SOME GEOMETRIC PERSPECTIVES IN CONCURRENCY THEORY

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(communicated by Gunnar Carsson and Rick Jardine)

Abstract

Concurrency, i.e., the domain in computer science which deals with parallel (asynchronous) computations, has very strong links with algebraic topology; this is what we are developing in this paper, giving a survey of “geometric” models for concurrency. We show that the properties we want to prove on concurrent systems are stable under some form of deformation, which is almost homotopy. In fact, as the “direction” of time matters, we have to allow deformation only as long as we do not reverse the direction of time. This calls for a new homotopy theory: “directed” or di-homotopy. We develop some of the geometric intuition behind this theory and give some hints about the algebraic objects one can associate with it (in particular homology groups). For some historic as well as for some deeper reasons, the theory is at a stage where there is a nice blend between cubical, \( \omega \)-categorical and topological techniques.

1. Introduction

Concurrency theory deals with systems in which several computational activities (called processes in general) can be performed at the same time, in an asynchronous manner. These were introduced in order to have increased computational power, so that computations can be achieved faster (essentially in scientific computing), or so that some concurrent transactions can be handled efficiently (user interfaces, embedded systems reacting to the external environment etc.) or just handled at all (mostly because of the amount of memory needed, as for concurrent databases).

The variety of applications that motivated the use of concurrent machines has led to many different architectures. The main problem about concurrency is to have processes cooperating for a common goal. Cooperation implies some form of synchronisation and information passing. This can be done through message passing for instance. In this class of models, processes have their own local memory, which cannot be accessed by other processes. The way to communicate values to other processes is by explicitly sending values to these other processes, which will have to explicitly ask for receiving values. One of the first of this class of models, is the rendezvous models (as used in most process algebra, like CCS [46], CSP [34] etc.) in which the action of sending is blocking the sender until the receiver actually receives the corresponding message. Symmetrically the action of receiving blocks the receiver until the message is actually sent. This is the simplest of all message-passing models (also called synchronous message-passing).

The variations of it include, non-blocking send but blocking receive, non-blocking send and receive (asynchronous message-passing), broadcasts to groups of processes instead of “point to point” communication etc.

Another important class of concurrent architectures is shared memory style. Here, processes have a local memory indeed, where they can perform their local computations, but also have
a common memory space, which is accessible to all. Communication between processes is essentially asynchronous and is realized by writing and reading values in this common space, as pictured in Figure 1. Processes $P_1, \ldots, P_n$ are writing and reading through shared locations such as scalar variables $x$ and $y$ (containing boolean or integer values for instance), and also through more complex structures, such as $z$, a “3-cell buffer” i.e. a variable consisting of a queue of 3 values. If concurrent reading of $a$ by several processes is not a problem in general, concurrent writing of scalar variables is not to be allowed. At the hardware level, this would mean at best, undefined behaviour, and at worst, short circuit. Therefore, it is necessary to “protect” the accesses to shared variables by some mechanism. A classical one is by using “semaphores” introduced by E. W. Dijkstra [10] in 1968.

Basically, before a process tries to write on a location in shared memory, it has to put a “lock” on it (through its associated semaphore), blocking the other processes which try to write at the same time and on the same location. Formally, the action of putting a lock on location $x$ is denoted by $P_x$ (using E. W. Dijkstra’s notation [10]). In case $x$ is some more complex structure than a read/write variable (such as $z$ above), at most $n \geq 1$ processes can hold a lock on $x$ (here with $z$, $n = 3$) before blocking the accesses by other processes. In this case we call the associated semaphore a $n$-semaphore. After some process has written what it needed to write on $x$, it can relinquish safely its lock by doing action $Vx$; this will allow another process to acquire a lock on $x$, i.e. to un-block it.

Let us forget about actual calculations on $x$, $y$, $z$ etc. and focus only on the locking, unlocking mechanism (the coordination of processes involved). We will then identify shared locations with their associated $n$-semaphores. This urges us to consider throughout this paper (except for some minor exceptions) a simplified programming language, in which processes are regular expressions on the alphabet $\{Pa, Va \mid a \in Loc\}$, where $Loc$ is a set of “locations”. Each of these locations are in fact $n$-semaphores, for some $n$, defined by a map $s : Loc \to \mathbb{N}$. Regular expressions are formed freely from the alphabet $\{Pa, Va \mid a \in Loc\}$ by application of the following algebraic operators: $+$ (which is associative and commutative), $.$ (which is associative), and unary operator $*$. “Elementary moves” (or actions) are elements of the alphabet, i.e. $Pa$, $Vb$ etc. $A + B$ means that sequences of actions that can be taken are those of $A$ or those of $B$ — this is non-deterministic choice. Sequences of actions of $A.B$ are concatenations of sequences of actions of $A$ and of actions of $B$ (this is the concatenation operation), and sequences of actions of $A^*$ are any number of concatenations of sequences of actions of $A$ (this is the Kleene star operation, or finite unbounded iteration).

What are we looking for now? We want to be able to derive properties of concurrent machines, even of such a simplified one. Of course, the theory of sequential computation is very much advanced and the properties of interest for sequential computation (what function of the arguments are we computing? Is the computation always terminating for all its arguments? How long will this take? etc.) are not the one we are dealing with here. The novelty in concur-
rent programming resides not in the fact we are computing another class of functions (which would contradict Turing's thesis) but is the fact that coordination between processes does matter. For instance, we might have forgotten to properly lock some locations, creating an unexpected behaviour of the program. On the contrary, we might have constrained too much the coordination, preventing the program to carry out normal computation. This is called a deadlocking situation. Another property of interest is to know whether a concurrent system can go into a "bad state" or not. Typically, we are trying to solve a "reachability problem", e.g., do we have an execution in our system which will go through such bad states? Also, we can ask for slightly more subtle properties: for some applications (we will see an example later on), some sequences of accesses to resources are considered right whether some others are not. It is therefore of primary importance to be able to classify such sequences; this will actually lead to homotopical arguments.

Before getting to this, let us briefly show how this would normally be dealt with. Of course to be able to prove things, one needs a mathematical model, in particular for the notion of execution (sequence of action) in a concurrent system.

There is a great variety of models for concurrency, as witnessed in [63] for instance. Transition systems are one of the oldest semantic models, both for sequential and concurrent systems:

**Definition 1.** A transition system is a structure \((S, i, L, Tran)\) where,
- \(S\) is a set of states with initial state \(i\)
- \(L\) is a set of labels, and
- \(Tran \subseteq S \times L \times S\) is the transition relation

Transition systems are nothing but discrete dynamical systems: in general the transition relation \(Tran\) is represented as a directed graph of actions. For instance the transition system depicted in Figure 2 gives semantics to the process \(Pa.Pb.c^*.(Va.Vb+Vb.Va)\), i.e. to a process which locks \(a\), then \(b\) then does some sequence \(c\) any finite number of times (this can be a computation on \(a\) and \(b\)), then unlocks \(a\) and \(b\) in any order. This behaviour can be seen by looking at paths (or executions) in this directed graph, from the leftmost state (the initial state) to the rightmost ones (the final states).

A simple way to look at processes in parallel is to build a transition system for each process and then to construct some kind of fibered product of all these graphs of actions (this has a formal sense, see for instance [2]): states of this transition system are now tuples of states of each individual process, and transitions from one to another are interleavings of transitions of each individual process. For instance, the graph of actions for \(T_1 = Pb.Pa.Vb.Va\) in parallel with \(T_2 = Pa.Pb.Va.Vb\) is shown in Figure 3. State 1 is the initial state, actions on parallel segments have the same label.

Now we can see that state 13 is not a "correct" final state. State 23 consists of the pair of endpoints of digraphs representing each process, but not 13, which has nevertheless no future. This is a deadlock. In this situation, the first process \(T_1\) has a lock on \(b\), waiting for a lock on \(a\) whether the second one \(T_2\) has a lock on \(a\) waiting for a lock on \(b\). This is typical of a "deadly embrace" as E. W. Dijkstra originally put it.

We can also ask ourselves whether this concurrent system can be in a state we do not want (which is rather artificial here); this would be a state in which \(T_1\) would have a lock on \(a\) and just released a lock on \(b\), whether \(T_2\) would have a lock on \(b\) and just released a lock on \(a\). Looking at the graph 3 one sees that this is precisely state 19, but there is not path from the initial state 1 to 19, so we never go through this "bad state".

Last but not least, we can also try to classify the different access orders to resources in this system. Looking at all the 10 paths from state 1 to state 23 in the directed graph of Figure 3, we see that we have only two such orders: \(T_1\) holds locks on \(b\) then \(a\) before \(T_2\) does, or \(T_2\) holds locks on \(a\) then \(b\) before \(T_1\) does. This situation is typical of concurrent databases, and is known under the name "serializability".
A distributed database can be seen as a shared-memory machine (containing items) on which processes (called transactions) act by reading and writing, getting permissions to do so by using the appropriate functions on attached semaphores. One of the main purposes of this area is to ensure coherence of the distributed database while ensuring good performance, through a definition of suitable policies (protocols) for transactions to perform their own actions (with $P$ and $V$). This entails of course that deadlock-freedom of transactions is of importance. Correctness of a distributed database is itself very often expressed by some form of a serializability condition.

Suppose we have two transactions $T_1 = Pb.Vb.Pa.Va$ and $T_2 = Pa.Va.Pb.Vb$ trying to modify two items $a$ and $b$. There is a path of execution in which $T_1$ acquires $b$, $T_2$ acquires $a$, then $T_1$ acquires $a$ and $T_2$ acquires $b$. Think of the database to represent airplane tickets (for instance $b$ is the return ticket corresponding to the one-way ticket $a$), and the two transactions to represent remote booking booths, the action between a $P$ and its corresponding $V$ is writing a name on the ticket. The situation here is that $T_1$ will have reserved its one-way ticket and $T_2$ will have reserved its return ticket only. This is not an allowed behaviour. It is not equivalent to a purely serial schedule which are the only ones that are specified as correct (only one of $T_1$ or $T_2$ gets the whole lot of tickets). Of course, this could be seen on the corresponding transition system, but if we have many complex processes running altogether, the “state-space” and therefore the path-space becomes enormous. Therefore it is important to have a way to “retract” this onto smaller transition systems (or shapes) where we can still observe similar state or path like properties. This is where some idea from algebraic topology sneaks in. We will see in the next sections how this can be made precise.

Organisation of the paper. In Sections 2, 3, 4 and 5, we show how to model these phenomena using, in a natural manner, concept from topology and combinatorial algebraic topology. This will give us a meaning for the terms used above, such as “retract”, first in a topological model (Section 2) and then in a combinatorial model (Sections 3 and 4). This is all based on a notion of deformation, or homotopy, which is slightly different from the usual homotopy of topological spaces. Here the direction of time should be preserved, meaning that no deformation can be done in the inverse direction of time. This is why the newly defined homotopy theory is called directed homotopy or “dihomotopy”.

To fully reflect the combinatorial model, we have to refine the topological model of Section 2; this is done in Section 5. Then we can attack in Section 6 a first geometric study of the notions such as deadlocks, schedules, serialisability conditions etc. This is only a first step, ideally one should try to find computable invariants of dihomotopy. Some leads are given in Section 7, but there again, this implies some refinement of the modeling, to have nice and “precise” functors; some of which are shown in Section 8. Then one can try to see if standard
results, such as van Kampen or exact sequence theorems still hold in the new theory. Some hints are given in Section 9. We conclude by some perspectives in Section 10.

Some further references. The “topological” formalization that follows is one of the most recent ones, and essentially dates back to [13] and [14], but is based on much older results [10].

The combinatorial (cubical) and homological calculations are older, and have been at the center of [25], starting with [20], [24] and [26].

I actually only realized the relationships between the combinatorial and the topological approaches quite recently, and have been aware of this line of research only after J. Gunawardena published his very enlightening paper [32].

There are some ideas about using n-categories in [49]. It is only quite recently that these have given their full flavour, see [26] for a start, where many algebraic invariants are also introduced. The “unification” of these approaches has led us to the concept of globular CW-complex [22] which I will briefly describe in Section 8.

The interested reader can find more references about this in [27] or [14].

2. A topological approach

The first “algebraic topological” model I am aware of is called progress graph and has appeared in operating systems theory, in particular for describing the problem of “deadly embrace” in “multiprogramming systems”. Progress graphs are introduced in [9], but attributed there to E. W. Dijkstra. In fact they also appeared slightly earlier (for editorial reasons it seems) in [55].

