Reachability in Two-Dimensional Vector Addition Systems with States is PSPACE-complete

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Abstract—Known to be decidable since 1981, there still remains a huge gap between the best known lower and upper bounds for the reachability problem for vector addition systems with states (VASS). Here the problem is shown PSPACE-complete in the two-dimensional case, vastly improving on the doubly exponential time bound established in 1986 by Howell, Rosier, Huynh and Yen. Coverability and boundedness for two-dimensional VASS are also shown PSPACE-complete, and reachability in two-dimensional VASS and in integer VASS under unary encoding are considered.

I. INTRODUCTION

Petri nets have a long history. Since their introduction by Petri in 1962, thousands of papers on Petri nets have been published. Nowadays, Petri nets find a variety of applications, ranging, for instance, from modeling of biological, chemical and business processes to the formal verification of concurrent programs, see e.g. [1], [4], [8], [21], [27]. For the analysis of their algorithmic properties, Petri nets are often equivalently viewed as vector addition systems with states (VASS), and we will adopt this view throughout this paper. A VASS comprises a finite-state controller with a finite number of counters ranging over the natural numbers. The number of counters is usually referred to as the dimension of the VASS, and we write d-VASS to denote VASS in dimension d. When taking a transition, a VASS can add or subtract an integer from a counter, provided that the resulting counter values are greater than or equal to zero; otherwise the transition is blocked. A configuration of a VASS is a tuple consisting of a control state and an assignment of natural numbers to the counters. The central decision problem for VASS is reachability: given two configurations, is there a path connecting them in the infinite graph induced by the VASS?

Resolving decidability of the VASS reachability problem required tremendous effort, extending until 1981. This was achieved by Mayr [18], who built upon an earlier partial proof by Sacerdote and Tenney [23]. Mayr’s argument was then polished and simplified by Kosaraju [11] in 1982, and Kosaraju’s argument was in turn simplified ten years later by Lambert [12]. More recently, beginning in 2009, Leroux began developing a fundamentally different approach to deciding the VASS reachability problem [15], [16]. Finally, at the time of writing of this paper, Leroux and Schmitz could establish the first explicit upper bound for VASS reachability and show that it can be decided in $F_{\omega^3}$ [13].

Milestones in the work on the complexity of the VASS reachability problem include Lipton’s 1976 proof that the problem, regardless of the choice of encoding for numbers but without fixed dimension, is EXPSPACE-hard [17]. Yet our knowledge of the situation for any fixed dimension $d$ is vastly lacking. For 1-VASS, reachability under unary encoding is easily seen to be NL-complete: the hardness is inherited from graph reachability and the upper bound follows from a simple pumping argument. Under binary encoding, 1-VASS reachability is known to be NP-complete [5]. As a substantial contribution towards showing decidability of the general problem, Hopcroft and Pansiot in 1979 showed the two-dimensional case decidable [9]. At the core of their proof lies an intricate algorithm that implicitly exploits the fact that the reachability set of a 2-VASS is semi-linear. Exhibiting a 3-VASS with a reachability set that is not semi-linear, Hopcroft and Pansiot could show that their method breaks down for $d$-VASS for any $d$ greater than 2. Further complexity aspects were left unanswered in [9]. In 1986, Howell, Rosier, Huynh and Yen [10] observed that Hopcroft and Pansiot’s algorithm runs in nondeterministic doubly-exponential time, under both unary and binary encoding. They then managed to improve this bound from nondeterministic to deterministic doubly-exponential time, and to identify a 2-VASS family on which Hopcroft and Pansiot’s algorithm requires this much. To summarize the state of the art today, 2-VASS reachability in 2-EXPSPACE has stood since 1986, with its NL-hardness and NP-hardness depending on number encodings. For any $d$ greater than 2, reachability is in $F_{\omega^d}$ [13].

The main contribution of this paper is to show that reachability in 2-VASS is PSPACE-complete when numbers are encoded in binary. The PSPACE lower bound follows as an easy consequence of a recent result by Fearnley and Jurdiński who showed PSPACE-completeness of reachability in bounded one-counter automata [3]. Our PSPACE upper
bound is obtained from showing that the length of a run witnessing reachability can be exponentially bounded in the size of the input, and consequently the existence of such a run can be decided by a PSPACE-algorithm. The difficult and main part of this paper is, of course, to establish the exponential upper bound on the length of witnessing runs. Our starting point is a careful analysis of an argument developed by Leroux and Sutre in [14] for the purpose of showing that reachability relations of 2-VASS can be captured by bounded languages, *i.e.*, speaking in the terminology of [14], 2-VASS can be flattened. More precisely, this means that for any 2-VASS there is a finite set $S$ of regular languages over the set of transitions, viewed as an alphabet, each of the form $u_0 \cdots u_n u_{n+1}$ such that for any two configurations reachable from one another there exists a witnessing run in the language defined by $S$. The paper of Leroux and Sutre reports that from any 2-VASS it is possible to construct such a bounded language; it has however not appeared as a fully refereed publication and omits some proof details. Thus, while we follow closely the proof strategy presented in [14], we provide a proof that 2-VASS can be flattened by small bounded languages. In doing so we develop new arguments setting the stage for the much deeper analysis of our constructions required for the purpose of establishing a PSPACE upper bound. In summary, we contribute:

1) a PSPACE-completeness proof for 2-VASS reachability,

2) a proof that 2-VASS can be flattened by bounded languages that have small presentations, and

3) remarks that reachability in 2-VASS with numbers encoded in unary is NL-hard and in NP.

Section II below fixes notation. Section III gives an overview of our main results. Section IV proves our main technical results for the sake of readability, we write configurations $p = (q, z_1, \ldots, z_d)$ and $(q, z)$ as $q(z_1, \ldots, z_d)$ and $q(z)$, respectively.

Due to space constraints, the proofs of some lemmas are only sketched, and full proofs can be found in the extended version of this paper.

II. Preliminaries

General notation. By $\mathbb{N} \triangleq \{0, 1, 2, \ldots\}$, $-\mathbb{N} \triangleq \{0, -1, -2, \ldots\}$ and $\mathbb{Z}$ we denote the sets of non-negative integers, non-positive integers and integers, respectively. By $\mathbb{Q}$ and $\mathbb{Q}_{\geq 0}$ we denote the set of rationals and non-negative rationals, respectively. For any $i, j \in \mathbb{Z}$, we define $[i, j] \triangleq \{i, i+1, \ldots, j\}$. For each $k \in \mathbb{Z}$ we write $[k, \infty)$ to denote $\{z \in \mathbb{Z} : z \geq k\}$. A *quadrant* is one of the four sets $\mathbb{N}^2$, $\mathbb{N} \times \mathbb{N}$, $\mathbb{N} \times -\mathbb{N}$ and $-\mathbb{N} \times -\mathbb{N}$. Given two vectors $u = (u_1, \ldots, u_d)$, $v = (v_1, \ldots, v_d) \in \mathbb{Z}^d$, we denote by $u + v \triangleq (u_1 + v_1, \ldots, u_d + v_d)$ their component-wise sum. Given two sets $U, V \subseteq \mathbb{Z}^d$, we let $U + V \triangleq \{u + v : u \in U, v \in V\}$. The *norm* of a vector $u = (u_1, \ldots, u_d)$ is defined as $\|u\| \triangleq \max\{|u_i| : i \in [1, d]\}$. The norm of a matrix $A = (a_{ij}) \in \mathbb{Z}^{m \times n}$ is defined as $\|A\| \triangleq n \cdot \max\{|a_{ij}| : i \in [1, m], j \in [1, n]\}$. For any word $w = a_1 \cdots a_n \in \Sigma^*$ over some alphabet $\Sigma$, $w[i, j]$ denotes $a_ia_{i+1}\cdots a_j$ for all $i, j \in [1, n]$.

**Graphs, Parikh Images and Linear Path Schemes.** For each set $\Sigma$, a *$\Sigma$-labeled directed graph* is a pair $G = (U, E)$, where $U$ is a set of vertices and $E \subseteq U \times \Sigma \times U$ is a set of edges. We say $G$ is finite if $U$ and $E$ are finite. Let $\pi = (u_1, u_2, \cdots, u_k, u_{k+1}) \in E^k$. The Parikh image $Parikh_{\pi}$ of $\pi$ is the mapping from $\Sigma$ to $\mathbb{N}$ such that $Parikh_{\pi}(a) = |\{(i \in [1, k] : a_i = a)\}|$ for each $a \in \Sigma$. If $X \subseteq E^*$, then $Parikh_X$ denotes the set of Parikh images of $X$, *i.e.* $Parikh_X = \{Parikh_{\pi} : \pi \in X\}$. We say $\pi$ is a path (from $u_1$ to $u_k$) if $u_i = u_{i+1}$ for all $i \in [1, k-1]$. A path $\pi$ is a cycle if $k \geq 1$ and $u_1 = u_k$, and cycle-free if no infix of $\pi$ is a cycle. A cycle $\pi$ is called simple if $\pi$ is the only infix of $\pi$ that is a cycle. A *linear path scheme* (from $u \in U$ to $u' \in U$) is a regular expression (whose language will be referred to implicitly) of the form

$$
\rho = \alpha_0 \beta_1^* \alpha_1 \cdots \beta_k^* \alpha_k,$$

where $\alpha_0 \beta_1 \cdots \beta_k \alpha_k$ is a path (from $u$ to $u'$) and each $\beta_i$ is a cycle. We define its length as $|\rho| \triangleq |\alpha_0 \beta_1 \cdots \beta_k \alpha_k|$ and its s-length as $|\rho|_{s} \triangleq k$. We call $\beta_1, \ldots, \beta_k$ the cycles of $\rho$. Note that every path is a linear path scheme by taking $k = 0$. The general structure of a linear path scheme is illustrated in Fig. 1.

