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# Gabriel's Theorem 

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#### Abstract

This thesis is devoted to the study of quiver representations and the proof of Gabriel's Theorem. Quiver representations are collections of vector spaces and linear maps. Just like with decomposing numbers into products of primes, one aims to decompose quiver representations into the smallest possible building blocks, called indecomposable representations. Gabriel's Theorem specifies the quivers that have a finite number of indecomposable representations. Furthermore, it classifies these indecomposable representations. I prove a decomposition theorem by Krull, Remak, and Schmidt in Chapter 2. This theorem relies heavily on indecomposable representations, which form the centerpiece of Chapter 3. In this chapter I introduce a variety of mathematical tools to analyze quivers and their representations, culminating in the proof of Gabriel's Theorem. I also show how Gabriel's Theorem can be used in topological data analysis to characterize the persistence of topological features. To do this, I present the necessary quiver representation theory in Chapter 4 and I connect this theory with the topological aspects of persistent homology in Chapter 5. The thesis is based on [1] and [2].


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#### Abstract

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# Acronyms and Abbreviations 

w.l.o.g. Without loss of generality<br>w.r.t With respect to<br>point cloud Finite metric space<br>$\bar{B}_{r}(x) \quad$ Closed ball with center $x$ and radius $r$<br>$S^{n} \quad$ Sphere of dimension $\mathrm{n},\left\{\left(x_{0}, x_{1}, \ldots x_{n}\right) \mid x_{0}^{2}+x_{1}^{2}+\ldots 1 x_{n}^{2}=1\right\} \subseteq \mathbb{R}^{n+1}$<br>$|X| \quad$ Cardinality of the set $X$<br>$Q \quad$ Quiver with vertex set $Q_{0}$ and arrow set $Q_{1}$

## Chapter 1

## Introduction

Topology is a mathematical discipline focused on investigating the characteristics of spaces that remain unchanged under deformations. These deformations involve actions like bending, squishing, shrinking and expanding, whereas tearing is prohibited. Homotopy equivalence emerges as a formalism to precisely describe this equivalence between topological spaces.

A homotopy between two continuous maps $f, g: X \rightarrow Y$ is a continuous map $H: X \times[0,1] \rightarrow Y$ such that $H(x, 0)=f(x)$ and $H(x, 1)=g(x)$ for all $x \in X$. The map H describes a family of functions interpolating continuously between $f$ and $g$ as the parameter $t$ varies from 0 to 1 . A pair of continuous functions $f: X \rightarrow Y$ and $g: Y \rightarrow X$ are called homotopy equivalences if there exists a homotopy $H$ between the maps $g \circ f$ and $\operatorname{id}_{X}$, and there exists a homotopy $K$ between the maps $f \circ g$ and $\mathrm{id}_{Y}$. In this case, the topological spaces $X$ and $Y$ are called homotopy equivalent. The circle $S^{1}$ and the punctured disk $D^{\prime}=D \backslash\{(0,0)\}$ are homotopy equivalent. The maps $f: D^{\prime} \rightarrow S^{1}, x \mapsto \frac{x}{\|x\|}$ and $g: S^{1} \rightarrow D^{\prime}, x \mapsto x$ are homotopy equivalences. Indeed, note that $f \circ g=\operatorname{id}_{S^{1}}$ and $H(x, t)=(1-t) \frac{x}{\|x\|}+t x$ is a homotopy between $g \circ f$ and $\operatorname{id}_{D^{\prime}}$.


Figure 1.1: The punctured disk $D^{\prime}$ (left) and the circle $S^{1}$ (right) are homotopy equivalent. The red and blue points represent $H(x, t)$ for two points in $D^{\prime}$, where $H(x, t)=(1-t) \frac{x}{\|x\|}+t x$ denotes the homotopy between the maps $(g \circ f)(x)=\frac{x}{\|x\|}$ and $\operatorname{id}_{D^{\prime}}$.

Properties of topological spaces that stay the same under deformations are called homotopy invariants. These invariants are the main objects of study in the field of algebraic topology. One such invariant is the presence of holes, voids and higher dimensional equivalents in a topological space and counting how many there are. In Figure 1.1, the punctured disk and the circle are homotopy equivalent. Observe that the volumes of the spaces differ, however, the middle hole is present throughout the deformation process. The mathematical formalization to measure holes and their higher dimensional analogues is called homology. Homology and other invariants measure the shape of topological spaces.

Recently, with a huge amount of data being generated, there have been efforts to adapt these methods to measure the shape of data. Topological data analysis offers one potential approach
to address this challenge. However, since data is given by point clouds, its topology does not reveal any information besides the number of points. To fix this problem, we consider triangulations of point clouds that are called Čech complexes. To compute a Čech complex, we replace the points with balls of a certain radius. If two balls intersect, we connect the corresponding points by an edge. The non-trivial intersection of three balls results in a triangle and so on. A crucial


Figure 1.2: A point cloud with associated balls on the left and its associated Čech complex on the right.
characteristic is that the Čech complex corresponding to a smaller radius is included in the Čech complex associated with a larger radius. Figure 1.2 illustrates the construction of a Čech complex from a point cloud. Using the inclusion property of Čech complexes of increasing radii, we get a sequence of Čech complexes and inclusion maps between them, called a filtration. We then apply homology to this filtration and track over which parameter values the topological features persist. This yields a collection of intervals. Each interval $[a, b]$ represents a topological feature born at time $a$ that disappears at time $b$. These intervals form a persistence module and they are visualized by a persistence barcode. Figure 1.3 shows the persistence barcode for the point cloud in Figure 1.2. Topological features that persist over a long parameter range are considered to be important with short-lived features as noise. This adaptation of homology to the setting of point clouds is called persistent homology.


Figure 1.3: A persistence barcode coming from a point cloud.
The existence of persistence barcodes relies on a theorem from quiver representation theory, called Gabriel's Theorem. This discipline focuses on quivers and their corresponding representations. Quivers are directed graphs consisting of points and arrows connecting the points. A quiver representation is an allocation of vector spaces to the points and linear maps are allocated to the arrows. Figure 1.4 shows an example of a quiver and a possible representation. This quiver is called a linear quiver. Quiver representation theory is built on a purely algebraic foundation which offers distinct advantages. Firstly, the broad algebraic framework allows for versatile application in various contexts without the necessity of introducing novel concepts. This flexibility will be used to extend the theory of persistent homology to include zigzag persistent homology. Additionally,
it is worth noting that quiver representation theory finds applications beyond topology, such as in the realms of Lie algebras and quantum groups [2]. However, these applications are beyond the scope of this thesis.


Figure 1.4: The linear quiver $L_{n}$ (left) and a quiver representation (right).
The goal of this thesis is to prove Gabriel's Theorem. In Chapter 2, we introduce quiver representations and observe that they always possess a unique decomposition into the smallest possible building blocks called indecomposable representations. This naturally leads to the question of which quivers only have a finite number of indecomposable representations. Such quivers are called finite-type quivers. The answer to this question is given by Gabriel's Theorem, whose proof is at the center of Chapter 3. It turns out that the persistence modules are quiver representations of the linear quivers. Therefore, we can apply Gabriel's Theorem in the setting of persistence modules. This is the reason why a persistence barcode always exists. Chapter 4 focuses on the algebraic aspects of the quiver theory of $A_{n}$-type quivers, which are a broader generalization of the linear quivers. In Chapter 5, we introduce the topological aspects of persistent homology while using the theory from the previous chapters.

## Chapter 2

## Quiver Representations

In this chapter, we introduce quiver representations and we consider their morphisms and direct sums of quiver representations. This culminates in a decomposition theorem by Krull, Remak, and Schmidt (Theorem 2.33). This chapter is primarily based on the book 'Persistence Theory: From Quiver Representations to Data Analysis' written by Steve Y. Oudot [1].

### 2.1 Quivers

Definition 2.1. A quiver $Q$ consists of two sets $Q_{0}, Q_{1}$ and two maps $h, t: Q_{1} \rightarrow Q_{0}$. The elements of $Q_{0}$ are called the vertices of $Q$, while those of $Q_{1}$ are called the arrows. The head map $h$ and the tail map $t$ assign $a$ head $h_{a}$ and a tail $t_{a}$ to every arrow $a \in Q_{1}$.

Example 2.2. Consider the quiver $Q$ with vertices $Q_{0}=\{1,2,3\}$ and arrows $Q_{1}=\{a, b, c, d, e\}$.


Its head map $h: Q_{1} \rightarrow Q_{0}$ is given by $a \mapsto 2, b \mapsto 2, c \mapsto 3, d \mapsto 1, e \mapsto 1$ and its tail map $t: Q_{1} \rightarrow Q_{0}$ is given by $a \mapsto 1, b \mapsto 2, c \mapsto 2, d \mapsto 3, e \mapsto 3$.
$Q$ is a directed graph, where the elements in $Q_{0}$ are the vertices and for every $a \in Q_{1}$, the pair $\left(t_{a}, h_{a}\right)$ is a directed edge in our graph. There are no restrictions on the sets $Q_{0}$ and $Q_{1}$, so there may be infinitely many points and edges and also multiple edges between two points, thus a quiver is graphically represented by a directed multigraph. We denote by $\bar{Q}$ the underlying undirected graph of $Q$.

Definition 2.3. A quiver $Q$ is called finite if both $Q_{0}$ and $Q_{1}$ are finite sets.
Definition 2.4. A quiver $Q$ is called a Dynkin quiver if its underlying graph $\bar{Q}$ is one of the graphs in Figure 2.1.

Dynkin quivers emerge in Gabriel's Theorem (Theorem 3.2). They also play a fundamental role in classifying semisimple Lie algebras [2, p. 29]. An important subset of the Dynkin quivers are the linear quivers.

Definition 2.5. For each $n \in \mathbb{N}$, the following quiver is called the linear quiver $L_{n}$ :

$$
\stackrel{\bullet}{1} \longrightarrow \stackrel{\bullet}{2} \longrightarrow \cdot \vec{n}-\stackrel{\bullet}{n}
$$

Definition 2.6. A quiver $Q$ is called acyclic if there is no oriented cycle in $Q$.
Every Dynkin quiver is acyclic. An example of a non-acyclic quiver is the loop quiver.


Figure 2.1: The Dynkin diagrams.

Definition 2.7. The loop quiver consists of a single vertex that loops to itself.


Sometimes, it is interesting to consider a part of a given quiver as it allows us to consider a property of the smaller quiver, and then extend it to the bigger quiver.
Definition 2.8. Let $Q=\left(Q_{0}, Q_{1}\right)$ be a quiver and let $A_{0} \subseteq Q_{0}$ and $A_{1} \subseteq Q_{1}$ such that for all arrows $a \in A_{1}$ we have that $h(a) \in A_{0}$ and $t(a) \in A_{0}$. Then $A=\left(A_{0}, A_{1}\right)$ defines a subquiver of $Q$ ( $A$ is a quiver). Moreover, if for all $i, j \in A_{0}$ we have that $\left\{a \in Q_{1} \mid t(a)=i, h(a)=j\right\} \subseteq A_{1}$, then $A$ is called a full subquiver with support $A_{0}$. In addition, if we say that we delete a vertex $i$ from a quiver $Q$, then the resulting (sub)quiver is the full subquiver with support $Q_{0} \backslash\{i\}$.
Example 2.9. Consider the quiver $Q$ on the left. The quivers $A$ (middle) and $B$ (right) are both subquivers of $Q$. However, $B$ is a full subquiver whereas $A$ is not full.


### 2.2 The Category of Quiver Representations

For a fixed quiver $Q$, we define quiver representations and morphisms between them. This is done in a way that turns them into a category. This category is called the category of quiver representations.

Definition 2.10. A quiver representation of a quiver $Q$ over a field $k$ is a pair $\mathbb{V}=\left(V_{i}, v_{a}\right)$, which consists of a set of $k$-vector spaces $\left\{V_{i} \mid i \in Q_{0}\right\}$ together with a set of $k$-linear maps $\left\{v_{a}: V_{t_{a}} \rightarrow V_{h_{a}} \mid a \in Q_{1}\right\}$. We often abbreviate quiver representation to representation.

There are no restrictions on the vector spaces and the maps. Thus the maps generally do not commute, i.e. if we have two different paths with the same starting point and end point, the compositions of the maps along the two paths need not be the same. Also, the vector spaces can be infinite-dimensional.

Definition 2.11. A quiver representation $\boxtimes$ of a quiver $Q$ is called finite-dimensional if the sum of the dimensions of the vector spaces $\sum_{i \in Q_{0}} \operatorname{dim} V_{i}$ is finite. If $Q$ is a finite quiver with vertex set $Q_{0}=\{1,2, \ldots, n\}$ and $\boxtimes$ is a finite-dimensional representation, then we define the dimension vector dim $\bigvee$ and its dimension $\operatorname{dim} \vee$ as:

$$
\begin{aligned}
& \underline{\operatorname{dim} \mathbb{V}}=\left(\operatorname{dim} V_{1}, \ldots, \operatorname{dim} V_{n}\right)^{T} \\
& \operatorname{dim} \mathbb{V}=\|\underline{\operatorname{dim} \mathbb{V}}\|_{1}=\sum_{i=1}^{n} \operatorname{dim} V_{i}
\end{aligned}
$$

Example 2.12. A representation of the linear $L_{2}$-quiver

is a set consisting of two vector spaces $V_{1}, V_{2}$ together with a linear map $v_{a}: V_{1} \rightarrow V_{2}$.
Example 2.13. [3] A representation of the loop quiver

is a vector space $V_{1}$ together with an endomorphism $v_{a}: V_{1} \rightarrow V_{1}$. If we restrict to the case where the field $k$ is algebraically closed and $V_{1}$ is finite-dimensional, then we know from linear algebra that the matrix of $v_{a}$ has a Jordan normal form (in a suitable basis)

$$
\left(\begin{array}{cccc}
J_{n_{1}, \lambda_{1}} & 0 & \cdots & 0 \\
0 & J_{n_{2}, \lambda_{2}} & & 0 \\
\vdots & & \ddots & \vdots \\
0 & 0 & \cdots & J_{n_{r}, \lambda_{r}}
\end{array}\right)
$$

where $J_{n, \lambda}$ denotes the $n \times n$-Jordan block:

$$
\left(\begin{array}{ccccc}
\lambda & 1 & & & \\
& \lambda & 1 & \emptyset & \\
& & \ddots & \ddots & \\
& \emptyset & & \lambda & 1 \\
& & & & \lambda
\end{array}\right)
$$

Definition 2.14. A representation $\mathbb{W}=\left(W_{i}, w_{a}\right)$ is a subrepresentation of a representation $\mathbb{V}=\left(V_{i}, v_{a}\right)$ if $W_{i}$ is a subspace of $V_{i}$ for all $i \in Q_{0}$ and if for all $a \in Q_{1}, w_{a}$ is the restriction of the map $v_{a}$ to the subspace $W_{t_{a}}$ of the domain and the subspace $W_{h_{a}}$ of the image, i.e. $w_{a}=\left.v_{a}\right|_{W_{t_{a}}} ^{W_{h_{a}}}$. We call $\mathbb{W}$ a proper subrepresentation of $\mathbb{V}$ if $0 \subsetneq \mathbb{W} \subsetneq \mathbb{V}$.

Definition 2.15. A representation of a quiver $Q$ is called simple if it is non-trivial and it has no proper subrepresentations.
Example 2.16. We fix a vertex $i \in Q_{0}$ of a quiver $Q$. We define the representation $\mathbb{S}_{i}=\left(S_{j}^{(i)}, s_{a}^{(i)}\right)$ to be

$$
S_{j}^{(i)}=\left\{\begin{array}{ll}
k & \text { if } j=i \\
0 & \text { if } j \neq i
\end{array} \quad \text { and } s_{a}^{(i)}=0\right.
$$

for $j \in Q_{0}$ and $a \in Q_{1}$. This representation is simple. Indeed if $\mathbb{W} \subsetneq \mathbb{S}_{i}$ is a subrepresentation, then $W_{j}=0$ for all $j \in Q_{0}$ and thus $\mathbb{W}=0$ is trivial.

Our goal is not only to define quiver representations but also to compare and classify them. To be able to do this, we need to define morphisms between quiver representations.

Definition 2.17. A morphism $\phi$ between two representations $\mathbb{V}, \mathbb{W}$ of a quiver $Q$ is a set of $k$-linear maps $\left\{\phi_{i}: V_{i} \rightarrow W_{i} \mid i \in Q_{0}\right\}$ such that the following diagram commutes for each arrow $a \in Q_{1}$ :


The morphism is called a monomorphism if every linear map $\phi_{i}$ is injective, an epimorphism if every $\phi_{i}$ is surjective, and an isomorphism (denoted by $\cong$ ) if every $\phi_{i}$ is bijective. If $\phi: \mathbb{V} \rightarrow \mathbb{V}$ maps onto itself, we call it an endomorphism. We denote the set of morphisms from $\mathbb{\vee}$ to $\mathbb{W}$ by $\operatorname{Hom}(\mathbb{V}, \mathbb{W})$. If we equip $\operatorname{Hom}(\mathbb{V}, \mathbb{W})$ with pointwise multiplication and addition of linear maps, then it is a $k$-vector space.

Example 2.18. A morphism between two representations $\mathbb{V}, \mathbb{W}$ of the $L_{2}$-quiver is given by two linear maps $\phi_{1}: V_{1} \rightarrow W_{1}, \phi_{2}: V_{2} \rightarrow W_{2}$ such that the following diagram commutes


Lemma 2.19. Every isomorphism $\phi: \mathbb{V} \rightarrow \mathbb{W}$ is invertible, meaning that the map $\psi=\phi^{-1}: \mathbb{W} \rightarrow \mathbb{V}$ is a morphism of quiver representations (where $\psi_{i}=\left(\phi_{i}\right)^{-1}$ ). Thus the expressions isomorphism and invertible morphism are exchangeable.

Proof. It is clear that $\psi$ is pointwise well-defined. We need to check that for each arrow $a \in Q_{1}$ the following diagram commutes


Now for each $w \in W_{i}$ we have $\phi_{j} v_{a} \psi_{i}(w)=w_{a}(w)$ since $\psi_{i}=\left(\phi_{i}\right)^{-1}$ and using that the diagram commutes with respect to the maps $\phi_{i}, \phi_{j}$. But then $v_{a} \psi_{i}(w)=\psi_{j} w_{a}(w)$. Thus the diagram commutes.

To turn quiver representations into a category, we need to define the composition of two morphisms in an associative way that guarantees the existence of an identity morphism. We can define such a composition by composing the maps $\phi_{i}$ at each point in our quiver.
Definition 2.20. The composition of two morphisms $\phi: \mathbb{U} \rightarrow \mathbb{V}$ and $\psi: \mathbb{V} \rightarrow \mathbb{W}$ is given by the maps $(\psi \circ \phi)_{i}=\psi_{i} \circ \phi_{i}$ at each point $i \in Q_{0}$ in our quiver.

This composition of morphisms is associative since it inherits the associativity of the composition of functions. Moreover, for each representation $\mathbb{V}$ we have the identity morphism $\mathbb{1}_{\vee}: \mathbb{V} \rightarrow \mathbb{V}$ which is the identity on each $V_{i}$. Thus for a fixed quiver $Q$ and a fixed field $k$, we get the category of the quiver representations of $Q$, denoted by $\operatorname{Rep}_{\mathrm{k}}(Q)$. If we restrict ourselves to the finite-dimensional representations we get the subcategory $\operatorname{rep}_{\mathrm{k}}(Q)$.

Example 2.21. We give a description of the category $\operatorname{rep}_{\mathrm{k}}(Q)$ for the loop quiver $Q$. A morphism between two representations $\mathbb{V}=\left(V_{1}, v_{a}\right), \mathbb{W}=\left(W_{1}, w_{a}\right)$ of the loop quiver is a map $\phi_{1}: V_{1} \rightarrow W_{1}$ such that $w_{a} \phi_{1}=\phi_{1} v_{a}$. In Example 2.13 we saw that every representation is given by a map $v_{a}$ which is in Jordan normal form in some suitable basis (if we restrict to an algebraically closed field $k$ and to finite-dimensional representations). We also know that two finite-dimensional endomorphisms $A, B \in \operatorname{End}\left(V_{1}\right)$ are equivalent if and only if their Jordan normal forms are the same (up to reordering the Jordan blocks). Equivalence of matrices means that there exists a change of basis matrix $C \in \operatorname{End}\left(V_{1}\right)$ such that $A=C^{-1} B C$. Notice that such a change of basis represents an isomorphism between two representations given by the matrices $A$ and $B$. Thus we have found that all representations are (up to isomorphism) given by Jordan normal forms and that two representations are isomorphic if and only if they have the same Jordan normal form (up to reordering the Jordan blocks). This is a complete description of all the isomorphism classes of the category $\operatorname{rep}_{\mathrm{k}}(Q)$ for the loop quiver $Q$ and an algebraically closed field $k$.

### 2.2.1 Direct Sums, Kernels and Cokernels

The categories $\operatorname{Rep}_{\mathrm{k}}(Q)$ and $\operatorname{rep}_{\mathrm{k}}(Q)$ have some properties that are useful towards our goal of classifying quiver representations. We list them here to give an overview and then we describe the second and third properties in detail.

ZO: Both categories contain a zero object, called the trivial representation, with all spaces and all maps equal to 0 .

DS: We can combine representations to get new representations (of the same quiver). This is called a direct sum.

KI: Every morphism between representations (of the same quiver) $\phi: \mathbb{V} \rightarrow \mathbb{W}$ has a kernel, an image and a cokernel.

DS: The direct sum of two representations is defined for any representations $\mathbb{V}, \mathbb{W}$ to be the representation $\mathbb{V} \oplus \mathbb{W}$ with spaces $V_{i} \oplus W_{i}$ for $i \in Q_{0}$ and maps $v_{a} \oplus w_{a}=\left(\begin{array}{cc}v_{a} & 0 \\ 0 & w_{a}\end{array}\right)$ for $a \in Q_{1}$. This definition explicitly shows what happens on vector spaces and linear maps, which is helpful if we look at a given direct sum of representations. We give another equivalent definition of the direct sum, which is especially useful in proofs.

Definition 2.22. [4, Def. 2.1] Let $\mathbb{V}_{1}, \ldots, \mathbb{V}_{r}$ be representations. A direct sum $\mathbb{V}=\mathbb{V}_{1} \oplus \ldots \oplus \mathbb{V}_{r}$ is a representation $\mathbb{V}$ together with morphisms $\iota_{i}: \mathbb{V}_{i} \rightarrow \mathbb{V}$ and $\pi_{i}: \mathbb{V} \rightarrow \mathbb{V}_{i}$ for $1 \leq i \leq r$ such that $\sum_{i=1}^{r} \iota_{i} \pi_{i}=\mathrm{id} \mathbb{V}$ and $\pi_{i} \iota_{i}=\mathrm{id}_{\vee_{i}}$. We write $\mathbb{V}^{r}=\mathbb{V} \oplus \ldots \oplus \mathbb{V}$ for the direct sum of $r$ copies of $\mathbb{V}$.
Remark 2.23. In the first definition of a direct sum, we can set the inclusion morphisms $\iota_{\checkmark, w}$ to be pointwise inclusion of subspaces and the projection morphisms $\pi_{v, w}$ are given pointwise by the projection onto a subspace. One can easily check that those morphisms satisfy the above conditions. Thus our second definition of a direct sum is equivalent to the first definition.

Definition 2.24. A non-trivial representation $\boxtimes$ is called decomposable if it is isomorphic to the direct sum of two non-trivial representations which are called summands. Else it is called indecomposable. The trivial representation is neither decomposable nor indecomposable.

Example 2.25. Again, we consider the loop quiver. Moreover, we assume that the field $k$ is algebraically closed and we consider finite-dimensional representations in $\operatorname{rep}_{\mathrm{k}}(Q)$. If $v_{a}, w_{a}$ are matrices in Jordan normal form then $v_{a} \oplus w_{a}=\left(\begin{array}{cc}v_{a} & 0 \\ 0 & w_{a}\end{array}\right)$ is also a matrix in Jordan normal form and thus the direct sum is a representation of the loop quiver (see Example 2.21). In addition, if $v_{a}$ has two or more Jordan blocks, then we can set $v_{a}^{\prime}$ to be the subrepresentation consisting only of the first (upper left) Jordan block and $v_{a}^{\prime \prime}$ to be all the other Jordan blocks (and in both cases restricting to the proper subspaces). Then $v_{a}^{\prime}, v_{a}^{\prime \prime}$ are actually summands and thus $\mathbb{V}=\left(V_{1}, v_{a}\right)$ is decomposable. This shows that the indecomposable representations of the loop quiver are given by the representations that only have one Jordan block [3]. We later give a more robust argument for this (see Example 2.39).

KI: Every morphism of quiver representations has a kernel, an image, and a cokernel.
Definition 2.26. Let $\phi: \mathbb{V} \rightarrow \mathbb{W}$ be a morphism of quiver representations. We define the kernel of $\phi$ to be $(\operatorname{ker} \phi)_{i}=\operatorname{ker} \phi_{i}$ for all $i \in Q_{0}$. Moreover, the image of $\phi$ is given by $(\operatorname{im} \phi)_{i}=\operatorname{im} \phi_{i}$ and the cokernel of $\phi$ is defined as $(\operatorname{coker} \phi)_{i}=\operatorname{coker} \phi_{i}$. The maps between the vector spaces are given by the induced subspace maps resp. quotient maps.

Lemma 2.27. For every morphism $\phi: \mathbb{V} \rightarrow \mathbb{W}$, $\operatorname{ker} \phi$ is a subrepresentation of $\mathbb{V}, \operatorname{im} \phi$ is a subrepresentation of $\mathbb{W}$ and coker $\phi$ is a representation.
Proof. Notice that for each point $i \in Q_{0}, \operatorname{ker} \phi_{i}$ is a subspace of $V_{i}$ and $\operatorname{im} \phi_{i}$ is a subspace of $W_{i}$. It remains to show that the maps restrict well. For each arrow $a \in Q_{1}$ we have the commutative diagram

and we get that $v_{a}\left(\operatorname{ker} \phi_{i}\right) \subseteq \operatorname{ker} \phi_{j}$, which shows that $\operatorname{ker} \phi$ is a subrepresentation. Similarly, we get that $w_{a}\left(\operatorname{im} \phi_{i}\right)=\phi_{j}\left(\operatorname{im} v_{a}\right) \subseteq \operatorname{im} \phi_{j}$ and thus $\operatorname{im} \phi$ is a subrepresentation. Since $w_{a}\left(\operatorname{im} \phi_{i}\right) \subseteq \operatorname{im} \phi_{j}$, we have a well-defined map $\tilde{w}_{a}: W_{i} / \operatorname{im} \phi_{i}=\operatorname{coker} \phi_{i} \rightarrow \operatorname{coker} \phi_{j}=W_{j} / \operatorname{im} \phi_{j}$, which is given by $\tilde{w}_{a}\left(w+\operatorname{im} \phi_{i}\right)=w_{a}(w)+\operatorname{im} \phi_{j}$. This shows that coker $\phi$ is a representation.

Remark 2.28. A morphism $\phi$ is a monomorphism if and only if $\operatorname{ker} \phi=0$, an epimorphism if and only if coker $\phi=0$, and an isomorphism if and only if $\phi$ is both a monomorphism and an epimorphism.

Since both morphisms and the direct sum of quiver representations are defined pointwise, there are many properties of vector spaces that carry over to quiver representations. E.g. if a morphism $\phi: \mathbb{V} \rightarrow \mathbb{W}$ is a monomorphism, then for each $i \in Q_{0}$ we have $\operatorname{dim} V_{i} \leq \operatorname{dim} W_{i}$. One important property that does not carry over is semisimplicity: while each subspace $W \subseteq V$ is a summand (i.e. there exists a subspace $W^{\perp}$ s.t. $W \oplus W^{\perp}=V$ ), not all subrepresentations of a given representation $\mathbb{V}$ are summands of $\mathbb{V}$. The lack of semisimplicity makes the classification of quiver representations a lot harder since it does not suffice to find all representations with no proper subrepresentations (which, for example, is enough to classify the complex finite-dimensional group representations of a finite group).

Example 2.29. We consider representations of the $L_{2}$-quiver from Example 2.12. Let $\mathbb{V}=k \xrightarrow{\mathbb{1}} k$ and $\mathbb{W}=0 \xrightarrow{0} k$ be two such representations. Then $\mathbb{W}$ is a subrepresentation of $\mathbb{V}$, but it is not a summand. Indeed, if $\mathbb{U}$ is a subrepresentation of $\mathbb{V}$ such that $\mathbb{V}=\mathbb{U} \oplus \mathbb{W}$, then $\mathbb{U}=k \rightarrow 0$, which can only be the zero map. But then $\mathbb{U} \oplus \mathbb{W}=k \xrightarrow{0} k$ which is not isomorphic to $\mathbb{V}$.

### 2.2.2 Properties of $\operatorname{Hom}(\mathbb{V}, \mathbb{W})$

We now decompose $\operatorname{Hom}(\mathbb{V}, \mathbb{W})$ and look at morphisms between indecomposable representations.
Lemma 2.30. [4, Lemma 2.1.1] Let $\mathbb{V}=\mathbb{V}_{1} \oplus \ldots \oplus \mathbb{V}_{r}$ and $\mathbb{W}=\mathbb{W}_{1} \oplus \ldots \oplus \mathbb{W}_{s}$ be two direct sums of representations (of the same quiver). The decomposition of the representations induces vector space decompositions

$$
\bigoplus_{i=1}^{r} \operatorname{Hom}\left(\mathbb{V}_{i}, \mathbb{W}\right) \stackrel{(1)}{\cong} \operatorname{Hom}(\mathbb{V}, \mathbb{W}) \stackrel{(2)}{\cong} \bigoplus_{j=1}^{s} \operatorname{Hom}\left(\mathbb{V}, \mathbb{W}_{j}\right)
$$

We associate the map $\phi \in \operatorname{Hom}(\mathbb{V}, \mathbb{W})$ with the collection of maps $\left(\phi_{i}\right)_{1 \leq i \leq r}$ where $\phi_{i} \in \operatorname{Hom}\left(\mathbb{V}_{i}, \mathbb{W}\right)$ and also with $\left(\phi_{j}\right)_{1 \leq j \leq s}$ where $\phi_{j} \in \operatorname{Hom}\left(\mathbb{V}, \mathbb{W}_{j}\right)$ through the isomorphisms in (1) and (2).

Proof. 1. Let $\iota_{i}: \mathbb{V}_{i} \rightarrow \mathbb{V}$ and $\pi_{i}: \mathbb{V} \rightarrow \mathbb{V}_{i}$ be the maps from the decomposition of $\mathbb{V}$. For $1 \leq i \leq r$ we define the map $\tilde{\iota_{i}}=\pi_{i}^{*}: \operatorname{Hom}(\mathbb{V}, \mathbb{W}) \rightarrow \operatorname{Hom}(\mathbb{V}, \mathbb{W}), \phi \mapsto \phi \pi_{i}$ and the map $\tilde{\pi}_{i}=\iota_{i}^{*}: \operatorname{Hom}(\mathbb{V}, \mathbb{W}) \rightarrow \operatorname{Hom}(\mathbb{V}, \mathbb{W}), \phi \mapsto \phi \iota_{i}$. For $\phi \in \operatorname{Hom}(\mathbb{V}, \mathbb{W})$ we have

$$
\left(\sum_{i=1}^{r} \tilde{\iota_{i}} \tilde{\pi}_{i}\right) \phi=\sum_{i=1}^{r} \pi_{i}^{*} \iota_{i}^{*} \phi=\sum_{i=1}^{r} \phi \iota_{i} \pi_{i}=\phi\left(\sum_{i=1}^{r} \iota_{i} \pi_{i}\right)=\phi \mathrm{id}_{\vee}=\phi .
$$

Thus we have $\sum_{i=1}^{r} \tilde{\iota_{i}} \tilde{\pi}_{i}=\operatorname{id}_{\operatorname{Hom}(\mathbb{V}, \mathbb{W})}$. In addition, for $\phi \in \operatorname{Hom}\left(\mathbb{V}_{i}, \mathbb{W}\right)$ we have

$$
\tilde{\pi}_{i} \tilde{\iota_{i}} \phi=\iota_{i}^{*} \pi_{i}^{*} \phi=\phi \pi_{i} \iota_{i}=\phi \mathrm{id}_{v_{i}}=\phi .
$$

So $\tilde{\pi}_{i} \tilde{\iota}_{i}=\operatorname{id}_{\operatorname{Hom}\left(\mathbb{V}_{i}, W\right)}$. This shows the first decomposition.
2. Let $\iota_{j}: \mathbb{W}_{j} \rightarrow \mathbb{W}$ and $\pi_{j}: \mathbb{W} \rightarrow \mathbb{W}_{j}$ be the maps from the decomposition of $\mathbb{W}$. For $1 \leq j \leq s$ we define the map $\tilde{\iota_{j}}=\iota_{j}^{*}: \operatorname{Hom}\left(\mathbb{V}, \mathbb{W}_{j}\right) \rightarrow \operatorname{Hom}(\mathbb{V}, \mathbb{W}), \phi \mapsto \iota_{j} \phi$ and we define the map $\tilde{\pi_{j}}=\pi_{j}^{*}: \operatorname{Hom}(\mathbb{V}, \mathbb{W}) \rightarrow \operatorname{Hom}\left(\mathbb{V}, \mathbb{W}_{j}\right), \phi \mapsto \pi_{j} \phi . \operatorname{For} \phi \in \operatorname{Hom}(\mathbb{V}, \mathbb{W})$ we have

$$
\left(\sum_{j=1}^{s} \tilde{\iota_{j}} \tilde{\pi}_{j}\right) \phi=\sum_{j=1}^{s} \iota_{j}^{*} \pi_{j}^{*} \phi=\sum_{j=1}^{s} \iota_{j} \pi_{j} \phi=\left(\sum_{j=1}^{s} \iota_{j} \pi_{j}\right) \phi=\operatorname{id}_{\mathbb{W}} \phi=\phi .
$$

Therefore, $\sum_{j=1}^{s} \tilde{\iota_{j}} \tilde{\pi_{j}}=\operatorname{id}_{\operatorname{Hom}(\mathbb{V}, \mathbb{W})}$. In addition, for $\phi \in \operatorname{Hom}\left(\mathbb{V}, \mathbb{W}_{j}\right)$ we have

$$
\tilde{\pi_{j} \tilde{\nu}_{j}} \phi=\pi_{j, *} \iota_{j, *} \phi=\pi_{j} \iota_{j} \phi=\operatorname{id}_{\mathbb{W}_{j}} \phi=\phi .
$$

So $\tilde{\pi_{j}} \tilde{\iota_{j}}=\operatorname{id}_{\operatorname{Hom}\left(\mathbb{V}, W_{j}\right)}$. This shows the second decomposition.
Lemma 2.31. [5, Lemma p. 112] Let $\mathbb{U}, \mathbb{V}, \mathbb{W}$ be indecomposable representations (of the same quiver) and let $\phi: \mathbb{U} \rightarrow \mathbb{V}$ and $\psi: \mathbb{V} \rightarrow \mathbb{W}$ be morphisms such that $\psi \phi$ is an isomorphism. Then both $\phi$ and $\psi$ are isomorphisms. In particular, the composition of a finite number of morphisms between indecomposable representations is invertible if and only if each morphism is invertible.

Proof. Let $\mathbb{U}, \mathbb{V}, \mathbb{W}$ and $\phi: \mathbb{U} \rightarrow \mathbb{V}, \psi: \mathbb{V} \rightarrow \mathbb{W}$ be as in Lemma 2.31 (i.e. $\psi \phi$ is an isomorphism). We show that both $\phi$ and $\psi$ are isomorphisms. We set $\sigma=(\psi \phi)^{-1} \psi: \mathbb{V} \rightarrow \mathbb{U}$ and observe that $\sigma \phi=\mathrm{id}_{U}$ and thus $\phi$ is injective. We notice that $\mathbb{U} \cong \operatorname{im} \phi$ and we claim that $\mathbb{V}=\operatorname{im} \phi \oplus \operatorname{ker} \sigma$. Indeed, for any $v \in \mathbb{V}$ (meaning $v \in V_{i}$ for some $i \in Q_{0}$ ), we have that $\sigma(v-\phi \sigma(v))=\sigma(v)-\mathrm{id}_{U} \sigma(v)=0$ and thus $v-\phi \sigma(v) \in \operatorname{ker} \sigma$. Also if $\phi(u) \in \operatorname{ker} \sigma$, then $u=\sigma \phi(u)=0$ and thus $\operatorname{im} \phi \cap \operatorname{ker} \sigma=0$. Thus we get the decomposition

$$
v=\phi \sigma(v)+(v-\phi \sigma(v)) \in \operatorname{im} \phi+\operatorname{ker} \sigma .
$$

This decomposition gives rise to the morphisms $\pi_{\phi}=\phi \sigma: \mathbb{V} \rightarrow \operatorname{im} \phi$ and $\pi_{\sigma}=\mathrm{id}_{\vee}-\phi \sigma: \mathbb{V} \rightarrow \operatorname{ker} \sigma$. We denote by $\iota_{\phi}: \operatorname{im} \phi \rightarrow \mathbb{V}$ and $\iota_{\sigma}: \operatorname{ker} \sigma \rightarrow \mathbb{V}$ the (pointwise) inclusion morphisms. Now we have $\pi_{\phi} \iota_{\phi}=\mathrm{id}_{\mathrm{im} \phi}$ and $\pi_{\sigma} \iota_{\sigma}=\mathrm{id}_{\mathrm{ker} \sigma}$. Further, this decomposition yields that $\iota_{\phi} \pi_{\phi}=\phi \sigma$ and $\iota_{\sigma} \pi_{\sigma}=\mathrm{id}_{\vee}-\phi \sigma$ and thus we get $\iota_{\phi} \pi_{\phi}+\iota_{\sigma} \pi_{\sigma}=\mathrm{id}_{\mathbb{V}}$. This shows that $\mathbb{V}=\operatorname{im} \phi \oplus \operatorname{ker} \sigma$. Since $\mathbb{V}$ is indecomposable and $\operatorname{im} \phi \cong \mathbb{U} \neq 0$ we get that $\operatorname{ker} \sigma=0$ and therefore $\operatorname{im} \phi=\mathbb{V}$ and thus both $\phi$ and $\psi$ are isomorphisms. The second claim follows by induction on the number of morphisms.

Remark 2.32. From Lemma 2.30 and Lemma 2.31 we have learned a lot about $\operatorname{Hom}(\mathbb{V}, \mathbb{W})$. The first lemma tells us that it suffices to consider $\operatorname{Hom}(\mathbb{V}, \mathbb{W})$ for indecomposable representations. Indeed, each morphism $\phi \in \operatorname{Hom}(\mathbb{V}, \mathbb{W})$ can be uniquely written as a block matrix $\left(\phi_{i, j}\right)$, where $\phi_{i, j}: \mathbb{V}_{j} \rightarrow \mathbb{W}_{i}$ are morphisms between indecomposable representations. Such morphisms (between indecomposable representations) are characterized by the second lemma.

### 2.3 The Krull-Remak-Schmidt Theorem

Towards our goal of classifying quiver representations, we can ask if a quiver representation can always be decomposed into a direct sum of finitely many indecomposable representations. We call such a decomposition a Remak decomposition and it turns out that such a decomposition always exists and it is unique up to isomorphism and permutation of the factors in the direct sum. This section is primarily based on the lecture notes 'Representations of quivers via reflection functors' written by Henning Krause [4].

Theorem 2.33 (Krull, Remak, Schmidt). Let $Q$ be a finite quiver. Then for any $\mathbb{V} \in \operatorname{rep}_{\mathrm{k}}(Q)$ there are indecomposable representations $\mathbb{V}_{1}, \ldots, \mathbb{V}_{r}$ such that $\mathbb{V} \cong \mathbb{V}_{1} \oplus \ldots \oplus \mathbb{V}_{r}$. Moreover, for indecomposable representations $\mathbb{W}_{1}, \ldots, \mathbb{W}_{s}$ such that $\mathbb{V} \cong \mathbb{W}_{1} \oplus \ldots \oplus \mathbb{W}_{s}, r=s$ and there is $a$ permutation $\sigma$ such that $\mathbb{V}_{i} \cong \mathbb{W}_{\sigma(i)}$ for $1 \leq i \leq r$.

Before we can prove Theorem 2.33, we observe a connection between the decomposability of a representation $\mathbb{V}$ and the structure and property of the ring of its endomorphisms $\operatorname{End}(\mathbb{V})$.

