On Termination of Integer Linear Loops

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Abstract

A fundamental problem in program verification concerns the termination of simple linear loops of the form:

$$\boldsymbol{x} \leftarrow \boldsymbol{u} \; ; \; while \; B \boldsymbol{x} \geq \boldsymbol{c} \; do \; \boldsymbol{x} \leftarrow A \boldsymbol{x} + \boldsymbol{a} \; ,$$

where \boldsymbol{x} is a vector of variables, \boldsymbol{u} , \boldsymbol{a} , and \boldsymbol{c} are integer vectors, and A and B are integer matrices. Assuming the matrix A is diagonalisable, we give a decision procedure for the problem of whether, for all initial integer vectors \boldsymbol{u} , such a loop terminates. The correctness of our algorithm relies on sophisticated tools from algebraic and analytic number theory, Diophantine geometry, and real algebraic geometry.

To the best of our knowledge, this is the first substantial advance on a 10-year-old open problem of Tiwari [38] and Braverman [8].

1 Introduction

Termination is a fundamental decision problem in program verification. In particular, termination of programs with linear assignments and linear conditionals has been extensively studied over the last decade. This has led to the development of powerful techniques to prove termination via synthesis of linear ranking functions [6, 7, 10, 13, 30], many of which have been implemented in software-verification tools, such as Microsoft's TERMINATOR [14].

A very simple form of linear programs are *simple linear loops*, that is, programs of the form

 $\mathsf{P1}: \boldsymbol{x} \leftarrow \boldsymbol{u} ; while \ B\boldsymbol{x} \geq \boldsymbol{c} \ do \ \boldsymbol{x} \leftarrow A\boldsymbol{x} + \boldsymbol{a},$

where \boldsymbol{x} is vector of variables, \boldsymbol{u} , \boldsymbol{a} , and \boldsymbol{c} are integer vectors, and A and B are integer matrices of the appropriate dimensions. Here the loop guard is a

conjunction of linear inequalities and the loop body consists of a simultaneous affine assignment to \boldsymbol{x} . If the vectors \boldsymbol{a} and \boldsymbol{c} are both zero then we say that the loop is *homogeneous*.

Suppose that the vector \boldsymbol{x} has dimension d. We say that P1 *terminates* on a set $S \subseteq \mathbb{R}^d$ if it terminates for all initial vectors $\boldsymbol{u} \in S$. Tiwari [38] gave a procedure to decide whether a given simple linear loop terminates on \mathbb{R}^d . Later Braverman [8] showed decidability of termination on \mathbb{Q}^d . However the most natural problem from the point of view of program verification is termination on \mathbb{Z}^d .

While termination on \mathbb{Z}^d reduces to termination on \mathbb{Q}^d in the homogeneous case (by a straightforward scaling argument), termination on \mathbb{Z}^d in the general case is stated as an open problem in [5, 8, 38]. The main result of this paper is a procedure to decide termination on \mathbb{Z}^d for simple linear loops when the assignment matrix A is diagonalisable. This represents the first substantial progress on this open problem in over 10 years.

Termination of more complex linear programs can often be reduced to termination of simple linear loops (see, e.g., [14] or [38, Section 6]). On the other hand, termination becomes undecidable for mild generalisations of simple linear loops, for example, allowing the update function in the loop body to be piecewise linear [5].

To prove our main result we focus on *eventual non*termination, where P1 is said to be eventually nonterminating on $u \in \mathbb{Z}^d$ if, starting from initial value u, after executing the loop body $x \leftarrow Ax + a$ a finite number of times while disregarding the loop guard we eventually reach a value on which P1 fails to terminate. Clearly P1 fails to terminate on \mathbb{Z}^d if and only if it is eventually non-terminating on some $u \in \mathbb{Z}^d$.

Given a simple linear loop we show how to compute a convex semi-algebraic set $W \subseteq \mathbb{R}^d$ such that the integer points $u \in W$ are precisely the eventually nonterminating integer initial values. Since it is decidable

^{*}Supported by EPSRC.

 $^{^{\}dagger}\text{Supported}$ by the ERC Advanced Grant 321171 (ALGAME) and by EPSRC.

whether a convex semi-algebraic set contains an integer point [21],¹ we can decide whether an integer linear loop is terminating on \mathbb{Z}^d .

Termination over the set of all integer points is easily seen to be **coNP**-hard. Indeed, if the update function in the loop body is the identity then the loop is non-terminating if and only if there is an integer point satisfying the guard. Thus non-termination subsumes integer programming, which is **NP**-hard. By contrast, even though not stated explicitly in [38] and [8], deciding termination on \mathbb{R}^d and \mathbb{Q}^d can be done in polynomial time.²

While our algorithm for deciding termination requires exponential space, it should be noted that the procedure actually solves a more general problem than merely determining the existence of a non-terminating integer point (or, equivalently, the existence of an eventually non-terminating integer point). In fact the algorithm computes a representation of the set of all eventually non-terminating integer points. For reference, the closely related problem of deciding termination on the integer points in a given convex polytope is **EX-PSPACE**-hard [5].

As well as making extensive use of algorithms in real algebraic geometry, the soundness of our decision procedure relies on powerful lower bounds in Diophantine approximation that generalise Roth's Theorem. (The need for such bounds in the inhomogeneous setting was conjectured in the discussion in the conclusion of [8].) We also use classical results in number theory, such as the Skolem-Mahler-Lech Theorem [22, 25, 36] on linear recurrences. Crucially the well-known and notorious ineffectiveness of Roth's Theorem (and its higher-dimensional and *p*-adic generalisations) and of the Skolem-Mahler-Lech Theorem are not a problem for deciding *eventual* non-termination, which is key to our approach.

1.1 Related Work Consider the termination problem for a homogeneous linear loop program

 $P2: \boldsymbol{x} \leftarrow \boldsymbol{u} ; while \ B\boldsymbol{x} \geq 0 \ do \ \boldsymbol{x} \leftarrow A\boldsymbol{x}$

on a single initial value $\boldsymbol{u} \in \mathbb{Z}^d$. Each row \boldsymbol{b}^T of matrix B corresponds to a loop condition $\boldsymbol{b}^T \boldsymbol{x} \geq 0$. For each such condition, consider the integer sequence $\langle x_n : n \in \mathbb{N} \rangle$ defined by $x_n = \boldsymbol{b}^T A^n \boldsymbol{u}$. Then P2 fails to terminate on an initial value \boldsymbol{u} if and only if each such sequence $\langle x_n \rangle$ is *positive*, i.e., $x_n \geq 0$ for all n. It is not difficult to show that each sequence $\langle x_n \rangle$ considered above is a *linear recurrence sequence*, thanks to the Cayley-Hamilton theorem. Thus deciding whether a homogeneous linear loop program terminates on a given initial value is at least as hard as the *Positivity Problem* for linear recurrence sequences, that is, the problem of deciding whether a given linear recurrence sequence has exclusively non-negative terms.

The Positivity Problem has been studied at least as far back as the 1970s [4, 17, 24, 33, 34]. Thus far decidability is known only for sequences satisfying recurrences of order 5 or less. It is moreover known that showing decidability at order 6 will necessarily entail breakthroughs in transcendental number theory, specifically significant new results in Diophantine approximation [27].

The key difference between studying termination of simple linear loops over \mathbb{Z}^d rather than a single initial value is that the former problem can be approached through eventual termination. In this sense the termination problem is related to the *Ultimate Positivity Problem* for linear recurrence sequences, which asks whether all but finitely many terms of a given sequence are positive [28]. This allows us to bring to bear powerful noneffective Diophantine-approximation techniques, specifically the *S*-units Theorem of Evertse, van der Poorten, and Schlickewei [16, 39]. Such tools enable us to obtain decidability of termination for matrices of arbitrary dimension, assuming diagonalisability.

The paper [11] studies higher dimensional versions of Kannan and Lipton's Orbit Problem [20]. These can be seen as versions of the termination problem for linear loops on a fixed initial value. That work uses substantially different technology from that of the current paper, including Baker's Theorem on linear forms in logarithms [2], and correspondingly relies on restrictions on the dimension of data in problem instances to obtain decidability.

Termination of P1 under the assumption that all eigenvalues of A are real was studied in [32, 31] using spectral techniques. However, as will become clear throughout the course of this paper, most of the machinery that we use is needed to tackle the case where there are both real and complex eigenvalues with the same absolute value. In the setting of [32, 31], the set of eventually non-terminating points is in fact a polytope, which can be effectively computed resorting only to straightforward linear algebra.

While we use spectral and number-theoretic techniques in this paper, another well-studied approach for

¹By contrast, recall that the existence of an integer point in an *arbitrary* (i.e., not necessarily convex) semi-algebraic set which is equivalent to Hilbert's tenth problem—is well-known to be undecidable.

