

# $n$ -cluster tilting subcategories for truncated path algebras

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Joint work in progress with Steffen Oppermann (NTNU)

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# Introduction

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A functorially finite subcategory  $\mathcal{C} \subseteq \text{mod } \Lambda$  is called an  *$n$ -cluster tilting (CT) subcategory* if

$$\begin{aligned}\mathcal{C} &= \{X \in \text{mod } \Lambda \mid \text{Ext}_{\Lambda}^i(\mathcal{C}, X) = 0 \text{ for } 0 < i < n\} \\ &= \{X \in \text{mod } \Lambda \mid \text{Ext}_{\Lambda}^i(X, \mathcal{C}) = 0 \text{ for } 0 < i < n\}.\end{aligned}$$

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- $n \leq \text{gl. dim.}(\Lambda)$ .

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For a subcategory  $\mathcal{C} \subseteq \text{mod } \Lambda$  we set

$\mathcal{C}_{\mathcal{P}} := \{\text{isoclasses of indecomposable non projective } \Lambda\text{-modules in } \mathcal{C}\}$

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- (b)  $\tau_n : \mathcal{C}_{\mathcal{P}} \longrightarrow \mathcal{C}_{\mathcal{I}}$  and  $\tau_n^- : \mathcal{C}_{\mathcal{I}} \longrightarrow \mathcal{C}_{\mathcal{P}}$  are mutually inverse bijections.

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- (c) Let  $M \in \mathcal{C}_{\mathcal{P}}$ . Then  $\Omega^i(M)$  is indecomposable for  $1 \leq i \leq n - 1$ .
- (d) Let  $M \in \mathcal{C}_{\mathcal{I}}$ . Then  $\Omega^{-i}(M)$  is indecomposable for  $1 \leq i \leq n - 1$ .

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Examples where  $n$ -cluster tilting subcategories exist:

- tensor products of  $l$ -homogeneous  $n$ -representation-finite algebras (if  $\mathbf{k}$  is perfect) [Herschend–Iyama]
- $n$ -APR tilts of  $n$ -representation-finite algebras [Iyama–Oppermann]
- higher Nakayama algebras [Jasso–Külshammer]
- many more...

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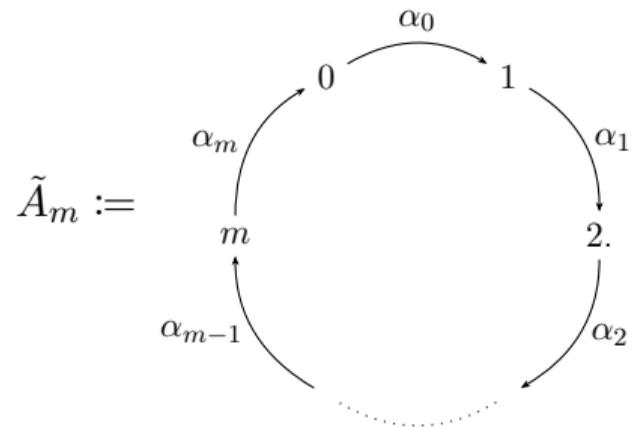
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## Question

For which  $Q$ ,  $L$  and  $n$  does there exist an  $n$ -CT subcategory/module of  $\text{mod } \Lambda$ ?

# The quivers $A_m$ and $\tilde{A}_m$

$$A_m := 1 \xrightarrow{\alpha_1} 2 \xrightarrow{\alpha_2} \dots \xrightarrow{\alpha_{m-1}} m.$$



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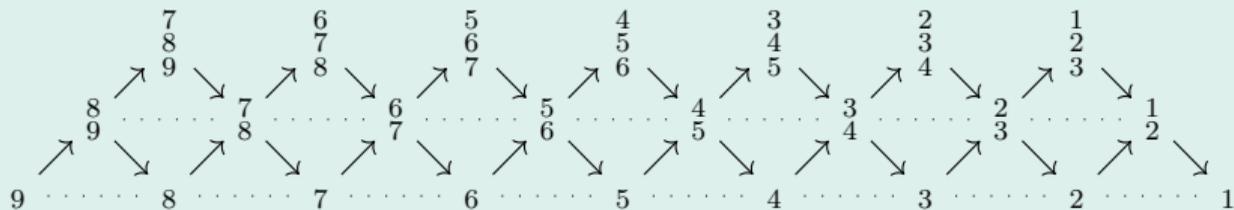
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$L = 2$  and  $L \geq 3$

