FRIEZE VARIETIES : A CHARACTERIZATION OF THE
FINITE-TAME-WILD TRICHOTOMY FOR ACYCLIC QUIVERS

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Abstract. We introduce a new class of algebraic varieties which we call frieze varieties. Each frieze variety is determined by an acyclic quiver. The frieze variety is defined in an elementary recursive way by constructing a set of points in affine space. From a more conceptual viewpoint, the coordinates of these points are specializations of cluster variables in the cluster algebra associated to the quiver.

We give a new characterization of the finite–tame–wild trichotomy for acyclic quivers in terms of their frieze varieties. We show that an acyclic quiver is representation finite, tame, or wild, respectively, if and only if the dimension of its frieze variety is 0, 1, or $\geq 2$, respectively.

1. Introduction

For every acyclic quiver $Q$, we define an algebraic variety $X(Q)$ which we call the frieze variety of $Q$. The terminology stems from the fact that for quivers of Dynkin type $A$ the coordinates of the points of the frieze variety are entries in Conway-Coxeter friezes [10]. The frieze variety gives a geometric interpretation of the quiver as well as concrete numerical invariants, for example the dimension, the number of components and the degree.

The construction of the variety $X(Q)$ is inspired from the theory of cluster algebras. It is defined as follows. Let $Q = (Q_0, Q_1)$ be an acyclic quiver (i.e., a directed graph without oriented cycles) with $n$ vertices. Then we can label the vertices by integers $1, \ldots, n$ such that $i > j$ if there is an arrow $i \to j$.

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For every vertex \( i \in Q_0 \) we define positive rational numbers \( f_i(t) (t \in \mathbb{Z}_{\geq 0}) \) recursively by
\[
(1.1) \quad f_i(t + 1) = \frac{1 + \prod_{j \to i} f_j(t) \prod_{j \leftarrow i} f_j(t + 1)}{f_i(t)}.
\]

We will see in Lemma 2.1 below that these \( f_i(t) \) are exactly the specializations at \( x_1 = \cdots = x_n = 1 \) of preprojective cluster variables in the cluster algebra of \( Q \). In particular the \( f_i(t) \) are integers.

For every \( t \), we thus obtain a point \( P_t = (f_1(t), \ldots, f_n(t)) \in \mathbb{C}^n \) in an affine space. We define the *frieze variety* \( X(Q) \) of the quiver \( Q \) to be the Zariski closure of the set of all points \( P_t (t \in \mathbb{Z}_{\geq 0}) \). If we choose a different labelling, then the coordinates of each new \( P_t \) are obtained from the old by permuting in the same way for every \( t \). So the new \( X(Q) \) is obtained from the old by permuting its coordinates, thus is isomorphic to the old one. In particular, \( \dim X(Q) \) is independent of the labelling.

Thus every acyclic quiver \( Q \) comes with an algebraic variety \( X(Q) \). At this point many natural questions arise. Does the geometry of the variety reflect the representation theory of the quiver? Which quivers have smooth frieze varieties? Are dimension, degree and number of components meaningful invariants of the quiver? Moreover, although we focus in this paper on acyclic quivers, one can easily generalize the definition of a frieze variety to quivers that are not acyclic themselves but that are mutation equivalent to an acyclic quiver. One may then ask how does the frieze variety behave under mutation?

In this paper, we show that the dimension of the frieze variety detects the representation type of the quiver. An acyclic quiver \( Q \) is either *representation finite*, *tame* or *wild*, depending on the representation theory of its path algebra. The quiver is representation finite if and only if its underlying graph is a Dynkin diagram of type \( A, D, \) or \( E \) \cite{17}, and it is tame if and only if the underlying graph is an affine Dynkin diagram of type \( \widetilde{A}, \widetilde{D}, \) or \( \widetilde{E} \) \cite{12, 26, 11}. All other acyclic quivers are wild.

We propose a new characterization of the finite–tame–wild trichotomy in terms of the frieze variety \( X(Q) \) of the quiver \( Q \).

**Theorem 1.1.** Let \( Q \) be an acyclic quiver.

(a) If \( Q \) is representation finite then the frieze variety \( X(Q) \) is of dimension 0.

(b) If \( Q \) is tame then the frieze variety \( X(Q) \) is of dimension 1.

(c) If \( Q \) is wild then the frieze variety \( X(Q) \) is of dimension at least 2.

If \( Q \) is representation finite then the cluster algebra has only finitely many cluster variables \cite{15} and hence \( X(Q) \) is a finite set of points. This shows part (a) of Theorem 1.1.

To prove part (b), we will specify linear recursions for the coordinates of the points \( P_t \) and then use a general argument to show that the projection of \( X(Q) \) to any coordinate plane is contained in the zero locus of a polynomial constructed from the linear recurrence. The key step here is to show that all the roots of the characteristic polynomials of all recursions are integral powers of a single complex number. Linear recursions for the sequences \( (f_i(t))_{t \geq 0} \) where already considered in \cite{11, 19}, where it is shown that there exists a linear recursion for \( (f_i(t))_{t \geq 0} \) for all \( i \) if and only if \( Q \) is representation finite or tame. In \cite{19}, explicit linear recursions were given in type \( \widetilde{D} \) for leaf vertices, and we give new proofs for these recursions here. For type \( \widetilde{A} \) as well as the non-leaf vertices in type \( \widetilde{D} \), we provide new explicit recursions.
To prove part (c) of the theorem, we use the fact that the points $P_t$ correspond to slices $\tau^{-t+1}kQ$ in the preprojective component of the Auslander-Reiten quiver of the path algebra $kQ$ of $Q$, as well as several known facts on the spectral theory of the Coxeter matrix of a wild quiver, see [27]. The key result, which we think interesting in its own right, is to show that, when $t$ goes to infinity, the natural logarithm of the coordinates $\ln f_i(t)$ grows in the same way as $\rho^t$, where $\rho$ is the largest eigenvalue, or spectral radius, of the Coxeter matrix. See Proposition 4.7.

There are several characterizations of the finite-tame-wild trichotomy. In [28], Ringel showed that $Q$ is wild if and only if the spectral radius of the Coxeter transformation is greater than 1. In [30], Skowroński and Weyman characterized tameness in terms of semi-invariants. Recently, Lorscheid and Weist characterized tameness using quiver Grassmannians [22]. To our knowledge, our characterization is the first one in terms of numerical invariants that are integers.

The paper is organized as follows. In Section 2, we recall several definitions and results from representation theory and cluster algebras that are needed later. We prove part (b) of Theorem 1.1 in Section 3 and part (c) in Section 4. We give several examples in section 5.

2. Preliminaries

Throughout the paper we work over the field of complex numbers $\mathbb{C}$.

2.1. Quivers and representations. We start by recalling a few basic facts about quivers and their representations. For further details we refer to [2] [29].

A quiver $Q = (Q_0, Q_1, s, t)$ consists of a set $Q_0$ of vertices, a set $Q_1$ of arrows, and two maps $s, t : Q_1 \to Q_0$ that send an arrow $\alpha$ to its starting point $s(\alpha)$ and terminal point $t(\alpha)$. We call $Q$ a finite quiver if $Q_0$ and $Q_1$ are both finite sets. We will always assume $Q$ to be finite in our paper.

A representation $M = (M_i, \varphi_\alpha)_{i \in Q_0, \alpha \in Q_1}$ of a quiver $Q$ is a collection of $\mathbb{C}$-vector spaces $M_i$ ($i \in Q_0$) together with a collection of $\mathbb{C}$-linear maps $\varphi_\alpha : M_{s(\alpha)} \to M_{t(\alpha)}$ ($\alpha \in Q_1$). A representation $M$ is called finite-dimensional if each $M_i$ is finite-dimensional. The representations considered in this paper are all finite-dimensional. Let $\dim M = (\dim M_i)_{i \in Q_0}$ be the dimension vector of $M$, and let $\text{rep}_{\mathbb{C}}Q$ denote the category of finite-dimensional representations of $Q$. Let $\mathbb{C}Q$ be the path algebra of the quiver $Q$ over $\mathbb{C}$ and mod $\mathbb{C}Q$ the category of finitely generated $\mathbb{C}Q$-modules. There is an equivalence of categories $\text{mod} \mathbb{C}Q \cong \text{rep}_{\mathbb{C}}(Q)$, and we use the notions of representations and modules interchangeably.

The projective representation $P(i)$ at vertex $i \in Q_0$ is defined as $(P(i)_j, \varphi_\alpha)$ where $P(i)_j$ is the vector space with basis the set of all paths from $i$ to $j$ in $Q$; for an arrow $j \xrightarrow{\alpha} \ell$ in $Q$, the map $\varphi_\alpha : P(i)_j \to P(i)_\ell$ is determined by composing the paths from $i$ to $j$ with the arrow $j \xrightarrow{\alpha} \ell$.

Let $D$ be the duality functor $\text{Hom}_\mathbb{C}(-, \mathbb{C})$, and let $A = \bigoplus_{j \in Q_0} P(j)$. The Nakayama functor is defined as $\nu = D\text{Hom}_A(-, A)$. Let $\tau, \tau^{-1}$ denote the Auslander-Reiten translations. Recall the definition of $\tau$. Let $M$ be an indecomposable, non-projective representation of an acyclic quiver $Q$, and

$$
P_1 \xrightarrow{f} P_0 \longrightarrow M \longrightarrow 0$$
be a minimal projective presentation. Then \( \tau M \) is defined by the following exact sequence

\[
0 \to \tau M \to \nu P_1 \overset{\nu f}{\to} \nu P_0 \to \nu M \to 0.
\]

The inverse Auslander-Reiten translation \( \tau^{-1} \) is defined dually for non-injective, indecomposable representations of \( Q \). For every indecomposable non-injective representation \( M \) there is a unique almost split sequence \( 0 \to M \to E \to \tau^{-1}M \to 0 \) starting at \( M \).

An indecomposable representation \( N \) is called \textit{preprojective} if there is a nonnegative integer \( t \) such that \( \tau^tN = P(i) \) for some \( i \in Q_0 \). The set of all indecomposable preprojective representations form the \textit{preprojective component} of the Auslander-Reiten quiver of \( Q \).

2.1.1. \textit{Admissible sequences}. A sequence of vertices \( (i_1, \ldots, i_n) \) (\( i_j \neq i_\ell \) if \( j \neq \ell \)) is called an \textit{admissible sequence} if the following conditions hold:

1. \( i_1 \) is a sink of \( Q \);
2. \( i_2 \) is a sink of the quiver \( s_{i_1}Q \) obtained from \( Q \) by reversing all arrows that are incident to the vertex \( i_1 \);
3. \( i_\ell \) is a sink of \( s_{i_{\ell-1}} \cdots s_{i_1}Q \) for \( t = 2, 3, \ldots, n \).

Note that the above definition is equivalent to saying that \( j < \ell \) if there is an arrow \( i_j \leftarrow i_\ell \) in \( Q \). Indeed, assuming \( (i_1, \ldots, i_n) \) is admissible, if there is an arrow \( \alpha : i_j \leftarrow i_\ell \) in \( Q \), then the sequence \( s_{i_{\ell-1}} \cdots s_{i_1} \) changes the orientation of \( \alpha \) if and only if exactly one of \( i_j \) and \( i_\ell \) is in \( \{i_1, \ldots, i_{\ell-1}\} \), or equivalently, \( j < \ell \) (because \( i_\ell \notin \{i_1, \ldots, i_{\ell-1}\} \)). Conversely, assume \( j < \ell \) if there is an arrow \( i_j \leftarrow i_\ell \) in \( Q \). Then any arrow of the form \( i_j \leftarrow i_\ell \) (thus \( j < \ell \)), changes its orientation under the sequence \( s_{i_{\ell-1}} \cdots s_{i_1} \), so we get a new arrow \( i_j \rightarrow i_\ell \). On the other hand, any arrow of the form \( i_\ell \leftarrow i_j \) (thus \( \ell < j \)), remains unchanged under the sequence \( s_{i_{\ell-1}} \cdots s_{i_1} \). Thus \( i_\ell \) becomes a sink in the quiver \( s_{i_{\ell-1}} \cdots s_{i_1}Q \).

