} = \lim x_p.$$

**Proposition 5.4.** Let c be a stopping criterion. Then, for every  $\varepsilon > 0$ , there exists a  $\delta > 0$  such that if  $c(x_k, \ldots, x_{k+d-1}) \leq \delta$ , and the sequence  $\{x_n\}$  is physically meaningful, then  $x_k$  is  $\varepsilon$ -accurate.

#### 6 Relation with Randomness

Towards relation with randomness. Intuitively, if a sequence s is random, it satisfies all the probability laws such as the law of large numbers. Vice versa, if a sequence satisfies all probability laws, then for all practical purposes we can consider it random. Thus, we can define a sequence to be random if it satisfies all probability laws.

**Definition 6.1.** A probability law is a statement S which is true with probability 1: P(S) = 1.

So, we arrive at the following definition:

**Definition 6.2.** A sequence is random if it belongs to all definable sets of measure 1.

A sequence belongs to a set of measure 1 if and only if it does not belong to its complement C = -S with P(C) = 0. So, we arrive at the following equivalent definition:

**Definition 6.3.** A sequence is random if it does not belong to any definable set of measure 0.

Randomness and Kolmogorov complexity. Different definabilities lead to different randomness. When definable means computable, the corresponding Kolmogorov-Martin-Löf randomness can be described in terms of Kolmogorov complexity [16], the smallest length of a program that generates a given string:

$$K(x) \stackrel{\text{def}}{=} \min\{\text{len}(p) : p \text{ generates } x\}.$$

Crudely speaking, an infinite string  $s = s_1 s_2 \dots$  is random if, for some constant C > 0, we have  $\forall n (K(s_1 \dots s_n) \ge n - C)$ .

Indeed, if a sequence  $s_1 
ldots s_n$  is truly random, then the only way to generate it is to explicitly print it:  $print(s_1 
ldots s_n)$ . In contrast, a sequence like 0101...01 generated by a short program is clearly not random.

From Kolmogorov-Martin-Löf theoretical randomness to a more physical one. The above definition means that (definable) events with probability 0 cannot happen. In practice, physicists also assume that events with a *very small* probability cannot happen.

For example, a kettle on a cold stove will not boil by itself – but the probability is non-zero. If a coin falls head 100 times in a row, any reasonable person will conclude that this coin is not fair.

It is not possible to formalize this idea by simply setting a threshold  $p_0 > 0$  below which events are not possible. Indeed, then, for N for which  $2^{-N} < p_0$ , no sequence of N heads or tails would be possible at all. We cannot have a universal threshold  $p_0$  such that events with probability  $\leq p_0$  cannot happen.

However, we know that for each decreasing  $(A_n \supseteq A_{n+1})$  sequence of properties  $A_n$  with  $\lim p(A_n) = 0$ , there exists an N above which a truly random sequence cannot belong to  $A_N$ . Here is a resulting definition:

**Definition 6.4.** We say that  $\mathcal{R}$  is a set of random elements if for every definable decreasing sequence  $\{A_n\}$  for which  $\lim P(A_n) = 0$ , there exists an N for which  $\mathcal{R} \cap A_N = \emptyset$ .

Random sequences and physically meaningful sequences. Let  $\mathcal{R}_K$  denote the set of all elements which are random in Kolmorogov-Martin-Löf sense. Then, the following two results hold:

**Proposition 6.1.** Every set of random elements consists of physically meaningful elements.

**Proposition 6.2.** For every set  $\mathcal{T}$  of physically meaningful elements, the intersection  $\mathcal{T} \cap \mathcal{R}_K$  is a set of random elements.

**Proof:** When  $A_n$  is definable, for  $D_n \stackrel{\text{def}}{=} \bigcap_{i=1}^n A_i - \bigcap_{i=1}^\infty A_i$ , we have  $D_n \supseteq D_{n+1}$  and  $\bigcap_{n=1}^\infty D_n = \emptyset$ , so  $P(D_n) \to 0$ . Therefore, there exists an N for which the set of random elements does not contain any elements from  $D_N$ . Thus, every

set of random elements indeed consists of physically meaningful elements.

## 7 Proofs of Results Not Proven in the Main Text

A formal definition of definable sets.

**Definition 7.1.** Let  $\mathcal{L}$  be a theory. Let P(x) be a formula from  $\mathcal{L}$  for which the set  $\{x \mid P(x)\}$  exists. We will then call the set  $\{x \mid P(x)\}$   $\mathcal{L}$ -definable.

Crudely speaking, a set is  $\mathcal{L}$ -definable if we can explicitly *define* it in  $\mathcal{L}$ .

