# Improved Algorithm for Reachability in d-VASS 

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#### Abstract

An $\mathrm{F}_{d}$ upper bound for the reachability problem in vector addition systems with states (VASS) in fixed dimension is given, where $\mathrm{F}_{d}$ is the $d$-th level of the Grzegorczyk hierarchy of complexity classes. The new algorithm combines the idea of the linear path scheme characterization of the reachability in the 2-dimension VASSes with the general decomposition algorithm by Mayr, Kosaraju and Lambert. The result improves the $\mathrm{F}_{d+4}$ upper bound due to Leroux and Schmitz (LICS 2019).


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## 1 Introduction

Petri nets, or equivalently vector addition system with states (VASS), are a well studied model of concurrency. A VASS consists of a finite state control where each state transition has as its effect an integer-valued vector, and its configurations are pairs of a state and a vector with natural number components. A transition may lead from one configuration to another by adding its effect component-wise, conditioned on that the components of the resulting vector remain non-negative. The reachability problem, which asks whether from one configuration there is a path reaching another configuration, lies in the center of the algorithmic theory of Petri nets and has found a wide range of applications due to its generic nature. Since the problem was shown to be decidable by Mayr [17] in 1981, its computational complexity had been a long-standing open problem in the field. In 2015, Leroux and Schmitz [14] presented the first complexity upper bound, stating that the reachability problem is cubicAckermannian. This was later improved to an Ackermannian upper bound in 2019, again by Leroux and Schmitz [15]. Regarding the hardness, in 2021 seminal works by Czerwiński and Orlikowski [6], and independently by Leroux [13], provided matching Ackermannian lower bounds, settling the exact complexity of the problem.

[^0]Concerning the parameterization by dimension, i.e. the reachability problem in $d$ dimensional VASSes where $d$ is fixed, there is still a gap in the known complexity bounds. Currently, we only have exact complexity bounds for dimension one and two [7, 1]. For dimension $d \geq 3$, the result of Leroux and Schmitz [15] shows that the problem is in $\mathrm{F}_{d+4}$, the $(d+4)$-th level of the Grzegorczyk hierarchy of complexity classes. Note that a recent work by Yang and Fu [22] points out that the problem for dimension 3 is in $F_{3}=$ TOWER. On the other hand, the best known lower bound by Czerwiński, Jecker, Lasota, Leroux and Orlikowski [5] states that reachability in (2d+3)-dimensional VASSes is $\mathrm{F}_{d}$-hard. Motivated by this gap, our paper focuses on the computational complexity of reachability problem in the fixed-dimensional VASSes.

## Our contribution

In this paper we show that the reachability problem in the $d$-dimensional VASS is in $\mathrm{F}_{d}$ for $d \geq 3$, improving the previous $\mathrm{F}_{d+4}$ upper bound by Leroux and Schmitz [15], and generalizing the tower upper bound for the reachability problem in 3-VASS [22]. The new upper bound is obtained with the help of two novel technical lemmas.

1. Our main technical tool (Theorem 3.4) is a generalization of the linear path scheme characterization for the reachability relation in the 2-dimensional VASSes [1]. By borrowing the key idea from the work of Yang and Fu [22], we show that as long as the "geometric dimension" of a VASS (that is, the dimension of the vector space spanned by the effects of cyclic paths) is bounded by 2 , its reachability relation can be characterized by short linear path schemes. We then apply the lemma to simplify the KLMST algorithm so that (i) a VASS is replaced by a short linear path scheme whenever its geometric dimension is no more than 2 and (ii) the linear path schemes will not be decomposed further. It is then routine [15], using the tools from [21], to show that the reachability problem in the $d$-dimensional VASS is in $\mathrm{F}_{d+1}$ for all $d \geq 3$.
2. Our second lemma (Lemma 6.3) allows us to improve further the bound from $F_{d+1}$ to $F_{d}$. This is done by a careful analysis of the properties of the fast-growing functions [21].
Due to space limitation the proofs of the two lemmas are placed in the appendices.

## Organization

Section 2 fixes notation, defines the VASS model and its reachability problem. Section 3 generalizes the linear path scheme characterization [1] to VASSes whose geometric dimension are bounded by 2. Section 4 recalls the characterization system of linear inequalities for linear path schemes. Section 5 makes use of the results of Section 3 to give an improved version of the classic KLMST decomposition algorithm. Section 6 analyzes the complexity of our modified algorithm, proving the main result. Section 7 concludes. Proofs omitted from the main text can be found in the appendices.

## 2 Preliminaries

We use $\mathbb{N}, \mathbb{Z}, \mathbb{Q}$ to denote respectively the set of non-negative integers, integers, and rational numbers. Let $n \in \mathbb{N}$ be a number, we write $[n]$ for the range $\{1,2, \ldots, n\}$. Let $\boldsymbol{u}, \boldsymbol{v} \in \mathbb{X}^{d}$ be $d$ dimensional vectors where $\mathbb{X}$ can be any one of $\mathbb{N}, \mathbb{Z}, \mathbb{Q}$. We write $\boldsymbol{v}(i)$ for the $i$-th component of $\boldsymbol{v}$ where $i \in[d]$, so $\boldsymbol{v}=(\boldsymbol{v}(1), \ldots, \boldsymbol{v}(d))$. The maximum norm of $\boldsymbol{v}$ is defined to be $\|\boldsymbol{v}\|:=\max _{i \in[d]}|\boldsymbol{v}(i)|$. We extend component-wise the order $\leq$ for vectors in $\mathbb{X}^{d}$, so $\boldsymbol{u} \leq \boldsymbol{v}$ if
and only if $\boldsymbol{u}(i) \leq \boldsymbol{v}(i)$ for all $i \in[d]$. Addition and subtraction of vectors are also componentwise, so $(\boldsymbol{u}+\boldsymbol{v})(i)=\boldsymbol{u}(i)+\boldsymbol{v}(i)$ for all $i \in[d]$. Define $\operatorname{supp}(\boldsymbol{v}):=\{i \in[d]: \boldsymbol{v}(i) \neq 0\}$ to be the set of indices of non-zero components of $\boldsymbol{v}$. This notation is extended to sets of vectors naturally, so $\operatorname{supp}(S)=\bigcup_{\boldsymbol{v} \in S} \operatorname{supp}(\boldsymbol{v})$ for any set $S \subseteq \mathbb{X}^{d}$.

For technical reasons we introduce the symbol $\omega$ that stands for the infinite element. Let $\mathbb{N}_{\omega}:=\mathbb{N} \cup\{\omega\}$. We stipulate that $n<\omega$ for all $n \in \mathbb{N}$, and $x+\omega=\omega+x=\omega$ for all $x \in \mathbb{Z}$. Define the partial order $\sqsubseteq$ over $\mathbb{N}_{\omega}$ so that $x \sqsubseteq y$ if and only if $x=y$ or $y=\omega$ for all $x, y \in \mathbb{N}_{\omega}$. The relation $\sqsubseteq$ is also extended component-wise to vectors in $\mathbb{N}_{\omega}^{d}$.

Let $\Sigma$ be a finite alphabet and $s, t \in \Sigma^{*}$ be two strings over $\Sigma$. We write st for the concatenation of $s$ and $t$, and $s^{n}$ for the concatenation of $n$ copies of $s$ where $n \in \mathbb{N}$. If $s=a_{1} a_{2} \ldots a_{\ell}$ where $a_{1}, \ldots, a_{\ell} \in \Sigma$, we write $|s|:=\ell$ for the length of $s$, and $s[i \ldots j]:=$ $a_{i} a_{i+1} \ldots a_{j}$ for the substring of $s$ between indices $i$ and $j$ where $1 \leq i \leq j \leq|s|$.

### 2.1 Vector Addition Systems with States

Let $d \geq 0$ be an integer. A $d$-dimensional vector addition system with states ( $d$-VASS) is a pair $G=(Q, T)$ where $Q$ is a finite set of states and $T \subseteq Q \times \mathbb{Z}^{d} \times Q$ is a finite set of transitions. Clearly a VASS can also be viewed as a directed graph with edges labelled by integer vectors. Given a word $\pi=\left(p_{1}, \boldsymbol{a}_{1}, q_{1}\right)\left(p_{2}, \boldsymbol{a}_{2}, q_{2}\right) \ldots\left(p_{n}, \boldsymbol{a}_{n}, q_{n}\right) \in T^{*}$ over transitions, we say that $\pi$ is a path from $p_{1}$ to $q_{n}$ if $q_{i}=p_{i+1}$ for all $i=1, \ldots, n-1$. It is a cycle if we further have $p_{1}=q_{n}$. The effect of $\pi$ is defined to be $\Delta(\pi):=\sum_{i=1}^{n} \boldsymbol{a}_{i}$, and the action word of $\pi$ is the word $\llbracket \pi \rrbracket:=\boldsymbol{a}_{1} \boldsymbol{a}_{2} \ldots \boldsymbol{a}_{n}$ over $\mathbb{Z}^{d}$. The Parikh image of $\pi$ is the function $\phi \in \mathbb{N}^{T}$ mapping each transition to its number of occurrences in $\pi$. Given a function $\phi \in \mathbb{N}^{T}$ we also define $\Delta(\phi):=\sum_{t=(p, \boldsymbol{a}, q) \in T} \phi(t) \cdot \boldsymbol{a}$. Note that $\Delta(\phi)=\Delta(\pi)$ if $\phi$ is the Parikh image of $\pi$. Let $L \subseteq T^{*}$ be a language (i.e. subset of words), we define its effect as $\Delta(L):=\{\Delta(\pi): \pi \in L\}$.