It is easy to see that if \( (i_1, \ldots, i_n) \) is an admissible sequence then \( s_{i_n} \cdots s_{i_1}Q = Q \). Indeed, since the admissible sequence contains each vertex exactly once, the reflection sequence reflects each arrow exactly twice.

Since we chose our vertex labels \( 1, \ldots, n \) such that \( i > j \) if there is an arrow \( i \rightarrow j \), we see that the sequence \( 1, \ldots, n \) is an admissible sequence.

2.1.2. \textit{Structure of the preprojective component of the Auslander-Reiten quiver}. The preprojective component of the Auslander-Reiten quiver has vertices \( \tau^{-t}P(i) \) with \( i \in Q_0, t \geq 0 \) and arrows \( \alpha_t : \tau^{-t}P(i) \to \tau^{-t}P(j) \) and \( \overline{\alpha}_t : \tau^{-t}P(j) \to \tau^{-t-1}P(i) \), for all \( \alpha : j \to i \in Q_1, t \geq 0 \).

For example, if \( Q \) is the quiver \( 1 \overset{\alpha}{\leftarrow} 2 \overset{\gamma}{\rightarrow} 3 \), then the beginning of the preprojective component is of the form

\[
\begin{array}{c}
P(1) \overset{\alpha_0}{\to} P(2) \overset{\tau^{-1}}{\to} P(1) \overset{\tau^{-1}}{\to} P(2) \overset{\tau^{-1}}{\to} P(1) \overset{\tau^{-1}}{\to} P(2) \overset{\tau^{-1}}{\to} P(1) \overset{\tau^{-1}}{\to} P(2) \overset{\tau^{-1}}{\to} P(1) \overset{\tau^{-1}}{\to} \cdots
\end{array}
\]
Each mesh of the preprojective component represents an almost split short exact sequence of the following form

\[ 0 \longrightarrow \tau^{-t+1}P(i) \longrightarrow \bigoplus_{j \rightarrow i} \tau^{-t+1}P(j) \oplus \bigoplus_{i \rightarrow j} \tau^{-t}P(j) \longrightarrow \tau^{-t}P(i) \longrightarrow 0 \]

2.2. **Cluster algebras.** Let \( Q \) be an acyclic quiver with \( n \) vertices. The cluster algebra \( A(Q) \) of the quiver \( Q \) is the \( \mathbb{Z} \)-subalgebra of the field of rational functions \( \mathbb{Q}(x_1, \ldots, x_n) \) generated by the set of all cluster variables obtained by mutation from the initial seed \( ((x_1, \ldots, x_n), Q) \).

For every vertex \( i \), the mutation \( \mu_i \) in direction \( i \) transforms a seed \( ((x_1, \ldots, x_n), Q) \) by replacing the \( i \)-th cluster variable \( x_i \) by the new cluster variable \( (\prod_{j \rightarrow i} x_j + \prod_{i \rightarrow j} x_j)/x_i \), where the first product runs over all arrows in \( Q \) that end at \( i \) and the second product over all arrows that start at \( i \). Moreover, the mutation also changes the quiver. We refer to [16] for further details on cluster algebras.

In this paper, we are only concerned with mutations at sinks. Recall that a vertex \( i \) is a sink if there is no arrow starting at \( i \). Thus in this case, the mutation formula becomes \( (\prod_{j \rightarrow i} x_j + 1)/x_i \). Moreover on the level of the quiver, the mutation \( \mu_i \) of \( Q \) at a sink \( i \) is the same as the reflection \( s_iQ \) of \( Q \).

Let \( i_1, \ldots, i_n \) be an admissible sequence for \( Q \), and denote the corresponding mutation sequence \( \mu = \mu_{i_n} \cdots \mu_{i_1} \). Then each mutation in this sequence is a mutation at a sink, and moreover \( \mu Q = Q \). Let \( x_0 = (x_1(0), \ldots, x_n(0)) \) denote the initial cluster and \( x_t = (x_1(t), \ldots, x_n(t)) = \mu^t(x_0) \) be the cluster obtained from it by applying the sequence \( \mu \) exactly \( t \) times, where \( x_j(t) \) is the unique cluster variable that appears for the first time after the mutations \( \mu_{i_1} \cdots \mu_{i_t} \mu^{-1} \).

A representation \( M \) is called **rigid** if \( \text{Ext}^1(M, M) = 0 \). For example, indecomposable preprojective representations are rigid, since \( \text{Ext}^1(\tau^{-t}P(i), \tau^{-t}P(i)) \cong \text{Ext}^1(P(i), P(i)) = 0 \).

The **cluster character**, or **Caldero-Chapoton map**, associates a cluster variable \( X_M \) to every indecomposable, rigid representation \( M \) of \( Q \) in such a way that the denominator of Laurent polynomial \( X_M \) is equal to \( \prod_{i \in Q_0} x_i^{d_i} \), where \( (d_1, \ldots, d_n) \) is the dimension vector of \( M \). This was shown in [5] for Dynkin quivers and in [6] for arbitrary acyclic quivers. This result applies in particular to all indecomposable preprojective representations \( \tau^{-t}P(i) \) with \( i \in Q_0, t \geq 0 \).

It was also shown in [4, 6] that if \( M \) is a rigid indecomposable representation with almost split sequence \( 0 \rightarrow M \rightarrow E \rightarrow \tau^{-1}M \rightarrow 0 \), then in the cluster algebra we have the exchange relation \( X_{\tau^{-1}M}X_M = X_E + 1 \). Using our description of the almost split sequences in the preprojective component in section 2.1.2 we see that

\[ X_{\tau^{-t}P(i)} = \left( \prod_{j \rightarrow i} X_{\tau^{-t+1}P(j)} \prod_{i \rightarrow j} X_{\tau^{-t}P(j)} + 1 \right)/X_{\tau^{-t+1}P(i)}, \]

and with our notation above this becomes

\[ x_i(t+1) = \left( \prod_{j \rightarrow i} x_j(t) \prod_{i \rightarrow j} x_j(t+1) + 1 \right)/x_i(t). \]

If we now specialize the initial cluster variables at 1, we obtain precisely the recursive definition of the coordinates \( f_i(t) \). We summarize the above results in the following Lemma.

**Lemma 2.1.** (1) \( f_i(t) \) is \( x_i(t) \) specialized at \( x_1 = \cdots = x_n = 1 \). In particular, \( f_i(t) \) is a positive integer.
(2) \[ x_i(t) = X_{\tau - t + 1}P(i). \]

(3) The denominator of \( X_{\tau - t + 1}P(i) \) is equal to \( \prod_{i=1}^{n} x_i^{d_i} \), where \( (d_1, \ldots, d_n) \) is the dimension vector of \( \tau - t + 1P(i) \).

2.3. Surface type. A special class of quivers are those associated to triangulations of surfaces with marked points. The cluster algebras of these quivers are said to be of surface type. The cluster algebra (with trivial coefficients) does not depend on the choice of triangulation of the surface. It was shown in [13] that there are precisely four types of surfaces that give rise to acyclic quivers.

(1) The disk with \( n + 3 \) marked points on the boundary corresponds to the finite type \( \mathbb{A}_n \). The quiver is acyclic if and only if the triangulation has no internal triangles.

(2) The disk with one puncture and \( n \) marked points on the boundary corresponds to the finite type \( \mathbb{D}_n \). The quiver is acyclic if and only if the triangulation has no internal triangles and exactly two arcs incident to the puncture.

(3) The annulus with \( p \) marked points on one and \( q \) marked points on the other boundary component corresponds to the affine type \( \tilde{A}_{p,q} \) with \( n = p + q \) vertices. The quiver is acyclic if and only if every arc in the triangulation connects two points on different boundary components.

(4) The disk with two punctures and \( n - 3 \) marked points on the boundary corresponds to the affine type \( \tilde{D} \) with \( n \) vertices (in the usual notation this would be type \( \tilde{D}_{n-1} \)). The quiver is acyclic if and only if

(i) for each of the two punctures \( p_i \) there are precisely two arcs \( \tau_{i1} \) and \( \tau_{i2} \) \( (i = 1, 2) \) that connect \( p_i \) to a boundary point \( a_{i1}, a_{i2} \), such that, either \( a_{i1} = a_{i2} \) or \( a_{i1} \) and \( a_{i2} \) are neighbors on the boundary. Therefore, either the arcs \( \tau_{i1}, \tau_{i2} \) form a selffolded triangle or they form a triangle together with the boundary segment \( a_{i1} - a_{i2} \).

(ii) letting \( B_1 \) and \( B_2 \) be the two parts of the boundary separated by the two triangles incident to the punctures, each of the remaining \( n - 4 \) arcs must connect a point of \( B_1 \) to a point of \( B_2 \).

It was also shown in [13] that there is a bijection between cluster variables and tagged arcs in the surface. Later, in [23], combinatorial formulas were given for cluster variables, and in [24] these formulas were used to associate elements of the cluster algebra to other curves in the surface including closed simple loops and bracelets. If \( L \) is a closed simple curve its \( k \)-bracelet \( \text{Brac}_k(L) \) is the \( k \)-fold concatenation of \( L \) with itself. Thus the 1-bracelet is just the loop \( L \) and the \( k \)-bracelet has \( k - 1 \) selfcrossings. These bracelets are essential in the construction of the canonical basis known as the bracelet basis in [24]. Bracelets satisfy the following Chebyshev recursion \( \text{Brac}_0(L) = 2, \text{Brac}_1(L) = L \) and

\[ \text{Brac}_k(L) = L \cdot \text{Brac}_{k-1}(L) - \text{Brac}_{k-2}(L). \]

All these elements satisfy the so-called skein relations, which are given on the level of curves by smoothing a crossing \( \times \) in two ways \( \simeq \) and \( \supset \subset \). The skein relations in the cluster algebra were proved in [25] using hyperbolic geometry and in [7, 8, 9] using only the combinatorial definition of the cluster algebra elements. The skein relations between bracelets and arcs play a crucial role in the proof of our main theorem in the affine types \( \tilde{A} \) and \( \tilde{D} \).

2.4. Linear recurrences. We recall the following result about linear recurrences.
Lemma 2.2. Let \((a_n)\) be a sequence given by the recurrence \(a_n = c_1a_{n-1} + c_2a_{n-2} + \ldots + c_da_{n-d}\), where the \(c_i \in \mathbb{C}\) are constant. Let \(p(x)\) be the characteristic polynomial of this recurrence, thus \(p(x) = x^d - c_1x^{d-1} - c_2x^{d-2} - \ldots - c_d\), and denote by \(r_1, r_2, \ldots, r_d\) the roots of \(p(t)\). If the roots of \(p(t)\) are all distinct then there exist constants \(\alpha_i \in \mathbb{C}\) such that \(a_n = \alpha_1r_1^n + \alpha_2r_2^n + \cdots + \alpha_dr_d^n\).

In particular, if there exists a complex number \(\rho\) such that the roots are \(\rho^i\), for \(d\) distinct integers \(i\), then \(a_n = \sum_i \alpha_i\rho^{ni}\).

