Covariant influences for finite discrete dynamical systems

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We develop a rigorous theory of external influences on finite discrete dynamical systems, going beyond the perturbation paradigm, in that the external influence need not be a small contribution. Indeed, the covariance condition can be stated as follows: if we evolve the dynamical system for \( n \) time steps and then we disturb it, it is the same as first disturbing the system with the same influence and then letting the system evolve for \( n \) time steps. Applying the powerful machinery of resource theories, we develop a theory of covariant influences both when there is a purely deterministic evolution and when randomness is involved. Subsequently, we provide necessary and sufficient conditions for the transition between states under deterministic covariant influences and necessary conditions in the presence of stochastic covariant influences, predicting which transitions between states are forbidden. Our approach, for the first time, employs the framework of resource theories, borrowed from quantum information theory, to the study of finite discrete dynamical systems. The laws we articulate unify the behavior of different types of finite discrete dynamical systems, and their mathematical flavor makes them rigorous and checkable.

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I. INTRODUCTION

Dynamical systems describe the evolution of several interesting situations. There are essentially two flavors: continuous dynamical systems, which are more often studied, and discrete dynamical systems, which are the subject here. In many cases, especially for continuous dynamical systems, when the evolution is particularly complex to deal with, one splits the evolution into two parts: the uninfluenced part and the perturbation. The latter is interpreted as a small correction to the uninfluenced evolution. In this article, we go beyond the perturbation paradigm by introducing the notion of covariant influence, which need not be a small contribution. Such a splitting can be done on the basis of time scale: we can think of the covariant
influence as acting on a significantly longer time scale than the uninfluenced evolution. A small influence is a perturbation, so to get a more interesting behavior we must also go beyond small influences. The covariance condition guarantees that, despite not being small, it preserves the underlying structure representing the evolution of states. In this way, the evolution of a discrete dynamical system can be written as the composition of two evolutions: one is understood as the basic evolution of the system, and the other is the covariant influence. The covariance requirement ensures that the order in which these two parts are applied does not matter.

In this article, for the first time, we develop a general theory of covariant influences in discrete dynamical systems with a finite number of states. Such a theory comes in two flavors. The first is a deterministic one, where randomness is completely forbidden both in the initial state and in the action of the covariant influence. In the second, instead, we allow the presence of randomness both in the initial state and in the covariant influence. We show that deterministic influences allow hopping between attractors whose length becomes smaller and smaller. Instead, in the presence of randomness all jumps between attractors become possible. In particular, we achieve a full characterisation of transitions between states in the deterministic setting, and in the random case, we predict which transitions between states are forbidden.

Our approach, for the first time, employs the framework of resource theories, borrowed from quantum information theory, to the study of discrete dynamical systems. This also constitutes the first application of resource theories outside the physics domain, to a field with countless applications to diverse areas of science, including genetic regulatory networks. The advantage of the resource-theoretic approach is the separation of the evolution of a discrete dynamical system into an “uninfluenced” contribution and an influence, as is standard practice when studying perturbations. However, the influence need not be a perturbative (i.e., small) contribution [1, 2].

Our results for discrete dynamical systems are obtained from first principles that are of a mathematical flavor, as they are grounded in category theory. This ensures that our analysis is rigorous, logical, and checkable. In particular, the use of resource theories allows us to get a unified picture of discrete dynamical systems under influences, regardless of the specifics of their evolutions, unlike most standard approaches to discrete dynamical systems. In this way, our results, in the form of simple mathematical laws, can be phrased in general terms, so they are applicable to a broad class of discrete dynamical systems. The key concept in our analysis is covariance, which can be thought of as a symmetry in time evolution, and can be expressed by a simple commutativity condition. The general laws we find are predictive for discrete dynamical systems where covariance is in force, and motivate novel experimental work.

The article is organized as follows: in §II we present the background of this work, starting with a presentation of discrete dynamical systems, and ending with a presentation of resource theories. In §III, we introduce the covariance condition, showing that it gives rise to a well-behaved resource theory. In §IV, we present our results both for deterministic dynamical systems and dynamical systems with randomness. A summary of the results is presented in §VI, and conclusions are drawn in §VII.

II. BACKGROUND

In this section we provide key background elements that underpin our contributions and advances in subsequent sections. We begin in §II A with an overview of a discrete dynamical system, which is a set of elements that evolve discretely in time. Then we present random Boolean networks as a special case of discrete dynamical systems in §II B. The subsequent subsection II C provides essential elements of resource theories.

A. Dynamical systems

In this subsection we elaborate on essential background concerning discrete dynamical systems. We begin with explaining general notions of discrete dynamical systems including the set of states and the mapping from states to states corresponding to discrete evolution. Then we discuss features of discrete dynamical systems such as attractors and basins of attraction, which arise in studies of dynamics. A discrete dynamical system can be deterministic or stochastic, with stochasticity in this case pertaining to states as the dynamics is always deterministic.

1. Introduction

Now we introduce the basics of discrete dynamical systems. Specifically, we define dynamical systems formally. Then we present a graphical representation of the dynamical system in the form of a graph.

A dynamical system is a pair \((S, \phi_t)\) comprising a set \(S\) of elements, with each element called a state \(s\) of the dynamical system, and a family of functions \(\{\phi_t\}\), which maps \(S\) to itself, and describe the evolution at time \(t\). If the label \(t\) is in \(\mathbb{N}\) or \(\mathbb{Z}\), the dynamical system is called discrete; if \(t\) is in \(\mathbb{R}\), the dynamical system is called continuous. The standard requirements for such a family of functions \(\{\phi_t\}\) are the following:

1. \(\phi_0 = \mathbb{I}\), where \(\mathbb{I}\) denotes the identity on \(S\);
2. \(\phi_{t+s} = \phi_t \circ \phi_s = \phi_s \circ \phi_t\);
3. $\phi^{-1} = \phi_t^{-1}$.

In this article, we focus on discrete dynamical systems, where time runs in $\mathbb{N}$, and the number of states is finite. Therefore, in this work, whenever we talk about discrete dynamical systems, we always mean finite discrete dynamical systems. For discrete dynamical systems, we have discrete time steps, and we define $\phi$ to be the function on $s$ describing the evolution of $s$ to itself in a single time step. With this in mind, we have

$$\phi_n := \phi^n,$$

(1)

where $\phi^n$ is the composition of $\phi$ with itself $n \in \mathbb{N}$ times and refers to the evolution of the system over $n$ time steps: the evolution of a state $s$ after $n$ time-steps is given by $\phi^n(s)$. For this reason we say that $\phi$ is the generator of the dynamics, and in this case we denote the dynamical system simply by $(S, \phi)$. By convention, $\phi^0 = 1$. From the second of the conditions above, note that evolving the system for $n$ time steps, and then for other $m$ time steps, is the same as evolving it for $n + m$ time steps. Therefore, the set $\{\phi^n\}_{n \in \mathbb{N}}$ of the powers of $\phi$ (1) has the algebraic structure of a commutative monoid, but, in general, not of a group because monoids are more general structures that include groups, which allow non-invertible evolution.

In many dynamical systems there are external influences that can change the evolution of the system. Think, e.g., of perturbation theory, where we view the perturbation as a small correction to the evolution of the system [1–3]. If we remove the requirement of the perturbation to be small, we have an influence. In mathematical terms, an influence is simply a mapping

$$f : S \rightarrow S.$$  
(2)

Influences form a monoid under function composition (the identity is trivially an influence), but it is different from the dynamical monoid because it is not generated by a single influence, and as such it is not commutative. Note that we take $f$ to be in a monoid and not a group because we want to allow general behavior. A dynamical graph is a representation of a dynamical system. A directed graph $G$ comprises an ordered set of $M$ vertices $S$ (also “nodes”), and directed edges $E$, which are ordered pairs $(s, s') \in S \times S$. The dynamical graph represents the dynamics corresponding to the evolution $\phi$: there is a directed edge from a vertex $s$ to a vertex $s'$ if $s' = \phi(s)$. As the dynamics are deterministic, there is only one outgoing edge for each vertex of the dynamical graph. An example of a dynamical graph is depicted in Fig. 1.

2. **Features**

We now describe features including cycles, attractors, transient states and basins of attractions, which are the fundamental features of a dynamical system. First we formally define each of these three features. Then we explain how these features are represented in the dynamical graph, which represents the dynamical system.

As the dynamics is deterministic on a finite number of states, at some point the dynamics must repeat itself. Therefore, there is at least one cycle. More formally, a cycle $C$ is a set of states $\{s\}$ such that there exists a positive integer $\ell$ such that, for each $s \in C$, we have $\phi^\ell(s) = s$. The minimum integer $\ell$ is called the period or length of the cycle. Note that the length of the cycle is also the number of states in the cycle. In general, in a dynamical system there are a number of cycles of different periods. Those cycles of period 1 are called fixed points. In this case,

$$\phi(s) = s.$$  
(3)

Cycles are loops in the dynamical graph such that the loop is oriented by directed edges pointing in one way.

An attractor is a set of states that is forward-invariant under the evolution $\phi$. In other words, it is a set of states where evolution “repeats itself”. As we are working with a finite dynamical system, attractors are actually cycles. Thus, all states end up in some cycle. Some states do not belong to any cycles; they are called transient because it is not possible to come back to them, not even after enough time steps.

In the dynamical graph, transient states are the vertices of the non-cyclic parts of the graph, i.e., everything that is not in a loop. Given an attractor, its basin of attraction is the set of all states that end up in that attractor after enough time steps. Notice that, since the dynamics is deterministic, basins of attractions relative to different attractors are disjoint.

![Dynamical graph](image-url)
i.e., there can be no state belonging to more than one basin of attraction. In other words, the set of states $s$ is partitioned into basins of attraction. The dynamical graph is made of several connected components: each of them represents a basin of attraction. With a little abuse of terminology, we often say that a state, even a transient one, has length $\ell$ if $\ell$ is the length of the attractor in its basin of attraction with the notion of basins made clear in Fig. 1.

3. Introducing randomness

Now we explain how randomness is introduced into the dynamical system. In some situations, the influence is best described as a stochastic process that activates jumps between the states of a dynamical system according to a certain probability distribution. In such a stochastic setting, we are no longer sure about the state of the dynamical system, which demands the introduction of randomness in the sense of a probability distribution over the states of the dynamical system. Despite randomness, evolution is still represented as a deterministic linear map over the set of state probability distributions.

We now represent randomness on states $s \in S$ by constructing the $\sigma$-algebra of the power set $\varphi(S)$. In this case, we first fix an ordering on the elements of $s$. To represent probability measures, it is enough to consider probability vectors $p \in P^M(\mathbb{R})$, whose entries represent the probability assigned to each of the $M$ singleton subsets of $S$, each corresponding to one of the states. Specifically, each canonical-basis vector $e_i$ represents the dynamical state $s_i$ (namely the $i$th state in the ordering) because $e_i$ as a probability vector represents the certainty of being in the state $s_i$.

Each probability vector $p$ can be visualized in the dynamical graph as follows: each entry of $p$ represents the probability of the corresponding vertex of the dynamical graph. In other words, we are labeling the vertices of the dynamical graph with the entries of $p$, i.e., vertex $s_i$ is labeled $p_i$, which we imagine as a continuous greyscale coloring of vertices in the graph.

The functions $\{\phi^n\}_{n \in \mathbb{N}}$ describing evolution in time can be extended by linearity

$$\phi^n : P^M(\mathbb{R}) \rightarrow P^M(\mathbb{R})$$

$$: e_i \mapsto e_j$$

(4)

such that $\phi^n(s_i) = s_j$, where $s_i$ is the state associated with the vector $e_i$, for every $n \in \mathbb{N}$. With a little abuse of notation, we also write the mapping in Eq. (4) as $e_i \mapsto e_{\phi(i)}$. As the generator of the dynamics $\phi$ is a linear map, it can be represented by a matrix $\Phi$, called dynamical matrix, once we fix the canonical basis in the domain and the codomain of $\phi$: its columns are given by the action of $\phi$ on the canonical basis. Note that $\Phi$ can be viewed as a column-stochastic matrix: every entry is non-negative and columns sum to 1 (there is only one 1 per column). In this spirit, $\Phi_{ij}$ is the probability for state $s_j$ to transition to state $s_i$. According to graph theory, the matrix $\Phi$ is nothing but the transpose of the adjacency matrix of the dynamical graph.

B. Random Boolean networks

Random Boolean networks, introduced in 1969 by Kauffman [4, 5], were the first successful theoretical model to explain genetic regulation in cells. They have been the privileged model to understand the expression of genes for fifty years [6], attracting the attention of biologists and physicists working on complex systems as well [7, 8]. A random Boolean network can be described as a discrete dynamical system; as such, it can be represented as a dynamical graph.

1. Introduction

We now discuss underlying concepts of random Boolean networks. First we explain how the network represents which genes are expressed and which are not expressed at a given time. Then we elaborate on how Boolean functions describe evolution from one configuration of gene expression to another. Finally, we explain what is random per se in a random Boolean network.

Random Boolean networks are directed graphs, whose $N$ vertices represent genes, and directed edges represent the mutual influence of genes on the expression of other genes.

Each vertex can be labeled black or white, representing an expressed or an unexpressed gene, respectively. We can represent the state of all genes, i.e. of the whole network, as a Boolean string

$$g = g_0g_1g_2 \cdots g_{N-1} \in \{0, 1\}^N$$

(5)

(once we fix a total ordering of the vertices) such that $g_i$ represents expression or not of the $i$th gene, whether suppressed ($g_i = 0$) or expressed ($g_i = 1$). An example of such a network is presented in Fig. 2.

The state of each gene $g_i$ evolves to the next time step according to the states of the $K_i$ parent genes in the previous time step, i.e. graphically, those genes from which a directed edge points to the $i$th gene. Mathematically, this random Boolean network describes evolution of the state $g$ of the whole network according to $N$ Boolean functions

$$f_i : \{0, 1\}^{K_i} \rightarrow \{0, 1\}$$

(6)

on the $i$th vertex with the output determining the black-or-white label for whether gene $g_i$ is expressed or not (black or white, respectively) at the next time step. All Boolean functions $\{f_i\}$ stay the same at every time step. In this setting, the expression of a gene at each time step
FIG. 2. This random Boolean network contains 4 genes (represented by the vertices of the graph). White vertices represent unexpressed genes; black vertices represent expressed ones. Directed edges represent which genes influence one another: the gene at the tip of an edge is influenced by the gene at the tail of the edge.

depends on which genes are expressed at the previous step.

The vertices of the directed graph, representing genes, are fixed. Instead, randomness is only initial, and comes in three ways:

1. randomness in which vertices are labeled black or white at the initial step;
2. randomness in the choice of directed edges between the vertices;
3. randomness in the choice of Boolean function associated with each vertex.

