Geocentric cosmology: a new look at the measure problem

Mahdiyar Noorbala\textsuperscript{\textcopyright} and Vitaly Vanchurin\textsuperscript{\textcopyright}
Department of Physics, Stanford University, Stanford, CA 94305

We show that most of cutoff measures of the multiverse violate some of the basic properties of probability theory when applied repeatedly to predict the results of local experiments. Starting from minimal assumptions, such as Markov property, we derive a correspondence between cosmological measures and quantum field theories in one lesser dimension. The correspondence allows us to replace the picture of an infinite multiverse with a finite causally connected region accessible by a given observer in conjunction with a Euclidean theory defined on its past boundary.

INTRODUCTION

The old problem of initial conditions in the big bang cosmology was replaced by a measure problem in the inflationary cosmology. In order to make predictions in eternally inflating multiverse\textsuperscript{\textcopyright} one usually postulates a regularization scheme, or a measure\textsuperscript{\textcopyright}, which could be repeatedly tested against observations. In this paper we show that most of the cutoff measures are inconsistent with the most basic properties of the probability theory when applied to local experiments. To resolve the issue we propose a new approach to the measure problem which allows us to explore more easily the space of consistent measures.

CUTOFF MEASURES

In this section we show an inconsistency in the predictions of cutoff measures about local experiments. Let us consider a semi-classical picture of eternal inflation. We can think of the multiverse as a realization of a random process of formation of pocket universes with rates determined by the tunneling instantons. In addition, other random processes take place within each pocket universe which in our discussion amount to local experiments in local labs. We define a local lab as a spacetime region \( R \) which is small enough compared with the curvature scale to be safely approximated by a flat metric. The boundary of \( R \) can be split into a past \( R^- \) and a future \( R^+ \) such that no causal signal can reach \( R \) without passing through \( R^- \). In other words, a knowledge of the initial state on \( R^- \) is enough to determine the state everywhere on \( R \). We concentrate on those (double) experiments in which a particular initial measurement is performed on \( R^- \) with outcome \( \phi_1 \) and another final measurement on \( R^+ \) with outcome \( \phi_2 \). Since the total number of such labs is infinite we use a cutoff measure which regularizes the infinite four-volume of spacetime and results in finite probabilities.

Suppose that a cutoff measure is given. We can divide the regularized four-volume of spacetime into box regions of fixed size \( \ell \) and count regions with various \( \phi_{1,2} \). We then calculate the following probabilities by taking the cutoff to infinity\textsuperscript{\textcopyright}:

- \( \mu(\phi_1) \): The probability that a local experiment, confined to a box of size \( \ell \) in the multiverse, has outcome \( \phi_1 \) in the initial measurement.
- \( M(\phi_2, \phi_1) \): The joint probability that a local experiment, confined to a box of size \( \ell \) in the multiverse, has outcomes \( \phi_1 \) and \( \phi_2 \) in the initial and final measurement, respectively.
- \( M(\phi_2|\phi_1) \): The conditional probability that a local experiment, confined to a box of size \( \ell \) in the multiverse and with initial measurement outcome \( \phi_1 \), has outcome \( \phi_2 \) in the final measurement.

Evidently, all three functions \( \mu \), joint \( M(\cdot, \cdot) \) and conditional \( M(\cdot|\cdot) \) depend on the cutoff measure from which they are derived and on \( \ell \). Naively one may expect

\[
M(\phi_2, \phi_1) = \mu(\phi_1)M(\phi_2|\phi_1), ~ (1)
\]

which implies

\[
\mu(\phi_1) = \int d\phi_2 M(\phi_2, \phi_1), ~ (2)
\]

to hold true; but actually neither of Eqs. (1) and (2) are correct in general. The reason is that the two \( M \) functions (joint and conditional) search for whole regions \( R \) whereas \( \mu \) searches for past boundaries \( R^- \) only. However, below any cutoff surface there are more \( R^- \)'s than \( R \)'s. Were there equal number of them the above relations would hold; but \( M \)'s and \( \mu \) give probabilities on different sample spaces so this naive expectation is wrong.