3. Proof of the main theorem part (b), the tame case

In this section, we prove that \(\dim X(Q) = 1\) for affine types. In each case we shall exhibit a linear recursion for the coordinates of the points of \(X(Q)\), which then will imply the desired result via the results in section 2.4. In the types \(\tilde{A}\) and \(\tilde{D}\), we use skein relations to obtain the recursions, and in type \(\tilde{E}\) we checked the recursion formulas by computer.

We start with two preparatory lemmas.

Lemma 3.1. Let \(e\) be a positive integer and \(\rho \in \mathbb{C}\). Let \((a_n)_{n \in \mathbb{Z}}\) and \((b_n)_{n \in \mathbb{Z}}\) be two sequences such that

\[
a_n = \sum_{i=-e}^{e} \alpha_i\rho^{ni} \quad \text{and} \quad b_n = \sum_{i=-e}^{e} \beta_i\rho^{ni}
\]

with \(\alpha_i, \beta_i \in \mathbb{C}\). Then there exists a nonzero polynomial \(g(x,y) \in \mathbb{C}[x,y]\) of degree at most \(4e - 2\) such that \(g(a_n, b_n) = 0\) for all \(n\).

Proof. Let \(d = 4e - 2\) and consider the general polynomial of degree \(d\)

\[
g(x,y) = \sum_{0 \leq j+k \leq d} c_{j,k}x^jy^k.
\]

This polynomial has \((d+2)(d+1)/2 = 8e^2 - 2e\) coefficients \(c_{j,k}\). We want to find coefficients \(c_{j,k}\) such that \(g(a_n, b_n) = 0\) for all integers \(n\), which is

\[
(3.1) \sum_{0 \leq j+k \leq d} c_{j,k}(\sum_{i=-e}^{e} \alpha_i\rho^{ni})^j(\sum_{i=-e}^{e} \beta_i\rho^{ni})^k = 0.
\]

If we write the left hand side as a polynomial in \(\rho^{\pm n}\) it is of the form

\[
\sum_{\ell=-ed}^{ed} \gamma_\ell(\rho^n)\ell,
\]

where \(\gamma_\ell\) is a linear function of the \(c_{j,k}, 0 \leq j+k \leq d\). Thus in order to show \((3.1)\) it suffices to show that the system of equations

\[
\gamma_\ell((c_{j,k})) = 0
\]

has a nontrivial solution, and this is true since the number of equations \(2ed+1 = 8e^2 - 4e + 1\) is strictly smaller than the number of variables \(8e^2 - 2e\). \(\square\)

Lemma 3.2. Let \(K\) be any field. Fix a nonnegative integer \(e\) and any element \(\rho \in K\). Let \(n\) be a positive integer, and for each \(p \in \{1, \ldots, n\}\), let \(\{a_j^{(p)}\}_{j \in \mathbb{Z}}\) be a sequence such that

\[
a_j^{(p)} = \sum_{i=-e}^{e} \alpha_i^{(p)}\rho^{ji},
\]
where $\alpha_i^{(p)} \in K$ is independent of $j$. Then the Zariski closure of \( \{(a_j^{(1)}, a_j^{(2)}, \cdots, a_j^{(n)}) \in \mathbb{A}_K^n : j \in \mathbb{Z}\} \) is of Krull dimension $\leq 1$.

**Proof.** We use induction on $n$. There is nothing to show for $n = 1$. The base case of $n = 2$ is proved in Lemma 3.1. Suppose that it holds for $n$, and we prove for $n + 1$. Let

\[
C := \left\{(a_j^{(1)}, a_j^{(2)}, \cdots, a_j^{(n)}, 0) \in \mathbb{A}_K^{n+1} : j \in \mathbb{Z}\right\},
\]

\[
D := \left\{(a_j^{(1)}, a_j^{(2)}, \cdots, a_j^{(n-1)}, 0, a_j^{(n+1)}) \in \mathbb{A}_K^{n+1} : j \in \mathbb{Z}\right\},
\]

\[
E := \left\{(a_j^{(1)}, \cdots, a_j^{(n-2)}, 0, a_j^{(n)}, a_j^{(n+1)}) \in \mathbb{A}_K^{n+1} : j \in \mathbb{Z}\right\},
\]

\[
Z := \left\{(a_j^{(1)}, a_j^{(2)}, \cdots, a_j^{(n-1)}, a_j^{(n)}, a_j^{(n+1)}) \in \mathbb{A}_K^{n+1} : j \in \mathbb{Z}\right\},
\]

where the bar denotes the Zariski closure. By induction each of $\text{dim}(C)$, $\text{dim}(D)$, and $\text{dim}(E)$ are $\leq 1$. If one of them is equal to 0, then $\text{dim}(Z) \leq 1$ since $Z \subset C \times \mathbb{A}^1$, $Z \subset D \times \mathbb{A}^1$, and $Z \subset E \times \mathbb{A}^1$. Suppose that $\text{dim}(C) = \text{dim}(D) = \text{dim}(E) = 1$. Aiming at contradiction, assume that $\text{dim}(Z) = 2$. Let $Z_1$ be an irreducible component of $Z$ with $\text{dim}(Z_1) = 2$. Then $Z_1 = C_1 \times \mathbb{A}^1 = D_1 \times \mathbb{A}^1 = E_1 \times \mathbb{A}^1$ for some irreducible component $C_1$ of $C$, some irreducible component $D_1$ of $D$, and some irreducible component $E_1$ of $E$. This implies that all of $C_1, D_1$, and $E_1$ are lines. Hence $Z_1$ is a linear plane in $\mathbb{A}_K^{n+1}$. Choose three non-collinear points in general position, say $(q_{w,1}, \ldots, q_{w,n+1})_{w \in \{1,2,3\}}$, on $Z_1$. Then the points $(q_{w,1}, \ldots, q_{w,n}, 0)_{w \in \{1,2,3\}}$ are distinct and collinear, because they are on the line $C_1$. Similarly $(q_{w,1}, \ldots, q_{w,n-1}, 0, q_{w,n+1})_{w \in \{1,2,3\}}$ are distinct and collinear, and $(q_{w,1}, \ldots, 0, q_{w,n}, q_{w,n+1})_{w \in \{1,2,3\}}$ are collinear as well. Then $(q_{w,1}, \ldots, q_{w,n+1})_{w \in \{1,2,3\}}$ become collinear, which is a contradiction. \qed

3.1. **Affine type A.** Let $Q$ be an acyclic quiver of type $\tilde{A}_{p,q}$. We will use the annulus with $p$ marked points on the inner boundary component and $q$ marked points on the outer boundary component as a model for mod $\mathbb{C}Q$ as described in section 2.3. Let

\[
k = \frac{p + q}{\gcd(p, q)} \quad \text{and} \quad m = \text{lcm}(p, q).
\]

Thus $k = (p + q)m/pq$. Let $L$ be the (isotopy class of the) closed simple curve formed by the equator of the annulus, and consider its $k$-bracelet $\text{Brac}_k(L)$. The crossing number $e(\gamma, L)$ between any two isoclasses of curves is defined to be the minimum number of crossings between a curve in the isocopy class of $\gamma$ and the isotopy class of $L$. We define the constant

\[
C(p, q) = X_{\text{Brac}_k(L)}|_{x_1 = 1}
\]

to be the positive integer obtained from the Laurent polynomial $X_{\text{Brac}_k(L)}$ of the bracelet by specializing the initial cluster variables $x_1 = \cdots = x_n = 1$. Note that, unless one of $p$ or $q$ is 1, the value of $C(p, q)$ depends on the orientation of the arrows of $Q$. For $q = 1$ we have $C(p, 1) = T_{p+1}(p + 2)$, where $T_p$ is the $p$-th Chebyshev polynomial with $T_0 = 2$. So $C(1, 1) = 7, C(2, 1) = 52, C(3, 1) = 527, C(4, 1) = 6726$. For $p = q = 2$, there are two possible values, $C(2, 2) = 34$ or 47.

We have the following linear recursion for the coordinates $f_i(t)$ of the points defining $X(Q)$.

**Theorem 3.3.** Let $Q$ be of type $\tilde{A}_{p,q}$ and $m = \text{lcm}(p, q)$. Then for all $i \in Q_0$ and all $t \geq m$

\[
f_i(t + m) = C(p, q)f_i(t) - f_i(t - m).
\]
Figure 1. The skein relation of equation (3.2) for $k = 2$.

**Proof.** The indecomposable representations in the preprojective component correspond to arcs that connect points on different boundary components. For each such arc $\gamma$, the crossing number $e(\gamma, L)$ with $L$ is 1, and the crossing number $e(\gamma, \text{Brac}_{k}(L))$ with the $k$-bracelet is $k$. Smoothing one of these crossings we obtain the following skein relation

\[(3.2) \quad \text{Brac}_{k}(L) \cdot \gamma = D_{k}(\gamma) + D_{-k}(\gamma),\]

where $D$ denotes the Dehn twist along $L$. We give an example for $k = 2$ in Figure 1. On the other hand, the inverse Auslander-Reiten translation $\tau^{-1}$ acts on the arc $\gamma$ by moving its endpoint on the outer boundary component to its counterclockwise neighbor and its endpoint on the inner boundary component to its clockwise neighbor. Since $m = \text{lcm}(p, q)$, we see that the arc $\tau^{-m}(\gamma)$ has the same endpoints as $\gamma$ but $\tau^{-m}(\gamma)$ wraps around the inner boundary component exactly \(m p + m q = m pq\) times more than $\gamma$. In other words, $\tau^{-m}(\gamma) = D_{k}(\gamma)$. Therefore, equation (3.2) becomes

\[(3.2) \quad \text{Brac}_{k}(L) \cdot \gamma = \tau^{m}(\gamma) + \tau^{-m}(\gamma),\]

and passing to the cluster algebra and specializing at $x_i = 1$, we have

\[C(p, q) f_i(t) = f_i(t - m) + f_i(t + m). \quad \Box\]

**Remark 3.4.** Similar looking recurrence relations involving bracelets but for arcs that have both endpoints on the same boundary component were found in [3, Theorem 2.5] and [18, Theorem 5.4].

**Corollary 3.5.** Let $Q$ be an acyclic quiver of affine type $\tilde{\mathbb{A}}$. Then the dimension of the frieze variety $X(Q)$ is equal to one.

**Proof.** It suffices to show that, for every pair $i, j \in Q_0$, the projection $\pi_{ij}(X(Q))$ of $X(Q)$ onto the $(i, j)$-plane is of dimension one. By Theorem 3.3 and Lemma 2.2, we have linear recurrences

\[f_h(t + m) = C(p, q) f_h(t) - f_h(t - m) = \sum_{e=-1}^{1} \alpha_{h, r, e} \rho^{e \lfloor t/m \rfloor},\]

where $\rho$ is one of the roots of the polynomial $x^2 - C(p, q)x + 1$, and $\alpha_{h, r, e}$ depends only on $h \in \{i, j\}$, $r = t - m \lfloor t/m \rfloor \in \{0, 1, ..., m - 1\} \mod m$, and $e$. Since $C(p, q) > 2$, we have $\rho \neq 1$.