**Vector Addition Systems with States.** A vector addition system with states (VASS) in dimension $d$ ($d$-VASS for short) is a finite $\mathbb{Z}^d$-labeled directed graph $V = (Q, T)$, where $Q$ will be referred to as the states of $V$, and where $T$ will be referred to as transitions of $V$. The size of $V$ is defined as $|V| \triangleq |Q| + |T| \cdot d \cdot \log_2 |T||$, where $|T|$ denotes the absolute value of the largest number that appears in $T$, *i.e.* $|T| \triangleq \max\{|\|z\|| : (p, z, q) \in T\}$. We say that $V$ is encoded in *binary* when we use this definition of $|V|$, which we will use as standard encoding in this paper. Alternatively, when we set $|V| \triangleq |Q| + |T| \cdot d \cdot |T||$, we say that $V$ is encoded in *unary*. Subsequently, $Q \times \mathbb{Z}^d$ denotes the set of configurations of $V$. Note that in the literature, the set of configurations is usually $Q \times \mathbb{N}^d$, however in this paper we will often deal with VASS whose counters can take integer values. For the sake of readability, we write configurations $(q(z_1, \ldots, z_d))$ and $(q, z)$ as $q(z_1, \ldots, z_d)$ and $q(z)$, respectively.

For every subset $A \subseteq \mathbb{Z}^d$, $p(u), q(v) \in Q \times A$ and every transition $t = (p, z, q)$, we write $p(u) \xrightarrow{t}^1 q(v)$ whenever $v = u + z$. We extend $\xrightarrow{t}^1$ to sequences of transitions $\pi \in T^*$ as follows: $\pi \xrightarrow{t}^1$ is the smallest relation satisfying the following conditions for all configurations $p(u)$, $q(v)$, $r(w) \in Q \times A$ and all $t \in T$,
Parikh extends to languages placement of π contains nested loops, V∗(as a linear path schemes. However, by carefully unraveling loops off a characterization of its reachability set by a finite set of

\[ q \bigcup \pi \rightarrow_q \delta \bigcup \epsilon \rightarrow \]

In order to determine the complexity of this problem, we • if \( p(u) \rightarrow_q \pi \rightarrow q(v) \) and \( q(v) \rightarrow_r r(w) \), then \( p(u) \rightarrow_r r(w) \).

We extend \( \rightarrow_q \) to languages \( L \subseteq T^* \) in the natural way, \( \rightarrow_q \subseteq L^* \). We write \( \rightarrow_q \) to denote \( \rightarrow_q \). An \( \mathcal{A} \)-run \( \pi = t_1 \cdots t_k \) is a sequence of configurations \( q_0(v_0) \rightarrow q_k(v_k) \) that we sometimes just abbreviate by \( q_0(v_0) \rightarrow q_k(v_k) \). When \( \mathcal{A} = \mathbb{N} \) we also refer to an \( \mathcal{A} \)-run as a run.

Throughout this paper, we refer to \( \rightarrow_{N^2} \) as the reachability relation, and \( \rightarrow_{Z^d} \) as the \( \mathbb{Z} \)-reachability relation. Let \( \pi = (p_1, z_1, p_1) \cdots (p_k, z_k, p_k) \in T^k \) for some \( k \geq 0 \). The displacement of \( \pi \) is \( \delta(\pi) = \sum_{i=0}^k z_i \), and the definition naturally extends to languages \( L \subseteq T^* \) as \( \delta(L) = \{ \delta(\pi) : \pi \in L \} \). Note in particular that if \( Parikh_p \subseteq Parikh_\pi \), then \( \delta(\pi) \leq \delta(\rho) \).

III. MAIN RESULTS

In this paper, our main interest is in the reachability problem for 2-VASS, formally defined as follows:

2-VASS REACHABILITY

INPUT: A 2-VASS \( V = (Q, T) \) and configurations \( p(u) \) and \( q(v) \) from \( Q \times \mathbb{N}^2 \).

QUESTION: Is there a run from \( p(u) \) to \( q(v) \), i.e. does \( p(u) \rightarrow_{N^2} q(v) \) hold?

In order to determine the complexity of this problem, we show that the reachability relation of any 2-VASS can be defined by a finite set of linear path schemes. In particular, we are able to show strong bounds on their lengths and \( \pi \)-lengths. For example, consider the 2-VASS \( V \) depicted in Fig. 2. Since \( V \) contains nested loops, e.g. \( (t_1 t_2 t_2)^* \), we cannot directly read off a characterization of its reachability set by a finite set of linear path schemes. However, by carefully unraveling loops we obtain the reachability set from the set of the subsequent linear path schemes, and in particular this means that \( V \) can be flattened (cf. Fig. 3):

\[ p(u) \rightarrow_{N^2} q(v) \iff p(u) \rightarrow_{N^2} r(v) \]

We will show that such a flattening exists for any 2-VASS. More precisely, our main technical result states that the global reachability relation of any 2-VASS \( V = (Q, T) \) can be defined via a set of linear path schemes whose lengths can be polynomially bounded in \( |Q| + |T| \), and a fortiori are at most exponential in \( |V| \), and whose \( \pi \)-lengths are at most quadratic in \( |Q| \):

Theorem 1. Let \( V = (Q, T) \) be a 2-VASS. There is a finite set \( S \) of linear path schemes such that

- \( p(u) \rightarrow_{N^2} q(v) \) if, and only if, \( p(u) \rightarrow_{N^2} q(v) \), and
- \( |\rho| \leq (|Q| + |T|)^{O(1)} \) and \( |\rho|_{\pi} \leq O(|Q|^2) \) for each \( \rho \in S \).

Having established Theorem 1, we can show that proving the existence of a run between two reachable configurations in a 2-VASS reduces to checking the existence of a solution for suitably constructed systems of linear Diophantine inequalities that depend on \( S \) and the properties listed in Theorem 1. The absence of nested cycles in linear path schemes in \( S \) is crucial to this reduction. By application of standard bounds from integer linear programming, this in turn enables us to bound the length of paths witnessing reachability, and to prove the upper bound of the the main theorem of this paper in Section V:

Theorem 2. 2-VASS REACHABILITY is PSPACE-complete.

IV. PROOF OF THEOREM 1

In this section, we prove Theorem 1 and show that runs of a 2-VASS \( V = (Q, T) \) are captured by a finite set of linear path schemes each of which has length at most \( (|Q| + |T|)^{O(1)} \).

The expanded technical meaning of this statement is that there are constants \( c_1 \) and \( c_2 \) such that for every 2-VASS \( V = (Q, T) \) there exists a finite set \( S \) of linear path schemes, each of length \( \leq (|Q| + |T|)^{c_1} \) and of \( \pi \)-length \( \leq c_2|Q|^{c_2} \), with the property that for every \( p(u), q(v) \in Q \times \mathbb{N}^2 \), \( p(u) \rightarrow_{N^2} q(v) \) if, and only if, \( p(u) \rightarrow_{N^2} q(v) \). The more familiar statements of this theorem and of lemmas of a similar nature in the rest of the paper were chosen to avoid clutter and to downplay the role of the precise constants.
Here, \( \rho \) if, types of runs set of linear path schemes, we consider the following three

- For all configurations of the run \( p(u_1, u_2) \) both counter values are sufficiently large.
- For all configurations of the run \( p(u_1, u_2) \) at least one counter value is small.

In Subsections IV-A, IV-B and IV-C, we will show how to construct linear path schemes for these three types of runs. Then, in Subsection IV-D, we prove Theorem 1 by showing that any run can be decomposed as finitely many runs of these types.

In more detail, the first step is to show in Section IV-A that Parikh images of finite labeled graphs can be captured by linear path schemes of polynomial size. This will allow us to prove that \( \mathbb{Z} \)-reachability, i.e. runs in which counter values may drop below zero, can be captured by linear path schemes of polynomial size. We then give in Section IV-B an effective decomposition of certain linear sets in dimension two into semi-linear sets with special properties, and use this decomposition in order to derive together with the results in Section IV-A linear path schemes of size \( (|Q| + |T|)^{O(1)} \) and constant \( * \)-length for runs of type (1). Linear path schemes for runs of type (2) will then be seen to follow from the type (1) case.