Remark 2.34. Note that the composition of morphisms gives a ring structure on $\operatorname{End}(\mathbb{V})$, which is the reason why we generally refer to $\operatorname{End}(\mathbb{V})$ as a ring. $\operatorname{Hom}(\mathbb{V}, \mathbb{W})$ on the other hand does not have a ring structure (with respect to the composition of morphisms).

Lemma 2.35 (Fitting's Lemma). [5, Lemma p. 113] Let $\vee$ be a finite-dimensional representation and let $\phi$ be an endomorphism in $\operatorname{End}(\mathbb{V})$.

1. For large enough $r$, we have $\mathbb{V}=\operatorname{im} \phi^{r} \oplus \operatorname{ker} \phi^{r}$.
2. If $\mathbb{V}$ is indecomposable, then $\phi$ is either an automorphism or nilpotent.

Proof. 1. Because $\mathbb{V}$ is finite-dimensional, we can choose r large enough, s.t. $\operatorname{im} \phi^{r}=\operatorname{im} \phi^{r+1}$. This is possible since all the maps $\phi_{i}$ are endomorphisms on finite-dimensional vector spaces and we can deduce the result pointwise: we have $\operatorname{im} \phi_{i}^{r_{i}+1} \subseteq \operatorname{im} \phi_{i}^{r_{i}} \subseteq \ldots \subseteq \operatorname{im} \phi_{i}^{2} \subseteq \operatorname{im} \phi_{i}$. If all the inclusions were strict, then $\operatorname{dim}\left(\operatorname{im} \phi_{i}\right)>\operatorname{dim}\left(\operatorname{im} \overline{\phi_{i}^{2}}\right)>\ldots$ is an infinite, strictly decreasing sequence of non-negative numbers, which is impossible since $\operatorname{dim}\left(\operatorname{im} \phi_{i}\right)<\infty$. This also shows that w.l.o.g. $r_{i} \leq \operatorname{dim}\left(\phi_{i}\right)$, where $\operatorname{dim}\left(\phi_{i}\right)$ is the dimension of the vector space on which $\phi_{i}$ operates. Thus we get the desired result for $r \geq \operatorname{dim} \mathbb{V}$. For the same $r$ it also holds that $\operatorname{ker} \phi^{r+1}=\operatorname{ker} \phi^{r}$. Thus $\phi^{r}: \operatorname{im} \phi^{r} \rightarrow \operatorname{im} \phi^{2 r}$ is an isomorphism and we denote by $\psi$ its inverse. Let $\iota_{1}: \operatorname{im} \phi^{r} \rightarrow \mathbb{V}$ and $\iota_{2}: \operatorname{ker} \phi^{r} \rightarrow \mathbb{V}$ denote the inclusions. We put

$$
\pi_{1}=\psi \phi^{r}: \mathbb{V} \rightarrow \operatorname{im} \phi^{r} \text { and } \pi_{2}=\operatorname{id}_{\mathbb{V}}-\psi \phi^{r}: \mathbb{V} \rightarrow \operatorname{ker} \phi^{r}
$$

This is well-defined since $\phi^{r} \pi_{2}=\phi^{r}-\phi^{r} \psi \phi^{r}=\phi^{r}-\mathrm{id}_{\mathrm{im}} \phi^{r} \phi^{r}=0$. Then $\iota_{1} \pi_{1}+\iota_{2} \pi_{2}=\mathrm{id} \vee$ and $\pi_{1} \iota_{1}=\mathrm{id}_{\mathrm{im} \phi^{r}}, \pi_{2} \iota_{2}=\mathrm{id}_{\mathrm{ker} \phi^{r}}$. Thus im $\phi^{r}$ and ker $\phi^{r}$ are summands and by the definition of the direct sum we have: $\mathbb{V}=\operatorname{im} \phi^{r} \oplus \operatorname{ker} \phi^{r}$.
2. If $\mathbb{V}$ is indecomposable, then one of the factors in the decomposition in part (1) needs to be 0 . If $\operatorname{ker} \phi^{r}=0$, then $\phi$ is an automorphism and if $\operatorname{im} \phi^{r}=0$, then $\phi$ is nilpotent.

Remark 2.36. Notice that $\vee$ needs to be finite-dimensional. Otherwise the decomposition in part (1) is not guaranteed, since there may not be any such $r$.

Definition 2.37. A ring is called local if the sum of two non-units is again a non-unit.
Proposition 2.38. [5, Proposition 3.1] A finite-dimensional representation $\mathbb{V}$ is indecomposable if and only if $\operatorname{End}(\mathbb{V})$ is local.
Proof. $\Longrightarrow$ : Let $\mathbb{V}$ be indecomposable and let $\phi, \phi^{\prime} \in \operatorname{End}(\mathbb{V})$ such that $\phi+\phi^{\prime}$ is invertible with inverse $\rho$. If $\phi$ is non-invertible, so is $\rho \phi$ and thus by Lemma 2.35, $\rho \phi$ is nilpotent, say $(\rho \phi)^{r}=0$. The summation formula for geometric series yields:

$$
\left(\mathrm{id}_{\vee}-\rho \phi\right)\left(\mathrm{id}_{\vee}+\rho \phi+\ldots+(\rho \phi)^{r-1}\right)=\mathrm{id}_{\vee}-(\rho \phi)^{r}=\mathrm{id}_{\vee} .
$$

Therefore $\rho \phi^{\prime}=\mathrm{id}_{\vee}-\rho \phi$ is invertible and thus $\phi^{\prime}$ is invertible, which shows that $\operatorname{End}(\mathbb{V})$ is local. $\Longleftarrow:$ If $\mathbb{V}=\mathbb{U} \oplus \mathbb{W}$ is decomposable, then the endomorphisms $\iota_{U} \pi_{U}, \iota_{\mathbb{W}} \pi_{\mathbb{W}}$ have image $\mathbb{U}$ respectively $\overline{\mathbb{W}}$ and thus are not invertible. But $\iota_{\cup} \pi_{\cup}+\iota_{\mathbb{W}} \pi_{\mathbb{W}}=\mathrm{id} \mathbb{V}_{\mathbb{V}}$ is invertible and therefore $\operatorname{End}(\mathbb{V})$ is not local.

Example 2.39. In Example 2. 25 we have stated that the finite-dimensional indecomposable representations of the loop quiver are given by the endomorphisms $v_{a}$, which only have one Jordan block. Using Proposition 2.38 we can now prove that every indecomposable representation consists of a single Jordan block (proof by contrapositive). Indeed if $v_{a}$ has two or more Jordan blocks, then we can define $\phi_{1}=\left(\begin{array}{cc}\text { id } & 0 \\ 0 & 0\end{array}\right)$ to be the identity on the subspace belonging to the first Jordan block. Analogously, let $\phi_{2}=\left(\begin{array}{ll}0 & 0 \\ 0 & \text { id }\end{array}\right)$ be the identity on the subspace belonging to the other Jordan blocks. Then $\phi_{1}, \phi_{2}$ are both non-invertible, but $\phi_{1}+\phi_{2}=\mathrm{id}_{\vee}$ is invertible. Thus $\operatorname{End}(\mathbb{V})$ is not local and therefore $\mathbb{V}$ is decomposable.

We use Proposition 2.38 to consider a subspace of the vector space $\operatorname{Hom}(\mathbb{V}, \mathbb{W})$. To do this, we use the Lemmas 2.30 and 2.31 from Section 2.2.

Definition 2.40. The radical of two representations $\mathbb{V}, \mathbb{W}$ is defined to be

$$
\operatorname{Rad}(\mathbb{V}, \mathbb{W})=\left\{\phi \in \operatorname{Hom}(\mathbb{V}, \mathbb{W}) \left\lvert\, \begin{array}{c}
\text { for every } \mathbb{U} \text { indecomposable, } \mathbb{U} \xrightarrow{\sigma} \mathbb{V} \text { and } \mathbb{W} \xrightarrow{\tau} \mathbb{U}, \\
\tau \phi \sigma \text { is non-invertible }
\end{array}\right.\right\}
$$

A morphism $\phi \in \operatorname{Rad}(\mathbb{V}, \mathbb{W})$ is called radical.
Lemma 2.41. Let $\mathbb{V}, \mathbb{W}$ be two representations.

1. $\operatorname{Rad}(\mathbb{V}, \mathbb{W})$ is a subspace of $\operatorname{Hom}(\mathbb{V}, \mathbb{W})$.
2. $\operatorname{Rad}\left(\mathbb{V}, \mathbb{W}_{1} \oplus \mathbb{W}_{2}\right) \cong \operatorname{Rad}\left(\mathbb{V}, \mathbb{W}_{1}\right) \oplus \operatorname{Rad}\left(\mathbb{V}, \mathbb{W}_{2}\right)$.
3. $\operatorname{Rad}\left(\mathbb{V}_{1} \oplus \mathbb{V}_{2}, \mathbb{W}\right) \cong \operatorname{Rad}\left(\mathbb{V}_{1}, \mathbb{W}\right) \oplus \operatorname{Rad}\left(\mathbb{V}_{2}, \mathbb{W}\right)$.
4. If $\mathbb{V}, \mathbb{W}$ are indecomposable, then $\operatorname{Hom}(\mathbb{V}, \mathbb{W}) \backslash \operatorname{Rad}(\mathbb{V}, \mathbb{W})$ is the set of isomorphisms $\mathbb{V} \rightarrow \mathbb{W}$.

Proof. 1. Let $\alpha, \beta \in k$ and let $\phi_{1}, \phi_{2} \in \operatorname{Rad}(\mathbb{V}, \mathbb{W})$ and let $\mathbb{U} \xrightarrow{\sigma} \mathbb{V}$ and $\mathbb{W} \xrightarrow{\tau} \mathbb{U}$ for $\mathbb{U}$ indecomposable as in the definition of $\operatorname{Rad}(\mathbb{V}, \mathbb{W})$. Then $\alpha \cdot \tau \phi_{1} \sigma, \beta \cdot \tau \phi_{2} \sigma$ are non-invertible and thus $\tau\left(\alpha \cdot \phi_{1}+\beta \cdot \phi_{2}\right) \sigma=\alpha \cdot \tau \phi_{1} \sigma+\beta \cdot \tau \phi_{2} \sigma$ is non-invertible since $\operatorname{End}(\mathbb{U})$ is a local ring. Thus $\alpha \cdot \phi_{1}+\beta \cdot \phi_{2} \in \operatorname{Rad}(\mathbb{V}, \mathbb{W})$.
2. Let $\mathbb{W}=\mathbb{W}_{1} \oplus \mathbb{W}_{2}$ and let $\phi=\left(\phi_{1}, \phi_{2}\right) \in \operatorname{Hom}\left(\mathbb{V}, \mathbb{W}_{1}\right) \oplus \operatorname{Hom}\left(\mathbb{V}, \mathbb{W}_{2}\right) \cong \operatorname{Hom}(\mathbb{V}, \mathbb{W})$. Let $\mathbb{U}$ be indecomposable and let $\sigma \in \operatorname{Hom}(\mathbb{U}, \mathbb{V})$ and $\tau=\left(\tau_{1}, \tau_{2}\right) \in \operatorname{Hom}\left(\mathbb{W}_{1}, \mathbb{U}\right) \oplus \operatorname{Hom}\left(\mathbb{W}_{2}, \mathbb{U}\right)$. Then $\tau \phi \sigma=\left(\tau_{1}, \tau_{2}\right)\left(\phi_{1}, \phi_{2}\right) \sigma=\tau_{1} \phi_{1} \sigma+\tau_{2} \phi_{2} \sigma$. This results directly from the decomposition of $\operatorname{Hom}\left(\mathbb{V}, \mathbb{W}_{1} \oplus \mathbb{W}_{2}\right)$ given in Lemma 2.30. If $\phi_{i} \in \operatorname{Rad}\left(\mathbb{V}, \mathbb{W}_{i}\right)$ for $i \in\{1,2\}$, then $\tau_{1} \phi_{1} \sigma$ and $\tau_{2} \phi_{2} \sigma$ are non-invertible and thus $\tau \phi \sigma=\tau_{1} \phi_{1} \sigma+\tau_{2} \phi_{2} \sigma$ is non-invertible since $\operatorname{End}(\mathbb{U})$ is local (since $\mathbb{U}$ is indecomposable). Thus $\phi \in \operatorname{Rad}(\mathbb{V}, \mathbb{W})$. On the other hand, if we fix $i \in\{1,2\}$ and if $\phi \in \operatorname{Rad}(\mathbb{V}, \mathbb{W})$, then we can choose $\tau=\left(\tau_{1}, \tau_{2}\right)$ such that $\tau_{j}=0$ for $j \neq i$ (and $\tau_{i}$ arbitrary) and we get that $\tau_{i} \phi_{i} \sigma=\tau \phi \sigma$ is non-invertible. Therefore $\phi_{i} \in \operatorname{Rad}\left(\mathbb{V}, \mathbb{W}_{i}\right)$.
3. Let $\mathbb{V}=\mathbb{V}_{1} \oplus \mathbb{V}_{2}$ and let $\phi=\left(\phi_{1}, \phi_{2}\right) \in \operatorname{Hom}\left(\mathbb{V}_{1}, \mathbb{W}\right) \oplus \operatorname{Hom}\left(\mathbb{V}_{2}, \mathbb{W}\right) \cong \operatorname{Hom}(\mathbb{V}, \mathbb{W})$. Let $\mathbb{U}$ be indecomposable and let $\sigma=\left(\sigma_{1}, \sigma_{2}\right) \in \operatorname{Hom}\left(\mathbb{U}, \mathbb{V}_{1}\right) \oplus \operatorname{Hom}\left(\mathbb{U}, \mathbb{V}_{2}\right) \cong \operatorname{Hom}(\mathbb{U}, \mathbb{V})$ and $\tau \in \operatorname{Hom}(\mathbb{W}, \mathbb{U})$. Then $\tau \phi \sigma=\tau\left(\phi_{1}, \phi_{2}\right)\left(\sigma_{1}, \sigma_{2}\right)=\tau_{1} \phi_{1} \sigma+\tau_{2} \phi_{2} \sigma$. This results directly from the decomposition of $\operatorname{Hom}\left(\mathbb{V}_{1} \oplus \mathbb{V}_{2}, \mathbb{W}\right)$ given in Lemma 2.30. If $\phi_{i} \in \operatorname{Rad}\left(\mathbb{V}_{i}, \mathbb{W}\right)$ for $i \in\{1,2\}$, then $\tau \phi_{1} \sigma_{1}$ and $\tau \phi_{2} \sigma_{2}$ are non-invertible and thus $\tau \phi \sigma=\tau \phi_{1} \sigma_{1}+\tau \phi_{2} \sigma_{2}$ is non-invertible since $\operatorname{End}(\mathbb{U})$ is local (since $\mathbb{U}$ is indecomposable). Thus $\phi \in \operatorname{Rad}(\mathbb{V}, \mathbb{W})$. On the other hand, if we fix $i \in\{1,2\}$ and if $\phi \in \operatorname{Rad}(\mathbb{V}, \mathbb{W})$, then we can choose $\sigma=\left(\sigma_{1}, \sigma_{2}\right)$ such that $\sigma_{j}=0$ for $j \neq i$ (and $\sigma_{i}$ arbitrary) and we get that $\tau \phi_{i} \sigma_{i}=\tau \phi \sigma$ is non-invertible. Therefore $\phi_{i} \in \operatorname{Rad}\left(\mathbb{V}_{i}, \mathbb{W}\right)$.
4. Let $\phi \in \operatorname{Hom}(\mathbb{V}, \mathbb{W})$ be an isomorphism. Then we can choose $\mathbb{U}=\mathbb{V}, \sigma=\operatorname{id} \mathbb{V}, \tau=\phi^{-1}$ and see that $\tau \phi \sigma=\mathrm{id}_{\mathbb{\vee}}$ is invertible and thus $\phi \notin \operatorname{Rad}(\mathbb{V}, \mathbb{W})$. Now let $\phi \in \operatorname{Hom}(\mathbb{V}, \mathbb{W}) \backslash \operatorname{Rad}(\mathbb{V}, \mathbb{W})$. W.l.o.g we can assume that $\phi$ is non-invertible (else we are done). Now choose $\mathbb{U}$ indecomposable and $\sigma \in \operatorname{Hom}(\mathbb{U}, \mathbb{V})$ and $\tau \in \operatorname{Hom}(\mathbb{W}, \mathbb{U})$ such that $\tau \phi \sigma$ is invertible. Since $\mathbb{V}, \mathbb{W}$ are indecomposable, this is impossible. Indeed, since $\phi$ is non-invertible, it follows from Lemma 2.31 that $\tau \phi \sigma$ is non-invertible.

Remark 2.42. If $\phi: \mathbb{V} \rightarrow \mathbb{W}$ is radical and $\phi=\left(\phi_{i, j}\right)$ is written as a block matrix of morphisms between indecomposable representations (see Lemma 2.30), then every such block $\phi_{i, j}$ is non-invertible. Indeed if $\phi_{i, j}: \mathbb{V}_{j} \rightarrow \mathbb{W}_{i}$ is invertible, then we can choose $\mathbb{U}=\mathbb{W}_{i}$ indecomposable, $\tau: \mathbb{W} \rightarrow \mathbb{W}_{i}$ the projection onto the summand $\mathbb{W}_{i}$ and $\sigma=\iota_{\mathbb{V}_{j}}\left(\phi_{i, j}\right)^{-1}: \mathbb{W}_{i} \rightarrow \mathbb{V}_{j} \rightarrow \mathbb{V}$ where $\iota_{\mathbb{V}_{j}}$ is the inclusion of the summand $\mathbb{V}_{j}$ into $\mathbb{V}$. But then $\tau \phi \sigma=\phi_{i, j}$ is invertible, which is a contradiction to the choice of $\phi$.

Proof of Theorem 2.33. By induction on $\operatorname{dim} \mathbb{V}$ we can show that a Remak decomposition exists: for $\operatorname{dim} \mathbb{V}=0$ all vector spaces are trivial and thus indecomposable. Therefore $\mathbb{V}$ is already indecomposable. For the induction step, let $\operatorname{dim} \mathbb{V}=n \geq 1$. If $\mathbb{V}$ is decomposable, then there exist proper summands $\mathbb{U}, \mathbb{W}$. Since both $\operatorname{dim} \mathbb{U}, \operatorname{dim} \mathbb{W}<n$ we know by the induction hypothesis that they each have a Remak decomposition. The direct sum of these decompositions gives a Remak decomposition for $\mathbb{V}$. For the proof of uniqueness, we order isomorphic summands together, i.e. let $\mathbb{V}=\mathbb{V}_{1}^{a_{1}} \oplus \ldots \oplus \mathbb{V}_{r}^{a_{r}}$ be a direct sum decomposition of $\mathbb{V}$ where the $\mathbb{V}_{i}$ are pairwise non-isomorphic representations and $a_{i} \geq 1$ for all $1 \leq i \leq r$. For $\mathbb{W}$ indecomposable, we can consider the number

$$
\frac{\operatorname{dim} \operatorname{Hom}(\mathbb{V}, \mathbb{W})-\operatorname{dim} \operatorname{Rad}(\mathbb{V}, \mathbb{W})}{\operatorname{dim} \operatorname{Hom}(\mathbb{W}, \mathbb{W})-\operatorname{dim} \operatorname{Rad}(\mathbb{W}, \mathbb{W})}
$$

From Lemma 2.30 and Lemma 2.41 we know that for $\mathbb{W}=\mathbb{V}_{i}$ this number is equal to:

$$
\begin{aligned}
\frac{\sum_{j=1}^{r}\left(\operatorname{dim} \operatorname{Hom}\left(\mathbb{V}_{j}^{a_{j}}, \mathbb{V}_{i}\right)-\operatorname{dim} \operatorname{Rad}\left(\mathbb{V}_{j}^{a_{j}}, \mathbb{V}_{i}\right)\right)}{\operatorname{dim} \operatorname{Hom}\left(\mathbb{V}_{i}, \mathbb{V}_{i}\right)-\operatorname{dim} \operatorname{Rad}\left(\mathbb{V}_{i}, \mathbb{V}_{i}\right)} & =\frac{\sum_{j=1}^{r} a_{j} \cdot \operatorname{dim}\left(\operatorname{Hom}\left(\mathbb{V}_{j}, \mathbb{V}_{i}\right) \backslash \operatorname{Rad}\left(\mathbb{V}_{j}, \mathbb{V}_{i}\right)\right)}{\operatorname{dim}\left(\operatorname{Hom}\left(\mathbb{V}_{i}, \mathbb{V}_{i}\right) \backslash \operatorname{Rad}\left(\mathbb{V}_{i}, \mathbb{V}_{i}\right)\right)} \\
& =\frac{a_{i} \cdot \operatorname{dim}\left(\operatorname{Hom}\left(\mathbb{V}_{i}, \mathbb{V}_{i}\right) \backslash \operatorname{Rad}\left(\mathbb{V}_{i}, \mathbb{V}_{i}\right)\right)}{\operatorname{dim}\left(\operatorname{Hom}\left(\mathbb{V}_{i}, \mathbb{V}_{i}\right) \backslash \operatorname{Rad}\left(\mathbb{V}_{i}, \mathbb{V}_{i}\right)\right)}=a_{i}
\end{aligned}
$$

Here we used part (4) from Lemma 2.41 to see that $\operatorname{dim} \operatorname{Hom}\left(\mathbb{V}_{j}, \mathbb{V}_{i}\right)-\operatorname{dim} \operatorname{Rad}\left(\mathbb{V}_{j}, \mathbb{V}_{i}\right)=0$ for $j \neq i$ since $\mathbb{V}_{j}$ and $\mathbb{V}_{i}$ are non-isomorphic. We notice that this number is independent of the decomposition of $\mathbb{V}$ and thus the decomposition is unique up to isomorphisms and reordering.

Remark 2.43. In this section we restricted our focus to finite quivers and finite-dimensional representations. Moreover, we assumed that the quivers are connected. Those restrictions are sensible. Indeed, if a quiver $Q$ is the disjoint union of two quivers $Q^{\prime}$ and $Q^{\prime \prime}$, then any representation of $Q$ is the same as a pair of representations, one of $Q^{\prime}$ and one of $Q^{\prime \prime}$, and any morphism acts on each component of the representation separately. Therefore, $\operatorname{rep}_{\mathrm{k}}(Q)$ is isomorphic to the product category $\operatorname{rep}_{\mathrm{k}}\left(Q^{\prime}\right) \times \operatorname{rep}_{\mathrm{k}}\left(Q^{\prime \prime}\right)$. From the viewpoint of topological data analysis, restricting to finite quivers and finite-dimensional representations is somewhat justified since one looks at persistence modules, which generally are finite-dimensional representations of finite quivers. This is done in Chapter 5.

## Chapter 3

## Gabriel's Theorem

In this chapter, we consider quivers that have a finite number of indecomposable representations.
Definition 3.1. Let $Q$ be a quiver and let $k$ be a field. The quiver $Q$ is of finite-type if it has finitely many isomorphism classes of indecomposable finite-dimensional representations.

This leads to the question of whether there are quivers of finite-type and if so, which quivers are of finite-type. The answer to this question is given by Gabriel's Theorem. Proving Gabriel's Theorem is our main goal in this chapter. Sections 3.1 to 3.4 introduce important results, which help understand quiver representations. These sections are written with the goal of proving Gabriel's Theorem. Therefore, we generally only introduce what is needed to prove Gabriel's Theorem.

Theorem 3.2 (Gabriel, version 1). Let $Q$ be a finite connected quiver and let $k$ be a field. Then $Q$ is of finite type if and only if $Q$ is a Dynkin quiver.

The Krull-Remak-Schmidt Theorem guarantees that every finite-dimensional quiver representation admits a unique decomposition into indecomposable summands. Therefore, for a quiver of finite-type, it suffices to characterize the finitely many indecomposable representations. Using the indecomposable representations, one can then classify all quiver representations. The general outline of this chapter follows the book 'Persistence Theory: From Quiver Representations to Data Analysis' written by Steve Y. Oudot [1]. Some of the more technical details and proofs are based on the book 'Finite dimensional algebras and quantum groups' written by Bangming Deng [2].

Remark 3.3. Note that Gabriel's Theorem is true for any field $k$. In addition, the classification of finite-type quivers happens on the level of the underlying graph, thus the orientations of the arrows do not matter.

### 3.1 Dynkin and Euclidean Diagrams

In this section, we introduce the Tits form. Using the Tits form, we then divide all quivers into three distinct classes.

### 3.1.1 Tits Form

For a finite quiver, we look at a special quadratic form on $\mathbb{Z}^{n}$, called the Tits form. We later apply the Tits form to the dimension vectors of finite-dimensional representations.

Definition 3.4. A vector in $\mathbb{Z}^{n}$ is called positive if it belongs to $\mathbb{N}_{0}^{n} \backslash\{0\}$. This means that every coordinate is non-negative and the vector is non-trivial (not all coordinates are 0). We write $x>0$ if $x$ is positive and we write $x>y$ if $x-y$ is positive. For example, the dimension vectors of nontrivial representations are positive. Dually, a vector $x \in \mathbb{Z}^{n}$ is called negative if $-x$ is positive. $A$ vector $x \in \mathbb{Z}^{n}$ is called sincere if $x_{i} \neq 0$ for all $1 \leq i \leq n$. Finally, we denote by $e_{i} \in \mathbb{Z}^{n}$ the $i$-th coordinate vector, i.e. $\left(e_{i}\right)_{j}=\delta_{i j}$ ( $\delta$ here denotes the Kronecker-delta).

Definition 3.5. The Euler form of a finite quiver $Q$ is the bilinear form $\langle\cdot, \cdot\rangle_{Q}: \mathbb{Z}^{n} \times \mathbb{Z}^{n} \rightarrow \mathbb{Z}$, given by:

$$
\langle x, y\rangle_{Q}=\sum_{i \in Q_{0}} x_{i} y_{i}-\sum_{a \in Q_{1}} x_{t_{a}} y_{h_{a}}
$$

The symmetrization of the Euler form is called the symmetric Euler form and is given by:

$$
(x, y)_{Q}=\langle x, y\rangle_{Q}+\langle y, x\rangle_{Q} .
$$

Lemma 3.6. If we view elements in $\mathbb{Z}^{n}$ as column vectors and if the finite number of edges (in $\bar{Q}$ ) joining the vertices $i$ and $j$ is denoted by $d_{i j}=d_{j i}$, then the symmetric Euler form can be expressed as:

$$
(x, y)_{Q}=x^{T} C_{Q} y
$$

where $C_{Q}=\left(c_{i j}\right)_{i, j \in Q_{0}}$ is the symmetric matrix with the entries

$$
c_{i j}= \begin{cases}2-2 \mid\{\text { loops at } i\} \mid=2-2 d_{i i} & \text { if } i=j, \\ -\mid\{\text { arrows between } i \text { and } j\} \mid=-d_{i j} & \text { if } i \neq j .\end{cases}
$$

Proof. We need to check that for all $1 \leq i, j \leq n$ we have: $\left(e_{i}, e_{j}\right)_{Q}=e_{i}^{T} C_{Q} e_{j}=c_{i j}$. For $i=j$ we have: $\left(e_{i}, e_{i}\right)_{Q}=2\left(1 \cdot 1-d_{i i}\right)=2-d d_{i i}$. For $i \neq j$ we have:

$$
\begin{aligned}
\left(e_{i}, e_{j}\right)_{Q} & =\left\langle e_{i}, e_{j}\right\rangle_{Q}+\left\langle e_{j}, e_{i}\right\rangle_{Q}=(0-\mid\{\text { arrows from } i \text { to } j\} \mid)+(0-\mid\{\text { arrows from } j \text { to } i\} \mid) \\
& =-\mid\{\text { arrows between } i \text { and } j\} \mid=-d_{i j} .
\end{aligned}
$$

Example 3.7. We consider the $L_{2}$-quiver


The Euler form of $L_{2}$ is $\langle x, y\rangle_{L_{2}}=x_{1} y_{1}+x_{2} y_{2}-x_{1} y_{2}$. Its symmetric Euler form is given by $(x, y)_{L_{2}}=2 x_{1} y_{1}+2 x_{2} y_{2}-x_{1} y_{2}-x_{2} y_{1}$.

Example 3.8. The Euler form of the loop quiver is $\langle x, y\rangle_{Q}=x_{1} y_{1}-x_{1} y_{1}=0$. Its symmetric Euler form is $(x, y)_{Q}=0$.

Definition 3.9. For a finite quiver $Q$, the Tits form of $Q$ is the quadratic form $q_{Q}: \mathbb{Z}^{n} \rightarrow \mathbb{Z}$ associated with the Euler form

$$
q_{Q}(x)=\langle x, x\rangle_{Q}=\frac{1}{2}(x, x)_{Q}=\frac{1}{2} x^{T} C_{Q} x=\frac{1}{2} \sum_{i, j} c_{i j} x_{i} x_{j}=\frac{1}{2} \sum_{i}\left(2-2 d_{i i}\right) x_{i}^{2}-\sum_{i<j} d_{i j} x_{i} x_{j}
$$

The radical of the quadratic form $q_{Q}$ is the set $\operatorname{rad} q_{Q}=\left\{x \in \mathbb{Z}^{n} \mid(x, \cdot)_{Q}=0\right\}$ and $x \in \operatorname{rad} q_{Q}$ is called a radical vector.

Definition 3.10. Let $q: \mathbb{Z}^{n} \rightarrow \mathbb{Z}$ be a quadratic form.

1. $q$ is called positive definite if $q(x)>0$ for all non-zero $x \in \mathbb{Z}^{n}$.
2. $q$ is called positive semi-definite if $q(x) \geq 0$ for all $x \in \mathbb{Z}^{n}$.
3. $q$ is called indefinite if it takes on both positive and negative values, i.e. there exist $x, y \in \mathbb{Z}^{n}$ such that $q(x)>0, q(y)<0$.

We use the same terminology for quadratic forms $\mathbb{Q}^{n} \rightarrow \mathbb{Q}$ respectively $\mathbb{R}^{n} \rightarrow \mathbb{R}$.

Example 3.11. The Tits forms of the Dynkin quivers (respectively Dynkin graphs) are:

$$
\begin{aligned}
& A_{n}: q_{Q}(x)=\sum_{i \in Q_{0}} x_{i}^{2}-\sum_{a \in Q_{1}} x_{t_{a}} x_{h_{a}}=\sum_{i=1}^{n} x_{i}^{2}-\sum_{i=1}^{n-1} x_{i} x_{i+1}=\frac{1}{2}\left[x_{1}^{2}+x_{n}^{2}+\sum_{i=1}^{n-1}\left(x_{i}-x_{i+1}\right)^{2}\right] \\
& D_{n}: q_{Q}(x)=\sum_{i=1}^{n} x_{i}^{2}-\sum_{i=1}^{n-2} x_{i} x_{i+1}-x_{n-2} x_{n}=\frac{1}{2}\left[x_{1}^{2}+x_{n}^{2}+\left(x_{n-2}-x_{n}\right)^{2}+\sum_{i=1}^{n-2}\left(x_{i}-x_{i+1}\right)^{2}\right] \\
& E_{6}: q_{Q}(x)=\sum_{i=1}^{6} x_{i}^{2}-x_{1} x_{3}-x_{2} x_{4}-x_{3} x_{4}-x_{4} x_{5}-x_{5} x_{6} \\
& E_{7}: q_{Q}(x)=\sum_{i=1}^{7} x_{i}^{2}-x_{1} x_{3}-x_{2} x_{4}-x_{3} x_{4}-x_{4} x_{5}-x_{5} x_{6}-x_{6} x_{7} \\
& E_{8}: q_{Q}(x)=\sum_{i=1}^{8} x_{i}^{2}-x_{1} x_{3}-x_{2} x_{4}-x_{3} x_{4}-x_{4} x_{5}-x_{5} x_{6}-x_{6} x_{7}-x_{7} x_{8}
\end{aligned}
$$

We observe that the Tits forms of the Dynkin graphs $A_{n}$ and $D_{n}$ are positive definite. We will later see that the Tits forms of the other Dynkin graphs are also positive definite, but this is not obvious from our calculations above.

Lemma 3.12. Let $Q$ be a finite quiver. Then neither its symmetric Euler form nor its Tits form depends on the orientations of the arrows in $Q$.
Proof. Let $Q$ be a finite quiver and let $a=\left(t_{a}, h_{a}\right) \in Q_{1}$ be an arrow. If we consider the quiver $\tilde{Q}$ resulting from $Q$ by exchanging $a$ with $\tilde{a}=\left(h_{a}, t_{a}\right)$ (reverse its direction), we get:

$$
\begin{aligned}
(x, y)_{\tilde{Q}} & =\sum_{i \in \tilde{Q}_{0}} x_{i} y_{i}-\sum_{b \in \tilde{Q}_{1}} x_{t_{b}} y_{h_{b}}+\sum_{i \in \tilde{Q}_{0}} y_{i} x_{i}-\sum_{b \in \tilde{Q}_{1}} y_{t_{b}} x_{h_{b}} \\
& =\sum_{i \in Q_{0}} x_{i} y_{i}-\sum_{\substack{b \in Q_{1} \\
b \neq a}} x_{t_{b}} y_{h_{b}}-x_{h_{a}} y_{t_{a}}+\sum_{i \in Q_{0}} y_{i} x_{i}-\sum_{\substack{b \in Q_{1} \\
b \neq a}} y_{t_{b}} x_{h_{b}}-y_{h_{a}} x_{t_{a}} \\
& =\sum_{i \in Q_{0}} x_{i} y_{i}-\sum_{\substack{b \in Q_{1} \\
b \neq a}} x_{t_{b}} y_{h_{b}}-y_{h_{a}} x_{t_{a}}+\sum_{i \in Q_{0}} y_{i} x_{i}-\sum_{\substack{b \in Q_{1} \\
b \neq a}} y_{t_{b}} x_{h_{b}}-x_{h_{a}} y_{t_{a}} \\
& =\sum_{i \in Q_{0}} x_{i} y_{i}-\sum_{b \in Q_{1}} x_{t_{b}} y_{h_{b}}+\sum_{i \in Q_{0}} y_{i} x_{i}-\sum_{b \in Q_{1}} y_{t_{b}} x_{h_{b}}=(x, y)_{Q} .
\end{aligned}
$$

The result for the Tits form is an immediate consequence: $q_{Q}(x)=(x, x)_{Q}=(x, x,)_{\tilde{Q}}=q_{\tilde{Q}}(x)$.

### 3.1.2 Dynkin, Tame and Wild Quivers

We want to understand for which quivers $Q$, its Tits form $q_{Q}$ is positive-definite respectively positive semi-definite. This divides all quivers into three distinct classes. The following lemma marks the first step towards this goal.
Lemma 3.13. [2, Claim p. 53] Let $Q$ be a finite connected quiver and let $y \in \operatorname{rad} q_{Q}$ be a positive radical vector. Then $y$ is sincere and $q_{Q}$ is positive semi-definite. Moreover for $x \in \mathbb{Z}^{n}$ we have

$$
q_{Q}(x)=0 \Leftrightarrow x \in \mathbb{Q} y \text {, i.e. } \exists a \in \mathbb{Q} \text { s.t. } x=a \cdot y \Leftrightarrow x \in \operatorname{rad} q_{Q} .
$$

Proof. We denote by $e_{i}$ the $i$-th coordinate vector, i.e. $\left(e_{i}\right)_{j}=\delta_{i j}$. The assumption on $y$ gives:

$$
\begin{equation*}
0=\left(y, e_{i}\right)_{Q}=y^{T} C_{Q} e_{i}=\sum_{j=1}^{n} y_{j} c_{j i}=\left(2-2 d_{i i}\right) y_{i}-\sum_{\substack{j=1 \\ j \neq i}}^{n} d_{j i} y_{j} \text { for } 1 \leq i \leq n \tag{3.1}
\end{equation*}
$$

If $y_{i}=0$, then $\sum_{j=1, j \neq 1}^{n} d_{j i} y_{j}=0$ and since each term is non-negative ( $y$ is positive), we get that $y_{j}=0$ for all $j$ s.t. $i$ and $j$ are joined by an edge (in $\bar{Q}$ ). Using that $\bar{Q}$ is connected, we get that
$y=0$. Indeed, let $i, j_{1}, \ldots, j_{k}, j$ be a path from $i$ to $j$ (meaning that there are edges in $\bar{Q}$ connecting the pairs $\left.\left(i, j_{1}\right),\left(j_{1}, j_{2}\right), \ldots,\left(j_{k}, j\right)\right)$. Since $i$ and $j_{1}$ are connected by an edge, we have $y_{j_{1}}=0$ by the observation above. Using the same observation again we get $y_{j_{2}}=0$. We conclude the claimed result using the observation inductively on the set $i, j_{1}, \ldots, j_{k}, j$. This is a contradiction because $y$ is positive and therefore sincere. We now show that $q_{Q}$ is positive semi-definite. For $x \in \mathbb{Z}^{n}$ we have:

$$
\begin{aligned}
q_{Q}(x) & =\frac{1}{2}(x, x)_{Q}=\sum_{i}\left(2-2 d_{i i}\right) y_{i} \frac{1}{2 y_{i}} x_{i}^{2}-\sum_{i<j} d_{i j} x_{i} x_{j} \\
& =\sum_{i} \sum_{j \neq i} d_{i j} y_{j} \frac{1}{2 y_{i}} x_{i}^{2}-\sum_{i<j} d_{i j} x_{i} x_{j} \\
& =\sum_{i<j} d_{i j} \frac{y_{j}}{2 y_{i}} x_{i}^{2}+\sum_{j<i} d_{i j} \frac{y_{j}}{2 y_{i}} x_{i}^{2}-\sum_{i<j} d_{i j} x_{i} x_{j} \\
& =\sum_{i<j} d_{i j} \frac{y_{j}}{2 y_{i}} x_{i}^{2}+\sum_{i<j} d_{i j} \frac{y_{i}}{2 y_{j}} x_{j}^{2}-\sum_{i<j} d_{i j} x_{i} x_{j} \\
& =\sum_{i<j} d_{i j} \frac{y_{i} y_{j}}{2}\left(\frac{x_{i}^{2}}{y_{i}^{2}}+\frac{x_{j}^{2}}{y_{j}^{2}}-2 \frac{x_{i} x_{j}}{y_{i} y_{j}}\right) \\
& =\sum_{i<j} d_{i j} \frac{y_{i} y_{j}}{2}\left(\frac{x_{i}}{y_{i}}-\frac{x_{j}}{y_{j}}\right)^{2} \geq 0
\end{aligned}
$$

In the second line, we used Equality (3.1). Now if $q_{Q}(x)=0$, then $\frac{x_{i}}{y_{i}}=\frac{x_{j}}{y_{j}}$ whenever $i$ and $j$ are joined by an edge. By a similar argument as above, we use that $\bar{Q}$ is connected and conclude that $\frac{x_{i}}{y_{i}}=\frac{x_{j}}{y_{j}}=a \in \mathbb{Q}$ (for all $i, j$ ). Thus $x \in \mathbb{Q} y$. Now if $x \in \mathbb{Q} y$ we use that $y$ is radical and get $(x, \cdot)_{Q}=a(y, \cdot)_{Q}=0$. Therefore, $x \in \operatorname{rad} q_{Q}$ and $x \in \operatorname{rad} q_{Q}$ implies $q_{Q}(x)=(x, x)_{Q}=0$.

Definition 3.14. A graph is called Euclidean if it is one of the graphs in Figure 3.1.


Figure 3.1: The Euclidean diagrams.

Lemma 3.15. [2, Observation p. 52] Let $\bar{Q}$ be a graph that is neither Dynkin nor Euclidean. Then $\bar{Q}$ contains an Euclidean graph $\Gamma$ as a subgraph.

Proof. Note that if $\bar{Q}$ contains a loop, then $\Gamma=\tilde{A}_{0}$ is the desired subgraph. If $\bar{Q}$ contains a cycle, then the desired subgraph is of the type $\tilde{A_{n}}$. Thus w.l.o.g., we can assume that $\bar{Q}$ does not contain any loops or cycles. We also notice that $\bar{Q}$ has at least one branch vertex (= vertex connected to at least 3 different vertices by edges). Else $\bar{Q}$ is of type $A_{n}$, which is a Dynkin quiver. We do a proof by cases on the number of branch vertices.