 $^{^{2}}$ This observation relies on the facts that one can compute Jordan canonical forms of integer matrices and solve instances of linear programming problems with algebraic numbers in polynomial time [9, 1].

proving termination of linear loops involves designing linear ranking functions, that is, linear functions from the state space to a well-founded domain such that each iteration of the loop strictly decreases the value of the ranking function. However, this approach is incomplete: it is not hard to construct an example of a terminating loop which admits no linear ranking function. Sound and relatively complete methods for synthesising linear ranking functions can be found in [30] and [6]. Whether a linear ranking function exists can be decided in polynomial time when the state space is \mathbb{Q}^d and is **coNP**complete when the state space is \mathbb{Z}^d .

2 Overview of Main Results

The main result of this paper is as follows:

THEOREM 2.1. The termination over the integers of simple linear loops of the form

$$P1: x \leftarrow u$$
; while $Bx \ge c$ do $x \leftarrow Ax + a$

is decidable using exponential space if A is diagonalisable and using polynomial space if A has dimension at most 4.

In this section we give a high-level overview of the proof of Theorem 2.1.

Let $f : \mathbb{R}^d \to \mathbb{R}^d$ be the affine function $f(\mathbf{x}) = A\mathbf{x} + \mathbf{a}$ computed by the body of the while loop in P1 and $P = \{\mathbf{x} \in \mathbb{R}^d : B\mathbf{x} \ge \mathbf{c}\}$ the convex polytope corresponding to the loop guard. We define the set of *non-terminating points* to be

$$NT = \{ \boldsymbol{u} \in \mathbb{R}^d : \forall n \in \mathbb{N}, f^n(\boldsymbol{u}) \in P \}.$$

Following Braverman [8], we moreover define the set of *eventually non-terminating points* to be

$$ENT = \{ \boldsymbol{u} \in \mathbb{R}^d : \exists n \in \mathbb{N}, f^n(\boldsymbol{u}) \in NT \}.$$

It is easily seen from the above definitions that both NT and ENT are convex sets.

By definition, P1 is non-terminating on \mathbb{Z}^d if and only if *NT* contains an integer point. It is moreover clear that *NT* contains an integer point if and only if *ENT* contains an integer point.

Recall that a subset of \mathbb{R}^d is said to be *semi-algebraic* if it is a Boolean combination of sets of the form $\{x \in \mathbb{R}^d : p(x) \ge 0\}$, where p is a polynomial with integer coefficients. Equivalently the semi-algebraic sets are those definable by quantifier-free first-order formulas over the structure $(\mathbb{R}, <, +, \cdot, 0, 1)$. In fact, since the first-order theory of the reals admits quantifier elimination [37], the semi-algebraic sets are precisely the first-order definable sets.

Define $W \subseteq \mathbb{R}^d$ to be a *non-termination witness* set (or simply a witness set) if it satisfies the following two properties (where \mathbb{A} denotes the set of algebraic numbers):

- (i) W is convex and semi-algebraic;
- (ii) $W \cap \mathbb{A}^d = ENT \cap \mathbb{A}^d$.

The integer points in a witness set W are precisely the integer points of ENT, and so P1 is non-terminating on \mathbb{Z}^d precisely when W contains an integer point. Our approach to solving the termination problem consists in computing a witness set W for a given program and then using the following theorem of Khachiyan and Porkolab [21] to decide whether W contains an integer point.

THEOREM 2.2. (KHACHIYAN AND PORKOLAB) Let $W \subseteq \mathbb{R}^d$ be a convex semi-algebraic set defined by polynomials of degree at most D and that can be represented in space S. In that case, if $W \cap \mathbb{Z}^d \neq \emptyset$, then W must contain an integral point that can be represented in space $SD^{O(d^4)}$.

Our approach does not attempt to characterise the set ENT directly, but rather uses the witness set W as a proxy. However, our techniques do allow us to establish that $\overline{ENT} = \overline{W}$, which in particular implies that \overline{ENT} is semi-algebraic, since the closure of a semi-algebraic set is semi-algebraic (see the Appendix for details). A natural question is whether the set ENT itself is semi-algebraic, which we leave as an open problem.

We next describe some restrictions on linear loops that can be made without loss of generality and that will ease our upcoming analysis.

We first reduce the problem of computing witness sets in the general case to the same problem in the homogeneous case. Note that Program P1 terminates on a given initial value $u \in \mathbb{Z}^d$ if and only if the homogeneous program P3 below terminates for the same value of u:

$$\mathsf{P3}: \boldsymbol{x} \leftarrow \begin{pmatrix} \boldsymbol{u} \\ 1 \end{pmatrix} \textit{while} \begin{pmatrix} B & -\boldsymbol{c} \end{pmatrix} \boldsymbol{x} \geq 0 \textit{ do } \boldsymbol{x} \leftarrow \begin{pmatrix} A & \boldsymbol{a} \\ 0 & 1 \end{pmatrix} \boldsymbol{x}$$

Note that if A is diagonalisable then all eigenvalues of $\begin{pmatrix} A & a \\ 0 & 1 \end{pmatrix}$ are simple, with the possible exception of the eigenvalue 1. (Recall that an eigenvalue is said to be simple if it has multiplicity one as a root of the minimal polynomial of A.) Now if W is a witness set for program P3 then $\left\{ u \in \mathbb{R}^d : \begin{pmatrix} u \\ 1 \end{pmatrix} \in W \right\}$ is a witness set for P1. We conclude that, in order to settle the inhomogeneous case with a diagonalisable matrix, it suffices to compute a witness set in the case of a homogeneous linear loop P2 in which the only repeated eigenvalues of the new matrix A are positive and real. Likewise, to handle the inhomogeneous case for matrices of dimension at most d, it suffices to be able to compute witness sets in the homogeneous case for matrices of dimension at most $d+1.^3$

We can further simplify the homogeneous case by restricting to loop guards that comprise a single linear inequality. To see this, first note that program P2 above is eventually non-terminating on \boldsymbol{u} if and only if for each row \boldsymbol{b}^T of B program P4 below is eventually nonterminating on \boldsymbol{u} :

$$\mathsf{P4}: \boldsymbol{x} \leftarrow \boldsymbol{u} ; while \ \boldsymbol{b}^T \boldsymbol{x} \geq 0 \ do \ \boldsymbol{x} \leftarrow A \boldsymbol{x}.$$

Noting that the finite intersection of convex semialgebraic sets is again convex and semi-algebraic, we can compute a witness set for P2 as the intersection of witness sets for each version of P4.

The final simplification concerns the notion of nondegeneracy. We say that matrix A is *degenerate* if it has distinct eigenvalues $\lambda_1 \neq \lambda_2$ whose quotient λ_1/λ_2 is a root of unity.

Given an arbitrary matrix A, let L be the least common multiple of all orders of quotients of distinct eigenvalues of A which are roots of unity. It is known that $L = 2^{O(d\sqrt{\log d})}$ [15]. The eigenvalues of the matrix A^L have the form λ^L for λ an eigenvalue of A, by the spectral mapping theorem. It follows that A^L is nondegenerate, since if λ_1, λ_2 are eigenvalues of A such that λ_1^L/λ_2^L is a root of unity then λ_1/λ_2 is a root of unity and hence $\lambda_1^L/\lambda_2^L = 1$. Note that all eigenvectors of A are still eigenvectors of A^L , thus A^L will be diagonalisable whenever A is.

Now program P4 is eventually non-terminating on $u \in \mathbb{Z}^d$ if and only if program P5 below is eventually non-terminating on the set $\{u, Au, \ldots, A^{L-1}u\}$:

$$\mathsf{P5}: \boldsymbol{x} \leftarrow \boldsymbol{v} ; while \ \boldsymbol{b}^T \boldsymbol{x} \ge 0 \ do \ \boldsymbol{x} \leftarrow A^L \boldsymbol{x}.$$

Thus if W is a witness set for P5 then $\bigcap_{i=0}^{L-1} \{ u \in \mathbb{Z}^d : A^i u \in W \}$ is a witness set for P4.

The main technical result of the paper is the following proposition:

PROPOSITION 2.1. Given a homogeneous simple linear loop

$$\mathsf{P4}: \boldsymbol{x} \leftarrow \boldsymbol{u} \; ; \; while \; \boldsymbol{b}^T \boldsymbol{x} \geq 0 \; do \; \boldsymbol{x} \leftarrow A \boldsymbol{x},$$

such that A is non-degenerate and either A has dimension at most 5 or all complex eigenvalues of A are simple, we can compute a witness set for P4 using exponential space if A is diagonalisable and using polynomial space if A has dimension at most 5.