The norm of a transition $t=(p, \boldsymbol{a}, q)$ is defined by $\|t\|:=\|\boldsymbol{a}\|$. The norm of a path $\pi=t_{1} t_{2} \ldots t_{n}$ is $\|\pi\|:=\max _{i \in[n]}\left\|t_{i}\right\|$. For a $d$-VASS $G=(Q, T)$ we write $\|T\|:=\max \{\|t\|:$ $t \in T\}$. The size of $G$ is defined by

$$
\begin{equation*}
|G|:=|Q|+|T|+d \cdot|T| \cdot\|T\| \tag{1}
\end{equation*}
$$

## Semantics of VASSes

Let $G=(Q, T)$ be a $d$-VASS. A configuration of $G$ is a pair of a state $p \in Q$ and a vector $\boldsymbol{v} \in \mathbb{Z}^{d}$, written as $p(\boldsymbol{v})$. Let $\mathbb{D} \subseteq \mathbb{Z}^{d}$, we define the $\mathbb{D}$-semantics for $G$ as follows. For each transition $t=(p, \boldsymbol{a}, q) \in T$, the one-step transition relation $\xrightarrow{t}{ }_{\mathbb{D}}$ relates all pairs of configurations of the form $(p(\boldsymbol{u}), q(\boldsymbol{v}))$ where $\boldsymbol{u}, \boldsymbol{v} \in \mathbb{D}$ and $\boldsymbol{v}=\boldsymbol{u}+\boldsymbol{a}$. Then for a word $\pi=t_{1} t_{2} \ldots t_{n} \in T^{*}$, the relation $\xrightarrow{\pi} \mathbb{D}$ is the composition $\xrightarrow{\pi_{\mathbb{D}}}:={\xrightarrow{t_{1}}}_{\mathbb{D}} \circ \cdots \circ{\xrightarrow{t_{n}}}_{\mathbb{D}}$. So $p(\boldsymbol{u}) \xrightarrow{\pi} q(\boldsymbol{v})$ if and only if there are configurations $p_{0}\left(\boldsymbol{u}_{0}\right), \ldots, p_{n}\left(\boldsymbol{u}_{n}\right) \in Q \times \mathbb{D}$ such that

$$
\begin{equation*}
p(\boldsymbol{u})=p_{0}\left(\boldsymbol{u}_{0}\right) \xrightarrow{t_{1}} \mathbb{D} p_{1}\left(\boldsymbol{u}_{1}\right) \xrightarrow{t_{2}} \mathbb{D} \cdots \xrightarrow{t_{n}} p_{n}\left(\boldsymbol{u}_{n}\right)=q(\boldsymbol{v}) . \tag{2}
\end{equation*}
$$

Also, when $\pi=\epsilon$ is the empty word, the relation $\xrightarrow{\epsilon_{\mathbb{D}}}$ is the identity relation over $Q \times \mathbb{D}$. Note that $\xrightarrow{\pi} \mathbb{D}$ is non-empty only if $\pi$ is a path. When $p(\boldsymbol{u}) \xrightarrow{\pi} q(\boldsymbol{v})$ we also say that $\pi$ induces (or is) a $\mathbb{D}$-run from $p(\boldsymbol{u})$ to $q(\boldsymbol{v})$. We emphasize that all configurations on this run lie in $\mathbb{D}$, and that they are uniquely determined by $p(\boldsymbol{u})$ and $\pi$. For a language $L \subseteq T^{*}$ we define $\xrightarrow{L}{ }_{\mathbb{D}}$ as $\bigcup_{\pi \in L} \xrightarrow{\pi} \mathbb{D}$. Finally, the $\mathbb{D}$-reachability relation of $G$ is defined to be $\xrightarrow{*} \mathbb{D}:=\xrightarrow{T^{*}} \mathbb{D}$.

For the above definitions, we shall often omit the subscript $\mathbb{D}$ when $\mathbb{D}=\mathbb{N}^{d}$.

We mention that in Section 5 we need to generalize the VASS semantics to configurations in $Q \times \mathbb{N}_{\omega}^{d}$, allowing $\omega$ components in vectors. The definitions of $\stackrel{t}{\rightarrow}_{\mathbb{N}_{\omega}^{d}}, \xrightarrow{\pi}_{\mathbb{N}_{\omega}^{d}}$, and $\xrightarrow{*}_{\mathbb{N}_{\omega}^{d}}$ are similar to the above.

## Reachability problem

The reachability problem in vector addition systems with states is formulated as follows:
Reachability in $d$-Dimensional Vector Addition System with States
Input: A $d$-dimensional VASS $G=(Q, T)$, two configurations $p(\boldsymbol{u}), q(\boldsymbol{v}) \in Q \times \mathbb{N}^{d}$. Question: Does $p(\boldsymbol{u}) \stackrel{*}{\rightarrow}_{\mathbb{N}^{d}} q(\boldsymbol{v})$ hold in $G$ ?

Note that we study the reachability problem for fixed-dimensional VASSes, where the dimension $d$ is treated as a constant to allow more fine-grained analysis. So we shall use the big- $O$ notation to hide constants that may depend on $d$. The general problem where the dimension can be part of the input was already shown to be Ackermann-complete [6, 13].

## Cycle spaces and geometric dimensions

One of the key insights of [15] is a new termination argument for the KLMST decomposition algorithm based on the dimensions of vector spaces spanned by cycles in VASSes, which yielded the primitive recursive upper bound of VASS reachability problem in fixed dimensions. The vector spaces spanned by cycles also play an important role in our work.

- Definition 2.1. Let $G$ be a d-VASS. The cycle space of $G$ is the vector space $\operatorname{Cyc}(G) \subseteq \mathbb{Q}^{d}$ spanned by the effects of all cycles in $G$, that is:

$$
\begin{equation*}
\operatorname{Cyc}(G):=\operatorname{span}\{\Delta(\beta): \beta \text { is a cycle in } G\} \tag{3}
\end{equation*}
$$

The dimension of the cycle space of $G$ is called the geometric dimension of $G$. We say $G$ is geometrically $k$-dimensional if $\operatorname{dim}(\operatorname{Cyc}(G)) \leq k$.

## 3 Flattability of Geometrically 2-dimensional VASSes

A VASS is flat if each of its states lies on at most one cycle. Flat VASSes form an important subclass of VASSes due to its connection to Presburger arithmetic, and we refer the readers to [12] for a survey. In dimension 2, it was proved in [1] that 2-VASSes enjoy a stronger form of flat representation, known as the linear path schemes. A linear path scheme is a regular expression of the form $\alpha_{0} \beta_{1}^{*} \alpha_{1} \ldots \beta_{k}^{*} \alpha_{k}$ where $\alpha_{0}, \ldots \alpha_{k}$ are paths and $\beta_{1}, \ldots, \beta_{k}$ are cycles of the VASS, such that they form a path when joined together. The results of [1] show that the reachability relation of every 2-VASS can be captured by short linear path schemes.

Linear path schemes are extremely useful as they can be fully characterized by linear inequality systems so that standard tools for linear or integer programming can be applied. In this section, we generalize the results in [1] and show that the reachability relation of any $d$-VASS whose geometric dimension is bounded by 2 can also be captured by short linear path schemes.

Our proof follows closely to the lines of [1]. Given a geometrically 2-dimensional VASS $G$, we first cover $\mathbb{N}^{d}$ by the following two regions: one for the region far away from every axis:

$$
\begin{equation*}
\mathbb{O}:=\left\{\boldsymbol{u} \in \mathbb{N}^{d}: \boldsymbol{u}(i) \geq D \text { for all } i \in[d]\right\} \tag{4}
\end{equation*}
$$





Figure 1 A geometrically 2-dimensional VASS $G$ and two possible projections of it.
where $D$ is some properly chosen threshold; the other one for the region close to some axis:

$$
\begin{equation*}
\mathbb{L}:=\left\{\boldsymbol{u} \in \mathbb{N}^{d}: \boldsymbol{u}(i) \leq D^{\prime} \text { for some } i \in[d]\right\} \tag{5}
\end{equation*}
$$

where $D^{\prime}$ is chosen slightly above $D$ to create an overlap with $\mathbb{O}$. Let $\pi$ be a run that we are going to capture by linear path schemes. We can extract its maximal prefix that lies completely in either $\mathbb{O}$ or $\mathbb{L}$, depending on where $\pi$ starts. This prefix must end at a configuration that lies in $\mathbb{L} \cap \mathbb{O}$, if we haven't touched the end of $\pi$. From this configuration we then extract a maximal cycle that also ends in $\mathbb{L} \cap \mathbb{O}$. Continuing this fashion, we can break $\pi$ into segments of runs that lie completely in $\mathbb{O}$ or $\mathbb{L}$, interleaved by cycles that start and end in $\mathbb{O}$ (actually in $\mathbb{L} \cap \mathbb{O})$. Note that the number of such cycles cannot exceed the number of states of $G$ since they are maximal. Now we only need to capture the following three types of runs by short linear path schemes:

1. Runs that are cycles starting and ending in $\mathbb{O}$.

This will be handled in Section A. 4 in the appendix. Briefly speaking, since the geometric dimension of $G$ is 2 , the effect of such a cycle must belong to a plane in $\mathbb{Z}^{d}$. We will find a clever way to project this plane to a coordinate plane, and then project the $d$-VASS $G$ onto this plane to get a 2-VASS. This is made possible by a novel technique called the "sign-reflecting projection" developed in Section A.3. Intuitively speaking, for any vector in a plane we are able to determine whether it belongs to a certain orthant by observing only 2 entries of this vector. The $d$-VASS can then be projected onto these 2 coordinates.
(See Example 3.1 for a more concrete demonstration.) Now we apply the results of [1] to obtain a linear path scheme that captures the projected cyclic run. Combined with a lemma in [16], the projection guarantees that we can safely project it back to get a linear path scheme for the run in $G$.
2. Runs that lie completely in $\mathbb{O}$.

This is just an easy corollary of the first type, and will also be handled in Section A.4. Just note that any run can be broken into a series of simple paths interleaved by cycles.
3. Runs that lie completely in $\mathbb{L}$.

This will be handled in Section A.5, by a long and tedious case analysis. In principle, we are going to argue that any such run essentially corresponds to a run in some ( $d-1$ )-VASS, so that we can use induction.

- Example 3.1. Consider the geometrically 2-dimensional 3-VASS $G$ as shown in Figure 1. It consists of a single state $p$ and two transitions with effects $(1,0,-1)$ and $(0,1,1)$. In order to apply the results of [1], one would like to derive a 2 -VASS from $G$ that reflects runs in $G$. A simple idea is to discard one coordinate of $G$. Two possibilities of this idea are shown in Figure 1 as $G^{\prime}$, which discards the third coordinate, and $G^{\prime \prime}$, which discards the second coordinate. However, not all of them are satisfactory. For example, the legal run $p(0,0) \xrightarrow{(1,0)} p(1,0)$ in $G^{\prime}$ reflects an illegal run $p(0,0,0) \xrightarrow{(1,0,-1)} p(1,0,-1)$ in $G$ where the third coordinate goes negative. On the other hand, all runs in $G^{\prime \prime}$ can be safely projected back to a run in $G$. To see this, just observe that for any vector $\boldsymbol{v}$ in the linear span of $(1,0,-1)$ and $(0,1,1), \boldsymbol{v} \geq \mathbf{0}$ if and only if $\boldsymbol{v}(1) \geq 0$ and $\boldsymbol{v}(3) \geq 0$. Thus we can safely discard the second coordinate.

In general, the "sign-reflecting projection" developed in Section A. 3 shows that any geometrically 2-dimensional VASS can be projected onto two coordinates so that the signs of these two coordinates reflects the signs of other coordinates.

In the rest of this section we just state formally our main technical results. The detailed proofs are placed in the appendix.