Case 1. Let us consider the case that $\bar{Q}$ has exactly one branch vertex. There must be another vertex attached to the graph, or else we have the Dynkin graph $D_{4}$. If the additional vertex is attached to our branch vertex, we have the Euclidean graph $\tilde{D}_{4}$ as a subgraph and we are done. If the additional vertex is not attached to the branch vertex, then we actually have at least two additional vertices, which are attached to different branches. Else $\bar{Q}=D_{n}$ is a Dynkin graph (for some $n$ ). Thus $\bar{Q}$ contains a subgraph that looks like this:


But this is the Dynkin graph $E_{6}$, thus there needs to be another vertex $d$ attached to $a, b$ or $c$ (notice that we only have one branch vertex!). If $d$ is attached to $a$, then our desired subgraph is $\tilde{E}_{6}$. If there is no additional vertex attached to $a$, then by gradually attaching vertices to the ends at $b$ or $c$ we get the graphs $E_{7}, \tilde{E}_{7}, E_{8}$ or $\tilde{E}_{8}$. Thus we eventually get an Euclidean subgraph as claimed. Case 2. $\bar{Q}$ has at least two different branch vertices $a$ and $b$ as illustrated in the following picture:


Since $\bar{Q}$ is connected, we find a path between $a_{1}$ and $b_{1}$. By taking $\Gamma$ to be this path (dotted line) and the 'forks' in the picture above, we get an Euclidean subgraph $\Gamma=\tilde{D_{n}}$ for some $n \geq 6$.

Theorem 3.16. [2, Theorem 1.11] Let $Q$ be a finite connected quiver and let $q_{Q}$ be its Tits form.

1. $Q$ is a Dynkin quiver if and only if $q_{Q}$ is positive definite.
2. The underlying graph $\bar{Q}$ is Euclidean if and only if $q_{Q}$ is positive semi-definite but not positive definite. In that case, there is a unique positive vector $\delta \in \mathbb{Z}^{n}$ s.t. $\operatorname{rad} q_{Q}=\mathbb{Z} \delta$.

Proof. Step 1. If $\bar{Q}$ is Euclidean, then $q_{Q}$ is positive semi-definite and $\operatorname{rad} q_{Q}=\mathbb{Z} \delta$ : For each Euclidean graph, we explicitly give a positive radical vector $\delta$ (we enumerate the vertices as in Definition 3.14).

$$
\begin{aligned}
& \tilde{A}_{0} \delta=(1), \\
& \tilde{A_{n}}: \delta=(1,1, \ldots, 1), \\
& \tilde{D_{n}}: \delta=(1,1,2, \ldots, 2,1,1), \\
& \tilde{E}_{6}: \delta=(1,1,2,2,3,2,1) \\
& \tilde{E}_{7}: \delta=(1,2,2,3,4,3,2,1), \\
& \tilde{E}_{8}: \delta=(1,2,3,4,6,5,4,3,2) .
\end{aligned}
$$

To check that $\delta$ is radical, we check that (see proof of Lemma 3.13):

$$
\begin{equation*}
0=\left(\delta, e_{i}\right)_{Q}=\left(2-2 d_{i i}\right) \delta_{i}-\sum_{\substack{j=1 \\ j \neq i}}^{n} d_{j i} \delta_{j} \text { for } 1 \leq i \leq n \tag{3.2}
\end{equation*}
$$

If $\bar{Q}$ has no loops or multiple edges this is equivalent to (use $d_{i i}=0$ and $d_{i j} \in\{0,1\}$ ):

$$
\begin{equation*}
0=\left(\delta, e_{i}\right)_{Q}=2 \delta_{i}-\sum_{\substack{j=1 \\ d_{i j} \neq 0}}^{n} \delta_{j} \text { for } 1 \leq i \leq n \tag{3.3}
\end{equation*}
$$

We do the calculation for $\tilde{A_{0}}, \tilde{A_{n}}, \tilde{D_{n}}$. The other cases can also be checked by calculation.
$\tilde{A}_{0}:\left(\delta, e_{1}\right) \stackrel{(3.2)}{=}\left(2-2 d_{11}\right) \delta_{1}=(2-2) \cdot 1=0$.
$\tilde{A}_{n}:\left(\delta, e_{i}\right) \stackrel{(3.3)}{=} 2 \delta_{i}-\delta_{i-1}-\delta_{i+1}=2-1-1=0($ we consider the indices $\bmod n+1)$.
師: For $i \in\{0,1, n-1, n\}:\left(\delta, e_{i}\right) \stackrel{(3.3)}{=} 2 \delta_{i}-2=2-2=0$.
For $i \in\{3, \ldots, n-3\}:\left(\delta, e_{i}\right) \stackrel{(3.3)}{=} 2 \delta_{i}-\delta_{i-1}-\delta_{i+1}=2 \cdot 2-2-2=0$.
For $i \in\{2, n-2\}:\left(\delta, e_{i}\right) \stackrel{(3.3)}{=} 2 \delta_{i}-2-1-1=2 \cdot 2-2-1-1=0$.
From Lemma 3.13 it follows that $q_{Q}$ is positive semi-definite. It also follows that rad $q_{Q}=\mathbb{Q} \delta \cap \mathbb{Z}^{n}$. Since every $\delta$ above has $\delta_{i}=1$ for some $i$, we know that for $a \in \mathbb{Q} \backslash \mathbb{Z}$ we have $a \cdot \delta_{i}=a \notin \mathbb{Z}$ and thus $\mathbb{Q} \delta \cap \mathbb{Z}^{n}=\mathbb{Z} \delta=\operatorname{rad} q_{Q}$. From Lemma 3.13 also follows that $q_{Q}(\delta)=0$ and thus $q_{Q}$ is not positive definite. The uniqueness of $\delta$ is a consequence of the property $\operatorname{rad} q_{Q}=\mathbb{Z} \delta$.

Step 2. If $Q$ is Dynkin, then $q_{Q}$ is positive definite: We first notice that for every Dynkin graph $\bar{Q}$, there exists an Euclidean graph $\bar{Q}^{\prime}$ s.t. $\bar{Q}$ results from $\bar{Q}^{\prime}$ by deleting the vertex 0 (and its incident edges). $A_{n}$ results from $\tilde{A}_{n}, D_{n}$ from $\tilde{D}_{n}$ and so on. We can extend $x \in \mathbb{Z}^{n}$ to $x^{\prime}=(0, x) \in \mathbb{Z}^{n+1}$ and from step 1 we get that $q_{\bar{Q}}(x)=q_{\bar{Q}^{\prime}}\left(x^{\prime}\right) \geq 0$. Thus $q_{\bar{Q}}$ is positive semidefinite. To see that it is actually positive definite, let $x \in \mathbb{Z}^{n}$ be such that $q_{\bar{Q}}(x)=q_{\bar{Q}^{\prime}}\left(x^{\prime}\right)=0$. From Lemma 3.13 we get that $x^{\prime}=a \cdot \delta$ for $a \in \mathbb{Q}$ and $\delta$ defined in step 1. But since $x_{0}^{\prime}=0$ we have that $a=0$ (since $\delta_{0} \neq 0$ for all the Euclidean graphs) and thus $x=0$ is trivial.

Step 3. If $\bar{Q}$ is neither Dynkin nor Euclidean, then $q_{Q}$ is indefinite, i.e. $q_{Q}(x)<0$ for some
 radical vector for $\Gamma$ (see step 1 ).

- If $\Gamma_{0}=\bar{Q}_{0}$ (the vertices coincide) and $\Gamma_{1} \subsetneq \bar{Q}_{1}$ (edges do not coincide), we put $x=\delta \in \mathbb{N}^{n}$ and have $q_{\bar{Q}}(x)<q_{\Gamma}(x)=0$. The inequality comes from the fact that there exists an edge $a=(i, j) \in \bar{Q}_{1} \backslash \Gamma_{1}$. Notice that this yields an additional term $-x_{i} x_{j}<0$ in $q_{\bar{Q}}(x)$.
- If $\Gamma_{0} \subsetneq \bar{Q}_{0}$, then let $i$ be a vertex in $\bar{Q}_{0} \backslash \Gamma_{0}$ which is connected to $\Gamma$ by an edge $a$ (such a vertex exists since $\bar{Q}$ is connected). Now let $\delta^{\prime}$ be an extension of $\delta\left(\delta_{j}^{\prime}=0\right.$ for $\left.j \in \bar{Q}_{0} \backslash \Gamma_{0}\right)$ and define $x=2 \delta^{\prime}+e_{i}$. We have $q_{\bar{Q}}(x) \leq 4 q_{\Gamma}(\delta)+x_{i}^{2}-x_{t_{a}} x_{h_{a}} \leq 0+1-2=-1<0$.

Definition 3.17. A quiver $Q$ is called a tame quiver if its underlying graph $\bar{Q}$ is Euclidean. If $Q$ is neither a Dynkin quiver nor a tame quiver, it is called a wild quiver.

This definition classifies all finite connected quivers into three different categories: Dynkin quivers, tame quivers, and wild quivers. Theorem 3.16 yields the same classification, but it considers the Tits forms $q_{Q}$ rather than classifying the quivers based on underlying graphs $\bar{Q}$. The Tits form of the Dynkin quivers is positive definite, the Tits form of the tame quivers is positive semi-definite and the Tits form of the wild quivers is indefinite.

Remark 3.18. In Gabriel's Theorem, there is no distinction between tame and wild quivers. But in the theory beyond Gabriel's Theorem, there are some interesting differences between the two categories, which is why we differentiate between the two categories.

Example 3.19. The loop quiver is a tame quiver since its underlying graph is the Euclidean graph $\tilde{A}_{0}$. This matches with our calculation in Example 3.8: $q_{Q}(x)=\frac{1}{2}(x, x)_{Q}=0$, which is a positive semi-definite form.

### 3.2 Weyl Groups and Root Systems

In this section, we introduce the Weyl group and the root system of a finite quiver. This section is based on the book 'Finite dimensional algebras and quantum groups' by B. Deng [2].

Definition 3.20. Let $Q$ be a finite quiver and let $i \in Q_{0}=\{1, \ldots, n\}$ be a vertex such that $Q$ has no loops at $i$, i.e. $d_{i i}=0$. Then we say that the vector $e_{i} \in \mathbb{Z}^{n}$ is a simple root. We denote by $\Pi_{Q}$ the set of all simple roots.

Definition 3.21. Let $Q$ be a finite quiver and let $i$ be a vertex such that $Q$ has no loops at $i$. The simple reflection at the vertex $i$ is the map

$$
\sigma_{i}: \mathbb{Z}^{n} \rightarrow \mathbb{Z}^{n}, x \mapsto x-\left(x, e_{i}\right)_{Q} e_{i}
$$

where $(x, y)_{Q}$ is the symmetric Euler form of $Q$ and $e_{i}$ is the $i$-th coordinate vector. Using that $\left(e_{j}, e_{i}\right)_{Q}=-d_{i j}$ for $j \neq i$ and $\left(e_{i}, e_{i}\right)_{Q}=2-2 d_{i i}=2$ we get $\sigma_{i}(x)_{j}=x_{j}$ for $j \neq i$ and

$$
\sigma_{i}(x)_{i}=x_{i}-\sum_{j=1}^{n} x_{j}\left(e_{j}, e_{i}\right)_{Q}=x_{i}-\sum_{\substack{j=1 \\ j \neq i}}^{n}-d_{j i} x_{j}-2 x_{i}=\sum_{\substack{j=1 \\ j \neq i}}^{n} d_{j i} x_{j}-x_{i}
$$

Example 3.22. We consider the $L_{2}$-quiver. We have $\left(e_{i}, e_{i}\right)_{Q}=2-2 d_{i i}=2$ for $i \in\{1,2\}$. Also $\left(x, e_{1}\right)_{Q}=x_{1}\left(e_{1}, e_{1}\right)_{Q}+x_{2}\left(e_{2}, e_{1}\right)_{Q}=2 x_{1}-d_{12} x_{2}=2 x_{1}-x_{2}$. Similarly, we get $\left(x, e_{2}\right)_{Q}=2 x_{2}-x_{1}$. Thus we get the maps $\sigma_{1}, \sigma_{2}: \mathbb{Z}^{2} \rightarrow \mathbb{Z}^{2}$ given by:

$$
\begin{aligned}
& \sigma_{1}\left(\left(x_{1}, x_{2}\right)^{T}\right)=\binom{x_{1}}{x_{2}}-\binom{2 x_{1}-x_{2}}{0}=\binom{x_{2}-x_{1}}{x_{2}}=\left(\begin{array}{cc}
-1 & 1 \\
0 & 1
\end{array}\right)\binom{x_{1}}{x_{2}}, \\
& \sigma_{2}\left(\left(x_{1}, x_{2}\right)^{T}\right)=\binom{x_{1}}{x_{2}}-\binom{0}{2 x_{2}-x_{1}}=\binom{x_{1}}{x_{1}-x_{2}}=\left(\begin{array}{cc}
1 & 0 \\
1 & -1
\end{array}\right)\binom{x_{1}}{x_{2}} .
\end{aligned}
$$

Lemma 3.23. The simple reflections $\sigma_{i}$ are automorphisms of order two and they preserve the symmetric Euler form $(\cdot, \cdot)_{Q}$ and the Tits form $q_{Q}$.
Proof. We show that $\sigma_{i}^{2}=\mathrm{id}_{\mathbb{Z}^{n}}$. Indeed, for $j \neq i$ we have $\sigma_{i}^{2}(x)_{j}=\sigma_{i}(x)_{j}=x_{j}$ and we also have

$$
\sigma_{i}\left(\sigma_{i}(x)\right)_{i}=\sum_{\substack{j=1 \\ j \neq i}}^{n} d_{j i} \sigma_{i}(x)_{j}-\sigma_{i}(x)_{i}=\sum_{\substack{j=1 \\ j \neq i}}^{n} d_{j i} x_{j}-\sigma_{i}(x)_{i}=\sum_{\substack{j=1 \\ j \neq i}}^{n} d_{j i} x_{j}-\left(\sum_{\substack{j=1 \\ j \neq i}}^{n} d_{j i} x_{j}-x_{i}\right)=x_{i}
$$

We now show that for any quiver $Q$ and for all $x, y \in \mathbb{Z}^{n}$, we have $\left(\sigma_{i}(x), \sigma_{i}(y)\right)_{Q}=(x, y)_{Q}$. First, we notice that

$$
\sigma_{i}\left(e_{k}\right)=e_{k}-\left(e_{i}, e_{k}\right)_{Q} e_{i}= \begin{cases}e_{i}-\left(2-2 d_{i i}\right) e_{i}=-e_{i} & \text { if } k=i \\ e_{k}+d_{i k} e_{i} & \text { if } k \neq i\end{cases}
$$

Thus we get $\left(\sigma_{i}\left(e_{i}\right), \sigma_{i}\left(e_{i}\right)\right)_{Q}=\left(-e_{i},-e_{i}\right)_{Q}=\left(e_{i}, e_{i}\right)_{Q}$ and for $k \neq i$ we get:

$$
\left(\sigma_{i}\left(e_{k}\right), \sigma_{i}\left(e_{i}\right)\right)_{Q}=\left(e_{k}+d_{i k} e_{i},-e_{i}\right)_{Q}=-\left(e_{k}, e_{i}\right)_{Q}-d_{i k}\left(e_{i}, e_{i}\right)_{Q}=d_{i k}-2 d_{i k}=-d_{i k}=\left(e_{k}, e_{i}\right)_{Q}
$$

Also for $k \neq i, j \neq i$ we get:

$$
\begin{aligned}
\left(\sigma_{i}\left(e_{k}\right), \sigma_{i}\left(e_{l}\right)\right)_{Q} & =\left(e_{k}+d_{i k} e_{i}, e_{l}+d_{i l} e_{i}\right)_{Q}=\left(e_{k}, e_{l}\right)_{Q}+d_{i k}\left(e_{i}, e_{l}\right)_{Q}+d_{i l}\left(e_{k}, e_{i}\right)_{Q}+d_{i k} d_{i l}\left(e_{i}, e_{i}\right)_{Q} \\
& =\left(e_{k}, e_{l}\right)_{Q}-d_{i k} d_{i l}-d_{i l} d_{i k}+2 d_{i k} d_{i l}=\left(e_{k}, e_{l}\right)_{Q}
\end{aligned}
$$

Now, using linearity of $\sigma_{i}$ and bilinearity of $(\cdot, \cdot)_{Q}$, we get that:

$$
\begin{aligned}
\left(\sigma_{i}(x), \sigma_{i}(y)\right)_{Q} & =\left(\sigma_{i}\left(\sum_{j} x_{j} e_{j}\right), \sigma_{i}\left(\sum_{k} y_{k} e_{k}\right)\right)_{Q}=\sum_{j} \sum_{k} x_{j} y_{k}\left(\sigma_{i}\left(e_{j}\right), \sigma_{i}\left(e_{k}\right)\right)_{Q} \\
& =\sum_{j} \sum_{k} x_{j} y_{k}\left(e_{j}, e_{k}\right)_{Q}=\left(\sum_{j} x_{j} e_{j}, \sum_{k} y_{k} e_{k}\right)_{Q}=(x, y)_{Q} .
\end{aligned}
$$

Finally, this implies:

$$
q_{Q}\left(\sigma_{i}(x)\right)=\frac{1}{2}\left(\sigma_{i}(x), \sigma_{i}(x)\right)_{Q}=\frac{1}{2}(x, x)_{Q}=q_{Q}(x)
$$

Definition 3.24. For a finite quiver $Q$, we define the Weyl group $W(Q)$ to be the subgroup of the automorphism group $\operatorname{Aut}\left(\mathbb{Z}^{n}\right)$ generated by the simple reflections $\sigma_{i}$. The set

$$
\Phi_{\mathrm{re}}(Q)=\bigcup_{w \in W(Q)} w\left(\Pi_{Q}\right)=\bigcup_{\substack{w \in W(Q) \\ \alpha \in \Pi_{Q}}} w(\alpha)
$$

is called the set of real roots of $Q$.

Because the symmetric Euler form is independent of the orientations of the arrows in $Q$, both the simple representations $\sigma_{i}$ and the Weyl group $W(Q)$ are also independent of the orientations of the arrows in $Q$.
Definition 3.25. For each $x=\sum_{i=1}^{n} x_{i} e_{i} \in \mathbb{Z}^{n}$ we define the support of $x$ to be the set

$$
\operatorname{supp} x=\left\{i \in Q_{0} \mid x_{i} \neq 0\right\}
$$

We say that $\operatorname{supp} x$ is connected if the full subquiver of $Q$ with vertex set $\operatorname{supp} x$ is connected.
Definition 3.26. For a finite quiver $Q$, the fundamental set $\mathcal{F}_{Q}$ is defined to be the set

$$
\mathcal{F}_{Q}=\left\{0 \neq x \in \mathbb{N}^{n} \mid\left(x, e_{i}\right)_{Q} \leq 0 \text { for all } e_{i} \in \Pi_{Q}, \text { and } \operatorname{supp} x \text { is connected }\right\} .
$$

Then the set

$$
\Phi_{\mathrm{im}}(Q)=\bigcup_{w \in W(Q)} w\left(\mathcal{F}_{Q}\right) \cup w\left(-\mathcal{F}_{Q}\right)=\bigcup_{\substack{w \in W(Q) \\ x \in \mathcal{F}_{Q}}} w(x) \cup w(-x)
$$

is called the set of imaginary roots of $Q$. Finally, the root system of $Q$ is defined to be the set

$$
\Phi(Q)=\Phi_{\mathrm{re}}(Q) \cup \Phi_{\mathrm{im}}(Q)
$$

An element $x \in \Phi(Q) \cap \mathbb{N}^{n}$ is called a positive root. We denote by $\Phi^{+}(Q)\left(\right.$ resp. $\left.\Phi_{\mathrm{re}}^{+}(Q), \Phi_{\mathrm{im}}^{+}(Q)\right)$ the set of all positive (resp. positive real, positive imaginary) roots.

Example 3.27. Let $Q$ be the loop quiver. Then no simple root exists and thus $\Pi_{Q}=\emptyset$ and $W(Q)=\{\mathrm{id}\}$. We also have that $\mathcal{F}_{Q}=\mathbb{N}$ and thus $\Phi(Q)=\Phi_{\mathrm{im}}(Q)=\mathbb{Z} \backslash\{0\}$.

For a better understanding of the definitions above, we prove some elementary properties of our root systems $\Phi_{\mathrm{re}}(Q), \Phi_{\mathrm{im}}(Q)$ and $\Phi(Q)$.

Lemma 3.28. Let $Q$ be a finite quiver.

1. If $Q$ has no loops at $i \in Q_{0}$, then $e_{i}$ is a real root.
2. If $Q$ has a loop at $i \in Q_{0}$, then $e_{i} \in \mathcal{F}_{Q}$ and thus $e_{i}$ is an imaginary root.
3. If $x \in \Phi_{\mathrm{re}}(Q)$, then $q_{Q}(x)=1$.
4. If $x \in \Phi_{\mathrm{im}}(Q)$, then $q_{Q}(x) \leq 0$.
5. $\Phi_{\mathrm{re}}(Q) \cap \Phi_{\mathrm{im}}(Q)=\emptyset$.
6. If $Q$ is a Dynkin or a tame quiver, then every root is either positive or negative. Therefore, we have $\Phi(Q)=\Phi^{+}(Q) \cup\left(-\Phi^{+}(Q)\right)$.

Proof. 1. Taking $w=\mathrm{id} \in w(Q)$, we get that $w\left(e_{i}\right)=e_{i} \in \Phi_{\mathrm{re}}(Q)$.
2. In this case, $\operatorname{supp} e_{i}$ is just one point and thus connected. Moreover, for all $e_{j} \in \Pi_{Q}$, we have:

$$
\left(e_{i}, e_{j}\right)_{Q}= \begin{cases}2-2 d_{i i} \leq 0 & \text { if } j=i \\ -d_{i j} \leq 0 & \text { if } j \neq i\end{cases}
$$

3. Each $w \in W(Q)$ is of the form $w=\sigma_{i_{k}} \ldots \sigma_{i_{1}}$ for $\sigma_{i_{j}}$ simple reflections (possibly $k=0$ and $w=\mathrm{id}$ ). By Lemma 3.23, we have:

$$
q_{Q}(x)=q_{Q}\left(\sigma_{i_{k}} \ldots \sigma_{i_{1}}\left(e_{j}\right)\right)=q_{Q}\left(\sigma_{i_{k-1}} \ldots \sigma_{i_{1}}\left(e_{j}\right)\right)=\ldots=q_{Q}\left(e_{j}\right)=\frac{1}{2}\left(e_{j}, e_{j}\right)_{Q}=1-d_{j j}=1
$$

4. We have $x=w(y)$ or $x=w(-y)$ for $w \in W(Q)$ and $y \in \mathcal{F}_{Q}$. In both cases, we have $q_{Q}(x)=q_{Q}(w( \pm y))=q_{Q}( \pm y)=q_{Q}(y)$. Since $y \in \mathbb{N}^{n}$, we have:

$$
q_{Q}(x)=q_{Q}(y)=\frac{1}{2} \sum_{i=1}^{n} y_{i}\left(y, e_{i}\right)_{Q}=\frac{1}{2} \sum_{d_{i i}=0} y_{i} \underbrace{\left(y, e_{i}\right)_{Q}}_{\leq 0}+\frac{1}{2} \sum_{d_{i i} \neq 0} y_{i} \sum_{j=1}^{n} y_{j} \underbrace{\left(e_{j}, e_{i}\right)_{Q}}_{\leq 0 \text { by part } 2} \leq 0 .
$$

5. This immediately follows from parts 3 and 4.
6. Let $x$ be a root and write $x=x^{+}-x^{-}$where $x^{+}$is the positive part of $x$ and $x^{-}$is the negative part of $x$, i.e. $x^{+}, x^{-} \geq 0$ and they have disjoint support. Then we have $\left(x^{+}, x^{-}\right)_{Q} \leq 0$ since they have disjoint support. We then have:

$$
\begin{aligned}
1 \geq q_{Q}(x) & =\frac{1}{2}\left(x^{+}-x^{-}, x^{+}-x^{-}\right)_{Q}=\frac{1}{2}\left[\left(x^{+}, x^{+}\right)_{Q}+\left(x^{-}, x^{-}\right)_{Q}-2\left(x^{+}, x^{-}\right)_{Q}\right] \\
& =q_{Q}\left(x^{+}\right)+q_{Q}\left(x^{-}\right)-\left(x^{+}, x^{-}\right)_{Q} \geq q_{Q}\left(x^{+}\right)+q_{Q}\left(x^{-}\right) \geq 0 .
\end{aligned}
$$

In the last inequality, we used that $q_{Q}$ is positive semi-definite since $Q$ is Dynkin or tame. We conclude that either $q_{Q}\left(x^{+}\right)=0$ or $q_{Q}\left(x^{-}\right)=0$. If $Q$ is Dynkin, then either $x^{+}=0$ or $x^{-}=0$ since $q_{Q}$ is positive definite. If $Q$ is tame, then we know from Lemma 3.13 that either $x^{+}$or $x^{-}$is a radical vector and thus it is sincere. But then the other vector is trivial, which shows that $x$ is either positive or negative.

Remark 3.29. Part (6) of Lemma 3.28 is true for any quiver $Q$. In addition, one can show that for any root $x \in \Phi(Q), \operatorname{supp} x$ is connected. We do not prove these statements because they are not needed to prove Gabriel's Theorem. A proof can be found on page 7 in Chapter 1 of [6].

An immediate consequence of this lemma is that $\Phi(Q) \subseteq\left\{0 \neq x \in \mathbb{Z}^{n} \mid q_{Q}(x) \leq 1\right\}$. We will later see that if $Q$ is a Dynkin quiver, then these two sets are equal. Although the characterization $q_{Q}(x) \leq 1$ looks a lot simpler than the definition of the root systems $\Phi_{\mathrm{re}}(Q)$ and $\Phi_{\mathrm{im}}(Q)$, the latter give a more explicit description of the roots and are therefore useful in proofs. Moreover, the latter descriptions of $\Phi_{\mathrm{re}}(Q)$ and $\Phi_{\mathrm{im}}(Q)$ are closely related to the Weyl group, which plays an important role in the proof of Gabriel's Theorem (see Lemma 3.68 and Proposition 3.69).

### 3.2.1 Root Systems of Dynkin, Tame and Wild Quivers

We show an explicit characterization of the roots of a Dynkin quiver. In particular, we show that the root system $\Phi(Q)$ is finite for all Dynkin quivers. We also mention some results about the root systems of tame and wild quivers.

Lemma 3.30. If $Q$ is a Dynkin quiver, then $\Phi(Q)$ is finite.
Proof. We prove that the set $\left\{x \in \mathbb{Z}^{n} \mid q_{Q}(x)=1\right\}$ is finite. First, we can view $q_{Q}$ also as a quadratic form on $\mathbb{Q}^{n}$ and $\mathbb{R}^{n}$. Since $q_{Q}$ is positive definite on $\mathbb{Z}^{n}$ it is also positive definite on $\mathbb{Q}^{n}$. Indeed, if $0 \neq x=\left(\frac{a_{i}}{b_{i}}\right)_{i=1}^{n} \in \mathbb{Q}^{n}$ such that $q_{Q}(x) \leq 0$, then we have that $q_{Q}(a x) \leq 0$, where $a$ is the least common multiple of the $b_{i}$ and $0 \neq a x \in \mathbb{Z}^{n}$. Now, we can take limits in $\mathbb{Q}^{n}$ to extend $q_{Q}$ to $\mathbb{R}^{n}$ and get that $q_{Q}$ is positive semi-definite on $\mathbb{R}^{n}$. However, since $q_{Q}$ is positive definite on $\mathbb{Q}^{n}$, we know that its matrix $C_{Q}$ is positive definite and symmetric (with values in $\mathbb{Q}$ ), thus it is invertible in $\mathbb{Q}$ and therefore also invertible in $\mathbb{R}$. Therefore, $q_{Q}$ is also positive definite on $\mathbb{R}^{n}$ (since it is positive semi-definite and $C_{Q}$ is invertible). Now consider the subset $S^{1}=\left\{x=\left(x_{i}\right) \in \mathbb{R}^{n} \mid\|x\|=1\right\}$. Here, $\|\cdot\|$ is the Euclidean norm. Since $S^{1}$ is compact and $q_{Q}$ is positive definite and continuous, there exists $c>0$ such that $q_{Q}(x) \geq c$ for all $x \in S^{1}$. Indeed, we have that $q_{Q}\left(S^{1}\right)$ is compact and thus $\inf _{x \in S^{1}}\left(q_{Q}(x)\right)=c>0$. Therefore, for all $0 \neq x \in \mathbb{R}^{n}$, we have:

$$
q_{Q}(x)=\|x\|^{2} q_{Q}\left(\frac{x}{\|x\|}\right) \geq c\|x\|^{2}
$$

Therefore, $x \in \Phi(Q)$ implies $q_{Q}(x)=1$, which again implies $\|x\|^{2} \leq \sqrt{\frac{1}{c}}$. But there are only finitely many $x \in \mathbb{Z}^{n}$ such that the last inequality holds.

Corollary 3.31. If $Q$ is a Dynkin quiver, then

1. the Weyl group $W(Q)$ is finite,
2. $\Phi(Q)=\left\{x \in \mathbb{Z}^{n} \mid q_{Q}(x)=1\right\}$.

Proof. 1. By the theorem above, $\Phi(Q)=\Phi_{\mathrm{re}}(Q)$ is finite. For any $w \in W(Q)$, the induced map $\Phi_{\mathrm{re}}(Q) \rightarrow \Phi_{\mathrm{re}}(Q), w^{\prime}\left(e_{i}\right) \mapsto w \circ w^{\prime}\left(e_{i}\right)$ is a permutation on the set $\Phi(Q)$ of roots since $w$ is an automorphism. Since $\Phi(Q)$ contains the basis $\left\{e_{i} \mid i \in Q_{0}\right\}$ of $\mathbb{Z}^{n}$, it follows that the permutation maps that are induced by $w_{1}, w_{2} \in W(Q)$ are equal if and only if $w_{1}=w_{2}$. Therefore, $W(Q)$ can be embedded into the permutation group of $\Phi(Q)$, which is finite since $\Phi(Q)$ is finite.
2. We have already shown that each $x \in \Phi(Q)=\Phi_{\text {re }}(Q)$ satisfies $q_{Q}(x)=1$. Now let $x \in \mathbb{Z}^{n}$ be such that $q_{Q}(x)=1$. We know from Lemma 3.28 that $x$ is either positive or negative. Since $\Phi(Q)=\Phi^{+}(Q) \cup\left(-\Phi^{+}(Q)\right)$, it suffices to show that $x$ or $-x$ is a root. Thus w.l.o.g., we assume that $x$ is positive. We use induction of the number $l_{x}=\sum_{i \in Q_{0}}$. If $l_{x}=1$, then $x=e_{i} \in \Pi_{Q} \subseteq \Phi(Q)$. Now let $l_{x} \geq 2$. We claim that there exists an $i \in Q_{0}$ such that $0<\sigma_{i}(x)<x$. Suppose this is not the case, i.e. for each $i \in Q_{0}$ we have either $0 \nless \sigma_{i}(x)$ or $\sigma_{i}(x) \nless x$. We consider the following formula:

$$
2=2 q_{Q}(x)=(x, x)_{Q}=\sum_{i \in Q_{0}} x_{i}\left(x, e_{i}\right)_{Q}
$$

It implies that there exists an $i_{0} \in Q_{0}$ such that $x_{i_{0}}>0$ and $\left(x, e_{i_{0}}\right)_{Q}>0$. Thus we get $\sigma_{i_{0}}(x)=x-\left(x, e_{i_{0}}\right)_{Q} e_{i_{0}}<x$. But this forces $0 \nless \sigma_{i_{0}}(x)$, i.e. $x_{i_{0}}-\left(x, e_{i_{0}}\right)_{Q}<0$. On the other hand,

$$
\begin{aligned}
0 & \leq q_{Q}\left(x-x_{i_{0}} e_{i_{0}}\right)=q_{Q}(x)+q_{Q}\left(x_{i_{0}} e_{i_{0}}\right)+\left(x,-x_{i_{0}} e_{i_{0}}\right)_{Q} \\
& =1+x_{i_{0}}^{2}-x_{i_{0}}\left(x, e_{i_{0}}\right)_{Q}=1+x_{i_{0}}\left(x_{i_{0}}-\left(x, e_{i_{0}}\right)_{Q}\right)
\end{aligned}
$$

This forces $x_{i_{0}}=1$ and $x_{i_{0}}-\left(x, e_{i_{0}}\right)_{Q}=-1$. But then $q_{Q}\left(x-x_{i_{0}} e_{i_{0}}\right)=0$, i.e. $x=x_{i_{0}} e_{i_{0}}=e_{i_{0}}$ which contradicts $l_{x} \geq 2$. Therefore, there exists an $i \in Q_{0}$ such that $0<\sigma_{i}(x)<x$. But then $q_{Q}\left(\sigma_{i}(x)\right)=q_{Q}(x)=1$ and $l_{\sigma_{i}(x)}<l_{x}$ and thus by the induction hypothesis, we have $\sigma_{i}(x) \in \Phi(Q)$. We conclude that $x=\sigma_{i}\left(\sigma_{i}(x) \in \Phi(Q)\right.$.

Remark 3.32. To get a sense of the size of the finite root system $\Phi(Q)$ for a Dynkin quiver $Q$, we list the sizes. If $Q$ is a Dynkin quiver of type $A_{n}(n \geq 1), D_{n}(n \geq 4), E_{6}, E_{7}$ or $E_{8}$, then the number of roots in $\Phi(Q)$ is $n(n+1), 2 n(n-1), 72,126$, or 240, respectively.

We now have proven everything about the root system $\Phi(Q)$ and the Weyl group $W(Q)$ that we need to prove Gabriel's Theorem. Thus, if you are only interested in the proof of Gabriel's Theorem, you can continue reading Section 3.3. However, we have already done (almost) everything to prove a really nice theorem that connects different types of quivers and their root systems. Therefore, we spend the remaining part of this section proving the following theorem.

Theorem 3.33. [2, Theorem 1.13] Let $Q$ be a finite connected quiver without loops. Then

1. $Q$ is a Dynkin quiver $\Longleftrightarrow \Phi(Q)$ is finite $\Longleftrightarrow \Phi_{\mathrm{im}}(Q)=\emptyset$.
2. $Q$ is a tame quiver $\Longleftrightarrow \Phi_{\mathrm{im}}(Q)=\mathbb{Z} \delta \backslash\{0\}$ for the unique positive radical vector $\delta \in \mathbb{N}^{n}$ from Theorem 3.16.
3. $Q$ is a wild quiver $\Longleftrightarrow$ there exists $x \in \Phi^{+}(Q)$ such that $\operatorname{supp} x=Q$ and such that $\left(x, e_{i}\right)_{Q}<0$ for all $i \in Q_{0}$.

To prove this theorem, we need to consider the following lemma.
Lemma 3.34. Let $Q$ be a finite connected quiver without loops. Then $Q$ is wild if and only if there exists a positive, sincere $x \in \mathbb{N}^{n}$ such that $C_{Q} x$ is negative and sincere, i.e. $\left(C_{Q} x\right)_{i}<0$. Here, $C_{Q}$ denotes the matrix belonging to the symmetric Euler form. This means that for all $x, y \in \mathbb{Z}^{n}$ we have $(x, y)_{Q}=x^{T} C_{Q} y$.
Proof. The proof of this lemma can be found on page 59 of [2]. We skip it since it mainly consists of algebraic manipulations.

Proof of Theorem 3.33. 1. We know that $Q$ is a Dynkin quiver if and only if its Tits form $q_{Q}$ is positive definite. We use this to show that $Q$ is Dynkin if and only if $\mathcal{F}_{Q}=\emptyset$. Indeed, if $Q$ is Dynkin, then $x \in \mathcal{F}_{Q}$ implies that $q_{Q}(x)=\frac{1}{2}(x, x)_{Q}=\frac{1}{2} \sum_{i=1}^{n} x_{i}\left(x, e_{i}\right) \leq 0$ and since $q_{Q}$ is positive definite, this implies that $x=0$ and thus $\mathcal{F}_{Q}=\emptyset$. We prove the other implication via contraposition. Therefore, we consider the cases for a tame resp. a wild quiver $Q$. If $Q$ is a tame quiver, then we know from Theorem 3.16 that the unique smallest positive radical vector $\delta$ is in $\mathcal{F}_{Q}$. In the case that $Q$ is a wild quiver, we know from the proof of Theorem 3.16 that there exists a positive vector $\delta$ (with connected support) such that $q_{Q}(\delta)<0$. We show that $\delta \in \mathcal{F}_{Q}$. First, notice that if $x \in \mathbb{N}^{n}$ and $j \notin \operatorname{supp} x$, then $\left(x, e_{j}\right)_{Q} \leq 0$. Indeed $\left(x, e_{j}\right)_{Q}=\sum_{k \in \operatorname{supp} x} x_{k}\left(e_{k}, e_{j}\right)_{Q}=\sum_{k \in \operatorname{supp} x}-x_{k} d_{k j} \leq 0$. Moreover, using the case distinction from the proof of Theorem 3.16, we get that for $i \in \operatorname{supp} x$, we either have $\left(\delta, e_{i}\right)_{Q}=\left(\delta^{\prime}, e_{i}\right)_{Q}=0$ or $\left(\delta, e_{i}\right)_{Q}=\left(e_{j}, e_{i}\right)_{Q} \leq 0$ (since $Q$ has no loops). Using that $\mathcal{F}_{Q}=\emptyset \Longleftrightarrow \Phi_{\mathrm{im}}(Q)=\emptyset$, we get that $Q$ is Dynkin $\Longleftrightarrow \Phi_{\mathrm{im}}(Q)=\emptyset$. From Lemma 3.30 we know that if $Q$ is Dynkin, then $\Phi(Q)$ is finite. Notice that $x \in \mathcal{F}_{Q}$ implies $a x \in \mathcal{F}_{Q}$ for all $0<a \in \mathbb{Z}$. Therefore, if $\Phi(Q)$ is finite, this implies that $\Phi_{\mathrm{im}}(Q)=\emptyset$ and therefore $Q$ is Dynkin.
2. Let $Q$ be a tame quiver. We show that $\Phi_{\mathrm{im}}(Q)=\operatorname{rad} q_{Q} \backslash\{0\}=\mathbb{Z} \delta \backslash\{0\}$. Let $0 \neq x \in \operatorname{rad} q_{Q}$. Then we have $\left(x, e_{i}\right)_{Q}=0$ for all $i \in Q_{0}$ and from Lemma 3.13, we know that $x$ is sincere and therefore supp $x=Q$ is connected. Thus we have $x \in \Phi_{\mathrm{im}}(Q)$. Now let $0 \neq x \in \Phi_{\mathrm{im}}(Q)$. Then we have $q_{Q}(x) \leq 0$ but since $q_{Q}$ is positive semi-definite ( $Q$ is tame), we have $q_{Q}(x)=0$ and from Lemma 3.13, we know that $x \in \operatorname{rad} q_{Q}$. Finally, from Theorem 3.16, we know that $\operatorname{rad} q_{Q}=\mathbb{Z} \delta$. Note that if $Q$ is Dynkin, we know that $\Phi_{\mathrm{im}}(Q)=\emptyset$. Therefore, it remains to show that for a wild quiver $Q$, we have that $\Phi_{\mathrm{im}}(Q) \neq \mathbb{Z} \delta \backslash\{0\}$. Let $Q$ be wild and let $x$ be as in Lemma 3.34. For any $y \in \mathbb{Z}^{n}$, there exists a suitably large positive integer $m$ such that $y-m x \in \mathcal{F}_{Q} \subseteq \Phi_{\mathrm{im}}(Q)$. Indeed, since $\operatorname{supp} x$ is connected, we have that $\operatorname{supp}(y-m x)$ is connected (for $m$ large enough). Moreover, we can find $m$ large enough such that for all $i \in Q_{0}$ we have (note that $Q_{0}$ is finite):

$$
\left(y-m x, e_{i}\right)_{Q}=\left(y, e_{i}\right)_{Q}-m\left(x, e_{i}\right)_{Q}=\left(y, e_{i}\right)_{Q}-m \underbrace{\left(C_{Q} x\right)_{i}}_{<0} \leq 0
$$

But then, $\Phi_{\mathrm{im}}(Q)$ cannot be equal to $\mathbb{Z} \delta \backslash\{0\}$ for a single $\delta$. Otherwise we had that for any $y, z \in \mathbb{Z}^{n}$ and all $m_{1}, m_{2}$ large enough, there would be $a_{1}, a_{2}$ such that $y-m_{1} x=a_{1} \delta$, $z-m_{2} x=a_{2} \delta$. This yields

$$
\left(m_{2} a_{1}-m_{1} a_{2}\right) \delta=m_{2} y-m_{2} m_{1} x-m_{1} z+m_{1} m_{1} x=m_{2} y-m_{1} z
$$

However, this is not possible if $y, z$ are linearly independent. Therefore, if we have that $\Phi_{\mathrm{im}}(Q)=\mathbb{Z} \delta \backslash\{0\}$, then $Q$ can only be a tame quiver.
3. This follows immediately from Lemma 3.34. Indeed, notice that $x$ is sincere if and only if $\operatorname{supp} x=Q$. In addition, $\left(x, e_{i}\right)_{Q}=\left(C_{Q} x\right)_{i}$.