The basic idea is to give a description of what can happen when several processes are modifying shared resources. Given a shared resource $a$, we see it as its associated semaphore that rules its behaviour with respect to processes. For instance, if $a$ is an ordinary shared variable, it is customary to use its semaphore to ensure that only one process at a time can write on it (this is mutual exclusion). Then, given $n$ deterministic sequential processes $Q_1,\ldots,Q_n$, abstracted as a sequence of locks and unlocks on shared objects, $Q_i = R^1a_1^i.R^2a_2^i\cdots R^nia_i^{n_i}$ ($R^k$ being $P$ or $V$ for respectively acquiring and releasing a lock on a semaphore), there is a natural way to understand the possible behaviours of their concurrent execution, by associating to each process a coordinate line in $\mathbb{R}^n$. The state of the system corresponds to a point in $\mathbb{R}^n$, whose $i$th coordinate describes the state (or “local time”) of the $i$th processor.

Consider a system with finitely many processes running altogether. We assume that each process starts at (local time) 0 and finishes at (local time) 1; the $P$ and $V$ actions correspond to sequences of real numbers between 0 and 1, which reflect the order of the $P$'s and $V$'s. The initial state is $(0,\ldots,0)$ and the final state is $(1,\ldots,1)$. An example consisting of the two processes $T_1 = Pa.Pb.Vb.Va$ and $T_2 = Pb.Pa.Va.Vb$ gives rise to the two dimensional progress graph of Figure 4.

The shaded area represents states which are not allowed in any execution path, since they correspond to mutual exclusion. Such states constitute the forbidden area. For instance, look at Figure 4 again and take a point whose abscissa is (strictly) between local times marked as $Pb$ and $Vb$ and whose ordinate is (strictly) between local times marked also as $Pb$ and $Vb$. Having these coordinates means that both $T_1$ and $T_2$ have acquired $b$ and not relinquished it, which is impossible since $b$ can only be shared by at most one process at a time. This point ought to be forbidden.

An execution path is a path from the initial state $(0,\ldots,0)$ to the final state $(1,\ldots,1)$ avoiding the forbidden area and increasing in each coordinate - time cannot run backwards. We call these paths directed paths or dipaths. This entails that paths reaching the states in the dashed square underneath the forbidden region, marked “unsafe” are deemed to deadlock, i.e. they cannot possibly reach the allowed terminal state which is $(1,1)$ here. Similarly, by
reversing the direction of time, the states in the square above the forbidden region, marked “unreachable”, cannot be reached from the initial state, which is $\langle 0, 0 \rangle$ here. Also notice that all terminating paths above the forbidden region are “equivalent” in some sense, given that they are all characterized by the fact that $T_2$ gets $a$ and $b$ before $T_1$ (as far as resources are concerned, we call this a schedule). Similarly, all paths below the forbidden region are characterized by the fact that $T_1$ gets $a$ and $b$ before $T_2$ does.

On this picture, one can already recognize many ingredients that are at the center of the main problem of algebraic topology, namely the classification of shapes modulo “elastic deformation”. As a matter of fact, the actual coordinates that are chosen for representing the times at which Ps and Vs occur are unimportant, and these can be “stretched” in any manner, so the properties (deadlocks, schedules etc.) are invariant under some notion of deformation, or homotopy. This is a particular kind of homotopy though, and this will be at the center of many difficulties in later work. We call it (in subsequent work) directed homotopy or dihomotopy in the sense that it should preserve the direction of time. For instance, the two homotopic shapes, all of which have two holes, of Figure 6 and Figure 7 have a different number of dihomotopy classes of dipaths. In Figure 6 there are essentially four dipaths up to dihomotopy (i.e. four schedules corresponding to all possibilities of accesses of resources $a$ and $b$) whereas in Figure 7, there are essentially three dipaths up to dihomotopy.

Before going to the formalization, we should ask ourselves if there is not a simpler way to model these properties.

There is another method to determine deadlocks in such situations, which was of course known long ago and was entirely graph-theoretic, known as the request graph. Figure 5 depicts the request graph corresponding to the progress graph of Figure 4. Nodes of this graph are resources of the concurrent system, i.e. here, semaphores. There is an directed edge from a resource $x$ to a resource $y$ if there is a process which has locked $x$ and needs to lock $y$ at a given time. A sufficient condition for such systems to be deadlock-free is that their corresponding request graphs be acyclic\(^1\). Unfortunately, this is not a necessary condition in general. For instance a request graph cannot capture the notion of $n$-semaphores, i.e. resources that can be shared by up to $n$ processes but not $n+1$ (for instance, asynchronous buffers of communication of size $n$ which can be “written” on by at most $n$ processes). This in fact really calls for some higher-dimensional versions of graphs.

There is no need to resort to fancy $n$-semaphores to see that request graphs are not enough.

\(^1\)Note that this is a very geometric condition indeed.
Consider the following concurrent program, which is composed of processes A, B and C in parallel (introduced for other reasons by Lipsky and Papadimitriou, mentioned in [32]), its request graph (Figure 8) and the corresponding progress graph viewed from different angles in Figure 9:

\[ A = \text{Px}. \text{Py}. \text{Pz}. \text{Vx}. \text{Pw}. \text{Vz}. \text{Vy}. \text{Vv} \]
\[ B = \text{Pu}. \text{Pv}. \text{Px}. \text{Vv}. \text{Pz}. \text{Vv}. \text{Vx}. \text{Vv} \]
\[ C = \text{Py}. \text{Vv}. \text{Px}. \text{Pv}. \text{Vz}. \text{Vv}. \text{Vu}. \text{Vv} \]

The request graph for this example contains cycles, but it can be proved that it does not deadlock.

A progress graph can be seen as a topological space - a sub-space of IR\(^6\) in fact. The topology is necessary to formally define the notion of path, which has to be continuous (executions cannot skip from one point to another in no time). We also need a partial order, allowing to characterize the “direction” of the time flow, i.e. to characterize future and past of points (which are states of the concurrent system). The two should be at least minimally compatible. At least we should be able to take (topological) limits under the partial order sign, leading to the following definition:

2The holes in the cube of states are in fact represented as plain shapes.
Definition 2. A po-space (or partially ordered space) is a pair $(X, \leq)$ formed by a topological space $X$ together with a partial order $\leq$ such that $\leq$ is closed (i.e. $\leq$ is a closed subset of $X \times X$ with the product topology).

This implies two natural properties: the sets $\uparrow x = \{ y \mid x \leq y \}$ and $\downarrow x = \{ y \mid y \leq x \}$ are closed subsets of $X$.

We need now to define suitable morphisms between local po-spaces, which in turn will give the notion of dipath.

Definition 3. Let $(X, \leq_X)$ and $(Y, \leq_Y)$ be two po-spaces.

A continuous function $f : X \to Y$ is called a dimap if for all $y, z \in X : y \leq_X z \Rightarrow f(y) \leq_Y f(z)$.

A dipath on $X$ is then a dimap $f : \tilde{I} \to X$ where $\tilde{I}$ is the topological space $I = [0, 1] \subseteq \mathbb{R}$ with the (global) partial order inherited from the one of $\mathbb{R}$. We write $P_1(X)$ for the set of dipaths on $X$ and $P_1^{\alpha, \beta}(X)$ for the set of dipaths on $X$ going from $\alpha$ to $\beta$.

Now, we can define more precisely what we mean by deformation of dipaths, which we call directed homotopy, or dihomotopy. It is very important here to fix the extremities of dipaths. The idea is that, contrarily to the classical case where it suffices\(^3\), in order to characterize a shape, to choose a basepoint and then to consider loops around this basepoint modulo homotopy, here we really need two basepoints\(^4\). As a matter of fact, it is very unlikely that we have lots of directed loops in our shapes (in fact, there are only trivial constant directed loops in a progress graph) so we have to choose a source basepoint and a target basepoint, and then study all dipaths between these two points modulo dihomotopy.

Definition 4. Let $f$ and $g$ be two dipaths on $X$ between an initial point $\alpha$ and a final point $\beta$. A dihomotopy between $f$ and $g$ is a dimap from $\tilde{I} \times I$ to $X$ such that for all $x \in \tilde{I}$, $t \in I$, $H(x, 0) = f(x)$, $H(x, 1) = g(x)$ and $H(0, t) = \alpha$, $H(1, t) = \beta$. We write $f \sim g$.

A first example of the kind of properties one wants to check on some systems, which involves a characterization of dihomotopy classes of dipaths is the “serialisation” property in some concurrent databases.

Look for instance at Figure 7, we have already been mentioning it in the introduction. All paths of execution above the left hole are equivalent to a serial execution of transaction $T_2$ then transaction $T_1$. All paths of execution below the right hole are equivalent to the serial execution of transaction $T_1$ then $T_2$. The third type of dipath is not a serial dipath; it describes several equivalent cases, for instance: $T_1$ acquires $k$, $T_2$ acquires $a$, then $T_1$ acquires $a$ and $T_2$ acquires $b$.

\(^3\)When we restrict to a connected component.

\(^4\)Or in fact as we will see later on when looking at higher-order homotopies, we need a base dipath.
Can we now make sense of the serializability condition of the introduction? A simple criterion can seem to do the job here. It seems on the example of Figure 7 that connectivity of the forbidden region is a necessary condition for a system to be serializable. But it is unfortunately not a sufficient condition. Consider again the Lipsky/Papadimitriou example (Figure 9). The forbidden region is connected but not simply connected (it is homeomorphic to a filled in torus). There is a dipath going through the center of this “torus” which cannot be deformed on one of the boundaries of the outer cube of states.

We definitely need some new theory here, developed in Section 6 in particular. But let us divert a little and look at another “geometrically flavoured” model for concurrency.

3. A Combinatorial Approach

Let us try to get back closer to transition systems now. In fact, there is another “geometric” model for concurrency which seems to relate more to transition systems than progress graphs, introduced in the article by Vaughan Pratt [49], and which has inspired a lot the following work on the subject (for instance [25]). It was essentially motivated by a defect in the duality between event structures and automata⁵, two well known mathematical models for concurrency.

The models which have been introduced to fix this defect, which can be attributed to the fact that the former semantics is a “truly-concurrent” one where the latter is a simulation by interleaving, were based on one form or another of CW-complex. Such objects are gluings of “elementary” shapes along their boundaries. The following explanation is inspired by [25].

Consider first transition systems. They allow to model concurrency with an interleaving semantics. They already are (1-dimensional) geometric objects. Many important semantic properties are actually geometric properties on their underlying graph of transitions. For instances, initial and final (or deadlocking) states, unreachables states, cycles, branchings and confluences, as seen in the introductory section. All these geometric notions are important for validating or proving correct concurrent systems. For instance, branchings are of importance for the so-called “branching-time” semantic equivalences such as bisimulation equivalence [45], and deadlocks and unreachable states are useful for static analysis (such as model-checking for instance).

In fact, the modelisation of concurrent systems by interleaving naturally constructs cubical shapes. For instance squares like $a \mid b$:

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which represents the asynchronous execution of actions $a$ and $b$ ($a'$ and $b'$ are transition which have respectively $a$ and $b$ as “labels”).

The natural idea is that this interleaving is an expansive encoding of the fact that $a$ and $b$ are independent indeed. Using a physicist’s image, we would like to represent all the ways we can mix together any number of sub-actions of $a$ and $b$ (all the shuffles of possibly infinitely many chunks of $a$ and $b$) as shown in Figure 10.

It is thus natural to put into our model all the possible subdivisions of the square, i.e. the interior $A$ of the square as well as its boundary. Adding in the model the concept “interior of a square” (as well as “interior of any $n$-cube” when we model $n$ concurrent processes in parallel) naturally leads to the notion of pre-cubical set.

⁵Which, later, will also motivate the introduction of Chu spaces [50].
The usual way to define a precubical set is to define the boundary operators; for instance, for the square, we have four boundary operators, respectively corresponding to \(a, b, a'\) et \(b'\). This is not enough since we want to encode the “direction of time” as well in this model. In dimension one, we use a “directed graph” of transitions; we would like something similar here but with “higher-dimensional” transitions.

The choice made in [25] is to divide the family of boundary operators into two families of operators. In the case of a square, we have a family of two source boundary operators \(d^0_0\) and \(d^0_1\), with \(d^0_0(A) = a\) and \(d^0_1(A) = b\), and a family of two target boundary operators \(d^1_0\) and \(d^1_1\), with \(d^1_0(A) = a'\) and \(d^1_1(A) = b'\). This naturally extends the notion of source and target of arcs of directed graphs.