Thus Lemma 3.1 with $e = 1$ implies there is a polynomial $g_h(x, y)$ of degree $4e - 2 = 2$ such that $g_h(f_i(t + sm), f_j(t + sm)) = 0$, for all $s$ such that $t + sm \geq 0$. Define $g = g_0 \cdot g_1 \cdots g_{m-1}$. Then $\pi_{ij}(X(Q))$ is contained in the zero locus of $g$, and hence has dimension one. \quad \Box
3.2. **Affine type** $D$. Let $Q$ be an acyclic quiver of type $\tilde{\mathbb{D}}$ with $n$ vertices. Thus the underlying graph of $Q$ is the following.

\[
\bigcirc \rightarrow \bullet \rightarrow \cdots \rightarrow \bullet \rightarrow \bigcirc
\]

Since that we use $n$ for the number of vertices, so in the usual notation this type is $\tilde{\mathbb{D}}_{n-1}$. Each of the vertices marked with the symbol $\bigcirc$ is called a *leaf* of $Q$ and each of the vertices marked with the symbol $\bullet$ is called a *non-leaf*. We will use the disk with two punctures and $n-3$ marked points on the boundary as a model for $\mathbb{C}Q$ as described in section 2.3. Let $L$ be the (isotopy class of the) closed simple curve around the two punctures and let $\text{Brac}_k(L)$ denote its $k$-bracelet. We denote by $D$ the (isotopy class of the) closed simple curve around the two punctures and let $\text{Brac}_k(L)$ denote its $k$-bracelet. We denote by $D$ the full Dehn twist along $L$ in counterclockwise direction, and by $D^{1/2}$ the half Dehn twist.

The indecomposable representations in the preprojective component correspond to two types of arcs, depending whether the vertex $i$ is a leaf of $Q$ or not.

- If $i$ is a leaf in $Q$, then $\tau^{-t}P(i)$ corresponds to an arc $\gamma$ that connects a boundary point to a puncture such that the crossing number $e(\gamma, L)$ is one.
- If $i$ is not a leaf in $Q$, then $\tau^{-t}P(i)$ corresponds to an arc $\gamma$ with both endpoints on the boundary and such that the crossing number $e(\gamma, L)$ is two.

We call an arc $\gamma$ *preprojective of orbit* $i$ if it corresponds to a preprojective representation of the form $\tau^{-t}P(i)$.

We have the following relation between the Dehn twist and the Auslander-Reiten translation.

**Lemma 3.6.** Let $\gamma$ be a preprojective arc of orbit $i$. Then

\[
\mathcal{D}^2(\gamma) = \tau^{-2(n-3)}(\gamma).
\]

Moreover, if $i$ is a non-leaf vertex or $n$ is odd, then also $\mathcal{D}(\gamma) = \tau^{-(n-3)}(\gamma)$.

**Proof.** Recall that $n - 3$ is the number of marked points on the boundary of the disk. If $i$ is a leaf then applying $\tau^{-1}$ moves the endpoint of $\gamma$ that lies on the boundary to its counterclockwise neighbor and changes the tagging at the puncture. Thus $\tau^{-(n-3)}(\gamma)$ is equal to the full Dehn twist $\mathcal{D}(\gamma)$ if $n - 3$ is even, and it is the Dehn twist with opposite tagging at the puncture if $n - 3$ is odd.

If $i$ is not a leaf, then applying $\tau^{-1}$ moves both endpoints of $\gamma$ to their counterclockwise neighbors on the boundary. Thus after applying $\tau^{-1}$ exactly $n - 3$ times, we have moved each endpoint of $\gamma$ counterclockwise around the whole boundary back to its initial position. On the other hand, the Dehn twist does exactly the same, since $\gamma$ crosses the loop $L$ twice in this situation.

If $\gamma$ is a preprojective arc of orbit $i$ with $i$ a non-leaf vertex, we let $\gamma_1, \gamma_2$ be the two arcs that have the same endpoints as $\gamma$ and do not cross the loop $L$, see Figure 2. Define $S_k(L)$ recursively by $S_1(L) = 1$, $S_2(L) = L + 2$ and $S_k(L) = LS_{k-1}(L) - S_{k-2}(L) + 2$.

**Lemma 3.7.** Let $\gamma$ be a preprojective arc of orbit $i$ and $k \geq 1$.

(a) If the vertex $i$ is a non-leaf then

\[
\text{Brac}_k(L) \cdot \gamma = \mathcal{D}^{k/2}(\gamma) + \mathcal{D}^{-k/2}(\gamma) + 2(\gamma_1 + \gamma_2)S_k(L).
\]
Figure 2. The skein relation of equation (3.3). In the cluster algebra the loops that are contractible to the puncture are equal to 2.

(b) If the vertex $i$ is a leaf then

$$\operatorname{Brac}_k(L) \cdot \gamma = D^k(\gamma) + D^{-k}(\gamma).$$

Proof. For $k = 1$ the results are the following skein relations illustrated in Figure 2

$$L \cdot \gamma = \begin{cases} D^{1/2}(\gamma) + D^{-1/2}(\gamma) + 2(\gamma_1 + \gamma_2) & \text{if } i \text{ is a non-leaf;} \\ D(\gamma) + D(\gamma) & \text{if } i \text{ is a leaf.} \end{cases}$$

Now suppose $k > 1$. In case (a), we may assume by induction that

$$\operatorname{Brac}_{k-1}(L) \cdot \gamma = D^{(k-1)/2}(\gamma) + D^{-(k-1)/2}(\gamma) + 2(\gamma_1 + \gamma_2) S_{k-1}(L).$$

Multiplying by $L$ and using equation (3.3) we have

$$L \cdot \operatorname{Brac}_{k-1}(L) \cdot \gamma = D^{k/2}(\gamma) + D^{(k-2)/2}(\gamma) + 2(\gamma_1 + \gamma_2) + D^{-k/2}(\gamma) + D^{-(k-2)/2}(\gamma) + 2(\gamma_1 + \gamma_2) + 2L(\gamma_1 + \gamma_2) S_{k-1}(L)$$

$$= D^{k/2}(\gamma) + D^{-k/2}(\gamma) + \operatorname{Brac}_{k-2}(L) \cdot \gamma - 2(\gamma_1 + \gamma_2) S_{k-2}(L)$$

$$+ 2(\gamma_1 + \gamma_2)(2 + L \cdot S_{k-1}(L)),$$

where the last equation holds by induction. Then

$$(L \cdot \operatorname{Brac}_{k-1}(L) - \operatorname{Brac}_{k-2}(L)) \cdot \gamma = D^{k/2}(\gamma) + D^{-k/2}(\gamma) + 2(\gamma_1 + \gamma_2)(L \cdot S_{k-1}(L) - S_{k-2}(L) + 2),$$

and using the recursions for the bracelets and for $S_k$ we get

$$\operatorname{Brac}_k(L) \cdot \gamma = D^{k/2}(\gamma) + D^{-k/2}(\gamma) + 2(\gamma_1 + \gamma_2) S_k(L).$$

In case (b), we assume by induction that

$$\operatorname{Brac}_{k-1}(L) \cdot \gamma = D^{k-1}(\gamma) + D^{-(k-1)}(\gamma).$$
Again multiplying by $L$ and using equation (3.3) we have
\[
L \cdot \text{Brac}_{k-1}(L) \cdot \gamma = D^k(\gamma) + D^{k-2}(\gamma) + D^{-k}(\gamma) + D^{-(k-2)}(\gamma)
\]
where the last equation holds by induction. Then
\[
(L \cdot \text{Brac}_{k-1}(L) - \text{Brac}_{k-2}(L)) \cdot \gamma = D^k(\gamma) + D^{-k}(\gamma)
\]
and the left hand side is equal to $\text{Brac}_k(L) \cdot \gamma$. \hfill \Box

For our next result we define the constants
\[
C_j(n) = X_{\text{Brac}_j(L)}|_{x_i=1} \quad j \in \mathbb{Z}_{\geq 0}.
\]
Note that $C_2(n) = C_1(n)^2 - 2$ and the value of the constant depends on the orientation of the edges in $Q$. We have the following linear recursion for the coordinates $f_i(t)$ of the points defining $X(Q)$.

**Theorem 3.8.** Let $Q$ be of type $\widetilde{D}$ with $n$ vertices.

(a) For all non-leaf vertices $i$ and all $t, p \in \mathbb{Z}_{\geq 0}$, we have
\[
f_i(t + 3p(n - 3)) = (C_{2p}(n) + 1) f_i(t + 2p(n - 3)) - (C_{2p}(n) + 1) f_i(t + p(n - 3)) + f_i(t).
\]

(b) For all leaf vertices $i$ and all $t \in \mathbb{Z}_{\geq 0}$, we have
\[
f_i(t + 2(n - 3)) = C_1(n) f_i(t + (n - 3)) - f_i(t) \quad \text{if } n \text{ is odd;}
\]
\[
f_i(t + 4(n - 3)) = C_2(n) f_i(t + 2(n - 3)) - f_i(t) \quad \text{for all } n.
\]

**Proof.** (a) Let $i$ be a non-leaf vertex. Combining Lemmas 3.6 and 3.7, we have for every preprojective arc $\gamma$ of orbit $i$
\[
\text{Brac}_{2p}(L) \cdot \gamma = \tau^{-p(n-3)}(\gamma) + \tau^{p(n-3)}(\gamma) + 2(\gamma_1 + \gamma_2) S_{2p}(L).
\]
Let $K$ denote the integer obtained by specializing the Laurent polynomial corresponding to $2(\gamma_1 + \gamma_2) S_{2p}(L)$ at $x_i = 1$. Then, if we let $\gamma = \tau^{-p(n-3)} \tau^{-t+1} P(i)$, the equation (3.4) yields
\[
C_{2p}(n) f_i(t + p(n - 3)) = f_i(t + 2p(n - 3)) + f_i(t) + K.
\]
Similarly, if we let $\gamma = \tau^{-2p(n-3)} \tau^{-t+1} P(i)$, the equation (3.4) yields
\[
C_{2p}(n) f_i(t + 2p(n - 3)) = f_i(t + 3p(n - 3)) + f_i(t + p(n - 3)) + K.
\]
Subtracting equation (3.5) from equation (3.6) and rearranging the terms we get
\[
f_i(t + 3p(n - 3)) = (C_{2p}(n) + 1) f_i(t + 2p(n - 3)) - (C_{2p}(n) + 1) f_i(t + p(n - 3)) + f_i(t).
\]
This completes the proof of (a).

(b) Let $i$ be a leaf vertex. Combining Lemmas 3.6 and 3.7, we have for every preprojective arc $\gamma$ of orbit $i$
\[
L \cdot \gamma = \tau^{-(n-3)}(\gamma) + \tau^{n-3}(\gamma) \quad \text{if } n \text{ is odd;}
\]
\[
\text{Brac}_{2}(L) \cdot \gamma = \tau^{-2(n-3)}(\gamma) + \tau^{2(n-3)}(\gamma) \quad \text{for all } n.
\]
We thus obtain
\[
C_1(n) \cdot f_i(t + n - 3) = f_i(t + 2(n - 3)) + f_i(t) \quad \text{if } n \text{ is odd;}
\]
\[
C_2(n) \cdot f_i(t + 2(n - 3)) = f_i(t + 4(n - 3)) + f_i(t) \quad \text{for all } n. \hfill \Box
\]

**Remark 3.9.** Part (b) of Theorem 3.8 was obtained in [10, Theorems 6.1 and 6.2] using cluster categories.
Corollary 3.10. Let $Q$ be an acyclic quiver of affine type $\widehat{\mathbb{D}}$. Then the dimension of the frieze variety $X(Q)$ is equal to one.

Proof. Suppose first that $n$ is odd. The characteristic polynomial of the recursions in Theorem 3.8 are

\[ x^3 - (C_2(n) + 1)x^2 + (C_2(n) + 1)x - 1 \] in case (a) with $p = 1;
\]
\[ x^2 - C_1(n)x + 1 \] in case (b) with $n$ odd;

Let $\rho$ be a root of $x^2 - C_1(n)x + 1$. Then $x^2 - C_1(n)x + 1 = (x - \rho)(x - \rho^{-1})$ and $C_1(n) = \rho + \rho^{-1}$. Moreover $\rho \neq \rho^{-1}$, since $C_1(n) > 2$. From the recursive formula of the bracelets we have $C_2(n) = C_1(n)^2 - 2$, thus $C_2(n) + 1 = \rho^2 + \rho^{-2} + 1$. Therefore the characteristic polynomial in case (a) is equal to $(x - \rho^2)(x - 1)(x - \rho^{-2})$. In particular the roots of the characteristic polynomials in case (a) and (b) are of the form $\rho^\ell$ with $\ell = -2, -1, 0, 1, 2$. Now the result follows from Lemma 2.2 and Lemma 3.2.