Example for  $L \geq 3$

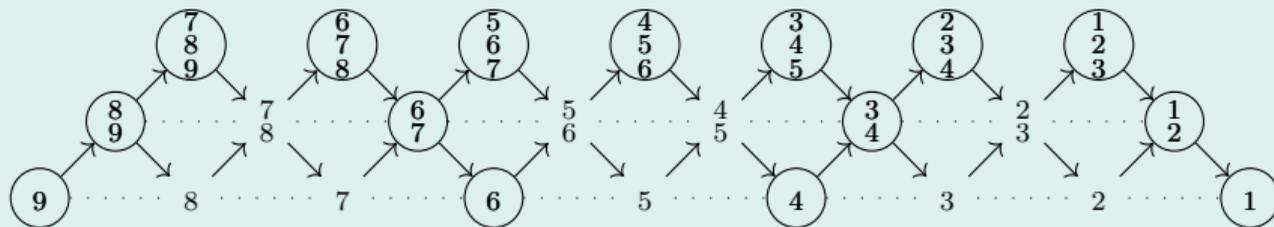
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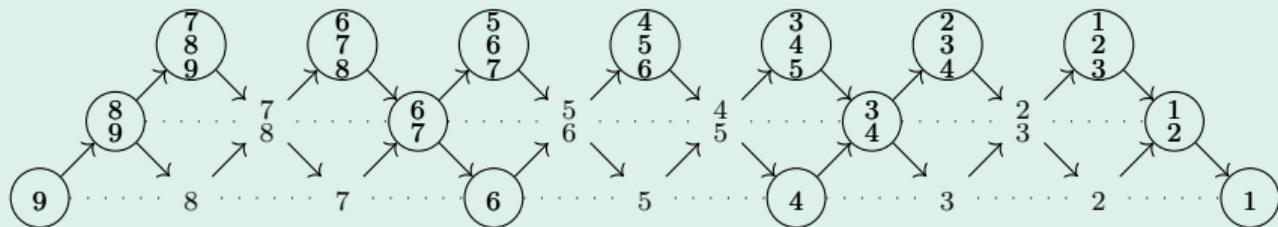


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In general, if  $Q = A_m$  and  $L \geq 3$ , and if there exists an  $n$ -CT subcategory, then  $n$  is even.

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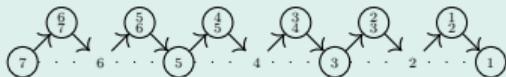
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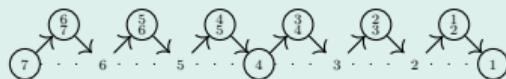
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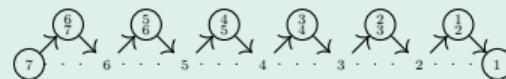
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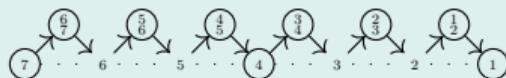
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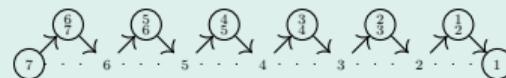
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In general, if  $L = 2$ , there is no restriction on the parity of  $n$ .

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### Theorem [Darpö–Iyama]

Let  $\Lambda = \mathbf{k}\tilde{A}_m/J^L$ . There exists an  $n$ -CT subcategory of  $\text{mod } \Lambda$  if and only if one of the following two conditions holds:

- (i)  $(2(\frac{n-1}{2}L + 1)) \mid 2(m + 1)$ , or
- (ii)  $(2(\frac{n-1}{2}L + 1)) \mid t(m + 1)$ , where  $t = \gcd(n + 1, 2(L - 1))$ .