### 3.3 Reflection Functors

In this section, we introduce reflection functors, which are a useful tool to analyze and compare representations.

Definition 3.35. A vertex $i \in Q_{0}$ is called a sink if there does not exist an edge starting at i, i.e. for all $a \in Q_{1}$ we have $t_{a} \neq i$. In contrast, a vertex $j \in Q_{0}$ is called a source if there is no edge ending at $j$, i.e. $h_{a} \neq j$. For any vertex $i \in Q_{0}$ we denote by $Q_{1}^{i} \subseteq Q_{1}$ the subset of all the arrows which are incident to $i$, i.e. $a \in Q_{1}^{i}$ if a either starts or ends at $i$.
Definition 3.36. Let $i \in Q_{0}$ be a vertex. We call $s_{i}$ the reflection at vertex $i$. This means that $s_{i} Q$ is the quiver obtained by reversing all the arrows that start or end at $i$.

Example 3.37. Consider the following quiver $Q$ :


The vertices 1 and 4 are sources, the vertices 2 and 5 are sinks and vertex 3 is neither a source nor a sink. We list all possible reflections of $Q$ :


Definition 3.38. Let $Q$ be a quiver and let $i \in Q_{0}$ be a sink. Notice that for each $a \in Q_{1}^{i}$ we have $h_{a}=i$. We define a covariant functor $\mathcal{R}_{i}^{+}$from the category $\operatorname{rep}_{\mathrm{k}}(Q)$ to the category $\operatorname{rep}_{\mathrm{k}}\left(s_{i} Q\right)$. For every representation $\mathbb{V}=\left(V_{i}, v_{a}\right) \in \operatorname{rep}_{\mathrm{k}}(Q)$, we define a representation $\mathcal{R}_{i}^{+} \mathbb{V}=\left(V_{i}^{\prime}, v_{a}^{\prime}\right)$ in the category $\operatorname{rep}_{\mathrm{k}}\left(\mathrm{s}_{\mathrm{i}} \mathrm{Q}\right)$ as follows. For all $j \neq i$ we set $V_{j}^{\prime}=V_{j}$, and define $V_{i}^{\prime}$ to be the kernel of the map

$$
\xi_{i}: \bigoplus_{a \in Q_{1}^{i}} V_{t_{a}} \rightarrow V_{i},\left(x_{t_{a}}\right)_{a \in Q_{1}^{i}} \mapsto \sum_{a \in Q_{1}^{i}} v_{a}\left(x_{t_{a}}\right) .
$$

For $\mathcal{R}_{i}^{+} \boxtimes$ to be a representation, we need to define linear maps belonging to the arrows in $s_{i} Q$. For each arrow $a \notin Q_{1}^{i}$ we set $v_{a}^{\prime}=v_{a}$. For $a \in Q_{1}^{i}$, let $b$ be the reverse arrow and we define $v_{b}^{\prime}$ to be the composition of maps

$$
V_{t_{b}}^{\prime}=V_{i}^{\prime}=\operatorname{ker} \xi_{i} \hookrightarrow \bigoplus_{c \in Q_{1}^{i}} V_{t_{c}} \rightarrow V_{t_{a}}=V_{t_{a}}^{\prime}=V_{h_{b}}^{\prime}
$$

Here the map before the direct sum is the canonical inclusion and the map after the direct sum is the canonical projection onto the component $V_{t_{a}}$. To better understand what happens, we look at the following picture. The red vertices belong to the summands of the space $\bigoplus_{a \in Q_{1}^{i}} V_{t_{a}}$. Outside of the circle, everything stays the same.


Now let $\phi: \mathbb{V} \rightarrow \mathbb{W}$ be a morphism between two representations $\mathbb{V}, \mathbb{W} \in \operatorname{rep}_{\mathrm{k}}(Q)$. We define the morphism $\phi^{\prime}=\mathcal{R}_{i}^{+} \phi: \mathcal{R}_{i}^{+} \boxtimes \rightarrow \mathcal{R}_{i}^{+} \mathbb{W}$ by $\phi_{j}^{\prime}=\phi_{j}$ for $j \neq i$ and $\phi_{i}^{\prime}$ to be the restriction of the map

$$
\bigoplus_{a \in Q_{1}^{i}} \phi_{t_{a}}: \bigoplus_{a \in Q_{1}^{i}} V_{t_{a}} \rightarrow \bigoplus_{a \in Q_{1}^{i}} W_{t_{a}}
$$

to the subspace $V_{i}^{\prime}=\operatorname{ker} \xi_{i}$.
Definition 3.39. Let $Q$ be a quiver and let $i \in Q_{0}$ be a source. Notice that for each $a \in Q_{1}^{i}$ we have $t_{a}=i$. We define a covariant functor $\mathcal{R}_{i}^{-}$from the category $\operatorname{rep}_{\mathrm{k}}(Q)$ to the category $\operatorname{rep}_{\mathrm{k}}\left(s_{i} Q\right)$. For every representation $\mathbb{V}=\left(v_{i}, v_{a}\right) \in \operatorname{rep}_{\mathrm{k}}(Q)$, we define a representation $\mathcal{R}_{i}^{-} \mathbb{V}=\left(V_{i}^{\prime}, v_{a}^{\prime}\right)$ in the category $\operatorname{rep}_{\mathrm{k}}\left(\mathrm{s}_{\mathrm{i}} \mathrm{Q}\right)$ as follows. For all $j \neq i$ we set $V_{j}^{\prime}=V_{j}$, and define $V_{i}^{\prime}$ to be the cokernel of the map

$$
\zeta_{i}: V_{i} \rightarrow \bigoplus_{a \in Q_{1}^{i}} V_{h_{a}}, x \mapsto\left(v_{a}(x)\right)_{a \in Q_{1}^{i}} .
$$

For $\mathcal{R}_{i}^{-} \boxtimes$ to be a representation, we need to define linear maps belonging to the arrows in $s_{i} Q$. For each arrow $a \notin Q_{1}^{i}$ we set $v_{a}^{\prime}=v_{a}$. For $a \in Q_{1}^{i}$, let $b$ be the reverse arrow and we define $v_{b}^{\prime}$ to be the composition of maps

$$
V_{t_{b}}^{\prime}=V_{h_{a}}^{\prime}=V_{h_{a}} \hookrightarrow \bigoplus_{c \in Q_{1}^{i}} V_{h_{c}} \rightarrow \operatorname{coker} \zeta_{i}=V_{i}^{\prime}=V_{h_{b}}^{\prime}
$$

Here the map before the direct sum is the canonical inclusion and the map after the direct sum is the canonical quotient map (i.e. the quotient map modulo the image of $\zeta_{i}$ ). Consider the following picture. The red vertices belong to the summands of the space $\bigoplus_{a \in Q_{1}^{i}} V_{h_{a}}$. Outside of the circle, everything stays the same.


Now let $\phi: \mathbb{V} \rightarrow \mathbb{W}$ be a morphism between two representations $\mathbb{V}, \mathbb{W} \in \operatorname{rep}_{\mathrm{k}}(Q)$. We define the morphism $\phi^{\prime}=\mathcal{R}_{i}^{-} \phi: \mathcal{R}_{i}^{-} \mathbb{V} \rightarrow \mathcal{R}_{i}^{-} \mathbb{W}$ by $\phi_{j}^{\prime}=\phi_{j}$ for $j \neq i$ and $\phi_{i}^{\prime}$ to map induced by

$$
\bigoplus_{a \in Q_{1}^{i}} \phi_{h_{a}}: \bigoplus_{a \in Q_{1}^{i}} V_{h_{a}} \rightarrow \bigoplus_{a \in Q_{1}^{i}} W_{h_{a}}
$$

on the quotient space $V_{i}^{\prime}=\operatorname{coker} \zeta_{i}$.
Remark 3.40. Note that we consider the category $\operatorname{rep}_{\mathrm{k}}(Q)$ of finite-dimensional representations. By doing this, we ensure that the sums in the definitions of the maps $\xi_{i}$ and $\zeta_{i}$ are finite sums (resp. there are only finitely many non-zero summands). Note that this is necessary for the welldefinition of the maps $\xi_{i}$ and $\zeta_{i}$ because else the order of summation may not be negligible. To fix this issue, one could also restrict to the case of finite quivers, as we sometimes do. Especially when considering dimension vectors of representations (which is important for Gabriel's Theorem) restricting to finite quivers is the right choice of restriction.

Before giving an example for $\mathcal{R}_{i}^{+}$and $\mathcal{R}_{i}^{-}$, we prove that they are covariant functors.
Lemma 3.41. $\mathcal{R}_{i}^{+}$and $\mathcal{R}_{i}^{-}$are well-defined covariant functors.
Proof. We show that $\mathcal{R}_{i}^{+}$is a functor. Let $\xi_{i}^{W}$ denote the map $\bigoplus_{a \in Q_{1}^{i}} W_{t_{a}} \rightarrow W_{i}$ from the definition of $\mathcal{R}_{i}^{+} \mathbb{W}$. First notice that $\mathcal{R}_{i}^{+} \mathbb{V}$ is indeed a representation of the quiver $s_{i} Q$. This means that for all arrows $b \in s_{i} Q_{1}$ we have that $v_{b}^{\prime}$ maps from $V_{t_{b}}^{\prime}$ to $V_{h_{b}}^{\prime}$. This is true by observing the definition. Similarly, we need to check that $\phi^{\prime}=\mathcal{R}_{i}^{+} \phi$ is a morphism between the representations $\mathcal{R}_{i}^{+} \mathbb{V}$ and $\mathcal{R}_{i}^{+} \mathbb{W}$. First, we observe that all the maps $\phi_{j}^{\prime}$ are well-defined, i.e. $\phi_{j}^{\prime}$ maps from $\mathcal{R}_{i}^{+} \mathbb{V}_{j}$ to $\mathcal{R}_{i}^{+} \mathbb{W}_{j}$. This is clearly true for $j \neq i$ and in the case $j=i$, we notice that for $x=\left(x_{t_{a}}\right)_{a \in Q_{1}^{i}} \in \operatorname{ker} \xi_{i}=V_{i}^{\prime}$ we have:

$$
\begin{aligned}
\xi_{i}^{W}\left(\bigoplus_{a \in Q_{1}^{i}} \phi_{t_{a}}(x)\right) & =\xi_{i}^{W}\left(\left(\phi_{t_{a}}\left(x_{t_{a}}\right)\right)_{a \in Q_{1}^{i}}\right)=\sum_{a \in Q_{1}^{i}} w_{a}\left(\phi_{t_{a}}\left(x_{t_{a}}\right)\right)=\sum_{a \in Q_{1}^{i}} \phi_{h_{a}}\left(v_{a}\left(x_{t_{a}}\right)\right) \\
& =\sum_{a \in Q_{1}^{i}} \phi_{i}\left(v_{a}\left(x_{t_{a}}\right)\right)=\phi_{i} \sum_{a \in Q_{1}^{i}}\left(v_{a}\left(x_{t_{a}}\right)\right)=\phi_{i}(0)=0 .
\end{aligned}
$$

In the third equality we used that $\phi$ is a morphism (i.e. $\phi_{h_{a}} v_{a}=w_{a} \phi_{t_{a}}$ ). Thus we have that $\operatorname{im} \phi_{i}^{\prime} \subseteq \operatorname{ker} \xi_{i}^{W}$ and this means that $\phi_{i}^{\prime}$ is well-defined. We also need to check that $\phi^{\prime}$ is a morphism, i.e. for each $b \in s_{i} Q_{1}$ we have $\phi_{h_{b}}^{\prime} v_{b}^{\prime}=w_{b}^{\prime} \phi_{t_{b}}^{\prime}$. Once again the case for $a \notin Q_{1}^{i}$ is clear since this part of the representations (and morphisms) stays the same (in the picture in the definition, this is the part on the outside of the circle). For $a \in Q_{1}^{i}$ (with inverse arrow b) we have that the following diagram commutes


Indeed, for $x=\left(x_{t_{c}}\right)_{c \in Q_{1}^{i}} \in \operatorname{ker} \xi_{i}=V_{i}^{\prime}$ we have:

$$
w_{b}^{\prime} \phi_{i}^{\prime}(x)=w_{b}^{\prime}\left(\left(\phi_{t_{c}}\left(x_{t_{c}}\right)\right)_{c \in Q_{1}^{i}}\right)=\phi_{t_{a}}\left(x_{t_{a}}\right)=\phi_{h_{b}}^{\prime}\left(x_{t_{a}}\right)=\phi_{h_{b}}^{\prime}\left(v_{b}^{\prime}\left(x_{t_{c}}\right)_{c \in Q_{1}^{i}}\right)=\phi_{h_{b}}^{\prime}\left(v_{b}^{\prime}(x)\right) .
$$

So far, we have shown that $\mathcal{R}_{i}^{+} \boxtimes \in \operatorname{rep}_{\mathrm{k}}\left(s_{i} Q\right)$ and $\mathcal{R}_{i}^{+} \phi: \mathcal{R}_{i}^{+} \boxtimes \rightarrow \mathcal{R}_{i}^{+} \mathbb{W}$ are well-defined. Thus it remains to show that $\mathcal{R}_{i}^{+}$preserves identity morphisms and composition of morphisms. Let $\mathbb{V} \in \operatorname{rep}_{\mathrm{k}}(Q)$ be a representation and let id $\boldsymbol{V}_{\vee}$ be the identity morphism. For $j \neq i$ we have that $\left(\mathcal{R}_{i}^{+} \mathrm{id}_{\vee}\right)_{j}=\operatorname{id}_{\vee_{j}}$ and for $j=i$ we have $\left(\mathcal{R}_{i}^{+} \mathrm{id}_{\vee}\right)_{i}=\left.\bigoplus_{a \in Q_{1}^{i}} \operatorname{id}_{V_{t_{a}}}\right|_{\operatorname{ker} \xi_{i}}=\operatorname{id}_{\operatorname{ker} \xi_{i}}=\mathrm{id}_{V_{i}}^{\prime}$. Thus we have that $\mathcal{R}_{i}^{+} \mathrm{id} \vee=\operatorname{id}_{\mathcal{R}_{i}^{+} \vee}$. Now let $\phi: \mathbb{U} \rightarrow \mathbb{V}$ and $\psi: \mathbb{V} \rightarrow \mathbb{W}$ be morphisms. We want to show that $\mathcal{R}_{i}^{+}(\psi \phi)=\left(\mathcal{R}_{i}^{+} \psi\right)\left(\mathcal{R}_{i}^{+} \phi\right)$. For $j \neq i$ we have that $\left(\mathcal{R}_{i}^{+}(\psi \phi)\right)_{j}=(\psi \phi)_{j}=\psi_{j} \phi_{j}=\left(\mathcal{R}_{i}^{+} \psi\right)_{j}\left(\mathcal{R}_{i}^{+} \phi\right)_{j}$. For $j=i$ we have:

$$
\begin{aligned}
\left(\mathcal{R}_{i}^{+}(\psi \phi)\right)_{i} & =\left.\bigoplus_{a \in Q_{1}^{i}}(\psi \phi)_{t_{a}}\right|_{\operatorname{ker} \xi_{i}^{U}} ^{\operatorname{ker} \xi_{i}^{W}}=\left.\bigoplus_{a \in Q_{1}^{i}} \psi_{t_{a}} \phi_{t_{a}}\right|_{\operatorname{ker} \xi_{i}^{U}} ^{\operatorname{ker} \xi_{i}^{W}}=\left.\left.\bigoplus_{a \in Q_{1}^{i}} \psi_{t_{a}}\right|_{\operatorname{ker} \xi_{i}^{V}} ^{\operatorname{ker} \xi_{i}^{W}} \phi_{t_{a}}\right|_{\operatorname{ker} \xi_{i}^{U}} ^{\operatorname{ker} \xi_{i}^{V}} \\
& =\left(\left.\bigoplus_{a \in Q_{1}^{i}} \psi_{t_{a}}\right|_{\operatorname{ker} \xi_{i}^{V}} ^{\operatorname{ker} \xi_{i}^{W}}\right)\left(\left.\bigoplus_{a \in Q_{1}^{i}} \phi_{t_{a}}\right|_{\operatorname{ker} \xi_{i}^{U}} ^{\operatorname{ker} \xi_{i}^{V}}\right)=\left(\mathcal{R}_{i}^{+} \psi\right)_{i}\left(\mathcal{R}_{i}^{+} \phi\right)_{i}
\end{aligned}
$$

This concludes the proof that $\mathcal{R}_{i}^{+}$is a covariant functor. The proof that $\mathcal{R}_{i}^{-}$is a functor is analogous.

Example 3.42. We consider the quiver $Q$ from Example 3.37. Let $\vee$ be a representation given by

$$
V_{1} \xrightarrow{v_{a}} V_{2} \stackrel{v_{b}}{\leftarrow} V_{3} \stackrel{v_{c}}{\leftarrow} V_{4} \xrightarrow{v_{d}} V_{5} .
$$

By applying $\mathcal{R}_{5}^{+}$to $\mathbb{V}$, we get the representation $\mathcal{R}_{5}^{+} \mathbb{V}$

$$
V_{1} \xrightarrow[a]{v_{a}} V_{2} \stackrel{v_{b}}{\leftarrow} V_{3} \stackrel{v_{c}}{\leftarrow} V_{4} \hookleftarrow \operatorname{ker}\left(v_{d}\right) .
$$

Notice that only the vector space at the vertex 5 and the incident edge between the vertices 4 and 5 change! This is true in general by looking at the pictures in the definition of $\mathcal{R}_{i}^{+}$and $\mathcal{R}_{i}^{-}$. There, only the vector space and the maps inside the (imaginary) circle change. Now we apply $\mathcal{R}_{5}^{-}$to $\mathcal{R}_{5}^{+} \boxtimes$ (the vertex 5 is now a source!) and get:

$$
V_{1} \xrightarrow{v_{a}} V_{2} \stackrel{v_{b}}{\longleftrightarrow} V_{3} \stackrel{v_{c}}{\longleftrightarrow} V_{4} \xrightarrow{\bmod \operatorname{ker}\left(v_{d}\right)} V_{4} / \operatorname{ker}\left(v_{d}\right)
$$

We conclude that $\mathcal{R}_{5}^{-} \mathcal{R}_{5}^{+} \mathbb{V} \cong \mathbb{V}$ if and only $V_{5}=\operatorname{im}\left(v_{d}\right)$. Indeed, by the first isomorphism theorem, if $V_{5}=\operatorname{im}\left(v_{d}\right) \cong V_{4} / \operatorname{ker}\left(v_{d}\right)$, then the vector spaces at vertex 5 coincide and the map $v_{d}$ and the projection map mod $\operatorname{ker}\left(v_{d}\right)$ are the same map! In the case that $\operatorname{im}\left(v_{d}\right) \subsetneq V_{5}$, we have that $V_{5} \cong \operatorname{im}\left(v_{d}\right) \oplus k^{r}$ for some $r \in \mathbb{N}$. Moreover, the projection map mod $\operatorname{ker}\left(v_{d}\right)$ is isomorphic to the restriction of $v_{d}$ to its image $\operatorname{im}\left(v_{d}\right) \cong V_{4} / \operatorname{ker}\left(v_{d}\right)$. Therefore, we have that $\mathbb{V} \cong \mathcal{R}_{5}^{-} \mathcal{R}_{5}^{+} \mathbb{V} \oplus \mathbb{S}_{5}^{r}$, where $\mathbb{S}_{5}$ is the simple representation defined in Example 2.16. Now we apply $\mathcal{R}_{2}^{+}$to $\mathbb{V}$ and get the top row of the following commutative diagram:


The maps $\pi_{1}$ and $\pi_{3}$ denote the canonical projections of $V_{1} \oplus V_{3}$ to its summands. Again we notice that only the vector space at the vertex 2 and the maps belonging to the edges a and b change. Now, by applying $\mathcal{R}_{2}^{-}$to $\mathcal{R}_{2}^{+} \vee$ we get the middle row of the following commutative diagram, where the lower vertical arrow is the canonical map from $V_{1} \oplus V_{3}$ onto $\frac{V_{1} \oplus \vee_{3}}{\operatorname{ker}\left(v_{a}+v_{b}\right)}=\operatorname{coker} \zeta_{2}$. By the first isomorphism theorem, we get an injective map $\phi_{2}: \frac{V_{1} \oplus V_{3}}{\operatorname{ker}\left(v_{a}+v_{b}\right)} \rightarrow V_{2}$.


We thus conclude that $\mathcal{R}_{2}^{-} \mathcal{R}_{2}^{+} \mathbb{V} \cong \mathbb{V}$ if and only if $V_{2}=\operatorname{im}\left(v_{a}+v_{b}\right)$. Indeed, by the first isomorphism theorem, if $V_{2}=\operatorname{im}\left(v_{a}+v_{b}\right) \cong\left(V_{1} \oplus \mathbb{V}_{3}\right) / \operatorname{ker}\left(v_{a}+v_{b}\right)$, then the vector spaces at the vertex 2 coincide. Moreover, the map $V_{1} \rightarrow \frac{V_{1} \oplus V_{3}}{\operatorname{ker}\left(v_{a}+v_{b}\right)}$ in the middle row and the map $v_{a}$ are isomorphic. Also the map $\frac{V_{1} \oplus \vee_{3}}{\operatorname{ker}\left(v_{a}+v_{b}\right)} \leftarrow V_{3}$ in the middle row and the map $v_{b}$ are isomorphic. Therefore, we have that $\mathcal{R}_{2}^{-} \mathcal{R}_{2}^{+} \mathbb{V} \cong \mathbb{V}$. In the case that $\operatorname{im}\left(v_{a}+v_{b}\right) \subsetneq V_{2}$, we have that $V_{2} \cong \operatorname{im}\left(v_{a}+v_{b}\right) \oplus k^{r}$ for some $r \in \mathbb{N}$. Moreover, the maps from $V_{1}$ resp. $V_{3}$ to $\mathcal{R}_{2}^{-} \mathcal{R}_{2}^{+} \mathbb{V} \cong \mathbb{V}$ (in the middle row) are isomorphic the restriction of $v_{a}$ to its image $\operatorname{im}\left(v_{a}\right)$ resp. to the restriction of $v_{b}$ to its image $\operatorname{im}\left(v_{b}\right)$. Therefore, we have $\mathbb{V} \cong \mathcal{R}_{2}^{-} \mathcal{R}_{2}^{+} \mathbb{V} \oplus \mathbb{S}_{2}^{r}$.

### 3.3.1 Composition of Reflection Functors

Since $s_{i} s_{i} Q=Q$ (for a sink or source $i$ ) it is sensible to ask if the functors $\mathcal{R}_{i}^{+}$and $\mathcal{R}_{i}^{-}$are inverse to each other. This is not the case! Indeed, Example 3.42 is a counterexample in the cases where $V_{5} \neq \operatorname{im}\left(v_{d}\right)$ resp. $v_{2} \neq \operatorname{im}\left(v_{a}+v_{b}\right)$. Moreover, one can see that for the simple representation $\mathbb{S}_{i}$ we have $\mathcal{R}_{i}^{+} \mathbb{S}_{i}=0$ (for a $\operatorname{sink} i$ ), respectively $\mathcal{R}_{i}^{-} \mathbb{S}_{i}=0$ (for a source $i$ ). In the following example, we show this for an explicit quiver $Q$, however, the proof for a general quiver $Q$ is analogous.

Example 3.43. We again consider the quiver $Q$ from Example 3.37. From Example 3.42 we get that $\left(\mathcal{R}_{5}^{+} \mathcal{S}_{5}\right)_{5}=\operatorname{ker}\left(v_{d}\right) \subseteq V_{4}$ but $V_{4}=0$. Thus we have that $\mathcal{R}_{5}^{+} \mathcal{S}_{5}=0$. For the sink vertex 2, we get that $\left(\mathcal{R}_{2}^{+} \mathcal{S}_{2}\right)_{2}=\operatorname{ker}\left(v_{a}+v_{b}\right) \subseteq V_{1} \oplus V_{3}=0$ and thus $\mathcal{R}_{2}^{+} \mathcal{S}_{2}=0$.

Even though $\mathcal{R}_{i}^{+}$and $\mathcal{R}_{i}^{-}$are not inverse to each other, we consider their compositions and observe what is missing to be equal to the identity. The following lemma marks a first step towards that.

Lemma 3.44. Let $Q$ be a quiver and let $\boxtimes$ be a finite-dimensional representation of $Q$.

1. Let $i$ be a sink of $Q$. There exists a canonical monomorphism $\iota_{i} \mathbb{V}: \mathcal{R}_{i}^{-} \mathcal{R}_{i}^{+} \mathbb{V} \rightarrow \mathbb{V}$.
2. Let $i$ be a source of $Q$. There exists a canonical epimorphism $\pi_{i} \mathbb{V}: \mathbb{V} \rightarrow \mathcal{R}_{i}^{+} \mathcal{R}_{i}^{-} \mathbb{V}$.

Proof. Since $s_{i} s_{i} Q=Q$ we have that $\mathcal{R}_{i}^{-} \mathcal{R}_{i}^{+} \mathbb{V}, \mathcal{R}_{i}^{+} \mathcal{R}_{i}^{-} \mathbb{V} \in \operatorname{rep}_{\mathrm{k}}\left(s_{i} s_{i} Q\right)=\operatorname{rep}_{\mathrm{k}}(Q)$. We now prove that the maps are indeed morphisms (of representations) with the desired properties (monomorphism resp. epimorphism).

1. Let $\left(\iota_{i} \vee\right)_{j}=\mathrm{id}_{V_{j}}$ for $j \neq i$ and consider

$$
\begin{aligned}
\left(\mathcal{R}_{i}^{-} \mathcal{R}_{i}^{+} \mathbb{V}\right)_{i} & =\operatorname{coker} \zeta_{i}=\operatorname{coker}\left(\operatorname{ker} \xi_{i} \rightarrow \bigoplus_{a \in\left(s_{i} Q\right)_{1}^{i}} V_{h_{a}}^{\prime}\right)=\operatorname{coker}\left(\operatorname{ker} \xi_{i} \hookrightarrow \bigoplus_{a \in Q_{1}^{i}} V_{t_{a}}\right) \\
& =\left(\bigoplus_{a \in Q_{1}^{i}} V_{t_{a}}\right) / \operatorname{ker} \xi_{i} \stackrel{\phi}{=} \operatorname{im} \xi_{i} \subseteq V_{i} .
\end{aligned}
$$

Using the canonical isomorphism $\phi$, we set $\left(\iota_{i} \vee\right)_{i}:\left(\mathcal{R}_{i}^{-} \mathcal{R}_{i}^{+} \vee\right)_{i} \xrightarrow{\phi} \operatorname{im} \xi_{i} \hookrightarrow V_{i}$. First notice that all the maps $\left(\iota_{i} \mathbb{V}\right)_{j}$ are well-defined. Indeed, for $j \neq i$ we have $\left(\mathcal{R}_{i}^{-} \mathcal{R}_{i}^{+} \mathbb{V}\right)_{j}=V_{j}$ and thus $\left(\iota_{i} \mathbb{V}\right)_{j}=\mathrm{id}_{V_{j}}$ is well-defined. For $j=i$, we have that $\left(\iota_{i} \mathbb{V}\right)_{i}$ is well-defined by definition. Notice that all the maps are injective. Therefore, $\iota_{i} \vee$ is a monomorphism, if we can show that it is a morphism. Let us now consider an arrow $a \in Q_{1}$. If $a \notin Q_{1}^{i}$ we have that the following commutative diagram.

$$
\begin{gathered}
\left(\mathcal{R}_{i}^{-} \mathcal{R}_{i}^{+} \mathbb{V}\right)_{t_{a}}=V_{t_{a}} \xrightarrow{v_{a}} V_{h_{a}}=\left(\mathcal{R}_{i}^{-} \mathcal{R}_{i}^{+} \mathbb{V}\right)_{h_{a}} \\
\begin{array}{c}
\operatorname{id}_{V_{t_{a}}} \downarrow \\
V_{t_{a}}
\end{array} \xrightarrow[v_{a}]{\text { id }_{V_{h_{a}}}}
\end{gathered}
$$

For $a \in Q_{1}^{i}$ we have the following diagram.

$$
\begin{gathered}
\left(\mathcal{R}_{i}^{-} \mathcal{R}_{i}^{+} \vee\right)_{t_{a}}=V_{t_{a}} \xrightarrow{f} \operatorname{im} \xi_{i} \cong\left(\mathcal{R}_{i}^{-} \mathcal{R}_{i}^{+} \vee\right)_{i} \\
\operatorname{id}_{V_{t_{a}}} \downarrow \\
V_{t_{a}} \\
\underset{v_{a}}{ } \xrightarrow{\left(\iota_{i} \vee\right)_{i}} \\
\end{gathered}
$$

Using the isomorphism $\left(\mathcal{R}_{i}^{-} \mathcal{R}_{i}^{+} \mathbb{V}\right)_{i} \stackrel{\phi}{\cong} \operatorname{im} \xi_{i},\left(x_{t_{c}}\right)_{c \in Q_{1}^{i}}+\operatorname{ker} \xi_{i} \mapsto \sum_{c \in Q_{1}^{i}} v_{c}\left(x_{t_{c}}\right)$ we get that $f\left(x_{t_{a}}\right)=v_{a}\left(x_{t_{a}}\right)+\sum_{c \in Q_{1}^{i}, c \neq a} v_{c}(0)=v_{a}\left(x_{t_{a}}\right) \in \operatorname{im} \xi_{i}$. Therefore, the above diagram commutes. This shows that $\iota_{i} \vee$ is a morphism (and thus a monomorphism).
2. Let $\left(\pi_{i} \vee\right)_{j}=\operatorname{id}_{V_{j}}$ for $j \neq i$. Now consider

$$
\begin{aligned}
\left(\mathcal{R}_{i}^{+} \mathcal{R}_{i}^{-} \mathbb{V}\right)_{i} & =\operatorname{ker} \xi_{i}=\operatorname{ker}\left(\bigoplus_{a \in\left(s_{i} Q\right)_{1}^{i}} V_{t_{a}}^{\prime} \rightarrow \operatorname{coker} \zeta_{i}\right)=\operatorname{ker}\left(\bigoplus_{a \in Q_{1}^{i}} V_{h_{a}} \rightarrow \operatorname{coker} \zeta_{i}\right) \\
& =\operatorname{ker}\left(\bigoplus_{a \in Q_{1}^{i}} V_{h_{a}} \rightarrow \bigoplus_{a \in Q_{1}^{i}} V_{h_{a}} / \operatorname{im} \zeta_{i}\right) \stackrel{\psi}{\cong} \operatorname{im} \zeta_{i} \subseteq \bigoplus_{a \in Q_{1}^{i}} V_{h_{a}} .
\end{aligned}
$$

Using the canonical isomorphism $\psi$, we set $\left(\pi_{i} \vee\right)_{i}: V_{i} \xrightarrow{\zeta_{i}} \operatorname{im} \zeta_{i} \xrightarrow{\psi^{-1}}\left(\mathcal{R}_{i}^{+} \mathcal{R}_{i}^{-} \mathbb{V}\right)_{i}$. First notice that all the maps $\left(\pi_{i} \vee\right)_{j}$ are well-defined. Indeed, for $j \neq i$ we have $\left(\mathcal{R}_{i}^{+} \mathcal{R}_{i}^{-} \vee\right)_{j}=V_{j}$ and thus $\left(\pi_{i} \vee\right)_{j}=\operatorname{id}_{V_{j}}$ is well-defined. For $j=i$, we have that $\left(\pi_{i} \vee\right)_{i}$ is well-defined by definition. Notice that all the maps are surjective. Therefore, $\pi_{i} \vee$ is an epimorphism, if we can show that it is a morphism. Let us now consider an arrow $a \in Q_{1}$. If $a \notin Q_{1}^{i}$ we have that the following commutative diagram.


For $a \in Q_{1}^{i}$ we have the following diagram.


Using the isomorphism $\left(\mathcal{R}_{i}^{+} \mathcal{R}_{i}^{-} \mathbb{V}\right)_{i} \stackrel{\psi}{\approx} \operatorname{im} \zeta_{i}$ we get that $g\left(\left(v_{c}(x)\right)_{c \in Q_{1}^{i}}\right)=v_{a}(x) \in V_{h_{a}}$ (for $x \in V_{i}$ ). Therefore, the above diagram commutes. This shows that $\pi_{i} \mathbb{V}$ is a morphism (and thus an epimorphism).

Example 3.45. We have already seen examples for the canonical monomorphism $\iota_{i} \bigvee$ in Example 3.42. Indeed, for the sink vertex 5, we have an inclusion $\phi_{5}: V_{4} / \operatorname{ker}\left(v_{d}\right) \hookrightarrow V_{5}$. By setting $\phi_{i}=\operatorname{id}_{V_{i}}$ for $i \in\{1,2,3,4\}$, we get that $\iota_{5} \mathbb{V}=\left(\phi_{i}\right)_{1 \leq i \leq 5}$. For the sink vertex 2, we also have $a$ canonical inclusion at the vertex 2, given by the map $\phi_{2}: \frac{V_{1} \oplus V_{3}}{\operatorname{ker}\left(v_{a}+v_{b}\right)} \rightarrow V_{2}$. Now we set $\phi_{i}=\mathrm{id}_{V_{i}}$ for $i \in\{1,3,4,5\}$ and get that $\iota_{2} \vee=\left(\phi_{i}\right)_{1 \leq i \leq 5}$.

To better understand the morphisms defined in Lemma 3.44, we observe, where $\iota_{i} \vee$ fails to be surjective. Notice that this is described by the representation coker $\left(\iota_{i} \mathbb{V}\right)$. Analogously we observe, where $\pi_{i} \vee$ fails to be injective by looking at the representation $\operatorname{ker}\left(\pi_{i} \vee\right)$.

Lemma 3.46. Let $Q$ be a quiver and let $\boxtimes$ be a finite-dimensional representation of $Q$.

1. Let $i$ be a sink of $Q$. Then the representation $\operatorname{coker}\left(\iota_{i} \vee\right)$ is supported at the vertex $i$. This means that $\operatorname{coker}\left(\iota_{i} \vee\right)_{j}=0$ for all $j \neq i$ and thus all maps $\tilde{v_{a}}=0$ are trivial (for $a \in Q_{1}$ ). Moreover, we have that $\operatorname{coker}\left(\iota_{i} \vee\right)_{i} \cong \operatorname{coker} \xi_{i}$, where $\xi_{i}$ is the map given in the definition of $\mathcal{R}_{i}^{+} \mathrm{V}$.
2. Let $i$ be a source of $Q$. Then the representation $\operatorname{ker}\left(\pi_{i} \mathbb{V}\right)$ is supported at the vertex $i$. This means that $\operatorname{ker}\left(\pi_{i} \mathbb{V}\right)_{j}=0$ for all $j \neq i$ and thus all maps $\tilde{v_{a}}=0$ are trivial (for $\left.a \in Q_{1}\right)$. Moreover, we have that $\operatorname{ker}\left(\pi_{i} \vee\right)_{i} \cong \operatorname{ker} \zeta_{i}$, where $\zeta_{i}$ is the map given in the definition of $\mathcal{R}_{i}^{-} \vee$.
Proof. 1. For $j \neq i$ we have $\left(\iota_{i} \vee\right)_{j}=\operatorname{id}_{V_{j}}$ and thus $\operatorname{coker}\left(\iota_{i} \vee\right)_{j}=V_{j} / V_{j}=0$. This shows that $\operatorname{coker}\left(\iota_{i} \mathbb{V}\right)$ is supported at the vertex $i$. Moreover, in the proof of Lemma 3.44 we showed that we have $\left(\mathcal{R}_{i}^{-} \mathcal{R}_{i}^{+} \mathbb{V}\right)_{i} \cong \operatorname{im} \xi_{i}$. Thus $\operatorname{coker}\left(\iota_{i} \mathbb{V}\right)_{i} \cong V_{i} / \operatorname{im} \xi_{i}=\operatorname{coker} \xi_{i}$.
3. For $j \neq i$ we have $\left(\pi_{i} \vee\right)_{j}=\operatorname{id}_{V_{j}}$ and thus $\operatorname{ker}\left(\pi_{i} \mathbb{V}\right)_{j}=V_{j} / V_{j}=0$. This shows that $\operatorname{ker}\left(\pi_{i} \mathbb{V}\right)$ is supported at the vertex $i$. Moreover, in the proof of Lemma 3.44 we showed that we have $\left(\mathcal{R}_{i}^{+} \mathcal{R}_{i}^{-} \vee\right)_{i} \cong \operatorname{im} \zeta_{i}$. Thus $\operatorname{ker}\left(\pi_{i} \mathbb{V}\right)_{i}=\operatorname{ker}\left(V_{i} \rightarrow \operatorname{im} \zeta_{i}\right)=\operatorname{ker} \zeta_{i}$.

Remark 3.47. Since $\operatorname{coker}\left(\iota_{i} \vee\right)$ and $\operatorname{ker}\left(\pi_{i} \vee\right)$ are supported at the vertex $i$, they are both direct sums of copies of the simple representation $\mathbb{S}_{i}$.

Example 3.48. Again, we consider the quiver from Example 3.42 and the induced monomorphisms $\iota_{2} \mathbb{V}$ and $\iota_{5} \mathbb{V}$ that are described in Example 3.45. Using Lemma 3.46, we get the following representations.

$$
\mathcal{R}_{5}^{-} \mathcal{R}_{5}^{+} \vee: \quad V_{1} \xrightarrow{v_{a}} V_{2} \stackrel{v_{b}}{\longleftrightarrow} V_{3} \stackrel{v_{c}}{\longleftrightarrow} V_{4}^{\bmod \operatorname{ker}\left(v_{d}\right)} \frac{V_{4}}{\operatorname{ker}\left(v_{d}\right)} \cong \operatorname{im}\left(v_{d}\right)
$$

$$
\operatorname{coker}\left(\iota_{5} \mathbb{V}\right): \quad 0 \xrightarrow{0} 0 \longleftarrow{ }^{0} 0 \longleftarrow 0<0 \xrightarrow{0} 0 \xrightarrow{0} \operatorname{coker}\left(v_{d}\right)
$$

We observe that $\mathbb{V} \cong \mathcal{R}_{5}^{-} \mathcal{R}_{5}^{+} \mathbb{V} \oplus \operatorname{coker}\left(\iota_{5} \mathbb{V}\right)$. Similarly, we get the representations.

$$
\mathcal{R}_{2}^{+} \mathcal{R}_{2}^{-} \mathbb{V}: \quad V_{1} \longrightarrow \frac{V_{1} \oplus V_{3}}{\operatorname{cer}\left(v_{a}+v_{b}\right)} \text { ( }
$$

Again, we have $\mathbb{V} \cong \mathcal{R}_{2}^{-} \mathcal{R}_{2}^{+} \mathbb{V} \oplus \operatorname{coker}\left(\iota_{2} \mathbb{V}\right)$. This shows that the representations $\operatorname{coker}\left(\iota_{2} \mathbb{V}\right)$ and $\operatorname{coker}\left(\iota_{5} \vee\right)$ are supported at the vertex 2 resp. 5 and that they are of the form described in Lemma 3.46. In addition, we see that they are summands of $\mathbb{V}$.

