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Citation for published version (APA):

DOI:
10.1109/CDC40024.2019.9029338

Document status and date:
Published: 12/03/2020

Document Version:
Accepted manuscript including changes made at the peer-review stage

Please check the document version of this publication:
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Generalized sensing and actuation schemes for local module identification in dynamic networks

Karthik R. Ramaswamy, Paul M.J. Van den Hof and Arne G. Dankers

Abstract—For the problem of identifying a target module that is embedded in a dynamic network with known interconnection structure, different sets of conditions are available for the set of node signals to be measured and the set of excitation signals to be applied at particular node locations. In previous work these conditions have typically been derived from either an indirect identification approach, considering external excitation signals as inputs, or from a direct identification approach, considering measured node signals as inputs. While both approaches lead to different sets of (sufficient) conditions, in this paper we extend the flexibility in the sufficient conditions for selection of excitation and measured node signals, by combining both direct and indirect approaches. As a result we will show the benefits of using both external excitation signals and node signals as predictor inputs. The provided conditions allow us to design sensor selection and actuation schemes with considerable freedom for consistent identification of a target module.

I. INTRODUCTION

In recent years increasing attention has been given to the identification of large-scale dynamically interconnected systems (modules), known as dynamic networks. Among the large amount of literature on this topic, there are three main research trends. The first one deals with the identification of the interconnection structure (topology) of systems in the dynamic network [1], [2], [3], [4]. The second deals with identification of the full network dynamics [5], [6], [7], while the third deals with identification of a target module in the dynamic network under the assumption of known topology (known as local module identification, see [8], [9], [10], [11], [12], [13]).

In this paper we focus on the local module identification problem. In [8], the classical direct-method [14] for closed-loop identification has been generalized to a dynamic network framework using a MISO identification setup. It introduces a method to achieve a consistent estimate of the target module when all the node signals in the MISO setup are measured. In [15], an extension has been made towards the situation where some of the node signals might be non-measurable. In [16] and [17], an approach has been introduced to consistently estimate the target module in the presence of confounding variables due to non-measured nodes and noise correlation. The direct method has been extended to a Bayesian setting in [9], where regularized kernel-based methods are used to reduce the mean-square error of the target module estimates. In [12] and [10], the situation has been addressed where the node measurements are affected by sensor noise.

An important condition in the works that use the direct method [8], [15], [16], [17] is that all parallel paths from the input of the target module to its output and all loops through the output node should pass through a measured node signal that is included as a predictor input. This requirement ensures that the identified module using the direct method is equal to the target module. However, in practical situations, there can be parallel paths and loops that might have all nodes non-measured. This creates a restriction for the selection of measured node signals.

In indirect method as in [11], [7], external excitation signals are used as predictor inputs for an open loop MIMO identification problem. These methods involve two steps: (1) First obtain consistent estimates of a transfer function from external signals to measured node signal; (2) Using these estimates obtain consistent estimates of the target module (we call this step as post-processing). In [7], the freedom in selection of measured node signal is exploited under the condition that all nodes are excited.

In this paper we extend the flexibility in the sufficient conditions for selection of excitation and measured node signals for consistent target module estimates and thereby generalizing the sensing and actuation schemes. We relax the above discussed condition on the parallel paths and loops around the output node. This relaxation in conditions are achieved by combining elements of both direct and indirect approaches. We use both the node signals and external excitation signals as predictor inputs, allow post-processing of module estimates, use MIMO identification setting and thereby mixing both direct and indirect methods. The provided conditions allow us to design sensor selection and actuation schemes with considerable freedom for consistent identification of a target module.

II. NETWORK AND IDENTIFICATION SETUP

A. Dynamic network setup

Following the basic setup of [8], a dynamic network is built up of \( L \) scalar internal variables or nodes \( w_j, j = 1, \ldots, L \), and \( K \) external variables \( r_k, k = 1, \ldots, K \). Each internal variable is described as:

This project has received funding from the European Research Council (ERC), Advanced Research Grant SYSDDYNET, under the European Unions Horizon 2020 research and innovation programme (Grant Agreement No. 694504).