If a 2-transition (square) is nothing but an independence relation between two 1-transitions, a 3-transition, or cube, is not merely a shortcut for 3 independence relations. The fact that action \(a\) is independent of action \(b\), \(b\) independent of \(c\), and \(c\) of \(a\) does not imply that \(a\), \(b\) and \(c\) can be executed altogether in an asynchronous manner. It is the case for instance of an abstraction of a print spooler with two printers, or of two floating-point units\(^6\), or of a communication buffer with two cells, i.e. of a semaphore \(s\) initialized to 2\(^7\).

If we use the notations of E. W. Dijkstra [16], it suffices to consider the three actions \(a = b = c = Ps; s\) can be shared by two but not by three processes at the same time (see Figure 11). This kind of object which synchronizes very weakly is of great importance for

\(^6\)For instance in Intel microprocessors.
\(^7\)What we call an \(n\)-semaphore (here with \(n = 2\)).
These properties cannot be expressed simply (if at all) in most of the other mathematical models for concurrency, as asynchronous transition systems, prime event structures etc. a notable exception being Petri nets. These have other drawbacks: they are not very compositional, which makes them clumsy for analyzing concurrent programs.

Of course this can be easily (and fruitfully) generalized. The concurrent access of \( n+1 \) processes to a given \( n \)-semaphore is represented by the boundary of an hypercube of dimension \( n+1 \). This means we need to put in our model \( n \)-transitions for all \( n \geq 0 \).

There again, we divide the family of boundary operators into two subfamilies: the set of \( n \)-transitions will have a family of \( n \) source boundary operators \( d_i^n, 0 \leq i \leq n-1 \) (all giving \( (n-1) \)-transitions), and a family of \( n \) target boundary operators \( d_j^n, 0 \leq j \leq n-1 \). For instance, for \( n = 3 \):

\[
\begin{align*}
(0,0,0) \quad & \xrightarrow{b} (1,0,0) \\
& \quad \downarrow c \\
& (0,1,0) \\
\quad & \xrightarrow{a} (1,1,0) \\
& \quad \downarrow c \\
& (0,0,1) \\
\quad & \xrightarrow{b} (1,0,1) \\
\quad & \quad \downarrow c \\
& (0,1,1) \\
\quad & \xrightarrow{a} (1,1,1)
\end{align*}
\]

The interior \( D \) of the cube has three source boundaries, the three faces containing \((0,0,0)\), and three target boundaries, the three faces containing \((1,1,1)\). Let \( A, B \) and \( C \) be the faces:

\[
\begin{align*}
& ((0,0,0),(1,0,0),(0,0,1),(1,0,1)) \\
& ((0,0,0),(0,1,0),(0,0,1),(0,1,1)) \\
& ((0,0,0),(1,0,0),(0,1,0),(1,1,0))
\end{align*}
\]

respectively. Let \( A', B' \) and \( C' \) be the parallel faces of \( A, B \) and \( C \) respectively.

Let \( d_0^3(D) = A, d_1^3(D) = B, d_2^3(D) = C \) and \( d_0^3(D) = A', d_1^3(D) = B', d_2^3(D) = C' \). Then \( d_0^3(A) = b, d_1^3(A) = c, d_0^3(B) = a, d_1^3(A) = c, d_0^3(C) = a, d_1^3(C) = b \).

More generally, we can prove what we see here, that is, the boundary operators can be defined so that they satisfy the following commutation rules (for \( i < j \) and \( k, l = 0,1 \)):

\[
d_i^k \circ d_j^k = d_{j-1}^k \circ d_i^k
\]

For instance, for a 2-transition, the relation with \( k = 0, l = 1 \) and \( i = 0, j = 1 \) means that the source of the target number one (i.e. of \( b' \)) is the same as the target of the source number zero (i.e. of \( a \)). This gives us the (classical, but presented in a slightly different manner) notion of precubical set:

**Definition 5.** A precubical set is a graded set \( M = (M_i)_{i \in \mathbb{N}} \) with two families of operators:

\[
M_n \xrightarrow{d_i^n} M_{n-1}
\]

\( ^8 \)A shared FIFO stack with two entries allows for instance to implement wait-free binary consensus for two processes, whereas a simple shared variable (with atomic reads and writes) does not allow it. A good reference for these problems is [37].

\( ^9 \)They are nevertheless very much used for analyzing boolean (telecommunication) protocols.
(i, j = 0, \ldots, n - 1) satisfying the relations
\[ d^i_k \circ d^j_j = d^i_{j-1} \circ d^i_i \]
\((i < j, k, l = 0, 1)\)

Of course, this is linked to progress graphs (just discretise using squares, in the trivial way, the picture of Figure 4), but there is more to it than one could suspect. This is partly developed in Section 5.

I used this formalization in my first article on the subject [29]. In fact, cubical (that we will see a bit later) and precubical sets have a quite old history. They have been used in the first developments of algebraic topology by D. Kan and later by J.-P. Serre in his thesis [53]. Nowadays, combinatorial algebraic topology uses simplicial sets \([40, 17]\), union of simplices of all dimensions, glued along their faces. In J.-P. Serre’s thesis, cubical sets were preferred to simplicial sets because, for studying fibrations, which are locally canonical projections from a cartesian product of two topological spaces to the first one, it is simpler to consider cubical sets which have good properties with respect to projections\(^\text{10}\).

We can also define a (combinatorial) notion of dipath and of dihomotopy. As we will see in Section 5, they are closely linked with the (topological) notions of dipath and dihomotopy we have seen in Section 2.

Let \(N\) be a precubical set. A dipath in \(N\) is a sequence \(p = (p_1, \ldots, p_k)\) of elements of \(N\) such that for all \(i, 1 \leq i < k\), \(d^i_i(p_i) = d^i_i(p_{i+1})\). \(d^i_i(p_1)\) is the initial point of \(p\). \(d^i_i(p_k)\) is the final point of \(p\).

We say that two combinatorial dipaths are dihomotopic if we can go from one to the other by a finite number of “transpositions” of two consecutive actions. It is in fact exactly the same notion as the “partial commutative monoids” used in Mazurkiewicz trace theory [41].

Consider the “contiguity” relation (or “combinatorial dihomotopy”) as follows. Let \(p = (p_1, \ldots, p_k)\) and \(q = (q_1, \ldots, q_l)\) be two non-empty dipaths with same initial and final states. Then, it is easy to see that \(k = l\). We say that \(p\) and \(q\) are elementary contiguous if there exists \(u, 1 \leq u < k\) such that,
- for all \(i < u\) and \(i > u + 1\), \(p_i = q_i\),
- there exists \(A \in M_2\) such that \(d^u_u(A) = p_u\), \(d^u_u(A) = p_{u+1}\), \(d^u_u(A) = q_u\) and \(d^u_u(A) = q_{u+1}\).

The contiguity relation is the equivalence relation (i.e. the reflexive, symmetric and transitive closure of) generated by the elementary contiguity. We will see in Section 5 that this is very close to dihomotopy in the topological space indeed.

To actually give semantics to concurrent systems, the use of cubical or precubical sets is natural as I already explained (for instance by starting with the “interleaving semantics” of transaction systems); this is fully exemplified in [13] for our little P, V language. The link can be made more formal as hinted at in next section.

4. Interpretation in terms of concurrency theory

Reciprocally, all “geometric shapes” built by glueing together hypercubes of any dimension along their faces can be presented as a precubical set\(^\text{11}\). For this to be clear, we need a number of definitions and lemmas.

Let \(K\) and \(L\) be two precubical sets. Then \(f = (f_n)_{n \in \mathbb{N}}\) is a morphism of precubical sets from \(K\) to \(L\) if for all \(n \in \mathbb{N}\), \(f_n\) is a function from \(K_n\) to \(L_n\) such that:
\[ f_n \circ \partial^0_i = \partial^0_i \circ f_{n+1} \]

\(^{10}\)Even if a simplicial construction was published later, see for instance [42].

\(^{11}\)This terminology was recently suggested to us by Ronnie Brown. Before this, we used the term “precubical set” by analogy with the old term “presimplicial set” or simplicial (formerly called complete presimplicial) sets without degeneracy operators.
(for all $i$, $0 \leq i \leq n$)

This forms a category called $\mathcal{Y}^S$. It is a presheaf category as follows. Let $\square^1$ be the free category whose objects are $[n]$, where $n \in \mathbb{N}$, and whose morphisms are generated by,

$$[n-1] \xrightarrow{\delta^0_i} [n] \xrightarrow{\delta^1_j} [n]$$

for all $n \in \mathbb{N}^*$ and $0 \leq i, j \leq n-1$, such that $\delta^0_i \delta^1_j = \delta^1_j \delta^0_{i-1}$ $(i < j)$.

Now, the category $\square^1 \text{op} \mathbf{Set}$ of contravariant functors from $\square^1$ to $\mathbf{Set}$ (morphisms are natural transformations) is isomorphic to the category of precubical sets. This implies, by general theorems [39], that $\mathcal{Y}^S$ is an elementary topos. Moreover it is complete and co-complete because $\mathbf{Set}$ is complete and co-complete.

The category of precubical sets of dimension less or equal than $n$ can be seen as the presheaf category $(\square^1 \leq n) \text{op} \mathbf{Set}$ where $\square^1 \leq n$ is the full subcategory of $\square$ where objects are $[p]$ with $p \leq n$.

**Lemma 1.** Let $T_n$ (respectively $T^S_n$) be the function from $\mathcal{Y}$ (respectively $\mathcal{Y}^S$) to $\mathcal{Y}_n$ (respectively $\mathcal{Y}^S_n$), which to every $M \in \mathcal{Y}$ (respectively $M \in \mathcal{Y}^S$) associates $N \in \mathcal{Y}_n$ (respectively $N \in \mathcal{Y}^S_n$) with,

$$N([k]) = M([k]) \text{ if } k \leq n$$

$$N(\epsilon_i : [k+1] \rightarrow [k]) = M(\epsilon_i) \text{ for } k < n$$

$$N(\delta^0_i : [k-1] \rightarrow [k]) = M(\delta^0_i) \text{ for } k < n$$

defines a functor, called the n-truncation functor.

Now, let $D_{[p]}$ be the representable functor from $\square^1$ to $\mathbf{Set}$ with $D_{[p]}([p]) = \square^1([p], [n])$. We define the singular n-cubes of a precubical set $M$ to be any morphism $\sigma : D_{[p]} \rightarrow M$.

**Lemma 2.** The set of singular n-cubes of a precubical $M$ is in one-to-one correspondence with $M_n$. The unique singular n-cube corresponding to a n-cube $x \in M_n$ is denoted by $\sigma_x : D_{[p]} \rightarrow M$. It is the unique singular n-cube $\sigma$ such that $\sigma(1_{D_{[p]}}) = x$.

**Proof.** The proof goes by Yoneda’s lemma [38].

There is only one morphism in $\square$ from a given $[n]$ to itself, the identity of $[n]$, hence $D_{[p]} \{1_d\}$ is a functor which has only as non-empty values the $D_{[p]}([p])$ with $p < n$ (“it is of dimension $n-1$”). We write $\partial D_{[p]}$ for this functor. For $\sigma$ a natural transformation from $D_{[p]}$ to $M$, we write $\partial \sigma$ for its restriction to $\partial D_{[p]}$.

**Proposition 1.** Let $M$ be a precubical set. The following diagram is co-cartesian (for $n \in \mathbb{N}$),

$$\bigcoprod_{x \in M_{n+1}} \partial D_{[p+1]} \xrightarrow{\bigcoprod_{x \in M_{n+1}} \partial \sigma_x} T_n(M)$$

$$\bigcup_{x \in M_{n+1}} D_{[p+1]} \xrightarrow{\bigcup_{x \in M_{n+1}} \sigma_x} T_{n+1}(M)$$

where $\partial D_{[p+1]} = T_{n}(D_{[p+1]})$ and $\partial \sigma_x = \sigma_x | \partial D_{[p+1]}$.

**Proof.** We mimic the proof of [17]; it suffices to prove that the diagram below (in the category
of sets) is cocartesian for all $p \leq n + 1$,

$$\prod_{x \in M_{n+1}} (\partial D_{[p+1]}|_p \coprod_{x \in M_{n+1}} (\partial \sigma_x|_p (T_n(M))_p \subseteq \prod_{x \in M_{n+1}} (D_{[p+1]}|_p \coprod_{x \in M_{n+1}} (\sigma_x|_p (T_{n+1}(M))_p \subseteq$$

since colimits (hence pushouts) are taken point-wise in a functor category into $\text{Set}$.

For all $p < n + 1$, the inclusions are in fact bijections, and the diagram is then obviously cocartesian.