Let us see this in a simple example. Suppose that there are only two possible outcomes: 1 and 2. We first calculate regularized \( M(1, 1) \) under a cutoff surface:

\[
M(1, 1) = \frac{N_\phi(1,1)}{N_\phi}, ~ (3)
\]

\footnote{More precisely, \( R^- \) is the smallest subset of the boundary of \( R \) such that its future domain of dependence \( D^+(R^-) \) contains \( R \).}

\footnote{We assume that these quantities do not depend on the arbitrariness that exists in tiling the spacetime with boxes.}
where $N_\Omega$ is the total number of regions $R$ and $N_\Omega(1,1)$ is the number of such regions with initial outcome 1 and final outcome 1. Similarly,

$$M(1|1) = \frac{N_\Omega(1,1)}{N_\Omega(1)} \cdot \mu(1) = \frac{N_\nu(1)}{N_\nu},$$

where $N_\nu$ is the total number of past boundaries $R^-$ and \(^{(1)}\) means that only regions with initial outcome 1 are counted. Clearly under any cutoff surface $N_\Omega \leq N_\nu$. We find from these relations that:

$$M(1,1) = \frac{N_\nu/N_\Omega}{N_\nu(1)/N_\Omega(1)} \cdot \mu(1)M(1|1),$$

which violates Eq. (I) unless in the limit that the cutoff goes to infinity the fraction is equal to 1, or if

$$\frac{N_\nu(1)}{N_\nu(1)} = \frac{N_\nu(2)}{N_\nu(2)}$$

But this can’t be true if, for some reason, the rate of motion of the cutoff surface is correlated with regions of high abundance of initial outcome 1 (or 2). Eq. (6) is satisfied if proper time is used.

It should be emphasized that there is no mathematical contradiction here; $\mu$ and $M$’s are three independent probability distributions that are not related to each other through Eqs. (I) and (2). However, if we use these three functions as predictions for the results of local experiments, then that would contradict basic properties of probability theory that we expect to hold for experimental probabilities.

**LOCAL EXPERIMENTS**

In this section we review local experiments and the process by which we can experimentally extract information from measurements. Unless explicitly stated, our discussion does not depend on whether the system is classical or quantum. A local observer performs experiments in a local lab and records the outcomes $\phi$ of the measurement. We assume that the same measurement can be done independently on copies of a system to yield probabilities on the possible outcomes. To begin, suppose that initially at $t_1 = 0$ the observer knows nothing about the state of the system, i.e., there is no previous measurement. Then he conducts a measurement and records the experimental probability distribution $W(\phi; t_1 = 0)$. He can perform a second measurement later and record the conditional probability $W(\phi_2; t_2|\phi_1; t_1 = 0)$. We will suppress the time label hereafter but we always take the time of the first measurement to be $t_1 = 0$. More measurements can be done in a similar fashion. Therefore the most general information an observer can acquire from successive measurements can be captured by weight functions of the form $W(\phi_n|\phi_{n-1}, \cdots, \phi_1)$, which satisfy the obvious normalization condition:

$$\int d\phi_n W(\phi_n|\phi_{n-1}, \cdots, \phi_1) = 1.$$  

It is possible to convey the same information via joint probability distributions. For example and for future reference, we can define the joint $W(\cdot, \cdot)$ for two measurements by

$$W(\phi_2, \phi_1) = W(\phi_1)W(\phi_2|\phi_1),$$

which has a different normalization condition:

$$\int d\phi_1 d\phi_2 W(\phi_2, \phi_1) = 1.$$  

Conversely, starting from a given joint $W(\cdot, \cdot)$ we can define

$$W(\phi) = \int d\phi' W(\phi', \phi)$$

and $W(\phi_2|\phi_1) = W(\phi_2, \phi_1)/W(\phi_1)$ which satisfy normalization conditions of the form (7). This establishes that the pair of functions $W(\cdot)$ and conditional $W(\cdot|\cdot)$ is equivalent to the single function joint $W(\cdot, \cdot)$.

Before proceeding, note that $\phi$ labels all possible outcomes of a measurement. Therefore $W$ depends on the sequence $\mathcal{E}_n$ of measurements performed at each step, although this is not explicit in our notation. So, for example, if we first measure the position and then the momentum of a particle we have $W(p_2|x_1)$, whereas if the second measurement is also a position measurement then we have $W(x_2|x_1)$. Here $\phi$ is either $x$ or $p$.

We say that $W$ has a “Markov property” (in time) associated to a sequence $(\mathcal{E}_i)_{i=1}^n$ of immediately successive experiments, if all experimental probabilities $W(\phi|\phi_{m+n}, \cdots, \phi_{m+1}, \cdots)$ are independent of the times $t_{m+1} \approx \cdots \approx t_{m+n}$ and of the number $m$, type and outcome of measurements prior to $t_{m+1}$. Here $\phi_{m+i}$ is the outcome of a measurement of type $\mathcal{E}_i$ (but the rest of the measurements are arbitrary). We call a minimal such $\mathcal{E}_i$ a “complete set of measurements” and say that its outcome $\Phi = \{\phi_i\}$ gives the “complete state” (or “state,” for short) of the system after the last measurement $\mathcal{E}_n$.

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3 It is important that the $R^+$ surface of each measurement constitutes the $R^-$ of the next one. Equivalently, we can say that the lab is shielded from the external effects so that experiments are not affected by any causal signal besides those of the previous ones.

