Learning Instrumental Variables with Structural and Non-Gaussianity Assumptions

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Abstract

Learning a causal effect from observational data requires strong assumptions. One possible method is to use instrumental variables, which are typically justified by background knowledge. It is possible, under further assumptions, to discover whether a variable is structurally instrumental to a target causal effect $X \rightarrow Y$. However, the few existing approaches are lacking on how general these assumptions can be, and how to express possible equivalence classes of solutions. We present instrumental variable discovery methods that systematically characterize which set of causal effects can and cannot be discovered under local graphical criteria that define instrumental variables, without reconstructing full causal graphs. We also introduce the first methods to exploit non-Gaussianity assumptions, highlighting identifiability problems and solutions. Due to the difficulty of estimating such models from finite data, we investigate how to strengthen assumptions in order to make the statistical problem more manageable.

Keywords: causality, causal discovery, instrumental variables

1. Contribution

Given observational data for a treatment variable $X$, an outcome $Y$, and a set of observed covariates $O$ that precede $X$ and $Y$ causally, we present methods to estimate the causal effect of $X$ on $Y$ when hidden common causes between $X$ and $Y$ cannot be blocked by conditioning on observed variables. The main methodology consists of discovering instrumental variables among $O$, which informally can be described as surrogate experiments on $X$ that do not affect $Y$ directly. This complements approaches where hidden common causes can be blocked. We will assume that the model is linear, although this assumption can be relaxed to some extent, as discussed in the conclusion. Much of the contribution is theoretical, and intended to describe what to the best of our knowledge is the first graphical account of the limits of what can be discovered about instrumental variables from constraints in the
observational distribution. We also discuss a pragmatic implementation of such ideas and clarify their practical limitations.

Consider a linear graphical causal model (Spirtes et al., 2000; Pearl, 2000) where, given a directed acyclic graph (DAG) $G$ with vertex set $V$, we define a joint distribution in terms of conditional relationships between each variable $V_i$ and its given parents in $G$:

$$V_i = \sum_{V_j \in \text{par}_G(i)} \lambda_{ij} V_j + e_i. \quad (1)$$

That is, each random variable $V_i$ corresponds to a vertex in $G$, where $\text{par}_G(i)$ are the parents of $V_i$ in $G$ and $e_i$ is an independent error term. Equation (1) is called a structural equation in the sense that it encodes a relationship that remains stable under a perfect intervention on other variables. Following the notation of Pearl (2000), we use the index "$\text{do}(V_k = v_k)$" to denote the regime under which variable $V_k$ is fixed to some level $v_k$ by an external agent. If $V_k$ is a parent of $V_i$, the (differential) causal effect of $V_k$ on $V_i$ is defined as:

$$\frac{\partial \mathbb{E}[V_i | \text{do}(V_k = v_k)]}{\partial v_k} = \lambda_{ik}. \quad (2)$$

Each $\lambda_{ik}$ will be referred to as a structural coefficient. Our goal is to estimate the differential causal effect of some treatment $X$ on some outcome $Y$ from observational data. If the common hidden causes of these two variables can be blocked by other observable variables, a formula such as the back-door adjustment of Pearl (2000) or the Prediction Algorithm of Spirtes et al. (2000) can be used to infer it. In general, unmeasured confounders of $X$ and $Y$ might remain unblocked.

When unmeasured confounding remains, and where it is reasonable to assume the linear structure (1), a possibility is to exploit an instrumental variable (or instrument, or IV; Morgan and Winship, 2015): some observable variable $W$ that is not an effect of either $X$ and $Y$, it is unconfounded with $Y$, and has no direct effect on $Y$.

Figure 1 illustrates one possible DAG containing an instrument. Let $\sigma_{ab}$ represent the covariance of two variables $A$ and $B$. The parameterization in (1) implies $\sigma_{wx} = \lambda_{xw} \sigma_{ww}$ and $\sigma_{wy} = \lambda_{yx} \lambda_{xw} \sigma_{ww}$. It follows that $\lambda_{yx} = \sigma_{wy}/\sigma_{wx}$. We can estimate $\sigma_{wy}$ and $\sigma_{wx}$ from observations, allowing for a consistent estimate of $\lambda_{yx}$. Notice that $\sigma_{wx} \neq 0$ is required. $W$ in this case is called an instrumental variable for the causal relationship $X \rightarrow Y$.

It is not possible to test whether some observable variable is an IV from its joint distribution with $X$ and $Y$ alone. IV assumptions can nevertheless be falsified by exploiting constraints in the joint distribution of multiple observable variables (Chu et al., 2001; Brito and Pearl, 2002; Kuroki and Cai, 2005), where such constraints are necessary but not sufficient to identify IVs and the corresponding causal effects. Our contribution are IV discovery algorithms for causal effect estimation. We characterize in which ways such algorithms can find the correct causal effect, and in which sense they will fail. We also introduce variations of the assumptions that are needed for practical reasons, complementing existing methods that rely on other sets of assumptions. Some of the ideas used in our methods are based on principles from causal discovery in linear non-Gaussian models (Shimizu et al., 2006).

The structure of the paper is as follows. In Section 2, we discuss the challenges of inferring causal effects when treatment and outcome are confounded by unobserved variables, and
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Figure 1: A graph illustrating a possible IV structure. $X$ and $Y$ have an unmeasured confounder $U$. $W$ is an instrument as it is unconfounded with $Y$, has no direct effect on it, and causes $X$. In this paper, variables named “$U$” will denote hidden variables.

provide an overview of our approach. In Section 3, we discuss the theory behind two classes of testable constraints that can be detected from data. The resulting algorithm has several practical issues, and a more realistic alternative is provided in Section 4, which is then validated experimentally in Section 5. Other related approaches are discussed in Section 6.

2. Outline of Methodology: Learning under Unmeasured Confounding

We assume that the system of interest is a linear causal model which is Markov with respect to an unknown DAG $\mathcal{G}$ with vertices $V \equiv O \cup U \cup \{X,Y\}$. Set $O$ is the set of observable variables besides treatment variable $X$ and outcome variable $Y$. Set $U$ is the set of unobservable (also called latent, or hidden) variables. $X$ and $Y$ do not precede any element of $V \setminus \{X,Y\}$. $Y$ does not precede $X$. The task is to estimate the differential causal effect of $X$ on $Y$. Our problems of interest are formulated so that all have a target causal effect $\lambda_{yx}$ corresponding to edge $X \rightarrow Y$. To avoid constant repetition, in several results that appear in this paper, we will often refer to vertices and edges of a (possibly unnamed) causal graph $\mathcal{G}$ while implicitly assuming the structural assumptions and notation introduced in this paragraph.

Focusing on a single $\lambda_{yx}$ is common in applied sciences: in particular, scenarios where we have a target causal effect $X \rightarrow Y$ to be estimated, and a set $O$ of covariates preceding $X$ and $Y$ is available. Morgan and Winship (2015) provide several examples. This is in contrast to the more familiar causal structure discovery tasks in the machine learning literature, where an equivalence class of a whole causal system is learned from data, and where some causal queries may or may not be identifiable (Spirtes et al., 2000). The focus here is on quantifying the strength of a particular causal effect $\lambda_{yx}$, as opposed to unveiling the directionality and connections of a causal graph. This allows more focused algorithms that bypass a full graph estimation. This philosophy was exploited by Entner et al. (2012) in the task of finding possible sets of observable variables that can block the effect of any hidden common cause of $X$ and $Y$.

However, the approach by Entner et al. will not provide a causal effect estimate if such a set does not exist. For instance, in Figure 1, if $U$ is a latent variable, their algorithm will provide no answer concerning the causal effect of $X$ on $Y$. Our goal is to cover this scenario, which complements approaches that require unmeasured confounding to be blocked. The methodological framework to accomplish this task is by discovering candidate instru-
Figure 2: (a) Variable W is an instrument for the relation X \rightarrow Y conditioning on \{Z_1, Z_2\}. This can be verified by noticing that W is d-separated from Y given \{Z_1, Z_2\} in the graph obtained by removing edge X \rightarrow Y from the given figure. (b) Both W_1 and W_2 are instruments conditioning on the empty set \emptyset. (c) The typical covariance constraints (“tetrads”) that are implied by instrumental variables also happen in the case where no instruments exist, implying that rank constraints in the covariance matrix are not enough information to discover IVs. (d) A case that is difficult even when considering information from non-Gaussian distributions.

mental variables without prior knowledge of the causal structure, besides the stated ordering assumptions about \{X, Y\}. The challenge is that we cannot guarantee which candidate instrumental variables are actual instruments without further assumptions.

We will make use of structural characterizations of causality using graphical models (Pearl, 2000; Spirtes et al. 2000). Prior exposure to causal graphical models is assumed, with key definitions summarized in Section 2.1. In Section 2.2, we outline the challenges and explain the general concept of equivalence class of causal effects, a concept adapted from the work of Maathuis et al. (2009) to the instrumental variable case. Finally, in Section 2.3 we provide a road map of the steps used in our approach.

2.1 Graphical Terminology

Directed acyclic graphs (DAGs) encode independencies among vertices, which correspond to independencies among random variables. This follows from the usual one-to-one relationship between vertices in a graph and random variables in a distribution. In a DAG, any two vertices \{V_i, V_j\} can be connected by up to one directed edge. Given an edge \(V_i \rightarrow V_j\), we say that \(V_i\) is the tail and \(V_j\) is the head of the edge. We also say that \(V_i\) and \(V_j\) are endpoints of edge \(V_i \rightarrow V_j\) where \(V_i\) is a parent of \(V_j\) and \(V_j\) is a child or \(V_i\).

A path in a graph is a sequence of edges, where any two consecutive edges in the sequence share a common endpoint. Paths are simple if no vertex appears twice, and non-simple
otherwise. For instance, \( W \rightarrow X \leftarrow U \rightarrow Y \) is a path in Figure 1. An example of a non-simple path is \( W \rightarrow X \leftarrow U \rightarrow X \leftarrow U \rightarrow Y \). For most of our presentation, it is sufficient to just assume that the paths referred to are simple unless stated otherwise. A collider in a path is a vertex that is common to two consecutive edges such that this vertex is the head of both edges. For instance, \( X \) is a collider in the path given in the previous example, while \( U \) is not. A vertex \( V \) is an endpoint of a path \( P \) if it is one of the endpoints of the first or last edge \( E \) in the path, and \( V \) is not shared with the (possible) edge next to \( E \) in \( P \). For instance, \( W \) and \( Y \) are the endpoints of \( W \rightarrow X \leftarrow U \rightarrow Y \).

A path \( P \) is between vertices \( V_i \) and \( V_j \) if \( V_i \) and \( V_j \) are the endpoints of \( P \). A trek is a path with no colliders. For instance, \( W \rightarrow X \leftarrow U \rightarrow Y \) is not a trek, but \( X \leftarrow U \rightarrow Y \) and \( W \rightarrow X \rightarrow Y \) are. A trek \( P \) must have an unique source, the vertex in \( P \) that is not the head of any edge in \( P \). A special case of a trek is a directed path, which is a trek where the source is one of the endpoints. For instance, \( W \rightarrow X \rightarrow Y \) is a directed path between \( W \) and \( Y \) where \( W \) is the source. We also say that this path is from source \((W)\) into the other endpoint \((Y)\). The source \( V_i \) in a directed path \( P \) is an ancestor of all elements \( V_j \) in \( P \), while \( V_j \) is a descendant of \( V_i \). It is possible that \( V_i = V_j \), so \( V_i \) is an ancestor and descendant of itself. A non-directed path is a path that is not directed.

A back-door (path) between \( V_i \) and \( V_j \) is a trek that is into \( V_i \) and \( V_j \). For instance, \( X \leftarrow U \rightarrow Y \) is a back-door between \( X \) and \( Y \).

A vertex \( V \) is active on a path with respect to some vertex set \( S \) if it is either (i) a collider in this path and itself or one of its descendants is in \( S \); or (ii) not a collider and not in \( S \). A path is active if all of its vertices are active, blocked otherwise. The notion of active and blocked paths will be important in the sequel, as activation has implications on which variables can be considered to be an instrument with respect to which conditioning sets.

These definitions lead to the concept of \( d \)-separation (Pearl, 2000). A vertex \( V_i \) is \( d \)-separated from a vertex \( V_j \) given a set \( S \) if and only if every path between \( V_i \) and \( V_j \) is blocked by \( S \). The interpretation of this definition is explained at length in the graphical modeling literature and we will assume familiarity with it. We say a probabilistic model \( \mathcal{M} \) is Markov with respect to a DAG \( \mathcal{G} \) if every \( d \)-separation in \( \mathcal{G} \) corresponds to a conditional independence constraint in \( \mathcal{M} \). A model \( \mathcal{M} \) is faithful to \( \mathcal{G} \) if every \( d \)-separation in \( \mathcal{G} \) corresponds to a conditional independence constraint in \( \mathcal{M} \) and vice-versa. We sometimes use the notation \( V_i \perp \!\!\!\perp V_j \mid S \) (Dawid, 1979) to denote conditional independence of \( V_i \) and \( V_j \) given \( S \) in a distribution, and also to denote \( d \)-separation when applied to vertex sets in a graph.

### 2.2 Scope and Fundamental Challenges

The identification of structural coefficients from given causal graphs is a classical problem in the structural equation modeling literature (Bollen, 1989). Much progress has been achieved on describing increasingly intricate combinations of structural features that lead to the identification of such coefficients (Brito and Pearl, 2002; Foygel et al., 2011; Chen et al., 2014). As these sophisticated criteria also lead to constraints which are hard to detect from data, we focus instead on the class of structures that corresponds to classical accounts of instrumental variables (Angrist and Pischke, 2009) as described by Brito and Pearl (2002).
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Brito and Pearl’s criteria are as follows. Given the causal graph $G$ of a system that includes an edge $X \rightarrow Y$, a vertex $W$ is a (conditional) instrument variable for $X \rightarrow Y$ given $Z$ if and only if:

1. $Z$ does not d-separate $W$ from $X$ in $G$, and

2. $Z$ d-separates $W$ from $Y$ in the graph obtained by removing the edge $X \rightarrow Y$ from $G$, and

3. $Z$ are non-descendants of $X$ and $Y$ in $G$.

For the rest of the paper, we will call the above conditions the graphical criteria for instrumental variable validity, or simply “Graphical Criteria.” Notice that the validity of a vertex as an IV is dependent on which set $Z$ we condition on. That is, if in the corresponding causal graph we find some set $Z$ that blocks all and only the paths relevant to the Graphical Criteria, then we can identify $\lambda_{yx}$ as $\sigma_{wy,x} / \sigma_{wx,z}$, where $\sigma_{ab,s}$ represents the partial covariance\(^1\) of two variables $A$ and $B$ given some set $S$. Figure 2(a) illustrates a case.

Unless strong background knowledge is available, the relevant structure needs to be learned from the data. As an example, consider Figure 1. When $U$ is unobserved and the distribution is faithful to the graph, the absence of the edge $W \rightarrow Y$ is not testable (Chu et al., 2001). However, in a situation such as Figure 2(b), the simultaneous lack of edges $W_1 \rightarrow Y$ and $W_2 \rightarrow Y$ has a testable implication, as the absence of each edge implies $\lambda_{yx} = \sigma_{w_1y} / \sigma_{w_1x}$ and $\lambda_{yx} = \sigma_{w_2y} / \sigma_{w_2x}$, respectively. This leads to a tetrad constraint,

$$
\sigma_{w_1y} \sigma_{w_2x} - \sigma_{w_1x} \sigma_{w_2y} = 0,
$$

which can be tested using observable data. Unfortunately, the tetrad constraint is necessary, but not sufficient, to establish that both elements in this pair of variables are instrumental. As an example, consider Figure 2(c). It is not hard to show that $\sigma_{w_1y,z} \sigma_{w_2z} - \sigma_{w_1x,z} \sigma_{w_2y,z} = 0$. However, the Graphical Criteria for IVs are not satisfied as $W_1$ is not d-separated from $Y$ given $Z$ if we remove edge $X \rightarrow Y$. This is because path $W_1 \leftarrow U_1 \rightarrow U_2 \rightarrow Y$ is active. Indeed, in this case $\lambda_{yx}$ can be much different from $\sigma_{w_1y,z} / \sigma_{w_1x,z}$. A major component of our contribution is to characterize graphically in which ways the solution is not unique.

We can summarize the main identification challenge as follows. In the Graphical Criteria, the challenging condition is the second, as the first is easily testable by faithfulness and the third is given by assumption. Another way of phrasing condition 2, with respect to conditioning set $Z$, is as follows:

2a. there is no active non-directed path between $W$ and $Y$ that does not include $X$ (that is, no active back-door path nor any active path that includes a collider), and

2b. there is no active directed path from $W$ to $Y$ that does not include $X$.

While we will be able to check 2a, we will not be able to verify 2b, learning to multiple candidate causal effects.

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\(^1\) Analogous to the (more widely used) notion of partial correlation, we define partial covariance as the covariance between the residuals of the least-squares regression of $A$ and $B$ on $S$, given by $\sigma_{ab,s} \equiv \sigma_{ab} - \sum_{i} \sum_{j} \Sigma_{i,j} a_i b_j$, where $\Sigma$ is the joint covariance matrix and $\sigma_{ij} \equiv \Sigma_{i,j}$ for scalars $i, j$. 
Figure 3: In this model, variables $W_1, \ldots, W_p$ are instrumental variables conditioning on the empty set. However, conditioning on $Z_0$ will introduce an active path from each $W_i$ to $Y$ via $U$, destroying their validity. This is particularly an issue for algorithms such as sisVIVE (Kang et al., 2016) and IV-BY-MAJORITY$_\infty$ (Section 3), where each variable is either deemed an IV or a conditioning variable. An algorithm such as IV-TETRAD$_\infty$ (Section 3) is able to avoid this mistake, although the computational cost can be high.