For runs of type (3), in Section IV-C we construct linear path schemes for 1-VASS and show that runs of a 2-VASS that stay within an “L-shaped band” are, essentially, runs of a 1-VASS. Our analysis of such runs of type (3) is a simple consequence of certain normal forms of shortest runs in one-counter automata, which 1-VASS are a subclass of, by Valiant and Paterson [26].

Similarities and differences in comparison with [14]. Our proof strategy of considering the three kinds of runs described above shares some similarities with [14]. There, the bounds on what we referred to above as “large” and “small” in the runs of type (1), (2) and (3) are not explicitly calculated. Our proofs for obtaining rather tight bounds require new insights. We capture runs of type (1) by linear path schemes of size \( (|Q| + |T|)^{O(1)} \), whereas in [14] the linear path schemes were of size at least exponential in \( |Q| \). To prove the former, we establish a new upper bound on the presentation size of Parikh images of finite automata in Lemma 4 below, which is a result of independent interest. The difference between our runs of type (2) and the ones analyzed in [14] is that our runs have to stay in the “outside region” entirely, whereas in [14] the set of displacements of paths from \( q \) to \( q' \) is analyzed. Runs of type (3) are treated as special cases of their runs of type (2) in [14], whereas we invoke a result by Valiant and Paterson on normal forms of minimal runs in one-counter automata.

Our final proof of Theorem 1 shows that each run can be factorized into segments of runs of types (1), (2) and (3) and requires a more careful treatment than in [14]. At every step, we have to ensure that the \( * \)-length of the linear path schemes we construct stays polynomial in the number of control states. This aspect is neglected in [14] as it is of no interest for the goal of [14], however, for us it is by far the technically most challenging part and one of the cornerstones of our PSPACE upper bound.

A. Parikh images of finite directed graphs and \( \mathbb{Z} \)-reachability of \( d \)-VASS

The purpose of this subsection is to prove the following proposition.

Proposition 3. Let \( V = (Q, T) \) be a \( d \)-VASS. There exists a finite set \( S \) of linear path schemes such that

1. Both counter values of \( p(u_1, u_2) \) and of \( q(v_1, v_2) \) are sufficiently large and \( p = q \), but intermediate configurations on the run \( p(u_1, u_2) \) may have arbitrarily small counter values.
2. For all configurations of the run \( p(u_1, u_2) \) both counter values are sufficiently large.
3. For all configurations of the run \( p(u_1, u_2) \) at least one counter value is small.

4. The \( * \)-length of the linear path schemes constructed stays polynomial in the number of control states.

In order to prove Proposition 3, we will prove suitable bounds on the representation size of the Parikh images of paths of a \( \Sigma \)-labeled finite graphs (or equivalently, nondeterministic finite automata) in terms of linear path schemes. Even though estimations on this size have been made in the literature (e.g. in [14] or [25, Prop. 7.3.4]), we are not aware of any in which the \( * \)-length is linear in the number of edges (in the aforementioned references, the \( * \)-length may be exponential in the number of control states).

Lemma 4. Let \( G = (U, E) \) be a finite \( \Sigma \)-labeled graph. There exists a finite set \( S \) of linear path schemes such that

1. \( \{ \text{Parikh}_{\pi} : \pi \text{ is a path in } G \} \) is a family of mappings and \( \{ \text{Parikh}_\rho : \rho \in S \} \), and
2. \( |\rho| \leq 2 \cdot |U| \cdot |E| \) and \( |\rho|_s \leq |E| \) for each \( \rho \in S \).

Proof. We first provide some additional definitions. Let \( \sigma, \sigma' \) : \( E \to \mathbb{N} \) be mappings and let \( X \) be a set of such mappings. We define \( \sigma + \sigma' \in \mathbb{N}^E \) as \( (\sigma + \sigma')(e) \) \( \sigma(e) + \sigma'(e) \) for each \( e \in E \) and \( \sigma + X \) \( \{ \tau + \sigma : \tau \in X \} \). For each \( u \in U \), let \( \text{in}(u) \) \( \{ (v', a, u''') \in E : u''' = u \} \) and \( \text{out}(u) \) \( \{ (u', a, u'') \in E : u' = u \} \) denote the set of incoming and outgoing edges of
Fig. 4. Example of the three types of runs. The region depicted in each case is the positive quadrant in the Cartesian plane. (1) left: run from \( q \) to \( q \) starting and ending sufficiently high; (2) middle: run staying sufficiently high; (3) right: run within an L-shaped band, i.e., running high on at most one component at a time.

\[ u, \text{ respectively. We say that } \sigma \text{ is flow-preserving if for every } u \in U \text{ we have} \]
\[ \sum_{e \in \text{in}(u)} \sigma(e) = \sum_{e \in \text{out}(u)} \sigma(e). \]

We will show the following claim:

**Claim.** Let \( \pi \in E^* \) be a path. There exists some \( h \geq 1 \), a sequence of linear path schemes \( \rho_1, \ldots, \rho_h \subseteq E^* \), and a sequence \( \sigma_1, \ldots, \sigma_h \in \mathbb{N}^E \) such that
(a) \( \rho_1 \) is a path of length at most \( |U| \cdot |E| \) that visits each vertex of \( \pi \) at least once,
(b) \( \sigma_1 \) is flow-preserving,
(c) \( \text{Parikh}_\pi = \text{Parikh}_{\rho_1} + \sigma_1 \),
and for every \( 1 < i \leq h \),
(1) \( \rho_i \) is a linear path scheme that can be obtained from \( \rho_{i-1} \)
by inserting \( i - 1 \) simple cycles (in the form \( \beta^* \)),
(2) \( \sigma_i \) is flow-preserving,
(3) \( \text{Parikh}_{\rho_{i-1}} + \sigma_{i-1} \subseteq \text{Parikh}_{\rho_i} + \sigma_i \),
(4) \( \sigma_{i-1}(e) \geq \sigma_i(e) \) for all \( e \in E \) and there exists some \( e \in E \) s.t. \( \sigma_{i-1}(e) > \sigma_i(e) = 0 \), and
(5) \( \sigma_i(e) = 0 \) for all \( e \in E \).

First observe that due to (4) we have \( |h| \leq |E| \), and due to (1) we have \( |\rho_1| \leq |\rho_1| + |U| \cdot (i-1) \). Thus, \( |\rho_h| \leq |\rho_1| + |U| \cdot |E| \leq 2 \cdot |U| \cdot |E| \), where the last inequality is due to (a). Moreover, \( \sigma_h(e) \leq |E| \) due to (1) and \( h \leq |E| \).

Before proving the claim, let us first see how it proves the lemma. We define
\[ S \overset{\text{def}}{=} \{ \rho : \rho \text{ is a linear path scheme, } |\rho| \leq 2 \cdot |U| \cdot |E| \text{ and } |\rho| \leq |E| \}. \]

Trivially, (ii) is satisfied. To establish (i), let us fix an arbitrary path \( \pi \) and obtain a linear path scheme \( \rho_h \in S \) for \( \pi \) from the above claim. We have \( \text{Parikh}_\pi \overset{(3)}{=} \text{Parikh} (\rho_1) + \sigma_1 \overset{(3)}{=} \text{Parikh}(\rho_2) + \sigma_2 \overset{(3)}{=} \cdots \overset{(3)}{=} \text{Parikh}_{\rho_h} + \sigma_h \overset{(5)}{=} \text{Parikh}_{\rho_h} \) as required.

We now prove the claim. Let \( \pi \) be a path and let us first define \( \rho_1 \) and \( \sigma_1 \) such that (a), (b) and (c) are satisfied. The path \( \pi \) can be decomposed as \( \pi = e_1 \pi_1 \cdots e_k \pi_k \) where \( k \leq |U| \) and each \( e_j = (u, a, u') \) is the first transition such that \( u \) or \( u' \) appears in \( \pi \). We define \( \rho_1 \) and \( \sigma_1 \) as the result of the following iterative process: We initially set \( \rho_1 \) to \( \pi \) and set \( \sigma_1(e) = 0 \) for all \( e \in E \); then we successively remove a simple cycle \( \beta \) from some \( \pi_j \), and add \( \text{Parikh}_\beta \) to \( \sigma_1 \). We repeat this process until no longer possible. The resulting \( \rho_1 \) is a path of length at most \( |U| \cdot |E| \). Moreover, \( \sigma_1 \) is flow-preserving since we successively removed cycles only, and clearly \( \text{Parikh}_\pi = \text{Parikh}_{\rho_1} + \sigma_1 \), by construction. Thus (a), (b) and (c) hold.

Let us prove (1) to (5) by induction on \( 1 < i \leq h \). We only prove the induction step, the base case can be proven analogously. Let \( E' \overset{\text{def}}{=} \{ e \in E : \sigma_{i-1}(e) > 0 \} \). If \( E' = \emptyset \), then (5) holds and we are done. Thus, we assume that \( E' \neq \emptyset \).

