Exogeneity tests, incomplete models, weak identification
and non-Gaussian distributions:
invariance and finite-sample distributional theory

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ABSTRACT

We study the distribution of Durbin-Wu-Hausman (DWH) and Revankar-Hartley (RH) tests for exogeneity from a finite-sample viewpoint, under the null and alternative hypotheses. We consider linear structural models with possibly non-Gaussian errors, where structural parameters may not be identified and where reduced forms can be incompletely specified (or nonparametric). On level control, we characterize the null distributions of all the test statistics. Through conditioning and invariance arguments, we show that these distributions do not involve nuisance parameters. In particular, this applies to several test statistics for which no finite-sample distributional theory is yet available, such as the standard statistic proposed by Hausman (1978). The distributions of the test statistics may be non-standard – so corrections to usual asymptotic critical values are needed – but the characterizations are sufficiently explicit to yield finite-sample (Monte-Carlo) tests of the exogeneity hypothesis. The procedures so obtained are robust to weak identification, missing instruments or misspecified reduced forms, and can easily be adapted to allow for parametric non-Gaussian error distributions. We give a general invariance result (block triangular invariance) for exogeneity test statistics. This property yields a convenient exogeneity canonical form and a parsimonious reduction of the parameters on which power depends. In the extreme case where no structural parameter is identified, the distributions under the alternative hypothesis and the null hypothesis are identical, so the power function is flat, for all the exogeneity statistics. However, as soon as identification does not fail completely, this phenomenon typically disappears. We present simulation evidence which confirms the finite-sample theory. The theoretical results are illustrated with two empirical examples: the relation between trade and economic growth, and the widely studied problem of the return of education to earnings.

Keywords: Exogeneity; Durbin-Wu-Hausman test; weak instrument; incomplete model; non-Gaussian; weak identification; identification robust; finite-sample theory; pivotal; invariance; Monte Carlo test; power.

JEL classification: C3; C12; C15; C52.
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1. Introduction

The literature on weak instruments is now considerable and has often focused on inference for the coefficients of endogenous variables in so-called “instrumental-variable regressions” (or “IV regressions”); see the reviews of Stock, Wright and Yogo (2002), Dufour (2003), Andrews and Stock (2007), and Poskitt and Skeels (2012). Although research on tests for exogeneity in IV regressions is considerable, most of these studies either deal with cases where instrumental variables are strong (thus leaving out issues related to weak instruments), or focus on the asymptotic properties of exogeneity tests. To the best of our knowledge, there is no study on the finite-sample performance of exogeneity tests when IVs can be arbitrary weak, when the errors may follow a non-Gaussian distribution, or when the reduced form is incompletely specified. The latter feature is especially important to avoid losing the validity of the test procedure when important instruments are “left-out” when applying an exogeneity test, as happens easily for some common “identification-robust” tests on model structural coefficients [see Dufour and Taamouti (2007)].

In this paper, we investigate the finite-sample properties (size and power) of exogeneity tests of the type proposed by Durbin (1954), Wu (1973), Hausman (1978), and Revankar and Hartley (1973), henceforth DWH and RH tests, allowing for: (a) the possibility of identification failure (weak instruments); (b) model errors with non-Gaussian distributions, including heavy-tailed distributions which may lack moments (such as the Cauchy distribution); and (c) incomplete reduced forms (e.g., situations where important instruments are missing or left out) and arbitrary heterogeneity in the reduced forms of potentially endogenous explanatory variables.

As pointed out early by Wu (1973), a number of economic hypotheses can be formulated in terms of independence (or “exogeneity”) between stochastic explanatory variables and the disturbance term in an equation. These include, for example, the permanent income hypothesis, expected profit maximization, and recursiveness hypotheses in simultaneous equations. Exogeneity (or “predetermination”) assumptions can also affect the “causal interpretation” of model coefficients [see Simon (1953), Engle, Hendry and Richard (1982), Angrist and Pischke (2009), Pearl (2009)], and eventually the choice of estimation method.

To achieve the above goals, we consider a general setup which allows for non-Gaussian distributions and arbitrary heterogeneity in reduced-form errors. Under the assumption that the distribution of the structural errors (given IVs) is specified up to an unknown factor (which may depend on IVs), we show that exact exogeneity tests can be obtained from all DWH and RH statistics [including Hausman (1978) statistic] through the Monte Carlo test (MCT) method [see Dufour (2006)]. The null distributions of the test statistics typically depend on specific instrument values, so “critical

---

values” should also depend on the latter. Despite this, the MCT procedure automatically controls the level irrespective of this complication, and thus avoids the need to compute critical values. Of course, as usual, the null hypothesis is interpreted here as the conjunction of all model assumptions (including “distributional” ones) with the exogeneity restriction.

The finite-sample tests built in this way are also robust to weak instruments, in the sense that they never over-reject the null hypothesis of exogeneity even when IVs are weak. This entails that size control is feasible in finite samples for all DWH and RH tests [including the Hausman (1978) test]. All exogeneity tests considered can also be described as identification-robust in finite samples. These conclusions stand in contrast with ones reached by Staiger and Stock (1997, Section D) who argue – following a local asymptotic theory – that size adjustment may not be feasible due to the presence of nuisance parameters in the asymptotic distribution. Of course, this underscores the fundamental difference between a finite-sample theory and an asymptotic approximation, even when the latter is “improved”.