2.3 Roadmap

In general, the best we can do is to provide a set of candidate causal effects, one of which will be correct if at least two instrumental variables (under the same conditioning set) are present in the true graph. In this case, the set can be used, for instance, to provide lower and upper bounds on the causal effect. We will discuss our methods in the context of different assumptions about the existence of instruments, particularly about the number of instruments that exist in the true unknown graph.

That is, we propose algorithms that in theory return equivalence classes of causal effects. Such a class of algorithms will be sound in this equivalence class sense of returning a set of candidate effects that includes the true effect if there are instrumental variables, but will be incorrect otherwise. Following Cooper (1997), Mani et al. (2006) and Entner et al. (2012), we will not need to discover full graphs in order to identify the causal effect, but we will not also find all causal effects that are identifiable from faithfulness assumptions. In particular, we will consider in which sense our algorithms are complete: that is, if there are instrumental variables satisfying the Graphical Criteria, we will characterize under which conditions we will find them. Computational considerations are relevant here.

In the next section, we introduce algorithms that return an equivalence class of causal effects using tetrad constraints, which are complemented by non-Gaussianity assumptions. Motivated by simplicity of presentation, all approaches in Section 3 assume that the distribution of the population is known and that computational resources are unbounded. We do not claim these algorithms are practical — the goal is to use them as a theoretical basis to choose and justify stronger assumptions that achieve practical learning. In Section 4, we discuss practical methods for learning from data and the computational and identification compromises we adopt.

A roadmap for all algorithms introduced in this paper is shown in Table 1. In the table, we provide the main assumptions used by each method and whether or not the algorithm
Table 1: All algorithms introduced in this paper, along with their main assumptions and whether they are computationally tractable. Algorithms 2-5 assume also a form of faithfulness discussed in Section 3. Further assumptions can be added on top of the assumptions listed below in order to reduce the size of the equivalence class returned, as discussed in Sections 3 and 4.

<table>
<thead>
<tr>
<th>Algorithm</th>
<th>Main Assumptions</th>
<th>Tractable</th>
</tr>
</thead>
<tbody>
<tr>
<td>1: IV-BY-MAJORITY_{\infty}</td>
<td>there exists a set ( W \subseteq O ), of size larger than (</td>
<td>O</td>
</tr>
<tr>
<td>2: IV-TETRAD_{\infty}</td>
<td>there exists a set ( W \subseteq O ), of size at least 2, and a set ( Z \subseteq O\setminus W ), such that each ( W \in W ) is a valid IV given ( Z ).</td>
<td>No</td>
</tr>
<tr>
<td>3: IV-TETRAD_{\infty}^{+}</td>
<td>the assumptions of IV-TETRAD_{\infty}, plus the assumption that no non-directed paths between each ( W ) and ( Y ) are active given ( Z \cup W \setminus {W} ), and that the distribution is non-Gaussian.</td>
<td>No</td>
</tr>
<tr>
<td>4: IV-TETRAD^{+}</td>
<td>the assumptions of IV-TETRAD_{\infty}^{+}, plus the assumptions that ( Z = O \setminus W ) and (</td>
<td>W</td>
</tr>
<tr>
<td>5: IV-TETRAD^{++}</td>
<td>the assumptions of IV-TETRAD^{+}.</td>
<td>Yes (by MCMC)</td>
</tr>
</tbody>
</table>

is tractable as a function of \( |O| \), the number of observed variables other than treatment \( X \) and outcome \( Y \). The relationship among these algorithms is as follows.

Algorithm 1, introduced at the beginning of Section 3, is inspired by current approaches for learning instrumental variables. One of its main requirements is the existence of a large number of valid instruments, which grows with \( |O| \). It does not, however, require faithfulness.

Algorithms 2 and 3, introduced in Sections 3.1 and 3.3 respectively, are motivated by the requirement of only two valid instruments. Also, there is no requirement concerning a special choice for the respective conditioning set: it can be a strict subset of the remaining variables. The two algorithms, however, assume a variant of faithfulness explained in the next section. Algorithm 3 requires stronger assumptions than Algorithm 2, but as it will be shown, it will return an equivalence class no larger than Algorithm 2 while sometimes guaranteeing a strictly smaller equivalence class.

Unlike the previous algorithms, Algorithm 4, introduced in Section 4, is aimed at practical learning with finite sample sizes and limited computational resources. That is, here computational tractability and the handling of statistical variability are important issues to be considered. We achieve a reduction on statistical variability by enforcing a stronger constraint: the unknown true model must contain a set of valid IVs no smaller than \( K \),
Algorithm 1  IV-BY-MAJORITY

1: Input: set of random variables $\mathbf{O} \cup \{X, Y\}$ with known joint distribution
2: Output: the causal effect of $X$ on $Y$, or a value (NA) indicating lack of knowledge
3: for each $W_i \in \mathbf{O}$ do
4: \hspace{1em} $\mathbf{Z}_i \leftarrow \mathbf{O} \setminus \{W_i\}$
5: \hspace{1em} $\beta_i \leftarrow \sigma_{w_i y,z_i} / \sigma_{w_i x,z_i}$
6: end for
7: if more than half of set $\{\beta_i\}$ is equal to the same finite value $\beta$ then
8: \hspace{1em} return $\beta$
9: end if
10: return NA

a user-defined parameter. We achieve tractability, and also a further reduction on statistical variability, by requiring the conditioning set to be automatically determined given a candidate choice of IVs, as in Algorithm 1.

Finally, Algorithm 5 is introduced late in Section 5.3.2 as an example of a Bayesian variation of Algorithm 4. The goal is to illustrate how variations of Algorithm 4 can be created to respond to particular needs of an individual application, such as the provision of Bayesian measures of uncertainty of our estimates. The algorithm is still tractable assuming Markov chain Monte Carlo (MCMC) can be reliably run on an acceptable amount of time.

Although it seems that we take a step backwards in Algorithm 4, by enforcing back assumptions about $\mathbf{Z}$ that are similar to Algorithm 1, we will explain in Section 4 that these are not fundamental. However, practical learning with finite sample sizes is hard to achieve otherwise.

3. From Structural Constraints to Instruments and Causal Effects

Consider Algorithm 1 as a method for learning causal effects given the distribution of the population (hence, the “$\infty$” symbol in the name of the algorithm, indicating that this is equivalent to having infinite sample sizes). If $W_i$ is an IV conditioned on $\mathbf{Z}_i \equiv \mathbf{O} \setminus W_i$, then $\beta_i \equiv \sigma_{w_i y,z_i} / \sigma_{w_i x,z_i} = \lambda_{yx}$, the true causal effect. Without knowing whether $W_i$ is a conditional IV with respect to $\mathbf{Z}_i$, we cannot make claims about the causal effect.

However, if more than half of the elements $W_i \in \mathbf{O}$ are conditional IVs given the respective $\mathbf{Z}_i$, then it follows that more than half of the elements in set $\{\beta_i\}$ will be equal, and equal to $\lambda_{yx}$. This is the same core assumption introduced by Kang et al. (2016). It sidesteps the problems introduced by models such as the one in Figure 2(c) by assuming that more than half of $\mathbf{O}$ are “valid” IVs. That is, we can partition $\mathbf{O}$ into two sets, $\mathbf{O} = \mathbf{W} \cup \mathbf{Z}$, such that each $W \in \mathbf{W}$ is a conditional IV given $\mathbf{Z} \cup \mathbf{W} \setminus \{W\}$. This is done without knowledge of which variables are valid and which are not. As discussed by Kang et al. (2016), there are situations where this assumption is plausible, or at least weaker than in standard approaches, as in some genetic studies where $\mathbf{O}$ are genetic features of a cell and $X$, $Y$ are phenotypes. Kang et al. (2016) introduced the SISVIVE algorithm ("some invalid, some
Algorithm 2 IV-TETRAD\(_\infty\)

1: **Input:** set of random variables \(O \cup \{X,Y\}\) with known joint distribution

2: **Output:** \(\mathcal{C}\), a set of candidate differential causal effects of \(X\) on \(Y\)

3: Initialize \(\mathcal{C} \leftarrow \emptyset\)

4: for each pair \(\{W_i, W_j\}\) \(\subseteq O\) do

5: \hspace{1em} for every set \(Z \subseteq O \setminus \{W_i,W_j\}\) do

6: \hspace{2em} if \(\sigma_{w_i,x,z} = 0\) or \(\sigma_{w_j,x,z} = 0\) then

7: \hspace{3em} next

8: \hspace{2em} end if

9: \hspace{2em} if \(\sigma_{w_i,x,z} \sigma_{w_j,y,z} \neq \sigma_{w_i,y,z} \sigma_{w_j,x,z}\) then

10: \hspace{3em} next

11: \hspace{2em} end if

12: \hspace{1em} \(\mathcal{C} \leftarrow \mathcal{C} \cup \{\sigma_{w_i,y,z} / \sigma_{w_j,x,z}\}\)

13: end for

14: end for

15: return \(\mathcal{C}\)

valid IV estimator") which, although very different\(^2\) from Algorithm 1, relies on the same fundamental assumptions and as such will be used in our benchmark methods in Section 5.

However, this assumption can be false even when nearly all of \(O\) are possibly valid instruments. Consider Figure 3 where we have an arbitrary number of IVs \(W_1, \ldots, W_p\) that are valid by conditioning on the empty set. None of them are valid by conditioning on \(Z_0\). In this situation, IV-BY-MAJORITY\(_\infty\) will return NA and sisVIVE loses its guarantees, potentially performing badly. We will characterize how this happens by using faithfulness.

3.1 Structural Signatures of Tetrad Constraints and the Graphical Criteria

Consider Algorithm 2 as a method for learning causal effects, disregarding for now its computational complexity. The idea is to find triplets \((W_i, W_j, Z)\) such that \((W_i, Z)\) and \((W_j, Z)\) both satisfy the Graphical Criteria. Under the assumption of linear faithfulness (Spirtes et al., 2000), Line 6 of the algorithm is equivalent to item 1 of the Graphical Criteria. To characterize what Algorithm 2 can say about item 2 of the Graphical Criteria, we will need other, less usual, graphical definitions introduced in the next section.

Algorithm 2 is simple, but to the best of our knowledge is has not been described in the literature. A possible reason for the absence of this algorithm in the literature is that characterizing its limitations is not straightforward. One of the main contributions of this paper is the analysis of IV-TETRAD\(_\infty\) in Section 3.2, for which we will first need the following theory.

\(^2\) sisVIVE uses a modified sparse regression method that is arguably more statistically elegant than using, for instance, hypothesis testing to decide which estimated \(\lambda_i\) are statistically the same. The method remains computationally efficient. We introduced Algorithm 1 because it is closely related to the other algorithms discussed in this paper and hence provides more direct insights on the interplay between different assumptions and identification issues. Notice, however, that sisVIVE never returns NA as it assumes there is always a majority.
3.1.1 t-Separation and Rank Constraints

In this section, we describe the notion of t-separation (Sullivant et al., 2010), the most important graphical concept in our approach. Readers familiar with this concept can skip this section.

Recall the basic definition of a trek from Section 2.1. Another way of describing a trek $T$ is by ordering its endpoints (such that a trek is “from” $V_i$ “to” $V_j$), implying an ordered pair of (possibly empty) directed paths $(P_1; P_2)$ where: $P_1$ has sink (vertex without children in $T$) $V_i$; $P_2$ has sink $V_j$; and $P_1, P_2$ have the same source (vertex in $T$ without parents in $T$).

Definition 1 (t-separation) The ordered pair of vertex sets $(C_A; C_B)$ t-separates vertex set $A$ from vertex set $B$ if, for every trek $(P_1; P_2)$ from a vertex in $A$ to a vertex in $B$, either $P_1$ contains a vertex in $C_A$ or $P_2$ contains a vertex in $C_B$.

Spirtes (2013) and Sullivant et al. (2010) discuss a generalization of this notion and further examples. As the definition is somewhat complex, later in this section we provide some basic examples of this concept in the context of instrumental variables.

Like $d$-separation, t-separation is relevant if it implies testable implications that decrease the set of graphical structures compatible with the data (Spirtes et al., 2000). While we will see that t-separation can also imply independence constraints among observable variables, this would be of limited interest as $d$-separation can also imply such constraints. The importance of t-separation is the possibility of implying testable consequence of $d$-separations given unobserved variables. This is done under the assumption of linearity, which will lead to rank constraints (Spirtes, 2013).

Definition 2 (rank constraint) A rank constraint in a matrix $M$ is any constraint of the type $\text{rank}(M) \leq r$, where $r$ is some constant and $\text{rank}(M)$ is the rank of matrix $M$.

If $M$ is the cross-covariance submatrix given by variables $\{V_i, V_j\}$ indexing the rows, and by $\{V_k, V_l\}$ indexing the columns, then the rank constraint $\text{rank}(M) \leq 1$ implies $\sigma_{ik}\sigma_{jl} - \sigma_{il}\sigma_{jk} = 0$. That is, the determinant of $M$ is zero, and a corresponding tetrad constraint holds.

Let $\Sigma_{AB}$ be the cross-covariance matrix of set $A$ (rows) and set $B$ (columns). The DAG Trek Separation Theorem of Sullivant et al. (2010) says:

Theorem 3 (Trek Separation for DAGs) Let $G$ be a DAG with vertex set $V$. Let $A$ and $B$ be subsets of $V$. We have $\text{rank}(\Sigma_{AB}) \leq r$ in all linear structural equation models with graph $G$ if and only if there exist subsets $C_A$ and $C_B$ of $V$ with $|C_A| + |C_B| \leq r$ such that $(C_A; C_B)$ t-separates $A$ from $B$.

To jump from (testable) rank constraints to (unobservable) structural constraints in $G$, we assume our model distribution $P$ is rank-faithful to a DAG $G$ (Spirtes, 2013):

Definition 4 ((linear) rank-faithfulness) Distribution $P$ is rank-faithful to a DAG $G$ if every rank constraint holding on a covariance (sub)matrix derived from $P$ is entailed by every linear structural equation model Markov with respect to $G$.
One implication of rank-faithfulness is that \(|\Sigma_{ZZ}| > 0\) for any set \(Z\) of vertices in \(G\), a fact which we will make use of in some proofs.

Unlike d-separation, t-separation is defined by a pair of conditioning sets. Moreover, when we say that \(V_i\) is t-separated from \(V_j\) given some ordered pair \((C_i; C_j)\), the order of the sets in the conditioning pair matters. These two sets do not need to be disjoint and can also contain \(V_i\) and \(V_j\).

**Example 1** Consider Figure 1 again. In this example, \(X\) alone cannot d-separate \(W\) from \(Y\), due to the active path \(W \rightarrow X \leftarrow U \rightarrow Y\). We need also to condition on \(U\), but we cannot directly test the respective conditional independence if \(U\) is a latent variable. It is however the case that if \(A \equiv \{W\}\) and \(B \equiv \{Y\}\), then \((\emptyset; \{X\})\) t-separates \(A\) from \(B\): using the notation of Definition 1, the only trek between \(W\) and \(Y\) is \((\emptyset; W \rightarrow X \rightarrow Y)\). \(C_A = \emptyset\) and \(C_B = \{X\}\) will suffice to imply the t-separation.

Moreover, by Theorem 3, cross-covariance matrix \(\Sigma_{AB}\) has rank at most \(|\emptyset| + |\{X\}| = 1\). Notice, however, that this is true regardless or not this t-separation holds. Therefore, without knowing the graph, we still cannot exclude edge \(W \rightarrow Y\).

If we try to redefine the problem such that \(A \equiv \{W, X\}\) and \(B \equiv \{X, Y\}\), hoping to show that \(\Sigma_{AB}\) will not be full rank, then this will not work either. This is because \((\emptyset, X)\) does not imply the t-separation anymore, as there is now the trek \((X \leftarrow U; U \rightarrow Y)\) between \(A\) and \(B\). Pair \((\emptyset; \{X\})\) does not block this trek. We need the conditioning set pair to be \((\{X\}, \{X\})\). By the notation of Theorem 3, this means \(C_A = C_B = \{X\}\) and all the theorem implies is that \(\Sigma_{AB}\) has rank at most \(|\{X\}| + |\{X\}| = 2\). Once again, this is true regardless of the absence of the edge \(W \rightarrow Y\). \(\square\)

**Example 2** Consider Figure 2(b) again. Pair \(A \equiv \{W_1, W_2\}\) is t-separated from \(B \equiv \{X, Y\}\) by \((\emptyset, \{X\})\).