We are now finally ready to investigate how $\mathcal{R}_{i}^{+}, \mathcal{R}_{i}^{-}$and the compositions $\mathcal{R}_{i}^{-} \mathcal{R}_{i}^{+}$and $\mathcal{R}_{i}^{+} \mathcal{R}_{i}^{-}$ operate on the summands of a representation $\mathbb{V}$. We have seen before that $\mathcal{R}_{i}^{+}$and $\mathcal{R}_{i}^{-}$are not inverse to each other. However, in Example 3.42, we have already seen that $\mathcal{R}_{i}^{-} \mathcal{R}_{i}^{+} \boxtimes$ (if i is a sink) resp. $\mathcal{R}_{i}^{+} \mathcal{R}_{i}^{-} \vee$ (if i is a source) is a summand of $\mathbb{V}$ and their 'difference' is a direct sum of copies of the simple representation $\mathbb{S}_{i}$. This suggests that $\mathcal{R}_{i}^{-} \mathcal{R}_{i}^{+}$resp. $\mathcal{R}_{i}^{+} \mathcal{R}_{i}^{-}$annihilates the summands isomorphic to $\mathbb{S}_{i}$. Since $\operatorname{coker}\left(\iota_{i} \mathbb{V}\right)$ and $\operatorname{ker}\left(\pi_{i} \mathbb{V}\right)$ are direct sums of copies of $\mathbb{S}_{i}$, we suspect that $\operatorname{coker}\left(\iota_{i} \mathbb{V}\right)$ and $\operatorname{ker}\left(\pi_{i} \mathbb{V}\right)$ are those missing copies of $\mathbb{S}_{i}$ in $\mathcal{R}_{i}^{-} \mathcal{R}_{i}^{+} \mathbb{V}$ resp. $\mathcal{R}_{i}^{+} \mathcal{R}_{i}^{-} \mathbb{V}$. Note that we have already seen this in Example 3.48. It turns out that this is true in general, which is the content of the following theorem.

Theorem 3.49. Let $Q$ be a finite connected quiver and let $\mathbb{V}, \mathbb{W}$ be representations of $Q$. Then for any sink or source $i$ we have:

1. $\mathcal{R}_{i}^{ \pm}(\mathbb{V} \oplus \mathbb{W}) \cong \mathcal{R}_{i}^{ \pm} \mathbb{V} \oplus \mathcal{R}_{i}^{ \pm} \mathbb{W}$.
2. If $i$ is a sink, then we have $\mathbb{V} \cong\left(\mathcal{R}_{i}^{-} \mathcal{R}_{i}^{+} \mathbb{V}\right) \oplus \operatorname{coker}\left(\iota_{i} \mathbb{V}\right)$.
3. If $i$ is a source, then we have $\mathbb{V} \cong\left(\mathcal{R}_{i}^{+} \mathcal{R}_{i}^{-} \mathbb{V}\right) \oplus \operatorname{ker}\left(\pi_{i} \mathbb{V}\right)$.
4. If $\vee$ is finite-dimensional and $\operatorname{coker}\left(\iota_{i} \mathbb{V}\right)=0$, then $\underline{\operatorname{dim}} \mathcal{R}_{i}^{+} \mathbb{V}=\sigma_{i}(\underline{\operatorname{dim} \mathbb{V}})$.
5. If $\mathbb{V}$ is finite-dimensional and $\operatorname{ker}\left(\pi_{i} \mathbb{V}\right)=0$, then $\underline{\operatorname{dim}} \mathcal{R}_{i}^{-} \mathbb{V}=\sigma_{i}(\underline{\operatorname{dim} \mathbb{V}})$.

Proof. 1. We refer to this property as 'preserving direct sums'. We show that it is true for vector spaces, but also for morphisms. For the vector spaces, we observe that for $j \neq i$ we have $\left(\mathcal{R}_{i}^{ \pm}(\mathbb{V} \oplus \mathbb{W})\right)_{j}=(\mathbb{V} \oplus \mathbb{W})_{j}=V_{j} \oplus W_{j}=\mathbb{V}_{j} \oplus \mathbb{W}_{j}=\left(\mathcal{R}_{i}^{ \pm} \mathbb{V}\right)_{j} \oplus\left(\mathcal{R}_{i}^{ \pm} \mathbb{W}\right)_{j}$. For $j=i$ we have

$$
\begin{aligned}
\left(\mathcal{R}_{i}^{+}(\mathbb{V} \oplus \mathbb{W})\right)_{i} & =\operatorname{ker}\left(\xi_{i}^{V \oplus W}: \bigoplus_{a \in Q_{1}^{i}}(\mathbb{V} \oplus \mathbb{W})_{t_{a}} \rightarrow(\mathbb{V} \oplus \mathbb{W})_{i}\right) \\
& =\operatorname{ker}\left(\xi_{i}^{V \oplus W}: \bigoplus_{a \in Q_{1}^{i}} V_{t_{a}} \oplus W_{t_{a}} \rightarrow V_{i} \oplus W_{i}\right) \\
& =\operatorname{ker}\left(\xi_{i}^{V}: \bigoplus_{a \in Q_{1}^{i}} V_{t_{a}} \rightarrow V_{i}\right) \oplus \operatorname{ker}\left(\xi_{i}^{W}: \bigoplus_{a \in Q_{1}^{i}} W_{t_{a}} \rightarrow W_{i}\right)=\left(\mathcal{R}_{i}^{+} \mathbb{V}\right)_{i} \oplus\left(\mathcal{R}_{i}^{+} \mathbb{W}\right)_{i} .
\end{aligned}
$$

Similarly, for $j=i$ we have

$$
\begin{aligned}
& \operatorname{im}\left(\zeta_{i}^{V \oplus W}:(\mathbb{V} \oplus \mathbb{W})_{i} \rightarrow \bigoplus_{a \in Q_{1}^{i}}(\mathbb{V} \oplus \mathbb{W})_{h_{a}}\right)=\operatorname{im}\left(\zeta_{i}^{V \oplus W}: V_{i} \oplus W_{i} \rightarrow \bigoplus_{a \in Q_{1}^{i}} V_{h_{a}} \oplus W_{h_{a}}\right) \\
& \operatorname{im}\left(\zeta_{i}^{V}: V_{i} \rightarrow \bigoplus_{a \in Q_{1}^{i}} V_{h_{a}}\right) \oplus \operatorname{im}\left(\zeta_{i}^{W}: W_{i} \rightarrow \bigoplus_{a \in Q_{1}^{i}} W_{h_{a}}\right) .
\end{aligned}
$$

This yields the equation

$$
\begin{aligned}
\left(\mathcal{R}_{i}^{-}(\mathbb{V} \oplus \mathbb{W})\right)_{i} & =\operatorname{coker} \zeta_{i}^{V \oplus W}=\left(\bigoplus_{a \in Q_{1}^{i}}(\mathbb{V} \oplus \mathbb{W})_{h_{a}}\right) / \operatorname{im} \zeta_{i}^{V \oplus W} \\
& =\left[\left(\bigoplus_{a \in Q_{1}^{i}} V_{h_{a}}\right) / \operatorname{im} \zeta_{i}^{V}\right] \oplus\left[\left(\bigoplus_{a \in Q_{1}^{i}} W_{h_{a}}\right) / \operatorname{im} \zeta_{i}^{W}\right] \\
& =\operatorname{coker} \zeta_{i}^{V} \oplus \operatorname{coker} \zeta_{i}^{W}=\left(\mathcal{R}_{i}^{-} \mathbb{V}\right)_{i} \oplus\left(\mathcal{R}_{i}^{-} \mathbb{W}\right)_{i}
\end{aligned}
$$

We now consider a pair of morphisms $\phi: \mathbb{V}_{1} \rightarrow \mathbb{W}_{1}, \psi: \mathbb{V}_{2} \rightarrow \mathbb{W}_{2}$. For $j \neq i$ we have $\left(\mathcal{R}_{i}^{ \pm}(\phi+\psi)\right)_{j}=(\phi+\psi)_{j}=\phi_{j}+\psi_{j}=\left(\mathcal{R}_{i}^{ \pm} \phi\right)_{j}+\left(\mathcal{R}_{i}^{ \pm} \psi\right)_{j}$. Now for $j=i$ we have

$$
\begin{aligned}
\left(\mathcal{R}_{i}^{+}(\phi+\psi)\right)_{i} & =\left.\bigoplus_{a \in Q_{1}^{i}}(\phi+\psi)_{t_{a}}\right|_{\operatorname{ker} \xi_{i}^{V_{1} \oplus V_{2}}}=\bigoplus_{a \in Q_{1}^{i}} \phi_{t_{a}}+\left.\psi_{t_{a}}\right|_{\operatorname{ker} \xi_{i}^{V_{1} \oplus V_{2}}} \\
& =\left.\bigoplus_{a \in Q_{1}^{i}} \phi_{t_{a}}\right|_{\operatorname{ker} \xi_{i}^{V_{1}}}+\left.\psi_{t_{a}}\right|_{\operatorname{ker} \xi_{i}^{V_{2}}}=\left.\bigoplus_{a \in Q_{1}^{i}} \phi_{t_{a}}\right|_{\operatorname{ker} \xi_{i}^{V_{1}}}+\left.\bigoplus_{a \in Q_{1}^{i}} \psi_{t_{a}}\right|_{\operatorname{ker} \xi_{i}^{V_{2}}} \\
& =\left(\mathcal{R}_{i}^{+} \phi\right)_{i}+\left(\mathcal{R}_{i}^{+} \psi\right)_{i} .
\end{aligned}
$$

Here, we used that $\operatorname{ker} \xi_{i}^{V_{1} \oplus V_{2}}$ splits into $\operatorname{ker} \xi_{i}^{V_{1}} \oplus \operatorname{ker} \xi_{i}^{V_{2}}$. We already showed this in the splitting part of the vector spaces. Similarly for $j=i$ we have

$$
\begin{aligned}
\left(\mathcal{R}_{i}^{-}(\phi+\psi)\right)_{i} & =\left.\bigoplus_{a \in Q_{1}^{i}}(\phi+\psi)_{h_{a}}\right|_{\operatorname{coker} \zeta_{i}^{V_{1} \oplus V_{2}}}=\bigoplus_{a \in Q_{1}^{i}} \phi_{h_{a}}+\left.\psi_{h_{a}}\right|_{\operatorname{coker} \zeta_{i}^{V_{1} \oplus V_{2}}} \\
& =\left.\bigoplus_{a \in Q_{1}^{i}} \phi_{h_{a}}\right|_{\operatorname{coker} \zeta_{i}^{V_{1}}}+\left.\psi_{h_{a}}\right|_{\operatorname{coker} \zeta_{i}^{V_{2}}}=\left.\bigoplus_{a \in Q_{1}^{i}} \phi_{h_{a}}\right|_{\operatorname{coker} \zeta_{i}^{V_{1}}}+\left.\bigoplus_{a \in Q_{1}^{i}} \psi_{h_{a}}\right|_{\operatorname{coker} \zeta_{i}^{V_{2}}} \\
& =\left(\mathcal{R}_{i}^{-} \phi\right)_{i}+\left(\mathcal{R}_{i}^{-} \psi\right)_{i} .
\end{aligned}
$$

In the third equality, we used that the coker $\zeta_{i}^{V_{1} \oplus V_{2}}$ splits into coker $\zeta_{i}^{V_{1}} \oplus \operatorname{coker} \zeta_{i}^{V_{2}}$ (already proven in the splitting part of the vector spaces).
2. To prove the direct sum decomposition, we first need to define two inclusion morphisms $\iota_{i} \mathbb{V}: \mathcal{R}_{i}^{-} \mathcal{R}_{i}^{+} \mathbb{V} \rightarrow \mathbb{V}$ (already defined) and $\rho: \operatorname{coker}\left(\iota_{i} \mathbb{V}\right) \rightarrow \mathbb{V}$. We also need to define two projection morphisms $\pi_{1}: \mathbb{V} \rightarrow \mathcal{R}_{i}^{-} \mathcal{R}_{i}^{+} \mathbb{V}$ and $\pi_{2}: \mathbb{V} \rightarrow \operatorname{coker}\left(\iota_{i} \mathbb{V}\right)$. From Lemma 3.46 we have $\operatorname{coker}\left(\iota_{i} \mathbb{V}\right)_{i} \cong \operatorname{coker} \xi_{i}$. We notice that the projection map $\rho_{i}^{\prime}: V_{i} \rightarrow \operatorname{coker} \xi_{i}=V_{i} /\left(\operatorname{im} \xi_{i}\right)$ has a section, i.e. a map $\rho_{i}: \operatorname{coker} \xi_{i} \rightarrow V_{i}$ s.t. $\rho_{i}^{\prime} \rho_{i}=\operatorname{id}_{\text {coker } \xi_{i}}$. Indeed, let $\left\{v_{j}+\operatorname{im} \xi_{i}\right\}_{j \in I}$ be a basis of $\operatorname{coker}\left(\iota_{i} \vee\right)$. Using the surjectivity of $\rho_{i}^{\prime}$, we can choose $\tilde{v_{j}} \in V_{i}$ s.t. $\rho_{i}^{\prime}\left(\tilde{v_{j}}\right)=v_{j}+\operatorname{im} \xi_{i}$ for all $j \in I$. Then the map induced by $v_{j}+\operatorname{im} \xi_{i} \mapsto \tilde{v_{j}}$ for $j \in I$ is a well-defined linear map and we have that $\rho_{i}^{\prime} \rho_{i}\left(v_{j}+\operatorname{im} \xi_{i}\right)=\rho_{i}^{\prime}\left(\tilde{v_{j}}\right)=v_{j}+\operatorname{im} \xi_{i}$. Using the map $\rho_{i}$, we get a morphism $\rho: \operatorname{coker}\left(\iota_{i} \mathbb{V}\right) \rightarrow \mathbb{V}$ by setting $\rho_{j}=0$ for $j \neq i$ and for $j=i$ we set $\rho_{i}: \operatorname{coker} \xi_{i} \rightarrow V_{i}$. Notice that $\rho$ defines a morphism since the following diagram commutes for every $a \in Q_{1}^{i}$ (for all other edges $a$ from $j_{1}$ to $j_{2}$, the commutativity condition follows directly from $\rho_{j_{1}}=0$ and $\rho_{j_{2}}=0$ ).


We also have projection morphisms $\pi_{1}: \mathbb{V} \rightarrow \mathcal{R}_{i}^{-} \mathcal{R}_{i}^{+} \vee$ and $\pi_{2}: \mathbb{V} \rightarrow \operatorname{coker}\left(\iota_{i} \vee\right)$ defined as follows. In the proof of Lemma 3.44 we have seen that $\left(\mathcal{R}_{i}^{-} \mathcal{R}_{i}^{+} \mathbb{V}\right)_{i} \cong \operatorname{im} \xi_{i} \subseteq V_{i}$. Thus we can view $\mathcal{R}_{i}^{-} \mathcal{R}_{i}^{+} \mathbb{V}$ as a subrepresentation of $\mathbb{V}$ and we can define $\pi_{1}$ to be the projection morphism. We define $\left(\pi_{2}\right)_{j}: V_{j} \rightarrow 0$ to be the trivial map for $j \neq i$ and $\left(\pi_{2}\right)_{i}=\rho_{i}^{\prime}: V_{i} \rightarrow \operatorname{coker} \xi_{i}=V_{i} / \operatorname{im} \xi_{i}$ is the canonical projection map. Now we check that this indeed yields a direct sum decomposition. We have that $\pi_{1} \iota_{i} \vee=\mathrm{id}_{\mathcal{R}_{i}^{-} \mathcal{R}_{i}^{+} \vee}$ since $\mathcal{R}_{i}^{-} \mathcal{R}_{i}^{+} \vee$ can be viewed as a subrepresentation of $\mathbb{V}$. Moreover, we have

$$
\left(\pi_{2} \rho\right)_{j}= \begin{cases}\rho_{i}^{\prime} \rho_{i}=\operatorname{id}_{\text {coker } \xi_{i}} & \text { if } j=i \\ 0=\operatorname{id}_{0}=\operatorname{id}_{\operatorname{coker}\left(\iota_{i} \vee\right)_{j}} & \text { if } j \neq i\end{cases}
$$

This shows that $\pi_{2} \rho=\operatorname{id}_{\text {coker }\left(\iota_{i} \vee\right)}$. It remains to show that $\iota_{i} \vee \pi_{1}+\rho \pi_{2}=\mathrm{id}_{\vee}$. Indeed, for $j \neq i$ we compute $\left(\iota_{i} \vee \pi_{1}\right)_{j}=\operatorname{id}_{V_{j}}$ and $\left(\rho \pi_{2}\right)_{j}=0$. We also have $\left(\iota_{i} \vee \pi_{1}\right)_{i}: V_{i} \rightarrow \operatorname{im} \xi_{i} \hookrightarrow V_{i}$ and $\left(\rho \pi_{2}\right)_{i}: V_{i} \rightarrow \operatorname{coker} \xi_{i} \hookrightarrow V_{i}$. Thus for all $v=v_{1}+v_{2} \in \operatorname{im} \xi_{i} \oplus\left(\operatorname{im} \xi_{i}\right)^{\perp}=V_{i}$ we have

$$
\left(\iota_{i} \mathbb{V} \pi_{1}+\rho \pi_{2}\right)_{i}(v)=\iota_{i}\left(v_{1}\right)+\rho\left(v_{2}+\operatorname{im} \xi_{i}\right)=v_{1}+v_{2}=v .
$$

Therefore, we conclude that $\mathbb{V} \cong\left(\mathcal{R}_{i}^{-} \mathcal{R}_{i}^{+} \mathbb{V}\right) \oplus \operatorname{coker}\left(\iota_{i} \vee\right)$.
3. Again, we need to define two inclusion morphisms and two projection morphisms. We have already defined the projection morphism $\pi_{i} \mathbb{V}: \mathbb{V} \rightarrow \mathcal{R}_{i}^{+} \mathcal{R}_{i}^{-} \mathbb{V}$. Using that $\operatorname{ker}\left(\pi_{i} \mathbb{V}\right)_{i} \cong \operatorname{ker} \zeta_{i}$ (Lemma 3.46), we can view $\operatorname{ker}\left(\pi_{i} \vee\right)$ as a subrepresentation of $\mathbb{V}$ and thus we can define the morphism $\pi_{2}: \mathbb{V} \rightarrow \operatorname{ker}\left(\pi_{i} \mathbb{V}\right)$ to be the canonical projection onto the subrepresentation. We now define two inclusion morphisms $\iota_{1}: \mathcal{R}_{i}^{+} \mathcal{R}_{i}^{-} \mathbb{V} \rightarrow \mathbb{V}$ and $\iota_{2}: \operatorname{ker}\left(\pi_{i} \vee\right) \rightarrow \mathbb{V}$ as follows. We know from the proof of Lemma 3.46 that $\left(\mathcal{R}_{i}^{+} \mathcal{R}_{i}^{-} \mathbb{V}\right)_{i} \cong \operatorname{im} \zeta_{i}$. We claim that the map $\zeta_{i}: V_{i} \rightarrow \operatorname{im} \zeta_{i}$ has a section, i.e. a map $\phi: \operatorname{im} \zeta_{i} \rightarrow V_{i}$ s.t. $\zeta_{i} \phi=\mathrm{id}_{\mathrm{im} \zeta_{i}}$. Similarly, as in part 2, we can define such a map on a basis of $\operatorname{im} \zeta_{i}$ using the surjectivity of $\zeta_{i}$ on the subspace $\operatorname{im} \zeta_{i} \cong\left(\mathcal{R}_{i}^{+} \mathcal{R}_{i}^{-} \mathbb{V}\right)_{i}$. Using the map $\phi$, we get a morphism $\iota_{1}: \mathcal{R}_{i}^{+} \mathcal{R}_{i}^{-} \mathbb{V} \rightarrow \mathbb{V}$ by setting $\left(\iota_{1}\right)_{j}=\operatorname{id}_{V_{j}}$ for $j \neq i$ and for $j=i$, we use the map $\phi$. This defines a morphism since the following diagram commutes for every $a \in Q_{1}^{i}$


Indeed, using that $v_{a}=\pi_{V_{h_{a}}} \zeta_{i}$ (where $\pi_{V_{h_{a}}}: \bigoplus_{c \in Q_{1}^{i}} V_{h_{c}} \rightarrow V_{h_{a}}$ is the projection map), we get that $v_{a} \phi=\pi_{V_{h_{a}}} \zeta_{i} \phi=\pi_{V_{h_{a}}} \operatorname{id}_{\operatorname{im} \zeta_{i}}=\left.\pi_{V_{h_{a}}}\right|_{\mathrm{im} \zeta_{i}}$. Finally, using that $\operatorname{ker}\left(\pi_{i} \vee\right)_{i} \cong \operatorname{ker} \zeta_{i}$, we
can view $\operatorname{ker}\left(\pi_{i} \mathbb{V}\right)$ as a subrepresentation of $\mathbb{V}$ and thus we can define $\iota_{2}: \operatorname{ker}\left(\pi_{i} \mathbb{V}\right) \rightarrow \mathbb{V}$ to be the canonical inclusion morphism. Now we check that these four morphisms indeed yield a direct sum decomposition. We have that

$$
\left(\pi_{i} \mathbb{V}\right)\left(\iota_{1}\right)= \begin{cases}\pi_{i} \mathbb{V} \phi=\zeta_{i} \phi=\operatorname{id}_{\mathrm{im}} \zeta_{i} & \text { if } j=i \\ \operatorname{id}_{V_{j}}\left(\iota_{1}\right)_{j}=\operatorname{id}_{V_{j}} \operatorname{id}_{V_{j}}=\operatorname{id}_{V_{j}} & \text { if } j \neq i\end{cases}
$$

This shows that $\left(\pi_{i} \vee\right)\left(\iota_{1}\right)=\operatorname{id}_{\mathcal{R}_{i}^{+} \mathcal{R}_{i}^{-} \vee}$. We also have that $\pi_{2} \iota_{2}=\operatorname{id}_{\operatorname{ker}\left(\pi_{i} \vee\right)} \operatorname{since} \operatorname{ker}\left(\pi_{i} \vee\right)$ is a subrepresentation of $\mathbb{V}$. It remains to show that $\left(\iota_{1}\right)\left(\pi_{i} \mathbb{V}\right)+\iota_{2} \pi_{2}=\mathrm{id}_{\mathbb{V}}$. Indeed, for $j \neq i$ we compute $\left(\iota_{1}\right)_{j}\left(\pi_{i} \vee\right)_{j}=\operatorname{id}_{V_{j}} \mathrm{id}_{V_{j}}=\mathrm{id}_{V_{j}}$ and $\left(\iota_{2} \pi_{2}\right)_{j}=0$. We have $\left(\iota_{2} \pi_{2}\right)_{i}: V_{i} \rightarrow \operatorname{ker} \zeta_{i} \hookrightarrow V_{i}$ and $\left(\iota_{1}\right)_{i}\left(\pi_{i} \vee\right)_{i}: V_{1} \rightarrow \operatorname{im} \zeta_{i} \hookrightarrow V_{i}$ and thus for all $v=v_{1}+v_{2} \in \operatorname{ker} \zeta_{i} \oplus\left(\operatorname{ker} \zeta_{i}\right)^{\perp}=V_{i}$ we have

$$
\left(\iota_{1}\right)_{i}\left(\pi_{i} \mathbb{V}\right)_{i}(v)+\left(\iota_{2} \pi_{2}\right)_{i}(v)=\phi\left(v_{2}\right)+\iota_{2}\left(v_{1}\right)=v_{2}+v_{1}=v
$$

where we used that $\left(\operatorname{ker} \zeta_{i}\right)^{\perp} \cong V_{i} / \operatorname{ker} \zeta_{i} \cong \operatorname{im} \zeta_{i}$. We conclude that $\mathbb{V} \cong\left(\mathcal{R}_{i}^{+} \mathcal{R}_{i}^{-} \mathbb{V}\right) \oplus \operatorname{ker}\left(\pi_{i} \mathbb{V}\right)$.
4. If $\operatorname{coker}\left(\iota_{i} \mathbb{V}\right)=0$, then we have $\operatorname{coker}\left(\iota_{i} \mathbb{V}\right)_{i} \cong \operatorname{coker} \xi_{i}=0$ and thus we have that $\operatorname{im} \xi_{i}=V_{i}$. Using the dimension formula for finite-dimensional linear maps we get

$$
\begin{aligned}
\operatorname{dim}\left(\mathcal{R}_{i}^{+} \mathbb{V}\right)_{i} & =\operatorname{dim} \operatorname{ker} \xi_{i}=\operatorname{dim}\left(\bigoplus_{a \in Q_{1}^{i}} V_{t_{a}}\right)-\operatorname{dim} \operatorname{im} \xi_{i}=\operatorname{dim}\left(\bigoplus_{a \in Q_{1}^{i}} V_{t_{a}}\right)-\operatorname{dim} V_{i} \\
& =\sum_{a \in Q_{1}^{i}} \operatorname{dim} V_{t_{a}}-\operatorname{dim} V_{i} \stackrel{(1)}{=} \sum_{\substack{j=1 \\
j \neq i}} d_{j i} \operatorname{dim} V_{j}-\operatorname{dim} V_{i}=\sigma_{i}(\underline{\operatorname{dim} \mathbb{V}})_{i},
\end{aligned}
$$

where $d_{j i}$ is the number of arrows between the vertices $i$ and $j$. In Equation (1), we used that there are no loops at $i$ since $i$ is a sink. Therefore, we have that $V_{t_{a}} \neq V_{i}$ for all $a \in Q_{1}^{i}$. For all $j \neq i$ we also have $\operatorname{dim}\left(\mathcal{R}_{i}^{+} \mathbb{V}\right)_{j}=\operatorname{dim} V_{j}=\sigma_{i}(\underline{\operatorname{dim}} \mathbb{V})_{j}$ and thus we have $\underline{\operatorname{dim}} \mathcal{R}_{i}^{+} \mathbb{V}=\sigma_{i}(\underline{\operatorname{dim}} \mathbb{V})$.
5. If $\operatorname{ker}\left(\pi_{i} \vee\right)=0$, we then have $\operatorname{ker}\left(\pi_{i} \mathbb{V}\right)_{i}=\operatorname{ker} \zeta_{i}=0$ and from the dimension formula, we get:

$$
\begin{aligned}
\operatorname{dim}\left(\mathcal{R}_{i}^{-} \mathbb{V}\right)_{i} & =\operatorname{dim} \operatorname{coker} \zeta_{i}=\operatorname{dim}\left(\bigoplus_{a \in Q_{1}^{i}} V_{h_{a}}\right)-\operatorname{dim} \operatorname{im} \zeta_{i} \\
& =\operatorname{dim}\left(\bigoplus_{a \in Q_{1}^{i}} V_{h_{a}}\right)-\left(\operatorname{dim} V_{i}-\operatorname{dim} \operatorname{ker} \zeta_{i}\right)=\sum_{a \in Q_{1}^{i}} \operatorname{dim} V_{h_{a}}-\operatorname{dim} V_{i} \\
& =\sum_{\substack{j=1 \\
j \neq i}} d_{j i} \operatorname{dim} V_{j}-\operatorname{dim} V_{i}=\sigma_{i}(\underline{\operatorname{dim} \mathbb{V}})_{i}
\end{aligned}
$$

For $j \neq i$ we have $\operatorname{dim}\left(\mathcal{R}_{i}^{-} \mathbb{V}\right)_{j}=\operatorname{dim} V_{j}=\sigma_{i}(\underline{\operatorname{dim} \mathbb{V}})_{j}$ and thus we get $\underline{\operatorname{dim}} \mathcal{R}_{i}^{-} \mathbb{V}=\sigma_{i}(\underline{\operatorname{dim} \mathbb{V}})$.

### 3.3.2 Reflection Functors on Indecomposable Representations

We always have a decomposition into indecomposable representations by Theorem 2.33. By part (1) of Theorem 3.49, it suffices to observe how $\mathcal{R}_{i}^{+}$and $\mathcal{R}_{i}^{-}$operate on indecomposable representations.

Corollary 3.50. Let $Q$ be a finite connected quiver and $\boxtimes$ be an indecomposable representation.

1. If $i$ is a sink, then we have two cases:

- If $\mathbb{V} \cong \mathbb{S}_{i}$, then $\mathcal{R}_{i}^{+} \mathbb{V}=0$ and $q_{s_{i} Q}\left(\underline{\operatorname{dim}} \mathcal{R}_{i}^{+} \mathbb{V}\right)=0$.
- If $\mathbb{V} \nsubseteq \mathbb{S}_{i}$, then $\mathcal{R}_{i}^{+} \mathbb{V}$ is non-zero and indecomposable, $\mathcal{R}_{i}^{-} \mathcal{R}_{i}^{+} \mathbb{V} \cong \mathbb{V}$ and the dimension vectors of $\mathbb{V}$ and $\mathcal{R}_{i}^{+} \mathbb{V}$ are related by $q_{s_{i} Q}\left(\underline{\operatorname{dim}} \mathcal{R}_{i}^{+} \mathbb{V}\right)=q_{Q}(\underline{\operatorname{dim} \mathbb{V}})$.

2. If $i$ is a source, then we have two cases:

- If $\mathbb{V} \cong \mathbb{S}_{i}$, then $\mathcal{R}_{i}^{-} \mathbb{V}=0$ and $q_{s_{i} Q}\left(\underline{\operatorname{dim} \mathcal{R}_{i}^{-}} \mathbb{V}\right)=0$.
- If $\mathbb{V} \nsubseteq \mathbb{S}_{i}$, then $\mathcal{R}_{i}^{-} \mathbb{V}$ is non-zero and indecomposable, $\mathcal{R}_{i}^{+} \mathcal{R}_{i}^{-} \mathbb{V} \cong \mathbb{V}$ and the dimension vectors of $\mathbb{V}$ and $\mathcal{R}_{i}^{-} \mathbb{V}$ are related by $q_{s_{i} Q}\left(\underline{\operatorname{dim}} \mathcal{R}_{i}^{-} \mathbb{V}\right)=q_{Q}(\underline{\operatorname{dim}} \mathbb{V})$.

Proof. If $\mathbb{V} \cong \mathbb{S}_{i}$, we have already seen that $\mathcal{R}_{i}^{+} \mathbb{V}=0$, respectively $\mathcal{R}_{i}^{-} \mathbb{V}=0$. Then clearly $q_{s_{i} Q}\left(\underline{\operatorname{dim}} \mathcal{R}_{i}^{ \pm} \mathbb{V}\right)=q_{s_{i} Q}(0)=0$. Now let $\mathbb{V} \nsubseteq \mathbb{S}_{i}$. If $i$ is a sink, $\mathbb{V} \cong\left(\mathcal{R}_{i}^{-} \mathcal{R}_{i}^{+} \mathbb{V}\right) \oplus \operatorname{coker}\left(\iota_{i} \mathbb{V}\right)$. Since $\mathbb{V}$ is irreducible, one of those summands needs to be trivial and since coker $\left(\iota_{i} \mathbb{V}\right)$ is isomorphic to a direct sum of $\mathbb{S}_{i}$, we conclude that $\operatorname{coker}\left(\iota_{i} \mathbb{V}\right)=0$ and $\mathbb{V} \cong \mathcal{R}_{i}^{-} \mathcal{R}_{i}^{+} \mathbb{V}$. This also proves that $\mathcal{R}_{i}^{+} \mathbb{V}$ is indecomposable. Indeed, if $\mathcal{R}_{i}^{+} \mathbb{V}=\mathbb{W}_{1} \oplus \mathbb{W}_{2}$, then $\mathbb{V} \cong \mathcal{R}_{i}^{-} \mathbb{W}_{1} \oplus \mathcal{R}_{i}^{-} \mathbb{W}_{2}$ is not indecomposable. Moreover, we have $q_{s_{i} Q}\left(\underline{\operatorname{dim}} \mathcal{R}_{i}^{+} \mathbb{V}\right)=q_{s_{i} Q}\left(\sigma_{i}(\underline{\operatorname{dim}} \mathbb{V})\right)=q_{Q}(\underline{\operatorname{dim} \mathbb{V}})$ (using Lemma 3.23). Analogously, if $i$ is a source, we use that $\mathbb{V} \cong\left(\mathcal{R}_{i}^{+} \mathcal{R}_{i}^{-} \mathbb{V}\right) \oplus \operatorname{ker}\left(\pi_{i} \mathbb{V}\right)$ and that $\operatorname{ker}\left(\pi_{i} \mathbb{V}\right)$ is isomorphic to a direct sum of $\mathbb{S}_{i}$ and thus $\operatorname{ker}\left(\pi_{i} \mathbb{V}\right)=0$ and $\mathbb{V} \cong \mathcal{R}_{i}^{+} \mathcal{R}_{i}^{-} \mathbb{V}$. Using Lemma 3.23, we get that $q_{s_{i} Q}\left(\underline{\operatorname{dim}} \mathcal{R}_{i}^{-} \mathbb{V}\right)=q_{s_{i} Q}\left(\sigma_{i}(\underline{\operatorname{dim} \mathbb{V}})\right)=q_{Q}(\underline{\operatorname{dim} \mathbb{V}})$.

Theorem 3.49 and Corollary 3.50 show that $\mathcal{R}_{i}^{+}$and $\mathcal{R}_{i}^{-}$are inverse everywhere except on the simple representation $\mathbb{S}_{i}$. To state this precisely, we consider the subcategory consisting of representations, which do not have a summand equal to $\mathbb{S}_{i}$.

Definition 3.51. Let $Q$ be a quiver and let $i \in Q_{0}$ be any vertex. We denote by $\operatorname{Rep}_{\mathrm{k}}(Q)\langle i\rangle$ the full subcategory of $\operatorname{Rep}_{\mathrm{k}}(Q)$ consisting of representations that do not have summands that are isomorphic to the simple representation $\mathbb{S}_{i}$. Analogously, we define the full subcategory $\operatorname{rep}_{\mathrm{k}}(Q)\langle i\rangle$ of $\operatorname{rep}_{\mathrm{k}}(Q)$.
Corollary 3.52. Let $Q$ be a quiver and let $i \in Q_{0}$ be a sink or a source. Then $\mathcal{R}_{i}^{+}$and $\mathcal{R}_{i}^{-}$ induce mutually inverse equivalences $\operatorname{rep}_{\mathrm{k}}(Q)\langle i\rangle \rightarrow \operatorname{rep}_{\mathrm{k}}\left(s_{i} Q\right)\langle i\rangle$ and $\operatorname{rep}_{\mathrm{k}}\left(s_{i} Q\right)\langle i\rangle \rightarrow \operatorname{rep}_{\mathrm{k}}(Q)\langle i\rangle$. In particular, there is a one-to-one correspondence between the isomorphism classes of indecomposable representations of $Q$ and those of $s_{i} Q$.

Proof. The first part follows directly from Corollary 3.50 and Theorem 3.49. Indeed, for a representation $\mathbb{V}=\mathbb{V}_{1} \oplus \ldots \oplus \mathbb{V}_{r} \in \operatorname{Rep}_{\mathrm{k}}(Q)\langle i\rangle$, where $\mathbb{V}_{i}$ are the indecomposable summands that are all non-isomorphic to $\mathbb{S}_{i}$. If $i$ is a sink, we have:

$$
\mathcal{R}_{i}^{-} \mathcal{R}_{i}^{+} \mathbb{V} \cong\left(\mathcal{R}_{i}^{-} \mathcal{R}_{i}^{+} \mathbb{V}_{1}\right) \oplus \ldots \oplus\left(\mathcal{R}_{i}^{-} \mathcal{R}_{i}^{+} \mathbb{V}_{r}\right) \cong \mathbb{V}_{1} \oplus \ldots \oplus \mathbb{V}_{r}=\mathbb{V}
$$

Similarly, if $i$ is a source, we have:

$$
\mathcal{R}_{i}^{+} \mathcal{R}_{i}^{-} \mathbb{V} \cong\left(\mathcal{R}_{i}^{+} \mathcal{R}_{i}^{-} \mathbb{V}_{1}\right) \oplus \ldots \oplus\left(\mathcal{R}_{i}^{+} \mathcal{R}_{i}^{-} \mathbb{V}_{r}\right) \cong \mathbb{V}_{1} \oplus \ldots \oplus \mathbb{V}_{r}=\mathbb{V}
$$

The one-to-one correspondence between the isomorphism classes of indecomposable representations in $\operatorname{Rep}_{\mathrm{k}}(Q)$ and $\operatorname{Rep}_{\mathrm{k}}\left(s_{i} Q\right)$ is given by:

$$
S_{i} \mapsto \mathbb{S}_{i}, \mathbb{V} \mapsto \begin{cases}\mathcal{R}_{i}^{+} \mathbb{V} & \text { if } i \text { is a sink } \\ \mathcal{R}_{i}^{-} \mathbb{V} & \text { if } i \text { is a source }\end{cases}
$$

where $\mathbb{V} \in \operatorname{Rep}_{\mathrm{k}}(Q)$ is indecomposable and $\mathbb{V} \neq \mathbb{S}_{i}$.
The subcategories $\operatorname{Rep}_{\mathrm{k}}(Q)\langle i\rangle$ and $\operatorname{rep}_{\mathrm{k}}(Q)\langle i\rangle$ are interesting because on those categories, $\mathcal{R}_{i}^{-} \mathcal{R}_{i}^{+}$and $\mathcal{R}_{i}^{+} \mathcal{R}_{i}^{-}$are the identity (up to isomorphism of representations). Moreover, if we classify all indecomposable representations in $\operatorname{Rep}_{\mathrm{k}}(Q)\langle i\rangle$, we also classify all indecomposable representations in $\operatorname{Rep}_{\mathrm{k}}(Q)$ since the only additional indecomposable representation is $\mathbb{S}_{i}$. However, the subcategory $\operatorname{Rep}_{\mathrm{k}}(Q)\langle i\rangle$ has the crucial disadvantage that morphisms in $\operatorname{Rep}_{\mathrm{k}}(Q)\langle i\rangle$ do (in general) not have kernels and cokernels. This means that for a morphism $\phi: \mathbb{V} \rightarrow \mathbb{W}$ between representations $\mathbb{V}, \mathbb{W} \in \operatorname{Rep}_{\mathrm{k}}(Q)\langle i\rangle$ it can happen that one (or both) of the representations ker $\phi$ or coker $\phi$ is not in $\operatorname{Rep}_{\mathrm{k}}(Q)\langle i\rangle$, meaning that the representation $\operatorname{ker} \phi$ or coker $\phi$ has a summand that is isomorphic to $\mathbb{S}_{i}$.

Example 3.53. We give an example, illustrating that the category $\operatorname{Rep}_{\mathrm{k}}(Q)\langle i\rangle$ does (in general) not have kernels and cokernels. Let $Q$ be the linear $L_{3}$-quiver and let $\phi: \mathbb{V} \rightarrow \mathbb{W}$ be the morphism represented by the following commutative diagram


Note that the representations $\mathbb{V}$ and $\mathbb{W}$ are indecomposable and thus do not have a summand that is isomorphic to a simple representation $\mathbb{S}_{i}$. The indecomposability of $\mathbb{V}$ and $\mathbb{W}$ is proven in Section 4.1. However, the representation coker $\phi$ is given by the following picture

$$
k \xrightarrow{0} 0 \xrightarrow{0} 0 .
$$

Therefore, $\operatorname{coker} \phi=\mathbb{S}_{1}$ which shows that $\operatorname{Rep}_{\mathrm{k}}(Q)\langle i\rangle$, in general, does not have cokernels. Analogously, we consider the 'reverse diagram':


Here, we get $\operatorname{ker} \psi=\mathbb{S}_{3}$, showing that $\operatorname{Rep}_{\mathrm{k}}(Q)\langle i\rangle$ does not necessarily have kernels.

### 3.4 Coxeter Functors

In this section, we introduce a notion of ordered vertices and we introduce the sequence of reflection functors belonging to this ordering. This notion of ordering coincides with the one given in [2], but is different from the one given in [1].

### 3.4.1 Admissible Orderings

We introduce a notion of a preferred ordering and show for which quivers such an ordering exists. The existence of such an ordering is closely related to the existence of cycles in the quiver.

Definition 3.54. An ordering $i_{1}, i_{2}, \ldots, i_{n}$ of the vertices of a quiver $Q$ is called (+)-admissible if $i_{1}$ is a sink in $Q$ and for each $2 \leq t \leq n$, $i_{t}$ is a sink in $s_{i_{t-1}} \ldots s_{i_{1}} Q$. Dually, an ordering $i_{1}, i_{2}, \ldots, i_{n}$ is called (-)-admissible if $i_{1}$ is a source in $Q$ and for each $2 \leq t \leq n$, $i_{t}$ is a source in $s_{i_{t-1}} \ldots s_{i_{1}} Q$. It is required that each vertex of $Q$ appears exactly once in our ordering.

Example 3.55. We consider the following quiver $Q$ :


The ordering $1,4,2,3,5$ is $(+)$-admissible. Indeed, 1 is a sink in $Q$, 4 is a sink in $s_{1} Q$, 2 is a sink in $s_{4} s_{1} Q$, 3 is a sink in $s_{2} s_{4} s_{1} Q$ and 5 is a sink in $s_{3} s_{2} s_{4} s_{1} Q$. One can also check that $3,2,1,5,4$ is a (-)-admissible ordering.