4 Note that with this definition there may exist “hidden variables” that the experimentalist is blind to. For example, if one studies an ideal gas and chooses to measure only $P$, $V$ and $T$ then $(P, V)$ determines the state of the system. He is simply blind to individual atom velocities and positions that, according to statistical mechanics, constitute the state of the system. However, a thermodynamicist would agree with the him on calling $(P, V)$ the state of the system.
We refer to any subset of a complete state as an incomplete state (but employ the same symbol φ as for a single outcome). We further assume that the outcomes of these incomplete measurements are independent of each other (commuting, in the context of quantum mechanics).

With this notation we can write:

\[ W(\phi_n|\Phi_{n-1}, \cdots) = W(\phi_n|\Phi_{n-1}). \] (11)

It follows that as far as the future evolution is concerned all that matters is the initial state of the system, known at the last complete measurement, not how it got to that state:

\[ W(\phi_n|\phi_{n-1}, \cdots, \Phi_1, \cdots) = W(\phi_n|\phi_{n-1}, \cdots, \Phi_i). \] (12)

Note that if \( \Phi_1 \neq \Phi_2 \) are obtained from the same set of complete measurements, then they describe distinct states. But it is possible that two sets of complete measurements lead to two totally different, but equivalent, sets of outcomes \( \Phi \) and \( \Phi' \) which describe the same state. \( \Phi \) and \( \Phi' \) are equivalent if they lead to identical subsequent evolutions:

\[ W(\phi|\Phi) = W(\phi|\Phi'), \quad \text{for all } \phi. \] (13)

The set of equivalence classes of \( \Phi \)'s defines the space of states of the system.

The preceding definition of state was from the point of view of an experimentalist who has access to \( W \). From the theoretical perspective, states and their Markovian evolution play a central role and \( W \) is derived from them as a prediction. This leads to numerous correlations among \( W \)'s and hence to a huge reduction in the information they carry. Let us review how it works. In classical mechanics the state of the system is a point \((x, p)\) in its phase space which evolves according to a Hamiltonian. Everything is deterministic except the initial state (just before the first measurement) which is determined by a probability distribution \( P(x, p) \). The result of all other experiments are found from the probability distribution induced on the measured observable. In quantum mechanics a closed system the state of the system is a vector in its Hilbert space. Again the evolution is governed by a Hamiltonian, but the outcomes of measurements are not deterministic: \( \lim_{t_2 \to t_1} W(\Phi_2|\Phi_1) = |\langle \Phi_2|\Phi_1 \rangle|^2 \), where \( |\Phi_{1,2} \rangle \) are vectors corresponding to the states \( \Phi_{1,2} \).

The initial state (just before the first measurement) is determined by a probability distribution on the Hilbert space which is practically equivalent to a density matrix: \( W(\Phi) = \langle \Phi|\rho|\Phi \rangle \).

In the next section we use these properties of \( W \), especially the Markov property, to address the problem encountered in the previous section.

**GEOCENTRIC MEASURES**

We can now summarize the problem with cutoff measures we found above: If we identify

\[ W(\phi) \equiv \mu(\phi), \] (14)
\[ W(\phi_2|\phi_1) \equiv M(\phi_2|\phi_1), \] (15)
\[ W(\phi_2, \phi_1) \equiv M(\phi_2, \phi_1), \] (16)

then Eqs. (8) and (10) imply Eqs. (1) and (2), but the latter are not necessarily satisfied for all cutoff measures. One way to escape this puzzle is to count only those initial outcomes for which the whole experiment is under the cutoff: \( \mu(1) = N_\delta(1)/N_\delta \). The problem with this solution is that we then have a measure which is modified in a different way for every type of experiment. It is also possible that this puzzle stems from the insufficiencies of semi-classical description of quantum cosmology and that we are using the measure more than once (see Appendix B in [14]).

On the other hand, we can simply forget about computing \( M(\cdot, \cdot) \) and \( M(\cdot|\cdot) \) from the measure and be content with what we get for \( \mu \). We can define \( M(\cdot, \cdot) \) such that Eq. (14) holds but we still need \( M(\cdot|\cdot) \). This motivates our approach in this section where we derive all probabilities from an initial state probability (corresponding to \( \mu \)) and some evolution or transition probability (corresponding to conditional \( M(\cdot|\cdot) \)). We study these generically and will not confine ourselves to \( \mu \)'s that are derived from the known cutoff measures.

The Markov property [11] suggests that a whole set of \( W \)'s can be produced by only two functions:

\[ \mu(\Phi) = W(\Phi), \] (17)

and

\[ P(\Phi_2|\Phi_1) = W(\Phi_2; t_2 = T|\Phi_1; t_1 = 0), \] (18)

which we call the “measure” \( \mu \) and the “local transition probability” \( P \).