If $n$ is even, we use $p = 2$ so that the characteristic polynomial of the recursions in Theorem 3.8 are

\[ x^3 - (C_4(n) + 1)x^2 + (C_4(n) + 1)x - 1 \] in case (a) with $p = 2;
\]
\[ x^2 - C_2(n)x + 1 \] in case (b) for all $n$.

Now we let $\rho$ be a root of $x^2 - C_2(n)x + 1$, use the Chebyshev relation $C_4(n) = C_2(n)^2 - 2$ and the proof is analogous to the previous case. \hfill \Box

3.3. Affine type $E$. Let $R$ be a commutative ring, and $R'$ be a subring of $R$. If a sequence $(a_\ell)$ of elements in $R$ satisfies a recurrence relation $c_0a_{\ell+k} + c_1a_{\ell+k-1} + \cdots + c_ka_\ell = 0$ for all $\ell \geq 0$ (where $c_0, \ldots, c_k \in R'$), then we call the polynomial $f(x) = c_0x^k + c_1x^{k-1} + \cdots + c_k \in R'[x]$ an annihilator of $(a_\ell)$. Recall the following fact:

Lemma 3.11. Let $f(x), g(x) \in R'[x]$ be annihilators of $(a_\ell)$ and $(b_\ell)$, respectively. Then:

(i) $f(x)g(x) \in R'[x]$ is an annihilator of $(a_\ell + b_\ell)$ of degree $\deg f + \deg g$.

(ii) Let $h(x) \in R'[x]$ be the characteristic polynomial of the tensor product of the companion matrices of $f$ and $g$. Then $h(x)$ is an annihilator of $(a_\ell b_\ell)$ of degree $(\deg f)(\deg g)$.

Moreover, if the leading and constant coefficients of $f, g$ are $\pm 1$, then the same holds for the above annihilators.

Proof. The parts (i) and (ii) are given in [19], Lemma 4.1. The last statement is obvious. \hfill \Box

Let $Q$ be of type $\tilde{E}_{n-1}$ with $n$ vertices, so $n = 7, 8, \text{ or } 9$. Let $\delta = (\delta_1, \ldots, \delta_n), \delta_i > 0$ be the unique imaginary Schur root for $Q$; see for example [20. 8.2.1] for the values of the $\delta_i$. Then there exists a one parameter family $(M_\lambda)$ of non-isomorphic indecomposable representations with dimension vector $\delta$ and such that $\text{End}(M_\lambda) \cong k$. Moreover each $M_\lambda$ is regular and non-rigid. Choose one such representation $M$ such that $M$ is a regular simple representation. This is equivalent to the condition $\tau M \cong M$, in other words, $M$ sits at the mouth of a homogeneous tube in the Auslander-Reiten quiver. (Note that $X_M$ is denoted $X_\delta$ in [19].)

Define

\[(3.7) \quad m_n = \begin{cases} \quad 6, & \text{if } n = 7; \\ \quad 12, & \text{if } n = 8; \\ \quad 30, & \text{if } n = 9. \end{cases} \quad d_n = \begin{cases} \quad 12, & \text{if } n = 7; \\ \quad 29, & \text{if } n = 8; \\ \quad 169, & \text{if } n = 9. \end{cases} \]
Lemma 3.12. Let $Q$ be of type $\tilde{E}_{n-1}$ where $n = 7, 8, 9$. For each $i \in Q_0$, and $0 \leq r \leq m_n - 1$, the sequence $\left( f_i(jm_n + r) \right)_{j=0}^{\infty}$ is annihilated by a polynomial in $\mathbb{Z}[x]$ of degree $d_n$.

Proof. We explain the conclusion for $n = 7$ in detail, since the other two cases are proved similarly. The idea is to give the bounds of degrees of the recursive relations described in [19].

First note that the cluster category does not depend on the orientation of $Q$. Once we prove a recurrence relation of $(X_{\tau^iP(j)})_t$ for one particular orientation, changing the orientation of the quiver will not change the recursion. We use the following orientation of $a$ recurrence relation of $(19, \S)$.

\[
\begin{array}{c}
1 \quad 2 \quad 7 \quad 6 \quad 5 \\
\downarrow \\
4 \\
\downarrow \\
3
\end{array}
\]

This is the same labelling, and the opposite orientation as in [19, \S 7]. (We use the opposite orientation because of a different convention used in [19].) For simplicity, we denote $m_n$ by $m$. All relations below can be found in [19, \S 7].

For $i = 1, 3, 5$, $(X_{\tau^iP(j)})_t$ is annihilated by $x^2 - X_Mx + 1$. For $i = 2$, since

\[ X_{\tau^iP(2)} = X_{\tau^iP(1)}X_{\tau^{i+1}P(1)} - 1, \]

we have

\[ X_{\tau^iP(2)} = X_{\tau^iP(1)}X_{\tau^{i+1}P(1)} - 1. \]

Since both sequences $(X_{\tau^iP(1)})_j$ and $(X_{\tau^{i+1}P(1)})_j$ are annihilated by a polynomial of degree 2, and since the constant sequence $(1)_j$ is annihilated by $x - 1$, we conclude that the sequence $(X_{\tau^iP(1)}X_{\tau^{i+1}P(1)})_j$ is annihilated by a polynomial of degree $2 \cdot 2 = 4$, thus $(X_{\tau^iP(2)})_j$ is annihilated by a polynomial of degree $4 + 1 = 5$, using Lemma 3.11

The same conclusion holds for $i = 4, 6$.

For $i = 7$, we have

\[ X_{\tau^iP(7)} = X_{\tau^iP(1)}X_{\tau^{i+1}P(2)} - X_{\tau^iP(2)}X_{\tau^{i+2}P(1)}. \]

By Lemma 3.11, $(X_{\tau^iP(7)})_j$ is annihilated by a polynomial of degree $2 \cdot 5 + 2 = 12$.

We illustrate these degrees as follows, where the notation $i^{(d)}$ means vertex $i$ corresponding to an annihilating polynomial of degree $d$ with coefficients in $\mathbb{Z}[X_M]$:

\[
\tilde{E}_6: \quad 1^{(2)} \quad 2^{(5)} \quad 7^{(12)} \quad 6^{(5)} \quad 5^{(2)}
\]

Now specialize the annihilating polynomial at $x_1 = \cdots = x_n = 1$. Since all coefficients are in $\mathbb{Z}[X_M]$, the specialization is well-defined. Moreover, since the leading and constant coefficients are $\pm 1$, the specialization is not trivial. This gives us the desired polynomial.
For $n = 8$, the degree of an annihilating polynomial for each vertex are illustrated below:

\[ \tilde{E}_7 : \quad \begin{array}{cccccccc}
1^{(2)} & 2^{(5)} & 3^{(12)} & 8^{(29)} & 6^{(12)} & 5^{(5)} & 4^{(2)} \\
\end{array} \]

The degrees at vertices $i = 1, 2, 3, 8$ are obtained similar as the case $n = 7$, using identities

\[ X_M X_{\tau^{-t}P(1)} = X_{\tau^{-t+m}P(1)} + X_{\tau^{-m}P(1)}, \]
\[ X_{\tau^{-t}P(2)} = X_{\tau^{-t}P(1)} X_{\tau^{-(t+1)}P(1)} - 1, \]
\[ X_{\tau^{-t}P(3)} = X_{\tau^{-t}P(1)} X_{\tau^{-(t+1)}P(2)} - X_{\tau^{-(t+2)}P(1)}, \]
\[ X_{\tau^{-t}P(4)} = X_{\tau^{-t}P(1)} X_{\tau^{-(t+1)}P(3)} - X_{\tau^{-(t+2)}P(2)}, \]
\[ X_{\tau^{-t}P(5)} = X_{\tau^{-t}P(1)} X_{\tau^{-(t+1)}P(4)} - X_{\tau^{-(t+2)}P(3)}, \]
\[ X_{\tau^{-t}P(6)} = X_{\tau^{-t}P(1)} X_{\tau^{-(t+1)}P(5)} - X_{\tau^{-(t+2)}P(4)}. \]

The degrees of $i = 4, 5, 6$ are obtained by symmetry. For the vertex 7, we use the first two exchange triangles in [19 page 1857]:

\[ P_1 \to \tau^{-1}P_7 \to \tau^{-4}P_4 \to \tau P_1, \quad \tau^{-4}P_4 \to N \to P_1 \to \tau^{-3}P_4 \]

(where $N$ is the indecomposable regular simple module of dimension vector $11100101$ which belongs to the mouth of the tube of width 4) to obtain

\[ X_{\tau^{-(t+1)}P(7)} = X_{\tau^{-t}P(1)} X_{\tau^{-(t+4)}P(4)} - X_{\tau^{-t}N}. \]

Note that the sequence $(\tau^tN)_t$ has period 4 which divides $m = 12$, thus $(X_{\tau^{-(jm+r-1)}N})_j$ is a constant sequence. Substituting $t = jm + r - 1$ and using the fact that $(X_{\tau^{-(jm+r-1)}P(1)})_j, (X_{\tau^{-(jm+r+1)}P(4)})_j, (X_{\tau^{-(jm+r-1)}N})_j$ are annihilated by polynomials of degrees 2, 2, 1, respectively, we conclude that the sequence $(X_{\tau^{-(jm+r)}P(7)})_j$ is annihilated by a polynomial of degree $2 \cdot 2 + 1 = 5$.

For $n = 9$, we obtain

\[ \tilde{E}_8 : \quad \begin{array}{ccccccccccc}
1^{(2)} & 2^{(5)} & 3^{(12)} & 4^{(29)} & 5^{(70)} & 9^{(169)} & 8^{(29)} & 7^{(5)} \\
\end{array} \]

Indeed, for vertices $i = 1, 2, 3, 4, 5, 9$, we use

\[ X_M X_{\tau^{-t}P(1)} = X_{\tau^{-t+m}P(1)} + X_{\tau^{-m}P(1)}, \]
\[ X_{\tau^{-t}P(2)} = X_{\tau^{-t}P(1)} X_{\tau^{-(t+1)}P(1)} - 1, \]
\[ X_{\tau^{-t}P(3)} = X_{\tau^{-t}P(1)} X_{\tau^{-(t+1)}P(2)} - X_{\tau^{-(t+2)}P(1)}, \]
\[ X_{\tau^{-t}P(4)} = X_{\tau^{-t}P(1)} X_{\tau^{-(t+1)}P(3)} - X_{\tau^{-(t+2)}P(2)}, \]
\[ X_{\tau^{-t}P(5)} = X_{\tau^{-t}P(1)} X_{\tau^{-(t+1)}P(4)} - X_{\tau^{-(t+2)}P(3)}, \]
\[ X_{\tau^{-t}P(6)} = X_{\tau^{-t}P(1)} X_{\tau^{-(t+1)}P(5)} - X_{\tau^{-(t+2)}P(4)}. \]

For vertex 7 we use

\[ X_{\tau^{-t}P(7)} = X_{\tau^{-(t-2)}P(1)} X_{\tau^{-(t+5)}P(1)} - X_{\tau^{-(t+2)}N}. \]

(where $N$ is the indecomposable regular simple module of dimension vector $001111001$ which belongs to the mouth of the tube of width 5) and that 5 divides $m = 30$. 
For vertex 6 we use
\[ X_{\tau-t}P(6) = X_{\tau-(t-1)}P(1)X_{\tau-(t+2)}P(7) - X_{\tau-(t+7)}P(1) \]
For vertex 8 we use
\[ X_{\tau-t}P(8) = X_{\tau-(t-1)}P(1)X_{\tau-(t+1)}P(6) - X_{\tau-(t+3)}P(7) \]
This completes the proof.