All usual sets are definable: the set of natural numbers  $\mathbb{N}$ , the set of real numbers  $\mathbb{R}$ , etc.

Not every set is  $\mathcal{L}$ -definable: indeed,

- every  $\mathcal{L}$ -definable set is uniquely determined by a text P(x) in the language of set theory;
- there are only countably many texts and therefore,
- there are only countably many  $\mathcal{L}$ -definable sets; so,
- some sets of natural numbers are not definable.

How to prove results about definable sets. Our objective is to be able to make mathematical statements about  $\mathcal{L}$ -definable sets. Therefore, in addition to the theory  $\mathcal{L}$ , we must have a stronger theory  $\mathcal{M}$  in which the class of all  $\mathcal{L}$ -definable sets is a countable set.

For every formula F from the theory  $\mathcal{L}$ , we denote its Gödel number by |F|. We say that a theory  $\mathcal{M}$  is stronger than  $\mathcal{L}$  if:

- $-\mathcal{M}$  contains all formulas, all axioms, and all deduction rules from  $\mathcal{L}$ , and
- $\mathcal{M}$  contains a predicate def(n, x) such that for every formula P(x) from  $\mathcal{L}$  with one free variable,

$$\mathcal{M} \vdash \forall y (\operatorname{def}(\lfloor P(x) \rfloor, y) \leftrightarrow P(y)).$$

**Existence of a stronger theory.** As  $\mathcal{M}$ , we take  $\mathcal{L}$  plus all above equivalence formulas.

Is  $\mathcal{M}$  consistent? Due to compactness property of first order logic, it is sufficient to prove that for any  $P_1(x), \ldots, P_m(x)$ ,  $\mathcal{L}$  is consistent with the equivalences corresponding to  $P_i(x)$ . Indeed, we can take

$$def(n, y) \leftrightarrow (n = |P_1(x)| \& P_1(y)) \lor \dots \lor (n = |P_m(x)| \& P_m(y)).$$

This formula is definable in  $\mathcal{L}$  and satisfies all m equivalence properties. Thus, the existence of a stronger theory is proven.

The notion of an  $\mathcal{L}$ -definable set can be expressed in  $\mathcal{M}$ : S is  $\mathcal{L}$ -definable if and only if

$$\exists n \in \mathbb{N} \, \forall y \, (\operatorname{def}(n, y) \leftrightarrow y \in S).$$

So, all the statements involving definability become statements from the  $\mathcal{M}$  itself, not from metalanguage.

## Consistency proof.

**Proposition 7.1.**  $\forall \varepsilon > 0$ , there exists a set  $\mathcal{T}$  of physically meaningful elements for which  $\underline{P}(\mathcal{T}) \geq 1 - \varepsilon$ .

**Proof.** Indeed, there are countably many definable sequences  $\{A_n\}$ :  $\{A_n^{(1)}\}$ ,  $\{A_n^{(2)}\}$ , ... For each k,  $P\left(A_n^{(k)}\right) \to 0$  as  $n \to \infty$ . Hence, there exists  $N_k$  for which  $P\left(A_{N_k}^{(k)}\right) \le \varepsilon \cdot 2^{-k}$ .

We take  $\mathcal{T} \stackrel{\text{def}}{=} - \bigcup_{k=1}^{\infty} A_{N_k}^{(k)}$ . Since  $P\left(A_{N_k}^{(k)}\right) \leq \varepsilon \cdot 2^{-k}$ , we have

$$\overline{P}\left(\bigcup_{k=1}^{\infty}A_{N_{k}}^{(k)}\right)\leq\sum_{k=1}^{\infty}P\left(A_{N_{k}}^{(k)}\right)\leq\sum_{k=1}^{\infty}\varepsilon\cdot2^{-k}=\varepsilon.$$

Hence, 
$$\underline{P}(\mathcal{T}) = 1 - \overline{P}\left(\bigcup_{k=1}^{\infty} A_{N_k}^{(k)}\right) \ge 1 - \varepsilon$$
.

**Proof of Proposition 3.2.** To compute the set  $R = \{x : f(x) = 0\}$  with accuracy  $\varepsilon > 0$ , let us take an  $(\varepsilon/2)$ -net  $\{x_1, \ldots, x_n\} \subseteq K$ .

For each i, we can compute  $\varepsilon' \in (\varepsilon/2, \varepsilon)$  for which  $B_i \stackrel{\text{def}}{=} \{x : d(x, x_i) \leq \varepsilon'\}$  is a computable compact set.