Once these features are chosen, they stay the same for all time steps. In practice, the evolution can still be regarded as deterministic. For each of these three forms of randomness, various prior probability distributions have been proposed [8]. To represent the randomness in the initial state, we resort to probability vectors, as discussed above. Here we have $M = 2^N$ states $\{g\}$ of the whole network. Therefore, randomness in the initial state is described by a probability vector $p \in \mathbb{P}^{2^N}(\mathbb{R})$. As the evolution can still be regarded as deterministic, for the scope of this article, we do not need a precise mathematical treatment of the other two types of randomness.

2. Random Boolean networks as discrete dynamical systems

Now we explain how random Boolean networks can be viewed as discrete dynamical systems. First, we show how to represent time evolution, and then how to construct the associated dynamical graph.

The evolution of a random Boolean network arising from the edges of the graph, and the choice of Boolean functions can be represented in a more high-level fashion. Take the set of Boolean strings $\{g\}$, which, as we saw above, represent the states of the network. The action on the Boolean functions of the network makes a bit string

$$g = g_0 g_1 g_2 \cdots g_{N-1}$$

(7)

evolve to a new bit string

$$g' = g'_0 g'_1 g'_2 \cdots g'_{N-1},$$

(8)

where $g'_i$ denotes the evolution of $g_i$ at the next time step under the action of $f_i$. We can represent the evolution of $g$ into $g'$ in a single time step as a function $\phi : \{g\} \rightarrow \{g\}$ on the set of bit strings of $N$ elements. Then the pair $((\{g\}, \phi^n))$ is a discrete dynamical system.

Given that a random Boolean network can be described as a dynamical system, one can construct the associated dynamical graph, which is also known as the "state space" in the literature on random Boolean networks [7, 8]. The rules to construct it are the same as presented for generic discrete dynamical systems: vertices are states of the network, and we use $\phi$ to draw directed edges.

3. Network of genes and mutual interaction

Now we present the biological significance of random Boolean networks, explaining how they have been studied in the literature. Specifically, different phases have been identified as far as the dynamics in the "thermodynamic limit" are concerned. In this setting, different attractors of the associated dynamical system are linked to different types of cells. Randomness arises because of the complexity in describing the noisy environment in which cells live.

Features of the discrete dynamical system associated with random Boolean network have a biological counterpart. For instance, it is argued that each attractor represents a different cell type [4–6]. This has also been verified experimentally [9].

Most of the analysis has been carried out in the so-called "thermodynamic limit", where one considers an "infinite" network, mimicking the very large number of genes in the human DNA, and studying the dynamics arising on it. In this scenario, a phase transition has been observed, separating an "ordered phase", where the dynamics are described by short and simple cycles, from a "chaotic phase", where the dynamics are described by complex and long cycles.

From the point of view of the adaptation of a cell to the environment, the ordered phase is too static and rigid, and does not provide any adaptive power to cells. On the other hand, the chaotic phase is too unstable to guarantee a successful adaptation to the environment. Therefore, Kauffman identified the "golden phase" to explain the
origin of the various cell types in the critical regime, separating the ordered phase from the chaotic one [4, 5]. According to the model of genetic regulation described by random Boolean networks, the number of attractors in the critical regime follows a power law: \( N^{0.63} \), where \( N \) is the number of genes [6, 10]. Empirical data show instead a power scaling given by \( N^{0.88} \) [6, 11]. This discrepancy in the predictions has led some to question the actual value of random Boolean networks as a model for genetic regulation [8]; nevertheless they still remain a strong and powerful model, and in any case, are interesting from the point of view of the physics of complex systems.

In the study of random Boolean networks, there is the implicit presence of an external environment, defined as everything that is outside the cell under consideration. We assume that the environment is computationally hard to describe; therefore, it is conveniently represented by the introduction of randomness in the form of probability vectors, as described above [7, 8, 12].

C. The framework of resource theories

In this subsection we introduce resource theories, which are a powerful mathematical framework for dealing with the notion of “resources” in a rigorous fashion, developed in quantum information [13–18] and beyond [19–26]. We start by explaining their relevance to quantum theory, and then we continue by formulating them in a theory-independent way, using process theories. This allows us to extend their scope to discrete dynamical systems in the rest of this article.

1. Motivation

Resources theories have been remarkably successful in quantum information, providing new insights into the manipulation of quantum resources, such as quantum entanglement. The notion of a resource theory is based on the fact that, in several cases, we cannot implement all possible operations because of some constraints. The processes that respect the constraints are called “free” because they do not require any resources to occur.

Resource theories were developed to manage resources of quantum information optimally, in particular quantum entanglement [27, 28]. However, they have been extended to several other quantum information areas, such as quantum thermodynamics [29, 30], the study of quantum coherence [31] and quantum reference frames [32–36], and quantum computing [37, 38]. In these areas they have produced many important results, making a huge contribution to the development of new research areas.

In many cases, constraints restrict what processes can occur, in the sense that breaking such constraints is possible, but requires additional resources. For instance, in a thermodynamic setting a system is often immersed in an environment at a certain temperature. In this situation, the system reaches thermal equilibrium spontaneously, according to the “minus first law” of thermodynamics [39]. This law restricts what thermodynamic processes can occur: they can only bring a state closer and closer to the thermal equilibrium with the environment. This can be argued to be the origin of the time-asymmetry in thermodynamics [39]. All the processes that go against the minus first law are a resource because they allow us to restore part of the time symmetry, and indeed they require some external work to be implemented. For this reason, whenever there is a restriction, the allowed processes are called “free operations”, and those that enable us to overcome the restriction are called “non-free”. Among processes, a particular type of them are those corresponding to the preparation or initialization of a system, which are identified with the states of the system itself. Again, free states correspond to preparations that can be done in the presence of the restriction, whereas non-free states require some resource to be implemented.

The goal of a resource theory is to describe which transitions between the states of the theory are possible in the presence of the restriction, namely under free operations. This is known as the “conversion problem”.

2. Process theories

Now we introduce process theories [40–42], which are a mathematical framework where the notion of process plays center stage. In these theories, one can describe systems as well as their evolution and interaction through processes. Their abstract formalisation is based on the notion of strict symmetric monoidal category [43, 44]. Process theories provide the necessary background to formulate resource theories outside the quantum setting.

In a process theory, systems are treated as labels to identify different inputs/outputs. Systems can be composed to yield a composite system. At this stage, it is useful to introduce a particular type \( I \), called the trivial system, which represents the lack of a system. Clearly composing any system \( a \) with \( I \) yields \( a \) itself. Example systems are chemical species, in chemistry, or Hilbert spaces arising in quantum physics.

As the name suggests, the core of a process theory are processes. Intuitively, a process is anything that happens between systems that has zero or more inputs and zero or more outputs [42]. The canonical example is a function, which has one or more inputs, and has an output. However, we can find more physical examples, such as chemical reactions, physical transformations. Therefore, we represent a process as \( f : A \rightarrow B \), where \( a \) is the input type, and \( B \) is the output type. When they share a type, two processes \( f : A \rightarrow B \) and \( g : B \rightarrow C \), like functions, can be composed sequentially to obtain a new process \( g \circ f \). This means feeding the output of \( f \) into the input of \( g \). It is natural to require such a composition to be associative, and to have an identity—the identity
process—which corresponds to “doing nothing”. This construction makes a process theory a category \[44\], with systems as objects and processes as morphisms. When we have a composite type, \(AB\) we can consider processes that run “independently” or in parallel on \(a\) and \(B\), i.e. two processes \(f : A \rightarrow C\) and \(g : B \rightarrow D\) that occur “at the same time” taking \(a\) and \(B\) as input, respectively. Clearly, these two processes occurring together must yield a valid process on the composite system \(AB\). This new process, called the parallel composition of \(f\) and \(g\), is denoted by

\[
f \otimes g : AB \rightarrow CD\tag{9}\]

It is natural to assume some properties of the parallel composition that make it interact nicely with the sequential composition \[42\].

So far, we have not touched on the issue of the order of composition of systems, namely whether \(AB = BA\). Intuitively, we do not expect these two composite systems to be very different. However, in some cases, there is a conceptual difference. For example, consider a flask divided into two parts, containing two different gases \(a\) and \(B\), one in each part. \(AB\) can stand for “gas \(a\) is in part 1 and gas \(B\) is in part 2”; therefore, \(BA\) stands for “gas \(B\) is in part 1 and gas \(a\) is in part 2”. Although these two settings are equivalent from the point of view of the chemical species involved, we cannot say they are exactly the same setting. Therefore, we can say that they are “equivalent”. To make this notion of equivalence mathematically precise, we introduce the notion of swap. A swap allows us to swap the order of two systems in parallel composition.

The properties just seen make a process theory a strict symmetric monoidal category \[44\]. In this setting, we can view states as a particular kind of a process: a process with no input, corresponding to the preparation of a system. Mathematically, they are processes with the trivial system as input; e.g. \(s : I \rightarrow A\) will be a state of \(A\). Therefore, in general, to understand what the states of a process theory are, we need to identify what the trivial system is, i.e. the identity of system composition.

3. Resource theories

From a formal point of view, a resource theory arises from a process theory when there is a restriction on the processes. In this setting, the conversion problem becomes relevant, which means determining whether, given two states, we can convert one into the other. The answer can be given with the aid of particular functions of the two states, called conversion witnesses. Finally, we present the example of the resource theory of asymmetry, which provides the mathematical inspiration for the approach in this article.

A resource theory is a process theory augmented by partitioning processes into “free” and “non-free” processes such that composing free processes and swapping systems are free as well as the identity process. These requirements make a resource theory equivalent to specifying a strict symmetric monoidal subcategory of the process theory. The morphisms in the subcategory are exactly the free operations of the resource theory. In particular, the states in this subcategory are the free states of the resource theory.

We can set up a hierarchy among the states of a theory, whereby a state is more valuable than another if, from the former, it is possible to reach a larger set of states using only free processes \[15, 19\]. More formally, we can set up a partial preorder, known as the resource preorder, defined as \(s' \preceq s\), if there exists a free process \(f\) such that \(s' = fs\), where the product denotes sequential composition of the processes \(f\) and \(s\). In this setting, we want to map the set of states with its partial preorder into a better known or simpler partially ordered set \((X, \leq)\) \[21\]. To this end, we need to introduce a function \(M\) from the set of states to a partially ordered set \(X\) such that \(M(s') \leq M(s)\) whenever \(s' \preceq s\). Such a function is called a monotone. Normally, the condition \(M(s') \leq M(s)\) is only a necessary condition for the conversion of \(s\) into \(s'\). However, in some special cases, combining possibly more than one monotone (even involving different partially ordered sets, as we show in §IV.B even involving the divisibility partial order, denoted by \(|\), can give us necessary and sufficient conditions for the convertibility under free operations. In this case, we can solve the conversion problem completely, and such a set of monotones is called a complete set.

A particularly successful example of a resource theory in quantum information is the resource theory of of asymmetry \[33–35\]. A symmetry of a system corresponds to it being unchanged under certain actions, such as time reversal, spacial reflections. A quantum system is associated with a Hilbert space, and its symmetries are represented by unitary or anti-unitary operators on it. In the resource theory of asymmetry, for every quantum system, there is a group representing its symmetry, whose action is given by unitary operators up to a global phase. Quantum states are described by trace-class positive semi-definite operators with unit trace. Quantum processes are completely positive and trace non-increasing linear maps on the set of self-adjoint operators on a Hilbert space. As seen above, free operations arise from some restriction. In this case it is a covariance condition: \(\mathcal{E}\) is free if \(U_\theta \circ \mathcal{E} = \mathcal{E} \circ U_\theta\), for every unitary operation \(U_\theta\) associated with the symmetry (in the context of group representation theory, such maps are known as equivariant maps). Thus, acting with the symmetry before and after the action of \(\mathcal{E}\) is the same. This fact can be illustrated with a commutative diagram.

\[
\begin{align*}
\mathcal{B}(\mathcal{H}) \xrightarrow{\varepsilon} \mathcal{B}(\mathcal{H}) \\
\downarrow_{U_\theta} & \downarrow_{U_\theta} \\
\mathcal{B}(\mathcal{H}) \xrightarrow{\varepsilon} \mathcal{B}(\mathcal{H})
\end{align*}\tag{10}
\]
The resource theory of asymmetry was used to give an informational reconstruction of Noether’s theorem [36].

III. APPROACH

In this section, we discuss how we develop a mathematical theory of evolutionary, which is based on the model of discrete dynamical systems as such systems represent all systems evolving in discrete systems. We begin by showing that discrete dynamical systems are special cases of process theories. Then we show how the covariance condition, discussed in §II C3, is applied to discrete dynamical systems, and it gives rise to a full-fledged resource theory. Finally, we explain how we treat influences, usually called perturbations, on evolving systems.

A. Formulating discrete dynamical systems

Now we present two new labels to introduce on states of a discrete dynamical systems (or, equivalently, on vertices of the dynamical graph) that are useful for stating our results. We also provide a formulation of a discrete dynamical system as a process theory both in the absence and in the presence of randomness.

1. New labels

Here we introduce two new labels for states of a dynamical system. The first, called transient progeny, looks at their evolution in the future, specifying how fast they reach an attractor, viz. dynamics that occurs in the long run. The second, called ancestry, looks at their evolution from the past, identifying how many steps were necessary to evolve to that state.

It is possible to quantify how transient a state $s$ is by considering its transient progeny (or progeny for short) the minimum natural number $d$ such that $\phi^d(s)$ is in a cycle. This also quantifies the number of successors of $s$ that are still transient, including $s$ itself. The bigger the progeny, the more transient the state is. This is the length of the shortest path from the vertex associated with that state to a vertex in a cycle in the dynamical graph. From this point of view, the states of a cycle can be viewed as transient states of progeny 0 because 0 time steps are necessary to evolve them to a state in a cycle.

We can also consider a “backwards distance” of a state from its farthest predecessor in the dynamical graph. More formally, we define the ancestry $a$ of a vertex as the maximum length of a path leading to that vertex. Equivalently, the ancestry of a state $s$, denoted as $a(s)$, is the maximum natural number for which there exist a state $t$ such that $\phi^{a(s)}(t) = s$. For states $s$ in an attractor of length $\ell$, since we have $\phi^\ell(s) = s$ for all $n \in \mathbb{N}$, the ancestry is formally infinite, so we set $a(s) = +\infty$. Instead for transient states, their ancestry will be a finite natural number.

2. Non-random case

In §II A, we explained that, in a discrete dynamical system $(S, \phi)$ with no randomness present at any step of its evolution, evolution is described by powers of $\phi$. It is natural to describe this situation with a process theory based on sets and functions [42].