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\[ w_j(t) = \sum_{i \neq j}^{L} G_{ji}(q)w_i(t) + u_j(t) + v_j(t) \]  

(1)

where \( q^{-1} \) is the delay operator, i.e. \( q^{-1}w_j(t) = w_j(t-1) \);

- \( G_{ji} \) is a proper rational transfer referred to as modules;
- \( u_j(t) \) is generated by the external variables \( r_k(t) \) that can directly be manipulated by the user and is given by \( u_j(t) = \sum_{k=1}^{K} R_{jk}(q)r_k(t) \) where \( R_{jk} \) are stable, proper rational transfer functions;
- \( v_j \) is process noise, where the vector process \( v = [v_1 \cdots v_L]^T \) is modelled as a stationary stochastic process with rational spectral density \( \Phi_v(\omega) \), such that there exists a white noise process \( e := [e_1 \cdots e_L]^T \), with covariance matrix \( \Lambda > 0 \) such that \( v(t) = He(t) \), where \( H \) is square, stable, monic and minimum-phase.

We will assume that the standard regularity conditions on the data are satisfied that are required for convergence results of prediction error identification method\(^1\). In this paper we consider the situation where \( u_j(t) = \sum_{k=1}^{K} R_{jk}(q)r_k(t) \), and \( R_{jk} = 1 \) if \( j = k \), \( R_{jk} = 0 \) if \( j \neq k \), and \( j = 1, \ldots, L \).

When combining the \( L \) node signals we arrive at the full network expression

\[
\begin{bmatrix}
  w_1 \\
  w_2 \\
  \vdots \\
  w_L \\
\end{bmatrix} =
\begin{bmatrix}
  G_{12} & \cdots & G_{1L} \\
  G_{21} & 0 & \cdots \\
  \vdots & \vdots & \ddots \\
  G_{L1} & \cdots & G_{L L-1} \\
\end{bmatrix}
\begin{bmatrix}
  w_1 \\
  w_2 \\
  \vdots \\
  w_L \\
\end{bmatrix} +
\begin{bmatrix}
  u_1 \\
  u_2 \\
  \vdots \\
  u_L \\
\end{bmatrix} +
He
\]

which results in the matrix equation:

\[ w = Gw + Rr + He. \]  

(2)

We will assume that the dynamic network is stable, i.e. \( (I - G)^{-1} \) is stable, and well posed (see [18] for details).

The identification problem to be considered is the problem of identifying one particular module \( G_{ji}(q) \) on the basis of a selection of measured variables \( w \), and possibly \( r \).

B. Direct method

Let us define \( N^-_j \) as the set of node indices \( k \) such that \( G_{kj} \neq 0 \), i.e. the node signals in \( N^-_j \) are the \( w \)-neighbors of the node signal \( w_j \). The nodes corresponding to the set of node indices \( k \) such that \( G_{kj} \neq 0 \) are called the \( w \)-out-neighbors of \( w_j \) (i.e. \( N^+_j \)). Let \( D_j \) denote the set of indices of the internal variables that are chosen as predictor inputs. Let \( Z_j \) denote the set of indices not in \( \{j\} \cup D_j \), i.e. \( Z_j = \{1, \ldots, L\} \setminus \{j\} \cup D_j \). Let \( w_0 \) denote the vector \( [w_{k1} \cdots w_{kn}]^T \), where \( \{k1, \ldots, kn\} = D_j \). Let \( w_0 \) denote the vector \( [w_{k1} \cdots w_{kn}]^T \), where \( \{k1, \ldots, kn\} = D_j \), and where the \( i \)th entry is zero if \( wi \) is not present in the network. The vectors \( w_2, v_1, v_2, \) and \( u_0 \) are defined analogously. The ordering of the elements in \( w_2, v_1, v_2, \) and \( u_0 \) is not important, as long as it is the same for all vectors. The transfer function matrix between \( w_0 \) and \( w_j \) is denoted \( G_{jo} \). The other transfer function matrices are defined analogously.