Lemma 3.56. There exists a (+)-admissible ordering of the vertices of a quiver $Q$ if and only if $Q$ is acyclic, i.e. there are no oriented cycles in $Q$.

Proof. We prove one implication by induction on the number of vertices. Let n be the number of vertices of $Q$ and suppose that $Q$ is acyclic. Since the base case $n=1$ is clear, we can assume that $n \geq 2$. Now let $i_{n}$ be the starting vertex of an oriented path of maximal length. Since $Q$ is acyclic there are no infinite ordered paths and thus such a path exists. Then $i_{n}$ is a source, else we have a bigger ordered path. Now we remove $i_{n}$ from the quiver (we also remove all its incident edges). By induction, the smaller quiver has an admissible ordering $i_{1}, \ldots, i_{n-1}$. Now $i_{1}, \ldots, i_{n-1}, i_{n}$ is an admissible ordering of the original quiver $Q$. Indeed, each incident edge $a=\left(i_{n}, i_{t}\right)$ once reverses direction by applying $s_{i_{t}}$ and then reverses the direction back by applying $s_{i_{n}}$. Therefore, $i_{t}$ is also a sink in $s_{i_{t-1}} \ldots s_{i_{1}} Q$ (notice that here, we apply the $s_{i_{j}}$ to the whole quiver $Q$ ) and also $i_{n}$ is a $\operatorname{sink}$ in $s_{i_{n-1}} \ldots s_{i_{1}} Q$.

In the induction step of the proof, we have chosen a vertex $i_{n}$. This hints at the fact that there is generally not a unique $(+) /(-)$-admissible ordering. A possible algorithm to choose a $(+)-$ admissible ordering is the following. For an acyclic quiver $Q$, there always exists a sink $i_{1}$. Now we can delete the vertex $i_{1}$ (and its incident edges) and the resulting quiver is again acyclic, which allows us to pick a sink $i_{2}$. Doing this inductively yields an $(+)$-admissible ordering $i_{1}, \ldots, i_{n}$.

Lemma 3.57. Let $i_{1}, i_{2}, \ldots, i_{n}$ be an ordering of the vertices of a quiver $Q$. Then $i_{1}, i_{2}, \ldots, i_{n}$ is $(+)$-admissible if and only if $i_{n}, i_{n-1}, \ldots, i_{1}$ is (-)-admissible. In particular, a quiver $Q$ admits a $(+)$-admissible ordering if and only if it admits an (-)-admissible ordering. Therefore, we often just refer to an admissible ordering, meaning a $(+)$-admissible ordering with induced $(-)$-admissible ordering.

Proof. Let $i_{1}, \ldots, i_{n}$ be (+)-admissible. Then, $i_{n}$ is a sink in $s_{i_{n-1}} \ldots s_{i_{1}} Q$. But in the process of going from $Q$ to $s_{i_{n-1}} \ldots s_{i_{1}} Q$, every incident arrow of $i_{n}$ is flipped once and thus $i_{n}$ is a source in $Q$. Similarly, because $i_{n-1}$ is a sink in $s_{i_{n-2}} \ldots s_{i_{1}} Q$, we know the following for all $1 \leq k \leq n-2$. If $a$ is an arrow (in $Q$ ) between $i_{n-1}$ and $i_{k}$, then $a$ must point away from $i_{n-1}$ (since it points towards $i_{n-1}$ in $s_{i_{n-1}} \ldots s_{i_{1}} Q$ ). Therefore, $a$ also points away from $i_{n-1}$ in $s_{i_{n}} Q$. In addition, if $b$ is an arrow (in $Q$ ) between $i_{n-1}$ and $i_{n}$, then $b$ must point away from $i_{n-1}$ in $s_{i_{n}} Q$ since $i_{n}$ is a source in $Q$. Repeating this argument inductively yields that $i_{n}, i_{n-1}, \ldots, i_{1}$ is $(-)$-admissible. The other implication can be proven analogously.

Example 3.58. We consider the quiver from Example 3.55. We observe that the sequence $5,3,2,4,1$ is a (-)-admissible ordering and that $4,5,1,2,3$ is a (+)-admissible ordering. In addition, this also shows that there can be multiple admissible orderings.

Remark 3.59. There is another way to think about admissible orderings. If a quiver $Q$ is acyclic, then $Q$ represents a partial order on the vertex set $Q_{0}$. Indeed, for two vertices $i$ and $j$, we can define $i \leq j$ if and only if there exists a directed path from $i$ to $j$ (we say that there is a path from each vertex $i$ to itself). This indeed defines a partial order since there are no oriented cycles in $Q$. Then we can choose a total order that extends our partial order and this yields an ordering of the vertices. It can be shown that this ordering is (-)-admissible. This point of view helps to understand why $Q$ needs to be acyclic for it to admit an admissible ordering. Indeed, if $Q$ would have an oriented cycle, then we would get a chain $i_{1} \leq i_{2} \leq \ldots \leq i_{r} \leq i_{1}$ which would imply that $i_{1}=i_{2}=\ldots=i_{r}$. Notice that an increasing ordering (with respect to our total order) yields a $(-)$-admissible ordering and a decreasing ordering yields a $(+)$-admissible ordering.

Example 3.60. We consider the quiver from Example 3.55. As explained in the remark above, we get the following partial order: $2<1,2<4,3<2,3<4,5<4$. We see that the $(+)$-admissible orderings $1,4,2,3,5$ and $4,5,1,2,3$ are total orders that extend our partial order (the total orderings are in decreasing order).

We notice that each arrow in $s_{i_{n}} \ldots s_{i_{1}} Q$ results from an arrow in $Q$ by changing its orientation exactly twice and therefore, we have that $s_{i_{n}} \ldots s_{i_{1}} Q=Q$. This motivates the following definition.
Definition 3.61. Let $Q$ be an acyclic quiver and let $i_{1}, \ldots, i_{n}$ be $a(+)$-admissible ordering of the vertices of $Q$. The two following functors are called Coxeter functors (with respect to this ordering).

$$
\begin{aligned}
& \mathcal{C}^{+}=\mathcal{R}_{i_{n}}^{+} \ldots \mathcal{R}_{i_{1}}^{+}: \operatorname{Rep}_{\mathrm{k}}(Q) \rightarrow \operatorname{Rep}_{\mathrm{k}}(Q) \\
& \mathcal{C}^{-}=\mathcal{R}_{i_{1}}^{-} \ldots \mathcal{R}_{i_{n}}^{-}: \operatorname{Rep}_{\mathrm{k}}(Q) \rightarrow \operatorname{Rep}_{\mathrm{k}}(Q)
\end{aligned}
$$

Lemma 3.62. The Coxeter functors $\mathcal{C}^{+}$and $\mathcal{C}^{-}$do not depend on the choice of the admissible ordering of the vertices of a quiver $Q$.

Proof. First, we claim that $\mathcal{R}_{i}^{+} \mathcal{R}_{j}^{+}=\mathcal{R}_{j}^{+} \mathcal{R}_{i}^{+}$if $i$ and $j$ are sinks that are not adjacent to each other. Indeed, notice that $\mathcal{R}_{i}^{+}$and $\mathcal{R}_{j}^{+}$change different vector spaces and linear maps, and therefore the order of operation does not matter. Now let $i_{1}, \ldots, i_{n}$ and $i_{1}^{\prime}, \ldots, i_{n}^{\prime}$ be two ( + )-admissible orderings of the vertices of $Q$. Let $i_{1}=i_{m}^{\prime}$. Then $i_{1}$ is a sink in $Q$ and in $s_{i_{m-1}^{\prime}} \ldots s_{i_{1}^{\prime}} Q$. Therefore, $i_{1}$ is not adjacent to any of the vertices $i_{1}^{\prime}, \ldots, i_{m-1}^{\prime}$ and we get that:

$$
\mathcal{R}_{i_{m}^{\prime}}^{+} \mathcal{R}_{i_{m-1}^{\prime}}^{+} \ldots \mathcal{R}_{i_{1}^{\prime}}^{+}=\mathcal{R}_{i_{m-1}^{\prime}}^{+} \mathcal{R}_{i_{m}^{\prime}}^{+} \ldots \mathcal{R}_{i_{1}^{\prime}}^{+}=\ldots=\mathcal{R}_{i_{m-1}^{\prime}}^{+} \ldots \mathcal{R}_{i_{1}^{\prime}}^{+} \mathcal{R}_{i_{m}^{\prime}}^{+}=\mathcal{R}_{i_{m-1}^{\prime}}^{+} \ldots \mathcal{R}_{i_{1}^{\prime}}^{+} \mathcal{R}_{i_{1}}^{+}
$$

Now let $i_{2}=i_{k}^{\prime}$. We denote by $\mathcal{R}_{i_{k}^{\prime}}^{+} \ldots \widehat{\mathcal{R}_{i_{1}}^{+}} \ldots \mathcal{R}_{i_{1}^{\prime}}^{+}$the sequence $\mathcal{R}_{i_{k}^{\prime}}^{+} \ldots \mathcal{R}_{i_{1}^{\prime}}^{+}$with $\mathcal{R}_{i_{1}}^{+}$omitted. If $\mathcal{R}_{i_{1}}^{+}$does not occur in the sequence, we just leave the sequence as is. We know that $i_{2}$ is a sink in $s_{i_{1}} Q$ and in $s_{i_{k-1}^{\prime} \ldots s_{i_{1}^{\prime}}} Q$. Therefore, $i_{2}$ is not adjacent to any of the vertices $i_{k-1}^{\prime}, \ldots, i_{1}^{\prime}$ apart from maybe $i_{1}=i_{m}^{\prime}$. We do not know if $i_{1}$ is a vertex in the list of $i_{k-1}^{\prime}, \ldots, i_{1}^{\prime}$, however, by possibly excluding $\mathcal{R}_{i_{1}}^{+}$we get that

$$
\mathcal{R}_{i_{k}^{\prime}}^{+} \ldots \widehat{\mathcal{R}_{i_{1}}^{+}} \ldots \mathcal{R}_{i_{1}^{\prime}}^{+}=\mathcal{R}_{i_{k-1}^{\prime}}^{+} \ldots \widehat{\mathcal{R}_{i_{1}}^{+}} \ldots \mathcal{R}_{i_{1}^{\prime}}^{+} \mathcal{R}_{i_{2}}^{+}
$$

Therefore, we can first move $\mathcal{R}_{i_{1}}^{+}$to the right and then move $\mathcal{R}_{i_{2}}^{+}$to the right and get

$$
\mathcal{R}_{i_{n}^{\prime}}^{+} \ldots \mathcal{R}_{i_{1}^{\prime}}^{+}=\mathcal{R}_{i_{n}^{\prime}}^{+} \ldots \widehat{\mathcal{R}_{i_{1}}^{+}} \ldots \widehat{\mathcal{R}_{i_{2}}^{+}} \ldots \mathcal{R}_{i_{1}^{\prime}}^{+} \mathcal{R}_{i_{2}}^{+} \mathcal{R}_{i_{1}}^{+}
$$

Doing this inductively yields

$$
\mathcal{R}_{i_{n}^{\prime}}^{+} \ldots \mathcal{R}_{i_{1}^{\prime}}^{+}=\mathcal{R}_{i_{n}}^{+} \ldots \mathcal{R}_{i_{1}}^{+} .
$$

The case of the Coxeter functor $\mathcal{C}^{-}$can be proven analogously.
Since only acyclic quivers admit admissible orderings, we restrict our attention to acyclic quivers in the rest of this section. Moreover, in the acyclic case, we have an admissible ordering and because the Coxeter functors $\mathcal{C}^{+}$and $\mathcal{C}^{-}$are independent of the choice of admissible ordering, we assume w.l.o.g. that the vertices in $Q_{0}=\{1, \ldots, n\}$ already depict an admissible ordering, i.e. that $1, \ldots, n$ is a $(-)$-admissible ordering. Now we use the results about reflection functors to better understand the Coxeter functors.

Corollary 3.63. Let $Q$ be a finite, connected, and acyclic quiver and let $\mathcal{C}^{ \pm}$be the Coxeter functors. Then for any indecomposable representation $\boxtimes$ of $Q$, either $\mathcal{C}^{ \pm}$is indecomposable or $\mathcal{C}^{ \pm} \bigvee=0$. In the first case, we have $q_{Q}\left(\underline{\operatorname{dim}} \mathcal{C}^{ \pm} \mathbb{V}\right)=q_{Q}(\underline{\operatorname{dim}} \mathbb{V})$, while in the second case we have $q_{Q}\left(\underline{\operatorname{dim}} \mathcal{C}^{ \pm} \mathbb{V}\right)=0$.
Proof. This proof directly follows from Corollary 3.50. Since the case $\mathcal{C}^{ \pm} \vee=0$ is trivial we can w.l.o.g. assume that $\mathcal{C}^{ \pm} \boxtimes 0$. From Corollary 3.50 we know that $\mathcal{R}_{i_{1}}^{+} \mathbb{V}$ is indecomposable. Using this argument inductively shows that $\mathcal{C}^{+} \vee$ is indecomposable. Similarly, we can conclude that $\mathcal{C}^{-} \vee$ is indecomposable. Moreover, we have

$$
\begin{aligned}
q_{Q}\left(\underline{\operatorname{dim}} \mathcal{C}^{+} \mathbb{V}\right) & =q_{s_{i_{n}} \ldots s_{i_{1} Q}}\left(\underline{\operatorname{dim}} \mathcal{R}_{i_{n}}^{+} \ldots \mathcal{R}_{i_{1}}^{+} \mathbb{V}\right)=q_{s_{i_{n-1}} \ldots s_{i_{1} Q}}\left(\underline{\operatorname{dim}} \mathcal{R}_{i_{n-1}}^{+} \ldots \mathcal{R}_{i_{1}}^{+} \mathbb{V}\right) \\
& =\ldots=q_{s_{i_{1}} Q}\left(\underline{\operatorname{dim}} \mathcal{R}_{i_{1}}^{+} \mathbb{V}\right)=q_{Q}(\underline{\operatorname{dim} \mathbb{V}})
\end{aligned}
$$

Analogously, it follows that $q_{Q}\left(\underline{\operatorname{dim}} \mathcal{C}^{-} \mathbb{V}\right)=q_{Q}(\underline{\operatorname{dim} \mathbb{V}})$.
Example 3.64. To see an example of the Coxeter functor, we refer to Section 4.2. There, it is demonstrated how the Coxeter functor operates on the linear $L_{n}$-quiver. Explicitly, it is calculated, what happens to the dimension vector of an indecomposable representation, when applying $\mathcal{C}^{+}$.
Definition 3.65. Let $Q$ be an acyclic quiver and fix a (+)-admissible ordering $i_{1}, \ldots, i_{n}$ of its vertices. The automorphism $c=\sigma_{i_{n}} \ldots \sigma_{i_{1}} \in W(Q) \subseteq \operatorname{Aut}\left(\mathbb{Z}^{n}\right)$ is called a Coxeter transformation.

The Coxeter transformation and the Coxeter functor are related by the following property.
Lemma 3.66. Let $Q$ be an acyclic quiver and let $\mathcal{C}^{+}$be the Coxeter functor and $c$ be the Coxeter transformation (w.r.t. any admissible ordering). Let $\boxtimes$ be an indecomposable representation of $Q$ that is not isomorphic to the simple representation $\mathbb{S}_{i}$ for any $i \in Q_{0}$. We then have that $\underline{\operatorname{dim}}\left(\mathcal{C}^{+}\right)^{r} \mathbb{V}=c^{r}(\underline{\operatorname{dim}} \mathbb{V})$.

Proof. This is an immediate consequence of $\underline{\operatorname{dim}} \mathcal{R}_{i}^{+} \vee=\sigma_{i}(\underline{\operatorname{dim}}) \vee$ from Theorem 3.49.
Remark 3.67. The Coxeter transformation is independent of the choice of admissible ordering. However, since the Coxeter functor $\mathcal{C}^{+}$is independent of the choice of admissible ordering and we only use the above property of the Coxeter transformation, it is not important that the Coxeter transformation is independent of the choice of admissible ordering. Therefore, we do not prove it.

Lemma 3.68. Let $Q$ be a Dynkin quiver.

1. The Coxeter transformation $c$ has no non-zero fixed vectors, that is $c x \neq x$ for all $0 \neq x \in \mathbb{Z}^{n}$.
2. For each positive $x \in \mathbb{Z}^{n}$, there exists an integer $r \geq 1$ such that $c^{r}(x)$ is not positive.

Proof. First notice that since $Q$ is Dynkin ( $\Longrightarrow$ acyclic), both the Coxeter functors $\mathcal{C}^{+}, \mathcal{C}^{-}$and the Coxeter transformation $c$ exist.

1. Suppose that $x \in \mathbb{Z}^{n}$ such that $c x=x$. Because $\sigma_{i_{2}}, \ldots, \sigma_{i_{n}}$ do not change the $i_{1}$-th coordinate, we must have $\sigma_{i}(x)_{i}=x_{i}$ and thus $\sigma_{i}(x)=x$. Doing this inductively, yields that $\sigma_{i_{j}}(x)=x$ for all $1 \leq j \leq n$. From the definition of $\sigma_{i}$, we get that $\left(x, e_{j}\right)_{Q}=0$ for all $1 \leq j \leq n$. This implies that $q_{Q}(x)=0$ and thus $x=0$.
2. We know from Corollary 3.31 that $W(Q)$ is finite. Therefore, there exists some $h \geq 1$ such that $c^{h}=\mathrm{id}$ (e.g. $h=|W(Q)|$ is the cardinality of $W(Q)$ ). If $x, c x, c^{2} x, \ldots, c^{h-1} x$ are all positive, then $y=x+c x+c^{2} x+\ldots+c^{h-1} x$ is positive $(\Longrightarrow y \neq 0)$ and $c y=y$ which contradicts part 1. But then one of the $c x, c^{2} x, \ldots, c^{h-1} x$ is not positive.
Proposition 3.69. Let $Q$ be a Dynkin quiver, and let $\vee$ be an indecomposable representation of $Q$. Then there is a finite $r$ such that $\underbrace{\mathcal{C}^{+} \ldots \mathcal{C}^{+}}_{r \text { times }} \boxtimes=0$.
Proof. Let $x=\underline{\operatorname{dim} \mathbb{V}}$. If $\mathbb{V}=\mathbb{S}_{i}$ (for some $i \in Q_{0}$ ), then we have $\mathcal{C}^{+} \mathbb{V}=0$ since $\mathcal{R}_{i}^{+} \mathbb{S}_{i}=0$. Thus let $\mathbb{V} \neq \mathbb{S}_{i}$ for all $i \in Q_{0}$. Then Lemma 3.68 yields an integer $r$ such that the vector $c^{r}(x)=c^{r}(\underline{\operatorname{dim} \mathbb{V}})=\underline{\operatorname{dim}}\left(\mathcal{C}^{+}\right)^{r} \boxtimes$ is not positive. Therefore, we have $\left(\mathcal{C}^{+}\right)^{r} \boxtimes=\underbrace{\mathcal{C}^{+} \ldots \mathcal{C}^{+}}_{r \text { times }} \mathbb{V}=0$.
Example 3.70. In Section 4.2, Proposition 3.69 is shown in the example of the linear $L_{n}$-quiver. Then, we can take $r=n$.

### 3.5 The Proof of Gabriel's Theorem

We first state the version of Gabriel's Theorem that we prove. This version of the theorem is superior to Theorem 3.2 since it gives a characterization of the isomorphism classes of the indecomposable representations for the Dynkin quivers. This section is based on [2].
Theorem 3.71 (Gabriel, version 2). Let $Q$ be a connected quiver and let $k$ be an arbitrary field. Then:

1. $Q$ is of finite-type if and only if $Q$ is a Dynkin quiver.
2. When the equivalent conditions in (1) are satisfied, the correspondence $\mathbb{V} \mapsto \underline{\operatorname{dim} \mathbb{V}}$ induces a bijection between the set of isomorphism classes of indecomposable representations of $Q$ and the set $\Phi^{+}(Q)=\Phi(Q) \cap \mathbb{N}^{n}$ of positive roots of $Q$.

Proof. Step 1: First, we assume that $Q$ is a Dynkin quiver. We show that $Q$ is of finite-type and that (2) holds. For simplicity, we set $Q_{0}=\{1,2 \ldots, n\}$ and we assume that $1,2, \ldots, n$ is a (+)admissible ordering. Let $\mathbb{V}$ be an indecomposable representation of $Q$. Then $\operatorname{dim} \mathbb{V}$ is positive. By Lemma 3.68, we know that there exists $r$ such that $c^{r}(\underline{\operatorname{dim} V})$ is non-positive which implies that $\left(\mathcal{C}^{+}\right)^{r} \boxtimes=0$ by Proposition 3.69. Here, $c=\sigma_{n} \ldots \sigma_{2} \sigma_{1}$ denotes the Coxeter transformation. We consider the sequence

$$
\left(i_{1}, i_{2}, \ldots, i_{n}, i_{n+}, \ldots, i_{r n}\right)=(1,2, \ldots, n, 1,2, \ldots, n, 1,2, \ldots, n)
$$

where the sequence $1,2, \ldots, n$ is repeated $r$ times. Let $0 \leq t<r n$ such that $\sigma_{i_{t}} \ldots \sigma_{i_{1}}(\underline{\operatorname{dim}} \mathbb{V})$ is positive but $\sigma_{i_{t+1}} \ldots \sigma_{i_{1}}$ (dim $\mathbb{V}$ ) is non-positive and thus 0 . This is equivalent to saying that $\mathcal{R}_{i_{t+1}}^{+} \ldots \mathcal{R}_{i_{1}}^{+}$is the smallest subsequence of $\left(\mathcal{C}^{+}\right)^{r}$, which sends $\mathbb{V}$ to the trivial representation. Since $\mathcal{R}_{i_{t}}^{+} \ldots \mathcal{R}_{i_{1}}^{+} \mathrm{V}$ is indecomposable and $\mathcal{R}_{i_{t+1}}^{+} \ldots \mathcal{R}_{i_{1}}^{+} \mathrm{V}=0$, we know from Corollary 3.50 that $\mathcal{R}_{i_{t}}^{+} \ldots \mathcal{R}_{i_{1}}^{+} \mathbb{V}=\mathbb{S}_{i_{t+1}}$, where $\mathbb{S}_{i_{t+1}}$ is the simple representation (at the vertex $i_{t+1}$ ) of $s_{i_{t}} \ldots s_{i_{2}} s_{i_{1}} Q$. Applying the same corollary inductively yields that

$$
\mathbb{V} \cong \mathcal{R}_{i_{1}}^{-} \mathcal{R}_{i_{2}}^{-} \ldots \mathcal{R}_{i_{t}}^{-} \mathcal{R}_{i_{t}}^{+} \ldots \mathcal{R}_{i_{1}}^{+\mathbb{V}}=\mathcal{R}_{i_{1}}^{-} \mathcal{R}_{i_{2}}^{-} \ldots \mathcal{R}_{i_{t}}^{-} \mathbb{S}_{i_{t+1}}
$$

Note that the sequence $\mathcal{R}_{i_{1}}^{-} \mathcal{R}_{i_{2}}^{-} \ldots \mathcal{R}_{i_{t}}^{-} \mathcal{R}_{i_{t}}^{+} \ldots \mathcal{R}_{i_{1}}^{+}$is well-defined since for all $1 \leq j \leq t$, we have that $i_{j}$ is a sink in $s_{i_{j-1}} \ldots s_{i_{1}} Q$ and $i_{j}$ is a source in $s_{i_{j+1}} \ldots s_{i_{t}} s_{i_{t}} \ldots s_{i_{1}} Q=s_{i_{j}} \ldots s_{i_{1}} Q$ because $1,2, \ldots, n$ is $(+)$-admissible. Moreover, Theorem 3.49 yields that

$$
\underline{\operatorname{dim} \mathbb{V}}=\underline{\operatorname{dim}} \mathcal{R}_{i_{1}}^{-} \mathcal{R}_{i_{2}}^{-} \ldots \mathcal{R}_{i_{t}}^{-} \mathbb{S}_{i_{t+1}}=\sigma_{i_{1}} \sigma_{i_{2}} \ldots \sigma_{i_{t}}\left(\underline{\operatorname{dim}} \mathbb{S}_{i_{t+1}}\right)=\sigma_{i_{1}} \sigma_{i_{2}} \ldots \sigma_{i_{t}}\left(e_{i_{t+1}}\right)
$$

Therefore, all dimension vectors of indecomposable representations are (positive) roots of $Q$. Moreover, let $\mathbb{V}, \mathbb{W}$ be two indecomposable representations such that $\operatorname{dim} \mathbb{V}=\underline{\operatorname{dim} W} \mathbb{W}$. We get that $\underline{\operatorname{dim}} \mathcal{R}_{i_{t}}^{+} \ldots \mathcal{R}_{i_{1}}^{+} \mathbb{W}=\sigma_{i_{t}} \ldots \sigma_{i_{1}}(\underline{\operatorname{dim} W})=\sigma_{i_{t}} \ldots \sigma_{i_{1}}(\underline{\operatorname{dim} \mathbb{V}})=e_{i_{t+1}}=\underline{\operatorname{dim}} \mathbb{S}_{i_{t+1}} . \overline{B u t} \mathbb{S}_{i_{t+1}}$ is the only representation with dimension vector $e_{i_{t+1}}$ and therefore $\mathcal{R}_{i_{t}}^{+} \ldots \mathcal{R}_{i_{1}}^{+} \mathbb{W}=\mathbb{S}_{i_{t+1}}$. Now, applying the mirrored sequence of reflection functors, we get

$$
\mathbb{W} \cong \mathcal{R}_{i_{1}}^{-} \ldots \mathcal{R}_{i_{t}}^{-} \mathcal{R}_{i_{t}}^{+} \ldots \mathcal{R}_{i_{1}}^{+} \mathbb{W} \cong \mathcal{R}_{i_{1}}^{-} \ldots \mathcal{R}_{i_{t}}^{-} \mathbb{S}_{i_{t+1}} \cong \mathbb{V}
$$

Therefore, we get a well-defined map

$$
\phi: \mathcal{I} \rightarrow \Phi^{+}(Q),[\mathbb{V}] \mapsto \underline{\operatorname{dim}} \mathbb{V}
$$

where $\mathcal{I}$ denotes the set of isomorphism classes of indecomposable representations $\mathbb{V}$ of $Q$ and [ $\mathbb{V}$ ] denotes the isomorphism class of $\mathbb{V}$. Note that we have shown that $\phi$ is an injective map. Since $\Phi^{+}(Q)$ is finite (Lemma 3.30) this shows that $Q$ is of finite-type. We now show that $\phi$ is surjective and thus proving (2). Let $x \in \Phi^{+}(Q)$ be a positive root of $q_{Q}$. From Lemma 3.68, we get an $m$ such that $c^{m}(x)$ is not positive. Similarly as before, let $0 \leq t<m n$ be minimal such that $\sigma_{i_{t}} \ldots \sigma_{i_{1}}(x)$ is positive, but $\sigma_{i_{t+1}} \sigma_{i_{t}} \ldots \sigma_{i_{1}}(x)$ is not positive. But since $\sigma_{i_{t+1}} \sigma_{i_{t}} \ldots \sigma_{i_{1}}(x)$ is a root, it is either negative or positive. However, $\sigma_{i_{t+1}}$ only changes the $i_{t+1^{-}}$th coordinate and thus we conclude that $\sigma_{i_{t}} \ldots \sigma_{i_{1}}(x)=e_{i_{t+1}}$ and therefore we have $x=\sigma_{i_{1}} \ldots \sigma_{i_{t}}\left(e_{i_{t+1}}\right)$. But the representation $\mathbb{V}=\mathcal{R}_{i_{1}}^{-} \ldots \mathcal{R}_{i_{t}}^{-} \mathbb{S}_{i_{t+1}}$ is indecomposable by Corollary 3.50 and has dimension vector $x$. This shows that $\phi$ is surjective.

Step 2: For any quiver $Q$ and a subquiver $Q^{\prime}$, we have that $\operatorname{Rep}_{\mathrm{k}}\left(Q^{\prime}\right)$ is a full subcategory of
 zero spaces and zero maps to vertices and arrows in $Q \backslash Q^{\prime}$. The subcategory is full because every morphism $\tilde{\phi}: \tilde{\mathbb{V}} \rightarrow \tilde{\mathbb{W}}$ induces a morphism $\phi: \mathbb{V} \rightarrow \mathbb{W}$. Indeed, if $i \in Q_{0} \backslash Q_{0}^{\prime}$, then $\tilde{\phi}_{i}=0: 0 \rightarrow 0$ and therefore $\phi=\left(\tilde{\phi}_{i}\right)_{i \in Q_{0}^{\prime}}$ is a morphism between representations of $Q^{\prime}$. If two representations $\tilde{\mathbb{V}}, \tilde{W}$ of $Q$ are isomorphic, then the induced representations $\mathbb{V}, \mathbb{W}$ of $Q^{\prime}$ are also isomorphic. Indeed, for an isomorphism $\tilde{\phi}: \tilde{\mathbb{V}} \rightarrow \tilde{W}$ we have that the induced morphism $\phi: \mathbb{V} \rightarrow \mathbb{W}$ is in $\operatorname{Rep}_{\mathrm{k}}\left(Q^{\prime}\right)$ (since the subcategory is full) and it is also an isomorphism. Therefore, if $Q$ is of finite-type, then $Q^{\prime} \subseteq Q$ is also of finite-type. But from Lemma 3.15, we know that each wild quiver $Q$ has a tame subquiver $Q^{\prime}$. Therefore, it suffices to show that all tame quivers are not of finite-type.

Step 3: We now go through the list of tame quivers resp. Euclidean graphs and show that they have infinitely many isomorphism classes of indecomposable representations.
$\tilde{A}_{0}$ : Let $Q$ be a quiver of type $\tilde{A}_{0}$. In Example 2.21 and Example 2.25 we have shown that (for $\bar{k}$ an algebraically closed field), the isomorphism classes of indecomposable representations are given by the Jordan normal blocks. For an arbitrary field $k$, we still get that the different Jordan blocks $J_{\lambda, n}$ for $\lambda \in k$ give pairwise non-isomorphic indecomposable representations. Indeed, from Example 2.21, we know that two representations $\mathbb{V}=k^{n} \xrightarrow{M} k^{n}, \mathbb{W}=k^{n} \xrightarrow{N} k^{n}$ are isomorphic if and only if there exists an invertible matrix $A$ such that $N A=A M$. By a generalized version of the Jordan normal form for an arbitrary field $k$, this implies that $M$ and $N$ have the same Jordan normal form. If $k$ is an infinite field, then for each $n \in \mathbb{N}$, we get an infinite number of isomorphism classes of indecomposable representations. If $k$ is finite, we can vary over $n \in \mathbb{N}$ and still get an infinite number of isomorphism classes.
$\tilde{A_{n}}, n \geq 1$ : Let $Q$ be a quiver of type $\tilde{A_{n}}$. For an arbitrary $\lambda \in k$, we define a representation $\mathbb{V}_{\lambda}$ by $\overline{\text { setting } V_{i}}=k$ for all vertices and by setting the maps $v_{a}=$ id for all $a \in Q_{1}$, except one map $v_{b}=\lambda$ id. $\mathbb{V}_{\lambda}$ is given by the following picture.


Since all the vector spaces are 1-dimensional and the quiver is connected, $\mathbb{V}_{\lambda}$ is indecomposable. Indeed if $\mathbb{V}_{\lambda}=\mathbb{V}_{1} \oplus \mathbb{V}_{2}$ and $\left(\mathbb{V}_{1}\right)_{i}=k$ for some vertex $i$, then we have a vertex $j$ such that the map between $i$ and $j$ is the identity, which implies that $\left(\mathbb{V}_{1}\right)_{j}=k$. Now, using that $Q$ is connected, we conclude that $\mathbb{V}=\mathbb{V}_{1}$. Moreover, for a morphism $\phi: \mathbb{V}_{\lambda_{1}} \rightarrow \mathbb{V}_{\lambda_{2}}$ we have that $\phi_{i}=\phi_{j}$ for all vertices $i, j$. This can be shown using a similar argument as for the part that $\mathbb{V}_{\lambda}$ is indecomposable. But then we have the following commutative diagram.


Therefore, if $\lambda_{1} \neq \lambda_{2}$ we have that $\phi_{0}=0$ which shows that $\mathbb{V}_{\lambda_{1}} \cong \mathbb{V}_{\lambda_{2}}$ if and only if $\lambda_{1}=\lambda_{2}$. More generally, for any $n \in \mathbb{N}$, we can define a representation by putting $k^{n}$ at each vertex and by letting all maps be the identity, apart from one map which is given by a Jordan block $J_{\lambda, n}$. Similarly, as in the first case $(n=1)$, it can be shown that this defines an indecomposable representation since each Jordan block is indecomposable. The same argument as above also shows that different choices of $\lambda \in k$ give pairwise non-isomorphic indecomposable representations. Therefore, we have an infinite number of indecomposable representations.
$\underline{\tilde{D}_{n}, n \geq 4}$ : We first consider the following quiver of type $\tilde{D}_{4}$ :


For an arbitrary $\lambda \in k$ we consider the following representation, called $\mathbb{V}_{\lambda}$ :


We claim that $\mathbb{V}_{\lambda}$ is indecomposable and that different choices of $\lambda \in k$ give pairwise non-isomorphic representations. Indeed, let $\mathbb{W}$ be a summand of $\mathbb{V}_{\lambda}$ and notice that $\mathbb{W}$ needs to have a non-trivial vector space at one of the outer vertices $0,1,3$ or 4 (else $\mathbb{W}$ is trivial). But this implies that the vector spaces at all the other outer vertices must be equal to $k$ (and thus $\mathbb{W}=\mathbb{V}$ ). To see this, consider the images of the four maps $\binom{1}{0},\binom{1}{1},\binom{0}{1},\binom{1}{\lambda}$. Notice that for any choice of three of the four maps, their images generate $k^{2}$. Also, for any partition into two pairs, the images of one pair (of maps) always generate $k^{2}$. Therefore, if we split the four copies of the vector spaces $k$ (at the outer vertices) into two different summands of $\mathbb{V}$ (either split it 2 to 2 or 3 to 1 ), we have that one summand must have $k^{2}$ at the middle vertex. But then the other summand has 0 at the middle vertex, which is impossible because none of the maps is trivial. This shows that $\mathbb{V}_{\lambda}$ is indecomposable. Moreover, if there is an isomorphism $\phi: \mathbb{V}_{\lambda_{1}} \rightarrow \mathbb{V}_{\lambda_{2}}$, we have

$$
\begin{aligned}
& \binom{1}{0} \phi_{2}=\left.\phi_{2}\right|_{k \oplus 0}=\phi_{1}\binom{1}{0}=\binom{\alpha}{0}, \\
& \binom{0}{1} \phi_{2}=\left.\phi_{2}\right|_{0 \oplus k}=\phi_{0}\binom{0}{1}=\binom{0}{\beta} .
\end{aligned}
$$

This implies that $\phi_{2}=\left(\begin{array}{cc}\alpha & 0 \\ 0 & \beta\end{array}\right)$. But since $\binom{1}{1} \phi_{3}=\phi_{2}\binom{1}{1}=\binom{\alpha}{\beta}$ we have $\phi_{2}=\left(\begin{array}{cc}\alpha & 0 \\ 0 & \alpha\end{array}\right)$. This yields $\binom{1}{\lambda_{1}} \phi_{4}=\phi_{2}\binom{1}{\lambda_{2}}=\binom{\alpha}{\alpha \lambda_{2}}$, which implies $\lambda_{1}=\lambda_{2}$. More generally, for any $n \in \mathbb{N}$, one can show that the following representations are indecomposable and are pairwise non-isomorphic for different choices of $\lambda \in k$ (here $J_{\lambda, n}$ denotes the $n \times n$-Jordan block).


Therefore, for any field $k$, the above quiver has infinitely many isomorphism classes of indecomposable representations. We can extend this argument to the following quiver $Q_{s p}$ of type $\tilde{D}_{n}$.


Indeed, for $n \in \mathbb{N}$, one can show that the following representations are indecomposable and pairwise non-isomorphic for different $\lambda \in k$.


Now we use Corollary 3.52 to conclude that any quiver of type $\tilde{D_{n}}$ has infinitely many isomorphism classes of indecomposable representations. Using the corollary, it remains to show that any quiver of type $\tilde{D}_{n}$ results from the quiver $Q_{s p}$ by applying a sequence of $\mathcal{R}_{i}^{+}$respectively $\mathcal{R}_{i}^{-}$(meaning that we apply $s_{i}$ to the quiver itself (but only if $i$ is a sink or a source) and we apply $\mathcal{R}_{i}^{ \pm}$to representations). Indeed, one can proceed as follows:

1. First, we fix the direction of the arrow between 2 and 3 (if it does not have the right direction) by applying $\mathcal{R}_{2}^{-} \mathcal{R}_{1}^{-} \mathcal{R}_{0}^{-}$.
2. Then we fix the arrows $0 \rightarrow 2$ and $1 \rightarrow 2$ by possibly applying $\mathcal{R}_{0}^{-}$resp. $\mathcal{R}_{1}^{-}$.
3. Now we sequentially fix the arrows $i \rightarrow i+1$ for $3 \leq i \leq n-3$ (i.e. the horizontal arrows). First notice that the sequence $\mathcal{R}_{i+1}^{+} \ldots \mathcal{R}_{n-3}^{+} \mathcal{R}_{n}^{+} \mathcal{R}_{n-1}^{+} \mathcal{R}_{n-2}^{+}$is well-defined and it only changes the direction of the arrow $i \rightarrow i+1$. Thus applying this sequence over and over (for increasing $i$, we can fix the direction of the horizontal arrows.
4. Finally, we can fix the direction of the arrows $n-1 \rightarrow n-2$ and $n \rightarrow n-2$ by possibly applying $\mathcal{R}_{n-1}^{-}$resp. $\mathcal{R}_{n}^{-}$.

To illustrate this algorithm, consider the following quiver $Q$ :


We have that $s_{n} s_{n-1} s_{0} s_{2} s_{1} s_{0} Q_{s p}=Q$ and therefore an indecomposable representation $\mathbb{V} \not \approx \mathbb{S}_{i}$ of $Q_{s p}$ yields an indecomposable representation $\mathcal{R}_{n}^{-} \mathcal{R}_{n-1}^{-} \mathcal{R}_{0}^{-} \mathcal{R}_{2}^{-} \mathcal{R}_{1}^{-} \mathcal{R}_{0}^{-} \vee$ of $Q$. In general, by Corollary 3.52, there is a one-to-one correspondence between the isomorphism classes of indecomposable representations of $Q_{s p}$ and any other quiver $Q$ of type $\tilde{D_{n}}$ which shows that all quivers of type $\tilde{D}_{n}$ are not of finite-type.
$\underline{\tilde{E}_{6}}$ : We first reduce the problem to the following quiver $Q_{s p}$ of type $\tilde{E}_{6}$.


Similarly as in the $\tilde{D}_{n}$-case, for any quiver $Q$ of type $\tilde{E}_{6}$, we can find a sequence of reflection functors such that applying this sequence to $Q_{s p}$ yields the quiver $Q$. Thus, using Corollary 3.52 , it is enough to prove that $Q_{s p}$ is not of finite-type. Now, instead of just giving a family of indecomposable representations, we show how we find such indecomposable representations. This is useful since this idea can then be extended to the cases $\tilde{E}_{7}$ and $\tilde{E}_{8}$. First, we note that by deleting the vertex 0 , we get a Dynkin quiver $Q^{\prime}$ (see proof of Theorem 3.16). We also notice that the unique smallest positive radical vector $\delta=(1,1,2,2,3,2,1) \in \mathbb{N}^{7}$ yields a vector $\delta^{\prime}=(1,2,2,3,2,1) \in \mathbb{N}^{6}$ (take $\delta_{i}=\delta_{i}^{\prime}$ for all $\left.i \in Q_{0}^{\prime}\right)$ such that $q_{Q^{\prime}}\left(\delta^{\prime}\right)=1$ and thus $\delta^{\prime}$ is a root of $Q^{\prime}$. From part (1) we know that there exists an indecomposable representation $\mathbb{V}$ of $Q^{\prime}$ that has dimension vector $\delta^{\prime}$. Now, for any $\lambda \in k$, we get an indecomposable representation $\mathbb{V}_{\lambda}$ of $Q$ by setting $V_{0}=k$ and the map between $V_{0}$ and $V_{2}=k^{2}$ to be $\binom{1}{\lambda}$.