### 3.1 Linear Path Schemes

A linear path scheme $(L P S)$ is a pair $(G, \Lambda)$ where $G$ is a VASS and $\Lambda$ is a regular expression of the form $\Lambda=\alpha_{0} \beta_{1}^{*} \alpha_{1} \ldots \beta_{k}^{*} \alpha_{k}$ such that $\alpha_{0}, \ldots, \alpha_{k}$ are paths in $G$ and $\beta_{1}, \ldots, \beta_{k}$ are cycles in $G$, and $\alpha_{0} \beta_{1} \alpha_{1} \ldots \beta_{k} \alpha_{k}$ is a path in $G$. We say an $\operatorname{LPS}(G, \Lambda)$ is compatible to a VASS $G^{\prime}$ if $G^{\prime}$ contains all states and transitions in $\Lambda$. Very often we shall omit the VASS $G$ and say that $\Lambda$ on its own is an LPS, with $G$ understood as any VASS to which $\Lambda$ is compatible. We write $|\Lambda|=\left|\alpha_{0} \beta_{1} \alpha_{1} \ldots \beta_{k} \alpha_{k}\right|$ for the length of $\Lambda,\|\Lambda\|=\left\|\alpha_{0} \beta_{1} \alpha_{1} \ldots \beta_{k} \alpha_{k}\right\|$ for its norm, and $|\Lambda|_{*}=k$ for the number of cycles in $\Lambda$.

We also treat $\Lambda$ as the language defined by it, and thus for two configurations $p(\boldsymbol{u})$ and $q(\boldsymbol{v})$ we write $p(\boldsymbol{u}) \xrightarrow{\Lambda}_{\mathbb{D}} q(\boldsymbol{v})$ if and only if there exists $e_{1}, \ldots, e_{k} \in \mathbb{N}$ such that

$$
\begin{equation*}
p(\boldsymbol{u}) \xrightarrow{\alpha_{0} \beta_{1}^{e_{1}} \alpha_{1} \ldots \beta_{k}^{e_{k}} \alpha_{k}} \mathbb{D} q(\boldsymbol{v}) . \tag{6}
\end{equation*}
$$

## Positive linear path schemes

A positive LPS is a regular expression of the form $\Lambda^{+}=\alpha_{0} \beta_{1}^{+} \alpha_{1} \ldots \beta_{k}^{+} \alpha_{k}$ which is similar to a linear path scheme except that we require each cycle to be used at least once. We write $p(\boldsymbol{u}) \xrightarrow{\Lambda^{+}} \mathbb{D} q(\boldsymbol{v})$ if and only if there are positive integers $e_{1}, \ldots, e_{k} \in \mathbb{N}_{>0}$ such that

$$
\begin{equation*}
p(\boldsymbol{u}) \xrightarrow{\alpha_{0} \beta_{1}^{e_{1}} \alpha_{1} \ldots \beta_{k}^{e_{k}} \alpha_{k}} \mathbb{D} q(\boldsymbol{v}) . \tag{7}
\end{equation*}
$$

A path $\pi$ is said to be admitted by $\Lambda^{+}$if $\pi=\alpha_{0} \beta_{1}^{e_{1}} \alpha_{1} \ldots \beta_{k}^{e_{k}} \alpha_{k}$ for some $e_{1}, \ldots, e_{k} \in \mathbb{N}_{>0}$. We prefer positive LPSes as they can be easily characterized by linear inequality systems (see Section 4 for details). In fact, positive LPSes can be easily obtained from LPSes:

- Lemma 3.2. For every linear path scheme $\Lambda$ there exists a finite set $S$ of positive linear path schemes compatible to the same VASSes with $\Lambda$, such that $\xrightarrow{\Lambda}=\bigcup_{\Lambda^{+} \in S} \xrightarrow{\Lambda^{+}}$and $\left|\Lambda^{+}\right| \leq|\Lambda|$ for every $\Lambda^{+} \in S$.
Proof. Suppose $\Lambda=\alpha_{0} \beta_{1}^{*} \alpha_{1} \ldots \beta_{k}^{*} \alpha_{k}$. For each cycle component $\beta_{i}^{*}$ in $\Lambda$ we replace it by either $\beta_{i}^{+}$or an empty word nondeterministically. Let $S$ be the set of all resulting positive LPSes. It is obvious that $S$ satisfies the desired properties.


### 3.2 Main Results

The main results of this section is stated as follows.

- Theorem 3.3. Let $G=(Q, T)$ be a geometrically 2-dimensional d-VASS. For every pair of configurations $p(\boldsymbol{u}), q(\boldsymbol{v}) \in Q \times \mathbb{N}^{d}$ with $p(\boldsymbol{u}) \xrightarrow{*} q(\boldsymbol{v})$ there exists a positive LPS $\Lambda^{+}$ compatible to $G$ such that $p(\boldsymbol{u}) \xrightarrow{\Lambda^{+}} q(\boldsymbol{v})$ and $\left|\Lambda^{+}\right| \leq|G|^{O(1)}$.

We remark that the big- $O$ term here and elsewhere in the paper hides constant that may depend on the dimension $d$, but does not depend on $G, \boldsymbol{u}, \boldsymbol{v}$ or anything else.

By Lemma 3.2 we know that positive LPSes can be obtained from LPSes. Thus theorem 3.3 follows easily from the following relaxed theorem, which will be proved in the appendix.

- Theorem 3.4. Let $G=(Q, T)$ be a geometrically 2-dimensional d-VASS. For every pair of configurations $p(\boldsymbol{u}), q(\boldsymbol{v}) \in Q \times \mathbb{N}^{d}$ with $p(\boldsymbol{u}) \xrightarrow{*} q(\boldsymbol{v})$ there exists an LPS $\Lambda$ compatible to $G$ such that $p(\boldsymbol{u}) \xrightarrow{\Lambda} q(\boldsymbol{v})$ and $|\Lambda| \leq|G|^{O(1)}$.


## 4 Characteristic Systems for Linear Path Schemes

The property that linear path schemes can be fully characterized by linear inequality systems is exploited in [2] to derive the PSPACE upper bound of the reachability problem in 2-VASSes. Here we recall this linear inequality system and its properties.

We mainly focus on positive linear path schemes. Fix $\Lambda=\alpha_{0} \beta_{1}^{+} \alpha_{1} \ldots \beta_{k}^{+} \alpha_{k}$ to be a positive LPS from state $p$ to $q$ that is compatible to some $d$-VASS $G=(Q, T)$, where $k=|\Lambda|_{*}$ is the number of cycles in $\Lambda$.

- Definition 4.1 (cf. [2, Lem. 14]). The characteristic system $\mathcal{E}_{\text {LPS }}(\Lambda)$ of the positive LPS $\Lambda$ is the system of linear inequalities such that a triple $\boldsymbol{h}=(\boldsymbol{u}, \boldsymbol{e}, \boldsymbol{v}) \in \mathbb{N}^{d} \times \mathbb{N}^{k} \times \mathbb{N}^{d}$ satisfies $\mathcal{E}_{\mathrm{LPS}}(\Lambda)$, written $\boldsymbol{h} \models \mathcal{E}_{\mathrm{LPS}}(\Lambda)$, if and only if the following conditions hold:

1. for every $i=1, \ldots, k, \boldsymbol{e}(i) \geq 1$;
2. for every $i=0, \ldots, k$ and every $j=1, \ldots,\left|\alpha_{i}\right|$,

$$
\begin{equation*}
\boldsymbol{u}+\Delta\left(\alpha_{0} \beta_{1}^{\boldsymbol{e}(1)} \alpha_{1} \ldots \alpha_{i-1} \beta_{i}^{\boldsymbol{e}(i)}\right)+\Delta\left(\alpha_{i}[1 \ldots j]\right) \geq \mathbf{0} \tag{8}
\end{equation*}
$$

3. for every $i=1, \ldots, k$ and every $j=1, \ldots,\left|\beta_{i}\right|$,

$$
\begin{array}{ll}
\boldsymbol{u}+\Delta\left(\alpha_{0} \beta_{1}^{\boldsymbol{e}(1)} \alpha_{1} \ldots \beta_{i-1}^{e(i-1)} \alpha_{i-1}\right) & +\Delta\left(\beta_{i}[1 \ldots j]\right) \geq \mathbf{0} \\
\boldsymbol{u}+\Delta\left(\alpha_{0} \beta_{1}^{\boldsymbol{e}(1)} \alpha_{1} \ldots \beta_{i-1}^{\boldsymbol{e}(i-1)} \alpha_{i-1} \beta_{i}^{e(i)-1}\right)+\Delta\left(\beta_{i}[1 \ldots j]\right) \geq \mathbf{0} \tag{10}
\end{array}
$$

4. and finally, $\boldsymbol{u}+\Delta\left(\alpha_{0} \beta_{1}^{\boldsymbol{e}(1)} \alpha_{1} \ldots \beta_{k}^{\boldsymbol{e}(k)} \alpha_{k}\right)=\boldsymbol{v}$.

The readers can easily verify that these constraints are indeed linear in terms of $\boldsymbol{u}, \boldsymbol{e}, \boldsymbol{v}$. The next lemma shows that $\mathcal{E}_{\text {LPS }}(\Lambda)$ indeed captures all runs admitted by $\Lambda$.

- Lemma 4.2. Let $G$ be a d-VASS and $\Lambda=\alpha_{0} \beta_{1}^{+} \alpha_{1} \ldots \beta_{k}^{+} \alpha_{k}$ be a positive LPS from state $p$ to $q$ compatible to $G$. Then for every $\boldsymbol{u}, \boldsymbol{v} \in \mathbb{N}^{d}, p(\boldsymbol{u}) \xrightarrow{\Lambda} q(\boldsymbol{v})$ if and only if there exists $\boldsymbol{e} \in \mathbb{N}^{k}$ such that $(\boldsymbol{u}, \boldsymbol{e}, \boldsymbol{v}) \models \mathcal{E}_{\mathrm{LPS}}(\Lambda)$. Moreover, for every $\boldsymbol{u}, \boldsymbol{v} \in \mathbb{N}^{d}$ and every $\boldsymbol{e} \in \mathbb{N}^{k}$ such that $(\boldsymbol{u}, \boldsymbol{e}, \boldsymbol{v}) \models \mathcal{E}_{\mathrm{LPS}}(\Lambda)$, we have $p(\boldsymbol{u}) \xrightarrow{\alpha_{0} \beta_{1}^{e(1)} \alpha_{1} \ldots \beta_{k}^{e(k)} \alpha_{k}} q(\boldsymbol{v})$.

We also need to introduce the homogeneous version of $\mathcal{E}_{\text {LPS }}(\Lambda)$ for technical reasons.