To check this, enumerate all simple treks \((P_1; P_2)\) such that path \(P_1\) has its sink in \(\{W_1, W_2\}\) and path \(P_2\) has its sink in \(\{X, Y\}\). We obtain:

- \((P_1^1; P_2^1) \equiv (\emptyset; W_1 \rightarrow X),\)
- \((P_1^2; P_2^2) \equiv (\emptyset; W_1 \rightarrow X \rightarrow Y),\)
- \((P_1^3; P_2^3) \equiv (\emptyset; W_1 \rightarrow W_2 \rightarrow X),\)
- \((P_1^4; P_2^4) \equiv (\emptyset; W_1 \rightarrow W_2 \rightarrow X \rightarrow Y),\)
- \((P_1^5; P_2^5) \equiv (\emptyset; W_2 \rightarrow X),\)
- \((P_1^6; P_2^6) \equiv (\emptyset; W_2 \rightarrow X \rightarrow Y),\)
- \((P_1^7; P_2^7) \equiv (W_2 \leftarrow W_1; W_1 \rightarrow X),\)
- \((P_1^8; P_2^8) \equiv (W_2 \leftarrow W_1; W_1 \rightarrow X \rightarrow Y).\)
We can then verify that, for each case, it is indeed true that \( \{X\} \) blocks all paths \( P^1_2, \ldots, P^8_2 \). Since \( |\emptyset| + |\{X\}| = 1 \), then \( \Sigma_{AB} \) has rank at most 1. This constraint is not satisfied under rank-faithfulness if, for instance, edge \( W_1 \rightarrow Y \) were to be present. That is, even though \( X \) does not d-separate \( \{W_1, W_2\} \) from \( Y \), it alone is responsible for a t-separation and no hidden variable is evoked in its testable claim of rank \( \Sigma_{AB} \leq 1 \).

Example 3 Consider Figure 2(c) again. Using standard d-separation, we can verify that \( Z \) alone does not d-separate \( \{W_1, W_2\} \) from \( \{X, Y\} \). Instead, \( \{U_1, Z\} \), or \( \{U_2, Z\} \), should be used instead. In the absence of the causal graph, we cannot discover this d-separation by directly testing the corresponding conditional independence given a latent variable. We need indirect evidence for this d-separation that uses observable variables only.

That is where conditional tetrad constraints play a role. Assuming \( Z \) is observable, we could de-activate all directed paths from \( \{W_1, W_2\} \) to \( \{X, Y\} \) that go through \( Z \) and consider the partial cross-covariance matrix of \( \{W_1, W_2\} \) against \( \{X, Y\} \) given \( Z \). The Trek Separation Theorem for DAGs, however, says nothing explicit about conditional cross-covariances. The independence model given by conditioning on a variable is sometimes a DAG itself, but this is not true in general (Richardson and Spirtes, 2002). Sullivant et al. (2010) present versions of the theorem for some classes of DAGs under conditioning, but not in a completely general way, as further explained in Appendix A.

Fortunately, for our purposes, conditioning can be easily dealt with: let us simply introduce the conditioning variables in both sets \( A \) and \( B \). In this example, define \( A \equiv \{W_1, W_2, Z\} \) and \( B \equiv \{X, Y, Z\} \). Then it is not hard to verify that \( (\emptyset; \{U_1, Z\}) \) t-separates \( A \) from \( B \): this follows by first verifying that all treks between \( \{W_1, W_2, Z\} \) and \( \{X, Y, Z\} \) include \( U_1 \) or \( Z \).3 Cross-covariance \( \Sigma_{AB} \) is now rank-deficient since \( |\emptyset| + |\{U_1, Z\}| = 2 < 3 \), which means its determinant is zero. Because its determinant can be written as \( \sigma_{zz}(\sigma_{w1x,z}\sigma_{w2y,z} - \sigma_{w1y,z}\sigma_{w2x,z}) \) and we assume \( \sigma_{zz} \neq 0 \), the (conditional) tetrad constraint will hold. This means that without observing \( U_1 \) we can make claims about the conditional structure between \( \{W_1, W_2\} \) and \( \{X, Y\} \).

Like d-separation, it is possible that more than one conditioning set pair implies t-separation between two sets \( A \) and \( B \). In the last example, it is also the case that \( (\emptyset; \{U_2, Z\}) \) implies t-separation.

In Appendix A, we provide further examples of t-separation. Linear faithfulness, the assumption that vanishing partial correlations hold in the distribution if and only if a corresponding d-separation also holds in \( G \) (Spirtes et al., 2000), is a special case of rank faithfulness since t-separation subsumes d-separation (Sullivant et al. 2010), as also discussed in Appendix A.

3.2 Analysis of Algorithm 2

The algorithm searches for rank constraints among quartets \( \{W_i, W_j\} \times \{X, Y\} \) where \( X \) is the treatment and \( Y \) is the outcome.4 We need to characterize which structures are

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3. The “trivial trek \( Z \),” composed of the single vertex \( Z \) and no edges, does not violate this statement, as \( Z \) is t-separated from itself by \( (\emptyset; \{U_1, Z\}) \) if we define this trivial trek to be \( (\emptyset; Z) \).

4. For simplicity, we assume that there is an edge \( X \rightarrow Y \) corresponding to a non-zero coefficient \( \lambda_{yx} \). Otherwise, if \( W_i \) satisfies the Graphical Criteria for \( X \rightarrow Y \) given \( Z \), then under unmeasured confounding...
otherwise, we can use exactly one (possibly latent) element out of $Z \subseteq a$ at-separation. Moreover, since Definition 6 (left conditional choke point) holds, as $|\Sigma_{Z|Z}| > 0$ by rank-faithfulness. Hence, $\Sigma_{AB}$ is not full rank, that is, its rank is strictly less than $\min(|A|, |B|) = |Z| + 2$. □

If the constraint holds, by Theorem 3 it follows that there is some pair $(C_A; C_B)$ that $t$-separates $\{W_i, W_j\} \cup Z$ from $\{X, Y\} \cup Z$. It is required that $Z \subseteq C_A \cup C_B$, since $Z$ is contained in both $A$ and $B$. This implies the lower bound $|C_A| + |C_B| \geq |Z|$ so that at most one element of $Z$ belongs to $C_A \cap C_B$, and at most one element not in $Z$ can belong to $C_A \cup C_B$. The main result of this section is a characterization of the ways this may happen such that we fail to discover the correct causal effect.

Consider the case where there exists a trek $T$, containing no element of $Z$, which is between some $W \in \{W_i, W_j\}$ and $\{X, Y\}$. $C_A \cup C_B$ must contain an element of $T$ for the $t$-separation to hold. Since by (4) and Theorem 3 we have that $|C_A| + |C_B| \leq |Z| + 1$, then we can use exactly one (possibly latent) element out of $Z$ but in $C_A \cup C_B$ to explain this $t$-separation. Moreover, since $Z \subseteq C_A \cup C_B$, no vertex can appear in both $C_A$ and $C_B$, or otherwise $|C_A| + |C_B| > |Z| + 1$. That is, the existence of $T$ implies $C_A \cap C_B = \emptyset$.

We call any vertex in $T$ satisfying the definition below a conditional choke point. A conditional choke point can be a “left” or “right” choke point:

**Definition 6 (left conditional choke point)** In a causal DAG $G$ with vertex set $V$, $V_0$ is a left conditional choke point for pairs $\{W_i, W_j\} \times \{X, Y\}$ given $Z$, where $\{V_0, W_i, W_j, X, Y\} \cup Z \subseteq V$, if

1. some $(C_A; C_B)$ $t$-separates $\{W_i, W_j\}$ from $\{X, Y\}$ in $G$, such that $V_0 \in C_A \setminus \{Z \cup C_B\}$, $C_A \cup C_B = Z \cup \{V_0\}$, $C_A \cap C_B = \emptyset$, and

2. for any trek $(P_1; P_2)$ from a vertex in $\{W_i, W_j\}$ to a vertex in $\{X, Y\}$ in $G$ that contains no member of $Z$, $P_1$ contains $V_0$, and

3. there exists at least one such a trek.

between treatment and outcome we have that the lack of an edge $X \rightarrow Y$ implies $W_i \perp Y \mid Z$ and $W_i \not\perp Y \mid Z \cup \{X\}$. This corresponds to the collider orientation rule of the FCI algorithm (Spirtes et al. 2000) that implies $\lambda_{yx} = 0$. If on top of the lack of edge $X \rightarrow Y$ we have that there is no unblocked unmeasured confounding between $X$ and $Y$, then $X \perp Y \mid Z$, which again by faithfulness will allow us to infer $\lambda_{yx} = 0$ (Spirtes et al., 2000). Algorithm 2 is assumed to be invoked only if the FCI algorithm and the method of Entner et al. (2012) do not provide any results.
A “right conditional choke point” follows an equivalent definition with respect to $C_B$ and $P_2$. The literature has characterizations of unconditional choke points (Shafer et al., 1993; Sullivant et al., 2010), relating them to unconditional tetrad constraints. To the best of our knowledge, this is the first time that conditional choke points are explicitly defined and used. The reason for the name “choke point” is explained later in this section.

A second category of choke point, now for vertices in $Z$, will be defined as follows:

**Definition 7 (collider choke point)** In a causal DAG $G$ with vertex set $V$, $Z_0$ is a collider choke point for pairs $\{W_i, W_j\} \times \{X, Y\}$ given $Z$, where $\{W_i, W_j, X, Y\} \cup Z \subseteq V$, if

1. some $(C_A; C_B)$ $t$-separates $\{W_i, W_j\}$ from $\{X, Y\}$ in $G$, such that $Z_0 \in Z$, $C_A \cup C_B = Z$, $C_A \cap C_B = Z_0$, and

2. for any (possibly non-simple) path $P$ from a vertex in $\{W_i, W_j\}$ to a vertex in $\{X, Y\}$ in $G$ that contains a collider in $Z$, $P$ contains $Z_0$, and

3. there exists at least one such a path.

The same vertex can be a conditional and a collider choke point depending on the conditioning set. For instance, $Z_0$ is a (right) conditional choke point in Figure 3 for each $\{W_i, W_j\} \times \{X, Y\}$ given the empty set, since $\emptyset, Z_0$ $t$-separates each pair in $\{W_1, ..., W_p\} \times \{X, Y\}$ and it lies on all of such treks. The tetrad constraint $\sigma_{w_i,x} \sigma_{w_j,y} = \sigma_{w_i,y} \sigma_{w_j,x}$ holds, and we can recover the correct causal effect from $\sigma_{w_j,y} / \sigma_{w_i,x}$ as the Graphical Criteria is satisfied.

On the other hand, if we take the conditioning set to be $Z_0$ itself, then the conditional tetrad constraint $\sigma_{w_i,x,0} \sigma_{w_j,y,0} = \sigma_{w_j,y,0} \sigma_{w_i,x,0}$ also holds, as well as the corresponding $t$-separation. However, there is no conditional choke point as all treks in $\{W_i, W_j\} \times \{X, Y\}$ contain an element of the conditioning set, $Z_0$ itself. $Z_0$ is now a collider choke point, as it lies as a collider on all active paths between $\{W_1, ..., W_p\}$ and $\{X, Y\}$. Notice that the Graphical Criteria does not hold as $W_1$ is not d-separated from $Y$ given $Z_0$ in the graph that modifies Figure 3 by removing edge $X \rightarrow Y$.

Choke points do not need to be unique, in the same way that t-separation between sets may be implied by more than one conditioning pair. For instance, both $U_1$ and $U_2$ are right conditional choke points in Figure 2(c), corresponding to the different choices of $C_B$ discussed in Example 3.

We are now ready to characterize when Algorithm 2 may provide an incorrect value for the target causal effect. The proof of this result is given in Appendix B.

**Theorem 8** Let $C$ be the outcome of Algorithm 2 as applied to a distribution which is rank-faithful to some causal DAG $G$. Let $\beta_i \in C$, and let $\{W_i, W_j, Z\}$ be the respective triplet that generated $\beta_i$ by passing the check in Line 9. Then $(W_i, W_j, Z)$ fails to satisfy the Graphical Criteria if and only if exactly one of these is true:

(i) there is at least one trek between $\{W_i, W_j\}$ and $Y$ that do not contain any member of $Z \cup \{X\}$, and all of those treks intersect at the same conditional choke point $V_0$;

(ii) there is at least one (possibly non-simple) active path between $\{W_i, W_j\}$ and $Y$ that contains a collider in $Z$ and does not contain $X$, and all of those paths intersect at the same collider choke point $Z_0$. 


Moreover, it is not possible that for \((W_i, W_j, Z)\) both a conditional choke point and a collider choke point exist in \(G\) regardless of whether the Graphical Criteria hold.

That is, if a candidate triplet \((W_i, W_j, Z)\) does not correspond to valid (conditional) instrumental variables, then all offending paths that avoid \(X\) and connect the candidate instruments to \(Y\) must pass through a single variable \(V \in V \setminus \{X\}\). If \(V\) is in the conditioning set \(Z\), then this variable must be a collider in all of the offending paths. For instance, \(U\) is a conditional choke point with respect to \(Z = \emptyset\) in Figure 2(d) and \(Z_0\) is a collider choke point for \(Z = \{Z_0\}\) in Figure 3. The difference between valid and invalid instruments is that, in the former, \(X\) is a conditional choke point.

We can formalize a sufficient assumption that guarantees that Algorithm 2 does not return an incorrect causal effect as follows:

**Corollary 9** Let \(G\) be a causal DAG with observable vertices \(O\), treatment vertex \(X\) and outcome vertex \(Y\). For all \(\{W_i, W_j\} \cup Z \subset O\) where a tetrad constraint \(\sigma_{w_i x, z} \sigma_{w_j y, z} = \sigma_{w_i y, z} \sigma_{w_j x, z}\) holds, if \(X\) is a corresponding conditional choke point in \(G\), then the output set of Algorithm TETRAD-IV \(\infty\) will not contain an incorrect causal effect.

The result above allows for an empty output set \(\mathcal{C}\), which will happen if no valid pair of instrumental variables exists. However, even if all elements of \(\mathcal{C}\) are equal to some \(\lambda\), it does not mean that \(\lambda_{\mathcal{C}} = \lambda_{yx}\) in case the assumptions of Corollary 9 fail. It is impossible to distinguish, based on tetrad relations, the graph in Figure 2(b) from the one in 2(d). In the former, \(X\) can take the role of a conditional choke point explaining the tetrad constraint. In the latter, \(X\) cannot take this role because of the active trek \(W_1 \rightarrow U \rightarrow Y\), but vertex \(U\) can, resulting on a tetrad constraint but on invalid IVs. In general, we will have problems in a graph which includes an intermediate latent variable that is a common ancestor of \(X\) and \(Y\) and blocks all treks from \(W_1\) and \(W_2\) to \(X\). As another example, vertex \(U_1\) takes the role of a conditional choke point given \(Z\) in Figure 2(c), while \(X\) cannot take this role because of the active trek \(W_1 \leftarrow U_1 \rightarrow U_2 \rightarrow Y\).

One might argue that the assumptions of Corollary 9 are not too restrictive, as failing them require that all offending paths intersect at an unique vertex. If two disjoint paths from \(W_1\) to \(Y\) via latent variables exist, for instance, then no tetrad constraint will hold and at worst we would be conservative and not return any causal effect. However, we argue that these assumptions should not be taken lightly. The whole motivation for including a conditioning set \(Z\) is to block paths such that a “desirable” choke point emerges: treatment \(X\). The assumptions that make the presence of an instrumental variable possible are the same that make the presence of undesirable choke points plausible.

Instead, the output of IV-TETRAD\(\infty\) is better interpreted as an equivalence class of causal effects, differing from the usual output of causal discovery algorithms such as the PC algorithm (Spirtes et al., 2000): the PC algorithm returns a set of graphs of which some graphical features might be the same (for instance, they all share the same directed edge),
Figure 4: Vertices named $U_*$ are latent variables, with $X$ being the treatment and $Y$ being the outcome. When given as input to algorithms for finding locally covariance equivalent solutions, the result will include a (potentially different) causal effect candidate for each of the pairwise groupings $\{W_{ia}, W_{ib}\}$, $i = 1, 2, \ldots, p$, where $p = |O|/2$.

and where some causal effects might be identified based on the common features. That is, the PC algorithm either gives (in the limit of infinite data) the correct answer, or it answers “I don’t know.” The same is true of the method introduced by Entner et al. (2012).

In contrast, IV-TETRAD$_\infty$ may return an empty set of solutions, but at the same time all solutions in a non-empty output set $C$ may be wrong in case the assumptions of Corollary 9 do not hold. The need to report an equivalence class is illustrated by the following result.

**Proposition 10** There exist problems where the output of IV-TETRAD$_\infty$ will contain $O(|O|)$ different elements, none of which is the correct solution $\lambda_{yx}$.

**Proof of Proposition 10.** This is illustrated by Figure 4: given $O = \{W_{1a}, W_{1b}, W_{2a}, \ldots, W_{pb}\}$, $p = |O|/2$, the structure shown in the Figure will generate solutions for all pairs $\{W_{ia}, W_{ib}\} \times \{X, Y\}$ with corresponding conditioning set $Z_i \equiv O \setminus \{W_{ia}, W_{ib}\}$. For parameters generated randomly and independently, it is clear each resulting $\lambda_i$ will be unique. The corresponding choke point in each case is $U_i$, and none of the returned causal effects will in general be correct. □

For practical purposes, one should be willing to to assume that at least one solution exists, that is, there exists at least one pair of IVs conditioned on the same set $Z$. Under this condition, we obtain the following guarantee:

**Corollary 11** Let $C$ be the output of IV-TETRAD$_\infty$ as applied to a distribution which is rank-faithful to a causal DAG $G$ with observable vertices $O \cup \{X, Y\}$ and target edge $X \rightarrow Y$. Let $\lambda_{yx}$ be the corresponding true causal effect. Assuming that there exists at least one set
{W_i, W_j} \cup Z \subseteq O \text{ where } \{W_i, W_j\} \text{ satisfies the Graphical Criteria conditioned on } Z, \text{ then the minimum and maximum elements of } C \text{ will provide lower and upper bounds on } \lambda_{yx}.

The proof of this Corollary is immediate from the previous discussions. Summaries other than lower bounds and upper bounds are possible: for instance, if the maximum of \( C \) is much larger than the second largest element, and background knowledge says that such a large value is unlikely, we may opt to exclude it or to reweight elements of \( C \) according to some prior. More of that is discussed in Section 5.