As we are dealing with discrete dynamical systems, we construct a process theory where systems are in fact dynamical systems, i.e. pairs $(S, \phi)$, where $S$ is a set and $\phi : S \to S$ is the generator of the dynamics on $s$. Processes between dynamical systems $(S_1, \phi_1)$ and $(S_2, \phi_2)$ are arbitrary functions between $S_1$ and $S_2$. In particular, when $S_1 = S_2$, we have what we called influences.

The composition of dynamical systems is by Cartesian product: if $(S_1, \phi_1)$ and $(S_2, \phi_2)$ are two dynamical systems, their composition is $(S_1 \times S_2, \phi_1 \times \phi_2)$, where

$$\phi_1 \times \phi_2 : S_1 \times S_2 \to S_1 \times S_2 : (s_1, s_2) \mapsto (\phi_1(s_1), \phi_2(s_2))$$

for every $s_1 \in S_1$, and every $s_2 \in S_2$. The sequential composition of processes is simply the composition of functions. Instead, the parallel composition of processes

$$f : S_1 \to S_2, \ f' : S_1' \to S_2'$$

is the function

$$f \times f' : S_1 \times S_1' \to S_2 \times S_2'$$

such that

$$(s_1, s_1') \mapsto (f(s_1), f'(s_1'))$$

for every $s_1 \in S_1$ and $s_1' \in S_1'$.

In this setting, it is clear that the trivial dynamical system $(I, \phi)$, which does not affect other systems when composed with them, is obtained by taking $I$ to be a singleton and $\phi$ to be the identity $I$. With this in mind, the process-theoretic states of $(S, \phi)$, viz. processes from $I$ to $S$ can be identified with the elements of $S$ [41, 44], i.e. with what we defined as the states of the dynamical system in §II A. In the following, these states are called “deterministic states” to differentiate them from the states permitting randomness.

3. Random case

In §II A 3, we explained that in a discrete dynamical system $(S, \phi)$ in the presence of randomness, we represent randomness with probability vectors and evolution is described by powers of a linearization of $\phi$. This conceptually distinct situation, can be described
naturally based on the process theory of stochastic maps [42].

The potential presence of randomness in a discrete dynamical system \((S, \phi)\) at any time step forces us to use a different process theory from the case of deterministic initial conditions. Now, if \(S\) has \(M\) elements, systems are pairs \((S, \phi)\), where \(S\) is the simplex of probability vectors in \(\mathbb{R}^M\), namely probability distributions over \(s\). Processes between systems \((S_1, \phi_1)\) and \((S_2, \phi_2)\) are stochastic maps between \(S_1\) and \(S_2\).

**Definition 1.** We call *stochastic influence* any stochastic map \(f : S \rightarrow S\).

Now a stochastic map

\[
f : S_1 \rightarrow S_2
\]

(15)
can be represented with a stochastic matrix once we fix the canonical bases in the domain \(S_1\) and the codomain \(S_2\) of \(f\):

\[
F = \begin{pmatrix}
p(1|1) & \cdots & p(1|M_1) \\
\vdots & \ddots & \vdots \\
p(M_2|1) & \cdots & p(M_2|M_1)
\end{pmatrix}.
\]

(16)
The \((i, j)\) entry is the probability \(p(i|j)\) of jumping to the state of underlying set \(S_2\) labeled by \(i\) from the state of the underlying set \(S_1\) labeled by \(j\). This probability is 1 if and only if there is a deterministic transition. The \(j\)th column of \(f\) is the evolution of the deterministic state labeled by \(j\) under the covariant stochastic map \(f\). If the influence is deterministic, as defined in §II A, we refer this influence now as 'deterministic influence'. Note that a deterministic influence only has zeros and ones as entries in its representation as a stochastic matrix.

The composition of a system \((S_1, \phi_1)\) with a system \((S_2, \phi_2)\) is the system \((S_1 \otimes S_2, \phi_1 \otimes \phi_2)\), where \(S_1 \otimes S_2\) is the simplex of probability vectors in the tensor product of the real vector spaces associated with \(S_1\) and \(S_2\). The sequential composition of processes is simply their composition as linear functions. Instead, the parallel composition of processes

\[
f : S_1 \rightarrow S_2, f' : S_1' \rightarrow S_2'
\]

(17)
is the tensor-product function

\[
f \otimes f' : S_1 \otimes S_1' \rightarrow S_2 \otimes S_2'.
\]

(18)

In every process theory, there is a trivial system corresponding to the system that does not affect other in composition.

In the presence of randomness, this system is \((\mathcal{S}_0, I)\), where \(\mathcal{S}_0\) is the simplex of probability vectors in \(\mathbb{R}\), i.e. the set containing only the number 1. With this in mind, the process-theoretic states of a generic dynamical system in the presence of randomness \((\mathcal{S}, \phi)\) can be identified with the elements of the simplex \(\mathcal{S}\) itself, i.e. with probability vectors. In the following, these states are called "stochastic states".

**B. Imposing covariance**

In this subsection we introduce the mathematical tools to describe the influence of the external environment on a discrete dynamical system. We focus on *covariant* influences, which are external influences that somehow respect the evolution of a dynamical system. We also show that restricting to covariant influences gives rise to a well-defined resource theory.

1. **Introducing covariance**

Now we give the formal definition of covariant maps in the process theory of discrete dynamical systems, both in the presence of deterministic and random initial conditions. We also present their foremost mathematical properties.

The definition of covariant maps can be given in a unified way in both process theories we defined above, namely for deterministic and random initial conditions. Recall that, in the case of no randomness allowed, systems are pairs \((S, \phi)\), where \(s\) is a set, and \(\phi\) is the generator of the dynamics. On the other hand, in the case of allowed randomness, systems are pairs \((\mathcal{S}, \phi)\), where \(\mathcal{S}\) is the set of probability vectors, and now \(\phi\) is the linear version of the generator of the dynamics, as described in §II A. We cover both situations by representing systems as \((A, \phi)\), where \(a\) is either \(S\) or \(\mathcal{S}\).

**Definition 2.** Let \((A_1, \phi_1)\) and \((A_2, \phi_2)\) be dynamical systems. A map \(f : A_1 \rightarrow A_2\) is called *covariant* if \(f \circ \phi_1 = \phi_2 \circ f\).

A covariant map in the absence of randomness is a function between sets, and, in the presence of randomness, is a stochastic map. In particular, in the case of a stochastic influence \(f\), the associated matrix \(F\) is square and the covariance condition takes a simple matrix form, namely \([F, \Phi] = 0\) for stochastic matrix \(\Phi\) associated with \(\phi\), the generator of dynamics.

As \(\Phi\) describes a deterministic dynamic, there will be only one non-zero entry per column, and such an entry will be equal to 1. The commutation condition \([F, \Phi] = 0\), signaling a covariant influence, will add further constraints to the entries of \(F\). Such constraints in general depend on \(\Phi\).

Note that, in particular, all powers of \(\phi\), namely,

\[
\phi^n : A \rightarrow A, n \in \mathbb{N},
\]

(19)
are covariant, as

\[
\phi^n \circ \phi = \phi \circ \phi^n = \phi^{n+1}.
\]

(20)

**Proposition 1.** Let \((A_1, \phi_1)\) and \((A_2, \phi_2)\) be dynamical systems. A map \(f : A_1 \rightarrow A_2\) is covariant if and only if (iff)

\[
f \circ \phi_1^n = \phi_2^n \circ f, n \in \mathbb{N}
\]

(21)
Proof. Sufficiency is straightforward: it is enough to take \( n = 1 \) to have the thesis. To prove necessity, we proceed by induction. If \( n = 0 \) or \( n = 1 \) the property is obviously true. Let us assume it is true for a generic \( n \), and let us prove is holds also for \( n + 1 \). We have
\[
f \circ \phi_1^n = f \circ \phi_1^n = \phi_2 \circ f \circ \phi_1^n
\]
where, in the second equality, we used the covariance of \( f \), and in the third equality the induction hypothesis. This proves the proposition.

2. Covariance yields a resource theory

Now we show that restricting ourselves to covariant maps yields strict symmetric monoidal subcategories of the process theories we introduced in to describe discrete dynamical systems. Thus, we think of covariant maps as free operations of a new resource theory.

**Proposition 2.** Covariant maps form a strict symmetric monoidal subcategory.

**Proof.** The identity, being part of the dynamics, is obviously covariant. The sequential composition of two covariant maps is still covariant. Indeed, if \( (A_1, \phi_1), (A_2, \phi_2), \) and \( (A_3, \phi_3) \) are dynamical systems, and if
\[
f : A_1 \to A_2, \ g : A_2 \to A_3
\]
are both covariant, then
\[
g \circ f : A_1 \to A_3
\]
satisfies
\[
g \circ f \circ \phi_1 = g \circ \phi_2 \circ f = \phi_3 \circ g \circ f,
\]
where we use the covariance of \( f \) in the first equality and the covariance of \( g \) in the second equality. Similarly, the parallel composition of covariant maps is still a covariant map. This is because the generator of dynamics on the composition of two discrete dynamical systems is given by the parallel composition of the generators of the respective dynamics. Hence, if \( (A_1, \phi_1), (A_2, \phi_2), (A_3, \phi_3) \), and \( (A_4, \phi_4) \) are dynamical systems, and if
\[
f : A_1 \to A_2, \ g : A_3 \to A_4,
\]
then
\[
(f \otimes g) \circ (\phi_1 \otimes \phi_3) = (f \circ \phi_1) \otimes (g \circ \phi_3)
\]
\[
= (\phi_2 \circ f) \otimes (\phi_1 \circ g) = (\phi_2 \otimes \phi_4) \circ (f \otimes g).
\]
The swap between two discrete dynamical systems \( (A_1, \phi_1), (A_2, \phi_2) \) is a covariant operation:
\[
\text{SWAP} \circ (\phi_1 \otimes \phi_2) = (\phi_2 \otimes \phi_1) \circ \text{SWAP}.
\]
as it immediately follows from the property of the \text{SWAP} itself.

This proposition shows that the definition of covariant operations is well posed, and it gives rise to a full-fledged resource theory. Therefore, we can rightfully consider them as free operations.

3. Covariant influence

Now we explain the notion of influence on a dynamical system, introduced in §II A and Definition 1, as the effective representation of the action of something external to the dynamical system itself. External influence happens in several situations, given the fact that truly isolated systems do not really exist. We also motivate the introduction of covariant influence, which arises due to the dynamical time scale of evolution being shorter than the time scale for changes due to influences external to the dynamical system.

Discrete dynamical systems are often studied as isolated systems, where one is interested in how the system evolves subject to its own dynamics. However, in Nature truly isolated system do not exist, and there is always something surrounding a dynamical system—an environment—that influences it. For this reason, despite the fact that a dynamical system and its environment can be jointly treated as a larger, isolated, dynamical system, it is still interesting to focus purely on the dynamical system and describe the presence of the environment through some effective action that disrupts the evolution of the dynamical system. This action is a function that maps states of the dynamical system, whether deterministic or stochastic, to other states. As such, we model this action by incorporating influence, whether deterministic or stochastic.

Among all possible influences, we focus on covariant ones, namely those obeying Definition 2. These influences, despite changing the evolution of the system, work in such a way that it is compatible with the dynamics. Indeed, the covariance condition can be stated as follows: if we evolve the dynamical system for \( n \) time steps and then we disturb it, it is the same as first disturbing the system with the same influence and then letting the system evolve for \( n \) time steps. The idea behind this is that a system has time to adapt to the external influence, perhaps because such an influence occurs on a longer time scale. In particular, as we show in §IV, a covariant influence cannot undo the evolution of state of a dynamical system under its dynamic.

**IV. RESULTS**

In this section we present the results of this article. First we identify free states in the resource theory of covariant influences, which allows to identify two different types of resources into play whether randomness is forbidden or not. In the former case, we have the resource theory of evolutionarity, and in the latter case,
the resource theory is of non-attractorlessness. Finally, we turn to the conversion problem, i.e. determining when one state can be converted into another by means of a covariant influence.

A. Free states

In this subsection we partition the states of a dynamical system into free and non-free ones according to the resource theory of covariant influences. Free states are very different whether randomness is forbidden or not. For this reason, we treat these two situations separately.

1. Deterministic case

Now we examine the situation in which randomness is completely forbidden, both in the initial conditions (i.e., the deterministic state) and in the changes induced by a covariant influence. This corresponds to a dynamical system whose complexity is low enough to allow us to use a fully deterministic description.

In the context of resource theories, free deterministic states of a system \( s \) are covariant maps \( f : I \to S \), where \( I \) is a singleton. We show that deterministic free states can be identified with fixed points of a discrete dynamical system (3).

**Proposition 3.** There is a bijection between free deterministic states of a discrete dynamical system \( (S, \phi) \) and its fixed points.

**Proof.** Consider a free deterministic state \( f : I \to S \); this sends the only element of \( I \), which we call 1 by convention, to a state \( s \) of \( S \). The dynamic on the trivial dynamical system \( I \) is given by the identity; therefore, the condition for covariance becomes \( f = \phi \circ f \). Applying \( f \) to 1, and recalling \( f(1) = s \), we have \( s = \phi(s) \). This means that \( s \) is a fixed point.

Conversely, given a fixed point \( s \) of the evolution of \( (S, \phi) \), i.e. \( \phi(s) = s \), we can construct the unique covariant function \( f : I \to S \) such that \( f(1) = s \). The covariance of \( f \) is guaranteed by the fact that \( s \) is a fixed point.

Not all discrete dynamical systems have fixed points; this depends on the dynamics. By Proposition 3, those systems with no fixed points, in the lack of randomness, have no free states.

2. Stochastic case

Now we analyze the situation of a dynamical system where randomness is permitted. Thus, we can have stochastic states at any time step and also the influence can be stochastic. This situation models dynamical systems where their behavior is too complicated to be described feasibly in purely deterministic terms, due to the presence of an influence that acts on the dynamical system as a stochastic process. We see that, for stochastic dynamical systems, the landscape of free states becomes richer.

As seen in §II A, states of a stochastic dynamical system can be represented as probability vectors \( p \). In this setting, free stochastic states are defined as maps

\[
\phi : \mathcal{S}_0 \to \mathcal{S},
\]

obeying the covariance condition of Def. 2, where now \( f(1) = p \) is still equivalent to \( \phi(p) = p \). By a similar argument to that for deterministic free states, the covariance condition on stochastic maps entails the probability vector must be invariant under the generator of the dynamics.

Below we prove equivalence between the existence of a bijection between free stochastic states and probability distributions over the attractors of the discrete dynamical system with the requirements that every deterministic state in an attractor sharing equal probability with all other states in the same attractor. Before formalizing this proposition, we introduce pertinent notation, building on the ability to partition deterministic states into basins of attraction. For \( b_i \) referring to block corresponding to the \( i \)-th basin of attraction, we write

\[
p = b_1 \oplus \cdots \oplus b_k.
\]

with \( k \) basins of attraction and \( p \) the probability vector.