- Definition 4.3. The homogeneous characteristic system $\mathcal{E}_{\text {LPS }}^{0}(\Lambda)$ of $\Lambda$ is the system of linear inequalities such that a triple $\boldsymbol{h}_{0}=\left(\boldsymbol{u}_{0}, \boldsymbol{e}_{0}, \boldsymbol{v}_{0}\right) \in \mathbb{N}^{d} \times \mathbb{N}^{k} \times \mathbb{N}^{d}$ satisfies $\mathcal{E}_{\mathrm{LPS}}^{0}(\Lambda)$, written $\boldsymbol{h}_{0} \models \mathcal{E}_{\mathrm{LPS}}^{0}(\Lambda)$, if and only if the following conditions hold:

1. for every $i=0, \ldots, k, \boldsymbol{u}_{0}+\Delta\left(\beta_{1}\right) \cdot \boldsymbol{e}_{0}(1)+\cdots+\Delta\left(\beta_{i}\right) \cdot \boldsymbol{e}_{0}(i) \geq \mathbf{0}$;
2. $\boldsymbol{u}_{0}+\Delta\left(\beta_{1}\right) \cdot \boldsymbol{e}_{0}(1)+\cdots+\Delta\left(\beta_{k}\right) \cdot \boldsymbol{e}_{0}(k)=\boldsymbol{v}_{0}$.

## 5 The Modified KLMST Decomposition Algorithm

In this section we apply our results of Section 3 to improve the notoriously hard KLMST decomposition algorithm for VASS reachability. Our narration will base on the work of Leroux and Schmitz [15]. For readers familiar with [15], the major modifications are listed below:

- The decomposition structure is now a sequence of generalized VASS reachability instances linked by (positive) linear path schemes rather than by single transitions.
- We introduce a new "cleaning" step that replaces all VASS instances which are geometrically 2 -dimensional by polynomial-length linear path schemes compatible to them.
- We do not guarantee the exact preservation of action languages at each decomposition step. Instead, we only preserve a subset of action languages. This is a compromise since linear path schemes capture only the reachability relation but not every possible run. Nonetheless, it is enough for the reachability problem.
In this section we focus on the effectiveness and correctness of the modified KLMST decomposition algorithm. Its complexity will be analyzed in Section 6.


### 5.1 Linear KLM Sequences

The underlying decomposition structure in the KLMST algorithm was known as KLM sequences, named after Mayr[17], Kosaraju[10], and Lambert[11].

- Definition 5.1. A KLM tuple of dimension d is a tuple $\langle p(\boldsymbol{x}) G q(\boldsymbol{y})\rangle$ where $G=(Q, T)$ is a d-VASS and $p(\boldsymbol{x}), q(\boldsymbol{y}) \in Q \times \mathbb{N}_{\omega}^{d}$ are two (generalized) configurations of $G$. A KLM sequence of dimension $d$ is a sequence of KLM tuples interleaved by transitions of the form

$$
\begin{equation*}
\xi=\left\langle p_{0}\left(\boldsymbol{x}_{0}\right) G_{0} q_{0}\left(\boldsymbol{y}_{0}\right)\right\rangle t_{1}\left\langle p_{1}\left(\boldsymbol{x}_{1}\right) G_{1} q_{1}\left(\boldsymbol{y}_{1}\right)\right\rangle t_{2} \ldots t_{k}\left\langle p_{k}\left(\boldsymbol{x}_{k}\right) G_{k} q_{k}\left(\boldsymbol{y}_{k}\right)\right\rangle \tag{11}
\end{equation*}
$$

where each tuple $\left\langle p_{i}\left(\boldsymbol{x}_{i}\right) G_{i} q_{i}\left(\boldsymbol{y}_{i}\right)\right\rangle$ is a KLM tuple of dimensiond and each $t_{i}$ is a transition of the form $\left(q_{i-1}, \boldsymbol{a}_{i}, p_{i}\right)$ from state $q_{i-1}$ to $p_{i}$ with effect $\boldsymbol{a}_{i} \in \mathbb{Z}^{d}$.

In this paper we generalize the definition of KLM sequences to allow (positive) linear path schemes to connect KLM tuples.

- Definition 5.2. A linear KLM sequence of dimension $d$ is a sequence

$$
\begin{equation*}
\xi=\left\langle p_{0}\left(\boldsymbol{x}_{0}\right) G_{0} q_{0}\left(\boldsymbol{y}_{0}\right)\right\rangle \Lambda_{1}\left\langle p_{1}\left(\boldsymbol{x}_{1}\right) G_{1} q_{1}\left(\boldsymbol{y}_{1}\right)\right\rangle \Lambda_{2} \ldots \Lambda_{k}\left\langle p_{k}\left(\boldsymbol{x}_{k}\right) G_{k} q_{k}\left(\boldsymbol{y}_{k}\right)\right\rangle \tag{12}
\end{equation*}
$$

where each tuple $\left\langle p_{i}\left(\boldsymbol{x}_{i}\right) G_{i} q_{i}\left(\boldsymbol{y}_{i}\right)\right\rangle$ is a KLM tuple of dimensiond and each $\Lambda_{i}$ is a positive linear path scheme from state $q_{i-1}$ to $p_{i}$.

One immediately sees that KLM sequences are just special cases of linear KLM sequences. Let $\xi$ be a linear KLM sequence given as (12), we write $\xi_{i}:=\left\langle p_{i}\left(\boldsymbol{x}_{i}\right) G_{i} q_{i}\left(\boldsymbol{y}_{i}\right)\right\rangle$ for the $i$ th KLM tuple occurring in $\xi$.

## Action languages

Let $\xi$ be a linear KLM sequence given as (12). We say a path $\pi$ from state $p_{0}$ to $q_{k}$ is admitted by $\xi$, written $\xi \vdash \pi$, if $\pi$ can be written as $\pi=\pi_{0} \rho_{1} \pi_{1} \ldots \rho_{k} \pi_{k}$ where $\pi_{i}$ is a path from $p_{i}$ to $q_{i}$ in $G_{i}$ for each $i=0, \ldots, k$, and $\rho_{i}$ is a path admitted by $\Lambda_{i}$ for each $i=1, \ldots, k$, such that there exist vectors $\boldsymbol{m}_{0}, \boldsymbol{n}_{0}, \ldots, \boldsymbol{m}_{k}, \boldsymbol{n}_{k} \in \mathbb{N}^{d}$ such that

$$
\begin{equation*}
p_{0}\left(\boldsymbol{m}_{0}\right) \xrightarrow{\pi_{0}} q_{0}\left(\boldsymbol{n}_{0}\right) \xrightarrow{\rho_{1}} p_{1}\left(\boldsymbol{m}_{1}\right) \xrightarrow{\pi_{1}} q_{1}\left(\boldsymbol{n}_{1}\right) \xrightarrow{\rho_{2}} \cdots \xrightarrow{\rho_{k}} p_{k}\left(\boldsymbol{m}_{k}\right) \xrightarrow{\pi_{k}} q_{k}\left(\boldsymbol{n}_{k}\right) \tag{13}
\end{equation*}
$$

and that $\boldsymbol{m}_{i} \sqsubseteq \boldsymbol{x}_{i}, \boldsymbol{n}_{i} \sqsubseteq \boldsymbol{y}_{i}$ for each $i=0, \ldots, k$.
The action language $L_{\xi}$ of $\xi$ is the language over $\mathbb{Z}^{d}$ defined by $L_{\xi}:=\{\llbracket \pi \rrbracket: \xi \vdash \pi\}$, where we recall that $\llbracket \cdot \rrbracket$ is the word morphism mapping each transition to its effect.

We are more interested in the action languages because in some decomposition steps we have to modify the set of transitions, and only the action word of admitted runs can be preserved. Notice that action languages preserve not only the effects of admitted runs, but also their lengths.

## Ranks and Sizes

Let $t$ be a transition in a $d$-VASS $G$, we define $\operatorname{Cyc}(G / t)$ to be the vector space spanned by the effects of all cycles in $G$ containing $t$. For the VASS $G$, let $r_{i}$ be the number of transitions $t$ in $G$ such that $\operatorname{dim}(\operatorname{Cyc}(G / t))=i$ for each $i=0, \ldots, d$. Then the rank of $G$ is defined as $\operatorname{rank}(G)=\left(r_{d}, \ldots, r_{3}\right) \in \mathbb{N}^{d-2}$. We also define the full rank of $G$ as $\operatorname{rank}_{\text {full }}(G)=\left(r_{d}, \ldots, r_{0}\right) \in \mathbb{N}^{d+1}$.

The following lemma was proved in [15], which shows that in a strongly connected VASS $G$, the space $\operatorname{Cyc}(G / t)$ corresponds to $\operatorname{Cyc}(G)$.

- Lemma 5.3 ([15, Lem. 3.2]). Let $t$ be a transition of a strongly connected VASS G. Then $\operatorname{Cyc}(G / t)=\operatorname{Cyc}(G)$.

The following corollary is immediate.

- Corollary 5.4. Let $G$ be a strongly connected d-VASS. Then $\operatorname{rank}(G)=\mathbf{0}$ if and only if $G$ is geometrically 2-dimensional.

Let $\xi$ be a linear KLM sequence given as (12). We define the rank of $\xi$ as $\operatorname{rank}(\xi)=$ $\sum_{i=0}^{k} \operatorname{rank}\left(G_{i}\right)$, and the full rank of $\xi$ as $\operatorname{rank}_{\text {full }}(\xi)=\sum_{i=0}^{k} \operatorname{rank}_{\text {full }}\left(G_{i}\right)$. We remark that the full rank corresponds to the rank defined in [15]. Ranks are ordered lexicographically: let $\boldsymbol{r}=\left(r_{d}, \ldots, r_{0}\right)$ and $\boldsymbol{r}^{\prime}=\left(r_{d}^{\prime}, \ldots, r_{0}^{\prime}\right)$, we write $\boldsymbol{r} \leq_{\text {lex }} \boldsymbol{r}^{\prime}$ if $\boldsymbol{r}=\boldsymbol{r}^{\prime}$ or the maximal $i$ with $r_{i} \neq r_{i}^{\prime}$ satisfies $r_{i}<r_{i}^{\prime}$.

Recall that for a VASS $G$ we write $|G|$ for its size as defined in (1). For a linear path scheme $\Lambda$, its length $|\Lambda|$ and norm $\|\Lambda\|$ are defined in Section 3.1. Let $\zeta=\langle p(\boldsymbol{x}) G q(\boldsymbol{y})\rangle$ be a KLM tuple of dimension $d$, its size is defined to be $|\zeta|=|G|+d \cdot(\|\boldsymbol{x}\|+\|\boldsymbol{y}\|+1)$. Let $\xi$ be a linear KLM sequence given as (12), we define its size as

$$
\begin{equation*}
|\xi|=\sum_{i=0}^{k}\left|\xi_{i}\right|+\sum_{i=1}^{k} d \cdot\left|\Lambda_{i}\right| \cdot\left(\left\|\Lambda_{i}\right\|+1\right) \tag{14}
\end{equation*}
$$

Note that the sizes defined in this paper reflect the sizes of unary encoding, thus have an exponential expansion in their binary encoding.