Since $\mathbb{V}$ is an indecomposable representation of $Q^{\prime}$, we conclude that $\mathbb{V}_{\lambda}$ is also indecomposable. Moreover, if $\phi: \mathbb{V}_{\lambda_{1}} \cong \mathbb{V}_{\lambda_{2}}$, we have that (w.l.o.g) $\phi_{i}=\operatorname{id}_{V_{i}}$ for $i \in\{1,2,3,4,5,6\}$. Thus, we get

$$
\binom{1}{\lambda_{2}} \phi_{0}=\phi_{2}\binom{1}{\lambda_{1}}=\operatorname{id}_{V_{2}}\binom{1}{\lambda_{1}}=\binom{1}{\lambda_{1}}
$$

which implies $\lambda_{1}=\lambda_{2}$. Therefore, $\mathbb{V}_{\lambda}$ are indecomposable, pairwise non-isomorphic (for different choices of $\lambda \in k$ ) representations of $Q$. More generally, for any $n \in \mathbb{N}$, we get that the following representations are indecomposable and pairwise non-isomorphic.


The tensor notation $v_{\alpha} \otimes \mathbb{1}_{k^{n}}$ comes from the fact that for all vertices $i \in\{1,2,3,4,5,6\}$ we have that the vector space at vertex $i$ is isomorphic to $V_{i} \otimes k^{n}$ and for all arrows $\alpha \in\{a, b, c, d, e\}$ we have that the map is given by $v_{\alpha} \otimes \mathbb{1}_{k^{n}}$. The maps are explicitly given by the block matrices

$$
v_{\alpha} \otimes \mathbb{1}_{k^{n}}=\left(\begin{array}{cccc}
v_{\alpha} & 0 & \ldots & 0 \\
0 & v_{\alpha} & \ldots & 0 \\
\vdots & \vdots & \ddots & 0 \\
0 & 0 & \ldots & v_{\alpha}
\end{array}\right) \text { for } \alpha \in\{a, b, c, d, e\} .
$$

$\tilde{E}_{7}$ and $\tilde{E}_{8}$ : The cases for quivers of type $\tilde{E}_{7}$ and $\tilde{E}_{8}$ are analogous to the case of $\tilde{E}_{6}$-quivers. Indeed, by deleting the vertex 0 , we get the roots $\delta^{\prime}=(2,2,3,4,3,2,1)$ resp. $\delta^{\prime}=(2,3,4,6,5,4,3,2)$ of a quiver of type $E_{7}$ resp. $E_{8}$. Then, we can use the same arguments as above.

Remark 3.72. In the proof of Gabriel's Theorem (step 1), we only needed that for a Dynkin quiver $Q$ and any indecomposable representation $\mathbb{V}$, there exists a sequence of indices $i_{1}, \ldots, i_{s}$ (possibly with repetitions) such that $\mathcal{R}_{i_{s}}^{+} \ldots \mathcal{R}_{i_{1}}^{+}$is well-defined and $\mathcal{R}_{i_{s}}^{+} \ldots \mathcal{R}_{i_{1}}^{+} \mathbb{V}=0$. This explains why one does not need to consider the Weyl group $W(Q)$ to prove the 'if part' of Gabriel's Theorem in the special case of the $A_{n}$-type quivers. Indeed, we mainly need the Weyl group to prove the existence of a finite number $r$ such that $\underbrace{\mathcal{C}^{+} \ldots \mathcal{C}^{+}}_{r \text { times }} \mathbb{V}=0$ for any indecomposable representation $\mathbb{V}$ of a Dynkin quiver $Q$. In Chapter 4, we explicitly show how we can construct such a sequence of indices for $A_{n}$-type quivers.

## Chapter 4

## The Special Case: $A_{n}$-type Quivers

In this chapter, we focus on the special case of $A_{n}$-type quivers $(n \geq 1)$, a subset of the Dynkin quivers. These quivers are crucial in topological data analysis as we will see in Chapter 5 . The theory in Chapters 2 and 3 is more explicit for $A_{n}$-type quivers. In this chapter we build indecomposable representations for $A_{n}$-type quivers, explore Coxeter functors on them, and establish the 'if part' of Gabriel's Theorem, demonstrating the finiteness of $A_{n}$-type quivers. This chapter is based on Appendix A of the book 'Persistence Theory: From Quiver Representations to Data Analysis' written by Steve Y. Oudot [1].

### 4.1 Interval Representations

We denote the set $\{b, b+1, \ldots, d\}$ by $[b, d]$. Recall that a quiver $Q$ is of type $A_{n}$ if its underlying graph $\bar{Q}$ is of the following form.


An important example is the linear quiver $L_{n}$.


A morphism $\phi=\left(\phi_{i}\right)_{1 \leq i \leq n}$ between two quiver representations $\mathbb{V}$ and $\mathbb{W}$ of an $A_{n}$-type quiver $Q$ is given by the following commutative diagram.


Lemma 4.1. For all $A_{n}$-type quivers $Q$, the set of positive roots of $q_{Q}$ is given by the vectors of the form $x=(0, \ldots, 0,1, \ldots, 1,0, \ldots, 0)^{T}$ with the first and last 1 's occurring at positions $b \leq d$ in the range $[1, n]$.

Proof. First, we recall that the Tits form $q_{Q}$ is independent of the orientations of the arrows of $Q$, which allows us to simultaneously find the roots of all $A_{n}$-type quivers (for a fixed $n$ ). From Example 3.11, we know the Tits form for $A_{n}$-type quivers $Q$ :

$$
q_{Q}(x)=\sum_{i \in Q_{0}} x_{i}^{2}-\sum_{a \in Q_{1}} x_{t_{a}} x_{h_{a}}=\sum_{i=1}^{n} x_{i}^{2}-\sum_{i=1}^{n-1} x_{i} x_{i+1}=\frac{1}{2}\left[x_{1}^{2}+x_{n}^{2}+\sum_{i=1}^{n-1}\left(x_{i}-x_{i+1}\right)^{2}\right]
$$

Since $q_{Q}$ is positive definite ( $Q$ is Dynkin), we know from Corollary 3.31, that the roots of $q_{Q}$ are exactly given by the set $\left\{x \in \mathbb{Z}^{n} \mid q_{Q}(x)=1\right\}$. Thus for any positive root $x \in \mathbb{Z}^{n}$, we have $x_{i} \in\{0,1\}$ for all coordinates/vertices $i$ and that there can be at most two pairs of consecutive vertices/coordinates $\left(x_{i}, x_{i+1}\right)$ with difference 1 . This shows that the roots are of the desired form $x=(0, \ldots, 0,1, \ldots, 1,0, \ldots, 0)^{T}$.

We know from Gabriel's Theorem that for every $A_{n}$-type quiver $Q$, we have a bijection between the isomorphism classes of indecomposable representations and the positive roots of $q_{Q}$ :

$$
\phi: \mathcal{I} \rightarrow \Phi^{+}(Q),[\mathbb{V}] \mapsto \underline{\operatorname{dim} V} .
$$

Therefore, we use Lemma 4.1 to identify the indecomposable representations of the $A_{n}$-type quivers.
Proposition 4.2. For all $A_{n}$-type quivers $Q$, each indecomposable representation is isomorphic to one of the representations in Figure 4.1 for some interval $[b, d] \subseteq[1, n]$.


Figure 4.1: The interval representation $\rrbracket_{Q}[b, d]$.

Definition 4.3. The representations in Figure 4.1 are called interval representations and are denoted by $\square_{Q}[b, d]$.

Proof of Proposition 4.2. Using Gabriel's Theorem, let $\mathbb{V}$ be the indecomposable representation of $Q$ that belongs to the positive root $x=(0, \ldots, 0,1, \ldots, 1,0, \ldots, 0)^{T}$ with the 1 's occurring in the interval $[b, d]$ (see Lemma 4.1). The representation $\mathbb{V}$ has the vector space $k$ at every vertex $i \in[b, d]$. The maps from and to 0 vector spaces are 0 .

Moreover, the maps between successive copies of $k$ need to be isomorphisms, because otherwise $\mathbb{V}$ could be decomposed further. Indeed, if $k \stackrel{v_{a}}{ } k$ is not an isomorphism, then $v_{a}=0$ and we have the decomposition $\mathbb{V}=\mathbb{V}_{1} \oplus V_{2}$, where the representations $\mathbb{V}_{1}$ and $\mathbb{V}_{2}$ are as follows.

$$
\begin{aligned}
& \mathbb{V}_{1}: \quad \cdots-0--V_{b}^{\prime}=k \xrightarrow{v_{b}} k----V_{a}^{\prime}=k-0=V_{a+1}^{\prime}--0--0=V_{d}^{\prime}--0 .-\cdots
\end{aligned}
$$

In addition, $\mathbb{V}$ is isomorphic to the interval representation $\square_{Q}[b, d]$. Indeed, for any isomorphism $k \xrightarrow{v_{a}} k$, consider the following commutative diagram.

Notice that this allows us to successively define the maps $\phi_{b}, \phi_{b+1}, \ldots, \phi_{d}$ so that the resulting morphism $\phi=\left(\phi_{i}\right)_{i=1}^{n}$ is an isomorphism between $\mathbb{V}$ and $\mathbb{\square}_{Q}[b, d]$.

Now, we can use the Krull-Remak-Schmidt Theorem to decompose any representation of an $A_{n}$-type quiver.

Definition 4.4. For an $A_{n}$-type quiver $Q$ and a finite-dimensional representation $\mathbb{V}$, the following decomposition is called the interval decomposition of $\mathbb{V}$ :

$$
\mathbb{V} \cong \bigoplus_{j=1}^{r} \rrbracket_{Q}\left[b_{j}, d_{j}\right]
$$

We have now found all isomorphism classes of indecomposable representations of an $A_{n}$-type quiver $Q$. Since there are only finitely many possible choices for $b \leq d$ in $[1, n]$, we have also shown that $A_{n}$-type quivers are of finite-type. However, we have used Gabriel's Theorem, thus our argument above is a circular argument. Therefore, we want to prove the 'if part' of Gabriel's Theorem for $A_{n}$-type quivers, i.e. that quivers of type $A_{n}$ are of finite-type and that the map $\phi: \mathcal{I} \rightarrow \Phi^{+}(Q),[\mathbb{V}] \mapsto \underline{\operatorname{dim} V}$ is an isomorphism. To do this, we first observe the action of reflection functors and Coxeter functors on $A_{n}$-type quivers.

### 4.2 Coxeter Functors on $A_{n}$-Type Quivers

In Remark 3.72, we noticed the following: in the proof of Gabriel's Theorem we only needed that for a Dynkin quiver $Q$ and any indecomposable representation $\mathbb{V}$, there exists a sequence of indices $i_{1}, \ldots, i_{s}$ (possibly with repetitions) such that $\mathcal{R}_{i_{s}}^{+} \ldots \mathcal{R}_{i_{1}}^{+}$is well-defined and $\mathcal{R}_{i_{s}}^{+} \ldots \mathcal{R}_{i_{1}}^{+} \mathbb{V}=0$. For an $A_{n}$-type quiver $Q$, we now construct such a sequence explicitly. We first consider the linear quiver $L_{n}$.


Proposition 4.5. [1, Ex. A.19] For every finite-dimensional representation $\mathbb{V}$ of the linear quiver $L_{n}$, we have $\underbrace{\mathcal{C}^{+} \ldots \mathcal{C}^{+}}_{n \text { times }} \mathbb{V}=0$. Here, $\mathcal{C}^{+}$is the Coxeter functor belonging to the $(+)$-admissible ordering $n, n-1, \ldots, 1$.

Proof. Notice that the natural ordering of the vertices is a ( - -admissible ordering and therefore, the ordering $n, n-1, \ldots, 1$ is $(+)$-admissible. Now, let $\mathbb{V}$ be an indecomposable representation of $L_{n}$ and let $x=\left(x_{1}, x_{2}, \ldots, x_{n}\right)^{T}=\underline{\operatorname{dim}} \mathbb{V}$. We now apply the Coxeter functor $\mathcal{C}^{+}$and by Theorem 3.49, we get:

$$
\begin{aligned}
\underline{\operatorname{dim}} \mathcal{R}_{n}^{+} \mathbb{V} & =0 \text { or } \sigma_{n}(x)=\left(x_{1}, x_{2}, \ldots, x_{n-1}, x_{n-1}-x_{n}\right)^{T}, \\
\underline{\operatorname{dim}} \mathcal{R}_{n-1}^{+} \mathcal{R}_{n}^{+} \mathbb{V} & =0 \text { or } \sigma_{n-1} \sigma_{n}(x)=\left(x_{1}, x_{2}, \ldots, x_{n-2}-x_{n}, x_{n-1}-x_{n}\right)^{T} .
\end{aligned}
$$

$$
\vdots
$$

$$
\begin{gathered}
\underline{\operatorname{dim}} \mathcal{R}_{2}^{+} \ldots \mathcal{R}_{n-1}^{+} \mathcal{R}_{n}^{+} \mathbb{V}=0 \text { or } \sigma_{2} \ldots \sigma_{n-1} \sigma_{n}(x)=\left(x_{1}, x_{1}-x_{n}, \ldots, x_{n-2}-x_{n}, x_{n-1}-x_{n}\right)^{T} . \\
\underline{\operatorname{dim}} \mathcal{R}_{1}^{+} \mathcal{R}_{2}^{+} \ldots \mathcal{R}_{n-1}^{+} \mathcal{R}_{n}^{+} \mathbb{V}=0 \text { or } \sigma_{1} \sigma_{2} \ldots \sigma_{n-1} \sigma_{n}(x)=\left(-x_{n}, x_{1}-x_{n}, \ldots, x_{n-2}-x_{n}, x_{n-1}-x_{n}\right)^{T} .
\end{gathered}
$$

Since dimension vectors are non-negative, we conclude that either $\mathcal{C}^{+} \boxtimes=\mathcal{R}_{1}^{+} \mathcal{R}_{2}^{+} \ldots \mathcal{R}_{n-1}^{+} \mathcal{R}_{n}^{+} \boxtimes=0$ or that $x_{n}=0$. We can iterate this process successively and get:

$$
\begin{aligned}
\begin{aligned}
& \operatorname{dim} \mathcal{C}^{+} \\
& V \\
& \underline{\operatorname{dim} \mathcal{C}^{+}} \mathcal{C}^{+} \\
& V \text { or }\left(0, x_{1}, x_{2}, \ldots, x_{n-2}, x_{n-1}\right)^{T}, \\
& \vdots \\
& \underline{\operatorname{dim}}\left(0,0, x_{1}, \ldots, x_{n-3}, x_{n-2}\right)^{T}, \\
& \underbrace{\mathcal{C}^{+} \ldots \mathcal{C}^{+}}_{n-1 \text { times }} \mathbb{V}=0 \text { or }\left(0,0,0, \ldots, 0, x_{1}\right)^{T}, \\
& \underline{\operatorname{dim}} \underbrace{\mathcal{C}^{+} \ldots \mathcal{C}^{+}}_{n \text { times }} \mathbb{V}=0 .
\end{aligned} .
\end{aligned}
$$

We now generalize the statement of Proposition 4.5 to be true for all $A_{n}$-type quivers.
Proposition 4.6. For an $A_{n}$-type quiver $Q$ and a finite-dimensional representation $\mathbb{V}$, there exists a sequence of reflection functors such that

$$
\underbrace{\mathcal{C}^{+} \ldots \mathcal{C}^{+}}_{n \text { times }} \mathcal{R}_{1}^{+} \ldots \mathcal{R}_{i_{r}}^{+} \mathcal{R}_{1}^{+} \ldots \mathcal{R}_{i_{r-1}}^{+} \ldots \mathcal{R}_{1}^{+} \ldots \mathcal{R}_{i_{1}}^{+} \mathbb{V}=0
$$

Proof. If $Q$ is not the linear quiver $L_{n}$, let $i_{1}<i_{2}<\ldots<i_{r}$ be the heads of the backward arrows.


If we apply the sequence $s_{1} \ldots s_{i_{1}}$ on $Q$, we get the same quiver, but the arrow between $i_{1}$ and $i_{1}+1$ has changed direction.


Also, notice that the sequence of reflection functors $\mathcal{R}_{1}^{+} \ldots \mathcal{R}_{i_{1}}^{+}$is well-defined. We can repeat this process for $i_{2}, \ldots, i_{r}$ and get that $s_{1} \ldots s_{i_{r}} s_{1} \ldots s_{i_{r-1}} \ldots s_{1} \ldots s_{i_{1}} Q$ is the linear quiver $L_{n}$. In addition, from Corollary 3.50, we get that for any indecomposable representation $\mathbb{V}$ of $Q$, the representation $\mathcal{R}_{1}^{+} \ldots \mathcal{R}_{i_{r}}^{+} \mathcal{R}_{1}^{+} \ldots \mathcal{R}_{i_{r-1}}^{+} \ldots \mathcal{R}_{1}^{+} \ldots \mathcal{R}_{i_{1}}^{+} \mathbb{V}$ is either 0 or an indecomposable representation of $L_{n}$. Using Proposition 4.5, we conclude that

$$
\underbrace{\mathcal{C}^{+} \ldots \mathcal{C}^{+}}_{n \text { times }} \mathcal{R}_{1}^{+} \ldots \mathcal{R}_{i_{r}}^{+} \mathcal{R}_{1}^{+} \ldots \mathcal{R}_{i_{r-1}}^{+} \ldots \mathcal{R}_{1}^{+} \ldots \mathcal{R}_{i_{1}}^{+} \mathbb{V}=0 .
$$

### 4.3 Proof of Gabriel's Theorem for $A_{n}$-Type Quivers

We are now ready to prove the following theorem.
Theorem 4.7 (Gabriel's Theorem for $A_{n}$-type quivers). Let $Q$ be a quiver of type $A_{n}$ and let $k$ be an arbitrary field. Then:

1. $Q$ is of finite-type.
2. The correspondence $\mathbb{V} \mapsto \underline{\text { dim } \sqrt{ } \text { induces a bijection between the set of isomorphism classes of }}$ indecomposable representations of $Q$ and the set $\Phi^{+}(Q)=\Phi(Q) \cap \mathbb{N}^{n}$ of positive roots of $Q$.

Proof. The proof is analogous to Step 1 of the proof of Theorem 3.71. Indeed, we have shown that there exists a minimal sequence of reflection functors $\mathcal{R}_{i_{t+1}}^{+} \ldots \mathcal{R}_{i_{t}}^{+}$, such that $\mathcal{R}_{i_{t+1}}^{+} \ldots \mathcal{R}_{i_{t}}^{+} \mathbb{V}=0$ and $\mathcal{R}_{i_{+}}^{+} \ldots \mathcal{R}_{i_{t}}^{+} \mathbb{V} \neq 0$. The same argument as in the proof of Theorem 3.71 now yields a welldefined and injective map $\phi: \mathcal{I} \rightarrow \Phi^{+}(Q),[\mathbb{V}] \mapsto \underline{\operatorname{dim} \mathbb{V}}$. This shows that $Q$ is of finite-type since we have shown that the set of positive roots $\Phi^{+}(Q)$ is finite for an $A_{n}$-type quiver. Notice that we do not need the result for a general Dynkin quiver (Lemma 3.30). In addition, proving the surjectivity of $\phi$ can be done explicitly in the case of $A_{n}$-type quivers. Indeed, we know that every positive root $x \in \Phi^{+}(Q)$ is of the form $x=(0, \ldots, 0,1, \ldots, 1,0, \ldots, 0)^{T}$ and we have that the interval representation $\rrbracket_{Q}[b, d]$ has this dimension vector for a suitable choice of $b \leq d$. It remains to show that the interval representations are indeed indecomposable. This can be either seen by direct observation or we can show that the endomorphism ring $\operatorname{End}\left(\square_{Q}[b, d]\right)$ is local and then apply Proposition 2.38. Now, let $\psi \in \operatorname{End}\left(\square_{Q}[b, d]\right)$ and observe the following commutative diagram.


By commutativity, we have $\psi_{b}=\psi_{b+1}=\ldots=\psi_{d}=\alpha \in k$. Therefore, we can embed $\operatorname{End}\left(\square_{Q}[b, d]\right)$ in the field $k$. However, since $\operatorname{End}\left(\square_{Q}[b, d]\right)$ is a $k$-vector space, it follows that $\operatorname{End}\left(\square_{Q}[b, d]\right)=k$ which is local because fields are local. This shows that $\phi: \mathcal{I} \rightarrow \Phi^{+}(Q)$ is an isomorphism.

Remark 4.8. We highlight the results from Chapter 3 that are needed in the proof above, and which results are not needed. We do this by sections:

[^0]3.2: We have used the central property of simple reflections: $\underline{\operatorname{dim}} \mathcal{R}_{i}^{ \pm} \mathbb{V}=\sigma_{i}(\underline{\operatorname{dim} \mathbb{V}})$ (for indecomposable representations $\mathbb{V} \neq \mathbb{S}_{i}$ ). However, we did not use the Weyl group and it was sufficient to consider the set $\left\{x \in \mathbb{Z}^{n} \mid q_{Q}(x)=1\right\}$ as our root system. Indeed, we were able to describe the roots explicitly, which showed that $\left\{x \in \mathbb{Z}^{n} \mid q_{Q}(x)=1\right\}$ is finite.
3.3: The proof relied heavily on reflection functors. Therefore, the whole Section 3.3 is needed in the proof above.
3.4: We have used Coxeter functors. However, we did not use its connection to the Coxeter transformation and the Weyl group.

### 4.4 The Diamond Principle

Let $Q$ be an $A_{n}$-type quiver and let $i \in Q_{0}$ be a sink. Then, any representation $\mathbb{V}$ and its reflection $\mathbb{W}=\mathcal{R}_{i}^{+} \boxtimes$ can be expressed by the following diagram, because $\mathcal{R}_{i}^{+}$only changes the vector space at the vertex $i$ and its incident arrows resp. the maps corresponding to those arrows.


Lemma 4.9. We have the following transformation rules for the interval representations:

$$
\begin{aligned}
\mathcal{R}_{i}^{+} \square_{Q}[i, i] & =0, \\
\mathcal{R}_{i}^{+} \rrbracket_{Q}[i, d] & =\square_{Q}[i+1, d] \text { if } i<d, \\
\mathcal{R}_{i}^{+} \square_{Q}[b, i] & =\square_{Q}[b, i-1] \text { if } b<i, \\
\mathcal{R}_{i}^{+} \rrbracket_{Q}[i+1, d] & =\square_{Q}[i, d] \text { if } i+1 \leq d, \\
\mathcal{R}_{i}^{+} \square_{Q}[b, i-1] & =\square_{Q}[b, i] \text { if } b \leq i-1, \\
\mathcal{R}_{i}^{+} \rrbracket_{Q}[b, d] & =\square_{Q}[b, d] \text { otherwise. }
\end{aligned}
$$

The following picture visualizes this result. The black interval represents the interval representation $\square_{Q}[b, d]$ and the red dot shows the position of the vertex $i$.


Proof. The result follows directly from the definition of the reflection functor $\mathcal{R}_{i}^{+}$. We have already seen the first rule and we prove the second and fifth rule, thus proving one case each, where the interval gets shorter resp. longer. The other cases are analogous. Using that $\mathcal{R}_{i}^{+} \square_{Q}[i, d]$ is indecomposable (Corollary 3.50) and thus isomorphic to an interval representation, it suffices to prove that $\left(\mathcal{R}_{i}^{+} \square_{Q}[i, d]\right)_{i}=0$. Indeed, we have:

$$
\left(\mathcal{R}_{i}^{+} \square_{Q}[i, d]\right)_{i}=\operatorname{ker} \xi_{i}=\operatorname{ker}(0 \oplus k \rightarrow k)=\operatorname{ker}(k \xrightarrow{\text { id }} k)=0 .
$$

For the fifth rule, it suffices to prove that $\left(\mathcal{R}_{i}^{+} \square_{Q}[b, i-1]\right)_{i}=k$. Indeed, we have:

$$
\left(\mathcal{R}_{i}^{+} \square_{Q}[b, i-1]\right)_{i}=\operatorname{ker} \xi_{i}=\operatorname{ker}(k \oplus 0 \rightarrow 0)=\operatorname{ker}(k \xrightarrow{0} k)=k .
$$

Now, we understand how reflection functors operate on interval decomposition. Using this, we consider the above diagram and we construct a general setting, which behaves similarly to the transformation rules in Lemma 4.9.

Definition 4.10. Let $Q$ be an $A_{n}$-type quiver and let $i$ be a sink. Given two finite-dimensional representations $\mathbb{V} \in \operatorname{rep}_{\mathrm{k}}(Q)$ and $\mathbb{W} \in \operatorname{rep}_{\mathrm{k}}\left(s_{i} Q\right)$, that differ only at the spaces $V_{i}, W_{i}$, and their incident maps, we get the following diagram, where the central rhombus is called a diamond:


The diamond is called exact if $\operatorname{im} f=\operatorname{ker} g$ in the following sequence

$$
W_{i} \xrightarrow{f} V_{i-1} \oplus V_{i+1} \xrightarrow{g} V_{i},
$$

where $f: x \mapsto\left(w_{a}(x), w_{b}(x)\right)$ and $g:(x, y) \mapsto v_{c}(x)+v_{d}(y)$.
Theorem 4.11 (Diamond principle). Let $Q$ be an $A_{n}$-type quiver and let $i$ be a sink. Let $\mathbb{V} \in \operatorname{rep}_{\mathrm{k}}(Q)$ and $\mathbb{W} \in \operatorname{rep}_{\mathrm{k}}\left(s_{i} Q\right)$ be two finite-dimensional representations, that differ only at the spaces $V_{i}, W_{i}$, and their incident maps. If we suppose that the diamond is exact, then the interval decompositions of $\mathbb{V}$ and $\mathbb{W}$ are related to each other through the following matching rules:

- summands $\square_{Q}[i, i]$ and $\square_{s_{i} Q}[i, i]$ are unmatched,
- summands $\square_{Q}[b, i]$ are matched with summands $\square_{s_{i} Q}[b, i-1]$, and $\square_{Q}[b, i-1]$ with $\square_{s_{i} Q}[b, i]$,
- summands $\square_{Q}[i, d]$ are matched with summands $\square_{s_{i} Q}[i+1, d]$, and $\square_{Q}[i+1, d]$ with $\square_{s_{i} Q}[i, d]$,
- every other summand $\square_{Q}[b, d]$ is matched with the summand $\square_{s_{i} Q}[b, d]$.


Figure 4.2: The matching rules from the diamond principle. The top row illustrates the second, third, and fourth matching rules. The bottom row shows the two unmatched cases.

The transformation result for the interval representations (Lemma 4.9) and the diamond principle look quite similar. Indeed, we prove that the results naturally imply each other.

Proof of the equivalence of Lemma 4.9 and the diamond principle. Step 1: We show Lemma 4.9 assuming the diamond principle. First, we prove the exactness of the diamond for the representations $\mathbb{V}=\rrbracket_{Q}[b, d]$ and $\mathbb{W}=\mathcal{R}_{i}^{+} \rrbracket_{Q}[b, d]$.


Notice that $g=\xi_{i}$. Thus, we get $\operatorname{im} f=W_{i}=\operatorname{ker} \xi_{i}=\operatorname{ker} g$ and thus the diamond is exact. Now observe that the first matching in the diamond principle implies the first transformation in the lemma. The second matching implies the second and fourth transformations. The third matching implies the third and fifth transformations and the last matching implies the last transformation.

Step 2: We now assume Lemma 4.9. Let $\mathbb{V}, \mathbb{W}$ be as in the diamond principle and such that the diamond is exact. We claim that $\mathbb{W} \cong \mathbb{U} \oplus \mathbb{K}$, where $\mathbb{U}=\mathcal{R}_{i}^{+} \mathbb{V}$ and $\mathbb{K}=\bigoplus_{j=1}^{r} \square_{s_{i} Q}[i, i]=\bigoplus_{j=1}^{r} \mathbb{S}_{i}$ is the representation of $s_{i} Q$ made from $r=\operatorname{dim} \operatorname{ker} f$ copies of the representation $\square_{s_{i} Q}[i, i]=\mathbb{S}_{i}$. Indeed, from the definition of $\mathcal{R}_{i}^{+} \vee$ and the exactness of the diamond, we have $U_{j}=V_{j}$ for all $j \neq i$ and $U_{i}=\operatorname{ker} \xi_{i}=\operatorname{ker} g=\operatorname{im} f$. Now let $C=\operatorname{ker} f^{\perp} \subseteq W_{i}$. From the first isomorphism theorem, we conclude that the map

$$
\left.f\right|_{C}: C \rightarrow \operatorname{im} f=U_{i}
$$

is an isomorphism. We can also take an arbitrary isomorphism $h$ : $\operatorname{ker} f \rightarrow K_{i}$ (here, we need $r=\operatorname{dim} \operatorname{ker} f$ ) and we can define an isomorphism of representations $\phi: \mathbb{W} \rightarrow \mathbb{U} \oplus \mathbb{K}$ as follows

$$
\phi_{j}= \begin{cases}\left.f\right|_{C} \oplus h=\left(\begin{array}{cc}
\left.f\right|_{C} & 0 \\
0 & h
\end{array}\right) & \text { if } j=i, \\
\mathbb{1}_{V_{j}} & \text { if } j \neq i\end{cases}
$$

One can explicitly check that this is an isomorphism of representations. Now, using the transformation rules in Lemma 4.9, we get the matching rules in the diamond principle.

Using the equivalence of Lemma 4.9 and Theorem 4.11, one may attempt to proof the 'if part' of Gabriel's Theorem for $A_{n}$-type quivers using the diamond principle. It turns out that this is possible [1, Section 4.4.2.]. It also turns out that the diamond principle is helpful because it is easier to check if im $f=\operatorname{ker} g$ rather than compute reflections of representations and check if they are isomorphic. We use the diamond principle in the setting of zigzag persistence in Section 5.3.

## Chapter 5

## Applications of Quiver Representation Theory to Persistent Homology

In this chapter, we explain the connection between topological data analysis and quiver representation theory. We first introduce simplicial complexes and homology groups. Using these concepts we establish a connection between the persistence of topological spaces and quiver representations. In Section 5.2 we focus on the topological aspect of this connection. Section 5.3 is devoted to the algebraic aspects of this connection. In particular, we show a connection between zigzag persistent homology and the diamond principle from Chapter 4.

### 5.1 Simplicial Complexes and Homology

In this section, we introduce simplicial complexes which are topological spaces that can be built from simple building blocks called simplices. We then introduce simplicial homology. This section roughly follows the outline of Chapters III. 1 and IV. 1 of the book 'Computational Topology: An Introduction' by Herbert Edelsbrunner and John Harer [7] and Chapter 1 of the paper 'Barcodes: The persistent topology of data' by Robert Ghrist [8].

Definition 5.1. Let $A=\left\{v_{0}, v_{1}, \ldots, v_{n}\right\} \subseteq \mathbb{R}^{m}$. We say that $A$ is in general position if $A$ is not contained in an affine hyperplane of dimension less than $n$. This is equivalent to the condition that the vectors $v_{1}-v_{0}, \ldots, v_{n}-v_{0}$ are linearly independent resp. to the condition that the subspace generated by $A$ has dimension at least $n$.

Definition 5.2. Let $A=\left\{v_{0}, v_{1}, \ldots, v_{n}\right\} \subseteq \mathbb{R}^{m}$ be points in general position. The $n$-simplex $\sigma=\left[v_{0}, v_{1}, \ldots, v_{n}\right]$ is defined to be the convex hull of $A$, i.e. the smallest subset of $\mathbb{R}^{m}$ that contains all the points $v_{0}, v_{1}, \ldots, v_{n}$. We can also represent the $n$-simplex using linear combinations:

$$
\left[v_{0}, \ldots, v_{n}\right]=\left\{\sum_{i=0}^{n} \lambda_{i} v_{i} \mid \sum_{i=0}^{n} \lambda_{i}=1, \forall 0 \leq j \leq n: \lambda_{j} \geq 0\right\}
$$

We say that the points $v_{i}$ span $\sigma$. The dimension of $\sigma$ is $\operatorname{dim} \sigma=n$.

Definition 5.3. For $A$ in general position and the associated $n$-simplex $\sigma$, the points $v_{i} \in A$ are called vertices and for each subset $B \subseteq A$, the convex hull of $B$, denoted by $\tau$, is called a face of $\sigma$. Notice that $B$ is also in general position and therefore, $\tau$ is a simplex itself. We often write $\tau \leq \sigma$. A face is called proper if $B \subsetneq A$ and is denoted by $\tau<\sigma$. Moreover, the boundary of $\sigma$ is denoted by $\partial \sigma$ and it is the union of all proper faces of $\sigma$. The interior is given by $\sigma-\partial \sigma$.

We observe that a point $x \in \sigma$ lies in the interior of $\sigma$ if and only if all its coefficients $\lambda_{i}$ are positive. Indeed, if some coefficient $\lambda_{i}=0$, then we have that $x$ lies in the face $\left[v_{0}, \ldots, \widehat{v_{i}}, \ldots, v_{n}\right]$,
-


Figure 5.1: The simplices of dimensions $0,1,2$, and 3 have special names and are called vertex, edge, triangle, and tetrahedron (from left to right).
where the hat indicates that the vertex $v_{i}$ is omitted in the list of vertices. We conclude that every point $x \in \sigma$ lies in the interior of exactly one of the faces of $\sigma$. The vertices corresponding to this face are exactly the vertices $v_{j}$, for which $\lambda_{j}>0$.

### 5.1.1 Geometric and Abstract Simplicial Complexes

We combine simplices to create a broader variety of topological spaces, called simplicial complexes. We do this in a way such that it is closed under taking faces and such that it has no improper intersections.

Definition 5.4. A (geometric) simplicial complex $K$ is a finite union of simplices such that for each simplex $\sigma \in K$ and any face $\tau \leq \sigma$, we have $\tau \in K$. Moreover, for two simplices $\sigma_{1}, \sigma_{2} \in K$, we have that $\sigma_{1} \cap \sigma_{2}$ is either empty of a face of both simplices $\sigma_{1}$ and $\sigma_{2}$. A subcomplex of $K$ is a simplicial complex $L \subseteq K$. The dimension of $K$ is given by the maximum dimension of its simplices.

Definition 5.5. For each $j \leq \operatorname{dim} K$, we define a particular subcomplex $K^{(j)}$, called the $\boldsymbol{j}$-skeleton which consists of all simplices of dimension $j$ or less. The 0-skeleton is usually called the vertex set $K^{(0)}=\operatorname{Vert} K$.

A simplicial complex $K$ is just a set of simplices and thus lacks a topology. We fix this by considering the topological space given by the union of the simplices together with the subspace topology inherited from the ambient Euclidean space $\mathbb{R}^{m}$ in which the simplices live. We denote this topological space by $|K|$. Because every point of a simplex belongs to the interior of one of its faces, we conclude that every point $x \in|K|$ belongs to the interior of exactly one simplex in $K$. This allows us to give the following description of $|K|$.

Definition 5.6. Let $K$ be a simplicial complex with vertices $v_{0}, v_{1}, \ldots, v_{n}$. Every point $x \in|K|$ belongs to the interior of exactly one simplex $\sigma=\left[v_{0}, v_{1}, \ldots, v_{k}\right] \in K$. Therefore, we have $x=\sum_{i=0}^{k} \lambda_{i} v_{i}$ with $\sum_{i=0}^{k} \lambda_{i}=1$ and $\lambda_{i}>0$ for all $i$. Setting $b_{i}(x)=\lambda_{i}$ for $0 \leq i \leq k$ and $b_{i}(x)=0$ for $k+1 \leq i \leq n$, we have $x=\sum_{i=0}^{n} b_{i}(x) v_{i}$ and we call the $b_{i}(x)$ the barycentric coordinates of $x$ in $K$.

Remark 5.7. If it is clear from the context that we consider a topological space, we often write $K$ instead of $|K|$.

It is often easier to describe or construct a simplicial complex abstractly instead of giving a full geometric description of its simplices.

Definition 5.8. An abstract simplicial complex $A$ is a finite collection of sets such that for all $\alpha \in A$ and all $\beta \subseteq \alpha$, we have $\beta \in A$. The sets in $A$ are the simplices and the dimension of $a$ simplex is $\operatorname{dim} \alpha=|\alpha|-1$. The dimension of the complex $A$ is given by the maximum dimension of its simplices. A face of $\alpha \in A$ is a non-empty subset $\beta \subseteq \alpha$ and it is proper if $\beta \subsetneq \alpha$. $A$ subcomplex is an abstract simplicial complex $B \subseteq A$.

We notice that this condition looks quite similar to the condition that for each simplex $\sigma \in K$ and each face $\tau \leq \sigma$, we have that $\tau \in K$. However, there is no explicit condition on the intersection of simplices. We now describe how one can switch from a geometric simplicial complex to an abstract simplicial complex (and vice versa).


Figure 5.2: Simplicial complex with vertices $A, B, C, D, E, F, G$ and simplices $[A, B],[A, C]$, $[A, D],[B, C],[B, D],[B, E],[C, D],[C, E],[F, G],[A, B, C],[A, B, D],[A, C, D],[B, C, D],[B, C, E]$ and $[A, B, C, D]$.

Definition 5.9. Let $K$ be a geometric simplicial complex. For each simplex $\sigma \in K$ we take its vertex set $\alpha$. Those sets form an abstract simplicial complex $A$. We call $A$ a vertex scheme of $K$. Analogously, we call $K$ a geometric realization of $A$.

Notice that for $\tau \leq \sigma$, we have an inclusion of vertex sets $\alpha_{\tau} \subseteq \alpha_{\sigma}$ and thus $\alpha_{\tau} \in A$ because $\tau \in K$. This proves that a vertex scheme is an abstract simplicial complex. We can also construct a geometric realization for an abstract simplicial complex if the dimension of the ambient space is high enough.
Theorem 5.10 (Geometric Realization Theorem). Every abstract simplicial complex of dimension d has a geometric realization in $\mathbb{R}^{2 d+1}$ [7, p. 64].

Since simplicial complexes are topological spaces, we can consider continuous maps between them. Favourable maps between simplicial complexes are the maps that respect the underlying simplicial complex structure.

Definition 5.11. A vertex map is a function $\phi$ : Vert $K \rightarrow \operatorname{Vert} L$ such that it sends the vertices of every simplex in $K$ to vertices of a simplex in $L$.

Definition 5.12. Every vertex map $\phi: \operatorname{Vert} K \rightarrow \operatorname{Vert} L$ can be extended to a continuous map $f:|K| \rightarrow|L|$ which is defined by:

$$
f(x)=\sum_{i=0}^{n} b_{i}(x) \phi\left(v_{i}\right) .
$$

This map is called the simplicial map induced by $\phi$. Note that $f$ is linear one each simplex in $K$. Therefore, we say that $f$ is piecewise linear. To abbreviate notation, we often write $f: K \rightarrow L$.

Since linear maps are continuous and $f: K \rightarrow L$ is piecewise linear, it is indeed continuous. Per definition, simplicial maps map simplices to simplices. For a simplex $\sigma \in K$, we now consider the dimension of $f(\sigma)$ and observe how the faces of $\sigma$ and $f(\sigma)$ are related.
Lemma 5.13. Let $\sigma$ be a p-simplex in $K$. Then $f(\sigma)$ is a simplex in $L$ that has dimension less or equal than $p$. Moreover, we have

- If $f(\sigma)$ has dimension $p$, then the $(p-1)$-dimensional faces of $\sigma$ map to the corresponding $(p-1)$-dimensional faces of $f(\sigma)$.
- If $f(\sigma)$ has dimension $p-1$, then exactly two of the $(p-1)$-dimensional faces of $\sigma$ map to $f(\sigma)$ and all the other $(p-1)$-dimensional faces of $\sigma$ map to faces of dimension $p-2$.
- If $f(\sigma)$ has dimension less than $p-1$, then the images of all faces of $\sigma$ also have dimension less than $p-1$.

Proof. Let $\sigma=\left[v_{0}, v_{1}, \ldots, v_{p}\right]$. We know that $f(\sigma)=\left[\phi\left(v_{0}\right), \phi\left(v_{1}\right), \ldots, \phi\left(v_{p}\right)\right]$ is a simplex in $L$, which has dimension less or equal than $p$ since it is spanned by $p+1$ vertices.