There are many different  $n$ -CT subcategories, all of the form  $\text{add}(M)$  for some  $M \in \text{mod } \Lambda$ .

# Shape of $Q$

For a vertex  $v$  in  $Q$  we denote

- $\delta^-(v)$  := number of arrows terminating at  $v$  (incoming degree)
- $\delta^+(v)$  := number of arrows starting at  $v$  (outgoing degree)
- $\delta(v)$  :=  $(\delta^-(v), \delta^+(v))$  (degree)

# Shape of $Q$

## Proposition [Oppermann-V]

Let  $\Lambda = \mathbf{k}Q/J^L$ . Assume there exists an  $n$ -CT subcategory  $\mathcal{C} \subseteq \text{mod } \Lambda$ . Then for every  $v \in Q_0$  we have

$$\delta(v) \in \{(0,0), (0,1), (1,0), (1,1), (1,2), (2,1), (2,2)\}.$$

Moreover, if  $L \geq 3$  or  $n \geq 3$ , then  $\delta(v) \neq (2,2)$ .

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Moreover, if  $L \geq 3$  or  $n \geq 3$ , then  $\delta(v) \neq (2,2)$ .

## Proof sketch

Assume that there are at least 3 arrows terminating at  $v$ . Show that  $\Omega(I(v))$  has at least two indecomposable summands using results of Huisgen-Zimmermann.

# Shape of $Q$

## Definition

Let  $Q$  be a quiver, let  $n \geq 2$  and let  $L \geq 2$ . We say that  $Q$  is  $(n, L)$ -pre-admissible if

- (i) every vertex of  $Q$  has at most two incoming and at most two outgoing arrows,
- (ii) no vertex of  $Q$  has degree  $(0, 2)$  or  $(2, 0)$ , and
- (iii) if  $L \geq 3$  or  $n \geq 3$ , then no vertex of  $Q$  has degree  $(2, 2)$ .

# Flow paths

## Definition

Let  $k \geq 2$ . A  **$k$ -flow path**  $\mathbf{v}$  in  $Q$  is a path

$$\mathbf{v} = v_1 \xrightarrow{\alpha_1} v_2 \xrightarrow{\alpha_2} \cdots \xrightarrow{\alpha_{k-2}} v_{k-1} \xrightarrow{\alpha_{k-1}} v_k$$

such that

- $\delta(v_1) \neq (1, 1)$ ,
- $\delta(v_k) \neq (1, 1)$ , and
- $\delta(v_i) = (1, 1)$  for all  $1 < i < k$ .

We define the **degree** of  $\mathbf{v}$  to be  $\delta(\mathbf{v}) = (\delta^-(\mathbf{v}), \delta^+(\mathbf{v})) := (\delta^-(v_1), \delta^+(v_k))$ .

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Note: if  $Q$  is  $(n, L)$ -pre-admissible, then there exists a  $k$ -flow path if and only if  $Q \neq A_1$  and  $Q \neq \tilde{A}_m$ .

## Length of flow paths

Let  $Q$  be  $(n, L)$ -pre-admissible and let  $\mathbf{v}$  be a  $k$ -flow path in  $Q$ . We define  $r(\mathbf{v}, L)$  depending on the degrees of  $v_1$  and  $v_2$  as in the following table:

$\delta(v_1) \backslash \delta(v_k)$	(1,0)	(2,1)	(1,2)	(2,2)
(0,1)	$\frac{L}{2}$	1	0	1
(1,2)	1	$2 - \frac{L}{2}$	$1 - \frac{L}{2}$	1
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### Example

Let  $\mathbf{v}$  be a  $k$ -flow path with  $\delta(v_1) = (1, 2)$  and  $\delta(v_k) = (2, 1)$ . Then  $r(\mathbf{v}, 4) = 2 - \frac{4}{2} = 0$ .