### 5.2 Characteristic Systems for Linear KLM Sequences

We define in this section the characteristic systems of linear KLM sequences, which are systems of linear inequalities that serve as an under-specification of admitted runs. Let $G=(Q, T)$ be a VASS, we first recall the Kirchhoff system $K_{G, p, q}$ of $G$ with respect to states $p, q \in Q$, which is a system of linear equations such that a function $\phi \in \mathbb{N}^{T}$ is a model of $K_{G, p, q}$, written $\phi \models K_{G, p, q}$, if and only if

$$
\begin{equation*}
\mathbf{1}_{q}-\mathbf{1}_{p}=\sum_{t=(r, \boldsymbol{a}, s) \in T} \phi(t) \cdot\left(\mathbf{1}_{s}-\mathbf{1}_{r}\right) \tag{15}
\end{equation*}
$$

where $\mathbf{1}_{p} \in\{0,1\}^{Q}$ is the indicator function defined by $\mathbf{1}_{p}(q)=1$ if $q=p$ and $\mathbf{1}_{p}(q)=0$ otherwise. We also need the homogeneous version of $K_{G, p, q}$, denoted by $K_{G, p, q}^{0}$, where a function $\phi \in \mathbb{N}^{T}$ is a model of it, written $\phi \models K_{G, p, q}^{0}$, if and only if

$$
\begin{equation*}
\mathbf{0}=\sum_{t=(r, a, s) \in T} \phi(t) \cdot\left(\mathbf{1}_{s}-\mathbf{1}_{r}\right) . \tag{16}
\end{equation*}
$$

- Definition 5.5. Let $\xi$ be a linear KLM sequence given by

$$
\begin{equation*}
\xi=\left\langle p_{0}\left(\boldsymbol{x}_{0}\right) G_{0} q_{0}\left(\boldsymbol{y}_{0}\right)\right\rangle \Lambda_{1}\left\langle p_{1}\left(\boldsymbol{x}_{1}\right) G_{1} q_{1}\left(\boldsymbol{y}_{1}\right)\right\rangle \Lambda_{2} \ldots \Lambda_{k}\left\langle p_{k}\left(\boldsymbol{x}_{k}\right) G_{k} q_{k}\left(\boldsymbol{y}_{k}\right)\right\rangle \tag{17}
\end{equation*}
$$

The characteristic system $\mathcal{E}(\xi)$ is a set of linear (in)equalities such that a sequence

$$
\begin{equation*}
\boldsymbol{h}=\left(\boldsymbol{m}_{0}, \phi_{0}, \boldsymbol{n}_{0}\right), \boldsymbol{e}_{1},\left(\boldsymbol{m}_{1}, \phi_{1}, \boldsymbol{n}_{1}\right), \boldsymbol{e}_{2}, \ldots, \boldsymbol{e}_{k},\left(\boldsymbol{m}_{k}, \phi_{k}, \boldsymbol{n}_{k}\right) \tag{18}
\end{equation*}
$$

where each $\left(\boldsymbol{m}_{i}, \phi_{i}, \boldsymbol{n}_{i}\right) \in \mathbb{N}^{d} \times \mathbb{N}^{T_{i}} \times \mathbb{N}^{d}$ and each $\boldsymbol{e}_{i} \in \mathbb{N}^{\left|\Lambda_{i}\right|_{*}}$, is a model of $\mathcal{E}(\xi)$, written $\boldsymbol{h} \models \mathcal{E}(\xi)$, if and only if

1. $\boldsymbol{m}_{i} \sqsubseteq \boldsymbol{x}_{i}, \phi_{i} \models K_{G, p, q}, \boldsymbol{n}_{i} \sqsubseteq \boldsymbol{y}_{i}$ and $\boldsymbol{n}_{i}=\boldsymbol{m}_{i}+\Delta\left(\phi_{i}\right)$ for every $i=0, \ldots, k$;
2. $\left(\boldsymbol{n}_{i-1}, \boldsymbol{e}_{i}, \boldsymbol{m}_{i}\right) \models \mathcal{E}_{\mathrm{LPS}}\left(\Lambda_{i}\right)$ for every $i=1, \ldots, k$.

Similarly, the homogeneous characteristic system $\mathcal{E}^{0}(\xi)$ is a set of linear (in)equalities such that a sequence $\boldsymbol{h}$ of the form (18) is a model of $\mathcal{E}^{0}(\xi)$, written $\boldsymbol{h} \models \mathcal{E}^{0}(\xi)$, if and only if

1. $\boldsymbol{m}_{i}(j)=0$ whenever $\boldsymbol{x}_{i}(j) \neq \omega, \phi_{i} \models K_{G, p, q}^{0}, \boldsymbol{n}_{i}(j)=0$ whenever $\boldsymbol{y}_{i}(j) \neq \omega$, and $\boldsymbol{n}_{i}=\boldsymbol{m}_{i}+\Delta\left(\phi_{i}\right)$ for every $i=0, \ldots, k$;
2. $\left(\boldsymbol{n}_{i-1}, \boldsymbol{e}_{i}, \boldsymbol{m}_{i}\right) \models \mathcal{E}_{\mathrm{LPS}}^{0}\left(\Lambda_{i}\right)$ for every $i=1, \ldots, k$.

The sequence $\xi$ is said to be satisfiable if $\mathcal{E}(\xi)$ has a model, otherwise it's unsatisfiable.
Let $\boldsymbol{h}$ be a model of $\mathcal{E}(\xi)$ (or $\mathcal{E}^{0}(\xi)$ ), we shall write $\boldsymbol{m}_{i}^{\boldsymbol{h}}, \phi_{i}^{\boldsymbol{h}}, \boldsymbol{n}_{i}^{\boldsymbol{h}}, \boldsymbol{e}_{i}^{\boldsymbol{h}}$ for the values of $\boldsymbol{m}_{i}, \phi_{i}, \boldsymbol{n}_{i}, \boldsymbol{e}_{i}$ assigned by $\boldsymbol{h}$, respectively. Recall that unsatisfiable linear KLM sequences have empty action languages.

- Lemma 5.6 (cf. [15, Lem. 3.5]). For any unsatisfiable linear $K L M$ sequence $\xi, L_{\xi}=\emptyset$.


### 5.2.1 Bounds on Bounded Values in $\mathcal{E}(\xi)$

We state here a lemma similar to [15, Lem. 3.7], which upper bounds the bounded values in the characteristic system $\mathcal{E}(\xi)$. Its proof can be found in the appendix, which is a straightforward application of tools in [3] and [18].

- Lemma 5.7. Assume that $\xi=\left\langle p_{0}\left(\boldsymbol{x}_{0}\right) G_{0} q_{0}\left(\boldsymbol{y}_{0}\right)\right\rangle \Lambda_{1} \cdots \Lambda_{k}\left\langle p_{k}\left(\boldsymbol{x}_{k}\right) G_{k} q_{k}\left(\boldsymbol{y}_{k}\right)\right\rangle$ is satisfiable. Then for every $0 \leq i \leq k$ we have:
- For every $1 \leq j \leq d$, the set of values $\boldsymbol{m}_{i}^{\boldsymbol{h}}(j)$ where $\boldsymbol{h} \models \mathcal{E}(\xi)$ is unbounded if, and only $i f$, there exists a model $\boldsymbol{h}_{0}$ of $\mathcal{E}^{0}(\xi)$ such that $\boldsymbol{m}_{i}^{\boldsymbol{h}_{0}}(j)>0$.
- For every $t \in T_{i}$, the set of values $\phi_{i}^{\boldsymbol{h}}(t)$ where $\boldsymbol{h} \models \mathcal{E}(\xi)$ is unbounded if, and only if, there exists a model $\boldsymbol{h}_{0}$ of $\mathcal{E}^{0}(\xi)$ such that $\phi_{i}^{h}(t)>0$.
- For every $1 \leq j \leq d$, the set of values $\boldsymbol{n}_{i}^{h}(j)$ where $\boldsymbol{h} \models \mathcal{E}(\xi)$ is unbounded if, and only if, there exists a model $\boldsymbol{h}_{0}$ of $\mathcal{E}^{0}(\xi)$ such that $\boldsymbol{n}_{i}^{\boldsymbol{h}_{0}}(j)>0$.
Moreover, every bounded value of $\mathcal{E}(\xi)$ is bounded by $(10|\xi|)^{12|\xi|}$.


### 5.3 Cleaning of Linear KLM Sequences

In this section we define three conditions that require a linear KLM sequence to be strongly connected, pure, and saturated. Together with the satisfiability condition, they make up the so-called "clean" condition of linear KLM sequences. Note that the purity condition is new compared to [15], which requires every geometrically 2-dimensional VASSes occur in a linear KLM sequence to be replaced by linear path schemes.

## Strongly Connected Sequences

A linear KLM sequence $\xi=\left\langle p_{0}\left(\boldsymbol{x}_{0}\right) G_{0} q_{0}\left(\boldsymbol{y}_{0}\right)\right\rangle \Lambda_{1} \cdots \Lambda_{k}\left\langle p_{k}\left(\boldsymbol{x}_{k}\right) G_{k} q_{k}\left(\boldsymbol{y}_{k}\right)\right\rangle$ is strongly connected if all the VASSes $G_{0}, \ldots, G_{k}$ are strongly connected (as they are understood as directed graphs). One can easily obtain strongly connected sequences by expanding the strongly connected components of each VASS:

Lemma 5.8 ([15, Lem. 4.2]). For any linear KLM sequence $\xi$ which is not strongly connected, one can compute in time $\exp (|\xi|)$ a finite set $\Xi$ of strongly connected linear KLM sequences such that $L_{\xi}=\bigcup_{\xi^{\prime} \in \Xi} L_{\xi^{\prime}}$ and that $\operatorname{rank}\left(\xi^{\prime}\right) \leq_{l e x} \operatorname{rank}(\xi)$ and $\left|\xi^{\prime}\right| \leq(2 d+1)|\xi|$ for every $\xi^{\prime} \in \Xi$.

## Pure Sequences

A KLM tuple $\langle p(\boldsymbol{x}) G q(\boldsymbol{y})\rangle$ is called trivial if $p(\boldsymbol{x})=q(\boldsymbol{y})$ and $G$ contains no transition and only a single state $p$. In this case we simply write $\langle p(\boldsymbol{x})\rangle$ for this tuple. Note that the action language of a trivial tuple contains exactly the empty word.