Let us focus on the first block as, for other blocks, the analysis is completely analogous. Note that, for transient deterministic states \( s \), by taking large enough \( m \), we can always ensure that \( \phi^m(s) \) is in an attractor. For example, it is enough to take \( M \) as the maximum of the transient progenies of transient states. In \( b_1 \) there are \( \ell_1 \) states corresponding to the deterministic states in the attractor (of length \( \ell_1 \)) and other \( m_1 \) transient deterministic states, i.e. states that are not in the attractor, but eventually goes into it.

**Definition 3.** A stochastic state is called uniform if two conditions hold:

1. all probabilities associated with transient deterministic states are 0;
2. deterministic states in the same attractor have the same probability.

In simpler terms, every attractor has a uniform distribution over its deterministic states (with uniform probabilities \( 1/\ell \), \( \ell \) being the length of the attractor), weighted by a probability distribution over the different attractors.

Below and henceforth we consider integers modulo some number, say \( m \). This means picking a convenient representative for each remainder class modulo \( M \) so
that expressions are well defined. Here we consider two representatives for each remainder class. The one we pick depends on the concrete setting where we consider integers modulo 1: in some cases it may mean picking a representative from 1 to m, in others from 0 to m − 1.

**Proposition 4.** There is a bijection between free stochastic states of a discrete dynamical system and uniform probability vectors.

**Proof.** We start by proving that if a probability vector p, representing a stochastic state, is free, that is, invariant under ϕ, then it must be uniform. Following Eq. (30), we arrange the deterministic states s_i corresponding to the first block b_1 so that the first ℓ_1 entries of b_1 are the deterministic states of its attractor, in such a way that ϕ(s_i) = s_{i+1}, where i + 1 is intended to be modulo ℓ_1: thus, we pick a representative in the remainder classes modulo ℓ_1 between 1 and ℓ_1. By Prop. 1, the request ϕ(p) = p is equivalent to ϕ^n(p) = p for every n ∈ N. In particular, we require covariance for n = m, where M is the minimum exponent for which ϕ^m(s) is in the attractor, for every deterministic state s in the first basin of attraction. We obtain

\[
b_1 = \begin{pmatrix}
p_1 \\
\vdots \\
\frac{p_{i_1}}{p_{i_1+m_1}}
\end{pmatrix} = \begin{pmatrix}
p_1 \cdots \\
\vdots \\
0
\end{pmatrix} = \phi^m(b_1), \quad (31)
\]

where p_i is either p, or p, plus a probability contribution coming from the transient states, and the subscripts of the entries of the vector in right-hand side are intended to be modulo ℓ_1, in the same way as discussed above. Equation (31) immediately tells us that the probabilities associated with transient states must all vanish. Now, if we require the invariance under ϕ of such a b_1 with 0 probabilities for transient states, we obtain

\[
b_1 = \begin{pmatrix}
p_1 \\
\vdots \\
\frac{p_{i_1}}{p_{i_1+m_1}}
\end{pmatrix} = \begin{pmatrix}
p_1 \cdots \\
\vdots \\
0
\end{pmatrix} = \phi(b_1). \quad (32)
\]

This yields the condition p_1 = ··· = p_{i_1}.

To conclude the proof, we repeat the same argument for each of the k blocks. Therefore, the probabilities associated with the deterministic states of an attractor must be equal for the states in that attractor (but possibly different between different attractors).

Conversely, by a similar argument as above, one can easily check that uniform states are invariant under ϕ.

The situation of a discrete dynamical system differs a lot whether randomness is allowed or not. Indeed, when randomness is forbidden, free states are only fixed points, whereas when randomness is permitted, free states are associated with attractors of any length. In particular, it means that free stochastic states always exist because any discrete dynamical system has attractors.

**B. The conversion problem: evolutionarity**

Now we analyze the conversion problem for discrete dynamical systems where randomness is forbidden. In this case, we call the resource “evolutionarity” as all valuable states (i.e. non-free) evolve in time because they are not fixed points. We are able to give a necessary and sufficient condition for convertibility under covariant operations, and we show that covariant influences shorten the length of attractors.

Denoting the length of a deterministic state s by ℓ, its transient progeny by d, and its ancestry by a(s), we have the complete characterization of the state transitions dividing it into different pieces.

Theorem 5. Let s and s’ be two deterministic states of a discrete dynamical system (S, ϕ). Then there exists a covariant influence converting s into s’ iff

\[
d' \leq d, \quad (33)
\]

\[
\ell' | \ell, \quad (34)
\]

\[
a(\phi^n(s')) \geq a(\phi^n(s)) \quad (35)
\]

for n = 0, . . . , d′ − 1.

In the following, we provide the proof of this result, dividing it into different pieces.

Theorem 5 gives necessary and sufficient conditions for the conversion of deterministic states of a discrete dynamical system under covariant influences. We can phrase these conditions in the language of monotones (see §II C). Indeed, we can view the transient progeny d of a deterministic state as a function from the set of deterministic states to the set N of natural numbers, seen as the (totally) ordered set (N, ≤). The length, instead, can be obtained from a function from the set of deterministic states and N, this time seen as the partially ordered set (N, |), where | denotes the divisibility partial order. Finally, the ancestry, is a function from the set of states to N ∪ {+∞}, totally ordered by the ≥ order.

According to Theorem 5, these functions are monotones in our resource theory. However, not only can they be used to formulate necessary conditions for the conversion of states, but also sufficient (33)–(35). Thus, they form a complete set of monotones.
1. Necessary conditions

As a first step to solve the conversion problem, we examine necessary conditions for the conversion of one deterministic state into another. These conditions involve labels of deterministic states: the lengths, the transient progenies and the ancestries of the states involved.

We saw that, when randomness is forbidden, there is a bijection between free states and fixed points of the discrete dynamical system. Thus, all non-free states, i.e. valuable states, are states that evolve in time. This makes us identify the relevant resource as “evolutionarity”, which is the property that a deterministic state evolves.

Let us consider two deterministic states $s$ and $s'$ of the same discrete dynamical system $(S, \phi)$ (the case of states of two distinct dynamical systems is treated in Appendix A). We want to determine when covariant influences can transform $s$ into $s'$. Here we formulate necessary conditions. The first involves the transient progeny and the length $\ell$ of a state.

**Lemma 6.** In a discrete dynamical system without randomness, states of transient progeny $d$ and of length $\ell$ are mapped to states of transient progeny $d' \leq d$ and of length $\ell'$, such that $\ell' \mid \ell$.

**Proof.** Let $s$ be a deterministic state of transient progeny $\ell$ and of length $\ell$ in a discrete dynamical system $(S, \phi)$. Then

$$\phi^{\ell+d}(s) = \phi^{\ell}(\phi^d(s)) = \phi^d(s). \quad (36)$$

Let $f : S \to S$ be covariant; then

$$\phi^d(f(s)) = f(\phi^d(s)) = f(\phi^{\ell+d}(s)) = \phi^{\ell+d}(f(s)). \quad (37)$$

This means that, after $\ell$ iterations, we come back to $\phi^d(f(s))$, which implies that $\phi^d(f(s))$ lies in an attractor of length $\ell'$ that divides $\ell$. This also means that $f(s)$ will be a state of transient progeny $d' \leq d$.

Thus, transient states become “less transient” (in particular they become attractor states if $d' = 0$).

**Remark 1.** If $d = 0$, the state $s$ is an attractor state, and Lemma 6 tells us how any state evolves, including attractor states, under covariant operations. Specifically, the constraint $0 \leq d' \leq d = 0$, implies $d' = 0$. If $s$ is an attractor state, then $s'$ must be an attractor state too, whose length divides the length of the original attractor. In other words, attractor states are sent to attractor states.

We can also say something about the ancestry of states: it increases under free operations. Clearly, this statement is meaningful when at least one of the states is transient.

**Lemma 7.** For every state $s$ of a deterministic dynamical system $(S, \phi)$, the ancestry of the successors of $s$ does not decrease under free operations. Formally, if $f : S \to S$ is covariant, then

$$a(\phi^n(s)) \leq a(\phi^n \circ f(s)), \quad (38)$$

for every $n \in \mathbb{N}$.

**Proof.** If $s$ is an attractor state, all its successors are attractor states too. This means, that for every $n \in \mathbb{N}$, $a(\phi^n(s)) = +\infty$. By Lemma 6, $f(s)$ is an attractor state, so, again its successors are attractor states. Therefore, $a(\phi^n \circ f(s)) = +\infty$ for every $n \in \mathbb{N}$, satisfying the condition $a(\phi^n(s)) \leq a(\phi^n \circ f(s))$ for every natural number $n \in \mathbb{N}$.

Suppose now $s$ is a transient state, and let $d > 0$ be its transient progeny. This means that, after $d$ steps, the successors of $s$ become attractor states. By Lemma 6, $f(s)$ is a state with transient progeny $d' \leq d$, which means that, after $d' \leq d$ steps the successors of $f(s)$ become attractor states. Therefore, for $n \geq d'$, we have $a(\phi^n \circ f(s)) = +\infty$, which clearly satisfies inequality (38). Instead, for every $n = 0, \ldots, d' - 1$, there exists a state $s_n^*$ such that $\phi^n(s) = \phi^n(\phi^n(s)) (s_n^*)$. By covariance of $f$,

$$\phi^n \circ f(s) = f(\phi^n(s)) = f \left( \phi^n(\phi^n(s)) (s_n^*) \right) = \phi^n(\phi^n(s)) (f(s_n^*)). \quad (39)$$

This shows that $f(s_n^*)$ is a predecessor of $\phi^n \circ f(s)$, and that there are $a(\phi^n(s))$ steps from $f(s_n^*)$ to $\phi^n \circ f(s)$. Thus, Eq. (38) holds even for $n = 0, \ldots, d' - 1$. In conclusion, we proved that, for any state $s$, Eq. (38) holds for every $n \in \mathbb{N}$.

In words, a (transient) state becomes farther and farther from its farthest predecessor. In particular, in the two lemmas above we can take $f$ to be the generator $\phi$ of the dynamics. When $f = \phi$, $d' \leq d$, $\ell = \ell'$, and the ancestry never decreases for all the successors of a state. Transitions from one state to another due to dynamics occur only inside a given basin of attraction. However, the statements of the two lemmas suggest that, when $f$ is not the dynamic, it may be possible to jump between different basins of attraction, provided the constraints are met, e.g. from an attractor of length 4 to an attractor of length 2. However, at the moment we do not know if this is the case because the two lemmas are only necessary conditions, but a priori not sufficient.

2. Sufficient conditions

The conditions expressed in Lemmas 6 and 7 are necessary, and they tell us how deterministic states change under covariant operations. Now we prove that these conditions are sufficient too: if they are satisfied,
we can conclude that a covariant map exists that sends a deterministic state \( s \) to \( s' \).

**Lemma 8.** Let \( s \) and \( s' \) be two deterministic states of discrete dynamical system \((S, \phi)\), with \( s' \) an attractor state, and \( \ell, \ell' \) their respective lengths. If \( \ell' \parallel \ell \), there exists a covariant operation converting \( s \) into \( s' \).

**Proof.** To prove sufficiency, we must construct a covariant function mapping \( s \) into \( s' \). To this end, we need to specify its action on all deterministic states of \((S, \phi)\). Now, for the states \( \{x\} \) not in the basin of attraction of \( s \), construct \( f(x) := x \). Such an \( f \) is covariant on these states by construction. Now we need to map the basin of attraction of \( s \) to the basin of attraction of \( s' \) in a covariant way. The starting point is to fix \( f(s) = s' \), and then we reconstruct the action of \( f \) on the other states of the basin of attraction of \( s \).

Observe that \( s \) has \( d + \ell \) successors \( s_n := \phi^n(s) \) (including itself \( s_0 = s \)), for \( n = 0, \ldots, d + \ell - 1 \), where \( d \) is the transient progeny of \( s \). For these states \( \{s_n\} \), construct \( f(s_n) := \phi^n(s') \), which yields in particular \( f(s) = s' \). Now, for \( n = 0, \ldots, d + \ell - 2 \), one has

\[
\begin{align*}
\phi \circ f(s_n) &= \phi \circ \phi^n(s') = \phi^{n+1}(s') \\
&= f(s_{n+1}) = f\left(\phi^{n+1}(s)\right) = f \circ \phi(s_n), \quad (40)
\end{align*}
\]

and, for \( n = d + \ell - 1 \),

\[
\begin{align*}
\phi \circ f(s_{d+\ell-1}) &= \phi \circ \phi^{d+\ell-1}(s') = \phi^{d+\ell}(s') = \phi^d(s')
\end{align*}
\]

(41)

because \( s' \) is in an attractor of length \( \ell' \parallel \ell \). On the other hand,

\[
\begin{align*}
f \circ \phi(s_{d+\ell-1}) &= f \circ \phi \circ \phi^{d+\ell-1}(s) \\
&= f \circ \phi^{d+\ell}(s) = f \circ \phi^d(s) = f(s_d) = \phi^d(s'),
\end{align*}
\]

(42)

where we have used the fact that \( \phi(s) = s' \) is a state of transient progeny \( d \) and length \( \ell \). We conclude that \( \phi \circ f(s_{d+\ell-1}) = f \circ \phi(s_{d+\ell-1}) \). In conclusion, \( f \) is covariant on every successor of \( s \).

Now we need to define \( f \) on all the other states in the same basin of attraction as \( s \). Note that every state in the same basin of attraction as \( s \) is a predecessor of \( s_d := \phi^d(s) \) because \( s_d \) is in the attractor of length \( \ell \). Let \( \mathcal{S} \) be the set of successors of \( s \) (cf. also Fig. 4).

\[
\mathcal{S} = \{s_n := \phi^n(s); n = 0, \ldots, d + \ell - 1\}. \quad (43)
\]

For the states \( t \) in the complement of \( \mathcal{S} \), if it is non-empty, define \( \delta(t; s_d) \) to be the minimum number of steps to reach \( s_d \) from \( t \). Note that these states are all transient because all the attractor states are among the successors of \( s \). On the states in the (non-empty) complement of \( \mathcal{S} \), set \( f \) to be \( f(t) := \phi^{\delta(t; s_d)}(s') \), where the exponent is regarded to be modulo \( \ell' \), where this time we pick a representative in the remainder classes modulo \( \ell' \) that is between 0 and \( \ell' - 1 \).