For $p = n + 1$, the complement of $\bigcup_{x \in M_{n+1}} (\partial D_{[p+1]}|_p$ in $\bigcup_{x \in M_{n+1}} (D_{[p+1]}|_p$ is the set of copies of cubes $\text{id}_{D_{[p+1]}}$, one for each cube of $M_{n+1}$. This means that the map $\bigcup_{x \in M_{n+1}} (\sigma_x|_p$ induces a bijection from the complement of $\bigcup_{x \in M_{n+1}} (\partial D_{[p+1]}|_p$ onto the complement of $(T_n(M))_p$. This implies that the diagram is cocartesian for $p = n + 1$ as well.

This lemma states that the truncation of dimension $n + 1$ of a precubical set $M$ is obtained from the truncation of dimension $n$ of $M$ by attaching some standard $(n+1)$-cubes $D_{[p+1]}$ along the boundary $\partial D_{[p+1]}$ of $(n+1)$-dimensional holes. In fact, any precubical set $M$ is the direct limit of the diagram consisting of all inclusions $T_{n+1}(M) \to T_n(M)$, hence is also the direct limit of the diagram consisting of all the cocartesian squares above. Computer-scientifically, this means that any precubical set can be seen as a (unlabelled) transition system, which is its 1-skeleton, on which we add independence relations. Filled-in squares specify that two actions commute, i.e. that they can be run asynchronously. This is exactly the asynchronous automata sort of models [5, 54] and [41]. Then in higher-dimensions, we fill in cubes etc. meaning that we add some extra ($n$-ary) independence relations. This is fully worked out in [28] in the form of adjunctions between suitable categories of transition systems and of asynchronous automata with (pre-)cubical sets. In fact, to do this properly, we have two steps to make. First, it is easy to relate $\square^{1\text{op}}\text{Set}$ with unlabelled transition systems only if we take as morphisms for transition systems, the “total” morphisms of [63], i.e. graph morphisms; this is unfortunately not quite enough, and to add the right morphisms is reflected on the geometric side by added degeneracy operators, i.e. by going from precubical to cubical sets. Secondly, we have to add up labels to cubical sets. This can easily be done using a comma category constructions; labelled cubical sets are just labelling morphisms from a “shape” cubical set to a “labelling” cubical set. In fact, it is a particular case of a sconing construction [16], and as a side result we automatically know that labelled cubical sets will also form a topos.

5. A Useful Generalization

Progress graphs are a very limited model for concurrency: in particular, we are unable to give a semantics to recursive processes other than unfold all loops, whereas we could give a semantics to loops in the combinatorial model. This is not very satisfying and motivates a more local definition: a first natural idea is to impose only a local partial order instead of a global one, on a topological space, leading to a definition very much alike differentiable manifolds. We recap here the main definitions, the full details can be found in [14].

**Definition 6.** Let $X$ be a topological space. A collection $U(X)$ of pairs $(U, \leq_U)$ of opens of $X$, covering $X$, and partially ordered by $\leq_U$ is called a local partial order on $X$ if for all $x \in X$ there exists a non-empty open neighbourhood $W(x) \subseteq X$ such that the restrictions of $\leq_U$ to $W(x)$ coincide for all $U \in U(X)$, i.e.,
for all $U_1, U_2 \in \mathcal{U}(X)$ such that $x \in U_1$ and for all $y, z \in W(x) \cap U_1 \cap U_2$:

$$y \leq_{U_1} z \Leftrightarrow y \leq_{U_2} z.$$ 

We call the collection of opens $U$ of the definition, an atlas for $X$ (by analogy with differentiable manifolds). Again by analogy with differentiable manifolds, there is a notion of equivalence of atlases:

**Definition 7.**

- Two local partial orders on $X$ are "equivalent" if their union is a local partial order on $X$.

- A locally partially ordered space consists of a topological space $X$ and of an equivalence class of local partial orders on $X$. If moreover there exists a covering $\mathcal{U}$ in this equivalence class such that all $(U, \leq_U) \in \mathcal{U}$ are po-spaces, then we say that $X$ is a local po-space.

Let us give a simple example. A "directed" loop $S^1 = \{e^{i\theta} \in \mathbb{C} \}$ is a local po-space: it suffices to take $U_1 = \{e^{i\theta} \in S^1 | 0 < \theta < 2\pi \}$ with the order induced by the natural order on the $\theta$ and $U_2 = \{e^{i\theta} \in S^1 | \pi < \theta < 3\pi \}$ again with the natural order on the $\theta$.

We need now to define suitable morphisms between local po-spaces, which in turn will give the notion of dipath.

**Definition 8.** Let $(X, \mathcal{U})$ and $(Y, \mathcal{V})$ be two local po-spaces.

A continuous function $f : X \to Y$ is called a dimap if for all $x \in X$ there exists a subset $V(f(x)) \subset Y$ on which $\leq_Y$ is well defined and $U(x) \subset f^{-1}(V(f(x)))$ on which $\leq_X$ is well defined, such that for all $y, z \in U(x): y \leq_X z \Rightarrow f(y) \leq_Y f(z)$.

There again, a dipath on $X$ is then a dimap $f : \tilde{I} \to X$ where $\tilde{I}$ is the topological space $I = [0, 1] \subset \mathbb{R}$ with the (global) partial order inherited from the one of $\mathbb{R}$. We still write $P_1(X)$ for the set of dipaths on $X$ and $P_1^{\alpha, \beta}(X)$ for the set of dipaths on $X$ going from $\alpha$ to $\beta$.

The notion of dihomotopy is exactly the same as for po-spaces; let $f$ and $g$ be two dipaths on $X$ between an initial point $\alpha$ and a final point $\beta$. A dihomotopy between $f$ and $g$ is a dimap from $\tilde{I} \times I$ to $X$ such that for all $x \in \tilde{I}, t \in I$, $H(x, 0) = f(x), H(x, 1) = g(x)$ and $H(0, t) = \alpha, H(1, t) = \beta$. We write once more $f \sim g$.

If we want to be more general, we should consider maximal dipaths (with respect to an obvious "extension" partial order) and not dipaths from an initial point to a final point. This is partially developed in [14] but there are still a number of open problems, in particular about infinite dipaths (on local po-spaces which are not compact).

Now, we can link the (new) topological model with the combinatorial one (the sequel is also taken from [14]).

Let $\Box_n$ be the "standard" $n$-cube in $\mathbb{R}^n$ ($n \geq 0$),

$$\Box_n = \{(t_1, \ldots, t_n) | \forall i, 0 \leq t_i \leq 1\}$$

and let $\delta^k : \Box_{n-1} \to \Box_n, 1 \leq i \leq n, k = 1, 2, \ldots$ be the continuous functions ($n \geq 1$),

![Diagram of the $n$-cube and its faces](image)

\[\delta^k : \Box_{n-1} \to \Box_n, 1 \leq i \leq n, \quad \delta^1 = \delta_{i, n-1}, \quad \delta^2 = \delta_{i-1, n-1} \]

\[\delta^2 \quad \delta^1 \]

\[\Box_n \quad \Box_{n-1} \]

In fact, the equivalence relation we need is the transitive closure of the relation we define here.
Figure 12: Illustration of the transitivity of $\leq_x$.

defined by,

$$\partial_i^{x}(t_1, \ldots, t_{n-1}) = (t_1, \ldots, t_{i-1}, k, t_i, \ldots, t_{n-1})$$

Given a precubical set $M$, consider now the set $R(M) = \bigsqcup M_n \times \Box_n$. The sets $M_n$ can be considered as topological spaces with the discrete topology and $\Box_n$ inherits the topological structure of $\mathbb{R}^n$. Thus $R(M)$ is a topological space with the disjoint union topology. Let now $\equiv$ be the equivalence relation induced by the identities:

$$\forall k, i, n, \forall x \in M_{n+1}, \forall t \in \Box_n, n \geq 0,$$

$$(\partial_i^{x}(x), t) \equiv (x, \partial_i^{t}(t))$$

Let $|M| = R(M) / \equiv$ with the quotient topology. The topological space $|M|$ is called the geometric realisation of $M$.

All this is rather classical as a direct imitation of the geometric realisation functor from simplicial sets to topological spaces. The only trouble here is to find how to interpret the direction of time as it is prescribed in precubical sets (as seen in the definition of dipaths for instance) in terms of local partial orders. For this, we restrict ourselves to non-pathological situations\(^{13}\).

**Definition 9.** Let $M$ be a precubical set. We say that $M$ is not self-linked if for all its $n$-cubes $x$, $\partial_i^{x}(x) = \partial_i^{x}(x)$ implies $k = k'$ and $l = l'$.

Let $x \in M$. The star of $x$ in $M$ is

$$St(x, M) = \{ y \mid \partial_{i_1}^{y} \ldots \partial_{i_u}^{y} y = x \}$$

(see for instance [60]).

Then we set, for $y \in St(x, M)$,

$$(x, t) \leq_{U^x} (y, u) \text{ if } \partial_{i_1}^{x} \ldots \partial_{i_u}^{x}(t) \leq u \text{ in } \Box_{n+1}$$

$$(y, u) \leq_{U^x} (x, t) \text{ if } \partial_{i_1}^{y} \ldots \partial_{i_u}^{y}(t) \geq u \text{ dans } \Box_{n+1}$$

Let $x$ be an element of $M$ and $(z, v)$ be a point in $U^x$ whose carrier is $z$. We set $(z, v) \leq_{x} (y, u)$ if there exists $t$ in the star of $x$ and $t$ such that $(z, v) \leq_{U^x} (b, t) \leq_{U^x} (y, u)$.

It is a partial order indeed; the only difficulty lies in the proof of transitivity (see Figure 12). As a matter of fact, both on the left and right hand sides of the figure, we have $z \leq_{b} y$ and $y \leq_{b} a$ but on the left hand side, $z \leq_{x} a$, and on the right hand side $z \leq_{x} a$ where $b$ is the segment going from the front faces to the back faces from $z$.

Then, the geometric realisation of a non self-linked precubical set $M$ defines a local po-space with atlas $\{ St(x, M) / x \in M_0 \}$ and local partial orders $\leq_{x}$ on each $St(x, M)$.

\(^{13}\) Which might well be non-necessary; this is currently worked out.
The geometric realisation is functorial, I refer the reader to [14]. We also have a right-adjoint to this functor, which is a “singular cube functor” defined very similarly as in simplicial sets theory.

The correspondence between homotopical properties in the topological case and in the combinatorial (cubical) case looks hopefully quite nice. This implies that we will be able to reason about concurrent systems both geometrically on local pospaces (for instance on progress graphs) and algebraically or combinatorially on cubical sets.

This has at least been proven in the simpler case of dimension two in [14]. The geometric realisation of a combinatorial dipath \( p \) of \( M \) induces a (topological) dipath \( |p| \) in \( |M| \). Every combinatorial dihomotopy between \( p \) and \( q \) in \( M \) induces a (topological) dihomotopy between \( |p| \) and \( |q| \).

We also have the inverse one could hope for. Let \( L \) be a finite precubical set and \( h \) be a dipath in \( |L| \) (i.e. a dipath from \( \square_1 \) to \( |L| \)). Then there exists a “cubical approximation” \( f : S_k \to L \) of \( h \) where \( S_k \) is a subdivision of \( \delta \). Moreover \( f \) defines a (combinatorial) dipath \( (f(u_1), \ldots, f(u_k)) \) which we call \( \tilde{f} \). Finally, \( |f| \) is homotopic to \( |\tilde{f}| \) in \( |L| \).

6. First study of dihomotopy

We have already seen that the equivalence classes of dipaths modulo dihomotopy are less numerous in general than the homotopy classes of dipaths. Since the article [49], as well as in [25], I had the intuition that studying the dihomotopy classes of dipaths with fixed extremities was equivalent to studying the homotopy classes of dipaths with fixed extremities. In fact this is not true. It suffices to consider the example of Figures 13 and 14 which give semantics to terms\(^{14}\)

\[\#\text{sem c 2}\]

\(^{14}\#\text{sem 2} \text{ means that c is a 2-semaphore. I have used in the sequel the syntax that my toy static analyzer uses (see http://www.dii.ens.fr/~goubault/analyse.html).}\]
The two dipaths that are represented on these pictures are homotopic but not dihomotopic. The two dipaths do correspond to real executions of a simple concurrent program (c being a very simple 2-place buffer, a and b being two shared scalar variables). This implies we need new tools.

In fact, it is easy to build an analogue of a fundamental groupoid, which will only be a fundamental category in fact. It is obviously linked to the construction of disconnected sets [14], and also to the recent constructions of S. Sokolowski (see [59]), but there is still some work to be done in that direction (see [52]).