Let us fix a choice function \( \chi : E' \to E' \) satisfying
\[ \chi(u_1, a, u_2) = (u_1', a, u_2') \implies u_2 = u_1'. \]

Note that \( \chi \) exists since \( \sigma_{i-1} \) is flow-preserving by induction hypothesis. By the pigeonhole principle, there exists some \( c \in E' \) and some \( \ell \geq 0 \) such that \( \beta \overset{\text{def}}{=} c \chi(c) \chi^2(c) \cdots \chi^\ell(c) \) is a simple cycle. Without loss of generality let us assume that \( c \overset{\text{def}}{=} \sigma_{i-1}(c) \overset{\min\{\sigma_{i-1}(\chi(e)) : h \in [0, \ell]\}}{=} \sigma_{i-1}(c) \overset{\text{is minimal among all edges that lie on the simple cycle } \beta}{} \). We define \( \sigma_1 \overset{\text{def}}{=} \sigma_{i-1} - \text{Parikh}_\beta \) and observe that \( \sigma_1 \) is flow-preserving because \( \beta \) is a cycle and \( \sigma_1 \in \mathbb{N}^E \) due to minimality of \( c \); thus (2) and (4) are shown. Let \( e = (u, a, u') \), hence \( \beta \) is a simple cycle from \( u \) to \( u' \). By (1) of induction hypothesis the linear path scheme \( \rho_{i-1} \) can be obtained from \( \rho_1 \) by inserting \( (i - 2) \) simple cycles and can hence be factorized as \( \rho_{i-1} = \alpha \gamma \), where \( \alpha \) is a linear path scheme from some state to \( u \). We set \( \rho_1 \overset{\text{def}}{=} \pi \alpha \beta^* \gamma \overset{\text{and hence (1) holds. Furthermore, (3) holds due to} \text{Parikh}_{\rho_{i-1}} + \sigma_{i-1} = \text{Parikh}_{\rho_{i-1}} + \text{Parikh}_\beta + \sigma_1 \subseteq \text{Parikh}_{\rho_1} + \sigma_i}{} \)

**Proof of Proposition 3.** We have \( T \subseteq Q \times \Sigma \times Q \) for some finite subset \( \Sigma \subseteq \mathbb{Z}_d \). Let \( S \) be the finite set of linear path schemes from Lemma 4, then (ii) of Proposition 3 is clear. Let us now prove (i). We have \( p(u) \xrightarrow{\Sigma} q(v) \) if, and only
if, there exists a path $\pi$ from $p$ to $q$ in $V$ such that $v - u = \sum_{z \in \Sigma} P_{\text{res}}(z) \cdot z$. By Lemma 4 (i), this is equivalent to the existence of some $\rho \in S$ from $p$ to $q$, and some $f \in P_{\text{res}}$ such that $v - u = \sum_{z \in \Sigma} f(z) \cdot z$. Now the latter existence of $f$ is equivalent to $v - u \in \delta(\rho)$. This shows that $p(u) \sim_{Z^d} q(v)$ if, and only if, $p(u) \rightarrow_{Z^d} q(v)$. 

B. Starting and ending in “sufficiently large” configurations

The goal of this subsection is to prove that, given a 2-VASS $V$, there exists a sufficiently small bound $D$ such that the reachability relation between any two configurations $q(u_1, v_1)$ and $q(u_2, v_2)$ with $u_1, u_2, v_1, v_2 \geq D$ can be captured by a finite set of small linear path schemes (in the sense of Theorem 1). In [14], this property is referred to as ultimately flat. As a consequence of this result, we can show that the reachability relation between arbitrary configurations for which there exists a run on which both counter values on all configurations stay above $D$ can be captured by a finite set of small linear path schemes as well.

**Proposition 5.** Let $V = (Q, T)$ be a 2-VASS. There exist $D \leq (|Q| + |T|)^{(O(1))}$ and finite sets of linear path schemes $R, X$ such that for $q \in \mathbb{Q}$ and $\rho \in R$.

(a) $q(u) \sim_{R \cap \mathbb{Q}^2} q(v)$ if, and only if, $q(u) \sim_{R \cap \mathbb{Q}^2} q(v)$, and $|\rho| \leq (|Q| + |T|)^{(O(1))}$ and $|\rho| < 2$ for every $\rho \in R$.

(b) $p(u) \rightarrow_{R} q(v)$ implies $p(u) \sim_{R \cap \mathbb{Q}^2} q(v)$, and $|\rho| \leq (|Q| + |T|)^{(O(1))}$ and $|\rho| < 2 \cdot |Q|$ for every $\rho \in X$.

The proof of this proposition requires two intermediate steps. First, in Lemma 6 below we prove an effective decomposition of certain linear sets in dimension two into semi-linear sets with nice properties. Similar decompositions have been the cornerstone of the results by Hopcroft and Pansiot [9] and Leroux and Sutre [14]. The contribution of Lemma 6 is to establish a new proof from which we can obtain sufficiently small bounds on this decomposition. Next, in Lemma 7 we show how this decomposition can be applied in order to capture reachability instances by linear path schemes of $s$-length two whose displacements all point into the same quadrant. This in turn enables us to prove Part (a) of Proposition 5, from which we can then prove Part (b).

Let us recall some definitions concerning semi-linear sets. Let $P = \{p_1, \ldots, p_n\} \subseteq \mathbb{Q}^2$ and $D \subseteq \mathbb{Q}_{\geq 0}$. The $D$-cone generated by $P$ is defined as

$$\text{cone}_D(P) = \left\{ \sum_{i \in [1,n]} \lambda_i \cdot p_i : \lambda_i \in \mathbb{D} \right\}.$$ 

A linear set $L(b, P)$ is determined by a base vector $b \in \mathbb{Z}^d$ and a finite set of period vectors $P \subseteq \mathbb{Z}^d$, where $L(b, P) = b + \text{cone}_D(P)$. A semi-linear set is a finite union of linear sets. The norm $\|P\|$ of a finite set $P \subseteq \mathbb{Z}^d$ is defined as $\|P\| = \max\{\|p\| : p \in P\}$. Recall that $u, v \in \mathbb{Z}^d$ are linearly dependent if $0 = \lambda_1 \cdot u + \lambda_2 \cdot v$ for some $\lambda_1, \lambda_2 \in \mathbb{Q} \setminus \{0\}$, and linearly independent otherwise.

We now show the following decomposition of linear sets.

**Lemma 6.** Let $b \in \mathbb{Z}^d$, let $P \subseteq \mathbb{Z}^d$ be finite with $b \in P$ and let $Z$ be a quadrant. Then there exists a finite set of indices $\{i_1, \ldots, i_k\}$ such that $\bigcup_{i \in I} L(c_i ; P_i)$, and for each $i \in I$ we have

- $|P_i| \leq 2$,
- $P_i \subseteq (P \cap L(c_i ; P_i)) \cap Z$,
- there exists $\varepsilon \leq \|P\|^{(O(1))}$ such that $\{c_i \} \cup (P_i \cap L(b; P_i) \subseteq b + \text{cone}_{\{0, \varepsilon\}}(P)$.

**Proof sketch.** It is sufficient to show the statement for $Z = \mathbb{N}^2$. Let $P = \{p_1, \ldots, p_n\}$ and let $r$ be a point in $L(b; P) \cap \mathbb{N}^2$. By definition,

$$r = b + \lambda_1 p_1 + \cdots + \lambda_n p_n$$

for some $\lambda_i \in \mathbb{N}$, $i \in [1,n]$. First, it is not difficult to show that we may assume that all but two of the $\lambda_i$ do not exceed $O(\|P\|^2)$. Hence, with no loss of generality we may assume that $r \in L(c; P')$, where

$$c \in b + \text{cone}_{\{0, \varepsilon\}}(P)$$

for some $\varepsilon \leq \|P\|^{(O(1))}$ and $P' = \{u, v\} \subseteq P$. The most interesting case is when $u$ and $v$ are linearly independent, which entails making a case distinction in which quadrant $u$ and $v$ lie. If both $u$ and $v$ lie in $\mathbb{N}^2$ then we are done.

If $u \notin \mathbb{N}^2$ and $v \notin \mathbb{N}^2$ then we can show that there exists some natural number $\alpha \leq \|P\|^{(O(1))}$ such that either

$$\begin{cases} (\alpha, 0) \in L(b; P) \cap \text{cone}_{\{\varepsilon\}}(P') \quad \text{or} \\ (0, \alpha) \in L(b; P) \cap \text{cone}_{\{\varepsilon\}}(P') \end{cases}$$

depending on the relative angle between $u$ and $v$. In the following, assume $(0, \alpha) \in L(b; P) \cap \text{cone}_{\{\varepsilon\}}(P')$, the other case follows symmetrically. A careful analysis allows to conclude that $r$ can equivalently be obtained as

$$r \in c + \omega \cdot v + \text{cone}_{\{\varepsilon\}}(P')$$

for some $\omega \leq \|P\|^{(O(1))}$ and $P'' = \{u, (0, \alpha)\}$, i.e., $r \in L(c + \omega \cdot v; P'')$, which fulfills the requirements of the lemma.