Bearing in mind that the transformation from P1 to P4 increases the dimension of A by one and does not introduce repeated complex eigenvalues, it follows from Proposition 2.1 that we can also compute witness sets for simple linear loops of the form P1 under the assumptions of Theorem 2.1, and thus we obtain the decidability part of Theorem 2.1. The exponential-space bound in Theorem 2.1 is obtained by bounding the representation of the witness set in Proposition 2.1 (see Section 5).

In the rest of this section we give a brief summary of the proof of Proposition 2.1.

To compute a witness set W for P4 we first partition the eigenvalues of the update matrix A by grouping eigenvalues of equal modulus. Correspondingly we write \mathbb{R}^d as a direct sum $\mathbb{R}^d = V_1 \oplus \ldots \oplus V_m$, where each subspace V_i is the sum of (generalised) eigenspaces of Aassociated to eigenvalues of the same modulus. Assume that V_1 corresponds to the eigenvalues of maximum modulus, V_2 the next greatest modulus, etc. Then there are two main steps in the construction of W:

- 1. By analysing multiplicative relationships among eigenvalues of the same modulus, we show that for each subspace V_i the set $ENT \cap V_i$ of eventually non-terminating initial values in V_i is semi-algebraic.
- 2. Given $\boldsymbol{v} \in \mathbb{R}^d$, we can write $\boldsymbol{v} = \boldsymbol{v}_1 + \ldots + \boldsymbol{v}_m$, with $\boldsymbol{v}_i \in V_i$. Using Theorem 7.2 on S-units, we show that if all entries of \boldsymbol{v} are algebraic numbers then the eventual non-termination of P4 on \boldsymbol{v} is a function of its eventual non-termination on each \boldsymbol{v}_i separately. More precisely we look for the first \boldsymbol{v}_i such that the sequence $\langle \boldsymbol{b}^T A^n \boldsymbol{v}_i : n \in \mathbb{N} \rangle$ is infinitely often non-zero. Then P4 is eventually non-terminating on \boldsymbol{v} if and only if it is eventually non-terminating on \boldsymbol{v}_i .

The computability of a witness set W easily follows from items 1 and 2 above. Our techniques require that the update matrix in the original linear loop P1 either be diagonalisable or have dimension at most 4. Eliminating these restrictions seems to require solving the Ultimate Positivity Problem for linear recurrence sequences of order greater than 5, which in turn requires solving hard open problems in the theory of Diophantine approximation [27].

³Note that whilst Braverman [8] shows how to decide termination over the integers for homogeneous programs with arbitrary update matrices, he does *not* compute a witness set for such programs—indeed this remains an open problem since it would enable one to solve termination over the integers for arbitrary inhomogeneous programs.

3 Groups of Multiplicative Relations

This section introduces some concepts concerning groups of multiplicative relations among algebraic numbers. Here we will assume some basic notions from algebraic number theory and the first-order theory of reals. We assume also a natural first-order interpretation of the field of complex numbers in the ordered field of real numbers (in which each complex number is encoded as a pair comprising its real and imaginary parts). Under this interpretation we refer to sets of complex numbers as being semi-algebraic and first-order definable. Details of the relevant notions can be found in the Appendix.

Let $\mathbb{T} = \{z \in \mathbb{C} : |z| = 1\}$. We define the *s*dimensional torus to be \mathbb{T}^s , considered as a group under componentwise multiplication.

Given a tuple of algebraic numbers $\boldsymbol{\lambda} = (\lambda_1, \ldots, \lambda_s)$, in this section we consider how to effectively represent the *orbit* $\{\boldsymbol{\lambda}^n : n \in \mathbb{N}\}$. More precisely, we will give an algebraic representation of the topological closure of that orbit in \mathbb{T}^s .

The group of multiplicative relations of λ , which is an additive subgroup of \mathbb{Z}^s , is defined as

$$L(\boldsymbol{\lambda}) = \{ \boldsymbol{v} \in \mathbb{Z}^s : \boldsymbol{\lambda}^{\boldsymbol{v}} = 1 \},\$$

where $\boldsymbol{\lambda}^{\boldsymbol{v}}$ is defined to be $\lambda_1^{v_1} \cdots \lambda_s^{v_s}$ for $\boldsymbol{v} \in \mathbb{Z}^s$, that is, exponentiation acts coordinatewise.

Since \mathbb{Z}^s is a free abelian group, its subgroups are also free. In particular, $L(\lambda)$ has a finite basis. The following powerful theorem of Masser [26] gives bounds on the magnitude of the components of such a basis.

THEOREM 3.1. (MASSER) The free abelian group $L(\lambda)$ has a basis $v_1, \dots, v_l \in \mathbb{Z}^s$ for which

$$\max_{1 \le i \le l, 1 \le j \le s} |v_{i,j}| \le (D \log H)^{O(s^2)}$$

where H and D bound respectively the heights and degrees of all the λ_i .

Membership of a tuple $v \in \mathbb{Z}^s$ in $L(\lambda)$ can be computed in polynomial space, using a decision procedure for the existential theory of the reals. In combination with Theorem 3.1, it follows that we can compute a basis for $L(\lambda)$ in polynomial space by brute-force search.

Corresponding to $L(\lambda)$, we consider the following multiplicative subgroup of \mathbb{T}^s :

$$T(\boldsymbol{\lambda}) = \{ \boldsymbol{\mu} \in \mathbb{T}^s : \forall \boldsymbol{v} \in L(\boldsymbol{\lambda}), \, \boldsymbol{\mu}^{\boldsymbol{v}} = 1 \}$$

If V is a basis of $L(\lambda)$ then we can equivalently characterise $T(\lambda)$ as $\{\mu \in \mathbb{T}^s : \forall v \in V, \mu^v = 1\}$. Crucially, this finitary characterisation allows us to represent $T(\lambda)$ as a semi-algebraic set. We will use the following classical lemma of Kronecker on simultaneous Diophantine approximation, in order to show that the orbit $\{\lambda^n : n \in \mathbb{N}\}$ is a dense subset of $T(\lambda)$.

LEMMA 3.1. Let $\boldsymbol{\theta}, \boldsymbol{\psi} \in \mathbb{R}^s$. Suppose that for all $\boldsymbol{v} \in \mathbb{Z}^s$, if $\boldsymbol{v}^T \boldsymbol{\theta} \in \mathbb{Z}$ then also $\boldsymbol{v}^T \boldsymbol{\psi} \in \mathbb{Z}$, i.e., all integer relations among the coordinates of $\boldsymbol{\theta}$ also hold among those of $\boldsymbol{\psi}$ (modulo \mathbb{Z}). Then, for each $\varepsilon > 0$, there exist $\boldsymbol{p} \in \mathbb{Z}^s$ and a non-negative integer n such that

$$\|n\boldsymbol{\theta} - \boldsymbol{p} - \boldsymbol{\psi}\|_{\infty} \leq \varepsilon$$
.

We now arrive at the main result of the section:

THEOREM 3.2. Let $\lambda \in \mathbb{T}^s$. Then the orbit $\{\lambda^n : n \in \mathbb{N}\}$ is a dense subset of $T(\lambda)$.

Proof. Let $\boldsymbol{\theta} \in \mathbb{R}^s$ be such that $\boldsymbol{\lambda} = e^{2\pi i \boldsymbol{\theta}}$ (with exponentiation operating coordinatewise). Notice that $\boldsymbol{\lambda}^v = 1$ if and only if $\boldsymbol{v}^T \boldsymbol{\theta} \in \mathbb{Z}$. If $\boldsymbol{\mu} \in T(\boldsymbol{\lambda})$, we can likewise define $\boldsymbol{\psi} \in \mathbb{R}^s$ to be such that $\boldsymbol{\mu} = e^{2\pi i \boldsymbol{\psi}}$. Then the premisses of Kronecker's lemma apply to $\boldsymbol{\theta}$ and $\boldsymbol{\psi}$. Thus, given $\varepsilon > 0$, there exist a non-negative integer nand $\boldsymbol{p} \in \mathbb{Z}^s$ such that $\|n\boldsymbol{\theta} - \boldsymbol{p} - \boldsymbol{\psi}\|_{\infty} \leq \varepsilon$. Whence

$$\|\boldsymbol{\lambda}^{n} - \boldsymbol{\mu}\|_{\infty} = \|e^{2\pi i (n\boldsymbol{\theta} - \boldsymbol{p})} - e^{2\pi i \boldsymbol{\psi}}\|_{\infty} \leq \\ \|2\pi (n\boldsymbol{\theta} - \boldsymbol{p} - \boldsymbol{\psi})\|_{\infty} \leq 2\pi\varepsilon \,.$$

4 Algorithm for Universal Termination

Our goal in this section is to prove the following proposition, which is restated from Section 2. We have already shown in Section 2 that the main result of this paper, Theorem 2.1, then follows.