In most physical theories, the local transition probability is determined by a unitary evolution via a Hamiltonian. But in this paper we are not interested in properties of \( P \) (which we assume to be given) and would like to study \( \mu \). It is clear that the shielded experiments that we discussed earlier do not help. Instead we have to let in external signals from outside of the already-explored initial surface \( R^- \) to probe information about \( \mu \) that has not yet been accessible.

So far we have only made use of the Markov property of \( W \) in time to break it into \( \mu \) and \( P \). As mentioned before, according to quantum mechanics the most general \( \mu \) is a density matrix. We restrict to a smaller class of \( \mu \)'s by considering a classical field whose state is given by \( \varphi(x) \) and its conjugate momentum \( \pi(x) \) defined on the initial surface \( R^- \). These field configurations correspond to coherent states in the quantum theory and hence reduce the
density matrix to a functional $\mu[\varphi, \pi]$. Now we impose an additional constraint on $\mu$ by demanding that the probabilities for field values in a region depend on those outside only through their common boundary. In other words, we can find the probability distribution everywhere by propagating our knowledge outward, hence possible names lab-centric, observer-centric or simply *geocentric*. For instance on a discrete 1D $R^-$:

$$
\mu(\varphi_1, \ldots, \varphi_N) = \mu(\varphi_1|\varphi_2)\mu(\varphi_2|\varphi_3)\ldots\mu(\varphi_N),
$$

where $\varphi_i$ stands for both $\varphi(x_i)$ and $\pi(x_i)$. We shall refer to this assumption as the Markov property in space, which by iteration enables us to write $\mu$ as the exponentiation of a local functional of fields. On a 3D $R^-$:

$$
\mu[\varphi, \pi] = \exp\left[-\int_{R^-} d^3x \mathcal{L}(\varphi(x),\nabla \varphi(x); \pi(x), \nabla \pi(x))\right],
$$

where

$$
\mathcal{L}(\varphi(x), \nabla \varphi(x); \pi(x), \nabla \pi(x))d^3x = -\log \mu(\varphi(x), \pi(x)|\varphi(x + dx), \pi(x + dx)).
$$

Therefore, we see that with the assumptions of Markov property in space and time the measures on classical fields is equivalent to the space of 3D Lagrangians, and the search for a measure is completely equivalent to a search for a 3D Lagrangian.

As an example, consider a classical statistical mechanics of electromagnetic field at finite temperature $T$. In temporal gauge the partition function is given by

$$
Z = \int e^{-\frac{1}{T} \int d^3x \mathcal{H}(P,A)} DADP,
$$

where $P_i = \dot{A}_i$, $F_{ij} = \partial_i A_j - \partial_j A_i$ and $\mathcal{H} = \frac{1}{4} F_{ij}^2 + \frac{1}{2} P_i^2$. It follows that the corresponding 3D Lagrangian for a thermal bath of photons could be written as

$$
\mathcal{L} = \frac{\mathcal{H}}{T} = \frac{1}{4T} F_{ij}^2 + \frac{1}{2T} P_i^2.
$$

Of course additional constraints must be added to the 3D path integrals of Eqs. (20) and (22) to impose the Gauss’s law. Therefore a geocentric measure determined by Eq. (22) can describe a thermal CMB radiation with $T \approx 2.7K$ given that appropriate additional fields and/or interactions are added to explain, for example, the acoustic peaks.

It is also possible to obtain a consistent measure from any of the cutoff measures by using them only once to compute $\mu$ on the initial surface at the time of the first measurement. For example, in the scale factor cutoff measure Eq. (6) and hence Eq. (11) are not satisfied if “1” and “2” refer to high and low cosmological constant. But if we use this measure only to compute $\mu$ but not $M(\cdot|\cdot)$ then it is consistent but may or may not have the Markov property in space. Similarly, from the causal diamond measure one should only obtain $\mu$, although there are no problems with Eq. (11) if the lab regions $R$ are themselves causal diamonds. On the other hand, the stationary measure does not suffer from this problem because different cutoff surfaces are used in making predictions for events with different stationarity times and Eq. (11) is not relevant.

**CONCLUSION**

We have shown that most of the cutoff measures of the multiverse suffer from severe inconsistencies and developed a new framework which allows us to study the measure problem from a completely different perspective. In the emerging picture an infinite multiverse is replaced with a finite geocentric region, and the search for the correct measure is replaced by a search for a 3D Lagrangian yet to be discovered.

There are two ways to look for the correct Lagrangian. One could either try to perform direct phenomenological searches or one could try to derive it from first principles. For the phenomenological approach one has to reinterpret the existing cosmological data from the geocentric view point. Although we are guaranteed to uncover some Lagrangian it is not a priori guaranteed that the corresponding Lagrangian will be local, simple or even useful. On the other hand it would be interesting to see whether one can derive the corresponding theory from some yet-unknown first principles.

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