\(^1\)See [14] page 249. This includes the property that \( e(t) \) has bounded moments of order higher than 4.

By this notation, the network equation (2) is rewritten as:

\[
\begin{bmatrix}
  w_j \\
  w_D \\
  w_z \\
\end{bmatrix} =
\begin{bmatrix}
  0 & G_{jo} & G_{jz} \\
  G_{Dj} & G_{zo} & G_{Dz} \\
  G_{Zj} & G_{zo} & G_{ZZ} \\
\end{bmatrix}
\begin{bmatrix}
  w_j \\
  w_D \\
  w_z \\
\end{bmatrix} +
\begin{bmatrix}
  u_j \\
  v_j \\
  u_0 \\
\end{bmatrix}, \quad (3)
\]

where \( G_{zo} \) and \( G_{zz} \) have zeros on the diagonal.

Identification of module \( G_{ji} \) can be done by selecting \( D_j \) such that \( i \in D_j \), and subsequently estimating a multiple-input single output model for the transfer functions in \( G_{jo} \). This can be done by considering the one-step-ahead predictor\(^2\) \( \hat{w}_j(t|t-1) = \mathbb{E}\{w_j(t) | w_j(t-1), w_D_j\} \), and the resulting prediction error (14) \( \varepsilon_j(t, \theta) = w_j(t) - \hat{w}_j(t|t-1, \theta) \), given by

\[ \varepsilon_j(t, \theta) = H_j(\theta)^{-1} \left( w_j - \sum_{k \in D_j} G_{jk}(\theta)w_k - u_j \right) \]  

(4)

where arguments \( q \) and \( \ell \) have been dropped for notational clarity. The parameterized transfer functions \( G_{jk}(\theta), k \in D_j \) and \( H_j(\theta) \) are estimated by minimizing the sum of squared (prediction) errors: \( V_j(\theta) = \frac{1}{N} \sum_{t=0}^{N-1} \varepsilon_j^2(t, \theta) \), where \( N \) is the length of the data set. We refer to this identification method as the direct method, [8]. Let \( \hat{\theta}_N \) denote the minimizing argument of \( V_j(\theta) \).

C. Indirect method

As an alternative approach, following the setup of [7], the network model (2) can be re-written as \( w = T_{wr}r + \bar{v} \) where \( T_{wr} = (I - G)^{-1}R \) and \( \bar{v} = (I - G)^{-1}He \). Using the known external references \( r \) as predictor inputs and measured signals \( w \) as predicted outputs, it is well known that, under appropriate conditions, a consistent estimate \( \hat{T}_{wr} \) of \( T_{wr} \) can be obtained using open loop MIMO identification methods. On the basis of \( \hat{T}_{wr} \), an estimate of \( G \) can be obtained by solving \( (I - G)\hat{T}_{wr} = R \). By solving only a subset of these equations, a target module embedded in the dynamic network can be identified. We refer to this type of identification method that uses external signals as predictor inputs as the indirect method.

III. BACKGROUND AND MOTIVATING EXAMPLE

In this section we highlight the motivation of the paper using a suitable example. In [15] it has been shown that we can identify the target module \( G_{ji} \) consistently provided that we choose the selection of predictor input signals to satisfy particular properties. One of the main conditions is formulated next.

Property 1: To identify a target module \( G_{ji} \), consider a set of internal variables \( w_k, k \in D_j \). Let \( D_j \) satisfy the following properties:

1) \( i \in D_j \) and \( j \notin D_j \);
2) Every path from \( w_i \) to \( w_j \), excluding the path through \( G_{ji} \), pass through a node \( w_k, k \in D_j \) (parallel path condition);
3) Every loop through \( w_j \) pass through a node \( w_k, k \in D_j \) (loop condition).