PROPOSITION 4.1. Given a homogeneous simple linear loop

$$\mathsf{P4}: \boldsymbol{x} \leftarrow \boldsymbol{u} ; while \ \boldsymbol{b}^T \boldsymbol{x} \ge 0 \ do \ \boldsymbol{x} \leftarrow A \boldsymbol{x},$$

such that A is non-degenerate, we can compute a witness set for P4 using exponential space if all complex eigenvalues of A are simple and using polynomial space if A has dimension at most 5.

Define the *index* of an eigenvalue of A to be its multiplicity as a root of the minimal polynomial of A. An eigenvalue is said to be *simple* if it has index 1 and *repeated* otherwise. We can write matrix A in the form $A = P^{-1}JP$ for some invertible matrix P and block diagonal Jordan matrix $J = Diag(J_1, \ldots, J_N)$, with each block J_i having the form

λ	1	0		$0 \rangle$	
0	λ	1		0	
:	÷	÷	۰.	:	
0	0	0	۰.	1	
0	0	$\begin{array}{c} 0 \\ 0 \end{array}$		λ)	

where λ is an eigenvalue of A whose index equals the the dimension of the block. The entries of P are all algebraic numbers lying in the extension field of \mathbb{Q} generated by the eigenvalues of A.

The *n*-th power of the matrix J has the form $J^n = Diag(J_1^n, \ldots, J_N^n)$, where each block J_i^n has the form

$$\begin{pmatrix} \lambda^{n} & n\lambda^{n-1} & \binom{n}{2}\lambda^{n-2} & \dots & \binom{n}{\nu-1}\lambda^{n-\nu+1} \\ 0 & \lambda^{n} & n\lambda_{i}^{n-1} & \dots & \binom{n}{\nu-2}\lambda^{n-\nu+2} \\ \vdots & \vdots & \vdots & \ddots & \vdots \\ 0 & 0 & 0 & \ddots & n\lambda^{n-1} \\ 0 & 0 & 0 & \dots & \lambda^{n} \end{pmatrix},$$

where λ is an eigenvalue of A of index ν , and $\binom{n}{k} = 0$ if n < k.

Let A have eigenvalues $\lambda_1, \ldots, \lambda_l$, with respective indices ν_1, \ldots, ν_l . Given $\boldsymbol{u} \in \mathbb{R}^d$, from our observations on the form of J^n , we can write

(4.1)
$$\boldsymbol{b}^T A^n \boldsymbol{u} = \sum_{j=1}^l \sum_{k=0}^{\nu_j - 1} \boldsymbol{\alpha}_{j,k}^T \boldsymbol{u} \, n^k \lambda_j^n \,,$$

where the $\alpha_{j,k}$ are vectors of algebraic numbers that do not depend on \boldsymbol{u} , and the equation holds for all $n \geq d$.

Since the characteristic polynomial of A has integer coefficients, the eigenvalues of A are all algebraic integers. Moreover, since for any positive integer t > 0 we have that $t \cdot \boldsymbol{b}^T A^n \boldsymbol{u} \ge 0$ if and only if $\boldsymbol{b}^T A^n \boldsymbol{u} \ge 0$, by rescaling we can assume that the vectors $\boldsymbol{\alpha}_{j,k}$ in (4.1) are comprised of algebraic integers.

Now let us partition the eigenvalues of A into sets S_1, \ldots, S_m by grouping eigenvalues of equal modulus. Assume that S_1 contains eigenvalues of maximum modulus, S_2 eigenvalues of the next greatest modulus, etc. Correspondingly we write \mathbb{R}^d as a direct sum of subspaces $\mathbb{R}^d = V_1 \oplus \ldots \oplus V_m$, where each subspace V_i is the sum of (generalised) eigenspaces of A associated to eigenvalues in S_i . By the assumption that A is nondegenerate, i.e., that no quotient of two distinct eigenvalues is a root of unity, S_i cannot have both a positive and a negative real eigenvalue of the same modulus. Thus each set S_i contains at most one real eigenvalue.

4.1 Eventual Non-Termination on Subspace V_i We first consider the eventual non-termination of P4 on initial vectors in the subspace V_i for a fixed $i \in$ $\{1, \ldots, m\}$. Writing $ENT_i := ENT \cap V_i$, our goal is to show that ENT_i is semi-algebraic.

Given $\boldsymbol{u} \in V_i$, membership of \boldsymbol{u} in ENT_i can be characterised in terms of the *ultimate positivity* of the sequence $\langle \boldsymbol{b}^T A^n \boldsymbol{u} : n \in \mathbb{N} \rangle$. More precisely, $\boldsymbol{u} \in ENT_i$ if and only if $\boldsymbol{b}^T A^n \boldsymbol{u} \geq 0$ for all but finitely many n. In particular, defining

$$ZERO := \{ \boldsymbol{u} \in \mathbb{R}^d : \forall n \ge d, \, \boldsymbol{b}^T A^n \boldsymbol{u} = 0 \}$$

and $ZERO_i := ZERO \cap V_i$, we have that $ZERO_i \subseteq ENT_i$.

It is easy to see that $ZERO_i$ is semi-algebraic. Indeed the uniqueness part of [18, Proposition 2.11] implies that $\boldsymbol{b}^T A^n \boldsymbol{u} = 0$ for all $n \geq d$ if and only if each term $n^k \lambda_j^n$ has coefficient zero in the expression (4.1). Thus

$$ZERO = \left\{ \boldsymbol{u} \in \mathbb{R}^d : \bigwedge_{j=1}^l \bigwedge_{k=0}^{\nu_j - 1} \boldsymbol{\alpha}_{j,k}^T \boldsymbol{u} = 0
ight\}.$$

is semi-algebraic. Since V_i is a semi-algebraic subset of \mathbb{R}^d , being spanned by a subset of the columns of P, it follows that $ZERO_i$ is semi-algebraic.

PROPOSITION 4.2. The set ENT_i is semi-algebraic for each $i \in \{1, \ldots, m\}$.

Proof. We consider three (overlapping) cases. Under the hypotheses of Proposition 4.1 at least one of these cases will apply.

Case I: A has dimension at most 5. Assume that A has dimension at most 5. The situations in which S_i does not contain a positive real eigenvalue, or all of the complex eigenvalues in S_i are simple, will be handled under Cases II and III, below. Otherwise, let $\lambda \in S_i$ be a complex eigenvalue of index at least 2. Since A has dimension at most 5, it must be the case that λ and its complex conjugate $\overline{\lambda}$ both have index exactly 2. Let $\rho \in S_i$ be the positive real eigenvalue. Since A has dimension at most 5, ρ must be simple. Thus $S_i = \{\rho, \lambda, \overline{\lambda}\}$ contains all the eigenvalues of A.

For $\boldsymbol{u} \in V_i$ we can write

$$\boldsymbol{b}^{T} A^{n} \boldsymbol{u} = \left(\boldsymbol{\alpha}_{0} \rho^{n} + (\boldsymbol{\beta}_{0} + \boldsymbol{\beta}_{1} n) \lambda^{n} + \overline{(\boldsymbol{\beta}_{0} + \boldsymbol{\beta}_{1} n) \lambda^{n}}\right)^{T} \boldsymbol{u},$$

for all $n \ge d$, where $\boldsymbol{\alpha}_0$ is a vector of real algebraic numbers, $\boldsymbol{\beta}_0, \boldsymbol{\beta}_1$ are vectors of complex algebraic numbers. If $\boldsymbol{\beta}_1^T \boldsymbol{u} \ne 0$, then as *n* tends to infinity the dom-

If $\beta_1^I \boldsymbol{u} \neq 0$, then as *n* tends to infinity the dominant terms on the right-hand side above are constant multiples of $n\lambda^n$ and $n\overline{\lambda^n}$. In this case it follows from [8, Lemma 4] that $\boldsymbol{b}^T A^n \boldsymbol{u}$ changes sign infinitely often as *n* grows, and hence $\boldsymbol{u} \notin ENT_i$.

The argument in case $\beta_1^T u = 0$ is a simple version of the approach in Case III, however we include details since the reader may find this special case instructive.

Define $f : \mathbb{T} \to \mathbb{R}$ by

$$f(z) = \boldsymbol{\alpha}_0^T \boldsymbol{u} + \boldsymbol{\beta}_0^T \boldsymbol{u} z + \overline{\boldsymbol{\beta}_0^T \boldsymbol{u} z}.$$

Then $\boldsymbol{b}^T A^n \boldsymbol{u} = \rho^n f(\lambda^n / \rho^n)$ for all $n \ge d$.