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Let  $Q$  be an  $(n, L)$ -pre-admissible quiver and  $\mathbf{v}$  be a  $k$ -flow path in  $Q$ . We say that  $\mathbf{v}$  is  $(n, L)$ -admissible if there exists an integer  $p_{\mathbf{v}} \geq 0$  such that

$$k = (p_{\mathbf{v}} + 1) \left( \frac{n-1}{2} L + 1 \right) + r(\mathbf{v}, L)$$

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- (ii)  $L \geq 3$ ,  $n$  and  $p_{\mathbf{v}}$  are both even and  $\delta(\mathbf{v}) = (0, 0)$ ,

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## Definition

Let  $Q$  be an  $(n, L)$ -pre-admissible quiver and  $\mathbf{v}$  be a  $k$ -flow path in  $Q$ . We say that  $\mathbf{v}$  is  $(n, L)$ -admissible if there exists an integer  $p_{\mathbf{v}} \geq 0$  such that

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- (iii)  $L \geq 3$ ,  $n$  and  $p_{\mathbf{v}}$  are both even,  $n + p_{\mathbf{v}} > 2$  and  $\delta(\mathbf{v}) \in \{(1, 1), (1, 2), (2, 1), (2, 2)\}$ ,  
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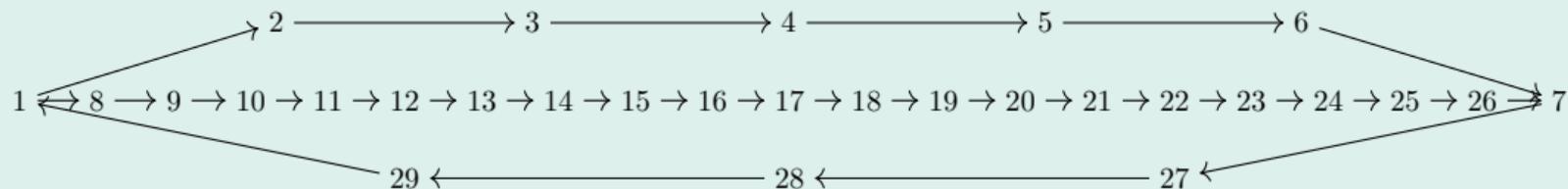
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or
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# Length of flow paths

## Example

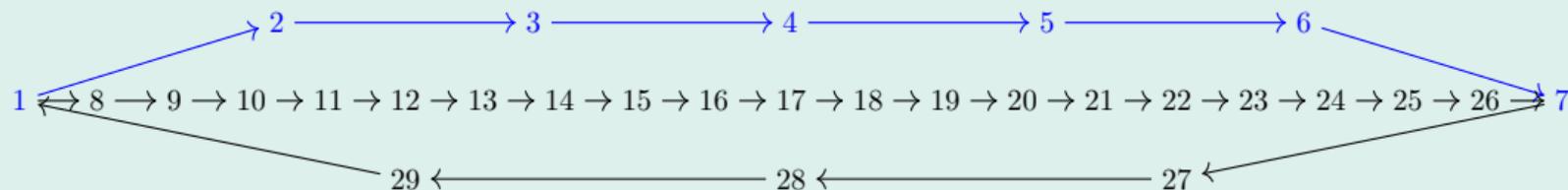
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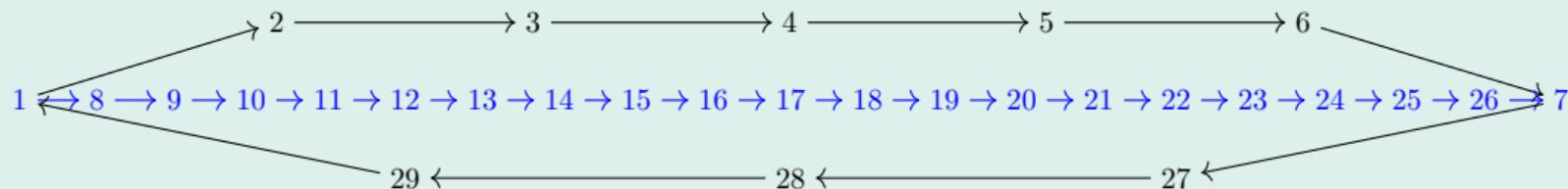


We have  $r(\mathbf{v}, 4) = 2 - \frac{4}{2} = 0$ . Since  $7 = (0 + 1) \left( \frac{4-1}{2} 4 + 1 \right) + 0$ ,  $\mathbf{v}$  is  $(4, 4)$ -admissible ( $p_{\mathbf{v}} = 0$ .)