We need to following simple fact.

**Lemma 3.13.** If a sequence \((a_j)_{j=0}^{\infty}\) is annihilated by a polynomial of degree \(d\), and \(c_0a_{n+k} + c_1a_{n+k-1} + \cdots + c_ka_n = 0\) holds for \(0 \leq n \leq d-1\), then the equality holds for every \(n\).

**Proof.** The sequence \((c_0a_{n+k} + c_1a_{n+k-1} + \cdots + c_ka_n)_n\) satisfies a recursive relation of degree \(d\) and its first \(d\) terms are 0, so it must be a constant 0 sequence. \(\square\)

**Proposition 3.14.** Let \(Q\) be an acyclic quiver of affine type \(\widetilde{E}\). Then the dimension of the frieze variety \(X(Q)\) is equal to one.

**Proof.** Define the constant \(C(n) = X_M|_{x_1=1}\) to be the specialization of the image of \(M\) under the Caldero-Chapoton map at \(x_1 = \cdots = x_n = 1\), and let
\[ \rho = (C(n) + \sqrt{C(n)^2 - 4})/2. \]
By definition of \(M\), we have \(C(n) > 2\) and therefore \(\rho \neq 1\).

We checked, by a computer, that for every \(i\) and \(r\) that the sequence \((f_i(jm_n + r))_{j=0}^{\infty}\) satisfies the (not-necessarily minimal) linear recurrence whose characteristic polynomial is
\[ \prod_{w=-6}^{6} (x - \rho^w). \]
Indeed, by Lemma 3.13 we only need to check that the linear recurrence holds for the first \(d_n\) instances, for each \(i \in Q_0, r = 0, \ldots, m_n - 1\), and each orientation.\(^1\)

The statement then follows from Lemma 3.2. \(\square\)

3.4. **A geometric remark.** Siegel’s theorem on integral points says that a smooth curve of genus at least one has only finitely many integral points. So in the affine case, each component of \(X(Q)\) is either of genus zero or singular. We conjecture that each component is a smooth curve of genus 0.

4. **Proof of the main theorem part (c), the wild case**

This section is divided into two subsections; in the first we recall facts on the Coxeter transformation and in the second we prove that \(\dim X(Q) > 1\) for wild type. We keep the notation of the previous sections.

4.1. **Coxeter transformation.** We recall some facts on the Coxeter matrix and its inverse, following the survey paper [27] (rewritten in our notation).

\(^1\)It took a few seconds for \(n = 7, 8\), and about half an hour for \(n = 9\) on an iMac.
4.1.1. The Coxeter matrix \( \Phi \). Let \( C = (c_{ij})_{1 \leq i,j \leq n} \) be the Cartan matrix of \( Q \), where \( c_{ij} \) is the number of paths from \( j \) to \( i \). Its inverse \( C^{-1} \) is the matrix \( (b_{ij})_{1 \leq i,j \leq n} \) where \( b_{ii} = 1 \) and if \( i \neq j \), then \( -b_{ij} \) is the number of arrows from \( j \) to \( i \) in \( Q \). Define the Coxeter matrix \( \Phi \) and its inverse \( \Phi^{-1} \) as

\[
\Phi = -C^T(C^{-1}) \quad \Phi^{-1} = -C(C^{-1})^T
\]

Then \( \Phi^{-1} \text{dim} \ M = \text{dim} \ (\tau^{-1}M) \) if \( M \) is not injective and \( \Phi^{-1} \text{dim} \ I(i) = -\text{dim} \ P(i) \). See for example [29, §3.1].

Let \( \rho_1, \ldots, \rho_n \) be the eigenvalues of \( \Phi^{-1} \) such that \( |\rho_1| \geq |\rho_2| \geq \cdots \geq |\rho_n| \), and \( \mathbf{v}_i = [v_{i1} \cdots v_{in}]^T \) a corresponding generalized eigenvector. The largest absolute value of the eigenvalues \( |\rho_1| \) is called the spectral radius of \( \Phi^{-1} \).

Recall that the characteristic polynomial of a matrix \( A \) is defined as \( \chi_A(x) = \det(xI - A) \). Now we recall some properties of the characteristic polynomial \( \chi_{\Phi^{-1}}(x) \) (which is called the Coxeter polynomial in [27]).

**Lemma 4.1.** (1) The following characteristic polynomials are equal:

\[
\chi_{\Phi^{-1}}(x) = \chi_{\Phi^T}(x) = \chi_{\Phi}(x).
\]

Therefore \( \Phi^{-1}, \Phi^T, \Phi \) have the same set of eigenvalues and the corresponding multiplicities. Moreover, the polynomials are monic, reciprocal and have integral coefficients; that is, if we write \( \chi_{\Phi^{-1}}(x) = \sum_{i=0}^n a_i x^i \), then \( a_n = 1 \) and \( a_i = a_{n-i} \in \mathbb{Z} \) for all \( i \).

In (2) and (3) we assume that \( Q \) is an acyclic wild quiver.

(2) The eigenvalue \( \rho_1 \) is equal to a real number \( \rho > 1 \), and has multiplicity 1. Moreover \( |\rho_1| < \rho \) for all \( i \neq 1 \). As a consequence, \( \mathbf{v}_1 \) is unique up to scale. Moreover, we can choose \( \mathbf{v}_1 \in \mathbb{R}_{>0}^n \), that is, all the coordinates of \( \mathbf{v}_1 = [v_{i1} \cdots v_{in}]^T \) are strictly positive.

(3) \( \rho_n = 1/\rho < 1 \) has multiplicity 1, and \( |\rho_n| > 1/\rho \) for all \( i \neq n \). As a consequence, \( \mathbf{v}_n \) is unique up to scale. Moreover, we can choose \( \mathbf{v}_n \in \mathbb{R}_{>0}^n \), that is, all the coordinates of \( \mathbf{v}_n = [v_{n1} \cdots v_{mn}]^T \) are strictly positive.

**Proof.** Most of the lemma is proved in [27]. The notation \( M, M^{-1}, C \) in [27] correspond to our \( (C^{-1})^T, C, \Phi^T \), respectively. (Below we shall also see that \( \rho(C), y^-, y^+ \) in [27] correspond to our \( \rho, \mathbf{v}_1^-, \mathbf{v}_n^+ \).)

(1) Since

\[
C^{-1} \Phi^{-1} C = C^{-1}(-C(C^{-1})^T)C = -(C^{-1})^T C = \Phi^T,
\]

we see that \( \Phi^{-1} \) and \( \Phi^T \) are similar, so \( \chi_{\Phi^{-1}}(x) = \chi_{\Phi^T}(x) \). Moreover, a matrix and its transpose have the same characteristic polynomial, so \( \chi_{\Phi^T}(x) = \chi_{\Phi}(x) \).

Moreover, note that \( \det(xI - \Phi^{-1}) \) has the leading coefficient \( a_n = 1 \), it has integral coefficients because all entries of \( \Phi^{-1} \) are integers (see §2.1), and the reciprocal property is proved in [27, §2.7].

(2) It is a result by Ringel [28] (Theorem 2.1 in [27]) that \( \rho > 1 \), that it is an eigenvalue of \( \Phi^T \) of multiplicity 1, and that other eigenvalues of \( \Phi^T \) have norm less than \( \rho \).

It is asserted in [27, §3.4] that there exists a (row) vector \( y^- \) with positive coordinates such that \( y^- \Phi^T = \rho^{-1} y^+ \). Thus \( \Phi(y^-)^T = \rho^{-1}(y^-)^T \), therefore \( \Phi^{-1}(y^-)^T = \rho(y^-)^T \). So we can take \( \mathbf{v}_1 = (y^-)^T \). This proves the last statement of (2).

(3) The first statement of (3) follows from (1) and (2); indeed, because \( \chi_{\Phi^{-1}}(x) \) being reciprocal is equivalent to \( \chi_{\Phi^{-1}}(x) = x^n \chi_{\Phi^{-1}}(x^{-1}) \) [27, §2.7], we have that \( \rho_i \) and \( \rho_i^{-1} \) are eigenvalues with the same multiplicity for every \( 1 \leq i \leq n \).
The proof of the second statement of (3) is similar to the proof of the second statement of (2), where $v_n = (y^+)^T$ for the vector $y^+$ defined in [27 §3.4].

**Lemma 4.2.** The eigenvalue $\rho$ is irrational.

**Proof.** The eigenvalue $\rho$ is a root of the characteristic polynomial $\chi_{\Phi^{-1}}(x)$ which, by Lemma 4.1, is an integral-coefficient polynomial whose leading coefficient and constant are both 1. The only possible rational roots of this polynomial are $\pm 1$ by the rational root theorem. But $\rho > 1$, so $\rho$ must be irrational. □

In the rest of the paper we require that, for $i = 1, n$,

$$v_i \in \mathbb{R}^n_{>0} \text{ and } ||v_i|| = 1.$$  
Such $v_1$ and $v_n$ exist uniquely by the above lemma.

**Lemma 4.3.** Let $Q$ be a wild acyclic quiver and $M$ an indecomposable, preprojective representation. Then $\lim_{t \to \infty} \frac{1}{\rho^t} \dim \tau^{-t} M = \lambda v_1$ for some real number $\lambda > 0$. As a consequence, there exist $c, N \in \mathbb{R}_{>0}$ such that, all components of $\dim \tau^{-t} P(i)$ are greater than or equal to $c\rho^t$ for every $t \geq N$ and every $i \in Q_0$.

**Proof.** The first statement is [27, Theorem 3.5]. The consequence is obvious. □

Indeed, a weaker version of Lemma 4.3 (replacing $\lambda > 0$ by $\lambda \geq 0$) is easy to prove, as shown in (1) of the lemma below.

Recall that the norm of a matrix $A$ is defined as

$$||A|| := \sup_{||x|| = 1} ||Ax||$$

and it satisfies the following inequalities, see for example [20, Theorem 14 on page 90].

(4.1) $||A + B|| \leq ||A|| + ||B||$ and $||AB|| \leq ||A|| \cdot ||B||$.

**Lemma 4.4.** (1) For any vector $v \in \mathbb{R}^n$, there exists $\lambda \in \mathbb{R}$ such that $\lim_{t \to \infty} \frac{1}{\rho^t} \Phi^{-t} v = \lambda v_1$.

(2) There exists a number $N$ such that $||\frac{1}{\rho^t} \Phi^{-t}|| \leq N$ for every $t \in \mathbb{Z}_{\geq 0}$.