Since A is assumed to be non-degenerate, λ/ρ is not a root of unity. Thus $\{\lambda^n/\rho^n : n \in \mathbb{N}\}$ is dense in \mathbb{T} . It follows that $\boldsymbol{u} \in ENT$ if and only if $f(z) \geq 0$ for all $z \in \mathbb{T}$. By inspection this last condition is equivalent to $\boldsymbol{\alpha}_0^T \boldsymbol{u} \geq 2|\boldsymbol{\beta}_0^T \boldsymbol{u}|$. We conclude that

$$ENT_i = \left\{ \boldsymbol{u} \in V_i : \boldsymbol{\beta}_1^T \boldsymbol{u} = 0 \land \boldsymbol{\alpha}_0^T \boldsymbol{u} \ge 2|\boldsymbol{\beta}_0^T \boldsymbol{u}| \right\},\$$

and hence ENT_i is semi-algebraic.

Case II: S_i does not contain a positive real eigenvalue. It follows from [8, Lemma 4] that if S_i does not contain a positive real eigenvalue then for $\boldsymbol{u} \in V_i$ the sequence $\boldsymbol{b}^T A^n \boldsymbol{u}$ is either identically zero for $n \geq d$ or is infinitely often strictly positive and infinitely often strictly negative. Thus in this case $ENT_i = ZERO_i$. But we have already shown that $ZERO_i$ is semi-algebraic.

Case III: all complex eigenvalues in S_i are simple. Suppose that all complex eigenvalues in S_i are simple. If S_i contains no positive real eigenvalue then Case II applies. Thus we may assume that S_i comprises a positive real eigenvalue ρ of index t and simple complex eigenvalues $\lambda_1, \overline{\lambda_1}, \ldots, \lambda_s, \overline{\lambda_s}$. Given $u \in V_i$ we can write

$$\boldsymbol{b}^{T} A^{n} \boldsymbol{u} = \boldsymbol{b}^{T} P^{-1} J^{n} P \boldsymbol{u}$$

$$(4.2) \qquad = \left[\sum_{j=0}^{t-1} \boldsymbol{\alpha}_{j} n^{j} \rho^{n} + \sum_{j=1}^{s} (\boldsymbol{\beta}_{j} \lambda_{j}^{n} + \overline{\boldsymbol{\beta}_{j} \lambda_{j}^{n}}) \right]^{T} \boldsymbol{u} ,$$

where the α_j and β_j are *d*-dimensional vectors of algebraic numbers, with all coefficients of each α_j being real.

Since $\rho = |\lambda_1| = \ldots = |\lambda_s|$, if $\alpha_j^T \boldsymbol{u} \neq 0$ for some strictly positive index j, then, for the largest such index j, the term $n^j \rho^n \alpha_j^T \boldsymbol{u}$ is dominating on the right-hand side of (4.2). In particular, if $\alpha_j^T \boldsymbol{u} > 0$ then the sequence $\boldsymbol{b}^T A^n \boldsymbol{u}$ is ultimately positive (hence $\boldsymbol{u} \in ENT_i$), and if $\alpha_j^T \boldsymbol{u} < 0$ then $\boldsymbol{b}^T A^n \boldsymbol{u}$ is not ultimately positive (hence $\boldsymbol{u} \notin ENT_i$). It follows that

(4.3)
$$\left\{ \boldsymbol{u} \in V_i : \bigvee_{j=1}^{t-1} \bigwedge_{k=j+1}^{t-1} (\boldsymbol{\alpha}_j^T \boldsymbol{u} > 0 \land \boldsymbol{\alpha}_k^T \boldsymbol{u} = 0) \right\}$$

is a subset of ENT_i .

The case that $\boldsymbol{\alpha}_j^T \boldsymbol{u} = 0$ for all $j = 1, \dots, t-1$ is more subtle since there is no single dominant term in (4.2); this is where we employ the results of Section 3 on multiplicative relations. In this case we rewrite (4.2) as

(4.4)
$$\boldsymbol{b}^T A^n \boldsymbol{u} = \rho^n f \left(\frac{\lambda_1^n}{\rho^n}, \dots, \frac{\lambda_s^n}{\rho^n}\right)^T \boldsymbol{u},$$

where $f: \mathbb{T}^s \to \mathbb{R}^d$ is defined by

$$f(z_1,\ldots,z_s) = \boldsymbol{\alpha}_0 + \sum_{j=1}^s \boldsymbol{\beta}_j z_j + \overline{\boldsymbol{\beta}_j z_j}.$$

Defining $\boldsymbol{\mu} = (\lambda_1/\rho, \dots, \lambda_s/\rho)$, we furthermore rewrite (4.4) as

(4.5)
$$\boldsymbol{b}^T A^n \boldsymbol{u} = \rho^n f(\boldsymbol{\mu}^n)^T \boldsymbol{u}.$$

By Theorem 3.2, $\{\boldsymbol{\mu}^n : n \in \mathbb{N}\}$ is a dense subset of the torus $T(\boldsymbol{\mu})$. Thus the right-hand side of (4.5) is non-negative for every n if and only if $f(\boldsymbol{z})^T \boldsymbol{u} \geq 0$ for all $\boldsymbol{z} \in T(\boldsymbol{\mu})$. It follows that

(4.6)
$$\left\{ \boldsymbol{u} \in V_{\boldsymbol{i}} : \forall \boldsymbol{z} \in T(\boldsymbol{\mu}), \, f(\boldsymbol{z})^{T} \boldsymbol{u} \ge 0 \right\}$$

is a subset of ENT_i .

In Section 3 we observed that the set $T(\boldsymbol{\mu})$ was (effectively) semi-algebraic. It follows that we can express the condition $\forall \boldsymbol{z} \in T(\boldsymbol{\mu}), f(\boldsymbol{z})^T \boldsymbol{u} \geq 0$ in the first-order theory of the reals. By the Tarski-Seidenberg theorem [37] on quantifier elimination, the set of $\boldsymbol{u} \in \mathbb{R}^d$ satisfying this condition is semi-algebraic. But now ENT_i is the union of the two semi-algebraic sets (4.3) and (4.6), and therefore ENT_i is itself semi-algebraic.

4.2 Definition of a Witness Set Having shown that $ZERO_i$ and ENT_i are semi-algebraic sets for $i = 1, \ldots, m$, we now define a witness set W for the loop P4.

Given $\boldsymbol{u} \in \mathbb{R}^d$, write $\boldsymbol{u} = \boldsymbol{u}_1 + \ldots + \boldsymbol{u}_m$, with $\boldsymbol{u}_1 \in V_1, \ldots, \boldsymbol{u}_m \in V_m$. Say that \boldsymbol{u}_i is the dominant component of \boldsymbol{u} if $\boldsymbol{u}_i \notin ZERO_i$ and $\boldsymbol{u}_j \in ZERO_j$ for all j < i. The intuition is that if \boldsymbol{u}_i is dominant then the eventual non-termination of P4 on \boldsymbol{u} is determined by its eventual non-termination on \boldsymbol{u}_i . However, to prove this we need to assume $\boldsymbol{u} \in (\mathbb{A} \cap \mathbb{R})^d$. Formally we have:

PROPOSITION 4.3. If u_i is the dominant component of $u \in (\mathbb{A} \cap \mathbb{R})^d$ then $u \in ENT$ if and only if $u_i \in ENT$.

Proof. From the fact that u_i is dominant we have:

(4.7)
$$\boldsymbol{b}^{T} A^{n} \boldsymbol{u} = \boldsymbol{b}^{T} A^{n} (\boldsymbol{u}_{1} + \dots + \boldsymbol{u}_{m})$$
$$= \boldsymbol{b}^{T} A^{n} (\boldsymbol{u}_{i} + \dots + \boldsymbol{u}_{m})$$

for all $n \ge d$. Moreover, for each j > i it is clear that $|\boldsymbol{b}^T A^n \boldsymbol{u}_j| = O(n^d \rho_j^n)$, where $\rho_j \ge 0$ is the modulus of the eigenvalues in S_j .

We now consider three cases, mirroring the proof of Proposition 4.2.

The first case is that A has dimension at most 5. As observed in the proof of Proposition 4.2, all instances of this case that are not already covered by the second and third cases are such that S_i contains all the eigenvalues of A, and hence $u_i = u$. In this situation the proposition holds trivially.

The second case is that S_i does not contain a positive real eigenvalue. Then it follows from [8, Lemma 4] that there is a constant c < 0 such that $\boldsymbol{b}^T A^n \boldsymbol{u}_i < c\rho_i^n$ for infinitely many n. In this case neither u_i nor u are elements of ENT.

It remains to consider the case that all complex eigenvalues in S_i are simple. Suppose that the dominant term in the expression for $\boldsymbol{b}^T A^n \boldsymbol{u}_i$ has the form $\alpha n^k \rho_i^n$ for some real constant $\alpha \neq 0$ and k > 0. If $\alpha > 0$ then both \boldsymbol{u} and \boldsymbol{u}_i are in ENT and if $\alpha < 0$ then neither \boldsymbol{u} or \boldsymbol{u}_i are in ENT.