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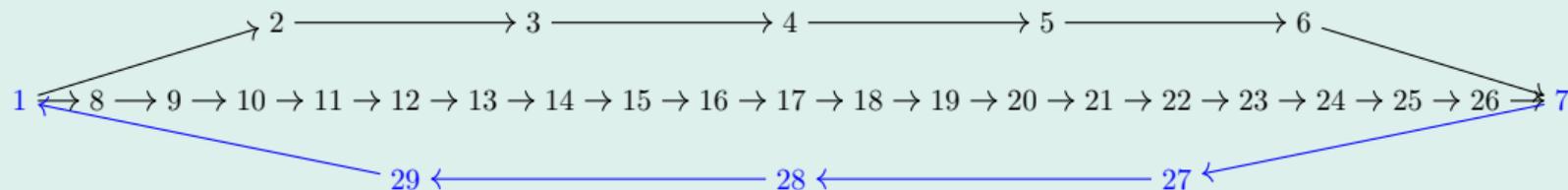


We have  $r(\mathbf{v}, 4) = 2 - \frac{4}{2} = 0$ . Since  $21 = (2 + 1) \left( \frac{4-1}{2} 4 + 1 \right) + 0$ ,  $\mathbf{v}$  is  $(4, 4)$ -admissible ( $p_{\mathbf{v}} = 2$ .)

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We have  $r(\mathbf{v}, 4) = -\frac{4}{2} = -2$ . Since  $5 = (0 + 1) \left( \frac{4-1}{2} 4 + 1 \right) - 2$ ,  $\mathbf{v}$  is  $(4, 4)$ -admissible ( $p_{\mathbf{v}} = 0$ .)

# Length of flow paths

## Proposition [Oppermann–V]

Let  $\Lambda = \mathbf{k}Q/J^L$ . Assume there exists an  $n$ -CT subcategory  $\mathcal{C} \subseteq \text{mod } \Lambda$ . Then every flow path in  $Q$  is  $(n, L)$ -admissible.

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To prove this, first we show the following.

## Lemma [Oppermann–V]

Let  $\Lambda = \mathbf{k}Q/J^L$  and let  $L \geq 3$ . Assume there exists an  $n$ -CT subcategory  $\mathcal{C} \subseteq \text{mod } \Lambda$ . If  $\mathbf{v}$  is a  $k$ -flow path in  $Q$ , then  $k \geq L + 1$ .

# Injective non-projective indecomposables

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Now let

$$\mathbf{v} = v_1 \xrightarrow{\alpha_1} v_2 \xrightarrow{\alpha_2} \cdots \xrightarrow{\alpha_{k-2}} v_{k-1} \xrightarrow{\alpha_{k-1}} v_k$$

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Then

$$\delta(v_1) \in \{(0, 1), (1, 2), (2, 1), (2, 2)\} \text{ and } \delta(v_k) \in \{(1, 0), (2, 1), (1, 2), (2, 2)\},$$

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We want to define  $L - 1$  indecomposable injective non-projective  $\Lambda$ -modules which depend on  $\delta(v_1)$ .