\(^2\)Refers to \( \lim_{N \to \infty} \frac{1}{N} \sum_{t=1}^{N} \mathbb{E} \), and \( w_j^0 \) and \( w_j^f \) refer to signal samples \( w_j(\tau) \) and \( w_k(\tau), k \in D_j \), respectively, for all \( \tau \leq \ell \).
When this property is satisfied, the direct method as discussed in section II-B provides a consistent estimate of the target module, if data informativity conditions are satisfied, and in addition there are no confounding variables for the estimation problem \( w_0 \to w_3 \). This restrictive property is required for the target module in the dynamic network to be invariant in an immersed network where all non-measured signals are being removed while keeping the remaining signals invariant \([15]\).

**Example 1:** Consider a dynamic network as represented in Figure 1 with all noises in \( v \) uncorrelated with each other. For identifying the target module \( G_{21} \) (in green box), we have \( j = 2 \), and in order to satisfy Property 1 we need \( D_j = \{1, 3, 5\} \) where \( w_3 \) is included to block the parallel path from \( w_1 \) to \( w_2 \), and \( w_5 \) is included to block the loop through \( w_2 \). Using this set of measured nodes, we arrive at an immersed network after removing the non-measured node \( w_4 \) as sketched in Figure 2. We can observe that the module between \( w_1 \) and \( w_2 \) (the green box) is \( G_{21} \) and remains invariant.

![Fig. 1. Example network](image)

If \( w_3 \) and \( w_5 \) are not selected in \( D_j \), and so \( D_j = \{1\} \), we arrive at an immersed network after removing all non-measured nodes, as depicted in Figure 3. We can now observe that the dynamic module between \( w_1 \) and \( w_2 \) (the green box in figure 3) is not equal to \( G_{21} \). The terms \( (1 - G_{25}G_{52})^{-1} \) and \( G_{23}G_{31} \) are due to the fact that in this situation the loop and parallel path condition in property 1 are not satisfied, respectively. In this paper we are going to relax these restrictive conditions in property 1 and increase the freedom in the selection of measured node signals.

For the approach based on the indirect identification method, in \([11]\) a method has been presented to identify a target module using external signals as predictor inputs, along the following reasoning:

**Proposition 1** (from \([11]\)) In order to identify a target module \( G_{ji} \), perform the following experiment:

1) Excite node \( w_i \) and all its \( w \)-out-neighbors with sufficiently rich signals. Include these excitation signals as predictor inputs;

2) measure the out-neighbors of \( w_i \). Include them as predicted outputs.

Under these conditions and using full order models for the elements of \( T_{wr} \), consistent estimates \( \hat{T}_{N_i^+N_i^+} \) of \( T_{N_i^+N_i^+} \) and \( T_{N_i^+} \) can be obtained using an open loop MIMO identification method. Then a consistent estimate of \( \hat{G}_{N_i^+} \) (which includes the target module) is obtained by,

\[
\hat{G}_{N_i^+} = [\hat{T}_{N_i^+N_i^+}]^{-1}\hat{T}_{N_i^+}
\]

A dual of this proposition with \( w \)-in-neighbors of \( w_j \) is also provided in \([11]\). It can be observed that a consistent estimate of the target module is obtained from consistent estimates of elements of \( T_{wr} \). We will refer to this step (5) of manipulating identified objects, as post-processing. Considering the earlier Example 1, we can now consistently identify our target module using an open loop MIMO identification setup with \( \{v_1, r_2, r_3\} \) as inputs and \( \{w_2, w_3\} \) as outputs. However this requires restrictive conditions on the nodes to be excited and nodes to be measured, i.e. measured excitation signals \( r_1, r_2, r_3 \). Further relaxations of these restrictive conditions on excitation and measured node signals will be addressed in the sequel.