Otherwise, specialising the expression (4.1) to the case at hand, we have that

(4.8)
$$\boldsymbol{b}^T A^n \boldsymbol{u}_i = \alpha_0 \rho_i^n + \sum_{j=1}^s \beta_j \lambda_j^n + \overline{\beta_j \lambda_j^n}$$

where α_0 and the β_j are algebraic-integer constants and $\rho_i, \lambda_1, \overline{\lambda_1}, \dots, \lambda_s, \overline{\lambda_s} \in S_i$. In this case one can use the S-units theorem of Evertse, van der Poorten, and Schlickewei [16, 39] to show that for all $\varepsilon > 0$ it is the case that $\boldsymbol{b}^T A^n \boldsymbol{u}_i = \Omega(\rho_i^n \Lambda^{-n\varepsilon})$, where Λ is an upper bound on the absolute value of eigenvalues of A (see the Appendix for details).

From this lower bound, taking ε suitably small, it follows that $|\boldsymbol{b}^T A^n \boldsymbol{u}_i| = o(|\boldsymbol{b}^T A^n \boldsymbol{u}_i|)$ for all j > i and hence that $\boldsymbol{u} \in ENT$ if and only if $\boldsymbol{u}_i \in ENT$.

Now we define a witness set W for program P4 by $W := \bigcup_{i=1}^{m} \{ \boldsymbol{u} \in \mathbb{R}^{d} : \boldsymbol{u}_{i} \text{ is the dominant component of } \boldsymbol{u}, \}$

 $u_i \in ENT \} \cup ZERO$.

From the fact that $ZERO_i$, ENT_i , and V_i are semialgebraic for $i = 1, \ldots, m$, it is easy to see that W is semi-algebraic. It moreover follows from Proposition 4.3 that $W \cap \mathbb{A}^d = ENT \cap \mathbb{A}^d$.

To conclude the proof of Proposition 4.1, it remains to observe that the witness set W, like the actual set ENT of eventually non-terminating points, is convex.

PROPOSITION 4.4. The witness set W is convex.

Proof. Suppose $\boldsymbol{y}, \boldsymbol{z} \in W$ and let $\boldsymbol{x} = \lambda \boldsymbol{y} + (1 - \lambda) \boldsymbol{z}$, where $0 < \lambda < 1$. Moreover, write $\boldsymbol{x} = \boldsymbol{x}_1 + \ldots + \boldsymbol{x}_m$, where $\boldsymbol{x}_1 \in V_1, \ldots, \boldsymbol{x}_m \in V_m$, and likewise for \boldsymbol{y} and \boldsymbol{z} .

If $y, z \in ZERO$ then $x \in ZERO$ since the latter is a convex set.

Suppose that $\boldsymbol{y} \in ZERO$ and $\boldsymbol{z}_i \in ENT$ is dominant for z for some index $i \in \{1, \ldots, m\}$. Then \boldsymbol{x}_i is dominant for \boldsymbol{x} , and $\boldsymbol{x}_i \in ENT$. Thus $\boldsymbol{x} \in W$.

Otherwise, let y_i be dominant for y and z_i be dominant for z for some $i, j \in \{1, \ldots, m\}$. Then $\boldsymbol{x}_k \in ZERO_k$ for all $k < \min\{i, j\}$ since $ZERO_k$ is convex. Moreover if $k = \min\{i, j\}$ then $y_k, z_k \in ENT_k$, and hence $x_k \in ENT_k$ by convexity of ENT_k . It follows that $\boldsymbol{x} \in W$.

This concludes the proof of Proposition 4.1. In the remaining part of this section we show that $\overline{ENT} = \overline{W}$.

The inclusion $\overline{W} \subseteq \overline{ENT}$ can be shown using the fact that the set algebraic points in any semi-algebraic set is dense in that set. (See the Appendix for details). From this we have:

$$\overline{W} = \overline{W \cap \mathbb{A}^d} = \overline{ENT \cap \mathbb{A}^d} \subseteq \overline{ENT} \cap \overline{\mathbb{A}^d} = \overline{ENT}$$

The reverse inclusion, $\overline{ENT} \subseteq \overline{W}$, can be shown in similar fashion but this time using the fact that $ENT \cap \mathbb{A}^d$ is dense in ENT. Our remaining goal is this last fact, which is established in Corollary 4.1 below.

We have previously shown that a vector of algebraic numbers $\boldsymbol{u} \in (\mathbb{A} \cap \mathbb{R})^d$ is eventually non-terminating if and only if its dominant component u_i is eventually non-terminating. We now prove a partial result of this nature for general vectors $\boldsymbol{u} \in \mathbb{R}^d$.

PROPOSITION 4.5. Suppose that $\boldsymbol{u} = \boldsymbol{u}_1 + \cdots + \boldsymbol{u}_m \in$ \mathbb{R}^d , where $\boldsymbol{u}_1 \in V_1, \ldots, \boldsymbol{u}_m \in V_m$. Then $\boldsymbol{u} \in ENT$ implies that its dominant component u_i is also in ENT.

Proof. The only non-trivial case corresponds to the situation in which $\boldsymbol{b}^T A^n \boldsymbol{u}_i$ is of the form (4.8). Let f and $\boldsymbol{\mu}$ be as in (4.4), that is, so that $\boldsymbol{b}^T A^n \boldsymbol{u}_i =$ $\rho_i^n f(\boldsymbol{\mu}^n)^T \boldsymbol{u}_i$. If $\boldsymbol{u}_i \notin ENT$, then there exists some constant c < 0 and some $\boldsymbol{z} \in T(\boldsymbol{\mu})$ such that $f(\boldsymbol{z})^T \boldsymbol{u}_i =$ c. Therefore, for any $\varepsilon > 0$, $\boldsymbol{b}^T A^n \boldsymbol{u}_i < (c+\varepsilon)\rho_i^n$ holds for infinitely many n, due to Proposition 3.2 and to continuity of f, and so $u \notin ENT$.

COROLLARY 4.1. $ENT \cap \mathbb{A}^d$ is dense in ENT.

Proof. At several points we will rely on the fact that if $X \subseteq \mathbb{R}^d$ is semi-algebraic, then the algebraic points in X are dense in X. (See Appendix for details.)

Fix $\boldsymbol{u} \in ENT$ and let $\varepsilon > 0$ be given. We will find $\boldsymbol{v} \in ENT \cap \mathbb{A}^d$ such that $||\boldsymbol{u} - \boldsymbol{v}|| < \varepsilon$.

The case in which $u \in ZERO$ is easy since ZERO is semi-algebraic and so we can take v to be an algebraic point in ZERO that is suitably close to \boldsymbol{u} .

Suppose now that $\boldsymbol{u} = \boldsymbol{u}_1 + \cdots + \boldsymbol{u}_m$, where $\boldsymbol{u}_1 \in$ $V_1,\ldots,\boldsymbol{u}_m \in V_m$, with \boldsymbol{u}_i the dominant component of u. By Proposition 4.5, $u \in ENT$ implies that $u_i \in ENT$. Since $ENT \cap V_i$ is semi-algebraic, we can pick $\boldsymbol{v}_i \in ENT \cap V_i \cap \mathbb{A}^d$ such that $\|\boldsymbol{v}_i - \boldsymbol{u}_i\| < \varepsilon/n$. For each j > i, we pick some $\boldsymbol{v}_j \in V_j \cap \mathbb{A}^d$ for which $\|\boldsymbol{v}_j - \boldsymbol{u}_j\| < \varepsilon/n$. For each j < i we pick some $v_j \in ZERO \cap V_j \cap \mathbb{A}^d$ for which $\|\boldsymbol{v}_j - \boldsymbol{u}_j\| < \varepsilon/n$.

Then, letting $\boldsymbol{v} = \boldsymbol{v}_1 + \cdots + \boldsymbol{v}_m \in \mathbb{A}^d$, it follows that $\|\boldsymbol{u} - \boldsymbol{v}\| < \varepsilon$. Finally, by Proposition 4.3 we have $\boldsymbol{v} \in ENT$ since \boldsymbol{v}_i is the dominant component of \boldsymbol{v} and $\boldsymbol{v}_i \in ENT$ by construction.

5 Complexity Analysis

The purpose of this section is to justify our previous claims about the complexity of the algorithm presented in this paper. We do this by proving the following result.

PROPOSITION 5.1. Our procedure requires space $poly(\log \max_{i,j} |A_{ij}|, d)^{poly(d)}$.