# Injective non-projective indecomposables

Case  $\delta(v_1) = (0, 1)$ : then we have

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# Injective non-projective indecomposables

Case  $\delta(v_1) = (2, 1)$ : then we have

$$\begin{array}{cccccccccccccccc} \dots & \xrightarrow{\alpha_{1-L}} & v_{2-L} & \xrightarrow{\alpha_{2-L}} & \dots & \xrightarrow{\alpha_{-1}} & v_0 & \xrightarrow{\alpha_0} & v_1 & \xrightarrow{\alpha_1} & v_2 & \xrightarrow{\alpha_2} & \dots & \xrightarrow{\alpha_{k-2}} & v_{k-1} & \xrightarrow{\alpha_{k-1}} & v_k \\ \dots & \xrightarrow{\beta_{1-L}} & u_{2-L} & \xrightarrow{\beta_{2-L}} & \dots & \xrightarrow{\beta_{-1}} & u_0 & \xrightarrow{\beta_0} & v_1 & & & & & & & & & \end{array}$$

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Dually we define  $P_{\mathbf{v}}(i)$  for  $1 \leq i \leq L-1$ .

## Length of flow paths

Now to show that a  $k$ -flow path  $\mathbf{v}$  must be  $(n, L)$ -admissible, we compute

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An explicit computation of the above isomorphism gives the condition on the length of  $\mathbf{v}$ .

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### Definition

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- (a)  $Q = \tilde{A}_m$  and  $(2 \binom{n-1}{2} L + 1) \mid 2(m + 1)$ , or
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## $(n, L)$ -admissible quivers

Theorem [case  $Q = \tilde{A}_m$  Darpö–Iyama, case  $L = 2$  V, case  $L \geq 3$  Oppermann–V]

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For  $Q \neq \tilde{A}_m$ :  $(\implies)$  has been motivated. For the other direction, we first show existence of an  $n$ -CT in a universal cover of  $Q$  via a direct computation. Then we use a result of Darpö–Iyama to induce an  $n$ -cluster tilting subcategory in  $\text{mod } \Lambda$ .

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  - they are supported in exactly one vertex with degree different than  $(1, 1)$ . If that vertex has degree  $(2, 1)$  then an indecomposable has the form

$$\begin{array}{ccccccc}
 0 & \hookrightarrow & M_{v_{2-L}} & \hookrightarrow & \dots & \hookrightarrow & M_{v_0} \\
 & & & & & & \searrow \\
 & & & & & & M_{v_1} & \twoheadrightarrow & M_{v_2} & \twoheadrightarrow & \dots & \twoheadrightarrow & M_{v_L} & \twoheadrightarrow & 0, \\
 & & & & & & \swarrow \\
 0 & \hookrightarrow & M_{u_{2-L}} & \hookrightarrow & \dots & \hookrightarrow & M_{u_0}
 \end{array}$$

and similarly in other cases.

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$$k = (p_{\mathbf{v}} + 1) \left( \frac{n-1}{2} L + 1 \right) + r(\mathbf{v}, L)$$

and there are exactly  $p_{\mathbf{v}}(L - 1)$  interval modules supported in  $\mathbf{v}$  which are direct summands of  $M$ .

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$$k = (p_{\mathbf{v}} + 1) \left( \frac{n-1}{2} L + 1 \right) + r(\mathbf{v}, L)$$

and there are exactly  $p_{\mathbf{v}}(L - 1)$  interval modules supported in  $\mathbf{v}$  which are direct summands of  $M$ . If  $L \geq 3$ , then these interval modules lie in diagonals as in the case  $Q = A_m$  and this is where the parity conditions come from.

## $(n, L)$ -admissible quivers

Assume  $Q \neq \tilde{A}_m$  is  $(n, L)$ -admissible. To find an  $n$ -CT module  $M$ :

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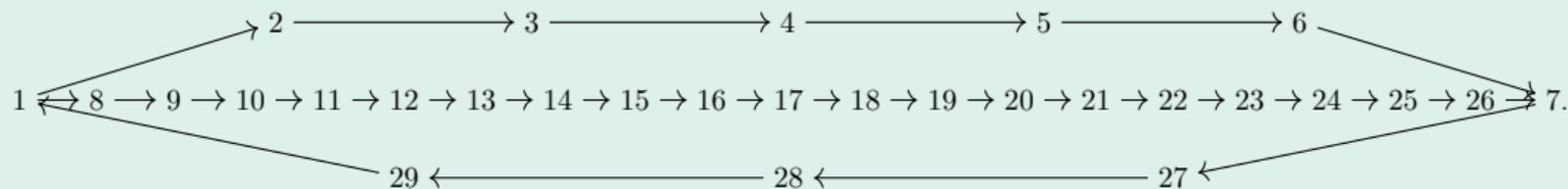
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# $(n, L)$ -admissible quivers