Finally, if both $u \notin \mathbb{N}^2$ and $v \notin \mathbb{N}^2$ then in the non-trivial case both $(\alpha, 0)$ and $(0, \alpha)$ can be obtained as in (1). Applying similar reasoning as above, it is then possible to show that

$$r \in c + \lambda \cdot u + \gamma \cdot v + \text{cone}_{\{\varepsilon\}}(P')$$

for some $\lambda, \gamma \leq \|P\|^{(O(1))}$ and $P'' = \{(\alpha, 0), (0, \alpha)\}$, which again fulfills the requirements of the lemma. 

Let us give an intuitive idea of how we can prove Proposition 5 (a) by an application of Lemma 6. Suppose we are given a run starting in $q(u_1, v_2)$ and ending in $q(u_1, v_2)$ such that w.l.o.g. $u_1 \leq v_1$ and $v_2 \leq v_2$. From Proposition 3 we know that the $Z$-reachability relation can be captured by a finite set of linear path schemes. Since we start and end in the same state one can show that any (slight modification of) such
a linear path scheme describes a set of displacements equal to a linear set \(L(b; P)\) such that \(b \in P\). An application of Lemma 6 then allows us to decompose such a linear set into a semi-linear set whose period vectors all point into the same \(\mathbb{N}^2\) direction. The crucial point is that any linear set in this semi-linear set can again be translated back into a linear path scheme of \(+\)-length at most two whose displacements point to \(\mathbb{N}^2\). Consequently, any path obtained from such a linear path scheme does not, informally speaking, drift away too much, and if \(u_1\) and \(u_2\) are sufficiently large then \(\mathbb{N}\)-reachability and \(\mathbb{Z}\)-reachability coincide.

Consequently, the first step is to interpret Lemma 6 in terms of linear path schemes. As in [14], subsequently we say that a linear path scheme \(\alpha_0\beta_1^*\alpha_1 \cdots \beta_k^*\alpha_k\) is zigzag-free if \(\{\delta(\beta_1), \ldots, \delta(\beta_k)\} \subseteq Z\) for some quadrant \(Z\).

**Lemma 7.** Let \(q \in Q\). For every linear path scheme \(\rho\) from \(q\) to \(q\), there exists a finite set \(R_\rho\) of zigzag-free linear path schemes such that

(i) \(\delta(\rho) \subseteq \delta(R_\rho)\), and

(ii) \(|\sigma| \leq (|\rho| + |T|)^{(O(1)}\) and \(|\rho|_\ast \leq 2\) for each \(\sigma \in R_\rho\).

**Proof sketch.** Let \(\rho = \alpha_0\beta_1^*\alpha_1 \cdots \beta_k^*\alpha_k\) be a linear path scheme. We have that \(\alpha_0 \cdots \alpha_k\) and \(\beta_1 \cdots \beta_k\) are cycles. The idea is to view the displacements of those cycles as linear sets, i.e., define \(b = \delta(\alpha_0 \cdots \alpha_k), p_i = \delta(\beta_i)\) for \(i \in [1, k]\), and \(P = \{b, p_1, \ldots, p_k\}\). It is not difficult to verify that

\[
\delta(\rho) \subseteq L(b; P) = \bigcup Z \in Z L(b; P) \cap Z.
\]

For a fixed quadrant \(Z\), by Lemma 6 we have \(\delta(\rho) \cap Z \subseteq \bigcup_{i \in I} L(c_i; P_i)\) such that for every \(i \in I\), \(|P_i| \leq 2\), \(P_i \subseteq (P \cup L(b; P)) \cap Z\), and \(u \in b + \text{cone}_{[a,b]}(P)\) for some \(e \leq |P|^{O(1)}\) for every \(u \in \{c_i\} \cup (P \cap L(b; P))\). The latter property allows for translating every such \(u\) into a path \(\pi_u \in \rho\) such that \(u = \delta(\pi_u)\). Hence, for every \(i \in I\) we can obtain from \(\rho\) some \(\rho_i\) such that \(\delta(\rho_i) = L(c_i; P_i)\), the length of \(\rho_i\) does not increase too much, and \(\rho_i\) only has the two cycles from \(P_i\) which point into \(Z\).

We are now fully prepared to give a proof of Proposition 5.

**Proof of Proposition 5.** Let us fix a 2-VASS \(V = (Q, T)\).

**Proof of (a):** Let \(S\) be the finite set of linear path schemes from Proposition 3 such that

- \(p(\rho) \xrightarrow{\tau} q(v)\) if, and only if, \(p(\rho) \xrightarrow{S} q(v)\), and
- \(|\rho| \leq 2 \cdot |Q| \cdot |T|\) and \(|\rho|_\ast \leq |T|\) for each \(\rho \in S\).

For each \(\rho \in S\), let \(R_\rho\) be the set of zigzag-free linear path schemes from Lemma 7, and define \(R \defeq \bigcup_{\rho \in S} R_\rho\). Hence, for each \(\sigma \in R\) we have \(|\sigma| \leq (2 \cdot |Q| \cdot |T| + |T||)^{(O(1)}\) = \((|Q| + |T||)^{(O(1)}\) by (ii) of Lemma 7. We see that set \(D\) required in Proposition 5 to \(D \defeq \max\{|\sigma| : \sigma \in R\} \cdot |T| \leq \((|Q| + |T||)^{(O(1)}\) the monotonicity of zigzag-free linear path schemes now provides the key ingredient for proving Proposition 5 (a). For the rest of the proof let us fix \(u, v \in [D, \infty)^2\) and some zigzag-free linear path scheme \(\sigma = \alpha_0\beta_1^*\alpha_1 \cdots \beta_k^*\alpha_k \in R\). Suppose \(q(u) \xrightarrow{\tau} q(v)\) for some \(\sigma = \alpha_0\beta_1^*\alpha_1 \cdots \beta_k^*\alpha_k\), then by definition of \(D\) and the fact that \(\sigma\) is zigzag-free it is clear that for every \(i \in [0, |\pi|]\),

\[
0 \leq u + \delta(\pi[1, i])
\]

It remains to prove \(q(u) \xrightarrow{\tau} q(v)\) if, and only if, \(q(u) \xrightarrow{\tau} q(v)\) for some \(\sigma \in R\), which follows from:

\[
q(u) \xrightarrow{\tau} q(v) \quad \Rightarrow \quad q(u) \xrightarrow{\tau} q(v)
\]

**Proposition 3**

\[
q(u) \xrightarrow{\tau} q(v) \quad \forall \rho \in S
\]

**Lemma 7 (i)**

\[
q(u) \xrightarrow{\tau} q(v)\) for some \(\sigma \in R_\rho, \rho \in S
\]

\[
(2) \quad q(u) \xrightarrow{\tau} q(v)\) for some \(\sigma \in R
\]

\[
q(u) \xrightarrow{\tau} q(v)\) for some \(\sigma \in R
\]

**Proof of (b):** Suppose that \(p(\rho) \xrightarrow{\tau} q(v)\). Then \(\pi\) can be factorized as \(\pi = \alpha_0\beta_1^*\alpha_1 \cdots \beta_k^*\alpha_k\) such that

\[
p(\rho) = \alpha_0 \beta_1^* \rho_1 \alpha_1 \beta_1 \cdots \beta_k^* \rho_k \alpha_k,
\]

where \(|\alpha_0|, |\alpha_1|, \ldots, |\alpha_k| \leq |Q|\), each \(\beta_i\) is a cycle from \(q_i\) to \(q_i\) for some \(q_i \in Q\), and \(k \leq |Q|\). Since \(u_1, u_2 \in \mathbb{O}\) for all \(i \in [1, k]\), by Part (a) of Proposition 5 we have \(q_i(u_i) \xrightarrow{\tau} q_i(u_i')\) for some linear path scheme \(\rho_i \in R\). Consequently, we define \(X\) as

\[
X \defeq \{\text{linear path scheme } \alpha_0\rho_1\alpha_1 \cdots \rho_k\alpha_k : k \leq |Q|, \alpha_i \in T^*, \alpha_i \leq |Q|, \rho_i \in R\}.
\]

Let \(\rho \in X\), then we have \(|\rho| \leq |Q|^2 + |Q| \cdot (|Q| + |T|)^{(O(1)}\) and \(|\rho|_\ast \leq 2 \cdot |Q|\).

**C. Reachability in 2-VASS with One Bounded Component**

The purpose of this section is to establish the following result on reachability within L-shaped bands, as illustrated in the right-most picture of Fig. 4.

**Proposition 8.** Let \(V = (Q, T)\) be a 2-VASS, \(D \in \mathbb{N}\) and \(\mathbb{L} = ([0, D] \times \mathbb{N}) \cup \mathbb{D} \times [0, D]\). There exists a finite set \(Y_L\) of linear path schemes such that

(i) \(p(\rho) \xrightarrow{\tau} q(v)\) implies \(p(\rho) \xrightarrow{Y_L} q(v)\), and

(ii) \(|\rho| \leq (|Q| + |T| + D)^{(O(1)}\) and \(|\rho|_\ast \leq 2\) for every \(\rho \in Y_L\).