**IV. ILLUSTRATION OF THE DEVELOPED METHOD**

In this section we illustrate the developed method in this paper with suitable examples. In this paper, we combine the ideas of both the direct and indirect method such that we introduce flexibility in the selection of excitation and measured node signals. We use both the measured node signals as well as the excitation signals as predictor inputs. In addition to that, we do not restrict to the situation of invariance of our target module after immersion as in the direct method, but use the mechanism of post-processing from the indirect method to consistently identify the target module.

**Example 2:** We now consider the same network as in Example 1 but with two constraints: (a) it is not possible
to measure $w_3$ and $w_5$; (b) it is not possible to excite node $w_1$ (i.e. no $r_1$). It can be inferred that it is not possible to consistently estimate $G_{31} = G_{21}$ using the direct method due to constraint (a). Similarly due to constraint (b), it is not possible with the indirect method either.

As shown in Example 1, if we do not measure $w_3$ and $w_5$ our target module changes to $(1 - G_{25}G_{52})^{-1}(G_{21} + G_{23}G_{31})$ in the immersed network. However, we can see that this module also contains the target module of interest $G_{21}$. Therefore we might extract the target module from this term if we know (or find) the other contributions.

Consider that we excite node $w_3$, $w_5$ and measure node $w_4$. After immersing the non-measured nodes (see [15]) we end up in a dynamic network setup as in Figure 4. Now consider the identification problem \{w_1, w_4, r_2, r_3\} → \{w_2, w_4\}. We can infer the following from the figure:

1) Identifying the transfer from $r_3 → w_4$ provides $G_{43}$ and the transfer from $w_1 → w_4$ provides $G_{43}G_{31}$. Thus we can identify $G_{31}$;

2) The transfer from $r_3 → w_2$ provides $(1 - G_{25}G_{52})^{-1}G_{23}$. The term $(1 - G_{25}G_{52})^{-1}$ is due to the fact that in the original network there is a loop around $w_2$ which is not “blocked” by a measured node. This term is given by the transfer from $r_2 → w_2$. Now, we can obtain $G_{23}$.

3) The term $G_{23}G_{31}$ is due to the fact that in the original network there is a path from $w_1$ to $w_2$ through $w_3$ which is not “blocked” by a measured node. Knowing $G_{23}$ and $G_{31}$ from the above two steps, we obtain the term $G_{23}G_{31}$. We also know $(1 - G_{25}G_{52})^{-1}$. Eventually we obtain our target module of interest from the transfer $w_1 → w_2$ (i.e. $(1 - G_{25}G_{52})^{-1}(G_{21} + G_{23}G_{31})$).

This shows that we can consistently identify the target module $G_{21}$ if we know or could consistently identify the transfer from \{w_1, w_4, r_2, r_3\} → \{w_2, w_4\}.

Remark 1: The consistency results may still require additional excitation conditions, which will be specified later on.

We can observe from Figure 4 that the noise at predictor input $w_1$ and at predicted outputs $w_2, w_4$ are correlated due to $v_3$. This is due to the fact that in the original network, $v_3$ (in turn $e_3$) has simultaneous paths to $w_1$ and $w_2$ (also $w_1$ and $w_4$), while these paths run through the unmeasured node $w_3$. Therefore $e_3$, which is a confounding variable, creates noise correlation between predictor inputs and predicted outputs. When using the prediction error framework with the MIMO setup as explained above (i.e. with \{w_2, w_4\} as predicted outputs), we only model the noise from \{w_2, e_4\} → \{w_2, w_4\} but not from the confounding variable $e_3$. This leads to a lack of consistency property of identified modules [16]. If we also predict $w_3$ (include it also as predicted output), we now model the noise from $e_3$ as well. This leads to consistent estimates. This has been studied in [20] for a two-node example network. Therefore for the Example 2, we need the MIMO identification setup \{w_1, w_4, r_2, r_3\} → \{w_1, w_2, w_4\}.