Let us check that this definition makes \( f \) covariant. We have \( \phi \circ f(t) = \phi^{\delta(t; s_d)\ell + 1}(s') \). Instead, to assess \( f \circ \phi(t) \), we must distinguish two cases. If \( \phi(t) \notin \mathcal{S} \), we have

\[
\begin{align*}
f \circ \phi(t) &= \phi^{\delta(t; s_d)}(s') = \phi^{\delta(t; s_d)\ell + 1}(s'), \quad (44)
\end{align*}
\]

where we have used the fact that \( \delta(t; s_d) = \delta(t; s_d) - 1 \) for transient states \( t \). Indeed, if \( \phi(t) = s_d \), then \( \phi^{k-1}(\phi(t)) = s \) (and for transient states \( k > 0 \)). Instead, if \( \phi(t) \in \mathcal{S} \), then \( \phi(t) = \phi^n(s) \), for some \( n \in \{0, \ldots, d + \ell - 1\} \). Therefore, by the above definition of \( f \), \( f \circ \phi(t) = \phi^n(s') \). Now, if \( n \leq d \), we have \( \delta(t; s_d) = d - n \), whence \( \delta(t; s) = d - n + 1 \). In this case

\[
\begin{align*}
f \circ \phi(t) &= \phi^n(s') = \phi^{\delta(\ell; s_d)\ell + 1}(s'). \quad (45)
\end{align*}
\]

Since \( s' \) is in an attractor of period \( \ell' \), the exponent must be regarded as modulo \( \ell \), as discussed above. On the other hand, if \( n \in \{d + 1, \ldots, d + \ell - 1\} \), we have \( \delta(\ell; s_d) = d - n \) (modulo \( \ell \), where for every remainder class modulo \( \ell \) we pick a representative between 0 and \( \ell \), whence \( \delta(t; s_d) = d - n + 1 \) (modulo \( \ell \), as before). Again,

\[
\begin{align*}
f \circ \phi(t) &= \phi^n(s') = \phi^{\delta(\ell; s_d)\ell + 1}(s'), \quad (46)
\end{align*}
\]

where the expression at exponent is modulo \( \ell \), as above. However, since \( s' \) is in a cycle of length \( \ell' \), the exponent is actually modulo \( \ell' \), in the sense above, because of the length \( \ell' \) of the cycle. Given \( \ell' \parallel \ell \), whether we take the exponent modulo \( \ell \) or \( \ell' \), the state \( \phi^{\delta(\ell; s_d)\ell + 1}(s') \) is the same. In conclusion, \( f \) is covariant on all states in the complement of \( \mathcal{S} \) too. This proves that there is a covariant influence \( f \) mapping \( s \) to \( s' \).

With reference to Lemma 8, if \( d \) and \( d' \) denote the transient progeny of \( s \) and \( s' \) respectively, we have \( d' = 0 \) because \( s' \) is an attractor state, and \( a(\phi^n(s')) = +\infty \) for every \( n \in \mathbb{N} \), because all successors of \( s' \) are attractor states too. Hence, note that this situation can be viewed as a special case of the more general conditions \( d' \leq d \) and \( a(\phi^n(s')) \geq a(\phi^n(s)) \) for \( n = 0, \ldots, d' - 1 \), which are examined for the case of transient states in the following lemma.

**Example 1.** The example of how to construct a free operation mapping an attractor of length 4 into an attractor of length 2 is shown in Fig. 3. ■

**Lemma 9.** Let \( s \) and \( s' \) be two deterministic states of a deterministic dynamical systems \((S, \phi)\), with \( s' \) transient. If \( d' \leq d \), \( \ell' \parallel \ell \) and \( a(\phi^n(s')) \geq a(\phi^n(s)) \) for \( n = 0, \ldots, d' - 1 \), then there exists a covariant influence converting \( s \) into \( s' \).

**Proof.** Again, we have to construct a covariant function mapping \( s \) into \( s' \), specifying its action on all deterministic states of \((S, \phi)\). As in the proof of Lemma 8, for the states \( s \) not in the basin of attraction of \( s \), define
f(x) := x. To construct the action of f on the states in the same basin of attraction as s, we proceed as in the proof of Lemma 8.

The conditions \( d' \leq d \) and \( a(\phi^n(s')) \geq a(\phi^n(s)) \) for \( n = 0, \ldots, d' - 1 \) imply that s is transient too. The situation is now more complicated than in the proof of Lemma 8, due to the presence of states in the basin of attraction of s that are neither successors nor predecessors of s.

Consider the set \( S \) of the \( d + \ell \) successors of s as in Eq. (43). As done in the proof of Lemma 8, construct \( f(s_n) := \phi^n(s') \), which yields \( f(s) = s' \). Now, for \( n = 0, \ldots, d + \ell - 2 \), one immediately proves the covariance of \( f \), as seen in the proof of Lemma 8. Instead, for \( n = d + \ell - 1 \), one has

\[
\phi \circ f(s_{d+\ell-1}) = \phi \circ \phi^d(s') = \phi^d(s') \quad (47)
\]

because \( s' \) has transient progeny \( d' \leq d \) and length \( \ell' | \ell \), so we know that \( \phi^d(s') \) is in the attractor in the same basin of attraction as \( s' \). Now, using the fact that s is a transient state of transient progeny d and length \( \ell \), as above, we have \( f \circ \phi(s_{d+\ell-1}) = \phi^d(s') \), because \( \phi(s_{d+\ell-1}) = s_d \). In conclusion, \( f \) is covariant on every successor of s.

Now consider the set of predecessors of s, namely,

\[
\mathcal{P}_0 = \{ t : \phi^k(t) = s \text{, for some } k \in \mathbb{N} \} \quad (48)
\]

which contains s itself (obtained for \( k = 0 \)), as represented in Fig. 4.

To respect the dynamics, these states must be mapped to predecessors of \( s' \). We are only interested in the states in \( \mathcal{P}_0 \) that are different from s (provided they exist). If they do not exist, we move to the next step, i.e. the predecessors of the successors of s. Instead, if \( \mathcal{P}_0 \setminus \{ s \} \) is non-empty, for the states \( t \in \mathcal{P}_0 \setminus \{ s \} \), let \( \delta(t; s) \) be the number of steps to go from t to s. We know that the maximum value \( \delta(t; s) \) is exactly \( a(s) \).

Now, by hypothesis \( a(s') \geq a(s) \), so it is possible to find a predecessor \( s'^* \) of \( s' \) such that

\[
\delta(s'^*; s) = a(s) \quad (49)
\]

Now, for \( t \in \mathcal{P}_0 \setminus \{ s \} \), construct

\[
f(t) := \phi^a(s) - \delta(t; s) (s'^*) \quad (50)
\]

Note that the exponent is always non-negative because \( \delta(t; s) \leq a(s) \).

Let us show that Definition (50) makes \( f \) covariant on \( \mathcal{P}_0 \setminus \{ s \} \) when such a set is non-empty. We have

\[
\phi \circ f(t) = \phi^{a(s) - \delta(t; s) + 1} (s'^*) \quad (51)
\]

On the other hand, to calculate \( f \circ \phi(t) \), we have to distinguish two cases. The first case is when \( \phi(t) = s' \), which means \( \delta(t; s) = 1 \). In this case, we have \( f \circ \phi(t) = s' \). Now, Eq. (51) reads

\[
\phi \circ f(t) = \phi^{a(s) - \delta(t; s) + 1} (s'^*) = \phi^{a(s) - 1 + 1} (s'^*) = \phi^{a(s)} (s'^*) = s' \quad (52)
\]

The second case is when \( \phi(t) \neq s' \). Observe that \( \phi(t) \) cannot be a state \( s_n \), for \( n > 0 \), i.e. a successor of s. Indeed, if this were the case, we would have \( \phi(t) = \phi^n(s) \), for some \( n \geq 1 \). Since t is in \( \mathcal{P}_0 \setminus \{ s \} \), we have \( \phi^k(t) = s \), for some \( k \geq 1 \). Combining these two statements we obtain

\[
s = \phi^{k-1} \circ \phi(t) = \phi^{k-1} \circ \phi^k(s) = \phi^{k-1+n}(s) \quad (53)
\]

Since s is transient, this is possible if \( k-1+n = 0 \). Since \( n > 0 \) and \( k \geq 1 \), this condition can never be satisfied. Additionally, \( \phi(t) \) can neither be one of the (possible) states \( w \) in the same basin of attraction as s that are neither successors nor predecessors of s. Indeed, if \( \phi(t) \) were such a w, we would have

\[
s = \phi^{k-1} \circ \phi(t) = \phi^{k-1} \circ \phi^k(w) \quad (54)
\]

which would mean that \( w \) is a predecessor of s, which is against the hypothesis. Therefore, if \( \phi(t) \neq s \), the only possibility is that \( \phi(t) \) is a predecessor of s. In this case,

\[
f \circ \phi(t) = \phi^{a(s) - \delta(\phi(t); s)} (s'^*) = \phi^{a(s) - \delta(t; s) + 1} (s'^*) \quad (55)
\]

All of this proves that the constructed \( f \) is covariant on \( \mathcal{P}_0 \setminus \{ s \} \).

In the same basin of attraction as s there can be some states \( w \) that are neither successors nor predecessors of s. In any case, they are predecessors of some of
FIG. 4. A state \( s \) in a dynamical graph, with \( d = 2 \) and \( \ell = 2 \). The set \( \mathcal{S} \) of the 2+2 = 4 successors of \( s \) is represented in blue. The set \( \mathcal{P}_0 \) of the predecessors of \( s \) is the purple set, along with \( s \) itself. The set \( \mathcal{P}_1 \) is depicted in red, and represents a side chain that ends up in the main blue chain. The set \( \mathcal{P}_2 \), the set of predecessors of \( s_1 \), is made up of the purple states, of the red states and \( s \) and \( s_1 \). By considering all sets \( \mathcal{P}_n \) for \( n = 1, \ldots, d + \ell - 1 \), in the proof of Lemma 9 we take care of all collateral chains ending up in \( \mathcal{S} \).

the successors of \( s \). Then, for every successor \( s_n \), with \( n = 1, \ldots, d + \ell - 1 \), consider the set of its predecessors

\[ \mathcal{P}_n = \{ t : \phi^n(t) = s_n, \text{ for some } k \in \mathbb{N} \}. \]  

(56)

Note that, for every \( n = 1, \ldots, d + \ell - 1 \), \( \mathcal{P}_m \subseteq \mathcal{P}_n \) for \( m < n \). Indeed, consider a state \( t \in \mathcal{P}_m \), with \( m < n \). Then \( \phi^k(t) = \phi^m(s), \) for some \( k \in \mathbb{N} \). However, as \( m < n \), we can consider

\[ \phi^{n-m+k}(t) = \phi^{n-m} \circ \phi^k(t) = \phi^{n-m} \circ \phi^m(s) = \phi^n(s). \]

(57)

Therefore, \( t \in \mathcal{P}_n \), showing the inclusion \( \mathcal{P}_m \subseteq \mathcal{P}_n \). To avoid defining \( f \) redundantly on the sets \( \mathcal{P}_n \), for \( n = 1, \ldots, d + \ell - 1 \), we have to prune the sets \( \mathcal{P}_n \), getting rid of the successors of \( s \) in \( \mathcal{S} \) (43), of the states in \( \mathcal{P}_0 \) (for which \( f \) has already been defined), and taking care of the inclusions \( \mathcal{P}_m \subseteq \mathcal{P}_n \) for \( m < n \). To this end, we define the sets

\[ \mathcal{P}_n := \mathcal{P}_n \setminus (\mathcal{P}_{n-1} \cup \mathcal{S}), \]

(58)

for \( n = 1, \ldots, d + \ell - 1 \). Now, we need to consider only the \( \mathcal{P}_n \)'s that are non-empty (see Fig 4).

Note that the maximum \( n \) for which \( \mathcal{P}_n \) can be non-empty is \( n = d \). Indeed, since, for \( n \geq d \), \( s_n \) is in the attractor, \( \mathcal{P}_n \) contains all the states of the basin, and so does \( \mathcal{P}_{n-1} \) if \( n > d \), so \( \mathcal{P}_n = \emptyset \) for every \( n > d \). We distinguish two cases: when \( 1 \leq n < d' \), and when \( d' \leq n < d \).

In the first case, if \( 1 \leq n < d' \), for the states \( w \in \mathcal{P}_n \), we know that \( \delta(w; s_n) \leq a(s_n) \). By hypothesis, we know that \( a(\phi^n(s')) \geq a(\phi^n(s)) \), so it is possible to find a predecessor \( s_n^* \) of \( \phi^n(s') \) such that

\[ \delta(s_n^*; \phi^n(s')) = a(s_n). \]

(59)

Therefore, if \( w \in \mathcal{P}_n \) (with \( n < d' \)), we construct \( f \) as

\[ f(w) := \phi(\phi^n(s') - \delta(w; s_n)) \]

(60)

Again, this function is well-defined because \( \delta(w; s_n) \leq a(s_n) \). Let us show that it is covariant. To this end, we have

\[ \phi \circ f(w) = \phi(\phi^n(s') - \delta(w; s_n) + 1) \]

(61)

To assess \( f \circ \phi(w) \), we need to distinguish two cases. The first is when \( \phi(w) \notin \mathcal{P}_n \). In this case,

\[ f \circ \phi(w) = \phi(\phi^n(s') - \delta(\phi(w); s_n) \]

\[ + 1(s_n^*) \] \]

(62)

The other case is when \( \phi(w) \notin \mathcal{P}_n \). We claim that the only possibility is \( \phi(w) = s_n \). Indeed, by a similar argument as above, \( \phi(w) \) cannot be a state \( s_m \) for \( m > n \). To see why, suppose by contradiction that \( \phi(w) = \phi^m(s) \) for some \( m > n \). Since \( w \in \mathcal{P}_n \), there exists \( k \geq 1 \) such that \( \phi^k(w) = \phi(s) \). Then, combining these two properties

\[ \phi^n(s) = \phi^{k-1} \circ \phi(w) = \phi^{k-1} \circ \phi^m(s) \]

\[ = \phi^{k-1+m}(s). \]

(63)

Here we assume \( n < d' \); therefore, \( \phi^n(s) \) is still a transient state. Then Eq. (63) is satisfied iff \( n = k-1+m \). This is equivalent to solving \((k-1)+(m-n) = 0 \). Since \( m > n \) and \( k > 1 \), this expression can never vanish, so eq. (63) cannot be satisfied. Additionally, \( \phi(w) \) cannot be in \( \mathcal{P}_{n-1} \), for otherwise there would exist \( k' \geq 0 \) such that \( \phi^{k'}(w) = s_{n-1} \), contradicting the assumption that \( w \notin \mathcal{P}_{n-1} \). Finally, \( \phi(w) \) cannot be a (possible) state \( w' \) that is neither a predecessor nor a successor of \( s_n \). Indeed, if this were the case, we would have \( \phi(w) = w' \), which combined with \( \phi^k(w) = s_n \), for some \( k \geq 1 \), yields

\[ s_n = \phi^{k-1} \circ \phi(w) = \phi^{k-1}(w'), \]

(64)

which is against the hypothesis that \( w' \) is not a predecessor of \( \phi^n(s) \). Therefore, we are only left with the possibility \( \phi(w) = s_n \).