We briefly sketch the proof of Proposition 8 here. In its essence, restricting the set of admissible values of one of the two counters of a 2-VASS to \([0, D]\) gives rise to a 1-VASS. This observation enables us to resort to techniques and results developed for 1-VASS. In particular, to the following lemma established by Valiant and Paterson.

**Lemma 9** (Lemma 2 in [26]). Let \(V = (Q, T)\) be a 1-VASS such that \(T \subseteq Q \times \{-1, 0, 1\} \times Q\) and let \(p(\rho) \xrightarrow{\tau} q(v)\)
for some configurations \( p(u) \) and \( q(v) \) such that \(|u - v| \geq |Q| + |Q|^2\). There exist \( \alpha, \beta, \gamma \in T^* \) and \( \pi \in T^* \) such that \( p(u) \xrightarrow{\pi} q(v) \) and \( \pi \) has the following properties,

(i) \( \pi = \alpha \beta \gamma \) for some \( \gamma > 0 \), 
(ii) \( \alpha \beta \gamma \) is a linear path scheme of *-length one, and 
(iii) \( |\alpha \gamma| < |Q|^2 \) and \( \beta \) is a cycle with \(|\beta| \leq |Q|\) and \(|\delta(\beta)| \in [1, |Q|]\).

As a consequence of Lemma 9, one can show that the reachability relation of 1-VASS can be captured by linear path schemes with the following properties.

**Lemma 10.** Let \( V = (Q, T) \) be a 1-VASS. There exists a finite set \( Y \) of linear path schemes such that

(i) \( p(u) \xrightarrow{Y} q(v) \) if, and only if, \( p(u) \xrightarrow{Y} q(v) \), and 
(ii) \(|\pi| \leq (|Q| + |T|)^{(O(1)} \) and \(|\pi|_+ \leq 1 \) for each \( \pi \in Y \).

This, in turn, allows us to prove Proposition 8.

**Proof sketch of Proposition 8.** Let \( V = (Q, T) \) be a 2-VASS and \( D \in \mathbb{N} \). Let \( B_1 = \{ u \in X \times [0, D] \} \times \mathbb{N} \) and \( B_2 = \{ u \in X \times \mathbb{N} \} \). Consider a run \( p(u) \xrightarrow{\pi} q(v) \) of minimal length. By making \( \pi \) slightly larger, we can decompose \( \pi \) as a sequence of runs that remain within either the vertical band, or the horizontal band. Formally, let \( B'_1 = \{ u \in X \times \mathbb{N} \} \times \mathbb{N} \), \( B'_2 = \{ u \in X \times \mathbb{N} \} \), and \( H = B'_1 \cap B'_2 \). Due to minimality of \( |\pi| \) we can factorize \( \pi = \pi_1 \cdots \pi_k \), where \( p_0 = p \), \( p_k = q \), \( u_0 = u \), \( u_k = v \), 

- \( \pi_i \in \{ B'_1, B'_2 \} \) for each \( i \in [1, k] \), 
- \( u_i \in H \) for each \( i \in [1, k-1] \), and 
- \( k \leq |H| \leq (D + |T|)^2 \).

Note that a 2-VASS in which the value of at least one counter is bounded by some \( E \in \mathbb{N} \) can be simulated by a 1-VASS with \(|Q|\cdot(E+1)\) states. This allows us, by applying Lemma 10, to replace \( \pi_0, \pi_1, \ldots, \pi_k \) by some linear path schemes \( p_0, p_1, \ldots, p_k \). However, this would yield a linear path scheme of \(*\)-length \( k \). In fact we can restrict only \( p_0 \) and \( p_k \) to have \(*\)-length at most one, where \( p_1, p_2, \ldots, p_{k-1} \) can be chosen to have \(*\)-length zero. Indeed, since \( \pi_1, \pi_2, \ldots, \pi_{k-1} \) are runs from \( H \) to \( H \), which is finite, and \( p_1, p_2, \ldots, p_{k-1} \) only have one cycle, then \( \pi_1, \pi_2, \ldots, \pi_{k-1} \) can each be replaced with runs of length at most \(|Q| + |T| + D)^{(O(1)} \). Therefore, \( p(u) \xrightarrow{\pi} q(v) \) where \( \sigma \in T^* \) and \(|\sigma| \leq (|Q| + |T| + D)^{(O(1)} \). \( \square \)

**D. Factorizing arbitrary runs: Proof of Theorem 1**

By application of the results established in Sections IV-B and IV-C, we will now prove Theorem 1. More precisely, we will show that any run can be factorized into few runs of types (1), (2) or (3). To this end, we will use a 2-VASS \( V = (Q, T) \). Let \( D \leq (|Q| + |T|)^{(O(1)} \) be the constant from Proposition 5. Informally speaking, we have hereby defined that “sufficiently large” means to be greater or equal to \( D \).

Moreover we set \( L \equiv ([0, D + |T|] \times \mathbb{N}) \cup X \times ([0, D + |T|] \times \mathbb{N}) \), 
\( \mathbb{N} \equiv \{ D, D + |T| \} \), and \( B \equiv \{ 0, \mathbb{N} \} \)

. Again, informally speaking, we have hereby defined that “small” means to be less or equal to \( D + |T| \).

Let us summarize what we have proven in Sections IV-B and IV-C:

- Runs of type (1) can be captured by a set of linear path schemes \( R \), where each \( \rho \in R \) has \(*\)-length at most two and length at most \(|Q| + |T|)^{(O(1)} \) by Proposition 5 (a).
- Runs of type (2) can be captured by a set of linear path schemes \( X \), where each \( \rho \in X \) has \(*\)-length at most two and length at most \(|Q| + |T|)^{(O(1)} \) by Proposition 5 (b).
- Runs of type (3) can be captured by a set of linear path schemes \( Y \), where each \( \rho \in Y \) has \(*\)-length at most two and length at most \(|Q| + |T|)^{(O(1)} \) by Proposition 8.

Given \( p(u) \) and \( q(v) \), let us fix an arbitrary run \( p(u) \xrightarrow{\pi} q(v) \), where \( \pi = t_1 \cdots t_k \in T^k \) and \( p(u) = q_0(u_0) \xrightarrow{t_1} q_1(u_1) \cdots \xrightarrow{t_k} q_k(u_k) = q(v) \).

We will be interested in the indices of configurations whose counter values lie in \( \mathbb{N} \) and define \( I \equiv \{ i \in \{ 0, k \} : u_i \in \mathbb{N} \} \).

Let us define the function \( x : I \to I \) that maps each index \( i \in I \) to the smallest element in \( I \) larger than \( i \) (and \( i \) if \( i = \max I \), i.e.

\[
x(i) \defeq \begin{cases} \min\{ j : j > i \} & \text{if } i < \max I, \\ i & \text{otherwise.} \end{cases}
\]

We also define the function \( \ell : \{ q_i \in Q : i \in I \} \to I \) that maps each state \( q \) that appears in a configuration in \( Q \times \mathbb{N} \) to the largest index in \( I \) where it appears, i.e. \( \ell(q) \defeq \max \{ i \in I : q = q_i \} \). We are now interested in factorizing the run \( p(u) \xrightarrow{\pi} q(v) \) into runs between configurations that start and end in \( B = \{ 0, \mathbb{N} \} \). More precisely, by the choice of \( \mathbb{N} \), \( \mathbb{L} \) and \( B \) and by the pigeonhole principle there exist indices \( i_1, \ldots, i_h \in I \) such that the run \( p(u) \xrightarrow{\pi} q(v) \) can be factorized as (cf. Fig. 5):

\[
q_0(u_0) \xrightarrow{t_1} q_1(u_1) \xrightarrow{t_2} q_2(u_2) \cdots \xrightarrow{t_h} q_h(u_h) = q(v)
\]

where

(a) \( h \leq |Q| \),
(b) \( i_t \in I \) and thus we have \( u_{i_t} \in \mathbb{N} \) and \( q_{i_t} = q_i \) for each \( t \in [1, h] \),
(c) \( D_t \in \{ \mathbb{N}, \mathbb{L} \} \) for each \( t \in [0, h] \), and 
(d) \( i_{t+1} = x(\ell(q_{i_t})) \) for each \( t \in [1, h - 1] \).

By (b) each run of the form \( q_i(u_{i_t}) \xrightarrow{t} q_i(\ell(q_i)) \) is a run of type (1) and can hence be replaced by some linear path scheme from \( R \) (recall that \( \mathbb{N} \subseteq \mathbb{O} \)). By (c) and (d), each run of the form \( \xrightarrow{t} \) is a run of type (2) or type (3) and can hence be replaced by some linear path scheme from \( X \cup Y \). In summary, the run \( p(u) \xrightarrow{\pi} q(v) \) can be replaced by a linear path scheme of \(*\)-length at most \((h + 1) \cdot 2 \cdot |Q| \leq O(|Q|^2) \) and length at most \((h + 1) \cdot (|Q| + |T|)^{(O(1)} = (|Q| + |T|)^{(O(1)} \).