Assuming that one solution exists is a much weaker assumption than the one in IV-BY-MAJORITY / sisVIVE. However, it does require rank-faithfulness and it can only provide bounds and other summaries of an equivalence class. If in the application of the algorithm we have that all elements of \( C \) are the same, then from the assumption that there is at least one valid tuple it follows that we found the true causal effect. A discussion of what further assumptions are necessary in practical learning is postponed to Section 4.

The soundness of Algorithm IV-TETRAD\(_\infty\) can be defined in terms of the bounds given by Corollary 11, or other summaries of interest, provably capturing the true \( \lambda_{yx} \). Soundness is however of limited value if the equivalence class is too large or trivial (that is, if there is a lack of specificity in the constraints chosen to distinguish possible models), or if the algorithm cannot find all elements of the equivalence class (a lack of completeness).

While soundness and completeness are a matter of mathematical proof, specificity must be argued in the context of the scientific plausibility of assumptions and practical matters of computational and statistical tractability.

The Markov equivalence class used by various causal discovery algorithms has a natural specificity, as independence constraints are nonparametric restrictions that require no further assumptions about the probability distribution underlying the causal model. The motivation for the specificity of our method is less clear and discussed by first formally defining the following concept:

**Definition 12 (locally covariance equivalent causal effect)** We say that two ratios \( \beta_i(z) \equiv \sigma_{w_i,y}z/\sigma_{w_i,x}z \) and \( \beta_j(z') \equiv \sigma_{w_j,y}z'/\sigma_{w_j,x}z' \) are locally covariance equivalent causal effects if there is some pair \( \{W_k, W_i\}, W_i \neq W_k, W_j \neq W_i, \) where \( \beta_i(z) = \beta_k(z) \) and \( \beta_j(z') = \beta_l(z') \).

IV-TETRAD\(_\infty\) is sound and complete in the sense that by construction it returns all and only locally covariance equivalent causal effects. The effects are “local” in the sense that it does not use covariance information from the joint distribution of \( \{W_i, W_j, W_k, W_l\} \cup Z \cup Z' \) beyond what is found in the marginals \( \{W_i, W_k\} \cup Z \) and \( \{W_j, W_l\} \cup Z' \). This is motivated by the following desiderata:

**Desiderata 1 (Local Covariance Discovery)** We say that an algorithm obeys the local covariance discovery desiderata if it follows these conditions:

1. only substructures where two or more variables follow the Graphical Criteria conditioned on a common set \( Z \subseteq O \) can be used to select conditional instrumental variables;
2. only covariance information is used.
Algorithm 3 IV-TETRAD\textsubscript{∞}

1: **Input:** set of zero-mean random variables \( O \cup \{ X, Y \} \) with known joint distribution
2: **Output:** \( C \), a set of candidate differential causal effects of \( X \) on \( Y \)
3: Initialize \( C \leftarrow \emptyset \)
4: for each pair \( \{ W_i, W_j \} \subseteq O \) do
5: for every set \( Z \subseteq O \setminus \{ W_i, W_j \} \) do
6: if \( \sigma_{wix.z} = 0 \) or \( \sigma_{wij.z} = 0 \) then
7: next
8: end if
9: if \( \sigma_{wix.z} \sigma_{wiy.z} \neq \sigma_{wiy.z} \sigma_{wix.z} \) then
10: next
11: end if
12: \( r_{Wi} \leftarrow \text{resProj}(W_i, Z \cup \{ W_j \}) \)
13: \( r_{Wj} \leftarrow \text{resProj}(W_j, Z \cup \{ W_i \}) \)
14: if \( r_{W_i} \perp \perp r_Y \text{ and } r_{W_j} \perp \perp r_Y \) then
15: \( C \leftarrow C \cup \{ \sigma_{wiy.z} / \sigma_{wix.z} \} \)
16: end if
17: end for
18: end for
19: return \( C \)

The requirement for covariance information only allows us to use only second moments, which is relatively statistically efficient. This will be relaxed in the next section. The need for “two or more” variables in condition 1 is necessary, since the Graphical Criteria provides no testable implications for a single instrument, as discussed in the introduction. It is also desirable, as it considers variables that are jointly instrumental by conditioning on the same set \( Z \). To drop this criterion is to require a search space where different conditioning sets are required for each IV, resulting in an algorithm that effectively reconstructs a “non-local” graphical substructure that will require the identification of structural coefficients other than \( \lambda_{yx} \). This is what is done, for instance, in one of the first published extensions of the Graphical Criteria (Brito and Pearl, 2002). However, identification of causal effects using given IVs is an easier problem than discovering IVs from limited structure information. Inferring structure from data is statistically challenging, and we claim that restricting a method to the desiderata above provides a baseline where stronger assumptions, and computational and statistical costs, can be exchanged for greater specificity of the solutions returned.

3.3 Exploiting Non-Gaussianity Assumptions

This Section introduces a variant of IV-TETRAD\textsubscript{∞} that can verify the validity of item 2a of the Graphical Criteria, as discussed in Section 2.3. We assume our causal model is a LiNGAM model, a linear structural equation model with independent, non-Gaussian error terms, which may include latent variables (Shimizu et al., 2006). We call this Algorithm IV-TETRAD\textsubscript{∞}, which is shown in Algorithm 3. We will also discuss the difficulties posed by 2b even under the assumption of non-Gaussianity. The motivation is to reduce the size of the
equivalence class by trading it off with the addition of more assumptions (non-Gaussianity) and a weakening of completeness (some solutions might be missed, as we will characterize).

In Algorithm 3, function \( \text{resProj}(V, S) \) is a function that returns the residual of the least-squares projection of \( V \) into the set of random variables \( S \) rearranged as a column vector,

\[
\text{resProj}(V, S) \equiv V - \mathbb{E}[S^T V] \mathbb{E}[SS^T]^{-1} S,
\]

which is itself a random variable. We assume that \( V \) and \( S \) have zero mean.

### 3.3.1 Main Result

The validity of Algorithm 3 holds “almost everywhere,” in the sense it holds for all but a (Lebesgue) measure zerosubset of the set of possible structural coefficients \( \lambda_{ij} \) and variances \( \omega_i \) of each respective error term \( e_i \). The motivation for this concept is analogous to the different variations of faithfulness, see the discussion on generic identifiability by Foygel et al. (2011) and Sullivant et al. (2010).

The main result of this section is the following:

**Theorem 13** Let \( O \cup \{Y\} \) be a subset of the variables in a zero-mean LiNGAM model with graph \( G \), where \( Y \) has no descendants in \( G \). Let \( O_i \in O \) and \( Z \subseteq O \setminus \{O_i\} \), and let \( r_i \equiv O_i - a^T Z \) be the residual of the least-squares regression of \( O_i \) on \( Z \), with \( a \) being the corresponding least-squares coefficients. Analogously, let \( r_y \equiv Y - b_i O_i - b^T Z \) be the residual of the corresponding least-squares regression. Then, almost everywhere, \( r_i \perp \perp r_y \) if and only if there are no active non-directed paths between \( O_i \) and \( Y \) conditioned on \( Z \).

The proof is given in Appendix B. This generalizes the main result of Entner et al. (2012), which considers the case where \( Z \) contains no descendant of \( O_i \).

Even under non-Gaussianity, we will see that it is still possible for conditional choke points in the true graph to add incorrect causal effect estimates to the output of Algorithm 3. However, only the following special case can occur:

**Definition 14 (downstream conditional choke point)** In a causal DAG \( G \) with vertex set \( V \), \( V_0 \) is a downstream conditional choke point for pairs \( \{W_i, W_j\} \times \{X, Y\} \) given \( Z \), where \( \{V_0, W_i, W_j, X, Y\} \cup Z \subseteq V \), if \( V_0 \) is a conditional choke point which is a descendant of \( \{W_i, W_j\} \) in \( G \), an ancestor of both \( X \) and \( Y \), such that no active back-door paths between \( V_0 \) and \( \{W_i, W_j\} \) exist in \( G \).

This implies that collider choke points are now ruled out, as shown in the following result proved in Appendix B:

**Theorem 15** Let \( C \) be the outcome of Algorithm 3 as applied to a distribution which is rank-faithful to some causal DAG \( G \). Let \( \beta_i \in C \), and let \( (W_i, W_j, Z) \) be the respective triplet that generated \( \beta_i \) by passing the check in Line 15. Then \( (W_i, W_j, Z) \) fails to satisfy the Graphical Criteria if and only if: there is at least one directed path between \( \{W_i, W_j\} \) and \( Y \) that does not contain any member of \( Z \cup \{X\} \), and all of those directed paths intersect at the same downstream conditional choke point \( V_0 \).

5. As error terms are non-Gaussian, other parameters are expected to exist, but our results are agnostic with respect to them.
By further restricting the set of possible choke points, we refine the explanation of possible disparities obtained by Algorithm 3 compared to Algorithm 2. By knowing that there is a single choke point per pair which is unconfounded with the candidate IVs, and which lies on all unblocked directed paths from \{W_i, W_j\} to Y, the output of the algorithm, combined with background knowledge, can be analyzed to exploit this piece of information. For example, any downstream conditional choke point must be a latent confounder of X and Y, which is itself caused by \{W_i, W_j\}: efforts might be directed to try to interpret and measure such confounders as possible explanations for the disparities in the output of Algorithm 3.

Examples of advantages and remaining limitations of IV-TETRAD⁺ are as follows.

Example 4 Consider Figure 2(c) again. As discussed before, t-separation between \{W_1, W_2, Z\} and \{X, Y, Z\} holds given (\emptyset; \{U_1, Z\}). However, there is an active path between \(W_1\) and \(Y\) that does not include \(X\). Algorithm 2 will include the candidate causal effect corresponding the invalid IV triplet \((W_1, W_2, Z)\) in its output. In contrast, Algorithm 3 will detect the active back-door path \(W_1 \leftarrow U_1 \rightarrow U_2 \rightarrow Y\) by the residual test of Line 15 and reject \((W_1, W_2, Z)\).

Example 5 Consider Figure 2(d) again. T-separation between \{W_1, W_2\} and \{X, Y\} holds given (\emptyset; U). However, there are active paths between \{W_1, W_2\} and \(Y\) that do not include \(X\) and intersect at the downstream conditional choke point \(U\). Algorithm 3 will not detect these offending paths and will incorrectly include a value different from \(\lambda_{yg}\) in its output set.

Example 6 Consider Figure 3 again. T-separation between \{W_1, W_2, Z_0\} and \{X, Y, Z_0\} holds given (\emptyset; \{Z_0\}). Algorithm 2 will include the candidate causal effect corresponding the invalid IV triplet \((W_i, W_j, Z_0)\) in its output. In contrast, Algorithm 3 will detect the active non-directed paths \(W_i \rightarrow Z_0 \leftarrow U \rightarrow Y\) by the residual test of Line 15.

Example 7 Figure 5 shows a more complex example where a latent variable \(U\) invalidates one particular set of candidate IVs \((\{V_1, V_2\} \text{ given } V_0)\). Meanwhile, in the same graph, conditioning on \(V_2\) is essential to validate another set of candidate IVs. These two cases cannot be distinguished by Algorithm 2 but can be by Algorithm 3.

Notice that Lines 12-14 in IV-TETRAD⁺∞ include \(W_j\) in the conditioning set for \(W_i\) and vice-versa. Including \(W_j\) in the conditioning set for \(W_i\) is acceptable, given that if \(W_i\) and \(W_j\) each satisfy the Graphical Criteria with respect to \(X \rightarrow Y\) given \(Z\), then \(W_i\) satisfies the Graphical Criteria with respect to \(X \rightarrow Y\) given \(Z \cup \{W_j\}\):

Proposition 16 If \(W_i\) and \(W_j\) each satisfy the Graphical Criteria with respect to \(X \rightarrow Y\) given \(Z\), then \(W_i\) satisfies the Graphical Criteria with respect to \(X \rightarrow Y\) given \(Z \cup \{W_j\}\).

Proof of Proposition 16. Given \(Z\), both \(W_i\) and \(W_j\) are d-separated from \(Y\) in the graph where edge \(X \rightarrow Y\) has been removed. If \(Z \cup \{W_j\}\) d-connected \(W_i\) to \(Y\) in this mutilated graph, then there should be an active path from \(W_i\) to \(Y\) where \(W_j\) is a collider. But this would imply an active path from \(W_j\) to \(Y\) given \(Z\), contrary to our assumption. \(\square\)
Figure 5: In this system, $W \equiv \{V_1, V_2\}$ satisfies the conditional tetrad constraints $\{V_1, V_2\} \times \{X, Y\}$ given $Z \equiv \{V_3, V_6\}$. The same is true for $W' \equiv \{V_3, V_4\}$ for $Z' \equiv \{V_2, V_5\}$. Only $W'$ represents a set of (conditionally) valid IVs. The non-Gaussianity criteria of $\text{IV-TETRAD}_\infty^+$ helps to rule out the former.

The motivation for including $W_j$ in the conditioning set of $W_i$ is that, if $W_i$ and $W_j$ share a hidden common cause that is independent of $X$ and $Y$ and $W_j$ is an unblocked ancestor of $X$, then not conditioning on $W_j$ will induce an active back-door between $W_i$ and $X$ that will make the algorithm unnecessarily remove $\sigma_{w_i|x,z}/\sigma_{w_i|x,z}$ from the output.

3.3.2 Implications to Completeness

The following result is a direct consequence of the assumptions and the exhaustive search done by $\text{IV-TETRAD}_\infty^+$. The proof of this result follows immediately from Theorem 13 and the existence of edge $X \rightarrow Y$:

**Proposition 17** If there is a pair of observable variables $\{W_i, W_j\}$ which are IVs conditioned on some $Z$ according to the Graphical Criteria, and some $W \in \{W_i, W_j\}$ has an active non-directed path with $X$ given $Z$, and $X$ is a parent of $Y$, then Algorithm 3 will not include $\{W_i, W_j\}$ in its output.

The implication of this is as follows. A positive property of $\text{IV-TETRAD}_\infty^+$ is that it will exclude from the equivalence class some candidate IVs that do not obey the Graphical Criteria but which are returned by $\text{IV-TETRAD}_\infty$. However, $\text{IV-TETRAD}_\infty$ was complete in the sense of always including in its output the correct causal effect, under the assumption that at least one pair of valid IVs existed given a common conditioning set. This assumption is not strong enough in the case of $\text{IV-TETRAD}_\infty^+$. As an example, suppose that there exists only one pair $\{W_i, W_j\}$ of valid IVs, and its conditioning set $Z$ does not block all back-doors into $X$. $\text{IV-TETRAD}_\infty^+$ may return a non-empty set but it will not include the correct causal effect.

**Example 8** $W$ is a valid IV in the model with paths $W \rightarrow X \rightarrow Y$, $W \leftarrow U_1 \rightarrow X \leftarrow U_2 \rightarrow Y$, but $W$ will be discarded due to the back-door path between $W$ and $Y$ that is unblocked by not conditioning on $X$. $\square$

As a result, the following is a modification of Corollary 11 to the alternative algorithm $\text{IV-TETRAD}_\infty^+$. The proof follows immediately from the previous results.
Algorithm 4 IV-TETRAD$^+$

1: **Input:** $D$, a sample from the joint distribution of random variables $O \cup \{X,Y\}$; $K$, number of instruments to consider
2: **Output:** two sets of estimates of the differential causal effect of $X$ on $Y$
3: Remove from $O$ any element that is not in the Markov blanket of either $X$ or $Y$, and also any element which is marginally uncorrelated with $X$ or $Y$
4: Let $\hat{\lambda}_i \leftarrow \hat{\sigma}_{wy,oi}/\hat{\sigma}_{wx,oi}$ for each $W_i \in O$
5: Initialize $C$ and $C_R$ as empty sets, and $W$ as $O$
6: while $|W| \geq K$ do
7:   Sort $W$ as $W^{(1)}, \ldots, W^{(|W|)}$ according to the corresponding $\{\hat{\lambda}\}$
8:   Find $j \in \{1, 2, \ldots, |W|\}$ that minimizes $\text{ScoreTetrads}(W^{(j+K-1)}, O, X, Y, D)$
9:   $W_j \leftarrow \text{ExpandTetrads}(j, K, W, O, X, Y, D)$
10:   if $\text{TestTetrads}(W_j, O, X, Y, D)$ then
11:      $\hat{\lambda}_{TSLS} \leftarrow \text{TSLS}(W_j, O \setminus W_j, X, Y, D)$
12:      if $\text{TestResiduals}(W_j, O, X, Y, D)$ then
13:         $C \leftarrow C \cup \hat{\lambda}_{TSLS}$
14:      else
15:         $C_R \leftarrow C_R \cup \hat{\lambda}_{TSLS}$
16:      end if
17:   end if
18: $W \leftarrow W \setminus W_j$
19: end while
20: return $(C, C_R)$

Corollary 18 Let $C$ be the output of IV-TETRAD$^+_{\infty}$ as applied to a distribution which is rank-faithful to a causal DAG $G$ with observable vertices $O \cup \{X,Y\}$ and target edge $X \rightarrow Y$. Let $\lambda_{yx}$ be the corresponding true causal effect. Assuming that there exists at least one set $\{W_i, W_j\} \cup Z \subseteq O$ where $\{W_i, W_j\}$ satisfies the Graphical Criteria conditioned on $Z$, and that no non-directed path between $\{W_i, W_j\}$ and $X$ is active given $Z$, then the minimum and maximum elements of $C$ will provide lower and upper bounds on $\lambda_{yx}$.