- If $f(\sigma)$ has dimension $p$, then the vertices $\phi\left(v_{0}\right), \phi\left(v_{1}\right), \ldots, \phi\left(v_{p}\right)$ are pairwise different and the ( $p-1$ )-dimensional faces of $\sigma$ map to the corresponding $(p-1)$-dimensional faces of $f(\sigma)$ :

$$
f\left(\left[v_{0}, v_{1}, \ldots, \widehat{v_{j}}, \ldots, v_{p}\right]=\left[\phi\left(v_{0}\right), \phi\left(v_{1}\right), \ldots, \widehat{\phi\left(v_{j}\right)}, \ldots, \phi\left(v_{p}\right)\right] .\right.
$$

- If $f(\sigma)$ has dimension $p-1$, then exactly two of the vertices in $\phi\left(v_{0}\right), \phi\left(v_{1}\right), \ldots, \phi\left(v_{p}\right)$ coincide. Let $\phi\left(v_{i}\right)=\phi\left(v_{j}\right)$. Then $f\left(\left[v_{0}, v_{1}, \ldots, \hat{v}_{i}, \ldots, v_{p}\right]\right)=f\left(\left[v_{0}, v_{1}, \ldots, \hat{v_{j}}, \ldots, v_{p}\right]\right)$ and the images of all other $(p-1)$-dimensional faces have dimension $p-2$ since they contain the coinciding vertices $\phi\left(v_{i}\right)$ and $\phi\left(v_{j}\right)$.
- If $f(\sigma)$ has dimension less than $p-1$, then the images of all faces of $\sigma$ also have dimension less than $p-1$.


### 5.1.2 Homology

Homology groups are the mathematical description of holes and higher dimensional analogues in a topological space. We introduce homology for simplicial complexes (called simplicial homology with $\mathbb{Z} / 2 \mathbb{Z}$-coefficients) and develop the necessary theory that is used in persistent homology. Using the simplices of a simplicial complex, we create abelian groups. We then combine them, creating a so-called chain complex, which forms the foundation for homology.

Definition 5.14. Let $K$ be a simplicial complex and let $p \in \mathbb{Z}$. A p-chain $c$ is a formal finite sum of $p$-simplices in $K$. We can write

$$
c=\sum a_{i} \sigma_{i},
$$

where $a_{i} \in \mathbb{Z} / 2 \mathbb{Z}$ are the coefficients and $\sigma_{i}$ are the $p$-simplices. Note that since the coefficients are either 0 or 1, we can assume w.l.o.g. that a p-chain is of the form $c=\sigma_{1}+\ldots+\sigma_{n}$. We can add two $p$-chains together componentwise: if $c=\sum a_{i} \sigma_{i}$ and $c^{\prime}=\sum b_{i} \sigma_{i}$, then $c+c^{\prime}=\sum\left(a_{i}+b_{i}\right) \sigma_{i}$, where the coefficients are added modulo 2, i.e. $1+1=0$. The $p$-chains together with this addition form an abelian group, called the group of p-chains, denoted by $C_{p}=C_{p}(K)$. We have such a group for all $p \in \mathbb{Z}$, however, for $p$ negative or bigger than $\operatorname{dim} K$, we have that $C_{p}(K)=0$ is trivial.

Definition 5.15. The boundary of a p-simplex is defined to be the sum of its ( $p-1$ )-dimensional faces. If we write $\sigma=\left[v_{0}, v_{1}, \ldots, v_{n}\right]$ for the $p$-simplex which is spanned by the listed vertices, then the boundary is given by

$$
\partial_{p} \sigma=\sum_{j=0}^{p}\left[v_{0}, v_{1}, \ldots, \widehat{v_{j}}, \ldots, v_{n}\right] .
$$

We extend this definition by linearity to define the boundary of a p-chain $c=\sum a_{i} \sigma_{i}$ to be $\partial_{p} c=\sum a_{i} \partial_{p} \sigma_{i}$. Therefore we get a map $\partial_{p}: C_{p}(K) \rightarrow C_{p-1}(K)$, which is called the boundary $\boldsymbol{m a p}$ for chains. Notice that the boundary map commutes with addition, i.e. $\partial_{p}\left(c+c^{\prime}\right)=\partial_{p} c+\partial_{p} c^{\prime}$, therefore $\partial_{p}$ is a homomorphism. The chain complex is the infinite sequence of chain groups connected by their boundary maps:

$$
\ldots \xrightarrow{\partial_{p+2}} C_{p+1}(K) \xrightarrow{\partial_{p+1}} C_{p}(K) \xrightarrow{\partial_{p}} C_{p-1}(K) \xrightarrow{\partial_{p-1}} \ldots
$$

We often drop the index of the boundary map because it is implied by the dimension of the chain to which it is applied.

We now consider two special types of chains called cycles and boundaries. We observe how they are related to each other.
Definition 5.16. A p-cycle is a p-chain with empty boundary, i.e. $\partial c=0$. Because $\partial$ commutes with addition, we have that the $p$-cycles form an abelian subgroup of $C_{p}(K)$, denoted by $Z_{p}=Z_{p}(K)$. Notice that $Z_{p}=\operatorname{ker} \partial_{p}$.

For $p=0$, we have that each vertex is mapped to 0 and thus $Z_{0}=\operatorname{ker} \partial_{0}=C_{0}$. For higher $p$, it is generally not true that $Z_{p}=C_{p}$. However, one finds that a $p$-chain $c=\sigma_{1}+\ldots+\sigma_{n}$ is a $p$-cycle if and only if each ( $p-1$ )-dimensional face of some $\sigma_{i}$ occurs to an even number in the set of the $(p-1)$-dimensional faces of $\sigma_{1}, \ldots, \sigma_{n}$.

Definition 5.17. A p-boundary is a p-chain that is the boundary of a $(p+1)$-chain, i.e. $c=\partial d$ for some $d \in C_{p+1}(K)$. Because $\partial$ commutes with addition, we have that the group of $p$-boundaries forms an abelian subgroup of $C_{p}(K)$ and it is denoted by $B_{p}=B_{p}(K)$. We have that $B_{p}=\operatorname{im} \partial_{p+1}$.

We now look at the following lemma that connects cycles and boundaries. This lemma is the fundamental property that makes homology work.

Lemma 5.18 (Fundamental Lemma of Homology). For each integer p and every $(p+1)$-chain $d$, we have $\partial_{p} \partial_{p+1} d=0$.

Proof. Because the boundary map is a homomorphism, it is enough to show $\partial_{p} \partial_{p+1} \tau=0$ for every $(p+1)$-simplex $\tau=\left[v_{0}, \ldots, v_{p+1}\right]$ :

$$
\begin{aligned}
\partial_{p} \partial_{p+1} \tau & =\partial_{p}\left(\sum_{j=0}^{p+1}\left[v_{0}, \ldots, \widehat{v_{j}}, \ldots, v_{p}\right]\right)=\sum_{\substack{i, j=0 \\
i \neq j}}^{p+1}\left[v_{0}, \ldots, \widehat{v_{i}}, \ldots, \widehat{v_{j}}, \ldots, v_{p}\right] \\
& =\sum_{\substack{i, j=0 \\
i<j}}^{p+1}\left[v_{0}, \ldots, \widehat{v_{i}}, \ldots, \widehat{v_{j}} \ldots v_{p}\right]+\sum_{\substack{i, j=0 \\
j<i}}^{p+1}\left[v_{0}, \ldots, \widehat{v_{j}}, \ldots, \widehat{v_{i}}, \ldots, v_{p}\right]=0 .
\end{aligned}
$$

Therefore, every $p$-boundary is also a $p$-cycle and we have that $B_{p}$ is a subgroup of $Z_{p}$. This allows us to take the quotients $Z_{p} / B_{p}$. Note that two elements $c+B_{p}, c^{\prime}+B_{p}$ are the same in $Z_{p} / B_{p}$ if and only if their difference $c-c^{\prime}$ is a boundary. We now use this connection between cycles and boundaries to define homology groups.

Definition 5.19. The p-th homology group of a simplicial complex $K$ is the quotient group, $H_{p}(K)=H_{p}=Z_{p} / B_{p}$. The p-th Betti number is the rank of this group, $\beta_{p}=\operatorname{rank} H_{p}$. A coset $[c]=c+B_{p}$ of $H_{p}$ is called a homology class. Two cycles that are in the same homology class are called homologous, denoted by $c \sim c^{\prime}$.

Note that $Z_{n}$ is an abelian group and thus the subgroup $B_{n}$ is a normal subgroup. Therefore, $H_{p}=Z_{p} / B_{p}$ is actually a group and it is abelian since $Z_{n}$ is. Note that $H_{p}$ is actually a $\mathbb{Z} / 2 \mathbb{Z}$-vector space.

Lemma 5.20. For every simplicial complex $K, H_{p}(K)$ is a finite-dimensional $\mathbb{Z} / 2 \mathbb{Z}$-vector space for each integer $p$.
Proof. Note that $C_{p}$ is isomorphic to the free abelian group over all the $p$-simplices $\sigma_{i}$ in $K$ :

$$
C_{p}=\bigoplus \mathbb{Z} / 2 \mathbb{Z} \cdot \sigma_{i}
$$

Because $K$ consists of only finitely many simplices, $C_{p}$ is finite-dimensional. Therefore the homomorphism $\partial_{p}$ is a a $\mathbb{Z} / 2 \mathbb{Z}$-linear map between the two finite-dimensional $\mathbb{Z} / 2 \mathbb{Z}$-vector spaces $C_{p}$ and $C_{p-1}$. This shows that $Z_{p}=\operatorname{ker} \partial_{p}$ and $B_{p}=i m \partial_{p+1}$ are $\mathbb{Z} / 2 \mathbb{Z}$-vector spaces and thus so is the quotient $H_{p}=Z_{p} / B_{p}$.

### 5.1.3 Maps in Homology and Homotopy Invariance

We have now defined the homology groups $H_{p}(K)$ for a simplicial complex $K$. Next, we observe how simplicial maps behave when passing to homology.

Definition 5.21. Let $K$ and $L$ be two simplicial complexes. We know that a simplicial map $f: K \rightarrow L$ takes each simplex of $K$ linearly to a simplex in $L$. For every integer $p$, it induces a map from the $p$-chains of $K$ to the $p$-chains in $L$. For a $p$-chain $c=\sum a_{i} \sigma_{i}$, we have that $f_{\#}(c)=\sum a_{i} \tau_{i}$, where $\tau_{i}=f\left(\sigma_{i}\right)$ if it has dimension $p$ and $\tau_{i}=0$ if $f\left(\sigma_{i}\right)$ has dimension less than $p$.

Lemma 5.22. For two simplicial complexes $K$ and $L, f: K \rightarrow L$ a simplicial map and $\partial_{K}$ and $\partial_{L}$ the boundary maps of the two chain complexes of $K$ and $L$, we have

$$
f_{\#} \circ \partial_{K}=\partial_{L} \circ f_{\#} .
$$

Proof. We use Lemma 5.13. If $f\left(\sigma_{i}\right)$ has dimension $p$, then all $(p-1)$-dimensional faces of $\sigma_{i}$ map to the corresponding $(p-1)$-dimensional faces of $\tau_{i}=f\left(\sigma_{i}\right)$. By linearity, we conclude that $f_{\#} \circ \partial_{K}\left(\sigma_{i}\right)=\partial_{L} \circ f_{\#}\left(\sigma_{i}\right)$. If $f\left(\sigma_{i}\right)$ has dimension less than $p$, then the ( $p-1$ )-dimensional faces of $\sigma_{i}$ map to simplices of dimension less than $(p-1)$, with the possible exception of exactly two ( $p-1$ )-dimensional faces that map to the same simplex and therefore their images cancel. Therefore, we have $f_{\#}\left(\partial_{K} \sigma_{i}\right)=\partial_{L} f_{\#}\left(\sigma_{i}\right)=0$. We conclude using the linearity of $f_{\#}, \partial_{K}$ and $\partial_{L}$.

Lemma 5.23. We have $f_{\#}\left(Z_{p}(K)\right) \subseteq Z_{p}(L)$ and $f_{\#}\left(B_{p}(K)\right) \subseteq B_{p}(L)$ for all simplicial maps $f: K \rightarrow L$. This induces a map in homology, denoted by $f_{*}: H_{p}(K) \rightarrow H_{p}(L)$.

Proof. For a $p$-cycle $c$ in $K$, the following map is well-defined:

$$
f_{*}\left([c]_{K}\right)=f_{*}\left(c+B_{p}(K)\right)=f_{\#}(c)+B_{p}(L)=\left[f_{\#}(c)\right]_{L} .
$$

This defines a well-defined map because $f_{\#}(c)$ is a $p$-cycle in $L$ and if $c$ and $c^{\prime}$ differ by a boundary in $K$, then their difference is mapped to a boundary in $L$.

Remark 5.24. The fact that simplicial maps induce maps in homology is often referred to as 'homology is functorial (for every p)'. This makes sense since the assignment

$$
K \mapsto H_{p}(K),(f: K \rightarrow L) \mapsto\left(f_{*}: H_{p}(K) \rightarrow H_{p}(L)\right.
$$

is a covariant functor from the category of simplicial complexes (with simplicial maps) to the category of $\mathbb{Z} / 2 \mathbb{Z}$-vector spaces (with $\mathbb{Z} / 2 \mathbb{Z}$-linear maps).

Considering induced maps in homology, one can ask, when (and if) such maps are isomorphisms. The following result answers a part of this question and explains, why homology is useful.

Theorem 5.25. [9, Corollary 2.11] If a simplicial map $f: K \rightarrow L$ is a homotopy equivalence, then the induced maps $f_{*}: H_{p}(K) \rightarrow H_{p}(L)$ are isomorphisms for all $p \geq 0$.

Therefore, if two simplicial complexes $K$ and $L$ are both triangulations of a space $X$, meaning that both $K$ and $L$ are homotopy equivalent to $X$, then we do not need to worry about which one is used since they yield isomorphic homology groups. This result is proven using a more general homology theory called singular homology. Singular homology can be applied to any topological space. In addition, singular homology can be defined with coefficients in any abelian group and not only for $\mathbb{Z} / 2 \mathbb{Z}$-coefficients. In the general case, homology groups are groups and not vector spaces. This explains why we talk about homology groups and not homology vector spaces (even though in our setting, they are always vector spaces). However, simplicial homology and singular homology are equivalent (on simplicial complexes) and therefore, we do not introduce singular homology here.

Example 5.26. We have claimed that homology groups measure the holes and higher-dimensional analogues of a topological space. To explain this claim, we consider the homology of the sphere $S^{n} \subseteq \mathbb{R}^{n+1}$. In particular, we look at $H_{n}\left(S^{n}\right)$. For $n=1, S^{1}$ is homotopy equivalent to the left simplicial complex in Figure 5.3. Its chain complex is ( $\langle a, b, c\rangle$ denotes $\mathbb{Z} / 2 \mathbb{Z} \cdot a \oplus \mathbb{Z} / 2 \mathbb{Z} \cdot b \oplus \mathbb{Z} / 2 \mathbb{Z} \cdot c$ ).

$$
\ldots \longrightarrow 0 \xrightarrow{\partial_{2}} C_{1}=\langle a, b, c\rangle \xrightarrow{\partial_{1}} C_{0}=\left\langle v_{0}, v_{1}, v_{2}\right\rangle \xrightarrow{\partial_{0}} 0 \longrightarrow \ldots
$$

We have $B_{1}=i m \partial_{2}=0, Z_{1}=\operatorname{ker} \partial_{1}=\langle a+b+c\rangle$. Therefore $H_{1}\left(S^{1}\right)=\mathbb{Z} / 2 \mathbb{Z}$ (using that homology is homotopy invariant). For $n=2, S^{2}$ is homotopy equivalent to the right simplicial complex in Figure 5.3. Notice that this simplicial complex is not a single 3-dimensional simplex, but four 2-dimensional simplices. We denote them by $S=\left[v_{0}, v_{1}, v_{2}\right], T=\left[v_{0}, v_{1}, v_{3}\right], U=\left[v_{0} \cdot v_{2}, v_{3}\right]$, $V=\left[v_{1}, v_{2}, v_{3}\right]$. Its chain complex is (all the other chain groups and boundary maps are zero).

$$
0 \xrightarrow{\partial_{3}} C_{2}=\langle S, T, U, V\rangle \xrightarrow{\partial_{2}} C_{1}=\langle a, b, c, d, e, f\rangle \xrightarrow{\partial_{1}} C_{0}=\left\langle v_{0}, v_{1}, v_{2}, v_{3}\right\rangle \xrightarrow{\partial_{0}} 0
$$

We have $B_{2}=i m \partial_{3}=0$ and to compute $Z_{2}=\operatorname{ker} \partial_{2}$, we look at its matrix (left) and compute its reduced row-echelon form (right) with $\mathbb{Z} / 2 \mathbb{Z}$-coefficients.

$$
\partial_{2}=\left(\begin{array}{llll}
1 & 1 & 0 & 0 \\
1 & 0 & 0 & 1 \\
1 & 0 & 1 & 0 \\
0 & 1 & 1 & 0 \\
0 & 1 & 0 & 1 \\
0 & 0 & 1 & 1
\end{array}\right) \cong\left(\begin{array}{llll}
1 & 1 & 0 & 0 \\
0 & 1 & 0 & 1 \\
0 & 0 & 1 & 1 \\
0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0 \\
0 & 0 & 0 & 0
\end{array}\right)
$$

We conclude that $Z_{2}=\operatorname{ker} \partial_{2}=\langle S+T+U+V\rangle$ and therefore $H_{2}\left(S^{2}\right)=\mathbb{Z} / 2 \mathbb{Z}$. It turns out that this is true for all $n \geq 1$, i.e. $H_{n}\left(S^{n}\right)=\mathbb{Z} / 2 \mathbb{Z}$. Therefore, the homology group $H_{n}$ detects the single $n$-dimensional hole inside $S^{n}$. We conclude that in general, the p-th Betti number $\beta_{p}=\operatorname{rank} H_{p}=\operatorname{dim} H_{p}$ counts the number of $p$-dimensional holes in our simplicial complex.


Figure 5.3: Two simplicial complexes that are homotopy equivalent to $S^{1}$ (left) resp. $S^{2}$ (right). The simplicial complex on the right has four 2-dimensional simplices and no 3-dimensional simplex.

### 5.2 Filtrations and Persistence

In this section, we use homology to finally show the connection between topological data analysis and quiver representation theory. First, we describe a method to construct a simplicial complex from a point cloud. We then apply homology to these simplicial complexes and obtain persistence modules, which are quiver representations of the linear $L_{n}$-quiver. We then analyze and visualize its interval decomposition. This section is based on Chapters III. 1 and VII. 1 of [7] and Chapter 2 of [8].

### 5.2.1 Point Cloud Triangulations

We describe a method to construct a simplicial complex from a point cloud.
Definition 5.27. Let $\left\{x_{i}\right\}=X$ be a finite set of points in $\mathbb{R}^{m}$ and let $\epsilon>0$. The Čech complex $\check{C}(X, \epsilon)$ is the abstract simplicial complex whose $n$-simplices are the sets $\left\{x_{i_{0}}, \ldots, x_{i_{n}}\right\} \subseteq X$ such that the closed balls of radius $\epsilon / 2$, centered at the $x_{i_{k}}$, have a non-trivial intersection, i.e.

$$
\bigcap_{j=0}^{n} \bar{B}_{\epsilon / 2}\left(x_{i_{j}}\right) \neq \emptyset .
$$

Note that the Čech complex is an abstract simplicial complex. However, since there always exists a geometric realization, we can consider the Čech complex to be a geometric simplicial complex.

We now have a method to construct simplicial complexes from point clouds. However, it requires a choice of some parameter $\epsilon$. For $\epsilon$ small enough, there are no non-trivial intersections, and thus $\check{C}(X, \epsilon)$ is a discrete set given by the vertices $\left\{x_{i}\right\}$. If $\epsilon$ is really big, all intersections are non-trivial, and therefore, $\check{C}(X, \epsilon)$ is a single high-dimensional simplex (and all of its faces). This leads to the question if there is a best-suited choice for $\epsilon$, such that $\check{C}(X, \epsilon)$ is a good representation of our point cloud. Consider the point cloud and the Cech complexes in Figure 5.5. This point


Figure 5.4: From a fixed set of points with balls of radius $\epsilon$ (left), one can construct the Čech complex $\check{C}(X, \epsilon)$, represented by its geometric realization (right).
cloud is a sample of points on an annulus. If $\epsilon$ is too small (left and middle), $\check{C}(X, \epsilon)$ does not yet represent the fact, that an annulus has a hole in the middle. For $\epsilon$ big enough (right), the hole suddenly appears. However, if $\epsilon$ is too big, then all balls intersect and the hole disappears. This example illustrates that there is generally no single preferred choice for the parameter $\epsilon$ and even if there was an optimal $\epsilon$, we would already need to have a rough understanding of our point cloud to figure out the optimal value for $\epsilon$. Instead of considering one value for $\epsilon$, we can choose two values for $\epsilon$ and observe how the associated Čech complexes differ. We do this by considering simplicial maps between Čech complexes.

Lemma 5.28. For $\epsilon_{1} \leq \epsilon_{2}$ there exists an injective simplicial map $\iota: \check{C}\left(X, \epsilon_{1}\right) \hookrightarrow \check{C}\left(X, \epsilon_{2}\right)$. For any simplex $\sigma \in \check{C}\left(X, \epsilon_{1}\right)$, its image $\iota(\sigma)$ is $\sigma$ itself, but now viewed in $\check{C}\left(X, \epsilon_{2}\right)$.

Proof. For a simplicial complex $L$ and a subcomplex $K$, the inclusion map Vert $K \hookrightarrow \operatorname{Vert} L$ is a vertex map and thus, there exists a simplicial map $\iota: K \rightarrow L$. Moreover, for any simplex $\sigma=$ $\left[v_{0}, v_{1}, \ldots, v_{p}\right]$, we have $\iota(\sigma)=\left[v_{0}, v_{1}, \ldots, v_{p}\right]$ and therefore this map is injective since $\iota$ is piecewise linear. Now we apply this to the simplicial complex $\check{C}\left(X, \epsilon_{2}\right)$ with subcomplex $\check{C}\left(X, \epsilon_{1}\right)$.


Figure 5.5: A sequence of Čech complexes for three different increasing choices of radius $\epsilon$, coming from a point cloud that is a sample of an annulus. The 2-dimensional simplices are yellow, the 3 -dimensional simplices are green and the 4 -dimensional simplices are violet.

### 5.2.2 Persistence Vector Spaces and Persistence Barcodes

Using the inclusion maps from Lemma 5.28, we can construct a sequence of Čech complexes (for a sequence of radii $\epsilon_{1}<\epsilon_{2}<\ldots<\epsilon_{n}$ ).

$$
\check{C}\left(X, \epsilon_{1}\right) \longleftrightarrow \check{C}\left(X, \epsilon_{2}\right) \longleftrightarrow \check{C}\left(X, \epsilon_{3}\right) \longleftrightarrow \ldots \longleftrightarrow \check{C}\left(X, \epsilon_{n}\right)
$$

Definition 5.29. A sequence of simplicial complexes $K_{1} \hookrightarrow K_{2} \hookrightarrow \ldots \hookrightarrow K_{n}$ is called a filtration of simplicial complexes.

Notice that the sequence above is a filtration. Now, we apply the homology functor(s) $H_{*}$. This means that we do not just consider the homology of the C Cech complexes $\check{C}\left(X, \epsilon_{i}\right)$, but for all the inclusion maps $\check{C}\left(X, \epsilon_{i}\right) \hookrightarrow \check{C}\left(X, \epsilon_{i+1}\right)$, we consider the corresponding map in homology $H_{*}\left(\check{C}\left(X, \epsilon_{i}\right)\right) \rightarrow H_{*}\left(\check{C}\left(X, \epsilon_{i+1}\right)\right)$. We know that all homology groups $H_{*}\left(\check{C}\left(X, \epsilon_{i}\right)\right)$ are finitedimensional vector spaces. Moreover, all maps between homology groups are linear maps. This yields the following sequence in homology

$$
H_{*}\left(\check{C}\left(X, \epsilon_{1}\right)\right) \longrightarrow H_{*}\left(\check{C}\left(X, \epsilon_{2}\right)\right) \longrightarrow H_{*}\left(\check{C}\left(X, \epsilon_{3}\right)\right) \longrightarrow H_{*}\left(\check{C}\left(X, \epsilon_{n}\right)\right) .
$$

As a sequence of vector spaces, this is a representation of the $L_{n}$-quiver!
Definition 5.30. The representation $H_{*}\left(\check{C}\left(X, \epsilon_{1}\right)\right) \rightarrow H_{*}\left(\check{C}\left(X, \epsilon_{2}\right)\right) \rightarrow \ldots \rightarrow H_{*}\left(\check{C}\left(X, \epsilon_{n}\right)\right)$ is called a persistence module. Note that it is a finite-dimensional representation since all the vector spaces are finite-dimensional.

Remark 5.31. By saying 'we apply the homology functor $H_{*}$ ', we mean that we apply the functor $H_{p}$ for some integer $p \geq 0$. We use the notation $H_{*}$ to illustrate that any choice of $p \geq 0$ is valid.

Applying Gabriel's Theorem for $A_{n}$-type quivers (Theorem 4.7) yields that the persistence module is isomorphic to its interval decomposition $\bigoplus_{j=1}^{r} \square_{Q}\left[b_{j}, d_{j}\right]$.
Remark 5.32. Note that this construction works for any filtration of simplicial complexes:

$$
K_{1} \xrightarrow{f_{1}} K_{2} \xrightarrow{f_{2}} K_{3} \xrightarrow{f_{3}} \ldots \xrightarrow{f_{n-1}} K_{n} .
$$

Indeed, applying homology yields a persistence module that can be decomposed using Gabriel's Theorem:

$$
H_{*}\left(K_{1}\right) \xrightarrow{f_{1, *}} H_{*}\left(K_{2}\right) \xrightarrow{f_{2, *}} H_{*}\left(K_{3}\right) \xrightarrow{f_{3, *}} \ldots \xrightarrow{f_{n-1, *}} H_{*}\left(K_{n}\right) .
$$

Moreover, we can also consider simplicial complexes that are not finite-dimensional (they consist of an infinite number of simplices). In this case, the vector spaces may not be finite-dimensional and thus the persistence module may not be finite-dimensional. However, there exists a more powerful version of Gabriel's Theorem, which works in the infinite-dimensional case (see Chapter 1.1 of [1]).

As we go from $\check{C}\left(X, \epsilon_{i}\right)$ to $\check{C}\left(X, \epsilon_{j}\right)$ in a filtration (for $i<j$ ), we gain new homology classes or we lose some when they become trivial or merge with each other [7, p. 179]. To observe this, consider the images of the induced maps in homology, $f_{p}^{i, j}: H_{p}\left(\check{C}\left(X, \epsilon_{i}\right)\right) \rightarrow H_{p}\left(\check{C}\left(X, \epsilon_{j}\right)\right)$.
Definition 5.33. The p-th persistent homology groups are the images of the homomorphisms induced by inclusion, $H_{p}^{i, j}=\operatorname{im} f_{p}^{i, j}$, for $1 \leq i \leq j \leq n$. The corresponding p-th persistent Betti numbers are the ranks of these groups, $\beta_{p}^{i, j}=\operatorname{rank} H_{p}^{i, j}$.

If we consider a homology class $[c]$ in $H_{p}\left(\check{C}\left(X, \epsilon_{b}\right)\right)$ that is not in any of the $p$-th persistent homology groups $H_{p}^{i, b}$ for all $i<b$, then this homology class appears for the first time at parameter $b$. Moreover, there exists a parameter $d$ for which this homology class is non-trivial for the last time. This means that $f_{p}^{b, j}[c] \neq 0$ for all $j \leq d$ and that $f_{p}^{b, d+1}[c]=0$ (or $d=n$ ). In the persistence module, this homology class (which describes a topological feature) corresponds to the interval representation $\square[b, d]$. We conclude that each interval representation in the interval decomposition $\bigoplus_{j=1}^{r} \square_{Q}\left[b_{j}, d_{j}\right]$ corresponds to the evolution of a homology class and therefore describes the persistence of some topological feature. Given a parameterized family of spaces, those topological features that persist over a significant parameter range are to be considered as signal (resp. important features) with short-lived features as noise [8, p. 5]. This motivates the following definition.

Definition 5.34. The lifetime of an interval representation $\square_{Q}[b, d]$ is given by $d-b$. Moreover, $b$ is called the birth time and $d$ is called the death time.

We conclude that the $p$-th persistent Betti number $\beta_{p}^{i, j}$ counts the number of $p$-dimensional holes (features) that persist over the parameter interval $[i, j]$. Using the interval decomposition, we present a visualization of the $p$-th persistent homology groups and the corresponding intervals.

Definition 5.35. If a persistence module has the interval decomposition $\bigoplus_{j=1}^{r} \square_{Q}\left[b_{j}, d_{j}\right]$, then its persistence barcode is the multiset of intervals $\left\{\left[b_{j}, d_{j}\right]\right\}_{j=1}^{r}$. The horizontal axis corresponds to the parameter $\epsilon$ and the vertical axis ranges over an arbitrary ordering of the homology generators of the p-th persistent homology groups of the persistence module.

Figure 5.6 gives an example of the barcode of the persistence module coming from the point cloud from Figure 5.5. This point cloud is a sample of points from a planar annulus. We observe that there are 18 different 0 -dimensional homology classes for $\epsilon$ small enough. They correspond to the 18 points in our sample. We also observe some small intervals in dimension 1 that represent noise. Moreover, we can see the big interval in dimension 1, which represents the 1-dimensional hole in the middle of the annulus. The long length of this interval corresponds to the fact that this middle hole is a significant feature of the planar annulus.


Figure 5.6: The persistence barcode for the point cloud from Figure 5.5.

Theorem 5.36. [7, p. 181] The p-th persistent Betti number $\beta_{p}^{i, j}$ is equal to the number of intervals in the barcode spanning the whole interval $[i, j]$ (or more).

This means that a barcode is not just a great tool to visualize the decomposition of a persistence module, it also encodes the entire information about persistent homology groups.

### 5.3 Zigzag Persistent Homology

In this section, we generalize the setting from persistent homology to introduce zigzag persistent homology. We use zigzag persistent homology to show a connection to the diamond principle from quiver representation theory. This section is based on the paper 'Zigzag Persistence' by Gunnar Carlsson and Vin Silva [10].

Similarly, as in Section 5.2, we start with a sequence of simplicial complexes.
Definition 5.37. The following sequence of simplicial complexes is called a zigzag diagram:

$$
K_{1} \stackrel{f_{1}}{\longleftrightarrow} K_{2} \stackrel{f_{2}}{\longleftrightarrow} K_{3} \stackrel{f_{3}}{\longleftrightarrow} \ldots \stackrel{f_{n-1}}{\longleftrightarrow} K_{n}
$$



Figure 5.7: A zigzag diagram $K \hookleftarrow K \cap K^{\prime} \hookrightarrow K^{\prime}$. The red lines form a 1-cycle that is non-trivial in homology in all three complexes.

The major difference to the setting of persistent homology is that here, the maps $f_{i}$ can go in either direction. Again, we apply the homology functors $H_{*}$ and get a sequence in homology:

$$
H_{*}\left(K_{1}\right) \stackrel{f_{1, *}}{\longleftrightarrow} H_{*}\left(K_{2}\right) \stackrel{f_{2, *}}{\longleftrightarrow} H_{*}\left(K_{3}\right) \stackrel{f_{3, *}}{\longleftrightarrow} \ldots \stackrel{f_{n-1, *}}{\longleftrightarrow} H_{*}\left(K_{n}\right) .
$$

This is again a sequence of vector spaces with linear maps connecting them. However, since the maps $f_{i, *}$ can go from left to right or the other way around, this is a representation of an $A_{n}$-type quiver. Specifically, it is the $A_{n}$-type quiver that matches the direction of the arrows.

Definition 5.38. The representation $H_{*}\left(K_{1}\right) \leftrightarrow H_{*}\left(K_{2}\right) \leftrightarrow \ldots \leftrightarrow H_{*}\left(K_{n}\right)$ is called a zigzag persistence module.

Using Gabriel's Theorem, this representation decomposes into its interval decomposition and yields a persistence barcode analogous to any persistence module. Therefore, the theory of zigzag persistence is a generalization of the persistence theory from Section 5.2. However, since we cannot observe homology classes from their birth time to their death time (because not all maps point in the same direction), understanding the intervals in the barcode is more involved. The following example illustrates this.

Example 5.39. Let $K$ and $K^{\prime}$ be the simplicial complexes shown in Figure 5.7 and let $K \cap K^{\prime}$ be the simplex-wise intersection. This defines a zigzag diagram

$$
K \longleftrightarrow K \cap K^{\prime} \longleftrightarrow K^{\prime}
$$

Applying homology in dimension 1 (i.e. $H_{1}$ ), we get a zigzag persistence module

$$
\mathbb{Z} / 2 \mathbb{Z} \cong H_{1}(K) \longleftarrow H_{1}\left(K \cap K^{\prime}\right) \cong \mathbb{Z} / 2 \mathbb{Z} \oplus \mathbb{Z} / 2 \mathbb{Z} \longrightarrow H_{1}\left(K^{\prime}\right) \cong \mathbb{Z} / 2 \mathbb{Z}
$$

The 1-cycle $[1,2]+[1,3]+[2,3]+[3,4]$ (red lines) is non-trivial in homology in all three complexes. Therefore, one may think that this cycle corresponds to an interval representation $\square_{Q}[1,3]$ that spans all three vertices. However, this is not the case. For $H_{1}\left(K \cap K^{\prime}\right)$ we choose the following basis

$$
\{[1,2]+[1,3]+[2,3],[2,3]+[2,3]+[3,4]\} .
$$

This gives the following decomposition, which consists of the two interval representations $\square_{Q}[1,2]$ and $\square_{Q}[2,3]$.

$$
\mathbb{Z} / 2 \mathbb{Z} \stackrel{[1,0]}{\longleftrightarrow} \mathbb{Z} / 2 \mathbb{Z} \oplus \mathbb{Z} / 2 \mathbb{Z} \xrightarrow{[0,1]} \mathbb{Z} / 2 \mathbb{Z}
$$

We now generalize the setting of Example 5.39 to observe two interesting zigzag persistence modules. For a finite collection of simplicial complexes $\mathcal{X}=\left\{K_{i}\right\}_{i=0}^{n}$, consider the following zigzag diagram


Applying homology yields a zigzag persistence module

$$
H_{*}\left(K_{0}\right) \longleftarrow H_{*}\left(K_{0} \cap K_{1}\right) \longrightarrow H_{*}\left(K_{1}\right) \longleftarrow H_{*}\left(K_{1} \cap K_{2}\right) \longrightarrow \ldots H_{*}\left(K_{n}\right) .
$$

We can also consider the zigzag diagram associated with the union of simplicial complexes


We obtain another zigzag persistence module

$$
H_{*}\left(K_{0}\right) \longrightarrow H_{*}\left(K_{0} \cup K_{1}\right) \longleftarrow H_{*}\left(K_{1}\right) \longrightarrow H_{*}\left(K_{1} \cup K_{2}\right) \longleftarrow \ldots \longleftarrow H_{*}\left(K_{n}\right)
$$

The motivation behind these zigzag persistence modules is the following: if the complexes in $\mathcal{X}$ correspond to different areas of a point cloud, these two zigzag persistence modules roughly correspond to the persistence of topological features over different parts of the point cloud. The following theorem gives a set of matching rules for the interval representations of their respective interval decomposition.

Theorem 5.40 (Strong diamond principle). Let $\mathbb{V}^{+}=H_{*}\left(\mathcal{X}_{\cup}\right)$ and $\mathbb{V}^{-}=H_{*}\left(\mathcal{X}_{\cap}\right)$ be the two zigzag persistence modules corresponding to the following zigzag diagrams.

$$
\begin{array}{ll}
\mathcal{X}_{\cup}: & K_{0} \longleftrightarrow \ldots \longleftrightarrow K_{k-1} \longleftrightarrow K_{k-1} \cup K_{k+1} \longleftrightarrow K_{k+1} \longleftrightarrow \ldots \longleftrightarrow K_{n} \\
\mathcal{X}_{\cap}: & K_{0} \longleftrightarrow \longleftrightarrow \longleftrightarrow K_{k-1} \longleftrightarrow K_{k-1} \cap K_{k+1} \longleftrightarrow K_{k+1} \longleftrightarrow \longleftrightarrow \longleftrightarrow K_{n} \longleftrightarrow
\end{array}
$$

The interval decompositions of $\mathbb{V}^{-}$and $\mathbb{V}^{+}$are related to each other through the following matching rules:

- summands $\square_{Q}[k, k]^{p+1}$ and $\square_{s_{k} Q}[k, k]^{p}$ are matched,
- summands $\square_{Q}[b, k]$ are matched with summands $\square_{s_{k} Q}[b, k-1]$, and $\square_{Q}[b, k-1]$ with $\square_{s_{k} Q}[b, k]$,
- summands $\square_{Q}[k, d]$ are matched with summands $\square_{s_{k} Q}[k+1, d]$, and $\square_{Q}[k+1, d]$ with $\square_{s_{k} Q}[k, d]$,
- every other summand $\square_{Q}[b, d]$ is matched with the summand $\square_{s_{k} Q}[b, d]$.

The superscripts $p+1$ and $p$ in the first matching rule denote $a \pm 1$ shift of homological dimension. This means that $\square_{Q}[k, k]$ in the zigzag persistence module of dimension $p+1$ (of $\mathbb{V}^{+}$) is matched with $\square_{S_{k} Q}[k, k]$ in the zigzag persistence module of dimension $p$ (of $\mathbb{V}^{-}$).


Figure 5.8: The matching rules from the strong diamond principle. The top row illustrates the second and third matching rules. The bottom row shows the first and last matching rules.

Proof of Theorem 5.40. For any $p \in \mathbb{N}_{0}$, the following diamond is exact:


This follows directly from the exactness of the Mayer-Vietoris sequence [9, p. 149ff]:

$$
\ldots \longrightarrow H_{p}\left(K_{k-1} \cap K_{k+1}\right) \longrightarrow H_{p}\left(K_{k-1}\right) \oplus H_{p}\left(K_{k+1}\right) \longrightarrow H_{p}\left(K_{k-1} \cup K_{k+1}\right) \longrightarrow \ldots
$$

Thus, the second, third, and fourth matching rules follow from the diamond principle (Theorem 4.11). It remains to prove the first matching rule. This follows from the Mayer-Vietoris sequence:

$$
\begin{gathered}
\ldots \longrightarrow H_{p+1}\left(K_{k-1}\right) \oplus H_{p+1}\left(K_{k+1}\right) \xrightarrow{f} H_{p+1}\left(K_{k-1} \cup K_{k+1}\right) \\
H_{p}\left(K_{k-1} \cap K_{k+1}\right) \stackrel{\partial}{\longleftrightarrow} H_{p}\left(K_{k-1}\right) \oplus H_{p}\left(K_{k+1}\right) \longrightarrow \ldots
\end{gathered}
$$

Using exactness and the first isomorphism theorem, we conclude that

$$
\operatorname{coker}(f)=H_{p+1}\left(K_{k-1} \cup K_{k+1}\right) / \operatorname{im}(f)=H_{p+1}\left(K_{k-1} \cup K_{k+1}\right) / \operatorname{ker}(\partial) \cong \operatorname{im}(\partial)=\operatorname{ker}(g)
$$

We notice that coker $(f)$ is spanned by homology classes that do not come from homology classes in $H_{p+1}\left(K_{k-1}\right)$ or $H_{p+1}\left(K_{k+1}\right)$. These homology classes (of dimension $p+1$ ) correspond exactly to the representations $\square_{Q}[k, k]$. Similarly, elements in $\operatorname{ker}(g)$ are homology classes (of dimension $p$ ) that are trivial when observed in $H_{p}\left(K_{k-1}\right)$ and $H_{p}\left(K_{k+1}\right)$. Therefore $\operatorname{ker}(g)$ corresponds to the interval representations $\rrbracket_{s_{k} Q}[k, k]$. This proves the first matching rule.

We have proven the diamond principle using results about reflection functors. Therefore, using results from quiver representation theory, we have shown a result about the interval decomposition (and thus the barcode) of zigzag persistence modules.

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## Titel der Arbeit:

## Gabriel's Theorem

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[^1]
[^0]:    3.1: We have used the explicit form of the Tits form $q_{Q}$ for an $A_{n}$-type quiver and the fact that the Tits form is independent of the orientation. We did not need radical vectors and results about tame and wild quivers.

[^1]:    ${ }^{1}$ z. B. ChatGPT, DALL E 2, Google Bard
    ${ }^{2}$ z. B. ChatGPT, DALL E 2, Google Bard
    ${ }^{3}$ z. B. ChatGPT, DALL E 2, Google Bard