Proof. There are three critical steps in our procedure for which a super-polynomial amount of space is required: when reducing to the case in which A is non-degenerate, when performing quantifier elimination, and when testing whether the witness set W intersects the integer lattice.

The last of these steps runs in space $SD^{O(d^4)}$, where S denotes the size of the representation of the quantifierfree formula defining the witness set W, D denotes the maximum degree of the polynomials occurring in that formula, and d denotes the dimension of the ambient space. Since d remains fixed throughout the procedure (apart from an increase by 1 in the reduction to the homogeneous case), it remains to show that S and D are bounded by an expression of the form $poly(\log \max_{i,j} |A_{ij}|, d)^{poly(d)}$.

The reduction to the case in which A is non-degenerate entails an increase by a factor of $poly(\log \max_{i,j} |A_{ij}|, d)^{poly(d)}$ in the size of the formula defining the witness set W, as the least common multiple of the orders of all ratios of eigenvalues of A that are roots of unity is $L = 2^{O(d\sqrt{\log d})}$ and $\log \max_{i,j} |A_{ij}^L| \le \log(d^L \max_{i,j} |A_{ij}|^L) = L \log(d \max_{i,j} |A_{ij}|).$

It remains to show that the quantifier-free formula defining the witness set W in the case where A is nondegenerate takes space $poly(\log \max_{i,j} |A_{ij}|, d)^{poly(d)}$ and involves exclusively polynomials of degree $poly(\log \max_{i,j} |A_{ij}|, d)^{poly(d)}$.

Let D_0, H_0 denote the maximum degree and height across all the eigenvalues of A, respectively. Then $D_0 \leq d$ and $\log H_0 \leq \log(d! \max_{i,j} |A_{ij}|^d) \leq d \log(d \max_{i,j} |A_{ij}|)$. Before performing quantifier elimination, the degree of any polynomial in the defining formula of the witness set W is bounded by $(D_0 \log H_0)^{O(d^2)}$, and the number of such polynomials is bounded by O(d), by Masser's theorem. Finally, after applying quantifier elimination, we know that $D \leq$ $(D_0 \log H_0)^{O(d^3)}$ and that $S \leq d^{O(d^2)} (D_0 \log H_0)^{O(d^4)}$, thanks to Theorem 7.4.

6 Conclusion

We have shown decidability of termination of simple linear loops over the integers under the assumption that the update matrix is diagonalisable, partially answering an open problem of [38, 8]. As we have explained before, the termination problem on the same class of linear loops, but for fixed initial values, seems to have a different character and to be more difficult. In this respect it is interesting to note that there are other settings in which universal termination is an easier problem than pointwise termination. For example, universal termination of Petri nets (also known as *structural boundedness*) is **PTIME**-decidable, but the pointwise termination problem is **EXPSPACE**-hard.

A natural subject for further work is whether our techniques can be extended to non-diagonalisable matrices, or whether, as is the case for pointwise termination [27], there are unavoidable number-theoretic obstacles to proving decidability. We would also like to further study the computational complexity of the termination problem. While there is a large gap between the **coNP** lower complexity bound mentioned in the Introduction and the exponential space upper bound of our procedure, this may be connected with the fact that our procedure computes a representation of the set of all integer eventually non-terminating points. Finally we would like to examine more carefully the question of whether the respective sets of terminating and nonterminating points are semi-algebraic. Note that an *ef*fective semi-algebraic characterisation of the set of terminating points would allow us to solve the termination problem over fixed initial values.

Acknowledgements

The authors would like to thank Elias Koutsoupias and Ventsislav Chonev for their advice and feedback.

7 Appendix

7.1 Algebraic Numbers The purpose of this section is threefold: to introduce the main concepts in Algebraic Number Theory, necessary to understanding the hypothesis for the *S*-units theorem, stated below; to justify the application of the aforementioned result in lower-bounding the dominant terms of linear recurrence sequences; to explain how one can effectively manipulate algebraic numbers.

Preliminaries A complex number α is said to 7.2be **algebraic** if it is the root of some polynomial with integer coefficients. Among those polynomials, there exists a unique one of minimal degree whose coefficients have no common factor, and it is said to be the **defining polynomial** of α , denoted by p_{α} , and it is always an irreducible polynomial. Moreover, if p_{α} is monic, α is said to be an **algebraic integer**. The degree of an algebraic number is defined as the degree of p_{α} , and its height as the maximum absolute value of the coefficients of p_{α} (also said to be the height of that polynomial). The roots of p_{α} are said to be the **Galois conjugates** of α . We denote the set of algebraic numbers by A, and the set of algebraic integers by $\mathcal{O}_{\mathbb{A}}$. For all $\alpha \in \mathbb{A}$, there exists some $n \in \mathbb{N}$ such that $n\alpha \in \mathcal{O}_{\mathbb{A}}$. It is well known that \mathbb{A} is a field and that $\mathcal{O}_{\mathbb{A}}$ is a ring.

A number field of dimension d is a field extension K of \mathbb{Q} whose degree as a vector-space over \mathbb{Q} is d. In particular, $K \subseteq \mathbb{A}$ must hold. Recall that, in that case, there are exactly d monomorphisms $\sigma_i : K \to \mathbb{C}$ whose restriction over \mathbb{Q} is the identity (and therefore these must map elements of K to their Galois conjugates). The **ring of integers** \mathcal{O} of a number field K is the set of elements of K that are algebraic integers, that is, $\mathcal{O} = K \cap \mathcal{O}_{\mathbb{A}}$. An ideal of \mathcal{O} is an additive subgroup of \mathcal{O} that is closed under multiplication by any element of \mathcal{O} . An ideal \mathfrak{P} is said to be prime if $ab \in \mathfrak{P}$ implies $a \in \mathfrak{P}$ or $b \in \mathfrak{P}$. The following theorem is central in Algebraic Number Theory [19]:

THEOREM 7.1. In any ring of integers, ideals can be uniquely factored as products of prime ideals up to permutation.

7.3 Lower-bounding simple linear recurrence sequences We are interested in lower-bounding expressions of the form

(7.9)
$$u_n = \sum_{j=1}^s \alpha_j \lambda_j^n$$

where the α_j are algebraic-integer constants and $\lambda_1, \ldots, \lambda_s$ have the same absolute value ρ . Any such

sequence must in fact be a simple linear recurrence sequence with algebraic coefficients and characteristic roots $\lambda_1, \ldots, \lambda_s$, as explained in Section 1.1.6 of [15].

The next theorem, by Evertse, van der Poorten, and Schlickewei, was established in [16, 39] to analyse the growth of linear recurrence sequences. It gives us a very strong lower bound on the magnitude of sums of S-units, as defined below. Its key ingredient is Schlickewei's padic generalisation [35] of Schmidt's Subspace theorem.

Let S be a finite set of prime ideals of the ring of integers \mathcal{O} of a number field K. We say that $\alpha \in \mathcal{O}$ is an S-unit if all the ideals appearing in the prime factorisation of (α) , that is, the ideal generated by α , are in S.

THEOREM 7.2. (S-UNITS) Let K be a number field, s be a positive integer, and S be a finite set of prime ideals of \mathcal{O} . Then for every $\varepsilon > 0$ there exists a constant C, depending only on s, K, S, and ε , with the following property. For every set of S-units $x_1, \ldots, x_s \in \mathcal{O}$ such that $\sum_{i \in I} x_i \neq 0$ for all non-empty $I \subseteq \{1, \ldots, s\}$, it holds that

$$|x_1 + \dots + x_s| \ge CYZ^{-\varepsilon}$$

where $Y = \max\{|x_j| : 1 \le j \le s\}$ and $Z = \max\{\sigma_i(x_j) : 1 \le j \le s, 1 \le i \le d\}$ and σ_i represent the different monomorphisms from K to \mathbb{C} .

In order to make use of this result, it is important to understand the set

(7.10)
$$\{n \in \mathbb{N} : \exists I \subseteq \{1, \dots, s\}, \sum_{j \in I} \alpha_j \lambda_j^n = 0\}$$

The following well-known theorem characterises the set of zeros of linear recurrence sequences. In particular, it gives us a sufficient condition for guaranteeing that the set of zeros of a non-identically zero linear recurrence sequence is finite. Namely, it suffices that the sequence is non-degenerate, that is, that no ratio of two of its characteristic roots is a root of unit.

THEOREM 7.3. (SKOLEM-MAHLER-LECH) Let $u_n = \sum_{j=1}^{l} \alpha_j \lambda_j^n$ be a linear recurrence sequence. The set $\{n \in \mathbb{N} : u_n = 0\}$ is always a union of a finite set and finitely

many arithmetic progressions. Moreover, if u_n is nondegenerate, this set is actually finite.