The trade-off between IV-TETRAD$^+_{\infty}$ and IV-TETRAD$^+$ can be implemented in different ways within a practical algorithm, as we shall see next.

4. Practical Learning from Data

Testing tetrad constraints from data is difficult in practice, particularly without assuming Gaussianity and under conditioning. In order to search for candidate IVs, we will first adopt the stronger (but falsifiable) assumption that, if any set of valid conditional IVs exist, then there is at least one such set $W$ of size $K$, which remains valid by conditioning on $O \setminus W$. One motivation is to avoid a combinatorial search for conditioning sets, while still having the option of rejecting a solution if confounding or collider biases remain by doing a test at the end. The other motivation is the statistical unreliability of candidate sets of small size: in a large system where the treatment may have many observed causes, instruments will in
general be weakly associated with the treatment, leading to high variance estimates. This issue of “weak instruments” is pervasive in real problems and one mitigation is to consider instrument sets of a minimum size.

The algorithm, IV-TETRAD$^+$, is shown in Algorithm 4. Its implementation diverges from a direct adaptation of IV-TETRAD$^{∞}$+ where hypothesis testing would be used to check the required constraints. This is because we found that, in simulations, directly testing conditional tetrad constraints is not a very robust strategy. Instead, IV-TETRAD$^+$ includes intermediate steps along the lines of a “score-based approach” (Cooper, 1999) for causal discovery. In the algorithm, we search for a window of $K$ estimated causal effects $\{\lambda\}$ within a sorted list, scoring this set according to a particular measure of variability of its elements. Ideally, a perfect score of 0 is obtained when all $\lambda$ in this window are equal. This corresponds to an alternative view of Line 9 of IV-TETRAD$^{∞}$+, where we find a set of $K$ variables $W_i$ with corresponding ratios $\hat{\sigma}_{w_i.y,o_i}/\hat{\sigma}_{w_i.z,o_i}$ as close as possible to each other, in a formal sense defined below.

The algorithm makes use of $TSLS(W, Z, X, Y) \quad \Delta \quad \sum ^{\infty}$, the two-stage least squares (TSLS) estimator for some set of conditional instruments $W$ and conditioning set $Z$, given a data set $D$. The TSLS estimator$^6 \lambda_{TSLS}$ is motivated by the redundant ways of estimating $\lambda_{g}$ if we have more than one instrument. For instance, we could have used either $\sigma_{w_i,y,z}/\sigma_{w_i,x,z}$ or $\sigma_{w_j,y,z}/\sigma_{w_j,x,z}$ in the previous algorithms without any difference. When finite sample covariance matrices are used instead of population covariance matrices, differences will appear. The TSLS is a classic way of combining all postulated instruments into a single estimator (Angrist and Pischke, 2009). This is particularly relevant in IV-TETRAD$^+$, where we use more than two instruments in each causal effect estimate inside function SCORETETRADs, as explained in the next paragraphs.

The output of the algorithm is a pair of two sets, $C$ and $C_R$, where the latter excludes from $C$ all elements that correspond to each $W_i$ detected to be connected to $Y$ via an active non-directed path, as in IV-TETRAD$^{∞}$+.

4.1 Interpretation

The details of the algorithm are as follows:

- In Line 3, we remove variables from $O$ that are considered to be conditionally independent of $X$ or $Y$. The deletion of variables outside the Markov blanket of $X$ correspond to Line 6 in IV-TETRAD$^{∞}$+. As discussed in the beginning of this section, we assume that if valid IVs exist, then $O$ can be partitioned into instrument set $W$ and conditioning set $Z \equiv O \backslash W$. By Proposition 16, $W_i \in W$ remains a valid IV if the conditioning set is $O_i \equiv O \backslash \{W_i\}$. We remove also variables which are statistically independent of $Y$ given the remaining variables in $O$. This is motivated by the instability of the ratio of estimates $\hat{\sigma}_{w_i,y,z}/\hat{\sigma}_{w_i,x,z}$ when both are close to zero. If $W$ is d-connected to $X$ and edge $X \rightarrow Y$ exists, then in theory we will not remove any valid IV by doing this test.

---

6. Assuming $D$ is centered to have zero empirical mean, this is defined as follows: let $R_W \quad \Delta \quad R_X \quad \Delta \quad R_Y$ denote the corresponding empirical residuals of the least-squares regression of $W, X$ and $Y$ on $Z$. The projection matrix $P_W$ is defined as $R_W(R_W^T R_W)^{-1} R_W$. Let $R_X \quad \Delta \quad F_W R_X$. The two-stage least squares coefficient is then given by the least-squares regression of $R_Y$ on $R_X$, resulting in $\lambda_{TSLS} \equiv (R_X^T P_W R_X)^{-1} R_X^T P_W R_Y$. 

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In our experiments in the next Section, Line 3 is implemented by running LARS (Efron et al., 2004), regressing separately X on O and Y on O. We then remove from O any element that is given a coefficient of value zero in either regression.

Moreover, we also remove from O any element marginally uncorrelated with X or Y by a test of vanishing Spearman correlations at a 0.05 level 7

- In Line 4, each \( \hat{\sigma}_{ab,z} \) is an estimated partial covariance as given by the empirical covariance matrix derived from \( D \), and \( O_{vi} \equiv O\{W_i\} \).

- The loop that starts in Line 6 will search for a subset of \( W \) of size at least \( K \) that respects desirable tetrad and residual independence constraints. Every time a set of minimum size is found, it is removed from \( W \) and the process repeats;

- In Line 8, function \( \text{SCORETETRADS}(W^{(j:j+K-1)}, V, X, Y, D) \) returns the median of the absolute differences between \( \hat{\lambda}_{\text{TSL}} \) and each \( \hat{\lambda}^{(j)} \), \( \hat{\lambda}^{(j+K-1)} \), where \( \hat{\lambda}_{\text{TSL}} \equiv \text{TSL} \{(W^{(j)}, \ldots, W^{(j+K-1)}), V\{W^{(j)}, \ldots, W^{(j+K-1)}\}, X, Y, D\} \);

- In Line 9, function \( \text{EXPANDTETRADS} \) performs a greedy search to optimize \( \text{SCORETETRADS} \) by expanding the current set \( \{W^{(j)}, \ldots, W^{(j+K-1)}\} \) as either \( \{W^{(j-1)}, \ldots, W^{(j+K-1)}\} \) or \( \{W^{(j)}, \ldots, W^{(j+K)}\} \). This is done until a local maximum for \( \text{SCORETETRADS} \) is found;

- At the end of Line 9, we have a set of proposed instrumental variables which are found to approximately imply the same causal effect, as defined by Line 4. This also implies approximately satisfying the corresponding tetrad constraints. In Line 10, we formally check the validity of this candidate set by performing Wishart tests of conditional tetrad constraints (Wishart, 1928; Spirtes et al., 2000) for every possible pair in \( W_j \) against pair \( \{X, Y\} \). 8 Function \( \text{TESTTETRADS} \) returns true if and only if more than half of the tested pairs return a test p-value greater than a particular pre-defined level. In our experiments in the next section, this level was set at 0.05. This step has therefore a direct correspondence to the tetrad checks in IV-TETRAD\(_{\infty}^+\), augmented with a pre-selection of pairs;

- Finally, in Line 12, for each \( W \in W_j \), function \( \text{TESTRESIDUALS} \) is defined such that we first obtain the residual \( r_W \) of the least-squares regression of \( W \) on \( O\{W\} \), as well as the residual \( r_X \) of the regression of \( X \) on \( O \) (alternatively we could calculate \( r_Y \), but we expect \( X \) to be more strongly associated with \( W \) as in theory there is no edge between \( W \) and \( Y \) ). We perform a marginal independence test of these two residuals using HSIC (Gretton et al., 2007). As with the tetrad tests, \( \text{TESTRESIDUALS} \) will return true if and only if more than half of the p-values are above a particular threshold (again 0.05 in our experiments).

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7. Although this rules out possible IVs which are dependent on \( X \) only if we condition in some other observed variable (for instance, if the only path between \( W_i \) and \( X \) is \( W_i \rightarrow Z \leftarrow U_z \rightarrow X \) for some \( Z \in Z \) and \( U_z \in U \)), in preliminary runs we found this to provide more stable outputs.

8. We ignore that the test distribution in Wishart’s test assumes Gaussianity and marginal covariances. As discussed by Spirtes (2013), there is evidence that tests without Gaussianity assumptions (Bollen, 1990) have no clear advantage compared to Wishart’s.
The rationale for Algorithm 4 is as follows. In the limit of infinite data, any set of $K$ variables which satisfies the corresponding tetrad constraints will return a value of zero for $\text{ScoreTetrads}$, which will then be expanded by $\text{ExpandTetrads}$ to include any remaining candidate IVs that imply the same ratio $\lambda$. Sets of size $K$ that satisfy the Graphical Criteria will contribute with the true causal effect to $C$, as will also sets of size $K$ or larger corresponding to locally covariance equivalent causal effects. The algorithm can fail to produce information that a more exhaustive algorithm would produce if subsets of conditioning sets $O \setminus W_i$ were considered instead (for instance, in the limit of infinite data, IV-TETRAD$^+$ will return an empty $C_R$ for data generated by the model in Figure 3, which is in theory avoidable).

The separation between $C$ and $C_R$ is due to the fact that full automation of effect discovery is not recommended, as there are unidentifiability trade-offs. Recall that in the graph where $W_1$ and $W_2$ are valid IVs, but where $W_1$ and $X$ share an unmeasured confounder, that the residual test will asymptotically reject this candidate pair. It is to remind the user of this possibility that the algorithm separates, in a second set, the outcomes filtered by a more stringent criteria. The price to be paid is the stronger assumption concerning the existence of a partition $(W, Z)$ of $O$ such that $W$ are valid IVs given $Z$ and that $|W| \geq K$. If this existence criteria also includes that no elements in $W$ has an active non-directed path with $Y$ given $Z$, then $C_R$ will be no larger and sometimes smaller than $C$. Both sets can be summarized in many ways, including using them to provide lower and upper bounds to $\lambda_{yx}$ in a way analogous to Corollaries 11 and 18.

4.2 Further Comments

IV-TETRAD$^+$ can in principle be expanded to include a greedy search over sets $Z$ which are strict subsets of $O \setminus \{W_i\}$. Some non-trivial theory on how to do it with formal guarantees and non-trivial assumptions is left as future work. The most immediate modification of IV-TETRAD$^+$ would be to include a further outer loop, where other conditioning sets are used instead of $O \setminus W_i$, starting from the empty set. The outer loop would increase this set by one element at a time, optimizing $\text{ScoreTetrads}$. As this is a relatively straightforward heuristic without any insightful theory behind it, and the statistical variance of the final result will just increase considerably, we do not discuss it further. Our preferred solution is to let the algorithm return an empty set if no tetrad constraints and residual tests are satisfied, ignoring the search for a more specific conditioning set. The theory behind IV-TETRAD$^+_\infty$ is used in IV-TETRAD$^+$ to provide guarantees about $C_R$. It also provides the explanation of what can go wrong by improperly conditioning on a covariate. This theory can be combined with background knowledge to manually select a pre-defined set of conditioning sets that are believed to be adequate for finding instruments, in case the default $O \setminus W_i$ is deemed insufficient.

The choice of $K$ is not a straightforward question. This decision is a bias-variance trade-off: smaller $K$ means weaker assumptions but higher variance. We do not provide a domain-free rule of thumb: our recommendation is to set $K$ as high as possible according to knowledge from the domain of interest. Notice that, unlike sisVIVE, this number does not need to grow with $|O|$. Sensitivity analysis, where different values of $K$ are attempted, can be a valuable tool to diagnose the impact of this choice.
5. Experiments

We assess how IV-TETRAD$^+$ compares to other methods in a series of simulations. We then provide an example on how the method can be used in an empirical study and assess its robustness.

5.1 Synthetic Studies for Finding a Single Causal Effect Assuming IVs Exist

We generate synthetic models of particular structures, intended to capture features such as confounding and collider bias that should be avoided by proper conditioning. In order to allow for a direct comparison of IV-TETRAD$^+$ against alternative methods that return a single causal effect estimate, we pragmatically force our algorithm to choose a single element in each output $C \cup C_R$ for the sake of simplifying comparisons.\(^9\) The choice is the effect estimate corresponding to the set $W_j$ that minimizes $\text{ScoreTetrads}$ at the end of Line 9, which is then contrasted against the true effect.\(^10\)

Simulations are performed as follows. Treatment $X$ and outcome $Y$ have two latent parents, $U_1$ and $U_2$. $U_1$ is connected to observed covariates as described below, so its variance decreases by conditioning. $U_2$ is not directly connected to the observed covariates. We then generate synthetic graphs by splitting the observed covariates in different groups.

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\(^9\) Also, in this setup, the non-Gaussianity test of Line 12 does not play a role in the final output.

\(^10\) It is theoretically possible that this criterion will make IV-TETRAD$^+$ return the wrong causal effect in the limit of infinite data, as by design we introduce locally covariance equivalent causal effects. However, with finite data this criterion typically chooses solutions corresponding to correct IVs: in our setup, the non-IV vertices that imply locally covariance equivalent effects tend to be more weakly associated with outcome $Y$ than the true IVs, as paths from such vertices to outcome $Y$ pass through latent choke points. This inflates the variance of the corresponding $\hat{\lambda}_i$, leading to higher disparity compared to the effect estimated by two-stage least squares.
### Table 2: Experimental results for 8 methods as described in the text, including three oracle competitors that illustrate idealized scenarios where parts of the structure are known. Two experimental parameters are varied: sample size $N$ and the amount $\lambda_{xu} = \lambda_{yu}$ of unmeasured confounding between $X$ and $Y$. For each pair of method and experimental condition, we report the summarized performance over 100 trials in two ways: the median absolute difference between the estimated causal effect and true $\lambda_{yx}$ (the smaller, the better); and the proportion in which the sign of the estimated effect and the sign of the true effect agree (the higher, the better). IV-TETRAD$^+$ beats all non-oracle competitors at sample size 10000, with a $p$-value $p < 0.001$ according to a binomial test (in bold). Notice that NAIVE2 is actually the closer competitor to IV-TETRAD$^+$, while in this particular setup sISVIVE does not improve even with large sample sizes.

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<th>$\lambda_{xu} = \lambda_{yu} = 0.25$</th>
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</tbody>
</table>

Group $W$ are variables which can be used as conditional IVs. This group has two subgroups, $W_{\text{best}}$ and $W_{\text{possible}}$. The former remains a set of valid IVs conditioned on all other observed covariates. The latter will have active collider paths with treatment $X$ that can be in principle deactivated by a combinatorial search, which is not performed by any algorithm in our benchmark. Group $Z_S$ is a set of covariates with a common latent child $U_S$ which is also a parent of $X$ and $Y$, making $U_S$ a spurious choke point between $Z_S$ and $\{X, Y\}$. Group $Z_D$ is a set of covariates which are parents of both $X$ and $Y$ and as such are not blocked by choke points, latent or not. Group $Z_C$ are variables which are children of $W_{\text{possible}}$ and $U_1$. Finally, all variables in $WZ \equiv W \cup Z_S \cup Z_D$ have a common latent parent $U_0$, making conditioning on $WZ \setminus W$ required for $W \in W_{\text{best}}$ to be a valid IV. Figure 6 shows an example where $|W_{\text{best}}| = 2$, $|W_{\text{possible}}| = 1$, $|Z_S| = 2$, $|Z_D| = 1$ and $|Z_C| = 1$.

The methods we compare against are:

1. NAIVE1, the least-squares regression coefficient of $Y$ on $X$;
2. NAIVE2, two-stage least squares (TSLS) of $Y$ on $X$ using all variables $O = W \cup Z_S \cup Z_C \cup Z_D$ as instruments;
3. NAIVE3, least-squares regression of $Y$ on $X$ and $O$;
4. ORACLE, TSLS estimation using the correct set of IVs $W$ and correct adjustment set $Z_S \cup Z_D$;

5. W-ORACLE, TSLS using $W$ as IVs, but conditioning on all of the other variables $Z_S \cup Z_D \cup Z_C$;

6. S-ORACLE, the algorithm of Kang et al. (2016) executed by first correctly removing the set $Z_C$ from the input;

7. SISVIVE, the algorithm of Kang et al. (2016) executed in the standard way, by taking all variables $O$ as input;

8. IV-TETRAD+, our method, using all variables $O$ as input and $K = 10$.

All error variables and latent variables are zero-mean Laplacian distributed, and coefficients are sampled from standard Gaussians and re-scaled such that the observed variables have a variance of 1. A simulation is rejected until $|\lambda_{yx}| > 0.05$. Coefficients between $\{X, Y\}$ and latent variables $\{U_1, U_2\}$ are set such that $\lambda_{xu_1} = \lambda_{yu_1} = \lambda_{yu_2} = \lambda_{yu_2}$, at two levels, $(0.125, 0.25)$. The difficulty of the problem increases with $\lambda_{xu_1}$, as this makes unmeasured confounding stronger. Comparisons are shown in Table 2, with the setup $|W_{\text{best}}| = 15$, $|W_{\text{possible}}| = 10$, $|Z_S| = 10$, $|Z_D| = 5$, $|Z_C| = 10$. This satisfies the criterion of $|W|$ being half the number of remaining variables, although only the 15 variables $Z_S \cup Z_D$ should be used as a conditioning set.