Now, if \( \phi(w) = s_n \), we have \( f \circ \phi(w) = \phi^n(s) \). On the other hand, as \( \phi(w) = s_n \), we have \( \delta(w; s_n) = 1 \). Thus, substituting into Eq. (61) yields

\[ \phi \circ f(w) = \phi^n(s) \]

(65)

Here we have used Eq. (59), which states that \( a(s_n) \) is the number of steps from \( s_n^* \) to \( \phi^n(s') \). This shows covariance of \( f \) on \( \mathcal{P}_n \), for \( n = 1, \ldots, d' - 1 \).

To conclude the proof, consider the case for which some \( \mathcal{P}_n \)'s are non-empty for \( n = d', \ldots, d \). For \( w \in \mathcal{P}_n \), construct

\[ f(w) := \phi^{d-\delta(w; s_n)}(s'), \]

(66)
where the exponent is modulo \( \ell \), and we pick a representative of the remainder classes modulo \( \ell \) from 0 to \( \ell - 1 \). Now we show that this definition makes \( f \) covariant on \( \overline{\mathcal{P}}_n \). Clearly, 
\[
\phi \circ f (w) = \phi^{n - \delta(w; s_n) + 1} (s') .
\] (67)
To assess \( f \circ \phi (w) \), as done above, we need to distinguish two cases. The first case is \( f (w) \in \overline{\mathcal{P}}_n \), in which case
\[
f \circ \phi (w) = \phi^{n - \delta \phi(w); s_n} (s') \\
= \phi^{n - \delta(w; s_n) + 1} (s') .
\] (68)
This case shows that \( f \) is covariant when \( \phi (w) \in \overline{\mathcal{P}}_n \), with \( n = d' \ldots , d \) now we proceed to examine the second case, which is when \( \phi (w) \notin \overline{\mathcal{P}}_n \).

To address the second case, i.e. \( \phi (w) \notin \overline{\mathcal{P}}_n \), for \( n = d', \ldots , d \), we further distinguish two cases. If \( n = d', \ldots , d - 1 \), \( s_n = \phi^n (s) \) (i.e. we exclude \( n = d \)), then \( s_n = \phi^n (s) \) is a transient state; thus, as seen above, if \( \phi (w) \notin \overline{\mathcal{P}}_n \), then the only possibility is \( \phi (w) = s_n \). In this case, \( f \circ \phi (w) = \phi^n (s') \). The fact \( \phi (w) = s_n \) implies that \( \delta (w; s_n) = 1 \). Then \( f \circ \phi (w) \) (cf. Eq. (68)) takes the form
\[
f \circ \phi (w) = \phi^{n - \delta(w; s_n) + 1} (s') \\
= \phi^{n - 1 + 1} (s') \\
= \phi^n (s') .
\] (69)
This concludes the proof of covariance when \( \phi (w) \notin \overline{\mathcal{P}}_n \) and \( n = d' \ldots , d - 1 \).

Now we tackle the final case, namely \( \phi (w) \notin \overline{\mathcal{P}}_n \) and \( n = d \). In this case, for \( n = d \), \( s_d \) is in the attractor, so \( \phi (w) \) this time can be any of the states in the attractor. Indeed Eq. (63), for \( n = d \) is equivalent to \( d \equiv k - 1 + m \mod \ell \), which has solutions even for \( m > d \) and \( k \geq 1 \) (it is enough to take an \( m > d \) such that \( m \equiv d - k + 1 \mod \ell \)). Therefore, we conclude that \( \phi (w) = s_m \), for some \( m \in \{d, \ldots , d + \ell - 1\} \), which is any of the states in the attractor. In this case \( f \circ \phi (w) = \phi^m (s') \). Now, \( \delta (\phi (w); s_d) = d - m \mod \ell \), where for every remainder class modulo \( \ell \) we pick a representative between 0 and \( \ell \), which means \( \delta (w; s_d) = d - m + 1 \mod \ell \). Therefore,
\[
f \circ \phi (w) = \phi^m (s') = \phi^{d - \delta(w; s_d) + 1} (s') ,
\] (70)
where the expression in the exponent is modulo \( \ell \), as above. However, since \( s' \) is in a cycle of length \( \ell' \), the exponent is actually modulo \( \ell' \), in the sense above, because of the length \( \ell' \) of the cycle. Given that \( \ell' \mid \ell \), whether we take the exponent modulo \( \ell \) or \( \ell' \), the state \( \phi^{d - \delta(t; s_d) + 1} (s') \) is the same.

With Lemmas 8 and 9 we finally obtain Theorem 5.

C. The conversion problem: non-attractorness

Now we analyze the conversion problem for discrete dynamical systems where randomness is allowed. In this case, we call the resource “non-attractorness” as all valuable states (i.e. non-free) are non-uniform probability vectors; i.e. they are are not associated with uniform probability distributions supported on attractor deterministic states. We show that randomness acts as an activator of transitions between deterministic states: with stochastic covariant influences it is possible to jump from one deterministic state to another even if their length do not satisfy the constraints of Theorem 5.

Given a stochastic dynamical system \((\mathcal{S}, \phi)\) and two stochastic states \(p\) and \(p'\), we want to establish whether there exists a stochastic matrix \(F\) commuting with the matrix \(\Phi\) associated with \(\phi\), such that \(p' = Fp\). From a computational point of view, the solution to this problem can be achieved with a set of constrained linear equations. The fact that all constraints are expressed by linear equations makes the problem tractable. In fact, we can show that this is a linear decision problem (see Appendix B).

We can still try to see if we can reduce the complexity of determining the solution to the conversion problem by reducing the number of linear conditions to impose. To achieve this, we can study if some constraints on the entries of the matrix \(F\) to be determined follow directly from general properties of commutation of \(F\) with some \(\Phi\), even without specifying what \(\Phi\) is actually like. For this reason, now we study how the covariance condition affects the entries of the most general square stochastic matrix \(F\), which, according to Eq. (16), are the jumping probabilities between deterministic states. For example, we can determine which entries of \(F\) are forced to be equal or to vanish by the covariance condition.

In the stochastic case we are mainly interested in the conversion problem between probability vectors. However, the structure of matrices of covariant influences can also provide a new angle on transitions between deterministic states. Indeed, the entries of such matrices are transition probabilities between deterministic states. In deterministic dynamical systems, influences map a deterministic state \(s\) to one deterministic state \(s'\). In stochastic dynamical systems, influences map a deterministic state \(s\) to deterministic states \(s'\). In stochastic dynamical systems, in general influences will allow multiple transitions from a state \(s\), each of which is weighted by a certain probability. In particular, we say that a transition from a deterministic \(s\) to a deterministic state \(s'\) is allowed if the corresponding transition probability \(p(s'|s)\) is non-zero. Therefore, in the light of Theorem 5, we are interested in seeing if the presence of randomness activates some types of transitions between deterministic states that would otherwise be forbidden by the conditions of Theorem 5. To answer this question, we need once more to understand which transition probabilities \(p(i|j)\) are compatible with the covariance condition expressed in its most general form.
Example 2. Consider a discrete dynamical system $(S, \phi)$ that features two cycles: one of length 2 and one of length 1 (i.e., a fixed point). Let us determine the constraints on a generic stochastic covariant influence acting on the probability vectors of that discrete dynamical system, arising from the commutation condition with the dynamical matrix $\Phi$. We call $s_1$ and $s_2$ the two deterministic states in the cycle of length 2, and $s_3$ the fixed point. With this convention, $p(i|j)$ ($i, j \in \{1, 2, 3\}$) is the transition probability from the deterministic state $s_j$ to the deterministic state $s_i$. The dynamical matrix is

$$\Phi = \begin{pmatrix} 0 & 1 & 0 \\ 1 & 0 & 0 \\ 0 & 0 & 1 \end{pmatrix},$$

where the first column tells us that $s_1$ gets mapped to $s_2$, the second column that $s_2$ gets mapped to $s_1$ (in agreement with the fact that we have a cycle of length 2), and the third column tells us that $s_3$ remains $s_3$ because it is a fixed point. Imposing the covariance condition $[F, \Phi] = 0$ with a generic stochastic matrix $F$ yields new constraints on the entries of $F$:

$$\begin{pmatrix} p(1|1) & p(1|2) & p(1|3) \\ p(2|1) & p(2|2) & p(2|3) \\ p(3|1) & p(3|2) & p(3|3) \end{pmatrix} = \begin{pmatrix} 0 & 1 & 0 \\ 1 & 0 & 0 \\ 0 & 0 & 1 \end{pmatrix}$$

This equality gives some indications on the relations between transition probabilities $p(i|j)$. For transitions within deterministic states in the attractor of length 2, we have:

- $p(1|1) = p(2|2)$
- $p(2|1) = p(1|2)$

There is no constraint on the transition probability $p(3|3)$. Instead if we want to jump from the attractor of length 2 to the fixed point with a covariant influence, we can do so in such a way that $p(3|1) = p(3|2)$. Conversely, if we want to transition from the fixed point to the attractor of length 2 with a covariant influence, we can do so provided $p(1|3) = p(2|3)$. Notice that a transition from a short attractor to a longer one is clearly was forbidden in the purely deterministic scenario described in Theorem 5, so we can say that the presence of randomness activates some transitions and allows us to go against some of the rules prescribed for a dynamical system under a deterministic covariant influence.

Motivated by Example 2, we want to see how the presence of stochastic covariant influences affects the statement Theorem 5. To this end, we focus on the columns of a stochastic matrix $F$ associated with a stochastic covariant influence $f$, where each column represents the evolution of a deterministic state (which can be transient or in an attractor) of the discrete dynamical system under $f$. This is stated in the following lemmas.

Lemma 10. Let $s$ and $s'$ be two deterministic states of a discrete dynamical system $(S, \phi)$. If a transition occurs from $s$ to $s'$ under a stochastic covariant influence, then $d' \leq d$.

Proof. We prove the statement by contradiction: we will assume $d' > d$, and we show that a stochastic transition from $s$ to $s'$ is not possible, i.e., $p(s'|s) = 0$. Let $e_i$ be the vector of the canonical basis associated with the state $s$. Then, if $F$ is the matrix of a covariant influence, the vector $p := Fe_i$ is the $i$th column of $F$, whose entries are the transition probabilities $p(t|s)$, where $t$ is a deterministic state. Among them we find $p(s'|s)$. Now, we have:

$$\Phi^{d}F e_i = \Phi^{d} e_i = \Phi^{d+1} e_i = \Phi^{d} p.$$  (73)

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Let us focus on the basin of attraction of the state \( c \) and which, without loss of generality, we can assume to be into basins of attraction. In the proof of Lemma 10, we partition where we have used the covariance of

\[
\Phi^\prime \phi_n = \left( \begin{array}{c}
\text{attractor} \\
0 \\
\vdots \\
p(s'|s) + \ldots \\
\vdots \\
\end{array} \right) = c_1.
\]

As all the terms in \( c_1 \) are non-negative, this equality implies that \( p(s'|s) = 0 \), which contradicts the hypothesis.

This lemma tells us that randomness in the covariant influence does not allow us to overcome the rule \( d' \leq d \) we derived for the fully deterministic case.

**Lemma 11.** Let \( s \) and \( s' \) be two deterministic states of a discrete dynamical system \( (S, \phi) \). If a transition occurs from \( s \) to \( s' \) under a stochastic covariant influence, then

\[
a(\phi^n(s')) \geq a(\phi^n(s)),
\]

for every \( n \in \mathbb{N} \).

**Proof.** Again, we prove the lemma by contradiction: we will assume there exists an \( n \in \mathbb{N} \) for which \( a(\phi^n(s')) < a(\phi^n(s)) \) and we show that a stochastic transition from \( s \) to \( s' \) is not possible, i.e. \( p(s'|s) = 0 \). Note that this inequality implies that \( \phi^n(s') \) is a transient state, which also means that \( s' \) is transient. If \( s \) is an attractor state, then a transition from \( s \) to \( s' \) cannot occur due to Lemma 10, because the transient progeny of \( s' \) is strictly greater than the transient progeny of \( s \).

We are left with the case when \( s \) and \( s' \) are both transient. Let \( e_i \) be the vector of the canonical basis associated with the state \( s \). Then consider the vector \( F \Phi^n e_i \), where \( F \) is the matrix of a covariant influence. The entries of \( F \Phi^n e_i \) are the transition probabilities from \( \phi^n(s) \), among which we find \( p(\phi^n(s')|\phi^n(s)) \).

As a first step, we will show that \( p(\phi^n(s')|\phi^n(s)) = 0 \). Now, recall we can always find a deterministic state \( s^*_n \) such that \( \phi^n(s) = \phi^a(\phi^n(s)) \). If \( e_{i_{s_n^*}} \) is the vector of the canonical basis associated with the state \( s_{n_{n}^*} \), then we can write \( \Phi^n e_i = \Phi^a(\phi^n(s)) e_{i_{s_n^*}} \). Then,

\[
F \Phi^n e_i = F \Phi^a(\phi^n(s)) e_{i_{s_n^*}} = \Phi^a(\phi^n(s)) F e_{i_{s_n^*}},
\]

where we have used the covariance of \( F \). Now, as done in the proof of Lemma 10, we partition \( F \Phi^n e_i \) and \( F e_{i_{s_n^*}} \) into basins of attraction.

\[
F \Phi^n e_i = b_1 + \ldots + b_k
\]

\[
F e_{i_{s_n^*}} = c_1 + \ldots + c_k.
\]

Let us focus on the basin of attraction of the state \( s' \), which, without loss of generality, we can assume to be \( b_1 \) and \( e_1 \), respectively. Specifically, let \( s_{n_{s_n^*}} \) be the state such that \( \phi^n(s') = \phi^a(\phi^n(s'))(s_{n_{s_n^*}}) \), which is in the same basin of attraction as \( \phi^n(s') \) and \( s' \). Let us list the entries of \( b_1 \) and \( c_1 \) starting from \( s_{n_{s_n^*}} \), then moving to \( \phi(s_{n_{s_n^*}}) \), and so on to the other successors of \( s_{n_{s_n^*}} \), until we come to \( \phi^n(s') \) after a \( \phi^n(s') \) steps. This means that the entry relative to \( \phi^n(s') \) is located to \( a(\phi^n(s')) \) below the first row:

\[
b_1 = \left( \begin{array}{c}
p(s_{n_{s_n^*}}|\phi^n(s)) \\
\vdots \\
p(\phi^n(s')|\phi^n(s)) \\
\vdots \\
\end{array} \right) \quad c_1 = \left( \begin{array}{c}
p(s_{n_{s_n^*}}|s_{n_{s_n^*}}) \\
\vdots \\
p(\phi^n(s')|s_{n_{s_n^*}}) \\
\vdots \\
\end{array} \right).
\]

Now, in the equality \( b_1 = \Phi^a(\phi^n(s)) c_1 \), the first \( a(\phi^n(s')) + 1 \) entries of the vector at the right-hand side are zero:

\[
b_1 = \left( \begin{array}{c}
p(s_{n_{s_n^*}}|\phi^n(s)) \\
\vdots \\
p(\phi^n(s')|\phi^n(s)) \\
\vdots \\
0 \\
\end{array} \right) = \Phi^a(\phi^n(s)) c_1.
\]

The reason for the presence of \( a(\phi^n(s')) + 1 \) zeros is that \( s_{n_{s_n^*}} \) is the farthest predecessor of \( \phi^n(s') \), so at each step (i.e. at each power of \( \Phi \)) we progressively remove the transition probabilities associated with the successors of \( s_{n_{s_n^*}} \). Therefore, \( \Phi \) removes the probability associated with \( s_{n_{s_n^*}} \), \( \Phi^2 \) removes the probability associated with \( \phi(s_{n_{s_n^*}}) \), and so on, until we come to \( \Phi^a(\phi^n(s')) + 1 \), which removes the probability associated with \( \phi^a(\phi^n(s')) \). The fact that \( a(\phi^n(s)) > a(\phi^n(s')) \) guarantees that we are always able to remove the probability associated with \( \phi^n(s') \). Therefore, we conclude that \( p(\phi^n(s')|\phi^n(s)) = 0 \).