## Example

Let  $Q$  be the quiver

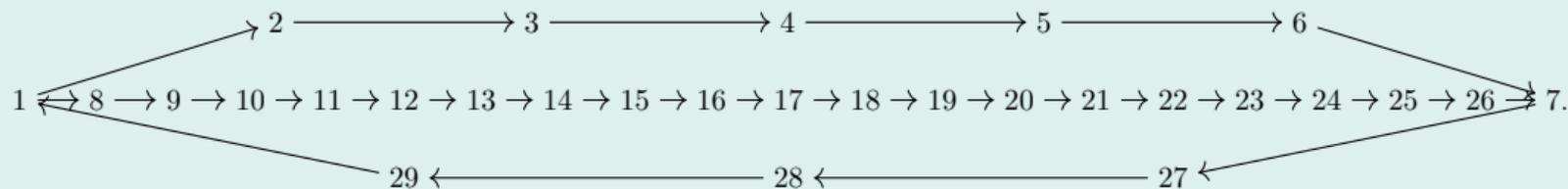


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Moreover,  $\mathcal{C} = \text{add}(M)$  where  $M$  is the direct sum of the projective modules, the injective modules, and the interval modules  $(13)$ ,  $(13, 14)$ ,  $(13, 14, 15)$ ,  $(19, 20, 21)$ ,  $(20, 21)$ ,  $(21)$ .

## How to find examples

It is easy to find  $(n, L)$ -admissible quivers such that  $\Lambda = \mathbf{k}Q/J^L$  is a wild algebra and admits an  $n$ -cluster tilting subcategory.

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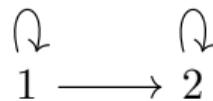
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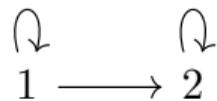
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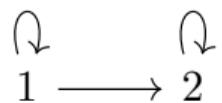


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Then extend each arrow in this graph to an  $(n, L)$ -admissible flow path. In this example, we may pick  $n = 2$  and  $p_v = 2$  for all arrows to obtain the  $(n, L)$ -admissible quiver



## $n\mathbb{Z}$ -cluster tilting subcategories

Definition [Iyama–Jasso]

An  $n$ -cluster tilting subcategory  $\mathcal{C} \subseteq \text{mod } \Lambda$  is called  $n\mathbb{Z}$ -cluster tilting if it is closed under  $\Omega^n$ .

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Theorem [Herschend–Kvamme-V, Oppermann-V]

Let  $\Lambda = \mathbf{k}Q/J^L$ . Then  $\Lambda$  admits an  $n\mathbb{Z}$ -cluster tilting subcategory if and only if one of the following conditions holds:

- (i)  $Q = A_m$  and  $L = 2$  or  $L \mid (m - 1)$ , and  $n = 2\frac{m-1}{L}$ , or
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### Corollary [Sandøy–Thibault]

Let  $\Lambda = \mathbf{k}Q/J^L$  and  $d = \text{gl. dim.}(\Lambda)$ . There exists a  $d$ -CT subcategory of  $\text{mod } \Lambda$  if and only if  $Q = A_m$  and either of  $L = 2$  or  $L \mid (m - 1)$  holds.

## A nice property for $L = 2$

### Theorem [V]

Let  $\Lambda = \mathbf{k}Q/J^2$  and let  $N$  be the largest integer for which  $Q$  is  $(N, 2)$ -admissible. Then the following hold.

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- (b) The set  $\{\mathcal{C}_n \mid n \text{ is a divisor of } N\}$  is a complete lattice with respect to inclusion isomorphic to the opposite of the lattice of divisors of  $N$ .