Let $X = (X, U, (\leq_U)_{U \in U_i})$ be a local po-space. We can define a composition operation between some of the dipaths of $X$, going from $P_1^{\alpha \beta}(X) \times P_1^{\beta \gamma}(X)$ to $P_1^{\alpha \gamma}(X)$:

**Definition 10.** Let $f \in P_1^{\alpha \beta}(X)$, $g \in P_1^{\beta \gamma}(X)$. We define $h : \bar{I} \to X$ as follows: for $x \in \bar{I}$,

$$h(x) = \begin{cases} 
    f(2x) & \text{if } 0 \leq x \leq \frac{1}{2} \\
    g(2x - 1) & \text{if } \frac{1}{2} \leq x \leq 1
\end{cases}$$

Then $h \in P_1^{\alpha \gamma}(X)$.

It is not a commutative nor associative operation in general, as it is not too in the classical case. Similarly to the classical case, this operation induces a composition operation of classes of dipaths modulo dihomotopy and then becomes associative.

We can then define the following category $\mathcal{C}(X)$:

- its objects are the points of $X$,
- its morphisms are the dihomotopy classes of dipaths; a morphism from $x$ to $y$ is a dihomotopy class $[f]$ of a dipath $f$ going from $x$ to $y$.

The composition defined earlier is an associative operation with identities (the dihomotopy classes of constant dipaths), we use it as the composition of morphisms.

In fact, this construction even defines a functor from the category of local po-spaces to the category of categories. Let $f : X \to Y$ be a diimap from a local po-space $X$ to a local po-space $Y$. We define $\mathcal{C}(f)$ as a morphism from $\mathcal{C}(X)$ to $\mathcal{C}(Y)$:
on objects $x$ of $C$, $C(f)(x) = f(x)$,
on morphisms $[\omega]$ of $C$, with $\omega$ any dipath, $C(f)([\omega]) = [f \circ \omega]$.

We have introduced in [14] the notion of “disconnected components” to study the dihomotopy classes of dipaths. This should be the natural counterpart of connected components in usual topology, but as the relation $x R y$ if there exists a dipath from $x$ to $y$ is certainly not an equivalence relation (and we certainly do not want to make it symmetric since this would mean we would study the arcwise connected components), this is more complicated.

**Definition 11.** 1. The homotopy history of a maximal dipath $\alpha : I \to X$ is

$$ha := \{ y \in X \mid \exists \text{ a dipath } \beta \text{ going through } y \text{ and } \alpha \sim \beta \}$$

2. Two points are homotopy history equivalent of

$$x \in ha \iff y \in ha \text{ for all } \alpha \in P_t(X).$$

3. The disconnected components of $X$ are the connected components (in the classical sense) of the equivalence classes of dipaths modulo the homotopy history equivalence in $X$.

For instance, the complement of the “Swiss flag” in $I^2$ (see Figure 15) has 10 disconnected components. This example gives semantics to the program having as parallel processes $T_1 = Pa.Pb.Vb.Va$ and $T_2 = Pb.Pa.Va.Vb$ (where $a$ and $b$ are 1-semaphores).

In region 1, we have all possible futures. In region 2, we can only go in the future to regions 4 and 6, i.e. the program will deadlock or $T_2$ will access to $a$ and $b$ before $T_1$. In region 6, we can only have come from 2 and go to 9; $T_2$ accesses $a$ and $b$ before $T_1$. In region 9, we can have “come” from the unreachable region 7 or from 6. In region 10, we can have come from any other region except 4.

The complement of the “two ordered holes” in $I^2$ (see Figure 16), which gives semantics to $Pa.Va.Pb.Vb \mid Pa.Va.Pb.Vb$ has 7 classes modulo the homotopy history equivalence. One of these contains both the initial point 0, the final point 1, and a region in the center of the Figure. This class is decomposed into 3 disconnected components.

The “room with 3 barriers” example in $I^3$ (see Ex. 13) has 8 classes modulo the homotopy history equivalence. One of the classes (in the center) is decomposed into two disconnected components.

This point of view has in particular made possible to prove that the “2-phase locking” protocol, which regulates the access to fields of a distributed database is correct (i.e. is sequentialisable). We can find this proof, by M. Raussen, based on ideas of [32], in the article [14]. The reader should notice that S. Sokolowski has defined in [59] a quite similar point of view (about disconnected components) but without discriminating the future of dipaths. This
can be more interesting in some situations (one of which might be when studying bisimulation equivalence, [45]). The interested reader can look at his other papers, [57], [58] and [56].

The problem now is to get to calculate or characterize somehow these dangerous regions etc. I discuss in Section 7 some ideas that have been used for this purpose. It is to be noted that in the case of the regular expressions we had at the beginning, we have a “critical point” approach to obstructions to dihomotopy, which is very algorithmic in nature. We refer the reader to [13], [12] and [15] (for the detection of deadlocks and unsafe regions), and [51] (for the classification of dipaths modulo dihomotopy in such models). Let us concentrate on the more general models in the following section.

7. Dihomotopy invariants

7.1. Homology

In algebraic topology, it is well known that homotopy is a subtle notion. It is in general very hard to prove that two topological spaces are homotopically equivalent, i.e. that one is an “elastic” deformation of the other. Even homotopy groups, whose isomorphism is necessary and sometimes sufficient to decide of the homotopy equivalence are fairly hard to compute.

Nevertheless, there exists so called “homotopy invariants” which can be computed. A homotopy invariant is a functor which to every topological space (or one in a suitable sub-category) X associates a mathematical object \( S(X) \) such that if X and Y are homotopically equivalent \( S(X) \) and \( S(Y) \) are isomorphic.

My first idea, expressed in [29], was that it was better to consider some homological invariants to compute important properties of concurrent and distributed systems.

To begin with, we can make a very simple remark: instead of starting with a concurrent program semantics expressed in the form of precubical sets \( M = (M_i)_{i \in \mathbb{N}} \), we can use “bi-graded” precubical sets, i.e. sets

\[
N = (N_{p,q})_{p,q \in \mathbb{N}}
\]

The start boundary operators \( d^0_i \) are now going from \( N_{p,q} \) to \( N_{p+1,q} \) and the end boundary operators \( d^1_i \) are going from \( N_{p,q} \) to \( N_{p,q+1} \). The sets \( N_{p,q} \) are disjoint only if the shape it models does not contain any “directed” cycle.

The crucial observation is:

**Lemma 3.** Consider the following diagram of \( R \)-modules (\( R \) being an integral domain, for instance \( \mathbb{Z} \) or \( \mathbb{Z}/2\mathbb{Z} \) as in [29]):

\[
\begin{array}{c}
\mathcal{A}(N_{p,q}) \xrightarrow{\partial_0} \mathcal{A}(N_{p-1,q}) \\
\mathcal{A}(N) = \mathcal{A}(N_{p,q-1}) \\
\mathcal{A}(N_{p,q-1}) \xrightarrow{\partial_1} \mathcal{A}(N_{p,q})
\end{array}
\]

where \( \mathcal{A}(N_{p,q}) \) is the free \( R \)-module generated by \( N_{p,q} \) and\(^{15}\)

\[
\partial_0 = \sum_{i=0}^{i=p+q-1} (-1)^i d^0_i \\
\partial_1 = \sum_{i=0}^{i=p+q-1} (-1)^i d^1_i
\]

\(^{15}\)When the isomorphism is induced by a continuous map between CW-complexes for instance.

\(^{16}\)In order to remember the dimension in which they act, we will write sometimes \( \partial_0^{p+q} \) and \( \partial_1^{p+q} \).
It is a ("weak") bicomplex of modules, i.e. it satisfies the equalities: $\partial_0 \circ \partial_0 = 0$, $\partial_1 \circ \partial_1 = 0$, $\partial_0 \circ \partial_1 + \partial_1 \circ \partial_0 = 0$.

For instance, the automaton:

```
1
|    |
| b  |
|    |
| a  |
```

is represented by the bicomplex of $\mathbb{Z}$-modules,

```
(a) \oplus (b) \xrightarrow{\partial_0} (1)
(a') \oplus (b') \xrightarrow{\partial_1} (a) \oplus (b)
(\gamma) \xrightarrow{\partial_0} 0
```

with $\partial_0(a) = \partial_0(b) = 1$, $\partial_1(a) = \partial_1(b') = a$, $\partial_1(b) = \partial_1(a') = \beta$ and $\partial_1(a') = \partial_1(b') = \gamma$.

The bicomplexes (or double complexes of modules) are very important objects in homology, they have in fact very many interesting properties.

Let,

- $H_i(N, \partial_0) = \frac{\text{Ker } \partial_0}{\text{Im } \partial_1}$
- $H_i(N, \partial_1) = \frac{\text{Ker } \partial_1}{\text{Im } \partial_0}$

where $\text{Ker } f$ (respectively $\text{Im } f$) is the kernel (respectively the image) of the linear application $f$. These form the "horizontal" (respectively "vertical") homology groups, which enable to determine the branchings (respectively confluenes) of the automata.

In the case of our example, we find easily, $H_0(M, \partial_0) = \{a\}$, $H_0(M, \partial_1) = \{1\}$, $H_1(M, \partial_0) = (b - a)$, $H_1(M, \partial_1) = (b' - a')$, and the other homology groups are trivial. The generator $(b - a)$ of the horizontal homology group of dimension one expresses the fact that there is a non-deterministic choice between actions $a$ and $b$. The generator of the first vertical homology group $(b' - a')$ shows that there is a confluence between actions $a'$ and $b'$.

A typical branching in dimension two is for instance:

```
    / \
   /   \ 
  /     \ 
```

where the three faces are filled in.

The homology functors have nice computational properties (colimits, tensor product [Kun- neth formula], Mayer-Vietoris long exact sequence etc.) which enables to compute them inductively on the syntax of a parallel program, as was done in [25] with the process algebra CCS.
The “total homology” functor, defined as being the homology functor for boundary operator \( \partial_0 - \partial_1 \) gives also an homological theory for dipaths with fixed extremities modulo dihomotopy (see [25] and [26]) - in fact unfortunately, it is also an invariant of dipaths with fixed extremities modulo homotopy. This implies for instance that these functors cannot separate the two dipaths of example 13.

To correct this defect, it is necessary, first of all, to start with a better category of cubical sets. Most of this part has been based on earlier work by R. Brown and P. Higgins [8] and [6], and has been developed later by P. Gaucher. I will briefly come back to this piece of work in Section 8.2.

The theory I proposed in [26] defines homology groups in all dimensions indeed. The goal was to be able to distinguish between the shape of Figure 11 representing a 2-semaphore which is accessed by three processes, with a 3-semaphore in the same situation. The difference between a 1-semaphore and a 2-semaphore which two processes try to access can be noticed by examining the dihomotopy classes of dipaths. In the first case, there is necessarily a mutual exclusion which serializes the access to the shared object. In the second case, there is no serialization. When we use \( n \)-semaphores with \( n > 1 \), we cannot distinguish the difference of behaviour by looking at whether two consecutive actions (locks or unlocks) commute or not. The only way is by looking at the difference of behaviour when there are at least \( n + 1 \) consecutive actions (accesses to the semaphore) in general.

In order to spot this difference geometrically, we must examine hypersurfaces of dimension \( n \) modulo dihomotopy instead of just dipaths modulo dihomotopy. There again, we must be cautious to fix the extremities of these hypersurfaces.

The idea of [26] was simple and can be found in different forms in more recent work by S. Sokolowski [59] and P. Gaucher [20]. As can be seen on Figure 17 (at the left hand side, in dimension 2, at the right hand side in dimension 3) by taking two dihomotopic dipaths having the same source and target, we can consider the surfaces on which we can deform one of these paths into the other by a homotopy. We then say that two such surfaces are dihomotopic if there exists a dihomotopy between each of the dihomotopies that define these surfaces. In Figure 17 the two surfaces (one above the hole, the other below) are not continuously deformable one into the other, whereas they would be if the cube were entirely filled in. Homologically (as in [26]), we can look at the surfaces with fixed boundaries modulo the total homology. We briefly get back to this, homotopically, in Section 9.

7.2. Achronal cuts

In the previous section, we have tried to go from a classification problem of dipaths modulo dihomotopy to a simpler classification problem of dipaths modulo homology.

There is of course another natural idea [14], that of taking “instant snapshots” of the
dynamics of dipaths and observe their evolution in time. In fact, this is fairly close to methods used in fault-tolerant distributed systems theory [33].

**Definition 12.** Let \((U, \leq)\) be a partially ordered set \(A \subseteq V \subseteq U\) is called achronal if for all \(x, y \in V: x \leq y \Rightarrow x = y\) (similarly to the notion in [48]).