Now we will show that \( p(\phi^n(s')|\phi^n(s)) = 0 \) implies \( p(s'|s) = 0 \). Now, \( p(\phi^n(s')|\phi^n(s)) \) is an entry of the vector \( F \Phi^n e_i \). By covariance, \( F \Phi^n e_i = F^n F e_i \). Let us partition \( F e_i \) into basins of attraction:

\[
F e_i = d_1 + \ldots + d_k,
\]

where we assume that \( d_1 \) is the basin of attraction of \( s' \). This time, we list the entries of \( b_1 \) and \( d_1 \) starting from \( s_{n_{s_n^*}} \), where \( s_{n_{s_n^*}} \) is a deterministic state such that \( s' = \phi^a(s'_{s_n^*}) \), similarly to what we have done above. Then we have \( b_1 = \Phi^a d_1 \). In this way, due to the presence of \( \Phi^n \), the term arising from \( p(s'|s) \) will be in the \( n + a(s') + 1 \) entry of \( \Phi^a d_1 \). On the other hand, the \( n + a(s') + 1 \) entry of \( b_1 \) is

\[
p(\phi^{n+a}(s)|\phi^n(s)) = p(\phi^n(s)|\phi^n(s)) = 0.
\]

Then the \( n + a(s') + 1 \) entry of \( d_1 \) must vanish; because it is a sum of non-negative terms (including \( p(s'|s) \)), each term must vanish, from which we finally obtain that \( p(s'|s) = 0 \).
We summarize the previous lemmas with the following theorem.

**Theorem 12.** Let $s$ and $s'$ be two deterministic states of a discrete dynamical system $(S, \phi)$. Then a transition from $s$ to $s'$ is allowed under stochastic covariant influences only if

\[
\begin{align*}
    d' &\leq d, \\
    a(\phi^n(s')) &\geq a(\phi^n(s))
\end{align*}
\]  

for $n \in \mathbb{N}$.

Notice that the divisibility constraint (33) is no longer present in the statement of Theorem 12. Indeed, Example 2 showed a concrete case where stochastic covariant influences overcome the divisibility constraint coming from deterministic covariant influences.

V. APPLICATION TO DISCRETE LOGISTIC MAP

We now investigate how our results translate into a popular case of discrete dynamical system, specifically, the discrete logistic map (DLM) [45], even though it does not have a finite number of states, in contrast with the discrete dynamical systems we have studied so far. The DLM originated in the context of studying biological systems not reaching a steady state, and the logistic map shows, within a simple mathematical model, not only systems not reaching a steady state, and the logistic map is a steady state not achieved for some parameters but also that chaos emerges, that a geometric convergence to chaos via period doubling occurs, and period-three cycles [45, 46].

A. Mathematical model

We employ the standard mathematical description for the DLM [45]. We can think of the DLM as representing the evolution of a certain population. To achieve a compact treatment of the DLM, it is convenient to phrase it by introducing a variable $x$ representing a fraction of the population with respect to the saturation level, i.e. the maximum population for the system. The DLM equation is

\[x \leftarrow rx(1 - x), \quad r \in [0, 4], \quad x \in [0, 1]\]

for $\leftarrow$ referring to the recursive replacement of $x$ by its new value at each successive generation. Here $r$ can be interpreted as a fertility rate. For sufficiently low $r$, i.e., $r \leq 1$, this recursive map yields the asymptotic value $x_{\infty} = 0$, which denotes the steady-state solution and extinction. However, larger $r$ delivers more complicated solutions in the asymptotic long-time limit.

For $r > 1$, the DLM (77) exhibits the phenomenon of period doubling until the onset of chaos at $r_{\text{ch}} \approx 3.56995$. Increasing $r$ leads to bifurcating the steady state, first into period 2, then period 4 and so on (powers of two). The spacing between doubling, with respect to $r$, asymptotically approaches a geometric progression with this progression quantified by Feigenbaum’s number 4.66920.

The onset of chaos starts at $r_{\text{ch}}$ and is characterized by exponentially sensitive dependence on initial conditions. Although chaos dominates for $r \gg r_{\text{ch}}$, these chaotic regions are punctuated by islands of stability including a period-three cycle, whose existence actually guarantees the presence of chaos [47].

We introduce an influence by modifying the recursion relation (77) to the form

\[x \leftarrow rx(1 - x) + f(x),\]

where $f(x)$ is any real function of $x$. Introducing our notation, as $\phi$ evolves the discrete dynamical system for one time step, in the case of the DLM (77), we have that $\phi(x) = rx(1 - x)$. Instead, the influence $f$ is given by the additive term added to the DLM. Imposing covariance amounts to imposing $f \circ \phi = \phi \circ f$. We note that, if $f$ is a covariant influence, then $f^n$ is also a covariant influence, where by $f^n$ we mean $f \circ \ldots \circ f$. This can be easily proven by induction. To find examples of covariant influences from which to take powers and generate other covariant influences, we start by examining the case of a quadratic influence, i.e. of the form

\[f(x) = ax^2 + bx + c,\]

where $a, b, c \in \mathbb{R}$.

B. Covariant influence for the DLM

Now we explain how to impose covariance on the DLM, and we describe the situations that arise in this setting. Specifically, we will see that an influence of the form $f(x) = ax^2 + bx + c$ is covariant only for a very restricted set of values for its parameters. Moreover, we see that covariant polynomial influences of higher degrees can be obtained by taking iterates of the quadratic covariant term.

To impose covariance, we need to calculate $f \circ \phi$ and $\phi \circ f$ separately and then impose equality between these two expressions. Specifically,

\[f \circ \phi = a(\phi x(1 - x))^2 + brx(1 - x) + c\]

and

\[\phi \circ f = r(ax^2 + bx + c)(1 - ax^2 - bx - c).\]

Imposing $f \circ \phi = \phi \circ f$ we get that $a, b, c$ can only take the following values. We find that $c$ must always be zero,
meaning that the covariant influence is forbidden from applying a constant shift. Given that \( c = 0 \), two solutions arise: either \( a = 0 \) and \( b = 1 \), in which case the influence is strictly linear in \( x \), or \( a = -r \) and \( b = r \), in which case the influence is strictly quadratic.

Given that the only covariant influence of polynomial degree one is \( x \), which is an identity map, and the case of polynomial degree two yields a rescaled logistic map, we investigate a higher polynomial degree. In this case, a natural choice of covariant influence is obtained by considering the \( n \)th iterate of the logistic map, i.e. \( f = \phi^n \).

C. Discussion of ‘influenced’ logistic map

The linear case is not a DLM, as the recursive relation is

\[
x \leftarrow rx (1-x) + x,
\]

which is not of the form in Eq. (77). Imposing the constraint that \( x \in [0,1] \) at all generations, as it should be, given that \( x \) represents a population fraction, we find that the new recursive relation is well-posed if and only if \( r \in [0,1] \). In this regime, the uninfluenced DLM yields population extinction for \( r \in [0,1] \). The linear term can be interpreted as proportional immigration term that is proportional to the population fraction. An easy analysis of the convergence of the new recursive relation yields that the presence of the linear migration terms always leads to population saturation (unless, of course, if the initial population is zero). With this interpretation, we can understand why the DLM with the covariant linear term is well posed only for \( r \in [0,1] \): the population decay, as prescribed by the uninfluenced DLM, balances the proportional immigration kick, making \( x \) not exceed the maximum population.

After straightforward calculations, in the case of a quadratic covariant influence, the conditions

\[
a = -r, \quad b = r, \quad c = 0,
\]

make Eq. (78) into the form

\[
x \leftarrow 2rx (1-x),
\]

which just yields a rescaling of the original DLM by a factor 2. This leads to a redefinition of \( \bar{r} \in [0,2] \) in order to keep \( x \in [0,1] \). In this case, the uninfluenced part \( \phi \) and the influence \( f \) coincide, so \( \phi \) and \( f \) actually describe the same evolution. We can interpret this situation as a splitting of the original population obeying the DLM into two halves. As \( x \) is still constrained to be in \([0,1]\) and \( r \) can be interpreted as the fertility rate of the population \([45]\), splitting the population into two means that both \( x \) and \( r \) must be renormalized against the reduced size of the population. In particular, after the splitting into two halves, \( r \) becomes \( r' := r/2 \). Therefore, when considering the evolution of the whole population, it is natural that the fertility rate will be twice the renormalized fertility rate of each half; in symbols \( r = 2r' \).

In the case of higher-degree polynomial influences the recursive relation is of the form

\[
x \leftarrow rx (1-x) + \circ^n rx (1-x)
\]

with \( \circ^n \) the \( n \)th-order composition of Eq. (77) with itself \( n \) times. We can interpret this expression as the presence of immigration that evolves exactly as the original system but \( n \) times as fast. Notice that this general case subsumes the two previous cases, for \( n = 0 \) and \( n = 1 \), respectively. In this approach, we need to take care that \( x \in [0,1] \), being \( x \) a population fraction. If \( r \) is small, the uninfluenced evolution will lead to a decay in fertility rate, and the immigration term, which evolves \( n \) times as fast, will produce an even smaller contribution, thus ensuring that \( x \) is always in the interval \([0,1] \).

Suppose an initial population of \( x = 1/2 \), which means half the saturation value. If \( r = 3/4 \), then the sequence of population fractions is \( 1/2 \rightarrow 3/16 \). With proportional immigration (80), we instead obtain \( 1/2 \rightarrow 11/16 \), which highlights the role of immigration in reaching a higher population fraction. In this case, the population fraction will reach a steady state of 1 due to the contribution from immigration. In the quadratic case, assuming again \( r = 3/4 \), the sequence of population fractions is \( 1/2 \rightarrow 3/8 \), which can be immediately seen as twice as fast the uninfluenced logistic evolution, approaching the 0 steady state. As an instance of higher-degree covariant influence, we consider the case \( n = 2 \). In this case, the immigration term produces a contribution that decays faster than the decay induced by the uninfluenced logistic map. Specifically, starting from \( 1/2 \) we reach \( 309/1024 \) in the next iteration. If we continue the iterations, we see that the population fraction approaches the steady state of 0.1975, which corresponds to the intersection between \( x \leftarrow x \) and the function on the right-hand side of \( \leftarrow \) in Eq. (83) for \( n = 2 \). Thus, our numerical illustration shows that the evolution leads to a steady-state population for \( r = 3/4 \) in each of the \( n \in \{0,1,2\} \) instances, and this steady-state population agrees with the theory expounded above.

VI. DISCUSSION

We study discrete dynamical systems over finite sets, which arise in several situations in different disciplines. In discrete dynamical systems states evolve in discrete steps according to a deterministic rule. Since we assume that the set of states is finite, after a finite time, the states must repeat themselves, ending up in a cycle. Discrete dynamical systems come in two flavors. The first is purely deterministic, where we know the initial state perfectly, and we are able to predict the evolution of each state with certainty. The second incorporates
some randomness, as we do not know the initial state precisely, but only a probability distribution over possible initial states. Despite the randomness in the states, the evolution remains deterministic: if we know the initial state perfectly, we will still be able to predict its evolution with certainty.

Randomness in the initial state arises to model our ignorance about the exact initial state. This happens because the precision with which we probe a system is normally limited, and an exact description of the initial state is thus not achievable or highly impractical. In such situations, we can still describe the system by shifting from the evolution of states to the evolution of probability distributions over states.

As no systems are truly isolated, we incorporate the environment in the description of a discrete dynamical system, where the environment is intended as everything that is outside the dynamical system. In particular, the environment disrupts the normal evolution of the discrete dynamical system, and it can do so in two ways: either deterministically, or randomly. In some cases, the external influence occurs over a longer time scale than the time scale of evolution of a dynamical system. In such situations, the dynamical system has time to adapt to the external influence, so that the influence itself does not become too disruptive, but it interplays nicely with the dynamic of the system. Such an influence is called covariant, and it is the subject of our analysis.

We examine the effect of a covariant influence on a discrete dynamical system with multiple attractors and an environment enacting a deterministic covariant influence. Without the covariant influence, any state will eventually reach an attractor. Turning on the covariant influence allows hopping between different attractors, where we show that hopping from one attractor to another one is possible if and only if the length of the final attractor divides the length of the initial one. Thus, covariant influences shorten the length of attractors, so that the system spontaneously reaches a state of “minimal evolution”. From the perspective of resource theories, a resource is something that goes against the spontaneous evolution induced by covariant influences. For this reason, the resource into play in this setting is called “evolutionarity”. From a biological point of view, this means that a system now can also increase the length of attractor. From a biological point of view, this means that a system can be adaptable to diverse changes induced by the environment.

The results of this article are derived using the general and powerful framework of resource theories, which has been used extensively in quantum information and foundations. For the first time, we exported resource theories outside the domain of physics. Besides the conceptual clarity resource theories brought to the study of discrete dynamical systems and their influences, resource theories also inspired our questions and informed our approach to finding an answer to those questions. For example, the question of how covariant influences affect the evolution of a discrete dynamical systems can be phrased as the conversion problem in resource theories, which in turn can be solved using resource monotones. Such monotones give rise precisely to the constraints we found for the behavior of covariant influences discussed above.