From the discussed example, we can now conjecture the following generalization:

1) Violating the parallel path condition can be handled by exciting a node in the parallel path, including the excitation signal in the predictor input, and by measuring a descendant node from the excited node, different from the output of the target module, and by including this descendant node in the predicted output;

2) Violating the loop condition can be handled by either

- exciting the output node and including the excitation signal in the predictor input; or
- exciting a node in the loop, including the excitation signal in the predictor input, and by measuring a descendant node from the excited node, different from the output of the target module, and by including this descendant node in the predicted output;

3) Confounding variables can be handled by including measured nodes as predicted outputs.

In the sequel of this paper, we will derive the formal results that underly the above conjectured statements.

V. CONCEPTS AND NOTATION

We will denote $w_j$ as the node signals in $w$ that serve as predicted outputs, and $w_o$ as the node signals in $w$ that serve as predictor inputs, and $r_p$ as the external excitation signals in $r$ that serve as predictor inputs. Next we decompose $w_o$ and $w_p$ in disjoint sets according to: $Y = Q ∪ O ∪ \{o\}$; $D = Q ∪ A$ where $w_o$ are the node signals that are common in $w_o$ and $w_p$; $w_o$ is the set of node signals that are only predicted outputs (excluding the output node of target module); $w_o$ is the output $w_j$ of the target module; if $j ∈ Q$ then \{o\} is void. Additionally we denote $w_2$ as the node signals in $w$ that are neither predicted output nor predictor input, i.e. $Z = L \setminus \{D ∪ Y\}$, where $L = \{1, 2, \cdots, L\}$. Next we define the set related to $r_p$ as $P ⊆ \mathcal{L} \setminus \mathcal{A}$.

VI. MIMO IDENTIFICATION SETUP

Consider a dynamic network defined by (2), however with\newline

\textit{cov}(e) = I$ and $H$ not necessarily monic. The identification that we need to perform refers to the estimation problem
\((w_\theta, r_\rho) \rightarrow w_y\). In order to analyze this problem, based on system equation of the considered network, we rewrite the system equations for the output variables \(w_y\).

**Proposition 2:** The systems equations for the output variables in \(w_y\) can always be written as,

\[
w_y = \tilde{G} w_\theta + H \xi_\theta + \tilde{R} r_\rho,
\]

where \(\xi_\theta\) a white noise process with dimensions conforming to \(w_\theta\), with \(\text{cov}(\xi_\theta) = \Lambda\) and with \(\tilde{H}\) being monic, stable and stably invertible.

**Proof:** The detailed proof is provided in [21]. \(\square\)

As a result we can set up a predictor model based on a parametrized model set determined by \(\mathcal{M} := \{(\tilde{G}(\theta), H(\theta), \tilde{R}(\theta), \Lambda(\theta)), \theta \in \Theta\}\), while the actual data generating system is represented by \(\mathcal{S} = (G(\theta_o), \tilde{H}(\theta_o), \tilde{R}(\theta_o), \Lambda(\theta_o))\). The corresponding identification problem is defined by considering the one-step-ahead prediction of \(w_y\), according to \(\hat{w}_y(t|t-1) := E\{w_y(t) \mid w_y(t-1), \xi_\theta, r_\rho\}\) where \(w_y^0, r_\rho^t\) denotes the past of \(w_\theta, r_\rho\) respectively, i.e. \(\{w_\theta(k)\) and \(r_\rho(k), k \leq t\). The resulting prediction error \(\varepsilon(t, \theta) := w_y(t) - \hat{w}_y(t|t-1; \theta)\) becomes:

\[
\varepsilon(t, \theta) = \tilde{H}(q, \theta)^{-1}[w_y(t) - \tilde{G}(q, \theta)w_\theta(t) - \tilde{R}(q, \theta)r_\rho(t)]
\]

and the weighted least squares identification criterion

\[
\hat{\theta}_N = \arg \min_{\theta} \frac{1}{N} \sum_{t=0}^{N-1} \varepsilon^T(t, \theta) W \varepsilon(t, \theta),
\]

with \(W\) any positive definite weighting matrix. This parameter estimate then leads to an estimated subnetwork \(\tilde{G}_{\theta_0}(q, \hat{\theta}_N)\), with the estimated module \(\tilde{G}_{ji}(q, \hat{\theta}_N)\) as one of its scalar entries.