This concludes the proof of Theorem 1.
A system of linear Diophantine equations is said to be feasible if there exists a solution \( e \in \mathbb{N}^k \) of \( \mathcal{E} \) such that \( \| e \| \leq (\| A \| + \| c \|)^{O(d)} \).

Consequently, obtaining a PSPACE upper bound for reachability reduces to bounding the binary representation of the \( e_i \) polynomially in the sizes of \( V \), \( u \) and \( v \).

Our approach is straightforward: we rephrase the existential question from (3) in terms of finding solutions to a system of linear Diophantine inequalities and then apply the aforementioned bounds from integer linear programming in order to bound the \( e_i \). For technical convenience we distinguish for each linear path scheme and every cycle of the linear path scheme whether the cycle is taken at least once or not at all. To this end, we use the function \( \text{sign} : \mathbb{N} \rightarrow \{0, 1\} \) as \( \text{sign}(n) = 1 \) if \( n \geq 1 \) and \( \text{sign}(n) = 0 \) if \( n = 0 \). Our approach is formalized by the following lemma.

**Lemma 14.** Let \( V = (Q, T) \) be a \( d \)-VASS, \( u \in \mathbb{N}^d \) and \( \rho = \alpha_0 \beta_1^e \cdot \alpha_1 \cdot \cdots \cdot \beta_k^e \cdot \alpha_k \) be a linear path scheme from \( p \) to \( q \) and let \( \chi : \{1, k\} \rightarrow \{0, 1\} \). Then there exists a system of linear Diophantine inequalities \( \mathcal{I} = \mathcal{I}(u, \rho, \chi) \) of the form \( \mathcal{I} : A \cdot x \geq e \) such that

- \( e \in [\mathbb{N}]^k \) if, and only if, \( \pi = \alpha_0 \beta_1^e \cdot \alpha_1 \cdot \cdots \cdot \beta_k^e \cdot \alpha_k \), \( p(u) \xrightarrow{\rho(r)} q(u + \delta(\pi)) \) and \( \chi(i) = \text{sign}(e_i) \) for every \( e = (e_1, \ldots, e_k) \in \mathbb{N}^k \).
- \( A \) is a \((d + 1) \cdot k \) \( \times k \)-matrix, and
- \( \| A \| \leq k \cdot \| \rho \| \cdot \| T \| \) and \( \| c \| \leq O(\| u \| + \| \rho \| \cdot \| T \|) \).

**Proof.** We only prove the lemma for the concrete function \( \chi : \{1, k\} \rightarrow \{0, 1\} \), where \( \chi(i) = 1 \) for all \( i \in [1, k] \). In the following, we write \( x = (x_1, \ldots, x_k) \). First, we assert that the solutions \( e_i \) are greater or equal to 1, i.e.,

\[
I_k \cdot x \geq 1,
\]

where \( I_k \) is the \( k \times k \) unit matrix and \( I = (1, \ldots, 1) \). Next, informally speaking, we have to construct \( \mathcal{I} \) in a way such that we assert that the counter value does not drop below zero on
any infix of $\rho$ in any dimension. For segments of $\rho$ between cycles, this can be ensured by the following constraints for every $j \in [0, k]$ and $\ell \in [1, |\beta_j|]$, which simply enforce the accumulated counter value to be non-negative:

$$u + \sum_{0 \leq i < j} (\delta(\alpha_i) + \delta(\beta_{i+1}) \cdot x_{i+1}) + \delta(\alpha_j[1, \ell]) \geq 0$$

if, and only if,

$$\sum_{1 \leq i \leq j} \delta(\beta_i) \cdot x_i \geq -u - \sum_{0 \leq i < j} \delta(\alpha_i) - \delta(\beta_j[1, \ell]) \quad (5)$$

For counter values which, informally speaking, occur along cycles $\beta_j$ of $\rho$, it is sufficient to only check whether their initial and final segments lead to counter values greater or equal to zero. Formally, we assert the following constraints for every $j \in [1, k]$ and $\ell \in [1, |\beta_j|]$:

$$u + \delta(\alpha_0) + \sum_{1 \leq i < j} (\delta(\beta_i) \cdot x_i + \delta(\alpha_i)) + \delta(\beta_j[1, \ell]) \geq 0$$

$$u + \delta(\alpha_0) + \sum_{1 \leq i < j} (\delta(\beta_i) \cdot x_i + \delta(\alpha_i)) + \delta(\beta_j) \cdot (x_j - 1) + \delta(\beta_j[1, \ell]) \geq 0$$

if, and only if,

$$\sum_{1 \leq i \leq j-1} \delta(\beta_i) \cdot x_i \geq -u - \sum_{0 \leq i < j} \delta(\alpha_i) - \delta(\beta_j[1, \ell]) \quad (6)$$

$$\sum_{1 \leq i \leq j} \delta(\beta_i) \cdot x_i \geq -u - \sum_{0 \leq i < j} \delta(\alpha_i) + \delta(\beta_j) - \delta(\beta_j[1, \ell]) \quad (7)$$

By our construction, it is easily verified that for every $e = (e_1, \ldots, e_k) \in \mathbb{N}^k$ we have $\chi(i) = 1$ for all $i \in [1, k]$ and $p(\rho) \sqsubseteq_{\text{Diophantine}} q(\rho + \delta(\pi))$ if, and only if, $e$ fulfills all constraints required in (4), (5), (6) and (7). It thus remains to, informally speaking, extract the required system $T$ of linear Diophantine inequalities from those constraints.

For every fixed $j \in [1, k]$, by combining the constraints from (5), (6) and (7), we obtain systems of linear Diophantine inequalities $T_j^* : B_j \cdot x \geq d_j$ such that $B_j$ consists of at most $d$ different rows, since every $x_i$ is multiplied by the same $\delta(\beta_i)$. Let $A_j$ be the following $(d \times k)$-matrix:

$$A_j \overset{\text{def}}{=} \begin{bmatrix} \delta(\beta_1) \cdots \delta(\beta_j) & 0 & \cdots & 0 \end{bmatrix}$$

For the $i$-th row of $A_j$, let $e_{j,i} \in \mathbb{Z}$ be the maximum value in $d_j$ of the rows with the same coefficients in $T_j^*$. We define $e_{j} \overset{\text{def}}{=} (e_{j,1}, \ldots, e_{j,d})$ and set $T_j : A_j \cdot x \geq e_j$. By construction, we now have that $e \in \mathbb{N}^k$ is a solution of $T_j$ if, and only if, $e$ is a solution to $T_j^*$ and in particular fulfills all relevant constraints in (5), (6) and (7).

In order to obtain the matrix $A$ and $e$ required in the lemma, we define

$$A \overset{\text{def}}{=} \begin{bmatrix} I_k \\ A_1 \\ \vdots \\ A_k \end{bmatrix}$$

and

$$e \overset{\text{def}}{=} \begin{bmatrix} 1 \\ e_1 \\ \vdots \\ e_k \end{bmatrix}$$

The dimension of $A$ and $c$ is as required. It thus remains to estimate the norm of $A$ and $c$. We have $\|A\| \leq \sum_{1 \leq i \leq k} |\delta(\beta_i)| \leq k \cdot |\rho| \cdot |T|$. For $c$, the following inequality bounds the norm of the right-hand sides of (5), (6) and (7):

$$\|c\| \leq \|u\| + 2 \cdot |\rho| \cdot |T|.$$
some $\pi = \alpha_0 \beta_1^e \cdot \alpha_1 \cdot \ldots \cdot \beta_k^e \cdot \alpha_k$ where $e_1, \ldots, e_k \in [0, \ell]$, for some $e \in \mathbb{N}$ that is bounded by
\[
e \leq 2^{2Q^{\Theta(1)} \cdot O_{\left(\right)} \left(||u|| + ||v|| + (|Q| + ||T||)^{O(1)} \cdot ||T||\right)} \\
\leq 2^{(|V| + \log ||u|| + \log ||v||)^{O(1)}}.
\]
Since $|\pi| \leq |p| \cdot e$, the run $p(u) \xrightarrow{\tau_{\mathcal{N}^2}} q(v)$ can be guessed nondeterministically in polynomial space by storing only the intermediate configurations in an on-the-fly manner. Consequently, reachability in 2-VASS in PSPACE.

In order to complete the proof of Theorem 2, it remains to show hardness for PSPACE. We reduce from reachability in bounded one-counter automata, which is known to be PSPACE-complete [3]. A bounded one-counter automaton is given by a tuple $V = (Q, T, b)$, where $(Q, T)$ is a 1-VASS and $b \in \mathbb{N}$ is a bound encoded in binary. Let $\mathbb{B} = \{0, 1\}$, given configurations $p(u), q(v)$ of $V$ such that $u, v \in \mathbb{B}$, reachability is to decide whether $p(u) \xrightarrow{\tau_{\mathcal{N}^2}} q(v)$.