VII. CONCLUSIONS

We introduced the notion of covariant influence, which is the backbone of our analysis. The notion of covariant influence is expressed in terms of the commutation of the influence $f$ with the influenced dynamics $\phi$. The uninfluenced dynamics is in general not reversible.

We developed the theory of discrete dynamical systems under covariant influences, giving a mathematical characterization of which transitions can take place in a discrete dynamical system due to the presence of a covariant influence. Specifically, we analyzed two situations: one where there is no randomness, and the other where the covariant influence is stochastic. These two situations are qualitatively and quantitatively different. For the case of the lack of randomness, we derived necessary and sufficient conditions for the evolution of states under covariant influences (Theorem 5). In this case, the covariant influence allows jumping between attractors whose final length divides the initial length. On the other hand, we showed that randomness permits transitions that violate this divisibility rule, and we were able to give some constraints on the entries of a generic stochastic
matrix representing a covariant influence. However, understanding if a stochastic state can be converted into another with a covariant influence is in general a rather complicated problem, which can be formulated as a linear decision problem.

The underpinning of our results is the framework of resource theories, which is a powerful tool developed in quantum information. In this work, for the first time, we exported it to another area of science, and we showed how it can be applied to the study of dynamical systems, which are widespread in science. The application of resource theories is the mathematical foundation that allowed us to split the evolution of a discrete dynamical system into an influence and an uninfluenced part, in the same spirit one does perturbation theory, but without the restriction that the influence be a small contribution.

Resource theories guarantee the solidity of our results, as they are expressed with a logical and rigorous mathematical framework that can be formalized using category theory. This ensures the validity of our findings for all discrete dynamical systems, as long as the covariance condition holds. The level of generality we obtained is particularly valuable, in that it allows one to understand the behavior of any discrete dynamical system under a covariant influence, without needing to know the exact details of the evolution of that dynamical system. Our results are in the form of laws that can be tested in concrete dynamical systems, paving the way for further analysis, even from an experimental perspective.

There are open questions left for future work. One direction involves formulating a generalized version of the covariant condition, perhaps involving powers of \( \phi \), and analyzing its consequences. Another direction further explores the action of stochastic covariant influences. For example, knowing the dynamical graph of a discrete dynamical system, we would like to be able to write the form of a stochastic covariant matrix immediately, just by looking at the dynamical graph. In particular, it is interesting to characterize the transition probabilities between states in attractors of different lengths. Indeed, we saw in the deterministic case that transitions under a covariant influence are permitted only if the length of the final attractor divides the length of the initial one. In the stochastic case, on the other hand, such transitions are allowed, but we may still observe some special behavior in the transition probabilities when the final state is in an attractor whose length divides the length of attractor in which the initial state is. A third direction involves studying applications of our framework based on covariant influences to concrete settings, such as genetic regulatory networks.

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**Appendix A: Conversion problem: evolutionarity between different dynamical systems**

In §IV B we presented necessary and sufficient conditions for the conversion of a state \( s \) into a state \( s' \) of the same dynamical system \((S, \phi)\). Now we generalize this situation by considering the conversion of a state \( s \) of \((S, \phi)\) into a state \( s' \) of \((S', \phi')\) under any covariant map \( f : S \rightarrow S' \).

As to necessary conditions, notice that the proofs of Lemmas 6 and 7 can be easily extended to the case when \( s \in S \) and \( s' \in S' \); it is enough to insert a prime where needed; therefore Lemmas 6 and 7 still hold when there are two dynamical systems involved. In fact, in some cases, we can use Lemma 6 to show that no free operations exist between \((S, \phi)\) and \((S', \phi')\), as shown in the following example.

**Example 3.** Suppose \((S, \phi)\) is a discrete dynamical system with only one attractor of length 3, and \((S', \phi')\) is a discrete dynamical system with only one attractor of length 2. If there were free operations \( f : S \rightarrow S' \), we would map states in the only attractor of \((S, \phi)\) to states in the only attractor of \((S', \phi')\), but the length of the attractor of \((S', \phi')\) does not divide the length of the attractor of \((S, \phi)\). Therefore, no covariant operations from \( S \) to \( S' \) exist, in agreement with Lemma 6. For the same reason, neither do free operations from \( S' \) to \( S \) exist.

Note that the possible presence of transient states associated with the attractors does not affect this result at all. This shows how strong the condition on divisibility of Lemma 6 is.

Now, let us turn to sufficient conditions when \((S, \phi)\) is not necessarily equal to \((S, \phi')\). In this case, given two states \( s \in S \) and \( s' \in S' \), the two conditions \( d' \leq d \) and \( \ell' \mid \ell \) in Theorem 5 no longer suffice as shown in the following example.

**Example 4.** Let \((S, \phi)\) be a discrete dynamical system with one attractor of length 3 and one of length 4. Let \((S', \phi')\) be a discrete dynamical system with only one attractor of length 2. Let \( s \) be a state in the attractor of length 4 of \((S, \phi)\), and let \( s' \) be a state in the attractor of length 2 of \((S', \phi')\). \( s \) and \( s' \) respect the conditions \( d' \leq d \) and \( \ell' \mid \ell \), but a transition is not possible between them because covariant maps do not exist between \((S, \phi)\) and \((S', \phi')\) because of incommensurate attractor lengths. Indeed, if they did exist, we must map the attractor
of length 3 in \((S, \phi)\) to an attractor with a length that
divides 3, but such an attractor does not exist in \((S', \phi')\).

Example 4 shows that, when \((S, \phi) \neq (S', \phi')\), the
conversion problem between \(s\) and \(s'\) does not depend
only on \(s\) and \(s'\), but involves the other deterministic
states. The reason is the lack of covariant maps from
one dynamical system to the other. Here we provide
a necessary and sufficient condition for the existence of
covariant maps between two different discrete dynamical
systems.

**Proposition 13.** Given two discrete dynamical systems,
\((S, \phi)\) and \((S', \phi')\), there exist covariant maps \(f : S \rightarrow S'\)
iff, for every attractor of \((S, \phi)\) of length \(\ell\), there exists
an attractor of \((S', \phi')\) of length \(\ell' \mid \ell\).

**Proof.** Necessity follows from Lemma 6. To prove
sufficiency, we need to show that, under the hypotheses
of the theorem, it is possible to construct a covariant
map \(f : S \rightarrow S'\). To this end, for every attractor of \(\ell\)
in \((S, \phi)\), let us pick an attractor state \(s\), and a state \(s'\)
in an attractor of \((S', \phi')\) whose length \(\ell'\) divides \(\ell\).
Then map all the states in the basin of attraction of \(s\) to states
in the basin of attraction of \(s'\), following the procedure
outlined in the proof of Lemma 8 (in this case \(d = 0\)).
This yields a covariant map on all the states in the basin
of attraction of \(s\). Repeating this procedure for all the
attractors of \((S, \phi)\), we can construct a covariant map \(f\)
on all states of \((S, \phi)\).

Now we can provide necessary and sufficient conditions
for the conversion of a deterministic state in \((S, \phi)\) into
a state of \((S', \phi')\) with covariant maps.

**Theorem 14.** Let \(s\) be a deterministic state of a discrete
dynamical system \((S, \phi)\) and \(s'\) be a deterministic state of
a discrete dynamical system \((S', \phi')\). Then there exists
a covariant map \(f : S \rightarrow S'\) converting \(s\) into \(s'\) iff

1. \(\ell' \mid \ell\);
2. \(d' \leq d\);
3. \(a(\phi'^n(s')) \geq a(\phi^n(s))\) for \(n = 0, \ldots, d' - 1\).

**Proof.** Again, necessity was proved in Lemmas 6 and
7. Sufficiency follows from combining the proofs of
Theorem 5 and Prop. 13. Indeed, under the hypotheses
of the theorem, following the proof of Theorem 5 we can
map the basin of attraction of \(s\) to the basin of attraction
of \(s'\) in a covariant way. Regarding the other basins of
attraction in \((S, \phi)\), following the proof of Prop. 13, we
can covariantly map every basin of attraction of \((S, \phi)\)
to a basin of attraction of \((S', \phi')\).

Notice that in Theorem 14 we have one additional
condition than in Theorem 5: the existence of covariant
maps from \((S, \phi)\) to \((S', \phi')\) is not present in Theorem 5.
Indeed, if \((S, \phi) = (S', \phi')\), the
existence of covariant maps is always guaranteed. To
understand why, note that, if \((S, \phi) = (S', \phi')\), then,
for every attractor of length \(\ell\), we can always find an
attractor of length \(\ell' \mid \ell\); it is enough to take the same
attractor. By Prop. 13, this is enough to guarantee the
existence of covariant maps. This is why Theorem 5 does
not mention this check at all.

**Appendix B: Linear programming**

In this appendix, we show that, in the presence of
randomness and stochastic influences, we can recast the
conversion problem defined in §IV C. Here, a square matrix
\(\Phi\), representing the dynamics, is given by the
specification of the dynamical system, and we want to
know if, for a pair of probability vectors \(p\) and \(q\), there
exists a stochastic matrix \(F\) such that \(Fp = q\) and
\([F, \Phi] = 0\).

Let \(\Phi\) be the dynamical matrix of a dynamical system
\(S\) with \(M\) states (see §I A). As such, \(\Phi\) can be expressed
in terms of the generator of the dynamics \(\phi : S \rightarrow S\) as

\[
\Phi = (e_{\phi(1)} \ e_{\phi(2)} \cdots e_{\phi(M)}) .
\]  \hspace{1cm} (B1)

Let

\[
F = \sum_{i=1}^{M} \sum_{j=1}^{M} f_{ij} e_i e_j^T
\]  \hspace{1cm} (B2)

be an \(M \times M\) column-stochastic matrix, and \(T\) denote the
transpose. Observe that the \(j\)th column of \(F\Phi\) is given by
\(Fe_{\phi(j)}\). On the other hand, the \(j\)th column of \(\Phi F\) is
given by \(\sum_{k=1}^{M} f_{kj} e_{\phi(k)}\). Note that the \(i\)th component
of \(\sum_{k=1}^{M} f_{kj} e_{\phi(k)}\) is given by \(\sum_{k=\phi^{-1}(1)} f_{kj}\), where \(\phi^{-1}(i)\)
denotes the pre-image of \(s_i\).

Hence, the condition \(F\Phi = \Phi F\) is equivalent to

\[
f_{i\phi(j)} = \sum_{k\in\phi^{-1}(i)} f_{kj},
\]  \hspace{1cm} (B3)

for every \(j, k \in \{1, \ldots, M\}\). This condition in turn can be
expressed in terms of Hilbert-Schmidt inner products

\[
\text{tr}[FG_{ij}^T] = 0
\]  \hspace{1cm} (B4)

for every \(i, j \in \{1, \ldots, M\}\), where

\[
G_{ij} = e_i e_{\phi(j)}^T - \sum_{k\in\phi^{-1}(i)} e_k e_j^T.
\]  \hspace{1cm} (B5)

The condition that the columns of \(F\) sums to one can be
expressed as

\[
Fu = u, \quad u := \begin{pmatrix} 1 \\ \vdots \\ 1 \end{pmatrix} .
\]  \hspace{1cm} (B6)
Define
\[ U_1 := \begin{pmatrix} 1 & \cdots & 1 \\ 0 & \cdots & 0 \\ \vdots & \ddots & \vdots \\ -1 & \cdots & -1 \\ 0 & \cdots & 0 \end{pmatrix} \] (B7)
the idea is to reshape the vector \( u \) into a matrix. For each \( i > 1 \), define the matrix
\[ U_i := \begin{pmatrix} 1 & \cdots & 1 \\ 0 & \cdots & 0 \\ \vdots & \ddots & \vdots \\ -1 & \cdots & -1 \\ 0 & \cdots & 0 \end{pmatrix} \] (B8)
where the first row is \( u^T \), the \( i \)th row is \( -u^T \), and the rest of the rows are zeros. Then, we can express the condition (B6) as
\[ \text{tr} \left[ F U_i^T \right] = \delta_{i1} \] (B9)
for all \( i \in \{1, \ldots, M\} \).

Let \( p, q \in \mathbb{P}^M (\mathbb{R}) \) be \( M \)-dimensional probability vectors of which we want to study the conversion under covariant influences. For each \( i \in \{1, \ldots, M\} \), define the matrix
\[ H_i := \begin{pmatrix} 0 & \cdots & 0 \\ \vdots & \ddots & \vdots \\ p_1 - q_i & \cdots & p_M - q_i \\ \vdots & \ddots & \vdots \\ 0 & \cdots & 0 \end{pmatrix} \] (B10)
whose \( i \)th row is \( p^T - q_i u^T \) and the rest of the rows are zeros. Then, we can express \( F p = q \) as
\[ \text{tr} \left[ F H_i^T \right] = 0 \] (B11)
for all \( i \in \{1, \ldots, M\} \).

We therefore arrive at the following SDP decision problem. Determine if there exists \( F \geq 0 \) (entry-wise) such that for all \( i, j \in \{1, \ldots, M\} \)
1. \( \text{tr} \left[ F U_i^T \right] = \delta_{i1} \);
2. \( \text{tr} \left[ F G_j^T \right] = 0 \);
3. \( \text{tr} \left[ F H_i^T \right] = 0 \).

Note that, if we remove the first condition \( \text{tr} \left[ F U_i^T \right] = 1 \), we just need to assume in addition that \( F \) is a non-zero matrix whose entries are non-negative. Therefore, there exists such an \( F \) in the orthogonal complement \( W^\perp \) to the subspace
\[ W = \text{span}_\mathbb{R} \{ U_2, \ldots, U_M, G_{11}, \ldots, G_{MM}, H_1, \ldots, H_M \} , \] (B12)
with respect to the Hilbert-Schmidt inner product.

Alternatively, let \( f \in \mathbb{R}^M \) be the vector obtained by reshaping \( F \). Similarly, let \( g_{ij}, u_i, \) and \( h_j \) be the vectors obtained by reshaping \( G_{ij}, U_i, \) and \( H_i \), respectively. With these notations, the inner products above take the form of the dot products
1. \( f \cdot u_i = \delta_{i1} \);
2. \( f \cdot g_{ij} = 0 \);
3. \( f \cdot h_i = 0 \).

Hence, the above conditions are equivalent to the existence of a non-trivial and non-negative solution to the homogeneous linear system of equations
\[ A f = 0 \] (B13)
where
\[ A = (u_2, \ldots, u_M, g_{11}, \ldots, g_{MM}, h_1, \ldots, h_M)^T \] (B14)
is an \((M^2 + 2M - 1) \times M^2\) matrix.

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