We generate 100 synthetic problems with a data set of 10,000 points each and assess methods using also subsets of size 1000 and 5000. For $\lambda_{xu} = 0.125$, the empirical distribution of effects $\lambda_{yx}$ had a median of 0.34, an upper quartile of 0.55 and a maximum of 1.05. For $\lambda_{yu} = 0.25$, a median of 0.35, an upper quartile of 0.75 and a maximum of 1.28. For assessment, we use two measures that summarize absolute deviance from $\lambda_{yx}$ and agreement with its sign.

The message in Table 2 is as follows. The parameter estimation problem is easy if one knows the correct IVs and background variables (ORACLE), but it shows already its challenges if one has potentially correct IVs but conditions on the wrong set (W-ORACLE). SISVIVE can work very well if one knows in advance which conditioning variables to discard (S-ORACLE). In general we will not have this knowledge, and SISVIVE can potentially behave badly, even worse than a naïve method (compare NAIVE3 to SISVIVE). The large contrast between S-ORACLE and SISVIVE is a warning against ignoring the effects of incorrect conditioning. Of all non-oracle methods, IV-TETRAD+ is the clear winner, although as expected it is a high variance estimator and should be advantageous only for relatively large sample sizes. In Section 5.3 we illustrate ways in which the uncertainty of this estimator can be assessed.

5.2 Residual Test Assessment

We assess the robustness of the residual test validations of IV-TETRAD+ in a series of simulations. We designed three scenarios as illustrated in Figure 7, where $\{U_0, U\}$ are latent variables:
Figure 7: The three different configurations used in the assessment performed in Section 5.2.

$S_1$: The graph is given by Figure 7(a). The ideal output is to falsify the null hypothesis of residual independence between $W_1$ and $X$ given $\{W_2, Z\}$, and $W_2$ and $X$ given $\{W_1, Z\}$.

$S_2$: The graph is given by Figure 7(b). The ideal output is to falsify the null hypothesis of residual independence between $W_1$ and $X$ given $W_2$, and $W_2$ and $X$ given $W_1$.

$S_3$: The graph is given by Figure 7(c). The ideal output is not to falsify the null hypothesis of residual independence between $W_1$ and $X$ given $W_2$, and $W_2$ and $X$ given $W_1$.

In each scenario, we simulate 100 synthetic problems and generate a data set of size 10,000 from each. The simulation of parameters is similar to Section 5.1, also including two levels of confounding by setting $\lambda_{xu} = \lambda_{yu}$ to either 0.125 or 0.25. We also increase the non-Gaussianity of the error terms: each error term is generated by sampling a Laplace distributed random variable as before, and then raising it to the power of 1.5, preserving the sign, and rescaling it back to have the same variance required so that each latent and observed variable has a marginal variance of 1. We do not claim that the tests proposed are particularly strong when the sample size is not large or the random variables are not clearly non-Gaussian. We do illustrate that, according to the theory, we should get better decisions with increasing sample sizes.

A joint decision is given for the pair of variables $\{W_1, W_2\}$, as in the definition of TESTRESIDUALS used in Algorithm 4. This is done by voting: for each candidate IV $W \in \{W_1, W_2\}$, we test the independence of $r_W$ against $r_X$ using HSIC (Gretton et al., 2007) conditioning on the remaining observed variables. If more than half of the candidate IVs result in a corresponding p-value greater than the chosen level of 0.05 (in our setup, this means that both $W_1$ and $W_2$ pass the test), then the model passes the test and the IVs are not falsified. Otherwise, the model is rejected. In our simulation, we are correct in each simulated instance if we reject the candidate instrument pair in scenarios $S_1$ and $S_2$, or if we do not reject it in scenario $S_3$.

We use a Shapiro-Wilk test of Gaussianity followed by the HSIC test as implemented in the R package dHSIC\textsuperscript{11} both at a level 0.05. That is, if for candidate instrument $W$ the Gaussianity assumption cannot be rejected for the residuals $r_W$ or $r_X$, then $W$ by default is

\textsuperscript{11} Available at the Comprehensive R Archive Network, CRAN, https://cran.r-project.org/
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Table 3: Experimental results for assessing the detection of invalid instrument candidates based on the dependencies between least-squares residuals of instruments and treatments on a common pool of covariates. Within each sample size and amount of confounding, we show the proportion of times we correctly decided whether the corresponding pair of candidate IVs are indeed valid or not (out of 100 simulations). The three problems $S_1$, $S_2$ and $S_3$ are discussed in the main text. The correct decisions are to reject the candidates in problems of type $S_1$ and $S_2$, and not reject them in $S_3$.

<table>
<thead>
<tr>
<th>Confounding</th>
<th>$N = 1000$</th>
<th>$N = 5000$</th>
<th>$N = 10000$</th>
</tr>
</thead>
<tbody>
<tr>
<td></td>
<td>$S_1$</td>
<td>$S_2$</td>
<td>$S_3$</td>
</tr>
<tr>
<td>$\lambda_{yu} = 0.125$</td>
<td>0.54</td>
<td>0.28</td>
<td>0.87</td>
</tr>
<tr>
<td>$\lambda_{yu} = 0.25$</td>
<td>0.63</td>
<td>0.59</td>
<td>0.79</td>
</tr>
</tbody>
</table>

As expected, correctly rejecting the model in scenarios $S_1$ and $S_2$ is easier when confounding is stronger, which is precisely when we would like to be more conservative regarding our choice of instruments. Scenario $S_2$ is particularly hard under smaller sample sizes, as correlations quickly go to zero given a single path of length 3 between candidate instruments and treatment/outcome. The performance in Scenario $S_3$ plateaus, meaning that in this particular simulated distribution of parameters we still erroneously reject the valid IVs approximately 17% of the time. The Type I error probability is not that straightforward to calculate due to the voting mechanism and the dependency of the least-squares rejection estimates and the residual estimates. However, as it is known, in any constraint-based search algorithm we should decrease the test level hyperparameter from 0.05 towards zero as sample size increases if we want consistency. This is analogous to the behavior of other algorithms such as the PC algorithm (Spirtes et al., 2000).

To conclude, despite the noted shortcomings of requiring detectable non-Gaussianity and large sample sizes, the tests do provide a validation of proposed instruments as advertised. Moreover, they are expected to be more effective in discarding possibly invalid IVs as unmeasured confounding gets stronger, which is a desired behavior. However, we still recommend that we output candidate effects both with and without the residual independence tests as done by IV-TETRAD$^+$, for the reasons explained.

\footnote{In theory the goal is to save computational time. In this cases the independence assumption is likely not to be rejected anyway, but the HSIC test can be costly if done many times so we want to avoid performing it whenever possible.}
5.3 Empirical Illustration

We consider the application of our method to the study by Sachs et al. (2005). That study collected cell activity measurements (concentration of proteins and lipids) for single cell data under a variety of conditions. While searching for instrumental variables under a combination of experimental conditions is an interesting topic, we will focus on a single condition (described in the paper as stimulation with anti-CD3 and anti-CD28). There are 11 variables and 853 data points, of which we selected the manipulation of concentration levels of molecule Erk as the treatment, and concentration of Akt as the outcome. We use as background knowledge the model inference result shown as Figure 3A presented by Sachs et al. (2005), encoding that the other 9 cell products are not causally affected by either Erk or Akt. Although previously unknown, evidence for this causal link was also given support by an experiment performed by Sachs et al. (2005).

The data shows some weak correlations, but we will assume for simplicity that no conditional independencies will hold between the treatment variable and the remaining 10 variables. The motivation for Step 3 in Algorithm 4 was primary as a variance reduction technique. Instead, here we will be primarily concerned about illustrating how we could assert uncertainty in our estimates using the modified Bayesian approach introduced by Silva and Kalaitzis (2015).

Running the standard regression adjustment NAIVE3 on this data, we get a differential causal effect of 1.36. sisVIVE reports an effect of 1.58, where all variables were chosen as instruments and as such this estimate here is the same as NAIVE2. This suggests that unmeasured confounding may be weakening the association between treatment and outcome, but sampling variability may be high as discussed in the next section. We run IV-TETRAD$^+$ as in Section 5.1 (and skipping Step 3), where a single effect of 1.43 is returned by minimizing SCORETetrads. For this run we chose a window size $K = 3$, as $K = 2$ is in general too noisy of a choice since it cannot exploit any redundancies of the tetrad constraints implied by the candidate instrument sets. TESTResiduals does not reject any candidate set proposed by the tetrad search. For $K = 4$ and $K > 4$ the results were 1.38 and 1.58 respectively, illustrating some of the difficulties of working with a relatively small sample size (853 points) using an algorithm that has been shown to require large samples.

To illustrate what a practitioner should do in this case, we describe ways of assessing the uncertainty of the output of IV-TETRAD$^+$.

5.3.1 Exploration with the Bootstrap

First, we show a simple comparison based on the bootstrap: we sample the data with replacement and show the corresponding (bootstrapped) sampling distribution of the estimates obtained by NAIVE3, sisVIVE and the single-output IV-TETRAD$^+$ with $K = 3$. A smoothed estimate of the respective bootstrapped outcomes is shown in Figure 8. It is evident that as the method gets more flexible, the entropy of the respective sampling distribution of the estimates also increases. There is some evidence of bimodality of the sampling distribution, with IV-TETRAD$^+$ being particularly more extreme.

Bootstrap can in principle be used to find confidence intervals for the corresponding causal effects. However, it is relevant also to discuss Bayesian approaches for quantifying uncertainty. This is particularly interesting for a method like IV-TETRAD$^+$, which po-
tentially returns sets of causal effects of varying sizes. It is not clear what a confidence interval (or confidence set) would be here. Returning a summary such as the minimum and maximum elements of the set, and building confidence intervals on bounds of the causal effect, may be too conservative and not very robust to outliers. In the next section, we will show a Bayesian solution to that problem. First, let us start by still considering the case where IV-TETRAD+ returns a single element as chosen by \textsc{ScoreTetrads}.

Obtaining a posterior distribution over the causal effect of interest requires a likelihood function that includes all observed variables. The usual Bayesian approach for graphical modeling would require exploring a space of conditional tetrad and residual independence constraints, where any chosen set of constraints defines a model. We would then evaluate the marginal likelihood of each model and search for the maximum a posteriori or list several models of high probability. Alternatively, we in principle could perform Markov chain Monte Carlo (MCMC) sampling in the space of parameters and model constraints. A very different algorithm from IV-TETRAD+ would be necessary, as it does not search among complete models. Such an algorithm would be computationally very intense, and it is not clear how to parameterize a joint likelihood for any given set of constraints, whether latent variables are explicitly introduced or not.

5.3.2 Learning by Bayesian Projections

Instead, we adopt the much simpler “Bayesian projections” approach suggested by Silva and Kalaitzis (2015), which has a relation to several other approaches for model selection that avoid introducing constraints directly in the likelihood function, such as the one by Goutis and Robert (1998). We define a black-box non-Gaussian likelihood function for the observed variables, generate posterior samples from it by MCMC, and perform a causal effect search algorithm similar to IV-TETRAD+. The main difference here is that instead of performing statistical tests, we instead reject constraints if they violate a threshold of misfit.

To make it more explicit, we define an algorithm which we will call IV-TETRAD++. In this algorithm, we are given $M$ posterior samples of the model parameters from a black-box
mixture of Gaussians likelihood. Within each sample $1 \leq m \leq M$, we run IV-TETRAD$^+$ with the following modifications:

1. Empirical covariances $\hat{\sigma}_{ab,s}$ are replaced with the corresponding model covariances $\sigma^{(m)}_{ab,s}$, which can be computed analytically for a mixture of Gaussians model;

2. We ignore Line 3 of IV-TETRAD$^+$, the one where we originally estimated Markov blankets. Instead, we discard every $\lambda_i$ such that $|\lambda_i| > \alpha_{cut}$. Hyperparameter $\alpha_{cut}$ is given as an input, implying the removal from the output of unusually high causal effects, including some of those resulting from variables $V_i$ such that $\sigma_{v_i \cdot v_i} \approx 0$.

3. TestTetrads is modified to reject a tetrad constraint if

$$ |\rho^{(m)}_{w_i x, o_i, j} \rho^{(m)}_{w_j y, o_i, j} - \rho^{(m)}_{w_j x, o_i, j} \rho^{(m)}_{w_j y, o_j, j}| > \tau, $$

where $\rho_{ab,s}$ is the partial correlation of $A$ and $B$ given $S$. Hyperparameter $\tau$ is given as an input, and provides the rule that decides on whether a tetrad constraint fits the model given by the $m$-th MCMC sample with model partial correlations $\rho^{(m)}_{ab,s}$.

4. For simplicity, we do not perform TestResiduals. It could in principle be implemented by sampling a large synthetic data set from model $m$, estimating a measure of dependence between the residuals, and rejecting independence with this measure is larger than a given hyperparameter. Due to its computational cost and the fact it does not change the message of this section, we will ignore this step.

A more formal description of the algorithm is given in Appendix C.

Silva and Kalaitzis (2015) discuss the shortcomings of algorithms based on Bayesian projections compared to a traditional Bayesian approach where the likelihood enforces constraints explicitly. However, the computational advantages are major, as sampling from the posterior of a more standard black-box model is much simpler than searching within a complex likelihood space. As a matter of fact, at the present time it is not known how to define such a likelihood function. A Bayesian projection approach has some similarities to frequentist bootstrap, where a black-box object (the empirical distribution, in the nonparametric bootstrap case) provides the source of variability (the sampling distribution, in the bootstrap case). This allows us to leverage constraint-based algorithms with few modifications.

The most delicate design choices are the choice of $\alpha_{cut}$ and $\tau$. One possibility to assess its impact is by sampling from the prior and plotting the implied distribution on differential causal effects. This is done in Figure 9, where the black-box distribution is a mixture of 5 Gaussians with the following priors: the mixing distribution is given a Dirichlet prior with hyperparameter $(2, 2, 2, 2, 2)$; the covariance matrix $\Sigma_c$ of each mixture component $c$ is given an inverse Wishart prior with $\nu = 11 + 3$ degrees of freedom and a $\nu I$ scale matrix, with $I$ being the $11 \times 11$ identity matrix; the mean vector $\mu_c$ of each component is given a multivariate Gaussian prior with zero mean and covariance matrix $\Sigma_c$.

Given the prior for the base distribution, Figure 9 shows a smoothed depiction of 1000 samples from the priors defined by sampling the mixture model parameters from the base distribution and passing them through IV-TETRAD$^{++}$, where we choose the single element that minimizes ScoreTetrads. To generate the figure, we used the following configurations: i. $(\alpha_{cut} = 10, \tau = 0.01)$; ii. $(\alpha_{cut} = 3, \tau = 0.01)$ and iii. $(\alpha_{cut} = 10, \tau = 0.001)$. 

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Figure 9: Examples of the implied prior on the differential causal effect using the Bayesian projections algorithm IV-TETRAD++ with three choices of hyperparameters \((\alpha_{cut} = 10, \tau = 0.01), (\alpha_{cut} = 3, \tau = 0.01)\) and \((\alpha_{cut} = 10, \tau = 0.001)\).

Figure 10: (a) The posterior distribution when a single value is chosen as the output of IV-TETRAD++. (b) The resulting posterior where we uniformly sample one of the elements of the equivalence class.

As expected, decreasing these two hyperparameters makes the prior narrower. At the same time, since the effect is given by the ratio of covariance entries, we expect the implied prior to be heavy-tailed. For illustration purposes, we proceed to analyze the Sachs et al. data by choosing \((\alpha_{cut} = 10, \tau = 0.01)\) as the prior configuration.

5.3.3 Results

Figure 10(a) shows the resulting posterior after running MCMC for the mixture of Gaussians for 10,000 iterations and thinning it down to 1000 samples. The evidence that 5 mixture components is enough is given by the fact that the component of smallest posterior probability was assigned almost zero mass. The figure illustrates the difficulty of the problem and
why a single point estimate might not be a good summary of the outcome even in the case where a single causal effect is returned.

More interestingly, the Bayesian approach allows us to generate posteriors over equivalence classes in a relatively simple way, without resorting to potentially uninformative summaries based on bounds. This however requires subject matter information, as the data cannot distinguish the elements in an equivalence class. Following the philosophy of Richardson et al. (2011), we separate priors for elements which are informed by the data (equivalence class) from priors for elements which cannot be distinguished (the choice within the equivalence class, given the equivalence class). Figure 10(b) shows the posterior distribution that, for each sample $m$ of the Markov chain, generates a sample causal effect by sampling uniformly at random from the elements in the corresponding equivalence class of sample $m$. That is, the extra information here is the uniform prior, which may be taken as a default choice. In principle, another prior can be used. For instance, if it is agreed that the choice should be made based, for instance, on a distribution that assigns mass in a way that is inversely proportional to the magnitude of each element in the class (this assumes that smaller effects are more likely than large effects). In any case, the important lesson is the understanding that this final component of the posterior is independent of the data given the equivalence class. The conclusion from Figure 10(b) is that the bimodality is explained not only by the data, but also by the indeterminacy of the solution, in contrast to the bimodality found in Figure 8, which is an artefact of the sampling distribution of the corresponding estimators.

6. Related Work

The notion of “equivalence class of causal effects” is not new, even if not previously named as such. Hoyer et al. (2008) discuss how causal effects in LiNGAM models can be partially identified if the number of latent confounders is known. Their method also returns sets of causal effects, but without any need for further covariates and instruments. However, the size of the equivalence class of causal effects grows rapidly with the number of assumed latent variables and it is unclear how to determine the number of unmeasured confounders. Maathuis et al. (2009) present extensions of the PC algorithm which return sets of causal effects that follow from different members of an equivalence class of DAGs compatible with the observable independence constraints.