**Proof.** (1) Let $v_1, \ldots, v_n$ be generalized eigenvectors corresponding to eigenvalues $\rho_1, \rho_2, \ldots, \rho_n$ such that

$$\Phi^{-1} V = VJ,$$

(denote $V := [v_1 \cdots v_n]$)

where $J$ is the Jordan normal form, thus

$$J = \begin{bmatrix} \rho & 0 & \cdots & 0 \\ 0 & J_2 & \cdots & 0 \\ \vdots & \vdots & \ddots & \vdots \\ 0 & 0 & \cdots & J_p \end{bmatrix}, \text{ where each } J_j \text{ is of the form } \begin{bmatrix} \rho_i & 1 & 0 & \cdots & 0 \\ 0 & \rho_i & 1 & \cdots & 0 \\ \vdots & \vdots & \ddots & \vdots & \vdots \\ 0 & 0 & \cdots & \rho_i \end{bmatrix} \text{ for some } i \neq 1. $$

Denote the size of $J_j$ as $m_j \times m_j$. We decompose the Jordan blocks

$$J_j = \rho_i I + K, \text{ where } I = \begin{bmatrix} 1 & 0 & 0 & \cdots & 0 \\ 0 & 1 & 0 & \cdots & 0 \\ \vdots & \vdots & \ddots & \vdots & \vdots \\ 0 & 0 & 0 & \cdots & 1 \end{bmatrix}, K = \begin{bmatrix} 0 & 1 & 0 & \cdots & 0 \\ 0 & 0 & 1 & \cdots & 0 \\ \vdots & \vdots & \ddots & \vdots & \vdots \\ 0 & 0 & 0 & \cdots & 1 \end{bmatrix}.$$
Since $IK = Kl = K$ and $K^\ell = 0$ for $\ell \geq m_j$, we have
\[
\frac{1}{\rho^\ell} J_j^\ell = \frac{1}{\rho^\ell} (\rho_i I + K)^\ell = \sum_{\ell=0}^{m_j-1} \frac{\rho_i^{\ell-\ell}}{\rho^\ell} \left( \frac{t}{\ell} \right) K^\ell
\]
Since $|\rho_i/\rho| < 1$, and the exponential grows faster than a polynomial, we have
\[
\lim_{t \to \infty} \left| \frac{\rho_i^{t-\ell}}{\rho^\ell} \left( \frac{t}{\ell} \right) \right| = \left| \frac{\rho_i}{\rho} \right|^t \lim_{t \to \infty} \left| \frac{\rho_i}{\rho} \right|^t = 0, \quad \text{for each } \ell = 0, \ldots, m_j - 1.
\]
Therefore every entry of the matrix $\frac{1}{\rho^\ell} J_j^\ell$ approaches 0 as $t \to \infty$. Thus
\[
\lim_{t \to \infty} \frac{1}{\rho^\ell} J_j^\ell = \begin{bmatrix} 1 & 0 & \cdots & 0 \\ 0 & 0 & \cdots & 0 \\ \vdots & \vdots & \ddots & \vdots \\ 0 & 0 & \cdots & 0 \end{bmatrix} =: E_{11}.
\]
Writing $v = \sum_i \lambda_i v_i = V\Lambda$ (where $\Lambda = \begin{bmatrix} \lambda_1 & \cdots & \lambda_n \end{bmatrix}^T$), and noticing that $\Phi^{-t} V = VJ^t$, we conclude
\[
\lim_{t \to \infty} \frac{1}{\rho^\ell} \Phi^{-t} v = \lim_{t \to \infty} \frac{1}{\rho^\ell} \Phi^{-t} V\Lambda = \lim_{t \to \infty} \frac{1}{\rho^\ell} VJ^t\Lambda = V(\lim_{t \to \infty} \frac{1}{\rho^\ell} J^\ell)\Lambda = V E_{11} \Lambda = \lambda_1 v_1.
\]
(2) It follows from (4.2) that $\|\frac{1}{\rho^\ell} J^\ell\| \to 1$ as $t \to \infty$, so there exists a number $N'$ such that
\[
\|\frac{1}{\rho^\ell} J^\ell\| \leq N', \quad \text{for every } t \in \mathbb{Z}_{\geq 0}.
\]
Then using the inequality $\|AB\| \leq \|A\| \cdot \|B\|$, we have
\[
\|\frac{1}{\rho^\ell} \Phi^{-t}\| \leq \|V\| \cdot \|\frac{1}{\rho^\ell} J^\ell\| \cdot \|V^{-1}\| \leq \|V\| \cdot N' \cdot \|V^{-1}\| =: N, \quad \text{for every } t \in \mathbb{Z}_{\geq 0}.
\]
\[
\square
\]
Similarly, we have

**Lemma 4.5.**  (1) For any vector $v \in \mathbb{R}^n$, there exists $\lambda \in \mathbb{R}$ such that $\lim_{t \to \infty} \frac{1}{\rho^t} \Phi^t v = \lambda v_n$.

(2) There exists a number $N$ such that $\|\frac{1}{\rho^t} \Phi^t\| \leq N$ for every $t \in \mathbb{Z}_{\geq 0}$.

**Proof.** Note that since $\rho_1, \ldots, \rho_n$ are eigenvalues of $\Phi^{-1}$, we see that $\rho_1^{-1}, \ldots, \rho_n^{-1}$ are eigenvalues of $\Phi$, and that $\rho_1^{-1} = 1/\rho$ and $\rho_n^{-1} = \rho$ have the largest and smallest norm. Moreover, since $\Phi^{-1} v_n = \rho^{-1} v_n$, we have $\Phi v_n = \rho v_n$, that is, $v_n$ is an eigenvector of $\Phi$ corresponding to the eigenvalue $\rho$. Then this lemma follows from Lemma 4.4. \[
\square
\]

4.2. **Proof of Theorem 1.1 (c).** We first need two results on the growth of the coefficients $f_i(t)$ in terms of the spectral radius $\rho$ of the inverse Coxeter matrix $\Phi^{-1}$.

**Lemma 4.6.** Let $d(t)$ be the largest coordinate in the vector $\dim \tau^{-t} P(i)$. Then $f_i(t) \geq 2^{d(t)}$. As a consequence, there exist $c, N_1 \in \mathbb{R}_{\geq 0}$ such that $f_i(t) \geq 2^{|c|}$ for every $t \geq N_1$, $i \in Q_0$. 

Proposition 4.7.

Recall that $v$ taking the initial cluster to be $x_i$ is defined in Lemma 4.1. Let $t$ be a real number $\eta > 0$ for some real number $\eta > 0$. The idea is to show that, for almost all $s$, the growth of $L(s + t)$ and $\Phi^{-t} L(s)$ are almost the same as $t \to \infty$, and the latter is well understood by Lemma 4.4.

Rewrite (1.1) as

$$f_i(t + 1) f_i(t) = 1 + \prod_{j \leftrightarrow i} f_{\tau_j}(t) \prod_{j \leftrightarrow i} f_j(t + 1)$$
Taking the logarithm on both sides and using the fact\footnote{To see $\ln(x+1) - \ln x < 1/x$, note the left side is $\ln(x+1) - \ln x = \ln(1+1/x)$. Replacing $1/x$ by $x$, it is equivalent to show $f(x) = x - \ln(1+x) > 0$ for all $x > 0$. Note that $f(0) = 0$. Since $f'(x) = 1-1/(1+x) > 0$ for $x > 0$, $f(x)$ is strictly increasing for $x \geq 0$, so the inequality follows.} that $0 < \ln(x + 1) - \ln x < 1/x$ for any positive real number $x$, we conclude

$$L_i(t + 1) + L_i(t) = \sum_{j \rightarrow i} L_j(t) + \sum_{j \leftarrow i} L_j(t + 1) + \delta_i(t),$$

where $\delta_i(t) > 0$, and by Lemma 4.6

$$(4.3) \quad \delta_i(t) < \frac{1}{\prod_{j \rightarrow i} f_j(t) \prod_{j \leftarrow i} f_j(t + 1)} < 2^{-cp'}, \text{ for } t \geq N_1.$$ 

Then

$$L_i(t + 1) - \sum_{j \leftarrow i} L_j(t + 1) = -L_i(t) + \sum_{j \rightarrow i} L_j(t) + \delta_i(t).$$

Rewriting these equations in matrix form brings up the inverse Cartan matrix as follows.

$$
\begin{bmatrix}
1 & 0 & \cdots & 0 \\
-b_{12} & 1 & 0 & \cdots & 0 \\
-b_{13} & -b_{23} & 1 & \cdots & 0 \\
\vdots & \vdots & \vdots & \ddots & \vdots \\
-b_{1n} & -b_{2n} & -b_{3n} & \cdots & 1
\end{bmatrix}\mathbf{L}(t + 1) =
\begin{bmatrix}
-1 & b_{12} & \cdots & b_{1n} \\
0 & -1 & b_{23} & \cdots & b_{2n} \\
0 & 0 & -1 & \cdots & b_{3n} \\
\vdots & \vdots & \vdots & \ddots & \vdots \\
0 & 0 & 0 & \cdots & -1
\end{bmatrix}\mathbf{L}(t) +
\begin{bmatrix}
\delta_1(t) \\
\delta_2(t) \\
\delta_3(t) \\
\vdots \\
\delta_n(t)
\end{bmatrix}
$$

which is

$$(C^{-1})^T\mathbf{L}(t + 1) = -(C^{-1})\mathbf{L}(t) + \mathbf{\delta}(t)$$

where $\mathbf{\delta}(t) = [\delta_1(t) \cdots \delta_n(t)]^T$ satisfying (by (4.3))

$$(4.4) \quad ||\mathbf{\delta}(t)|| < 2^{-cp'}\sqrt{n} < \sqrt{n}, \quad \text{for } t \geq N_1.$$ 

Left-multiplying the above equality by $C^T$, we obtain

$$\mathbf{L}(t + 1) = \Phi\mathbf{L}(t) + C\mathbf{\delta}(t)$$

It follows that, for $s,t \in \mathbb{Z}_{\geq0}$,

$$\mathbf{L}(s + t) = \Phi^{s+1}\mathbf{L}(s) + \Phi^{s+1-1}C\mathbf{\delta}(s) + \Phi^{s+1-2}C\mathbf{\delta}(s + 1) + \cdots + C\mathbf{\delta}(s + t - 1)$$
Assume $s \geq N_1$. Then
\[
\left\| \frac{1}{\rho^{s+t}} \left( L(s+t) - \Phi^t L(s) \right) \right\| = \frac{1}{\rho^{s+t}} \left\| \sum_{i=1}^{t} \Phi^{t-i} C \delta(s+i-1) \right\|
\leq \sum_{i=1}^{t} \frac{1}{\rho^{s+i}} \left\| \frac{1}{\rho^{i-1}} \Phi^{t-i} \right\| \cdot \left\| C \delta(s+i-1) \right\| \quad \text{by (4.1)}
\leq \sum_{i=1}^{t} \frac{1}{\rho^{s+i}} N \left\| C \right\| \cdot \left\| \delta(s+i-1) \right\| \quad \text{(This $N$ is the constant in Lemma 4.5 (2))}
\leq \sum_{i=1}^{t} \frac{1}{\rho^{s+i}} N \left\| C \right\| \cdot \sqrt{n} \quad \text{(by (4.4) and $s \geq N_1$)}
\leq N \left\| C \right\| \sqrt{n} \sum_{i=1}^{\infty} \frac{1}{\rho^{s+i}} = N \left\| C \right\| \sqrt{n} \frac{\rho^{-s-1}}{1 - \rho^{-1}} \to 0 \text{ as } s \to \infty.
\]

Therefore for any $\epsilon > 0$, we can find an integer $s_0$ sufficiently large such that
\[
(4.5) \quad \left\| \frac{1}{\rho^{s_0+t}} \left( L(s_0+t) - \Phi^t L(s_0) \right) \right\| < \epsilon, \quad \forall t \geq 0.
\]
By Lemma 4.5 (1), there exists $N_2 \in \mathbb{R}_{>0}$ and $\lambda \in \mathbb{R}$ such that
\[
(4.6) \quad \left\| \frac{1}{\rho^t} \Phi^t \left( \rho^{-s_0} L(s_0) \right) - \lambda v_n \right\| < \epsilon, \quad \forall t \geq N_2
\]
Adding the inequalities (4.5) and (4.6) and using the triangle inequality (4.1), we have
\[
\left\| \frac{1}{\rho^{s_0+t}} L(s_0+t) - \lambda v_n \right\| < 2\epsilon, \quad \forall t \geq N_2
\]
Now replace $s_0 + t$ by $t$. We have that, for any $\epsilon > 0$, there exist $\lambda \in \mathbb{R}$, $N_4 \in \mathbb{R}_{>0}$ (we can just take $N_4 = s_0 + N_2$), such that
\[
(4.7) \quad \left\| \frac{1}{\rho^t} L(t) - \lambda v_n \right\| < 2\epsilon, \quad \forall t \geq N_4
\]
thus for any $\epsilon > 0$, there exists $N_4 \in \mathbb{R}_{>0}$ such that
\[
\left\| \frac{1}{\rho^t} L(t) - \frac{1}{\rho^{t'}} L(t') \right\| < 4\epsilon, \quad \forall t, t' \geq N_4
\]
Therefore $\left\{ \frac{1}{\rho^t} L(t) \right\}$ is a Cauchy sequence, so must converge. Denote $u := \lim_{t \to \infty} \frac{1}{\rho^t} L(t)$. If $u$ is not on the line spanned by $v_n$, choose $\epsilon > 0$ such that $3\epsilon$ is less than the distance from $u$ and that line. Taking $t \to \infty$ in (4.7) gives the contradiction
\[
3\epsilon < \| u - \lambda v_n \| < 2\epsilon.
\]
Therefore $u = \eta v_n$ for some $\eta \in \mathbb{R}$.