A linear KLM sequence $\xi=\left\langle p_{0}\left(\boldsymbol{x}_{0}\right) G_{0} q_{0}\left(\boldsymbol{y}_{0}\right)\right\rangle \Lambda_{1} \cdots \Lambda_{k}\left\langle p_{k}\left(\boldsymbol{x}_{k}\right) G_{k} q_{k}\left(\boldsymbol{y}_{k}\right)\right\rangle$ is said to be pure if $\xi$ is strongly connected and for every $i=0, \ldots, k, \operatorname{rank}\left(G_{i}\right)=\mathbf{0}$ implies that the tuple $\left\langle p_{i}\left(\boldsymbol{x}_{i}\right) G_{i} q_{i}\left(\boldsymbol{y}_{i}\right)\right\rangle$ is trivial. By Corollary 5.4, a rank- $\mathbf{0}$ strongly connected VASS is geometrically 2 -dimensional, and thus can be replaced by linear path schemes in case it is not trivial.

- Lemma 5.9. Let $\xi$ be a strongly connected linear KLM sequence. Whether $\xi$ is pure is in PSPACE. If $\xi$ is not pure, one can compute in space poly $(|\xi|)$ a finite set $\Xi$ of pure linear KLM sequences such that $\bigcup_{\xi^{\prime} \in \Xi} L_{\xi^{\prime}} \subseteq L_{\xi}$ and $\bigcup_{\xi^{\prime} \in \Xi} L_{\xi^{\prime}} \neq \emptyset$ whenever $L_{\xi} \neq \emptyset$, and such that $\operatorname{rank}\left(\xi^{\prime}\right)=\operatorname{rank}(\xi)$ and $\left|\xi^{\prime}\right| \leq|\xi|^{O(1)}$ for all $\xi^{\prime} \in \Xi$.


## Saturated Sequences

Let $\xi=\left\langle p_{0}\left(\boldsymbol{x}_{0}\right) G_{0} q_{0}\left(\boldsymbol{y}_{0}\right)\right\rangle \Lambda_{1} \cdots \Lambda_{k}\left\langle p_{k}\left(\boldsymbol{x}_{k}\right) G_{k} q_{k}\left(\boldsymbol{y}_{k}\right)\right\rangle$ be a linear KLM sequence. We say $\xi$ is saturated if for every $0 \leq i \leq k$ and every $j \in[d]$, we have

- $\boldsymbol{x}_{i}(j)=\omega$ implies the set of values $\boldsymbol{m}_{i}^{\boldsymbol{h}}(j)$ where $\boldsymbol{h} \models \mathcal{E}(\xi)$ is unbounded; and
- $\boldsymbol{y}_{i}(j)=\omega$ implies the set of values $\boldsymbol{n}_{i}^{\boldsymbol{h}}(j)$ where $\boldsymbol{h} \models \mathcal{E}(\xi)$ is unbounded.
- Lemma 5.10 ([15, Lem. 4.4]). From any pure linear KLM sequence $\xi$, one can compute in time $\exp \left(|\xi|^{O(|\xi|)}\right)$ a finite set $\Xi$ of saturated pure linear KLM sequences such that $L_{\xi}=$ $\bigcup_{\xi^{\prime} \in \Xi} L_{\xi^{\prime}}$, and such that $\operatorname{rank}\left(\xi^{\prime}\right)=\operatorname{rank}(\xi)$ and $\left|\xi^{\prime}\right| \leq|\xi|^{O(|\xi|)}$ for every $\xi^{\prime} \in \Xi$.

Proof. By Lemma 5.7, if a variable $\boldsymbol{m}_{i}(j)$ or $\boldsymbol{n}_{i}(j)$ is bounded in $\mathcal{E}(\xi)$, we can replace the corresponding $\omega$ component in $\xi$ by all possible values bounded by $(10|\xi|)^{12|\xi|} \leq|\xi|^{O(|\xi|)}$.

## The Cleaning Lemma

A linear KLM sequence $\xi$ is called clean if it is satisfiable, strongly connected, pure and saturated. The lemmas 5.8 through 5.10 show how to make a linear KLM sequence clean.

- Lemma 5.11. From any linear KLM sequence $\xi$, one can compute in time $\exp (g(|\xi|))$ a finite set clean $(\xi)$ of clean linear KLM sequences such that $\bigcup_{\xi^{\prime} \in \operatorname{clean}(\xi)} L_{\xi^{\prime}} \subseteq L_{\xi}$ and $\bigcup_{\xi^{\prime} \in \operatorname{clean}(\xi)} L_{\xi^{\prime}} \neq \emptyset$ whenever $L_{\xi} \neq \emptyset$. Moreover, for every $\xi^{\prime} \in \operatorname{clean}(\xi)$ we have $\operatorname{rank}\left(\xi^{\prime}\right) \leq{ }_{\text {lex }}$ $\operatorname{rank}(\xi)$ and $\left|\xi^{\prime}\right| \leq g(|\xi|)$ where $g(x)=x^{x^{O(1)}}$.


### 5.4 Decomposition of Linear KLM Sequences

In this section we recall three conditions that require a linear KLM sequence to be unbounded, rigid, and pumpable. If any one of them is violated, a decomposition into a set of linear KLM sequences with strictly smaller ranks can be performed. Essentially there is nothing new in this section compared to [15]. The decomposition operations in [15] can be directly applied here, since they operate on a single KLM tuple and produce KLM sequences that are just special cases of linear KLM sequences. The proofs in [15] can also be adapted easily, and we will omit the details here. Especially, the next lemma shows that the arguments of strict decrease in ranks are still valid even though we discard the lower three components of ranks.

- Lemma 5.12. Let $\xi^{\prime}$ be a pure linear KLM sequence. For any linear KLM sequence $\xi^{\prime}$ with $\operatorname{rank}_{\text {full }}\left(\xi^{\prime}\right)<_{\text {lex }} \operatorname{rank}_{\text {full }}(\xi)$, we have $\operatorname{rank}\left(\xi^{\prime}\right)<_{\text {lex }} \operatorname{rank}(\xi)$.


## Unbounded Sequences

Let $\xi=\left\langle p_{0}\left(\boldsymbol{x}_{0}\right) G_{0} q_{0}\left(\boldsymbol{y}_{0}\right)\right\rangle \Lambda_{1} \cdots \Lambda_{k}\left\langle p_{k}\left(\boldsymbol{x}_{k}\right) G_{k} q_{k}\left(\boldsymbol{y}_{k}\right)\right\rangle$ be a linear KLM sequence. We say $\xi$ is unbounded if for all $i=0, \ldots, k$ and every transition $t \in T_{i}$ where $T_{i}$ is the set of transitions of $G_{i}$, the set of values $\phi_{i}^{h}(t)$ where $\boldsymbol{h} \models \mathcal{E}(\xi)$ is unbounded. Bounded transitions can be expanded exhaustively according to the bounds given by Lemma 5.7.

- Lemma 5.13 ([15, Lem. 4.6]). Whether a linear KLM sequence $\xi$ is unbounded is decidable in NP. Moreover, if $\xi$ is pure and bounded, one can compute in time $\exp \left(|\xi|^{O(|\xi|)}\right)$ a finite set $\Xi$ of linear KLM sequences such that $L_{\xi}=\bigcup_{\xi^{\prime} \in \Xi} L_{\xi^{\prime}}$ and such that $\operatorname{rank}\left(\xi^{\prime}\right)<_{\text {lex }} \operatorname{rank}(\xi)$ and $\left|\xi^{\prime}\right|<|\xi|^{O(|\xi|)}$ for every $\xi^{\prime} \in \Xi$.


## Rigid Sequences

A coordinate $j \in[d]$ is said to be fixed by a VASS $G=(Q, T)$ if there exists a function $f_{j}: Q \rightarrow \mathbb{N}$ such that $f_{j}(q)=f_{j}(p)+\boldsymbol{a}(j)$ for every transition $(p, \boldsymbol{a}, q) \in T$. We also say that $f_{j}$ fixes $G$ at coordinate $j$ in this case.

A KLM tuple $\langle p(\boldsymbol{x}) G q(\boldsymbol{y})\rangle$ is said to be rigid if for every coordinate $j$ fixed by $G=(Q, T)$, there exists a function $g_{j}: Q \rightarrow \mathbb{N}$ that fixes $G$ at coordinate $j$ and such that $g_{j}(p) \sqsubseteq \boldsymbol{x}(j)$ and $g_{j}(q) \sqsubseteq \boldsymbol{y}(j)$.

A linear KLM sequence $\xi=\left\langle p_{0}\left(\boldsymbol{x}_{0}\right) G_{0} q_{0}\left(\boldsymbol{y}_{0}\right)\right\rangle \Lambda_{1} \cdots \Lambda_{k}\left\langle p_{k}\left(\boldsymbol{x}_{k}\right) G_{k} q_{k}\left(\boldsymbol{y}_{k}\right)\right\rangle$ is said to be rigid if every tuple $\left\langle p_{i}\left(\boldsymbol{x}_{i}\right) G_{i} q_{i}\left(\boldsymbol{y}_{i}\right)\right\rangle$ in $\xi$ is rigid.

- Lemma 5.14 ([15, Lem. 4.9]). From any pure linear KLM sequence $\xi$ one can decide in polynomial time whether $\xi$ is not rigid. Moreover, in that case one can compute in polynomial time a linear $K L M$ sequence $\xi^{\prime}$ such that $L_{\xi}=L_{\xi^{\prime}}, \operatorname{rank}\left(\xi^{\prime}\right)<_{\text {lex }} \operatorname{rank}(\xi)$, and $\left|\xi^{\prime}\right| \leq|\xi|$.


## Pumpable Sequences

Given a KLM tuple $\langle p(\boldsymbol{x}) G q(\boldsymbol{y})\rangle$, recall the forward and backward acceleration vectors $\operatorname{FACC}_{G, p}(\boldsymbol{x}), \operatorname{BACC}_{G, q}(\boldsymbol{y}) \in \mathbb{N}_{\omega}^{d}$ defined by

$$
\begin{align*}
& \operatorname{FACC}_{G, p}(\boldsymbol{x})(j)= \begin{cases}\omega & \text { if } p(\boldsymbol{x}) \xrightarrow{*} p\left(\boldsymbol{x}^{\prime}\right) \text { for some } \boldsymbol{x}^{\prime} \text { with } \boldsymbol{x}^{\prime} \geq \boldsymbol{x}, \boldsymbol{x}^{\prime}(j)>\boldsymbol{x}(j) \\
\boldsymbol{x}(j) & \text { otherwise }\end{cases}  \tag{19}\\
& \operatorname{BACC}_{G, q}(\boldsymbol{y})(j)= \begin{cases}\omega & \text { if } q\left(\boldsymbol{y}^{\prime}\right) \xrightarrow{*} q(\boldsymbol{y}) \text { for some } \boldsymbol{y}^{\prime} \text { with } \boldsymbol{y}^{\prime} \geq \boldsymbol{y}, \boldsymbol{y}^{\prime}(j)>\boldsymbol{y}(j) \\
\boldsymbol{y}(j) & \text { otherwise }\end{cases} \tag{20}
\end{align*}
$$

A tuple $\langle p(\boldsymbol{x}) G q(\boldsymbol{y})\rangle$ is said to be pumpable if $\operatorname{FACC}_{G, p}(\boldsymbol{x})(j)=\operatorname{BACC}_{G, q}(\boldsymbol{y})(j)=\omega$ for every coordinate $j$ not fixed by $G$.