Didelez and Sheehan (2007) present a survey on instrumental variables from the point of view of applications in Mendelian randomization. A complementary method to sisVIVE based on similar assumptions, but performing estimation by the reweighted median of covariance ratios, is discussed by Bowden et al. (2016). The difficulty of estimating effects with weak instruments is well-known in the statistics literature. A review that covers some background on the estimation with weak instruments is provided by Burgess et al. (2017).

A different way of inducing “approximate” instrumental variables was introduced by Silva and Evans (2016). In that case, only (approximate) independence constraints are exploited. The same machinery in principle could include approximate tetrad constraints, although

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13. This idea could also be adapted to the case where the empty set is returned in some of the MCMC samples: the posterior would be a mixture model with positive mass on the empty set and the remaining mass spread over the values found.
further research on that would be necessary. The problem of equivalence class of causal effects is also present in that approach. In such a scenario, however, there are infinitely many alternatives, which can be represented as an interval. The width of the interval is typically large and uninformative in classical bounding methods for instrumental variables (Didelez and Sheehan, 2007), but the approach introduced by Silva and Evans (2016) allows for the introduction of stronger assumptions as a way of reducing the size of the equivalence class. This borrows some concepts from earlier work such as the one Ramsahai (2012).

The use of tetrad constraints for testing the validity of particular edge exclusions in linear causal models has a long history, dating back at least to Spearman (1904). Although the generalizations discussed by Sullivant et al. (2010) were presented in the context of Gaussian models, this distributional assumption is not necessary, with their rank constraint results exploited even in the context of partially non-linear models by Spirtes (2013). More recently, tetrad constraints have been used in the discovery of latent variable model structure (Silva et al., 2006; Spirtes, 2013), where structures such as Figure 2(c) emerge but no direct relationships among observables (such as $X \rightarrow Y$) are discoverable. The combination of tetrad constraints and non-Gaussianity assumptions has been exploited by Shimizu et al. (2009), again with the target being relationships among latent variables. Tetrad tests for the validity of postulated IVs were discussed by Kuroki and Cai (2005).

The literature on learning algorithms allowing for latent variables has been growing steadily, including the Fast Causal Inference algorithm of Spirtes et al. (2000) and more recent methods that exploit constraints other than independence constraints (Tashiro et al., 2014; Nowzohour et al., 2015), but none of these methods allow for the estimation of the causal effect of $X$ and $Y$ when there is an unblocked unmeasured confounder between them. Phiromswad and Hoover (2013) introduced an algorithm for IV discovery, but it does not take into account unidentifiability issues that can be solved by exploring constraints other than covariance matrix constraints. It also returns equivalence classes of graphs and requires searching for multiple causal effects at the same time, contrary to our goal of focusing on a given causal effect $\lambda_{yx}$.

7. Conclusion

Finding instrumental variables is one of the most fundamental problems in causal inference. To the best of our knowledge, this paper provides the first treatment on how this can be systematically achieved by exploiting rank constraints and clarifying to which extent an equivalence class of solutions remains. We then proceeded to show how non-Gaussianity can be exploited in a pragmatic way, by adapting a state-of-the-art algorithm. Finally, we illustrated how empirical improvement can be obtained.

We expect that theoretical challenges in instrumental variable discovery can be further tackled by building on the findings shown here. In particular, as also hinted by Kang et al. (2016), some of the ideas here raised extend to non-linear (additive), heterogeneous effects and binary models. Methods developed by Peters et al. (2014) can potentially provide a starting point on how to allow for non-linearities in the context of instrumental variables. As discussed by Spirtes (2013), linearity is only really needed “downstream” of the choke point: that is, it would be enough that only the structural equation for outcome $Y$ is
linear. Theoretical details need to be sorted out, but IV-TETRAD$^+$ should in principle be practically applicable to non-linear models with linearity for $Y$ only.

More sophisticated graphical criteria for the identification of causal effects in linear systems were introduced by Brito and Pearl (2002). Further work has led to rich graphical criteria to identify causal effects in confounded pairs (Foygel et al., 2011), going beyond the IV criteria discussed here. This opens up the possibility elaborated discovery algorithms where back-door blocking (Entner et al., 2012) and the methods in this paper cannot provide a solution. How to perform this task in a computationally and statistically tractable way remains an open question.

Code for all experiments is available at http://www.homepages.ucl.ac.uk/~ucgtrbd/code/iv_discovery.

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Appendix A. More about T-separation

We start with a few more examples of t-separation. Consider Figure 11. Vertex $V_0$ does not d-separate $V_{i1}$ from $V_{j2}$. However, pair $(V_0; V_0)$ t-separates $V_{i1}$ from $V_{j2}$. To see that, let us list all three (simple) treks between $V_{i1}$ and $V_{j2}$:
• \((P_1^{[1]}; P_2^{[1]}) \equiv (V_{i1} \leftarrow V_1; V_1 \rightarrow V_0 \rightarrow V_{j2})\)

• \((P_1^{[2]}; P_2^{[2]}) \equiv (V_{i1} \leftarrow V_0; V_0 \rightarrow V_{j2})\)

• \((P_1^{[3]}; P_2^{[3]}) \equiv (V_{i1} \leftarrow V_0 \leftarrow V_2; V_2 \rightarrow V_{j2})\)

In the first trek, \(V_0\) is contained in \(P_2^{[1]}\); in the second case, both paths in the trek satisfy the criterion of containing \(V_0\); in the final trek, \(P_1^{[3]}\) plays this role. Notice that \((V_0; \emptyset)\) does not t-separate \(V_{i1}\) from \(V_{j2}\) (first trek remains “unblocked”) and neither does \((\emptyset; V_0)\) (third trek remains “unblocked”).

That \(V_{i1}\) is t-separated from \(V_{j2}\) given \((V_0; V_0)\), however, brings out no useful implication about the cross-covariance matrix \(\Sigma_{AB}\) of \((A \equiv \{V_{i1}\}, B \equiv \{V_{j2}\})\): all it says is that \(\text{rank}(\Sigma_{\{V_{i1}\}\{V_{j2}\}}) \leq 2\), which is a vacuous claim. We may attempt to introduce \(V_0\) in the two sets \(A\) and \(B\), assuming \(V_0\) is observable, as this is compatible with the notion that \(V_0\) can t-separate itself from itself. Unfortunately, there is nothing to be gained, as \(\text{rank}(\Sigma_{\{V_{i1}\}V_0\{V_{j2}\}V_0}) \leq 2\), which is again a vacuous claim.

Pair \((\emptyset; V_0)\) does t-separate \(V_{i1}\) and \(V_{j2}\) with a testable implication \(\text{rank}(\Sigma_{\{V_{i1}\}V_0\{V_{j2}\}V_0}) \leq 1\), but this does not provide anything useful to a structure learning algorithm, as d-separation of \(V_{i1}\) and \(V_{j2}\) given \(V_0\) can be tested directly instead.

Consider now \(A \equiv \{V_{i1}, V_{j2}\}, B \equiv \{V_{j1}, V_{j2}\}\). To be useful, we need a conditioning pair \((C_A; C_B)\) where \(|C_A| + |C_B| \leq 1\). Following the previous examples, it should be clear that no such a pair will exist out of the given variables. However, we could consider de-activating the path \(V_{i1} \leftarrow V_1 \rightarrow V_0\) by considering the partial cross-covariance matrix of \(\{V_{i1}, V_{j2}\}\) against \(\{V_{j1}, V_{j2}\}\) given \(V_1\). The Trek Separation Theorem for DAGs, however, says nothing explicit about partial cross-covariances. Sullivant et al. (2010) present versions of the theorem for some classes of DAGs under conditioning, but not in a completely general way. Theorem 2.17 of Sullivant et al. (2010) assumes that conditioning can be encoded by undirected edges so that no vertices adjacent to an undirected edge can also be at the arrowhead endpoint of a directed edge. Although one can show this does not affect the Markov properties of general DAGs under conditioning, as ancestral graphs (Richardson and Spirtes, 2002) are a special case of this, it is not immediately obvious how this affects parametric constraints.

Fortunately, even in the conditioning case we do not need to consider anything but DAGs: simply introduce the conditioning variables in both sets \(A\) and \(B\). In this example, define \(A \equiv \{V_{i1}, V_{j2}, V_1\}\) and \(B \equiv \{V_{j1}, V_{j2}, V_1\}\). Then it is not hard to verify that \((V_0; V_1)\) t-separates \(A\) from \(B\) (recall that \(V_1\) “t-separates itself from itself” using the conditioning set pair \((\emptyset; V_1)\)). Cross-covariance \(\Sigma_{AB}\) is rank-deficient, which means its determinant is zero. Because its determinant can be written as \(\sigma_{i1} (\sigma_{i1,j1,1} \sigma_{j2,j2,1} - \sigma_{i1,j2,1} \sigma_{j1,j1,1})\) and we assume \(\sigma_{i1} \neq 0\), the (conditional) tetrad constraint will hold.

The relationship between t-separation and d-separation is given by Theorem 2.11 of Sullivant et al. (2010). Namely, d-separation between set \(A\) and set \(B\) given set \(C\) holds if and only if we can partition \(C\) as two sets \(C_A\) and \(C_B\) such that \((C_A; C_B)\) t-separates \(A \cup C_A\) and \(B \cup C_B\). That is, any d-separation can be cast as some type of t-separation. However, as t-separation of \(A \cup C_A\) from \(B \cup C_B\) implies t-separation of \(A\) from \(B\), this allows us to derive rank constraints using the marginal distribution of \(A \cup B\) only, unlike the conditional independence constraints directly implied by conditioning on \(C\). If \(C\) contains
hidden variables, this is another way of understanding how t-separation implies testable constraints that cannot be achieved directly by testing for conditional independencies.

Appendix B. Proofs

We being by proving the main result of Section 3.2

Proof of Theorem 8. If \((W_i, W_j, Z)\) satisfies the Graphical Criteria, then there should be no active path between any element of \(\{W_i, W_j\}\) and \(Y\) that does not include \(X\). Hence, neither point (i) nor point (ii) in the theorem statement are possible.

Assume now that \((W_i, W_j, Z)\) does not satisfy the Graphical Criteria. Then there must be at least one active path, conditioned on \(Z\), connecting some \(W \in \{W_i, W_j\}\) to \(Y\) which does not contain \(X\). Otherwise, if all active paths contained \(X\), then they would all be interrupted by the removal of edge \(X \rightarrow Y\), as all other edges adjacent to \(X\) are into \(X\), and \(X\) is not in the conditioning set \(Z\).

Suppose one of these active paths is a trek \(T\), which should not contain any member of \(Z\) or otherwise it would not be active. \(\{W_i, W_j, Z\}\) must be t-separated from \(\{X, Y, Z\}\) by some \((C_A; C_B)\), as implied by Theorem 3 and the conditional tetrad constraint detected in Line 9. As argued in the discussion leading to the definition of conditional choke point, \(C_A \cup C_B = Z \cup \{V_0\}\) for some conditional choke point \(V_0\) not in \(Z\). If there is some other trek that is active and contains no member of \(Z \cup \{V_0\}\), then this would imply the existence of another element of \(C_A \cup C_B\) not in \(Z \cup \{V_0\}\), contrary to the fact that \(|C_A| + |C_B| \leq |Z| + 1\).

Hence, scenario (i) in the theorem statement will happen if we do not satisfy the Graphical Criteria and there exists an active trek between \(\{W_i, W_j\}\) and \(Y\).

Assume that there exists a active path \(C'\) with endpoints \(W \in \{W_i, W_j\}\) and \(Y\), and which contains at least one collider in \(Z\). Allowing this path to be non-simple, we have that all colliders in this path must be in \(Z\), and no non-colliders are in \(Z\). Let \(Z_1, Z_2, \ldots, Z_K\) be the colliders in \(C'\) as we move from \(W\) towards \(Y\). Then the trek \((P_1^K; P_2^K)\) linking \(Z_K\) to \(Y\) is such that \(Z_K\) is the sink of \(P_1^K\) and hence \(Z_K \in C_A\). Likewise, \(Z_1\) will be the sink of \(P_1^1\) in the trek \((P_1^1; P_2^1)\) between \(W\) and \(Z_1\), so that \(Z_1 \in C_B\). As we move from \(Z_1\) towards \(Z_K\) in \(C'\), we will be forced to make one element of \(Z\) to belong to both \(C_A\) and \(C_B\), as and we cannot have all consecutive treks \(Z_i \leftarrow \cdots \rightarrow Z_{i+1}\) such that \(Z_i\) and \(Z_{i+1}\) are both only in \(C_A\) or only in \(C_B\). This is illustrated in Figure 12. Therefore, there exists some \(Z' \in C_A \cap C_B\). Since \(|C_A| + |C_B| \leq |Z| + 1\), we cannot have both (i) and (ii) being true, as \(Z'\) must be counted twice and \(Z \subseteq C_A \cup C_B\) (as we have seen, if there is a conditional choke point then \(C_A \cap C_B = \emptyset\)). Moreover, we cannot have another path \(C''\) satisfying the properties of \(C'\) but which does not contain \(Z'\), as the existence of another \(Z'' \neq Z'\) would again violate \(|C_A| + |C_B| \leq |Z| + 1\).

As all active paths connecting \(\{W_i, W_j\}\) to \(Y\) must be either treks or paths where all colliders are ancestors of \(Z\), then these two cases cover all possibilities.

We now proof the other core result from Section 3.3.

Proof of Theorem 15. By the same reasoning of Theorem 8, passing Line 9 leaves us with the two possible scenarios of that theorem. However, scenario (ii) cannot take place since, by Theorem 13 and Line 15, there must be no active non-directed paths between \(\{W_i, W_j\}\)

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Figure 12: Path $W_1 \rightarrow Z_1 \leftarrow U_1 \rightarrow U_4 \rightarrow Z_2 \leftarrow U_4 \leftarrow U_5 \rightarrow Y$ is an example of a (non-simple) active path from some vertex in $\{W_1, W_2\}$ to $Y$ given $Z \equiv \{Z_1, Z_2\}$ where both $Z_1$ and $Z_2$ are colliders in this path. Here, $Z_2$ is on the “left side” of the trek $Z_2 \leftarrow U_4 \leftarrow U_5 \rightarrow Y$ than links an element of $A \equiv \{W_1, W_2, Z_1, Z_2\}$ to $B \equiv \{X, Y, Z_1, Z_2\}$. $Z_1$ is on the “right side” of the trek $W_1 \rightarrow Z_1$ that links an element of $A$ to an element of $B$. Trek $Z_1 \leftarrow U_1 \rightarrow U_4 \rightarrow Z_2$ also links an element of $A$ to an element of $B$. If we use $Z_1$ to block this trek, then $Z_1$ is now on the “left side” of a trek. If we use $Z_2$, then $Z_2$ is now on the “right side” of a trek. Hence, in order to t-separate $A$ from $B$, we need either $(\{Z_2\}; \{Z_1, Z_2\})$ or $(\{Z_1, Z_2\}; \{Z_1\})$, both of size 3 even though only two vertices are used.

and $Y$. For the same reasons, choke point $V_0$ cannot lie on an active non-directed path between $W \in \{W_i, W_j\}$ and $Y$, nor it can be in an active back-door path with endpoints in $W$ and $Y$. As $Y$ is not an ancestor of any vertex, $V_0$ must be a descendant of $\{W_i, W_j\}$ and an ancestor or $Y$. What is left is showing that $V_0$ is also an ancestor of $X$.

There must be an active path between $W$ and $X$ given $Z$ by linear faithfulness and Line 6 of the algorithm. Assume $V_0$ is not an ancestor of $X$. There must be an active path between $W$ and $X$ which is not a trek or otherwise it would be blocked by $Z \cup \{V_0\}$. However, using the same argument as in Theorem 8, this path with colliders in $Z$ would also require a member of $Z$ to be in both sets of the conditioning set pair $(C_A; C_B)$, contradicting that $|C_A| + |C_B| \leq |Z| + 1$. Therefore, this makes $V_0$ also an ancestor of $X$, fulfilling all requirements of a downstream choke point. □

The result for Theorem 13 depends on this standard theorem (Darmois, 1953; Skitovitch, 1953):

**Theorem 19 (Darmois-Skitovitch Theorem)** Let $e_1, \ldots, e_n$ be independent random variables, $n \geq 2$. Let $v_1 = \sum_{i=1}^{n} \alpha_i e_i$, $v_2 = \sum_{i=1}^{n} \beta_i e_i$ for some coefficients $\{\alpha_i\}, \{\beta_i\}$. If $v_1$ and $v_2$ are independent, then those $e_j$ for which $\alpha_j \neq 0$, $\beta_j \neq 0$ are Gaussian.

The idea is that if we assume $\{e_i\}$ are not Gaussian, $\{V_1, V_2\}$ share a common source if and only if they are dependent. Shimizu et al. (2006) and Entner et al. (2012) present a deeper discussion on how this theorem is used in causal discovery.

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For the main results, we will assume that particular algebraic (polynomial) identities implied by the model graph do not vanish at the particular parameter values of the given model (which we called “almost everywhere” results in the theorem). We will in particular consider ways of “expanding” the structural equations of each vertex according to exogenous variables, that is, any variable which is either an error term or latent variable (assuming without loss of generality that latent variables have no parents).