It is possible to algorithmically compute the minimum of a function on a computable compact set. Thus, we can compute  $m_i \stackrel{\text{def}}{=} \min\{|f(x)| : x \in B_i\}$ .

Since  $f \in T$ , similarly to the proof that equality of typical real numbers is decidable, we can prove that

$$\exists N \, \forall f \in T \, \forall i \, (m_i = 0 \vee m_i \ge 2^{-N}).$$

Computing  $m_i$  with accuracy  $2^{-(N+2)}$ , we can check whether  $m_i = 0$  or  $m_i > 0$ .

Let's prove that  $d_H(R, \{x_i : m_i = 0\}) \le \varepsilon$ , i.e., that

$$\forall i (m_i = 0 \Rightarrow \exists x (f(x) = 0 \& d(x, x_i) \le \varepsilon))$$

and

$$\forall x (f(x) = 0 \Rightarrow \exists i (m_i = 0 \& d(x, x_i) < \varepsilon)).$$

Indeed,  $m_i = 0$  means that  $\min\{|f(x)| : x \in B_i \stackrel{\text{def}}{=} B_{\varepsilon'}(x_i)\} = 0$ .

Since the set K is compact, this value 0 is attained, i.e., there exists a value  $x \in B_i$  for which f(x) = 0. From  $x \in B_i$ , we conclude that  $d(x, x_i) \le \varepsilon'$  and, since  $\varepsilon' < \varepsilon$ , that  $d(x, x_i) < \varepsilon$ . Thus,  $x_i$  is  $\varepsilon$ -close to the root x.

Vice versa, let x be a root, i.e., let f(x) = 0. Since the points  $x_i$  form an  $(\varepsilon/2)$ -net, there exists an index i for which  $d(x, x_i) \le \varepsilon/2$ . Since  $\varepsilon/2 < \varepsilon'$ , this means that  $d(x, x_i) \le \varepsilon'$  and thus,  $x \in B_i$ . Therefore,

$$m_i = \min\{|f(x)| : x \in B_i\} = 0.$$

So, the root x is  $\varepsilon$ -close to a point  $x_i$  for which  $m_i = 0$ .

**Proof of Proposition 4.1.** As A, given an instance i, we simply produce the result  $\omega_i$  of the i-th experiment.

Let us prove, by contradiction, that for every  $\varepsilon > 0$  and for every n, there exists an integer  $N \ge n$  for which

$$\#\{i \le N : i \in S_{\mathcal{P}} \& \omega_i = s_{\mathcal{P},i}\} > (1 - \varepsilon) \cdot \#\{i \le N : i \in S_{\mathcal{P}}\}.$$

The assumption that this property is not satisfied means that for some  $\varepsilon > 0$  and for some integer n, we have

$$\forall N_{>n} \# \{ i \le N : i \in S_{\mathcal{P}} \& \omega_i = s_{\mathcal{P},i} \} \le (1 - \varepsilon) \cdot \# \{ i \le N : i \in S_{\mathcal{P}} \}.$$

Let

$$T \stackrel{\text{def}}{=} \{x : \#\{i \le N : i \in S_{\mathcal{P}} \& x_i = s_{\mathcal{P},i}\} \le (1 - \varepsilon) \cdot \#\{i \le N : i \in S_{\mathcal{P}}\} \text{ for all } N \ge n\}.$$

We will prove that this set T is a physical theory (in the sense of the above definition); then  $\omega \notin T$ .

By definition, a physical theory is a set which is non-empty, closed, nowhere dense, and definable.

- Non-emptiness is easy: the sequence  $x_i = \neg s_{\mathcal{P},i}$  for  $i \in S_{\mathcal{P}}$  belongs to T.
- One can prove that T is closed, i.e., if  $x^{(m)} \in T$  for which  $x^{(m)} \to \omega$ , then  $x \in T$ .
- Nowhere dense means that for every finite sequence  $x_1 \dots x_m$ , there exists a continuation  $x \notin T$ . Indeed, for such an extension, we can take  $x_i = s_{\mathcal{P},i}$  if  $i \in S_{\mathcal{P}}$ .
- Finally, we have an explicit definition of T, so T is definable.