A linear KLM sequence given by $\xi=\left\langle p_{0}\left(\boldsymbol{x}_{0}\right) G_{0} q_{0}\left(\boldsymbol{y}_{0}\right)\right\rangle \Lambda_{1} \cdots \Lambda_{k}\left\langle p_{k}\left(\boldsymbol{x}_{k}\right) G_{k} q_{k}\left(\boldsymbol{y}_{k}\right)\right\rangle$ is said to be pumpable if every tuple $\left\langle p_{i}\left(\boldsymbol{x}_{i}\right) G_{i} q_{i}\left(\boldsymbol{y}_{i}\right)\right\rangle$ in $\xi$ is pumpable.

- Lemma 5.15 ([15, Lem. 4.15]). Whether a linear KLM sequence $\xi$ is pumpable is decidable in EXPSPACE. Moreover, if $\xi$ is pure and unpumpable, one can compute in time $\exp \left(|\xi|{ }^{O(1)}\right)$ a finite set $\Xi$ of linear $K L M$ sequences such that $L_{\xi}=\bigcup_{\xi^{\prime} \in \Xi} L_{\xi^{\prime}}$ and such that $\operatorname{rank}\left(\xi^{\prime}\right)<$ lex $\operatorname{rank}(\xi)$ and $\left|\xi^{\prime}\right|<|\xi|^{O(1)}$ for every $\xi^{\prime} \in \Xi$.

Note that the $O(1)$ term here hides a constant depending on $d$, which essentially arises from a result on the coverability problem by Rackoff [19]. The $O(1)$ term also captures the difference between the sizes of linear KLM sequences defined here and that in [15].

## The Decomposition Lemma

A linear KLM sequence is normal if it is clean, unbounded, rigid, and pumpable. The lemmas 5.13 through 5.15 show that when a clean linear KLM sequence is not normal, we are able to decompose it into a finite set of linear KLM sequences of strictly smaller ranks.

- Lemma 5.16. From any clean linear KLM sequences $\xi$ that is not normal, one can compute in time $\exp (h(\xi))$ a finite set $\operatorname{dec}(\xi)$ of clean linear $K L M$ sequence such that $\bigcup_{\xi^{\prime} \in \operatorname{dec}(\xi)} L_{\xi^{\prime}} \subseteq$ $L_{\xi}$ and $\bigcup_{\xi^{\prime} \in \operatorname{dec}(\xi)} L_{\xi^{\prime}} \neq \emptyset$ whenever $L_{\xi} \neq \emptyset$. Moreover, for every $\xi^{\prime} \in \operatorname{dec}(\xi)$ we have $\operatorname{rank}\left(\xi^{\prime}\right)<_{\text {lex }} \operatorname{rank}(\xi)$ and $\left|\xi^{\prime}\right| \leq h(\xi)$ where $h(x)=x^{x^{x^{O(1)}}}$.


### 5.5 Normal Sequences

The following lemma shows that a normal linear KLM sequence is guaranteed to have non-empty action language, thus one can terminate the decomposition process once a normal sequence is produced.

- Lemma 5.17. Let $\xi$ be a normal linear KLM sequence, then there is a word $\sigma \in L_{\xi}$ whose length is bounded by $|\sigma| \leq \ell(|\xi|)$ where $\ell(x) \leq x^{O(x)}$.


### 5.6 Putting All Together: The Modified KLMST Algorithm

Here we describe the modified KLMST decomposition algorithm for VASS reachability problem. Suppose we are given a $d$-VASS $G=(Q, T)$ and two configurations $p(\boldsymbol{m}), q(\boldsymbol{n}) \in$ $Q \times \mathbb{N}^{d}$. To decide whether $p(\boldsymbol{m}) \xrightarrow{*} q(\boldsymbol{n})$ holds in $G$, it is enough to decide whether $L_{\xi}$ is non-empty where $\xi=\langle p(\boldsymbol{m}) G q(\boldsymbol{n})\rangle$. To start with, we use Lemma 5.11 to clean the sequence $\xi$, and then choose $\xi^{0} \in \operatorname{clean}(\xi)$ non-deterministically. If $\xi^{0}$ is normal then we are done by Lemma 5.17. Otherwise, we decompose $\xi^{0}$ using Lemma 5.16 and choose $\xi^{1} \in \operatorname{dec}\left(\xi^{0}\right)$ non-deterministically. The procedure continues to produce a series of linear KLM sequences $\xi^{0}, \xi^{1}, \xi^{2}, \ldots$ where $\xi^{i+1} \in \operatorname{dec}\left(\xi^{i}\right)$, until either we finally obtain a normal sequence $\xi^{n}$, or at some point we have to abort because the decomposition of a linear KLM sequence is the empty set. The procedure terminates because $\operatorname{rank}\left(\xi^{0}\right)>_{\text {lex }} \operatorname{rank}\left(\xi^{1}\right)>_{\text {lex }} \operatorname{rank}\left(\xi^{2}\right)>_{\text {lex }} \ldots$ form a decreasing chain of the well-order $\left(\mathbb{N}^{d-2},<_{l e x}\right)$, which must be finite. If $L_{\xi}=\emptyset$ then we cannot get a normal sequence since the action languages $L_{\xi} \supseteq L_{\xi^{0}} \supseteq L_{\xi^{1}} \supseteq \cdots$ are all empty. On the other hand, if $L_{\xi} \neq \emptyset$ then there are non-deterministic choices that always choose the linear KLM sequences with non-empty action languages, which finally lead to a normal sequence. This shows the correctness of the algorithm.

## 6 Complexity Upper Bound

The termination of the modified KLMST decomposition algorithm is guaranteed by a ranking function that decreases along a well-ordering. In order to analyze the length of this decreasing chain, we recall the so-called "length function theorems" by Schmitz [20] in Section 6.1. After that, we can locate the complexity upper bound of the algorithm in the fast-growing complexity hierarchy [21] which we recall in Section 6.2. Readers familiar with [15] may realize that the complexity upper bound for $d$-VASS can be improved to $\mathrm{F}_{d+1}$, i.e. the $(d+1)$-th level in the fast-growing hierarchy, with our ranking function. In fact, by a careful analysis on a property of fast-growing functions, we further improve this bound to $F_{d}$.

In this section we assume some familiarity with ordinal numbers (see, e.g. [9]). We write $\omega$ here for the first infinite ordinal, not to be confused with the infinite element in previous sections.

### 6.1 Length of Sequences of Decreasing Ranks

Let $\xi$ be a linear KLM sequence of dimension $d$ with $\operatorname{rank}(\xi)=\left(r_{d}, \ldots, r_{3}\right)$, we define the ordinal rank $\alpha_{\xi}$ of $\xi$ as the ordinal number given by

$$
\begin{equation*}
\alpha_{\xi}:=\omega^{d-3} \cdot r_{d}+\omega^{d-4} \cdot r_{d-1}+\cdots+\omega^{0} \cdot r_{3} . \tag{21}
\end{equation*}
$$

Note that $\operatorname{rank}(\xi)<_{\operatorname{lex}} \operatorname{rank}\left(\xi^{\prime}\right)$ if and only if $\alpha_{\xi}<\alpha_{\xi^{\prime}}$. With this reformulation, we now focus on the decreasing chain of ordinal ranks.

Let $\alpha<\omega^{\omega}$ be an ordinal given in Cantor Normal Form as $\alpha=\omega^{n} \cdot c_{n}+\cdots+\omega^{0} \cdot c_{0}$ where $n, c_{0}, \ldots, c_{n} \in \mathbb{N}$, we define the size of $\alpha$ as $N \alpha:=\max \left\{n, \max _{0 \leq i \leq n} c_{i}\right\}$. For the linear KLM sequence $\xi$ with $\operatorname{rank}(\xi)=\left(r_{d}, \ldots, r_{3}\right)$, we have $N \alpha_{\xi}=\max \left\{d-3, \max _{3 \leq i \leq d} r_{i}\right\} \leq|\xi|$.

Given a number $n_{0} \in \mathbb{N}$ and a function $h: \mathbb{N} \rightarrow \mathbb{N}$ that is monotone inflationary (that is, $x \leq h(x)$ and $h(x) \leq h(y)$ whenever $x \leq y$ ), we say a sequence of ordinals $\alpha_{0}, \alpha_{1}, \ldots$ is $\left(n_{0}, h\right)$-controlled if $N \alpha_{i} \leq h^{i}\left(n_{0}\right)$ for all $i \in \mathbb{N}$, where $h^{i}\left(n_{0}\right)$ is the $i$ th iteration of $h$ on $n_{0}$.

Let $\xi^{0}, \xi^{1}, \ldots$ be the linear KLM sequences produced in the modified KLMST algorithm, by Lemma 5.16 we know that the sequence of ordinal ranks

$$
\begin{equation*}
\alpha_{\xi^{0}}>\alpha_{\xi^{1}}>\cdots \tag{22}
\end{equation*}
$$

is $\left(\left|\xi^{0}\right|, h\right)$-controlled where $h$ is defined in Lemma 5.16. Recall that $\xi^{0} \in \operatorname{clean}(\xi)$ where $\xi=\langle p(\boldsymbol{m}) G q(\boldsymbol{n})\rangle$ corresponds to the input reachability instance. Then $\left|\xi^{0}\right| \leq g(|\xi|)$ where $g$ is defined in Lemma 5.11, and (22) is indeed $(g(n), h)$-controlled where $n:=|\langle p(\boldsymbol{m}) G q(\boldsymbol{n})\rangle|$.