For each vertex $V_k$ in the model, and each exogenous ancestor $E_m$ of $V_k$, let $P_{km}$ be the set of all directed paths from $E_m$ to $V_k$. For each path $p \in P_{km}$, define $\phi_{km} = \prod_j \lambda_{jj'}$, the product of all coefficients along this path for $V_j \in V \cap p$ where $V_{jj'} = p \cap \text{par}_{G}(j)$ (that is, $V_{jj'}$ is the parent of $V_j$ in this path). We multiply coefficients following a sequence $E_m \to \cdots \to V_{jj'} \to V_j \to \cdots \to V_k$). From this,

$$V_k = \sum_{E_m \in A_G(k)} \sum_{p \in P_{km}} \phi_{km} E_m,$$

(5)

where $A_G(k)$ is the set of exogenous ancestors of $V_k$, where for $E_m = e_k$ we have $\phi_{km} \equiv 1$ and path $p$ is given by the single edge $e_k \to V_k$. We refer to the idea of expansion a few times in the proofs as a way of describing how models can be written as polynomial functions of the parameters $\theta_G = \{\lambda_{ij} \mid V_j \in \text{par}_G(i)\} \cup \{\omega_i\}$.

Overall, for a LiNGAM model $M$ with DAG $G$, we denote by $X_G$ the set of exogenous variables of $M$, and by the expanded graph of $M$ the graph $G$ augmented with the error terms and the corresponding edges $e_i \to V_i$ for all observable vertices $V_i$ in $G$.

The main result used in the proof of Theorem 13 comes from the following Lemma. Notice that the non-Gaussianity assumption and the Darmois-Skitovitch Theorem are not necessary for its proof.

**Lemma 20** Let $O \cup U$ be the set of variables in a zero-mean LiNGAM model $M$ with graph $G$, where $U$ are the latent variables of the model. Let $O_{\setminus i} \equiv O \setminus \{O_i\}$ for any $O_i \in O$. Let $r_i \equiv O_i - a^i O_{\setminus i}$ be the residual of the least-squares regression of $O_i$ on $O_{\setminus i}$, with $a$ being the corresponding least-squares coefficients. Then, almost everywhere, $r_i$ can be written as a linear function of the exogenous variables of $M$, $r_i = \sum_{E_m \in X_G} c_m E_m$, where $c_m \neq 0$ if and only if $O_i$ is d-connected to $E_m$ given $O_{\setminus i}$ in the expanded graph of $M$.

**Proof of Lemma 20.** Without loss of generality, assume that each latent variable in $U$ has no parents. We will sometimes use $X_k$ as another representation of any particular model variable (observable, latent or error term), with the index $k$ indicating particular variables in $O \cup U$ and error terms, depending on the context.

Choose some $O_i \in O$. One way of obtaining $r_i$ is by first performing least-squares regression of each model variable $X_k$ on $O_j$, for some $O_j \neq O_i$ in $O$, and calculating residuals $X_k^{(1)}$. Among model variables in $O$, define $O^{(1)}$ as the set of all residuals $\{O^{(1)}_k\}$, $k \neq j$. That is, the superscript $(1)$ here denotes that we have regressed all observable variables on one arbitrary observable $O_j \neq O_i$, collecting the residuals $O^{(1)}_k$ of the regression of $O_k$ on $O_j$ for all variables other than $O_j$. We will repeat the process while incrementing the superscript and removing one element out of $O^{(m)}$ at a time for $m = 2, 3, \ldots, n - 1$, where $n$ is the size of $O$. This is done by the process of picking some $O^{(1)}_j \in O^{(1)} \setminus \{O^{(1)}_i\}$ and regressing
\{O_k^{(1)}\} on \{O_j^{(1)}\} to obtain \{O^{(2)}\}, iterating until we are left with \{O^{(n-1)}\} containing the single element \{O_i^{(n-1)}\}, resulting in \{O_i^{(n-1)}\} = r_i. The elimination sequence can be arbitrary.

Let \(O_j\) be a vertex in \(O_\setminus \{i\}\). Let \(\lambda_{km}\) be the structural coefficient between \(O_k\) and any \(X_m \in O \cup U\). We define \(\lambda_{km} \equiv 0\) if \(X_m\) is not a parent of \(O_k\) in \(G\). Since

\[O_k = \lambda_{kj} O_j + \sum_{X_m \in \text{par}_G(k) \setminus O_j} \lambda_{km} X_m + e_k,\]

we have

\[\sigma_{kj} = \lambda_{kj} \sigma_{jj} + \sum_{X_m \in \text{par}_G(k) \setminus O_j} \lambda_{km} \sigma_{mj} + \sigma_{ek,j},\]

where \(\sigma_{ek,j}\) is the covariance of \(e_k\) and \(O_j\), and \(\sigma_{mj}\) here represents the covariance of \(X_m\) and \(O_j\). Dividing both sides by \(\sigma_{jj}\) implies

\[a_{kj}^{(1)} = \lambda_{kj} + \sum_{X_m \in \text{par}_G(k) \setminus O_j} \lambda_{km} a_{mj}^{(1)} + a_{ek,j}^{(1)},\]

where \(a_{ek,j}^{(1)}\) is the least-squares regression coefficient of \(e_k\) on \(O_j\). This means \(O_k^{(1)} = O_k - a_{kj}^{(1)} O_j\) can be written as

\[O_k^{(1)} = \sum_{X_m \in \text{par}_G(k) \setminus O_j} \lambda_{km} X_m^{(1)} + e_k^{(1)},\]

with \(X_m^{(1)}\) and \(e_k^{(1)}\) defined analogously.

We can iterate this process until we are left with \(r_i\):

\[r_i = \sum_{U_k \in \text{par}_G(i) \cap U} \lambda_{ik} U_k^{(n-1)} + e_i^{(n-1)}.\]

Variable \(U_k^{(n-1)}\) is the residual of the regression of \(U_k\) on \(O_\setminus \{i\}\), similarly for \(e_i^{(n-1)}\).

What we will show next is that within (8) each \(U_k^{(n-1)}\) and \(e_i^{(n-1)}\) can be expanded as polynomial functions of \(\theta_G\) and \(X_G\), and the end result will contain non-vanishing monomials that are a (linear) function of only the exogenous variables \(E_m\) which are d-connected to \(O_i\) given \(O \setminus O_i\) in the expanded graph of \(M\). Since the monomials cannot vanish except for a strict subset of lower dimensionality than that of the set of possible \(\theta_G\), the result will hold almost everywhere.

Since we are free to choose the elimination ordering leading to \(r_i\), as they all lead to the same equivalent relation (8), let us define it in a way that a vertex can be eliminated at stage \(t\) only when its has no ancestors in \(O^{(t-1)}\) (where \(O^{(0)} \equiv O_\setminus \{i\}\). We will define the meaning of the ancestral relationship for elements of \(O^{(t)}\) shortly in terms of DAGs \(G^{(t)}\), but for \(t = 0\) it is the one implied by the expanded version of \(G\).

For \(t = 1\), the only exogenous variables which will have a non-zero coefficient multiplying \(O_j\) in the least-squares regression are the parents of \(O_j\) in the expanded graph, since \(O_j\) has
no other ancestors.\footnote{14} Let $U_{k}^{(1)}$ be the residual of some latent parent of $O_{j}$,\footnote{15} 

$$U_{k}^{(1)} = U_{k} - a_{kj}\left(\sum_{E_{m}\in A_{\delta}(j)}\sum_{p\in P_{jm}}\phi_{jmp}E_{m}\right), \quad \text{(9)}$$

where $\phi_{jmp} \equiv \lambda_{jm}$ if $E_{m}$ is a latent variable, or 1 if $E_{m} = e_{j}$. Moreover, $a_{kj} = \lambda_{jk}\omega_{kk}/\sigma_{jj}$, where $\omega_{kk}$ is the variance parameter of $U_{k}$ and $\sigma_{jj}$ is a polynomial function of $\theta_{G}$. We can multiply both sides of the equation above by $\sigma_{jj}$ (as well all equations referring to any $O_{k}^{(1)}$ or $X_{m}^{(1)}$ such as (7)) to get a new system of variables that is polynomial in $\theta_{G}$. We will adopt this step implicitly and claim that from (9) we have that $U_{k}^{(1)}$ can be expanded as parameters that are polynomial functions of $\theta_{G}$. Moreover, it is clear from (9) that there will be at least one non-vanishing monomial containing each $E_{m}$. In what follows, we refer to any expression analogous to (9) as the expansion of $U_{k}^{(1)}$ for $t = 1, 2, \ldots, n - 1$.

We define a DAG $G^{(1)}$ with vertices $\{X_{m}^{(1)}\}$, where the children of $X_{m}^{(1)}$ are defined as all and only the $O_{k}^{(1)}$ such that $X_{m}$ is a parent of $O_{k}$ in the original extended graph of the model. That is, $G^{(1)}$ is the extended graph over residuals after the first regression. The respective model $M^{(1)}$ is given by equations of type (7) with parameters coming from $\theta_{G}$.

For any $t > 1$, let $O_{j}^{(t)}$ be the vertex being eliminated. Each $U_{k}^{(t)}$ in which $U_{k}^{(t-1)}$ is a parent of $O_{j}^{(t-1)}$ in $G^{(t-1)}$ will be a polynomial function of $\theta_{G}$ and a linear function of the union of the exogenous variables present in the expansion of each parent of $O_{j}^{(t-1)}$: the expansion analogous to (9) in the new model will always introduce new symbols $\lambda_{j\ast}$ into existing monomials, or create new monomials with $e_{j}$, as vertex $O_{j}^{(t-1)}$ had no eliminated descendants up to iteration $t$. As such, no exogenous variable will be eliminated from the algebraic expansion of the respective $U_{k}^{(t)}$.

Finally, the expansion of $\lambda_{ik}U_{k}^{(n-1)}$ in (8) will not cancel any monomial in the expansion of some other $\lambda_{ik}U_{k'}^{(n-1)}$: since $U_{k}$ and $U_{k'}$ are both parents of $O_{i}$, no monomial in the expansion of $U_{k}$ can differ from a monomial in the expansion of $U_{k'}$ by a factor of $\lambda_{ik}\lambda_{ik'}$. So (8) will depend algebraically on the union of the exogenous terms leading to each $U_{k}^{(n-1)}$.

To prove the Lemma, we start by pointing out that $U_{k}^{(n-1)}$ will have a latent/error parent of some $O_{j}$ in its expansion if and only if there is at least one sequence of vertices $(O_{c}, \ldots, O_{j})$ where $O_{c}$ is an observable child of $U_{k}$ and any two consecutive elements in this sequence have at least one common latent parent in $G$ (the sequence can be a singleton, $O_{c} = O_{j}$). To see this, notice that the different $U_{k}^{(t)}$ form an equivalence relation: each $U_{k}^{(t)}$ with a $O_{j}^{(t-1)}$ child which is being eliminated at iteration $t$ will include into its expansion the exogenous variables found in the expansion of the other parents of $O_{j}^{(t-1)}$. This partitions $O_{i}$ into sets in which each vertex $O_{j}$ can “reach” some other vertex $O_{k}$ by first moving to some $O_{j'}$ which shares a latent parent with $O_{j}$ and which can “reach” $O_{k}$. The latent parents of $O$ are then partitioned according to their observed children.
To finalize the proof, suppose $O_i$ is d-separated from a latent/error parent $E_m$ of $O_j$ given $O \setminus O_i$. This happens if and only if all latent parents of $O_i$ (and $e_i$) are d-separated from $E_m$ given $O \setminus O_i$. Let $U_k$ be a latent parent of $O_i$ (or its error term). Then $U_k^{(n-1)}$ cannot have $E_m$ in its expansion. If this was the case, by the previous paragraph $U_k$ would be d-connected to all latent parents of $O_j$, meaning $O_i$ would be d-connected to them. This implies $c_m = 0$. Conversely, suppose $O_i$ is d-connected to the error term or a latent parents of $O_j$ given $O \setminus O_i$. Then again by the previous paragraph, for any latent parent $U_k$ of $O_i$, $U_k^{(n-1)}$ will have the latent parents of $O_j$ as terms in its expansion, implying $c_m \neq 0$ almost everywhere.

We can now prove Theorem 13.

**Proof of Theorem 13.** Considering the system for $\{O_i, Y\} \cup Z$, we can represent the model in an equivalent way where all other variables are latent variables, and all latent variables are exogenous. Applying Lemma 20 to both $r_i$ and $r_y$, and by Theorem 19, these variables will be dependent if and only if they are a non-trivial linear function of at least one common exogenous variable $E_m$ in the model. Again by Lemma 20, this happens if and only if $O_i$ is d-connected to $E_m$ given $Z$, and $Y$ is d-connected to $E_m$ given $O_i$ and $Z$. If $Y$ is d-connected to $E_m$ given $Z$ only, we can create an active non-directed path between $O_i$ and $Y$ given $Z$ by the concatenation of the $(O_i, E_m)$ path with $(Y, E_m)$ path: since $E_m$ is exogenous, the path is non-directed as it should include the subpath $V \leftarrow E_m \rightarrow V'$ for some $V, V'$. If $Y$ is not d-connected to $E_m$ given $Z$ only — that is, we must condition on $O_i$ also to d-connect them — then $Y$ must be d-connected to $O_i$ given $Z$ by a path that is into $O_i$. This path is non-directed as $Y$ is not an ancestor of $O_i$ by assumption, and the claim again follows. □

**Remarks:** The assumptions are stronger than, for instance, the ones used in the proofs of Tashiro et al. (2014). A closely related result in that paper is its Lemma 2, a result identifying the dependence between the residual of the regression of a variable on its children. It does not use any variation of the faithfulness assumption. This is because, in their context, it is enough to detect the dependence between the residual and some children. So if some path cancellations take place, some other path cancellations cannot occur. But we need the dependence of our $r_i$ and every relevant error term, because we cannot claim that $r_i$ depends on some error terms or latent variables, while $r_y$ depends on some error terms or latent variables, if these two sets do not overlap. Although some of the ideas by Tashiro et al. (2014) could be used in our context to build partial models and from them deduce instrumental variables, it goes against our framework of solving a particular prediction problem (causal effect of a target treatment-outcome) directly, instead of doing it by recovering parts of a broader causal graph.

Finally, we have not provided an explicit discussion on how to validate the non-Gaussianity assumption by testing the non-Gaussianity of the residuals, as done by Entner et al. (2012). Or, more precisely, showing which assumptions are necessary so that testing non-Gaussianity of the residuals is equivalent to testing non-Gaussianity of the error terms. This is left as future work.
Algorithm 5 IV-TETRAD++

1: Input: \( D \), a sample from the joint distribution of random variables \( O \cup \{X, Y\} \); \( K \), number of instruments to consider; \( S \), the number of mixture components; \( M \), the number of MCMC iterations; \( \pi \), a prior for a mixture of Gaussian model for \( O \cup \{X, Y\} \) with \( M \) components; \( \alpha_{\text{cut}} \) and \( \tau \), two threshold parameters
2: Output: a distribution over equivalence classes of differential causal effect of \( X \) on \( Y \)
3: Initialize \( C_{\text{Bayes}} \) as the empty set
4: Run a MCMC algorithm to obtain \( M \) samples from the posterior distribution of a Bayesian mixture of \( S \) Gaussians model with data \( D \) and prior \( \pi \)
5: for \( m = 1, 2, \ldots, M \) do
6: Let \( \Sigma^{(m)} \) be the \( m \)-th MCMC sample of the covariance matrix of \( O \cup \{X, Y\} \)
7: Let \( \hat{\lambda}_i \leftarrow \sigma^{(m)}_{w_i y, o_i} / \sigma^{(m)}_{w_i x, o_i} \) for each \( W_i \in O \)
8: Discard all \( \hat{\lambda}_i \) such that \( |\hat{\lambda}_i| > \alpha_{\text{cut}} \)
9: Initialize \( C^{(m)} \) as the empty set, and \( W \) as \( O \)
10: while \( |W| \geq K \) do
11: Sort \( W \) as \( W^{(1)}, \ldots, W^{(|W|)} \) according to the corresponding \( \{\hat{\lambda}\} \)
12: Find \( j \in \{1, 2, \ldots, |W|\} \) that minimizes \( \text{SCORETETRADS}^*(W^{(j:K+1)}, O, X, Y, \Sigma^{(m)}) \)
13: \( W_j \leftarrow \text{EXPANDTETRADS}^*(j, K, W, O, X, Y, \Sigma^{(m)}) \)
14: IF \( \text{TESTTETRADS}^*(W_j, O, X, Y, \Sigma^{(m)}) \) THEN
15: \( \hat{\lambda}_{\text{TSLS}} \leftarrow \text{TSLS}^*(W_j, O \setminus W_j, O, X, Y, \Sigma^{(m)}) \)
16: \( C \leftarrow C \cup \hat{\lambda}_{\text{TSLS}} \)
17: END IF
18: \( W \leftarrow W \setminus W_j \)
19: end while
20: Add set \( C^{(m)} \) as an element of \( C_{\text{Bayes}} \)
21: end for
22: return \( C_{\text{Bayes}} \)

Appendix C. A Bayesian Projections Variation of IV-TETRAD+

Algorithm 5 is a method that relies on substituting hypothesis testing by a thresholding mechanism over posterior samples of an unconstrained distribution. This algorithm is referred to in Section 5.3.2.

\( \text{TESTTETRADS}^* \) takes the correlation matrix implied by \( \Sigma^{(m)} \) and returns false if
\[
|\rho^{(m)}_{w_i x, o_i} \rho^{(m)}_{w_j y, o_j} - \rho^{(m)}_{w_i y, o_i} \rho^{(m)}_{w_j x, o_j}| > \tau.
\]

\( \text{EXPANDTETRADS}^* \) and \( \text{TSLS}^* \) are just the straightforward variations of the original functions that take a covariance matrix instead of a data set as input.

References