**Definition 13.** Let \((X, \leq)\) be a po-space.

1. \((X, \leq)\) is a parameterized po-space if there exists a dimap \(F : X \rightarrow \mathbb{R}\) such that \(X_t := F^{-1}(t)\) be achronal for all \(t \in \mathbb{R}\).
2. \(\leq\) is Euclidian, if there exists a finite number of dimaps \(f_i : X \rightarrow \mathbb{R}\) such that
   \[
   \forall x, y \in X: \quad x < y \Leftrightarrow \forall i: \quad f_i(x) \leq f_i(y);
   \]
   \[
   \exists i: \quad f_i(x) < f_i(y).
   \]
3. A local partial order on a topological space \(X\) is parameterized, respectively Euclidian if (one of its refinements) is a parameterized po-space, respectively, Euclidian po-space.

A Euclidian partial order is in fact a transcription of the natural componentwise ordering on an \(\mathbb{R}^n\) (like the progress graphs we saw at the beginning): given two points \(x = [x_1, \ldots, x_n], y = [y_1, \ldots, y_n] \in \mathbb{R}^n\),

\[
 x \leq y \Leftrightarrow \forall i: \quad x_i \leq y_i.
\]

In a parameterized po-space, we can always reparameterize the dipaths and the dihomotopies, such that, for any dipath \(p\) and any \(t \in I\), \(p(t) \in F^{-1}(t)\).

Let \(H : J \times I \rightarrow X\) be a well-parameterized dihomotopy between two well-parameterized dipaths \(a, a' : I \rightarrow X\). Then for all \(t \in I\), \(a(t)\) and \(a'(t)\) are in the same connected component of \(X_t\) (which is the “cut at instant \(t\) of \(X\)).

This gives us a means, by the study (with standard homotopy theory) of cuts, to determine the possible obstructions to dihomotopy, thus to find a subset of the possible schedules.

Unfortunately, this only gives us an approximation in general; let \(X\) be the subset \([0,3] \times [0,3] \times [1,2] \times [0,3] \subseteq \mathbb{R}^3\) with the natural partial order. There are two dihomotopy classes of dipaths going from \((0,0,0)\) to \((3,3,3)\), but the cuts induced by the “height function” \(F(x, y, z) = x + y + z\) are all connected.

Thus, to find all information about dihomotopy classes of dipaths, it is not enough to consider only one family of cuts. In fact, it seems that we need all possible families of cuts, in the general case of precubical sets. On the computer scientific side, this only means that some asynchronous systems have no global clock.

8. Refinements of the framework

8.1. The \(\omega\)-categorical point of view

The idea goes back to the article [49], and has been improved by P. Gaucher.

A \(\omega\)-category is a category with morphisms and compositions in all dimensions, somehow coherent with one another.

The idea for the modelisation of concurrency is that objects, or 0-morphisms, represent states of an HDA, the 1-morphisms represent all possible paths of executions (all dipaths), and the higher-dimensional morphisms represent the dihomotopies between morphisms of lesser dimension\(^\dagger\). In particular, the 2-morphisms represent the dihomotopies between paths of execution. The composition laws between higher-dimensional morphisms characterize the compositions of dihomotopies of higher dimension.

\(^\dagger\)Precisely those introduced in Section 7.1.
As V. Pratt already noticed in [49], the axioms of $\omega$-categories encode the composition properties of dipaths and of dihomotopies in an HDA. The interested reader can find the exploitation of these ideas in [19], [18] and [21].

We will give the formal definition of an $\omega$-category in three steps (see for instance [7], [62] and [61] for more details):

A 1-category is a pair $\langle A, (\ast, s, t) \rangle$ satisfying the following properties:

1. $A$ is a set,
2. $s$ and $t$ are functions from $A$ to $A$ called source and target respectively,
3. for all $x, y \in A$, $x * y$ is defined as soon as $tx = sy$,
4. $x * (y * z) = (x * y) * z$ as soon as the two members of the equality are well defined,
5. $sx * x = x * tx = x$, $s(x * y) = sx$ and $t(x * y) = ty$ (so $ssx = sx$ and $ttx = tx$).

A 2-category is a triple

$$\langle A, (\ast, s_0, t_0), (s_1, s_1, t_1) \rangle$$

such that,

1. The two pairs $\langle A, (\ast, s_0, t_0) \rangle$ and $\langle A, (\ast, s_1, t_1) \rangle$ are 1-categories,
2. $s_0s_1 = s_0t_1 = s_0$, $t_0s_1 = t_0t_1 = t_0$, and for $i \geq j$, $s_is_j = t_it_j = s_j$ and $s_it_j = t_it_j = t_j$ (globular axioms)
3. $(x *_0 y) *_1 (z *_0 t) = (x *_1 z) *_0 (y *_1 t)$ (Godement axiom)
4. if $i \neq j$, then $s_i(x *_j y) = s_ix *_j s_1y$ and $t_i(x *_j y) = t_ix *_j t_1y$.

A globular $\omega$-category $\mathcal{C}$ is composed of a set $A$ and of a family $(\ast_n, s_n, t_n)_{n \geq 0}$ such that

1. for all $n \geq 0$, $\langle A, (\ast_n, s_n, t_n) \rangle$ is a 1-category
2. for all $m, n \geq 0$ with $m < n$,
   $\langle A, (\ast_m, s_m, t_m), (\ast_n, s_n, t_n) \rangle$ is a 2-category
3. for all $x \in A$, there exists $n \geq 0$ such that $s_nx = t_nx = x$ (the smallest of these $n$ is called the dimension of $x$).

An $n$-dimensional element of $\mathcal{C}$ is called $n$-morphism. A 0-morphism is also called a state of $\mathcal{C}$, and a 1-morphism, an arrow. If $x$ is a morphism of an $\omega$-category $\mathcal{C}$, we call $s_n(x)$ the $n$-source of $x$ and $t_n(x)$ the $n$-target of $x$. The category of all $\omega$-categories is denoted by $\omega Cat$. The corresponding morphisms are called $\omega$-functors.

Now we can give, as in [20] some examples of $\omega$-categories coming from cubical sets. For instance, the automaton generated by a unique 2-transition is represented as the $\omega$-category of Figure 18$^{\text{18}}$.

---

$^{18}$The double arrow is a 2-morphism in the corresponding $\omega$-category. This is more generally due to a result in [1].
P. Gaucher in [21] and [20] uses the $\omega$-category generated by a cubical set to construct three homological theories, corresponding respectively to branchings, confluences and globes (or, computer-scientifically, mutual exclusions). This is constructed through suitable nerve functors. This is made possible because of the choice of a nice category of cubical sets first, and also because simplicial sets can be represented as $\omega$-categories as well (see [62]). In some ways, the branching and confluence nerves describe simplicially all achronal cuts of HDA, as hinted at in Section 7.2.

These constructions have a number of advantages over the ones of [25]:

- They are more discriminating (for instance, the “room with three barriers” example should be fully described by these homology theories)
- Concerning the branching and confluence homologies, they are not sensitive to subdivision.

We should mention two other points: $\omega$-categories seem to give the right structure to the “dihomotopy sets” or at least what should be the algebraic counterpart in the directed theory to the homotopy groups. The equivalence of categories between the category of cubical categories with connections and compositions and the category of $\omega$-categories (see [11]) is certainly a step in this way. Other papers by R. Brown and P. Higgins (see [8] and [6]) also pave the way toward van Kampen theorems in the directed theory, probably in weaker forms though (I develop this a little in Section 9).

### 8.2. Cubical sets

There exists several types of cubical sets. The first important remark is that, in the category of precubical sets, morphisms are somehow too rigid.

Mathematically, it is easily seen that morphisms respect the dimension of cells and length of dipaths (the number of segments they are composed of); this means in particular that the orthogonal projection of a filled-in (and also of a hollow) square onto one of its segments is not a morphism in this category.

Computer-scientifically, this implies that some important properties are not “natural”. We have to introduce some degeneracy operators; basically, what we need here is to be able to consider any $m$-transition or hypercube of dimension $m$ as a $n$-transition ($n \geq m$). In computer scientific terms, the degeneracy operators will allow us to consider any execution of $m$ busy processors as an execution of $n$ busy processors where $n - m$ processors are busy... doing nothing. In dimension one, this is directly connected to the notion of “idle transition” in transition systems theory, see [63] and [28].

**Definition 14.** A cubical set $K$ is a precubical set $(K, \partial^n)$ with degeneracy operators $\epsilon_i : K_{n-1} \to K_n$ ($0 \leq i \leq n - 1$) verifying the relations:

\[
\begin{align*}
\epsilon_i \epsilon_j &= \epsilon_{i+1} \epsilon_i, & (i \leq j) \\
\partial^n_i \epsilon_j &= \begin{cases} 
\epsilon_{j-1} \partial^n_i & (i < j) \\
\epsilon_j \partial^n_{i+1} & (i > j) \\
\text{Id} & (i = j)
\end{cases}
\end{align*}
\]

R. Brown et P. Higgins have added later on [8] other special degeneracy operators called connections:

**Definition 15.** $K$ is a cubical set with connections if it has also functions called connections
\(\Gamma^\alpha : K_{n-1} \to K_n\) \((0 \leq i \leq n - 2, \alpha = +, -)\) satisfying the relations:

\[
\begin{align*}
\Gamma_i^\alpha \Gamma_j^\beta & = \Gamma_{i+j}^\alpha \Gamma_i^\beta \quad (i \leq j) \\
\Gamma_i^\alpha \epsilon_j & = \begin{cases} 
\epsilon_j \Gamma_i^\alpha & (i > j) \\
\epsilon_j \Gamma_i^\beta & (i = j) \\
\epsilon_j = \epsilon_j \epsilon_j & (i = j)
\end{cases} \\
\delta_i^\alpha \Gamma_j^\beta & = \begin{cases} 
\Gamma_j^\beta \delta_i^\alpha & (i < j) \\
\Gamma_j^\beta \delta_i^{\alpha-1} & (i > j + 1)
\end{cases} \\
\delta_j^\alpha \Gamma_i^\beta & = \delta_j^{\alpha+1} \Gamma_i^\beta = 1d \\
\delta_i^\alpha \Gamma_j^\beta & = \delta_j^{\alpha+1} \Gamma_i^\beta = \epsilon_j \delta_i^\alpha
\end{align*}
\]

Now, we can define gluings of \(n\)-cubes or compositions, which are also necessary to the computer-scientific modelisation, if we want to be able to consider dipoles algebraically (which are gluings of \(n\)-cubes indeed):

**Definition 16.** \(K\) is a cubical set with connections and compositions if it is a cubical set with connections and it has also \(n\) operations of composition in each dimension \(n\), \(+j\) \((0 \leq j \leq n - 1)\), such that,

If \(a,b \in K_n\) then \(a + j b\) is well-defined if \(\delta_j^a b = \partial_j^a a\). When the terms in the following equalities are well-defined then we have:

\[
\begin{align*}
\partial_j^a (a + j b) & = \partial_j^a a \\
\partial_j^b (a + j b) & = \partial_j^b b \\
\partial_j^\alpha (a + j b) & = \begin{cases} 
\partial_j^\alpha a + j \epsilon_j b & (i < j) \\
\partial_j^\alpha a + j \epsilon_j b & (i > j)
\end{cases} \\
(a + i b) + j & = (a + j c) + (b + j d) \\
\epsilon_i (a + j b) & = \begin{cases} 
\epsilon_i a + j + 1 \epsilon_i b & (i \leq j) \\
\epsilon_i a + j + 1 \epsilon_i b & (i > j)
\end{cases} \\
\Gamma_i^\alpha (a + j b) & = \begin{cases} 
\Gamma_i^\alpha a + j + 1 \Gamma_i^\alpha b & (i < j) \\
\Gamma_i^\alpha a + j + 1 \Gamma_i^\alpha b & (i > j)
\end{cases} \\
\Gamma_j^\alpha (a + j b) & = (\Gamma_j^\alpha a + j + 1 \epsilon_j a) + j + 1 (\epsilon_j + 1 a + j \Gamma_j^\alpha b) \\
\Gamma_j^- (a + j b) & = (\Gamma_j^- a + j + 1 \epsilon_j b + j \Gamma_j^- b)
\end{align*}
\]

A \(\omega\)-cubical category is a cubical set with connections and compositions such that each \(+j\) gives a categorical structure to \(K_n\), with identities \(\epsilon_j y (y \in G_{n-1})\) and with the following conditions:

\[
\begin{align*}
\Gamma_i^+ + 1 \Gamma_i^- x & = \epsilon_{i+1} x \\
\Gamma_i^+ + 1 \Gamma_i^- x & = \epsilon_i x
\end{align*}
\]

I will not go further in this presentation, which would need a deeper study. One should note though that the category of \(\omega\)-cubical categories has been shown equivalent quite recently to (see [11]) the category of \(\omega\)-categories, whose use for the modelisation of concurrent processes has been finally put together by P. Gaucher [21].