**Proof of Proposition 4.2.** Let us assume that  $P \neq NP$ ; we want to prove that for every feasible ph-algorithm A, it is not possible to have

$$\forall N (\#\{i \leq N : i \in S_{\mathcal{P}} \& \mathcal{A}(\omega, i) = s_{\mathcal{P}, i}\} = \#\{i \leq N : i \in S_{\mathcal{P}}\}).$$

Let us consider, for each feasible ph-algorithm  $\mathcal{A}$ ,  $T(\mathcal{A}) \stackrel{\text{def}}{=}$ 

$$\{x : \#\{i \le N : i \in S_{\mathcal{P}} \& \mathcal{A}(x, i) = s_{\mathcal{P}, i}\} = \#\{i \le N : i \in S_{\mathcal{P}}\} \text{ for all } N\}.$$

Similarly to the proof of the main result, we can show that this set T(A) is closed and definable.

To prove that T(A) is nowhere dense, we extend  $x_1 \dots x_m$  by 0s; then  $x \in T$  would mean P=NP.

If  $T(A) \neq \emptyset$ , then T(A) is a theory, so  $\omega \notin T(A)$ .

If  $T(A) = \emptyset$ , this also means that A does not solve all instances of the problem  $\mathcal{P}$  – no matter what  $\omega$  we use.

**Proof of Proposition 4.3.** Let us assume that no non-oracle feasible algorithm  $\delta$ -solves the problem  $\mathcal{P}$ . Let's consider, for each  $N_0$  and feasible phalgorithm  $\mathcal{A}$ ,

$$T(\mathcal{A}, N_0) \stackrel{\text{def}}{=} \{x : \#\{i \le N : i \in S_{\mathcal{P}} \& \mathcal{A}(x, i) = s_{\mathcal{P}, i}\} > \delta \cdot \#\{i \le N : i \in S_{\mathcal{P}}\} \text{ for all } N \ge N_0\}.$$

We want to prove that  $\forall N_0 (\omega \notin T(\mathcal{A}, N_0))$ .

– Similarly to the proof of the Main Result, we can show that  $T(A, N_0)$  is closed and definable.

- To prove that  $T(A, N_0)$  is nowhere dense, we extend  $x_1 \dots x_m$  by 0s.
- If  $T(A, N_0) \neq \emptyset$ , then  $T(A, N_0)$  is a theory hence  $\omega \notin T(A, N_0)$ .
- If  $T(A, N_0) = \emptyset$ , then also  $\omega \notin T(A, N_0)$ .

**Proof of Proposition 5.2.** It is known that if a f is continuous and 1-1 on a compact, then the inverse function  $f^{-1}$  is also continuous.

Let us recall that S is compact if and only if it is closed and for every  $\varepsilon$ , it has a finite  $\varepsilon$ -net, i.e., a finite set such that each element of S is  $\varepsilon$ -close to one of the elements from the set S.

We assume that the set X is definably separable, i.e., that there exists a definable sequence  $s_1, \ldots, s_n, \ldots$  which is everywhere dense in X.

The solution is to take  $A_n \stackrel{\text{def}}{=} - \bigcup_{i=1}^n B_{\varepsilon}(s_i)$ . Since  $s_i$  are everywhere dense, we have  $\cap A_n = \emptyset$ . Hence, there exists N for which  $A_N \cap \mathcal{T} = \emptyset$ . Since

$$A_N = -\bigcup_{i=1}^N B_{\varepsilon}(s_i),$$

this means  $\mathcal{T} \subseteq \bigcup_{i=1}^{N} B_{\varepsilon}(s_i)$ . Hence  $\{s_1, \dots, s_N\}$  is an  $\varepsilon$ -net for  $\mathcal{T}$ . So, the set  $\mathcal{T}$  is pre-compact. Q.E.D.

**Proof of Proposition 6.1.** Let T consist of physically meaningful elements. Let us prove that  $T \cap \mathcal{R}_K$  is a set of random elements.

If 
$$A_n \supseteq A_{n+1}$$
 and  $P\left(\bigcap_{n=1}^{\infty} A_n\right) = 0$ , then for  $B_m \stackrel{\text{def}}{=} A_m - \bigcap_{n=1}^{\infty} A_n$ , we have  $B_m \supseteq B_{m+1}$  and  $\bigcap_{n=1}^{\infty} B_n = \emptyset$ .

Thus, by definition of a set consisting of physically meaningful elements, we conclude that  $B_N \cap T = \emptyset$ .

Since 
$$P\left(\bigcap_{n=1}^{\infty} A_n\right) = 0$$
, we also know that  $\left(\bigcap_{n=1}^{\infty} A_n\right) \cap \mathcal{R}_K = \emptyset$ . Thus,  $A_N = B_N \cup \left(\bigcap_{n=1}^{\infty} A_n\right)$  has no common elements with the intersection  $T \cap \mathcal{R}_K$ . Q.E.D.

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