**VII. MAIN RESULTS**

In this section we first present the consistency results for the above considered identification problem.

**Theorem 1:** Consider a (MIMO) network identification setup with predictor inputs \((w_\theta, r_\rho)\) and predicted outputs \(w_y\) as in (6). Then a prediction error identification method according to (7)-(8), applied to a parametrized model set \(\mathcal{M}\) will provide a consistent estimate of \(\tilde{G}\) and \(\tilde{R}\), if the following conditions on the sets of nodes are satisfied:

a) there are no confounding variables for the estimation problem \(w_\theta \rightarrow w_y\);

b) Every measured node signal that has an unmeasured path to a node signal in \(w_y\) is included in \(w_\theta\); and additionally:

1) \(\mathcal{M}\) is chosen to satisfy \(\mathcal{S} \in \mathcal{M}\);

2) \(\Phi_\kappa(\omega) > 0\) for a sufficiently high number of frequencies, where \(\kappa(t) := [w_\theta^T \xi_\theta^T \tilde{G}_{ji}^T \tilde{R}_i^T]^T\);

3) All \(r_\rho, k \in \mathcal{P}\) are uncorrelated to all \(\xi_\theta, \ell \in \mathcal{Y} \cup \mathcal{A}\);

4) All the elements in \(G_{\theta_0}, G_{\theta}, G_{\rho}, G_{\alpha}\) are strictly proper (or) all existing paths/loops from \(w_\theta, w_\rho, w_{\alpha}\) to \(w_{\theta}, w_{\rho}, w_{\alpha}\) have at least a delay. \(\square\)

**Proof:** The detailed proof is provided in [21]. \(\square\)

The estimate \(\tilde{G}\) contains the estimate of \(\tilde{G}_{ji}\) as one of its elements since \(w_i \in w_\theta\). However our final goal is to estimate our target module \(G_{ji}\) which will be present in \(\tilde{G}\) but will need to be extracted from this matrix through post-processing. For this post-processing step we will require two additional sets:

- A set \(\mathcal{Z}_r \subseteq \mathcal{Z} \cap \mathcal{P}\) which represents externally excited nodes in unmeasured paths\(^5\) from \(w_i\) to \(w_j\) and in loops around \(w_j\); and
- A set \(\mathcal{T} \subseteq \mathcal{Y} \setminus \{j\}\) which, for each of the nodes in \(\mathcal{Z}_r\) represents a measured descendant node that has an unmeasured path from \(w_{z_r}\), while \(w_j\) is excluded from \(\mathcal{T}\). Note that for each node \(w_k, k \in \mathcal{Z}_r\), the corresponding element in \(\mathcal{T}\) is a measured node, and therefore cannot be in the corresponding unmeasured path from \(w_i\) to \(w_j\) or loops around \(w_j\), that passes through \(w_k\). Therefore the descendant in \(\mathcal{T}\) typically breaks out of these unmeasured parallel paths/loops, as illustrated in Figure 5.

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\(^5\) An unmeasured path is a path that runs through nodes in \(w_\theta\) only. Analogously, we can define an unmeasured loop through a node \(w_j\).
f) For every subset of $\mathcal{Z}_r$ (i.e., $S_r \subseteq \mathcal{Z}_r$) with cardinality $C$, there are unmeasured paths to at least $C$ nodes in $w_T$ with each node in $w_{S_r}$ having at least one path.