**Lemma 17.** Reachability in 2-VASS is PSPACE-hard.

**Proof.** Let $V = (Q, T, b)$ be a bounded one-counter automaton, and let $V' \overset{\text{def}}{=} (Q, T, b)$ be the 2-VASS obtained from $V$ by setting $T' \overset{\text{def}}{=} \{h(t) : t \in T\}$, where $h(p, z, q) \overset{\text{def}}{=} (p, (z, -z), q)$. We define an injection $\varphi$ from configurations of $V$ to configurations of $V'$ as $\varphi(q(z)) \overset{\text{def}}{=} q(z, b - z)$. It is now easily verified that $p(u) \xrightarrow{\tau_{\mathcal{N}^2}} q(v)$ if, and only if, $\varphi(p(u)) \xrightarrow{\tau_{\mathcal{N}^2}} \varphi(q(v))$.

**B. Reachability in 2-VASS with Unary Updates**

For unary 2-VASS we can show that reachability is in NP and NL-hard.

Given a unary 2-VASS $V$, whenever $p(u) \xrightarrow{\tau_{\mathcal{N}^2}} q(v)$ then by Theorem 1 there exists a linear path scheme $p = \alpha_0 \beta_1^e \cdot \alpha_1 \cdot \ldots \cdot \beta_k^e \cdot \alpha_k$ whose length is polynomial in $|V|$ such that $p(u) \xrightarrow{\tau_{\mathcal{N}^2}} q(v)$. Moreover, the proof of Corollary 16 shows that there exist $e_1, \ldots, e_k \leq 2^{(|V| + \log ||u|| + \log ||v||)^{O(1)}}$ such that for $\pi = \alpha_0 \beta_1^e \cdot \alpha_1 \cdot \ldots \cdot \beta_k^e \cdot \alpha_k$, we have $p(u) \xrightarrow{\tau_{\mathcal{N}^2}} q(v)$. In particular, every $e_i$ can be represented using a polynomial number of bits. Hence, $(p, e_1, \ldots, e_k)$ may serve as a certificate that can be guessed in polynomial time. It remains to show that this certificate can be verified in polynomial time. Checking that $p$ is a linear path scheme is easily verified in polynomial time. In order to check if $p(u) \xrightarrow{\tau_{\mathcal{N}^2}} q(v)$ in polynomial time we can construct the system of linear Diophantine equations from Lemma 14 and verify that $e = (e_1, \ldots, e_k)$ is a solution to this system. This shows that reachability in unary 2-VASS is in NP.

NL-hardness of reachability trivially follows from NL-hardness of reachability in directed graphs. Here, we wish to slightly strengthen this result and remark that reachability is NL-hard already for unary 2-VASS, whose underlying graph corresponds structurally to a linear path scheme (formally, every vertex lies on at most one cycle and the deletion of all cycles yields a union of isolated vertices and a cycle-free path, cf. Fig. 1 at the beginning of this paper). Let $G = (U, E)$ be a directed graph such that $U = \{u_0, \ldots, u_{m-1}\}$ and $E = \{e_0, \ldots, e_{n-1}\} \subseteq U \times U$. We define an injection $h : U \rightarrow [0, m - 1]^{\ell}$ as $h(u_i) = (i, m - 1 - i)$ that relates vertices of $G$ with vectors from bounded intervals. Let $\ell \overset{\text{def}}{=} m \cdot n - 1$. The flat unary 2-VASS $V = (Q, T)$ can now be defined as $Q \overset{\text{def}}{=} \{q_0, q_0', \ldots, q_n, q_n'\}$ and $T \overset{\text{def}}{=} \{(q_j, 0, q_{j+1}) : j \in [0, \ell - 1]\} \cup \{(q_j, -h(u_n), q_j), (q_j, h(u_n) - q_j, e_j) : n = (u_n, u_{n}) \text{ of } G\}$. The transition from $q_j$ to $q_j'$ can only be traversed if the vertex encoded into the current counter values corresponds to $u_n$. If we are able to reach $q_j'$, the transition back to $q_j$ then updates the currently visited vertex to $u_n$. Since a path from $u_0$ to $u_{m-1}$ of minimal length in $G$ traverses at most $m$ vertices, $\ell + 1 = m \cdot n$ states $q_j$ suffice.

**Theorem 18.** Reachability in unary 2-VASS is in NP and NL-hard.

**C. Derived Results**

Here, we explicitly remark some results that can additionally be derived from the technical results of this paper.

1) $\mathbb{Z}$-Reachability in Unary d-VASS is $\mathsf{NL}$-complete: The complexity of $\mathbb{Z}$-reachability in d-VASS depends on the encoding of numbers as well as the dimension $d$. When numbers are encoded in binary, reachability is NP-complete even when $d = 1$ [5], [7], and reachability is also NP-complete when numbers are encoded in unary and $d$ is part of the input to the problem [7]. We solve the case of reachability under unary encoding of numbers for each fixed dimension $d$.

**Theorem 19.** $\mathbb{Z}$-reachability in unary d-VASS is $\mathsf{NL}$-complete for any fixed $d \geq 1$.

**Proof.** NL-hardness trivially follows from NL-hardness of reachability in directed graphs. Let $d \geq 1$ be fixed and $V = (Q, T)$ be a unary d-VASS. Suppose $p(u) \xrightarrow{\tau_{\mathcal{Z}^d}} q(v)$, then by Proposition 3, there exists a linear path scheme $p = \alpha_0 \beta_1^e \cdot \alpha_1 \cdot \ldots \cdot \beta_k^e \cdot \alpha_k \in S$ with $k \leq |T|$ and $|p| \leq 2 \cdot |Q| \cdot |T|$ such that $p(u) \xrightarrow{\tau_{\mathcal{Z}^d}} q(v)$.

Let $E : A \cdot x = c$ be the system of linear Diophantine equations such that $A \overset{\text{def}}{=} \left[\delta(\beta_1) \ldots \delta(\beta_k)\right]$ and $c \overset{\text{def}}{=} v - (\alpha_0 \alpha_1 \ldots \alpha_k)$. Then, $p(u) \xrightarrow{\alpha_0 \beta_1^e \cdot \alpha_1 \cdot \ldots \beta_k^e \cdot \alpha_k \tau_{\mathcal{Z}^d}} q(v)$ if, and only if, $(e_1, e_2, \ldots, e_k) \in [E]$. By Corollary 13, if $[E] \neq \emptyset$ then $E$ has a solution $e$ such that $||e|| \leq (|T| + |T|)^{O(1)} + ||u|| + ||v||^{O(d)}$. Since $|T|$, $||u||$ and $||v||$ are encoded in unary and $d$ is fixed, minimal runs are bounded by some $b \leq |V|^{O(1)}$. Thus, reachability can be
decided by guessing a path of polynomial length on-the-fly in logarithmic space.

2) **Boundedness and Coverability in d-VASS:** For the sake of completeness, here we wish to discuss some consequences of PSPACE-hardness of reachability in 2-VASS to the complexity of coverability and boundedness in d-VASS that were left open in the literature. The boundedness problem is to determine, given \( p(u) \), whether \( \{ q(v) : p(u) \rightarrow_{\text{seq}} q(v) \} \) is infinite. The coverability problem is to determine, given \( p(u) \) and \( q(v) \), whether there exists \( w \geq v \) such that \( p(u) \rightarrow_{\text{seq}} q(w) \).

The complexity of boundedness and coverability for d-VASS in a fixed dimension \( d \) has been studied by Rosier and Yen in [22]. They show that both problems are PSPACE-complete for any fixed \( d \geq 4 \). Chan [2] later noted that boundedness is already PSPACE-complete for \( d = 3 \), leaving the case \( d = 2 \) as an open problem. It is moreover known that for \( d = 1 \) those problems are NP-complete [6]. From the results in [3] and Lemma 17, it is easy to show that both problems are PSPACE-complete for every fixed \( d \geq 2 \). An instance of reachability between \( p(u) \) and \( q(v) \) in a bounded one-counter automaton with bound \( b \) can be reduced to boundedness and coverability in 2-VASS by using the construction of Lemma 17 as a gadget and adding an extra transition \( (v, v - b) \) from \( q \) to a fresh control state \( r \) which has a self-loop \((1, 1)\).

**Corollary 20.** Boundedness and coverability in d-VASS are PSPACE-complete for any fixed \( d \geq 2 \).

**VI. Conclusion**

This paper established the precise complexity, i.e., PSPACE-completeness, of the reachability problem for 2-VASS. We also noted that the coverability and boundedness problems for 2-VASS are PSPACE-complete. When numbers are encoded in unary we showed that \( \mathbb{Z} \)-reachability in d-VASS is NL-complete for fixed \( d \). Reachability for unary 2-VASS was shown to be NL-hard and in NP. Our approach does not immediately lead to a better upper bound than NP mainly due to the following reason. Our proof showed that the reachability relation can be captured by a set of linear path schemes whose \(*\)-length is quadratic in the number of control states. The matrix of the resulting system of linear Diophantine inequalities thus has quadratically many columns and its smallest solutions — which corresponds to the exponents of the cycles of the linear path scheme and hence of the length of the path — can thus become exponentially large.

**References**