In order to link the \(\omega\)-categorical formulation of P. Gaucher, we have been compelled to restrict the category of local po-spaces to consider. In fact, as in standard algebraic topology, the category of topological spaces is far too big and contains far too pathological elements to be the right object of study. In general, we restrict ourselves to topological spaces which have the homotopy type of a CW-complex [36]. P. Gaucher and myself have introduced in [22] a particular kind of CW-complex, which we called globular CW-complex, and which is
essentially a CW-complex which $n$-cells are all directed. A $n$-cell is homeomorphic to $I^{n-1}$ quotiented by relations
\[(k, x_1, \ldots, x_{n-1}) = (k, y_1, \ldots, y_{n-1})\]
\[(k = 0, 1).\] The advantages of this category of local po-spaces are:
- It contains only the classical homotopy types,
- It allows to construct the homology theories of P. Gaucher (defined originally on $\omega$-categories) in a topological framework,
- It permits to define a notion of dihomotopy equivalence (which refine the usual homotopy equivalence, i.e. the homotopy types). We hope that for the globular CW-complexes, this should coincide with a notion of weak dihomotopy equivalence (as in the classical case).

8.3. Domain theory

There are links between the theories briefly described before and domain theory\(^{19}\) (for instance as developed in chapter VII of [35]), or of other older topological considerations (the book [47] for instance).

L. Nachbin in [47] has studied some particular kind of topological spaces, the so-called compact order-Hausdorff topological spaces. In fact, this is nothing but compact po-spaces (for instance, finite progress graphs). One of the very interesting results in the theory is that there is an adjunction between these compact po-spaces and another type of topological space (no order there!), the stably-compact spaces.

We will write $PO$ for the category of compact po-spaces. We now define stably-compact spaces:

**Definition 17.** A stably-compact space is a set $X$ together with a topology $\tau$ on $X$ such that there exists another topology on $X$, $\tau^*$ sur $X$ satisfying the following conditions:
- $\tau \cup \tau^*$ is compact,
- for all $x \neq y$ in $X$, there exists an element $O \in \tau$, an element $O^* \in \tau^*$ such that $x \in O$, $y \in O^*$ and $O \cap O^* = \emptyset$.

In some ways, $(X, \tau)$ is compact Hausdorff with the help of topology $\tau^*$.

We write $SK$ for the category whose objects are the stably-compact spaces and whose morphisms are continuous functions.

**Proposition 2.** Let $(X, \tau, \leq)$ be a compact po-space ($\tau$ is the topology, $\leq$ is the partial order). We build $(X', \tau')$ a topological space from $(X, \tau, \leq)$ as follows:
- $X' = X$,
- $\tau'$, the set of opens, is composed of elements $U$ of $\tau$ which are such that $\forall x \in U$, $\forall y \geq x$, $y \in U$ ("upper sets").

Then $(X', \tau')$ is a stably-compact space.

**Sketch of proof.** It is a direct consequence of the local convexity theorem (see [47] or [35], Theorem 1.4, Chapter VII, page 272) which states that sets of the form $U \cap V$, where $U$ is an upper open set and $V$ is a lower open set, form a base for the topology of $X$. Thus it suffices to take for $\tau'$ (required by definition 17) the set of lower open sets. The axiom of "weak separation" of definition 17 is exactly corollary 1.2, Chapter VII, page 271 of [35].

Of course, $(X', \tau')$ is in general not at all Hausdorff.

Dimaps between compact po-spaces are naturally mapped under this transformation onto continuous maps between there stably-compact counterparts. This defines a functor $a$ from $PO$ to $SK$ which has a right-adjoint we will briefly describe:

\(^{19}\)This is part of a talk delivered by the author in Dagstuhl seminar 00231/1 “Topology in Computer Science”.
**Definition 18.** Let $(X, \tau)$ be a topological space, $A$ the set of its “compact opens” (i.e. the opens of $A$ whose adherence is compact in $X$), $A^*$ the set of complements (in $\phi(X)$) of elements of $A$. The “patch” topology on $X$ is the topology $\kappa(\tau)$ generated by the base $C = \{U \cap V \mid U \in A, V \in A^*\}$.

**Proposition 3.** Let $(Y, \sigma)$ be a stably-compact space. We can associate with it, a structure
\[(X, \tau, \preceq) = \gamma(Y, \sigma)\]
with,
- $X = Y$,
- $\tau = \kappa(\sigma)$,
- for all $x, y \in X$, $x \preceq y$ if for all $U \in \sigma$ with $x \in U$, $y$ is also in $U$.

Then $(X, \tau, \preceq)$ is a compact po-space.

Moreover, $\gamma$ define a functor from $\text{SK}$ to $\text{PO}$ (transforming continuous functions $\text{SK}$ into dimaps of $\text{PO}$).

Consider now a dihomotopy $H$ between two dipaths $f$ and $g$ with the same source and target in $X$. It is simply a dimap from $I \times I$ to $X$ such that $H(0, \cdot) = f$ and $H(1, \cdot) = g$.

Therefore $\alpha(H)$ is a continuous function from $\alpha(f)$ to $\alpha(g)$, which are themselves continuous functions from $\alpha(I \times I)$ to $\alpha(X)$. But $\alpha(I \times I) = I \times \alpha(I)$, because $\alpha(I) = I$ (all open set of $I$ is upper). We conclude that $\alpha(H)$ is a (classical) homotopy between $\alpha(f)$ and $\alpha(g)$. We thus have proved a simple obstruction criterion to dihomotopy:

**Proposition 4.** Let $f$ and $g$ be two dipaths in the compact po-space $X$. Then $\alpha(f)$ and $\alpha(g)$ are not homotopic (in the classical sense) implies that $f$ and $g$ are not dihomotopic.

Next question is, do we have a reciprocal to this? Let $H'$ be a homotopy between $\alpha(f)$ and $\alpha(g)$. At what condition is there a dihomotopy $H$ between $f$ and $g$? Do we have $\alpha(H) = H'$?

The first natural idea to characterize the homotopies between $\alpha(f)$ and $\alpha(g)$, given $f$ and $g$ two dipaths in a compact po-space $X$, is to study the group homomorphisms between the homotopy groups of $\alpha(I)$ to the homotopy groups of $\alpha(X)$. First of all, we notice that $\alpha(I)$ is a compact and connected topological space.

As $\alpha(I)$ is connected, we can define its fundamental group $\pi_1(\alpha(I))$ by choosing any basepoint, for instance 0. Unfortunately, the study of continuous functions from $I$ to $\alpha(I)$, reveals that they are the lower semi-continuous maps from $I$ to $I$ and that $\alpha(I)$ is a simply-connected topological space. This means that there is no interesting group homomorphism to look for there from $\pi_1(\alpha(I))$ to $\pi_1(\alpha(X))$.

Even worse, we can show that we can deform in a continuous manner (for the topology of $\alpha(I \times I)$ any maximal continuous path (for the topology of $I \times I$) in any other one, by going through discontinuous ones for the topology of $I \times I$. This entails that the homotopy between functions from $\alpha(I)$ to $\alpha(I \times I)$ does not even enable us to see the presence of a hole in the surface $I \times I$!

In fact, the adjunction between stably-compact spaces and compact po-spaces can be changed into an equivalence of categories. The way to do this is to consider the subcategory of stably-compact spaces where we impose that the morphisms be “perfect”. The question here is, can we develop a practical homotopy theory on stably-compact spaces, with morphisms somewhere between continuous functions and perfect maps, that would give us enough information on dihomotopy? The other important question is: is there a similar counterpart to the local pospaces? I suspect that there are also strong links with work by M. Grandis (see for instance [31]).
9. Some Desired Properties

Continuing our tour of classical notions that would be useful for concurrency theory, it is natural to ask ourselves what should be the counterpart of homotopy exact sequences etc. in the directed case.

9.1. van Kampen

One of the very useful theorems in the classical theory is van Kampen’s theorem which relates the fundamental group of a space which is the (non necessarily disjoint) union of two subspaces with the fundamental groups of these two subspaces, under some mild hypotheses. This would have very interesting application for the “modular” or “compositional” analysis of concurrent systems.

Theorem 1. Let $X$ be a local po-space, $X_1$ and $X_2$ two local po-spaces such that,

- $X = \hat{X}_1 \cup \hat{X}_2$, 
- All continuous paths (not dipaths in general) in $\hat{X}_1 \cap \hat{X}_2$ are concatenations of a finite number of dipaths and anti-dipaths (“zig-zag paths”).

Let $j_1 : X_1 \cap X_2 \to X_1$ (respectively $j_2 : X_1 \cap X_2 \to X_2$) and $i_1 : X_1 \to X$ (respectively $i_2 : X_2 \to X$) be the natural inclusion morphisms. Then the following diagram,

$$
\begin{array}{ccc}
C(X_1 \cap X_2) & \xrightarrow{C(j_1)} & C(X_1) \\
\downarrow{C(j_2)} & & \downarrow{C(i_1)} \\
C(X_2) & \xrightarrow{C(i_2)} & C(X)
\end{array}
$$

is co-cartesian in the category of categories.

Proof. Consider first the following commutative diagram,

$$
\begin{array}{ccc}
C(X_1 \cap X_2) & \xrightarrow{C(j_1)} & C(X_1) \\
\downarrow{C(j_2)} & & \downarrow{C(i_1)} \\
C(X_2) & \xrightarrow{C(i_2)} & C(X) \\
& & \downarrow{f_1} \\
& & G
\end{array}
$$

where $G$ is a category. The question is whether one can complete this commutative diagram.
into the following one,

\[
\begin{align*}
\mathcal{C}(X_1 \cap X_2) & \xrightarrow{\mathcal{C}(j_1)} \mathcal{C}(X_1) \\
\mathcal{C}(j_2) & \xrightarrow{\mathcal{C}(j_2)} \mathcal{C}(X_2) \\
\mathcal{C}(X_2) & \xrightarrow{\mathcal{C}(j_2)} \mathcal{C}(X) \\
f & \xrightarrow{f} \mathcal{C}(X) \\
 & \xrightarrow{f_2} \mathcal{C}(X_1) \\
 & \xrightarrow{f_1} \mathcal{C}(X_1) \\
G &
\end{align*}
\]

Looking at the diagram on objects of these categories, everything boils down to the existence of a push-out in the category of sets, which of course holds. On the objects, define \( f \) to be (for \( x \in X \)),

\[
f(x) = \begin{cases} 
  f_1(x) & \text{if } x \in X_1 \\
  f_2(x) & \text{if } x \in X_2
\end{cases}
\]

Now, let \( u : \bar{I} \to X \) be a dimap, \( u^{-1}(\bar{X}_1) \), \( u^{-1}(\bar{X}_2) \) form an open covering of \( \bar{I} \). Let \( \delta \) be its Lebesgue number, and let \( 0 = t_0 < t_1 < \cdots < t_{k+1} = 1 \) be a subdivision of the unit interval such that we have \( |t_{\alpha+1} - t_\alpha| < \delta \) for all \( \alpha \in \{0, \cdots, k\} \). By definition of Lebesgue numbers, for all \( \alpha \in \{0, \cdots, k\} \) we have an integer \( \epsilon_\alpha \) (being 1 or 2) such that \( u([t_\alpha, t_{\alpha+1}]) \subseteq X_{\epsilon_\alpha} \).

By reparameterisation, we consider \( u_\alpha : [t_\alpha, t_{\alpha+1}] \to X \). Now by definition of the composition in \( \mathcal{C}(X) \) we have,

\[
[u] = \mathcal{C}(i_{\epsilon_\alpha})[u_\epsilon_\alpha] \circ \cdots \circ \mathcal{C}(i_{\epsilon_\alpha})[u_0]
\]

Therefore the morphim \( f \) we are looking, if it exists, must necessarily satisfy:

\[
f([u]) = f_{\epsilon_\alpha}[u_\epsilon_\alpha] \circ \cdots \circ f_0[u_0]
\]

Let us define \( f \) this way and check that this is a correct definition.

First, we have to prove that this definition does not depend on the subdivision \( 0 = t_0 < t_1 < \cdots < t_{k+1} = 1 \) that has been chosen.