Finally, we show that $\eta > 0$. By Lemma 4.6, there exists $c, N_1 \in \mathbb{R}_{>0}$ such that
\[
L_i(t) = \ln f_i(t) \geq c \rho^t \ln 2, \quad \text{for every } t \geq N_1 \text{ and every } i \in Q_0.
\]
\[ \eta \mathbf{v}_n = \lim_{t \to \infty} \frac{1}{\rho^t} \mathbf{L}(t) \geq [c \ln 2 \ c \ln 2 \ \cdots \ c \ln 2]^T \]

(here we mean that “\( \geq \)” holds componentwise), which implies \( \eta > 0 \).

We also need the following simple result.

**Lemma 4.8.** For all \( a_1, \ldots, a_m, x_1, \ldots, x_m > 0 \)
\[ \ln(a_1x_1 + \cdots + a_mx_m) \leq \max(\ln(x_i)) + \max(\ln(a_i)) + \ln m. \]

**Proof.** The left hand side is at most \( \ln(m \max(a_i \max(x_i))) \) which is equal to the right hand side. \( \square \)

We are now ready for the proof of our main result.

**Proof of Theorem 1.1(c).** Suppose \( X(Q) \) is contained in a 1-dimensional variety.

Consider the projection \( \pi : \mathbb{C}^n \to \mathbb{C}^2, (x_1, \ldots, x_n) \mapsto (x_1, x_2) \). Then \( \pi(X(Q)) \) is at most 1-dimensional. So there exists a nonzero polynomial \( g(x, y) = \sum_{(i,j) \in S} a_{ij}x^i y^j \in \mathbb{C}[x, y] \) (where \( a_{ij} \neq 0 \) for every \( (i,j) \in S \subset \mathbb{Z}_+^2 \)) such that \( g(f_1(t), f_2(t)) = 0 \) for every \( t \).

For convenience, denote the \( i \)-th coordinates \( v_{ni} \) of the eigenvector \( \mathbf{v}_n \) by \( y_i \) for \( i = 1, \ldots, n \), that is, \( \mathbf{v}_n = [y_1 \ \cdots \ y_n]^T \). Let \((i_0, j_0) \in S\) such that \( iy_1 + jy_2 \) is maximal. Replacing \( g \) by \( g/a_{i_0j_0} \) if necessary, we may assume \( a_{i_0j_0} = 1 \). Then
\[ f_1(t)^{i_0} f_2(t)^{j_0} = \sum_{(i,j) \in S \setminus (i_0, j_0)} (-a_{ij}) f_1(t)^i f_2(t)^j \]

Taking the logarithm on both sides, we get
\[ i_0 \ln L_1(t) + j_0 L_2(t) = \ln \left( \sum_{(i,j) \in S \setminus (i_0, j_0)} (-a_{ij}) f_1(t)^i f_2(t)^j \right) \]
and according to Lemma 4.8 we get
\[ i_0 \ln L_1(t) + j_0 L_2(t) \leq \max_{(i,j) \in S \setminus (i_0, j_0)} (i \ln L_1(t) + j L_2(t)) + \max_{(i,j) \in S \setminus (i_0, j_0)} (\ln |a_{ij}|) + \ln(|S| - 1). \]

Now we use Proposition 4.7. Dividing the above inequality by \( \rho^t \) and letting \( t \to \infty \), we conclude that there exists \((i, j) \in S \setminus (i_0, j_0) \) such that
\[ i_0 \eta y_1 + j_0 \eta y_2 \leq i \eta y_1 + j \eta y_2. \]

Thus
\[ i_0 y_1 + j_0 y_2 \leq i y_1 + j y_2. \]

By the choice of \((i_0, j_0)\), the equality must hold. Therefore \((i_0 - i)(y_1/y_2) = (j - j_0)\). Since \( y_1/y_2 \neq 0 \) and \((i, j) \neq (i_0, j_0)\), we must have \( j - j_0 \) and \( i_0 - i \) both been nonzero. Thus \( y_1/y_2 = (j - j_0)/(i_0 - i) \) is rational.

By a similar argument, \( y_i/y_j \) is rational for any \( 1 \leq i < j \leq n \). So there is a constant \( c \in \mathbb{R}_{>0} \) such that \( c \mathbf{v}_n \in \mathbb{Q}^n \). Then it follows from
\[ \rho(c \mathbf{v}_n) = \Phi(c \mathbf{v}_n) \in \mathbb{Q}^n \]
that \( \rho \) is rational, which contradicts Lemma 4.2.

This completes the proof. \( \square \)
Remark 4.9. It is natural to ask what is the exact dimension of \( X(Q) \) in the wild type. One might be tempted to conjecture that it is always equal to the number of vertices \( n \). However, this is not true, because whenever the quiver has a non-trivial automorphism \( \phi \in \text{Aut}(Q) \), then we have \( f_i(t) = f_{\phi(i)}(t) \), for all \( t \), so the dimension cannot be \( n \).

An upper bound for the dimension is the number of orbits under the action of \( \text{Aut}(Q) \) on \( Q_0 \).

5. Examples

5.1. Type \( A_{1,1} \). For the Kronecker quiver \( 1 \xrightarrow{2} 2 \) the points \( P_t = (f_1(t), f_2(t)) \) are \((1, 1), (2, 5), (13, 34), (89, 233), \ldots \) whose coordinates consist of every other Fibonacci number. The frieze variety \( X(Q) \subset \mathbb{C}^2 \) is given by the polynomial \( x^2 - 3xy + y^2 + 1 \). This is a smooth curve of genus zero.

5.2. Type \( A_{2,1} \). In this case the first few points are

\[(1, 1, 1), (2, 3, 7), (11, 26, 41), (97, 153, 362), (571, 1351, 2131), \ldots \]

The frieze variety \( X(Q) \) has two components, and each is a planar curve of degree 2. \( V(x_1 - 2x_2 + x_3, 2x_2^2 - 6x_2x_3 + 3x_3^2 + 1) \) and \( V(x_1 - 3x_2 + x_3, 3x_2^2 - 6x_2x_3 + 2x_3^2 + 1) \). Both are smooth curves of degree 2, hence of genus zero.

5.3. Type \( D_5 \). Consider the following quiver \( Q \):

\[\begin{array}{c}
2 \\
\downarrow \\
1 \leftrightarrow 3 \leftrightarrow 4 \leftrightarrow 6 \\
\end{array}\]

The first five points are

\[(1, 1, 1, 1, 1, 1), (2, 2, 5, 6, 7, 7), (3, 3, 11, 90, 13, 13), (4, 4, 131, 246, 19, 19), (33, 33, 2045, 3001, 158, 158).\]

The frieze variety \( X(Q) \) has three components.

\[V(x_5 - x_6, x_1 - x_2, x_6^2 + x_3 - 2x_4, 5x_2x_6 - x_3 - x_4 - 3, x_2^2 - 2x_3 + x_4),\]

\[V(x_5 - x_6, x_1 - x_2, x_6^2 + x_3 - 9x_4, 5x_2x_6 - x_3 - 8x_4 - 17, x_2^2 - 2x_3 + x_4),\]

\[V(x_5 - x_6, x_1 - x_2, x_6^2 + x_3 - 2x_4, 5x_2x_6 - 8x_3 - x_4 - 17, x_2^2 - 9x_3 + x_4).\]

Each component is a smooth curve of degree 2 and genus zero by the same argument as in the previous example.

5.4. A wild example. Consider the following quiver \( Q \):

\[\begin{array}{c}
2 \\
\downarrow \\
1 \xrightarrow{\quad} 3
\end{array}\]
Then
\[
\begin{array}{c|ccc|ccc}
 t & f_1(t) & f_2(t) & f_3(t) & \mathbf{L}(t)^\dagger \\
 1 & 2 & 3 & 13 & 0.693 & 1.099 & 2.565 \\
 2 & 254 & 1101 & 5464009 & 5.537 & 7.004 & 15.514 \\
 3 & 1.294 \times 10^{14} & 6.422 \times 10^{17} & 1.969 \times 10^{39} & 32.49 & 41.00 & 90.48 \\
 4 & 1.923 \times 10^{82} & 5.895 \times 10^{103} & 1.107 \times 10^{229} & 189.47 & 238.94 & 527.39 \\
 5 & 3.759 \times 10^{479} & 7.063 \times 10^{604} & 9.012 \times 10^{1334} & 1104.26 & 1392.72 & 3073.85 \\
\end{array}
\]

The Cartan matrix and its inverse, the Coxeter matrix and its inverse are:
\[
C = \begin{bmatrix} 1 & 1 & 3 \\
0 & 1 & 1 \\
0 & 0 & 1 \end{bmatrix}, \quad C^{-1} = \begin{bmatrix} 1 & -1 & -2 \\
0 & 1 & -1 \\
0 & 0 & 1 \end{bmatrix}, \quad \Phi = \begin{bmatrix} -1 & 1 & 2 \\
-1 & 0 & 3 \\
-3 & 2 & 6 \end{bmatrix}, \quad \Phi^{-1} = \begin{bmatrix} 6 & 2 & -3 \\
3 & 0 & -1 \\
2 & 1 & -1 \end{bmatrix},
\]

The characteristic polynomial \(\chi_{\Phi^{-1}}(x) = x^3 - 5x^2 - 5x + 1 = (x + 1)(x^2 - 6x + 1)\), so \(\rho = 3 + \sqrt{8} \approx 5.8284\) and \(1/\rho = 3 - \sqrt{8} \approx 0.1716\) are irrational, and the corresponding eigenvectors are
\[
\mathbf{v}_1 \approx [0.866, 0.392, 0.311]^T, \quad \mathbf{v}_n \approx [0.311, 0.392, 0.866]^T
\]
(In this example, \(\mathbf{v}_n\) happens to be the “reverse” of \(\mathbf{v}_1\). This is not the case in general.)
Computation shows
\[
\lim_{t \to \infty} \frac{1}{\rho^t} \mathbf{L}(t) \approx \begin{bmatrix} 0.164 \\
0.207 \\
0.457 \end{bmatrix} \approx 0.528 \mathbf{v}_n.
\]
So roughly we can describe the growth of \(P_i(t)\) as
\[
P_i(t) = (f_1(t), f_2(t), f_3(t)) \approx (e^{0.164\rho^t}, e^{0.207\rho^t}, e^{0.457\rho^t}).
\]

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