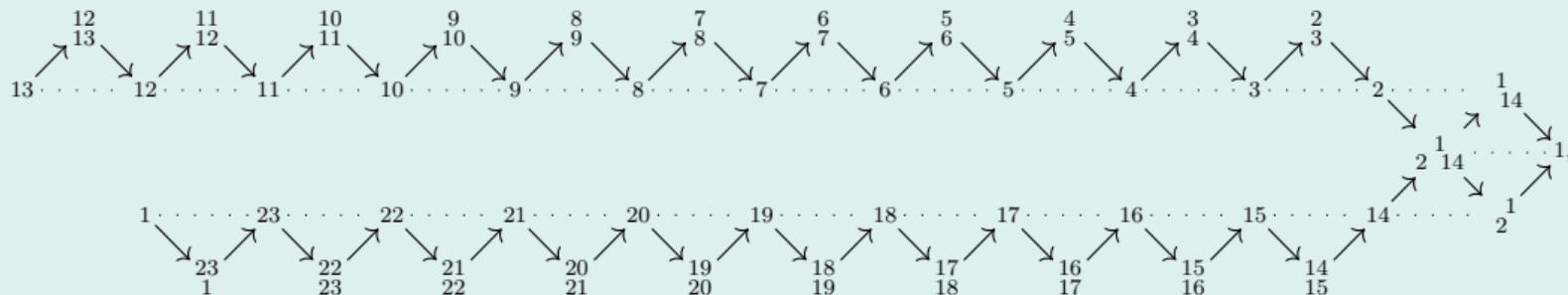




# A nice property for $L = 2$

## Example

The Auslander-Reiten quiver of  $\Lambda = \mathbf{k}Q/J^2$  is



where the simple module  $S(1)$  appears twice. Then we have

$$\mathcal{C}_1 = \text{mod } \Lambda, \quad \mathcal{C}_2 = \text{add}\{\Lambda, 11, 9, 7, 5, 3, \frac{1}{14}, 23, 21, 19, 17, 15, \frac{1}{2}\},$$

$$\mathcal{C}_3 = \text{add}\{\Lambda, 10, 7, 4, \frac{1}{14}, 22, 19, 16, \frac{1}{2}\}, \quad \mathcal{C}_4 = \text{add}\{\Lambda, 9, 5, \frac{1}{14}, 21, 17, \frac{1}{2}\},$$

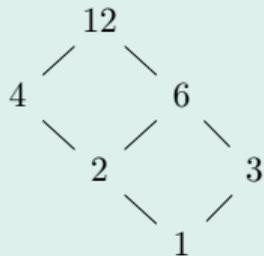
$$\mathcal{C}_6 = \text{add}\{\Lambda, 7, \frac{1}{14}, 19, \frac{1}{2}\}, \quad \mathcal{C}_{12} = \text{add}\{\Lambda, \frac{1}{14}, \frac{1}{2}\},$$

and  $\mathcal{C}_n$  is an  $n$ -cluster tilting subcategory of  $\text{mod } \Lambda$ .

# A nice property for $L = 2$

## Example

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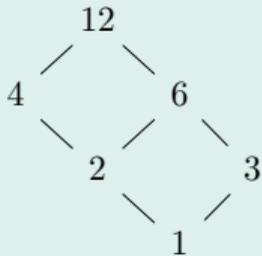


of divisors of 12

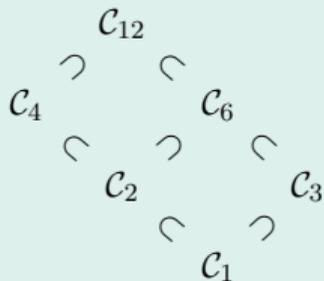
# A nice property for $L = 2$

## Example

Then the lattice



of divisors of 12 corresponds to the lattice



of inclusions of  $n$ -cluster tilting subcategories of  $\text{mod } \Lambda$ .

Thank You!