## Length function theorem

The length of the controlled sequence of ordinals (22) can be bounded in terms of the hierarchies of fast-growing functions of Hardy and Cichon [4]. First recall that given a limit ordinal $\lambda \leq \omega^{\omega}$, the standard fundamental sequence of $\lambda$ is a sequence $(\lambda(x))_{x<\omega}$ defined inductively by

$$
\begin{equation*}
\omega^{\omega}(x):=\omega^{x+1}, \quad\left(\beta+\omega^{k+1}\right)(x):=\beta+\omega^{k} \cdot(x+1) \tag{23}
\end{equation*}
$$

where $\beta+\omega^{k+1}$ is in Cantor Normal Form. Now given a function $h: \mathbb{N} \rightarrow \mathbb{N}$ that is monotone inflationary, we define the Hardy hierarchy $\left(h^{\alpha}\right)_{\alpha \leq \omega^{\omega}}$ and the Cichon hierarchy $\left(h_{\alpha}\right)_{\alpha \leq \omega^{\omega}}$ as two families of functions $h^{\alpha}, h_{\alpha}: \mathbb{N} \rightarrow \mathbb{N}$ indexed by ordinals $\alpha \leq \omega^{\omega}$ given inductively by

$$
\begin{array}{lll}
h^{0}(x):=x, & h^{\alpha+1}(x):=\quad h^{\alpha}(h(x)), & h^{\lambda}(x):=h^{\lambda(x)}(x), \\
h_{0}(x):=0, & h_{\alpha+1}(x):=1+h_{\alpha}(h(x)), & h_{\lambda}(x):=h_{\lambda(x)}(x) . \tag{25}
\end{array}
$$

Observe that Cichon hierarchy counts the number of iterations of $h$ in Hardy hierarchy, that is, $h^{h_{\alpha}(x)}(x)=h^{\alpha}(x)$. Also note that as $h$ is monotone inflationary, by induction on $\alpha$ we have $h_{\alpha}(x) \leq h^{\alpha}(x)$. Now we give the length function theorem as follows.

- Theorem 6.1 (Length function theorem, [20, Thm. 3.3]). Let $n_{0} \geq d-2$, then the maximal length of $\left(n_{0}, h\right)$-controlled decreasing sequences of ordinals in $\omega^{d-2}$ is $h_{\omega^{d-2}}\left(n_{0}\right)$.


## Small witness property

By Theorem 6.1 we can bound the length of (22), which then yields a bound on the minimal length of runs witnessing reachability.

- Lemma 6.2 (Small witnesses). Let $G=(Q, T)$ be a d-VASS where $d \geq 3$, let $p(\boldsymbol{m}), q(\boldsymbol{n}) \in$ $Q \times \mathbb{N}^{d}$ be two configurations, and let $n:=|\langle p(\boldsymbol{m}) G q(\boldsymbol{n})\rangle|$. If $p(\boldsymbol{m}) \xrightarrow{*} q(\boldsymbol{n})$ holds in $G$, then there is a path $\sigma$ such that $p(\boldsymbol{m}) \xrightarrow{\sigma} q(\boldsymbol{n})$ and $|\sigma| \leq \ell\left(h^{\omega^{d-2}}(g(n))\right)$, where $g, h, \ell$ are defined in lemmas 5.11, 5.16, and 5.17.
Proof. Suppose $p(\boldsymbol{m}) \xrightarrow{*} q(\boldsymbol{n})$, then there is a sequence of linear KLM sequence $\xi^{0}, \xi^{1}, \ldots, \xi^{L}$ produced in the modified KLMST algorithm, such that $\xi^{L}$ is normal. We have discussed that the sequence of their ordinal ranks $\alpha_{\xi^{0}}>\alpha_{\xi^{1}}>\cdots>\alpha_{\xi^{L}}$ is $(g(n), h)$-controlled, so by Theorem 6.1 we have $L \leq h_{\omega^{d-2}}(g(n))$. From Lemma 5.16 and the fact that $h^{h_{\alpha}(x)}=h^{\alpha}(x)$, the size of $\xi^{L}$ is bounded by

$$
\begin{equation*}
\left|\xi^{L}\right| \leq h^{L}\left(\left|\xi^{0}\right|\right) \leq h^{h_{\omega^{d-2}}(g(n))}(g(n))=h^{\omega^{d-2}}(g(n)) \tag{26}
\end{equation*}
$$

Now Lemma 5.17 bounds the length of the minimal witnesses by $\ell\left(h^{\omega^{d-2}}(g(n))\right)$.

### 6.2 Fast-Growing Complexity Hierarchy

We recall the fast-growing hierarchy formally introduced by Schmitz [21] that captures the complexity class high above elementary. Define $H(x):=x+1$, we shall use the Hardy hierarchy $\left(H^{\alpha}\right)_{\alpha<\omega^{\omega}}$, where for example $H^{\omega^{2}}(x)=2^{x+1}(x+1)$ and $H^{\omega^{3}}(x)$ grows faster than the tower function. First we define the family $\mathcal{F}_{\alpha}:=\bigcup_{\beta<\omega^{\alpha}} \operatorname{FDTIME}\left(H^{\beta}(n)\right)$ which contains functions computable in deterministic time $O\left(H^{\beta}(n)\right)$. Observe that, for example, $\mathcal{F}_{3}$ contains exactly the Kalmar elementary functions. Now we define

$$
\begin{equation*}
\mathrm{F}_{\alpha}:=\bigcup_{p \in \mathcal{F}_{\alpha}} \operatorname{DTIME}\left(H^{\omega^{\alpha}}(p(n))\right) \tag{27}
\end{equation*}
$$

which is the class of decision problems solvable in deterministic time $O\left(H^{\omega^{\alpha}}(p(n))\right)$. Note that non-deterministic time Turing machines can be made deterministic with an exponential overhead in $\mathcal{F}_{3}$, thus for $\alpha \geq 3$, we have equivalently that $\mathrm{F}_{\alpha}=\bigcup_{p \in \mathcal{F}_{\alpha}} \operatorname{NDTIME}\left(H^{\omega^{\alpha}}(p(n))\right)$. Observe that $\mathrm{F}_{\alpha}$ is closed under reductions in $\mathcal{F}_{\alpha}$.

### 6.2.1 Relativized Fast-Growing Functions

In order to express the complexity of the modified KLMST algorithm in terms of the hierarchy $\left(\mathrm{F}_{\alpha}\right)_{\alpha<\omega}$, one needs to locate the function $h^{\omega^{d-2}}$ in the Hardy hierarchy $\left(H^{\alpha}\right)_{\alpha<\omega^{\omega}}$ where $h \in \mathcal{F}_{3}$ is the elementary function from Lemma 5.16. Previously we can upper bound $h^{\omega^{d-2}}$ by $H^{\omega^{d+1}}$ with the help of [21, Lem. 4.2]. Here we show a slightly better result, from which we can bound $h^{\omega^{d-2}}(x)$ by $H^{\omega^{d}}(O(x))$.

- Lemma 6.3 (cf. [21, Lem. A.5]). Let $h: \mathbb{N} \rightarrow \mathbb{N}$ be a monotone inflationary function, let $a, b, c \geq 1$ and $x_{0} \geq 0$ be natural numbers. If for all $x \geq x_{0}, h(x) \leq H^{\omega^{b} \cdot c}(x)$, then $h^{\omega^{a}}(x) \leq H^{\omega^{b+a}}((c+1) x)$ for all $x \geq \max \left\{2 c, x_{0}\right\}$.


### 6.3 Upper Bounds for VASS Reachability

Now we analyze the time complexity of the modified KLMST algorithm. Given as input the $d$-VASS $G=(Q, T)$ and two configurations $p(\boldsymbol{m}), q(\boldsymbol{n})$, let $\xi:=\langle p(\boldsymbol{m}) G q(\boldsymbol{n})\rangle$ and $n:=|\xi|$. The initial sequence $\xi^{0} \in \operatorname{clean}(\xi)$ can be computed in (non-deterministic) time elementary in $n$ by Lemma 5.11 . Then the algorithm produces $\xi^{0}, \xi^{1}, \ldots, \xi^{L}$ with $L \leq h_{\omega^{d-2}}(g(n))$, where $g, h$ are defined in lemmas 5.11, 5.16. Note that in each step, the sequence $\xi^{i+1} \in \operatorname{dec}\left(\xi^{i}\right)$ can be computed in time elementary in $\left|\xi^{i}\right|$ by Lemma 5.16, and the sizes $\left|\xi^{i}\right|$ are all bounded by $h^{\omega^{d-2}}(g(n))$ as we have discussed above in the proof of Lemma 6.2. To sum up, the entire algorithm finishes in non-deterministic time elementary in $h^{\omega^{d-2}}(g(n))$.

- Lemma 6.4. On input a d-VASS $G=(Q, T)$ where $d \geq 3$ and $p(\boldsymbol{m}), q(\boldsymbol{n}) \in Q \times \mathbb{N}^{d}$, the modified KLMST algorithm finishes in non-deterministic time $e\left(h^{\omega^{d-2}}(g(n))\right)$ where $n=|\langle p(\boldsymbol{m}) G q(\boldsymbol{n})\rangle|, g$, $h$ are defined in lemmas 5.11, 5.16, and $e \in \mathcal{F}_{3}$ is some fixed function.

Since $h$ is an elementary function, there is a number $c \in \mathbb{N}$ such that $h$ is eventually dominated by $H^{\omega^{2} \cdot c}$. By Lemma 6.3 we can upper bound $h^{\omega^{d-2}}(x)$ by $H^{\omega^{d}}((c+1) x)$. Observe that the inner part $g(n)$ is elementary in the binary encoding size of the input $G, p(\boldsymbol{m}), q(\boldsymbol{n})$, thus can be captured by a function $p \in \mathcal{F}_{3}$. Finally, [21, Lem. 4.6] allows us to move the outermost function $e$ to the innermost position. Hence we have the following upper bound.

- Theorem 6.5. Reachability in d-dimensional VASS is in $\mathrm{F}_{d}$ for all $d \geq 3$.

Also, by Lemma 6.2 there is a simple combinatorial algorithm for $d$-VASS reachability. We fist compute the bound $B:=\ell\left(h^{\omega^{d-2}}(g(n))\right)$, which can be done in time elementary in $B$ by [21, Thm. 5.1]. Then we can decide reachability by just enumerate all possible paths in $G$ with length bounded by $B$.

## 7 Conclusion

We have shown that the reachability problem in $d$-dimensional vector addition system with states is in $\mathrm{F}_{d}$, improving the previous $\mathrm{F}_{d+4}$ upper bound by Leroux and Schmitz [15]. By capturing reachability in geometrically 2 -dimensional VASSes with linear path schemes, we are able to reduce significantly the number of decomposition steps in the KLMST decomposition algorithm. Combined with a careful analysis on fast-growing functions, we finally obtained the $F_{d}$ upper bound. It should be noticed though, that our algorithm avoids computing the "full decomposition" [14] of KLM sequences, thus cannot improve the complexity of problems that essentially rely on the full decomposition, e.g., the VASS downward language inclusion problem $[8,23,15]$.

It has been shown that the reachability problem in $(2 d+3)$-VASS is $\mathrm{F}_{d}$-hard [5]. In the case of 3-VASS, it is known that the reachability problem is PSPACE-hard. The gap between the lower bound and the upper bound $\mathrm{F}_{3}=$ TOWER [22] is huge. It is very unlikely that the problem is PSPACE-complete. Effort to uplift the lower bound is called for.

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