Let \( 0 = t'_0 < t'_1 < \cdots < t'_{k+1} = 1 \) be another possible choice of subdivision. Consider the new one, \( 0 < t''_0 < t''_1 < \cdots < t''_{k'} = 1 \) with \( t''_0 < t_0 < \cdots < t_{k+1} \). Let \( t''_0 < t''_1 < t''_2 < \cdots < t''_{k'} = 1 \) be another possible choice of subdivision with \( 0 < t''_0 < 1 \).

So, \( t''_0 < t''_1 < t''_2 < \cdots < t''_{k'} = 1 \) is a subdivision of \( 0 < t_0 < \cdots < t_{k+1} \). Now we prove that \( f_{\epsilon_\alpha}(u_\beta) = f_{\epsilon_\alpha}(u_\beta) \circ \cdots \circ f_{\epsilon_\alpha}(u_\beta) \) with their dihomotopic counterparts (by a reparameterization) \( f_{\epsilon_\alpha}(u_\beta) : [0, 1] \to X_{\epsilon_\alpha} \). Notice that \( u''_\alpha \) has its values in \( X_{\epsilon_\alpha} \). But as \( f \circ j_0 = f \circ j_1 \), \( f_{\epsilon_\alpha}(u''_\alpha) = f_{\epsilon_\alpha}(u''_\alpha) \circ \cdots \circ f_{\epsilon_\alpha}(u''_\alpha) \) is obviously equal to \( f_{\epsilon_\alpha}(u'_\alpha) \). This proves (by induction) that the definition of \( f \) does not depend on the subdivision.

Now we prove that \( f([u]) \) only depends on the class of \( u \) and not of \( u \). Let \( h : \bar{I} \times \bar{I} \to X \) be a dihomotopy linking dipath \( u \) to dipath \( v \). Let \( \delta \) be the Lebesgue number of the covering of \( \bar{I} \) by \( h^{-1}(X_1) \) and \( h^{-1}(X_2) \). Consider subdivisions

\[
0 = t_0 < \cdots < t_{k+1} = 1 \\
0 = t_0 < \cdots < t_{k} = 1
\]

so that all squares \( [t_i, t_{i+1}] \times [t_j, t_{j+1}] \) have diameter less than \( \delta \). Further subdivide so that each \( \tau \to h(t_i, \tau) \) for \( \tau \in [t_j, t_{j+1}] \) is a dipath or an anti-dipath (a dipath with the reverse
ordering on \( \tilde{f} \). This is possible since these are continuous paths which can be decomposed into a finite number of dipaths or anti-dipaths by hypothesis on \( X \).

Let \( u'(t) = h(t, \tau) \). If we prove \( f([u]) = f([u']) \) then it follows by an easy induction that \( f([u]) = f([v]) \). Let \( P_\alpha \) be the rectangle between lines \( t = t_\alpha, t = t_{\alpha + 1}, \tau = 0 \) and \( \tau = \tau \). Call,

\[
\begin{align*}
  u_\alpha : [t_\alpha, t_{\alpha + 1}] & \to X_{\tau}\nn_\alpha : [t_\alpha, t_{\alpha + 1}] & \to X_{\tau}\nw_\alpha : [0, \tau] & \to X_{\tau}
\end{align*}
\]

the applications \( u_\alpha(t) = h(t, 0) \), \( n_\alpha(t) = h(t, \tau) \) and \( w_\alpha(\tau) = h(t_\alpha, \tau) \). We use the same names for any reparameterization of these applications from \([0, 1]\) to \( X_{\tau}\). Of course \( u_0 = [w_0] = Id \) and \( w_k + 1 = [w_{k+1}] = Id \).

We now notice the following in \( C(X) \) (since we have dihomotopies \( h \) on each \( P_\alpha \)):

- Suppose \( w_\alpha \) is a dipath,
  - suppose \( w_{\alpha + 1} \) is a dipath, then it is easy to see that \( [u_\alpha'] \circ [w_\alpha] = [w_{\alpha + 1}] \circ [u_\alpha] \)

\[
h'(x, t) = \begin{cases} 
  h(0, 3x(1 - t)) & \text{if } x < \frac{1}{3} \text{ is a dihomotopy} \\
  h(3x - 1, 1 - t) & \text{if } \frac{1}{3} \leq x \leq \frac{2}{3} \text{ between } u_\alpha' \circ w_\alpha \text{ and } w_{\alpha + 1} \circ u_\alpha \\
  h(1, 3(x - 1)(t + 1)) & \text{if } \frac{2}{3} \leq x \leq 1 \end{cases}
\]

- suppose \( w_{\alpha + 1} \) is an anti-dipath, then \( w_{\alpha + 1}^{-1} \) defined as being \( w_{\alpha + 1}^{-1}(t) = w_{\alpha + 1}(1 - t) \), for \( t \in [0, 1] \), is a dipath and \( [u_\alpha'] \circ [w_\alpha] = [w_{\alpha + 1}^{-1}] \circ [u_\alpha] \).
  - consider dihomotopy \( h'(x, t) = \begin{cases} 
  h(0, 3x(1 - t)) & \text{if } x < \frac{1}{3} \\
  h(3x - 1, 1 - t) & \text{if } \frac{1}{3} \leq x \leq \frac{2}{3} \text{ if } \frac{2}{3} \leq x \leq 1 \end{cases} \)

- Similarly, suppose \( w_\alpha \) is an anti-dipath,
  - suppose \( w_{\alpha + 1} \) is a dipath, then \( [u_\alpha] \circ [w_{\alpha + 1}^{-1}] = [w_{\alpha + 1}] \circ [u_\alpha] \).
  - assume \( w_{\alpha + 1} \) is an anti-dipath, then \( [u_\alpha] \circ [w_{\alpha + 1}^{-1}] = [w_{\alpha + 1}^{-1}] \circ [u_\alpha] \).

Consider now:

\[
\begin{align*}
f([u]) &= f_*([u_k]) \circ \cdots \circ f_*([u_0]) \\
f([u']) &= f_*([u_k']) \circ \cdots \circ f_*([u_0'])
\end{align*}
\]

We also have,

\[
\begin{align*}
f([u]) &= f_*([w_{k+1}]) \circ f_*([u_k]) \circ \cdots \circ f_*([u_0]) \\
f([u']) &= f_*([u_k']) \circ \cdots \circ f_*([u_0']) \circ f_*([w_0])
\end{align*}
\]

We prove by induction on \( \alpha \) that,

\[
f_*([u_{\alpha + 1}']) \circ \cdots \circ f_*([u_k']) \circ f_*([w_0]) = f_*([w_{\alpha + 1}]) \circ f_*([u_{\alpha - 1}]) \circ \cdots \circ f_*([u_0])
\]

with \( j_\alpha \) being \(-1\) if \( w_\alpha \) is an anti-dipath, \(+1\) otherwise.

This is a direct consequence of the remark above.

As such, this theorem is not quite the one we would like to have for applications (since on the intersection in general, the hypothesis of the theorem does not hold). In fact, this restricted form of van Kampen, works much better in the combinatorial case. In order to do this, we have to define an analogous to the fundamental category, the edge-path category. We sketch the construction in the easier case of precubical sets, but this can be done as well for cubical sets.

A dipath in a precubical set \( (M, d^0, d^1) \) is a finite sequence \( (k \geq 1) \)

\[
p = (p_1, \cdots, p_k)
\]

where all \( p_i \) are 1 dimensional such that \( d^0_1(p_i) = d^0_1(p_{i+1}) \) or the empty sequence \( \emptyset \).

The initial state of a non-empty dipath \( p \) is \( s(p) = d^0_1(p_1) \) and its final state is \( t(p) = d^1_0(p_k) \).
The composition $\ast$ (or concatenation) of dipaths is as follows. Let $p = (p_1, \cdots, p_k)$ and let $q = (q_1, \cdots, q_l)$ be two non-empty dipaths such that the initial state of $q$ is the final state of $p$. Then $p \ast q = (p_1, \cdots, p_k, q_1, \cdots, q_l)$ (is a dipath indeed).

This composition is associative as has neutral element the empty dipath. Let $[p]$ denote the contiguity class of dipath $p$. Concatenation induces an operation, still denoted by $\ast$, with $[p] \ast [q] = [p \ast q]$.

Define the edge-path category $\mathcal{E}(M)$ of $M$ as follows. Its objects are elements of $M_0$. Its morphisms from $a \in M_0$ to $b \in M_0$ are contiguity classes of dipaths from $a$ to $b$, or the empty dipath. The composition between morphisms is $\ast$ (we write now $[q] \circ [p] = [p] \ast [q]$). This forms a category.

Let now $f : M \to N$ be a morphism of precubical sets. It is easy to see that $f$ maps dipaths of $M$ to dipaths of $N$ and that if $p$ and $q$ are two contiguous dipaths in $M$, then $f(p)$ and $f(q)$ are two contiguous dipaths in $N$. Similarly, $f(p \ast q) = f(p) \ast f(q)$. Therefore $f$ induces a transformation $\mathcal{E}(f)$ transforming objects of $\mathcal{E}(M)$ into objects of $\mathcal{E}(N)$, morphisms of $\mathcal{E}(M)$ into $\mathcal{E}(N)$, respecting composition. Hence $\mathcal{E}$ is a functor from the category of precubical sets to the category of categories.

**Proposition 5.** Let $X = X_1 \cup X_2$ be a finite precubical set, union of two precubical sets. Call $j_1 : X_1 \cap X_2 \to X_1$ (respectively $j_2 : X_1 \cap X_2 \to X_2$) and $i_1 : X_1 \to X$ (respectively $i_2 : X_2 \to X$) the natural inclusion morphisms. Then the following diagram

$$
\begin{array}{ccc}
\mathcal{E}(X_1 \cap X_2) & \xrightarrow{\mathcal{E}(j_1)} & \mathcal{E}(X_1) \\
\downarrow & & \downarrow \\
\mathcal{E}(X_2) & \xleftarrow{\mathcal{E}(i_2)} & \mathcal{E}(X) \\
\end{array}
$$

is co-cartesian.

**Proof.** Same as above, but simpler in that, there is no Lebesgue number argument (replaced by the finiteness argument), and all un-directed paths in $X_1 \cap X_2$ are now always compositions of a finite number of di- and anti-di- paths (no achronal part!).

9.2. Higher-order homotopies

We are going to define inductively the notion of higher-order dihomotopy. We call dihomotopy of order 0 any dipath from. A dihomotopy of order 1 is any dihomotopy between two dipaths with the same ends. Now, suppose we have defined dihomotopies of order up to $n$ ($n \geq 1$) between dihomotopies of order $n - 1$ with end fixed (gathered in a set $P_n(X)$):

**Definition 19.** A dihomotopy of order $n + 1$ ($n \geq 1$) between dihomotopies $H, G$ of order $n$ with equal ends is a dimap $A : I^{n+1} \to X$ such that for all $x \in I^n \times I$, $A(x, 0) = H(x)$ and $A(x, 1) = G(x)$. The source of $A$ is $s_n(A) = H$ and its target is $t_n(A) = G$. The set of such dihomo- topies is noted $P_n(X)$.

We can define compositions on the sets $P_n(X)$ ($n \geq 1$) as follows; $*_{n-1} : P_n(X) \times P_n(X) \to P_n(X)$ is defined for $(f, g)$ such that $t_n(f) = s_n(g)$:

$$(f *_{n-1} g)(x_0, \cdots, x_n) = \begin{cases} 0 \leq x_n \leq \frac{1}{2} & f(x_0, \cdots, x_{n-1}, 2x_n) \\ \frac{1}{2} \leq x_n \leq 1 & g(x_0, \cdots, x_{n-1}, 2x_n - 1) \end{cases}$$

This naturally gives groupoids $\pi_n$ for all dihomotopies of higher-dimension $n \geq 2$ modulo dihomotopies of dimension $n + 1$. Of course, if we take a base path and look at loops around this base paths we have $\pi_2$ becoming a group, and $\pi_n$ ($n \geq 3$) becoming abelian groups.
The question one has to solve now is: do we have exact sequences such as the homotopy exact sequence of a pair? Do we have interesting spectral sequences? Does all this come from a closed-model structure?

10. Conclusion and Future Work

The dihomotopy theory and applications to concurrency gradually developed together, but there is much left to do as I tried to show throughout this text.

We would like to consider also other potential applications. In fact, there are a certain number of other “geometric” theories which apply to computation models. For instance, there are algebraic topological considerations in linear logics (see for instance [43], [44], [3] and [4]) and in modal logics [36] which might be related to the subject of this paper. More generally, linear logic [23] has a strong geometric flavour and can be understood as a logic dealing with resources; I believe that some fruitful cross-fertilizations should bring new results in the years to come.

References


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