Theorem 2: Consider the situation of Theorem 1. Let $i \in \mathcal{D}$ and let the sets $\mathcal{Z}_r$ and $\mathcal{T}$ satisfy Property 2. Then a consistent estimate of target module $G_{ji}$ is obtained as

$$G_{ji}(\hat{\theta}_N) = \tilde{R}_{ij}(\hat{\theta}_N)\left(\tilde{G}_{ji}(\hat{\theta}_N) - \tilde{R}_{ij}(\hat{\theta}_N)\tilde{R}_{ij}(\hat{\theta}_N)\tilde{G}_{ji}(\hat{\theta}_N)\right)^{-1}\tilde{G}_{ji}(\hat{\theta}_N)$$

where\(^6\)

1) $\tilde{R}_{ij} = \tilde{R}_{ij}$ if $w_j$ is excited by an external signal $u_j$;
2) $\tilde{R}_{ij} = \left(1 - (1 + \tilde{R}_{ij}, \tilde{R}_{ij}, \tilde{G}_{ij}, \tilde{G}_{ij})\right)^{-1}\tilde{R}_{ij}, \tilde{G}_{ij}, \tilde{G}_{ij}$ if $w_j$ is not excited by an external signal $u_j$;
3) $\tilde{G}_{ij} = \tilde{G}_{ij}$ if $i \notin \mathcal{T}$;
4) $\tilde{G}_{ij} = (\tilde{G}_{ij} + \tilde{R}_{ij})$ if $i \in \mathcal{T}$, where $\tilde{R}_{ij}$ is a column vector with every element as zero except the element corresponding to node $w_i$ which is $\tilde{R}_{ij}(1 - \tilde{R}_{ij}^{-1})$.

Proof: The detailed proof for the above theorem is provided in [21].

Here $[\cdot]_i^T$ correspond to the left inverse of the matrix. The left inverse exists if set $\mathcal{Z}_r$ and $\mathcal{T}$ have Property 2.

We interpret Property 2 using the network in Figure 5. We have one unmeasured parallel path from $w_i$ to $w_j$ and one unmeasured loop through $w_j$. Considering the parallel path, the excited node $w_2$ and its measured descendant $w_1$ ensures that Property 2a) and 2b) are satisfied with $w_2$ in $w_{\mathcal{Z}_1}$ and $w_1$ in $w_{\mathcal{Z}_2}$. Similarly, considering the unmeasured loop through $w_j$, the excited node $w_5$ and its measured descendant $w_5$ ensures that Property 2d) is satisfied with $w_5$ in $w_{\mathcal{Z}_2}$. Property 2e) is satisfied with both $w_2$ and $w_5$ being excited by external signals. However, Property 2f) is not satisfied if $w_2 = \{w_2, w_5\}$ and $w_5 = w_1$. For the subset $S_r = \mathcal{Z}_r$ with cardinality equal to 2, there are unmeasured paths to only 1 node in $w_r$ (i.e., $w_1$). Hence we choose $w_1$ in $w_r$, which is a descendant of $w_5$ and ensure that Property 2f) is satisfied. Property 2c) is redundant for this case since $i \notin \mathcal{T}$. It is important to note that $w_r$ can be any node in the network that satisfies the Property 2 and thus relaxes the sensor placement scheme.

VIII. CONCLUSIONS

A new local module identification method has been introduced that consistently identifies the target module under known topology, with a generalized scheme for selection of measured node signals and excitation of nodes. We provide flexibility in the sufficient conditions to identify a target module which creates considerable freedom in sensor selection and actuation schemes. This is achieved by combining elements of the direct and indirect identification approaches. We use both node signals and external excitation signals as predictor inputs, allow post-processing of module estimates, and use a MIMO identification setting. With this step we remove restrictive conditions on measured node signals and excitation signals that are present in the currently available methods, e.g., concerning parallel paths and loops around the output.

\(^6\)notation $(\hat{\theta}_N)$ is dropped in the following expressions.