Therefore, it follows from the Skolem-Mahler-Lech theorem that if u_n is non-degenerate then (7.10) must be finite, assuming without loss of generality that $\sum_{j \in I} \alpha_j \lambda_j^n$

is never eventually zero.

We can now apply the S-units theorem in order to get a lower bound on (7.9) that holds for all but finitely many n, by letting K be the splitting field of the characteristic polynomial of u_n , S be the set of prime ideals of the ring of integers of K that appear in the factorisation of each of the algebraic integers α_j and λ_j , and $x_j = \alpha_j \lambda_j^n$ for each j, making (7.9) a sum of S-units.

In the notation of the theorem, we have $Y = \Omega(\rho^n)$. If Λ is an upper bound on the absolute value of the Galois conjugates of each λ_j (that is, each $\sigma_i(\lambda_j)$), then $Z = O(\Lambda^n)$. Thus, for any $\varepsilon > 0$, we know that

$$\sum_{j=1}^{s} \alpha_j \lambda_j^n = \Omega(YZ^{-\varepsilon}) = \Omega\left(\rho^n \Lambda^{-n\varepsilon}\right)$$

Finally, we note that by picking ε to be sufficiently small we can get $\rho \Lambda^{-\varepsilon}$ arbitrarily close to ρ .

7.4 Manipulating algebraic numbers The following separation bound allows us to effectively represent an arbitrary algebraic number by keeping its defining polynomial, a sufficiently accurate estimate for the root we want to store, and an upper bound on the error. We call this its standard/canonical representation.

LEMMA 7.1. (MIGNOTTE) Let
$$f \in \mathbb{Z}[x]$$
. Then
(7.11)
$$f(\alpha_1) = 0 = f(\alpha_2) \Rightarrow |\alpha_1 - \alpha_2| > \frac{\sqrt{6}}{d^{(d+1)/2}H^{d-1}}$$

where d and H are respectively the degree and height of f.

It is well known that arithmetic operations and equality testing on these numbers can be done in polynomial time on the size of the canonical representations of the relevant numbers, since one can:

- compute polynomially many bits of the roots of any polynomial $p \in \mathbb{Q}[x]$ in polynomial time, due to the work of Pan in [29]
- find the minimal polynomial of an algebraic number by factoring the polynomial in its description in polynomial time using the LLL algorithm [23]
- use the sub-resultant algorithm (see Algorithm 3.3.7 in [12]) and the two aforementioned procedures to compute canonical representations of sums, differences, multiplications, and divisions of canonically represented algebraic numbers

Moreover, we need to know how to decide whether a given canonically represented algebraic number α is a root of unity, that is, whether $\alpha^r = 1$ for some r. If that is the case, then its defining polynomial will be the r-th cyclotomic polynomial, which has degree $\phi(r)$, if r is taken to be minimal, that is, if α is a primitive r-th root of unity. The following (crude) lower bound on $\phi(r)$ allows us to decide this in polynomial time, assuming that the degree of α is given in unary.

LEMMA 7.2. Let ϕ be Euler's totient function. Then $\phi(r) \geq \sqrt{(r/2)}$. Therefore, if α has degree n and is a r-th root of unity, then $r \leq 2n^2$.

Therefore, in order to decide whether an algebraic number α of degree n is a root of unity, we check whether it is a r-th root of unity, for each $r \leq 2n^2$. In order to test whether α is a r-th root of unity, it suffices to see whether $gcd(p_{\alpha}, x^r - 1) = p_{\alpha}$, since we know that $x^r - 1$ is the product of each d-th cyclotomic polynomial, with d ranging over the divisors of n.

7.5 First-Order Theory of Reals Let $\boldsymbol{x} = (x_1, \ldots, x_m)$ be a list of m real-valued variables, and let $\sigma(\boldsymbol{x})$ be a Boolean combination of atomic predicates of the form $g(\boldsymbol{x}) \sim 0$, where each $g(\boldsymbol{x})$ is a polynomial with integer coefficients in the variables \boldsymbol{x} , and \sim is either > or =. Tarski has famously shown that we can decide the truth over the field \mathbb{R} of sentences of the form $\phi = Q_1 x_1 \cdots Q_m x_n \sigma(\boldsymbol{x})$, where Q_i is either \exists or \forall . He did so by showing that this theory admits quantifier elimination (Tarski-Seidenberg theorem [37]).

All sets that are definable in the first-order theory of reals without quantification are by definition semialgebraic, and it follows from Tarski's theorem that this is still the case if we allow quantification. We also remark that our standard representation of algebraic numbers allows us to write them explicitly in the firstorder theory of reals, that is, given $\alpha \in \mathbb{A}$, there exists a sentence $\sigma(x)$ such that $\sigma(x)$ is true if and only if $x = \alpha$. Thus, we allow their use when defining semi-algebraic sets, for simplicity.

It follows from the undecidability of Hilbert's Tenth Problem that, in general, we cannot decide whether a given semi-algebraic set has an integer point.

We shall make use of the following result by Basu, Pollack, and Roy [3], which tells us how expensive it is, in terms of space usage, to perform quantifier elimination on a formula in the first-order theory of reals:

THEOREM 7.4. Given a set $\mathcal{Q} = \{q_1, \ldots, q_s\}$ of s polynomials each of degree at most D, in h + dd variables, and a first-order formula $\Phi(\mathbf{x}) = Qy_1 \ldots Qy_h F(q_1(\mathbf{x}, \mathbf{y}), \ldots, q_s(\mathbf{x}, \mathbf{y})),$ where $Q \in \{\exists, \forall\},$ F is a quantifier-free Boolean combination with atomic elements of the form $q_i(\mathbf{x}, \mathbf{y}) \sim 0$, then there exists a quantifier-free formula $\Psi(\mathbf{x}) = \bigwedge_{i=1}^J \bigvee_{j=1}^{J_i} q_{ij}(\mathbf{x}) \sim 0$, where $I \leq (sD)^{O(hd)}$, $J \leq (sD)^{O(d)}$, the degrees of **References** the polynomials q_{ij} are bounded by D^d , and the bitsizes of the heights of the polynomials in the quantifierfree formula are only polynomially larger than those of $q_1,\ldots,q_s.$

We also make use of the following lemmas:

LEMMA 7.3. If $X \subseteq \mathbb{R}^d$ is semi-algebraic and nonempty, $X \cap \mathbb{A}^d \neq \emptyset$.

Proof. We prove this result by strong induction on d. Since X is semi-algebraic, there exists a quantifier-free sentence in the first-order theory of reals σ such that $X = \{ x \in \mathbb{R}^d \mid \sigma(x) \}.$

Suppose that d > 1. Letting $X_1 = \{x_d \in \mathbb{R} \mid \exists x_1, \ldots, x_{d-1} \in \mathbb{R}^{d-1}, \sigma(x_1, \ldots, x_d)\}$ and since $X_1 \neq \emptyset$ is semi-algebraic, by the induction hypothesis, there must be $x_d^* \in \mathbb{A} \cap X_1$. Moreover, we can define $X_2 = \{(x_2, \ldots, x_d) \in \mathbb{R}^{d-1} \mid \sigma(x_1^*, x_2, \ldots, x_n)\}$, which is non-empty and semi-algebraic, and again by induction hypothesis there exists some $(x_2^*, \ldots, x_d^*) \in \mathbb{A}^{d-1} \cap X_2$.

It remains to prove this statement for d = 1. When d = 1, X must be a finite union of intervals and points, since semi-algebraic sets form an o-minimal structure on \mathbb{R} [37]. Clearly A is dense in any interval, and each of these isolated points x corresponds to some constraint q(x) = 0, which implies that x must be algebraic, since g has integer coefficients.

LEMMA 7.4. If $X \subseteq \mathbb{R}^d$ is semi-algebraic, then $X \cap \mathbb{A}^d$ is dense in X.

Proof. Pick $x \in X$ and $\varepsilon > 0$ arbitrarily. Let $y \in \mathbb{Q}^d$ be such that $||x - y|| < \varepsilon/2$. Since $B(y, \varepsilon/2)$ is semialgebraic, so must be $X \cap B(y, \varepsilon/2)$, and so this set must contain an algebraic point, since it is nonempty (x is in it), and that point must therefore be at distance at most ε of x, by the triangular inequality. By letting $\varepsilon \to 0$, we get a sequence of algebraic points which converges to x.

LEMMA 7.5. If $X \subseteq \mathbb{R}^d$ is semi-algebraic, so is \overline{X} .

Proof. Let σ be a sentence in the first-order theory of reals such that $X = \{x \in \mathbb{R}^d \mid \sigma(x)\}$. Whence

$$\overline{X} = \{ x \in \mathbb{R}^d \mid \forall \varepsilon > 0, \exists y \in \mathbb{R}^d, \sigma(y) \land y \in B(x, \varepsilon) \}.$$

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