More importantly, we show that the proposed Monte Carlo test procedure remains valid even if the right-hand-side (possibly) endogenous regressors are heterogenous and the reduced-form model is incompletely specified (missing instruments). Because of the latter property, we say that the DWH and RH tests are robust to incomplete reduced forms. For example, robustness to incomplete reduced forms is relevant in macroeconomic models with structural breaks in the reduced form: this shows that exogeneity tests remain applicable without knowledge of break dates. In such contexts, inference on the structural form may be more reliable than inference on the reduced form. This is of great practical interest, for example, in inference based on IV regressions and DSGE models. For further discussion of this issue, see Dufour and Taamouti (2007), Dufour, Khalaf and Kichian (2013) and Doko Tchatoka (2015b).

We study analytically the power of the tests and identify the crucial parameters of the power function. In order to do this, we first prove a general invariance property (block triangular invariance) for exogeneity test statistics – a result of separate interest, e.g. to study how nuisance parameters may affect the distributions of exogeneity test statistics. This property yields a convenient exogeneity canonical form and a parsimonious reduction of the parameters on which power depends. In particular, we give conditions under which exogeneity tests have no power, and conditions under which they have power. We show formally that the tests have little power when instruments are weak. In particular, the power of the tests cannot exceed the nominal level if all structural parameters are completely unidentified. Nevertheless, power may exist as soon as one instrument is strong (partial identification).

We present a Monte Carlo experiment which confirms our theoretical findings. In particular, simulation results confirm that the MCT versions of all exogeneity statistics considered allow one to control test size perfectly, while usual critical values (under a Gaussian error assumption) are either exact or conservative. The conservative property is visible in particular when the two-stage-least-squares (2SLS) estimator of the structural error variance is used in covariance matrices. In such cases, the MCT version of the tests allows sizable power gains.

The results are also illustrated through two empirical examples: the relation between trade and economic growth, and the widely studied problem of the return of education to earnings.

The paper is organized as follows. Section 2 formulates the model studied, and Section 3 de-
scribes the exogeneity test statistics, including a number of alternative formulations (e.g., linear-regression-based interpretations) which may have different analytical and numerical features. In Section 4, we give general characterizations of the finite-sample distributions of the test statistics and show how they can be implemented as Monte Carlo tests, with either Gaussian or non-Gaussian errors. In Section 5, we give the general block-triangular invariance result and describe the associated exogeneity canonical representation. Power is discussed in Section 6. The simulation experiment is presented in Section 7, and the empirical illustration in Section 8. We conclude in Section 9. Additional details on the formulation of the different test statistics and the proofs are supplied in Appendix.

Throughout the paper, \( I_m \) stands for the identity matrix of order \( m \). For any full-column-rank \( T \times m \) matrix \( A \), \( \bar{P}[A] = A(A'A)^{-1}A' \) is the projection matrix on the space spanned by the columns of \( A \), and \( \bar{M}[A] = I_T - \bar{P}[A] \). For arbitrary \( m \times m \) matrices \( A \) and \( B \), the notation \( A > 0 \) means that \( A \) is positive definite (p.d.), \( A \geq 0 \) means \( A \) is positive semidefinite (p.s.d.), and \( A \leq B \) means \( B - A \geq 0 \). Finally, \( \|A\| \) is the Euclidian norm of a vector or matrix, i.e., \( \|A\| = \sqrt{\text{tr}(A'A)} \).

2. Framework

We consider a structural model of the form:

\[
y = Y \beta + X_1 \gamma + u, \tag{2.1}
\]

\[
Y = g(X_1, X_2, X_3, V, \bar{\Pi}), \tag{2.2}
\]

where (2.1) is a linear structural equation, \( y \in \mathbb{R}^T \) is a vector of observations on a dependent variable, \( Y \in \mathbb{R}^{T \times G} \) is a matrix of observations on (possibly) endogenous explanatory variables which are determined by equation (2.2), \( X_1 \in \mathbb{R}^{T \times k_1} \) is a matrix of observations on exogenous variables included in the structural equation (2.1), \( X_2 \in \mathbb{R}^{T \times k_2} \) and \( X_3 \in \mathbb{R}^{T \times k_3} \) are matrices of observations on exogenous variables excluded from the structural equation, \( u = (u_1, \ldots, u_T)' \in \mathbb{R}^T \) is a vector of structural disturbances, \( V = [V_1', \ldots, V_T']' \in \mathbb{R}^{T \times G} \) is a matrix of random disturbances, \( \beta \in \mathbb{R}^G \) and \( \gamma \in \mathbb{R}^k \) are vectors of unknown fixed structural coefficients, and \( \bar{\Pi} \) is a matrix of fixed (typically unknown) coefficients. We suppose \( G \geq 1, k_1 \geq 0, k_2 \geq 0, k_3 \geq 0 \), and denote:

\[
X = [X_1, X_2] = [x_1, \ldots, x_T]', \quad \bar{X} = [X_1, X_2, X_3] = [\bar{x}_1, \ldots, \bar{x}_T]', \tag{2.3}
\]

\[
\bar{Y} = [Y, X_1], \quad Z = [Y, X_1, X_2] = [z_1, \ldots, z_T]', \quad \bar{Z} = [Y, X_1, X_2, X_3] = [\bar{z}_1, \ldots, \bar{z}_T]', \tag{2.4}
\]

\[
U = [u, V] = [u_1, \ldots, u_T]', \tag{2.5}
\]

Equation (2.2) usually represents a reduced-form equation for \( Y \). The form of the function \( g(\cdot) \) may be nonlinear or unspecified, so model (2.2) can be viewed as “nonparametric” or “semiparametric”. The inclusion of \( X_3 \) in this setup allows for \( Y \) to depend on exogenous variables not used by the exogeneity tests. This assumption is crucial, because it characterizes the fact that we consider here “incomplete models” where the reduced form for \( Y \) may not be specified and involves unknown exogenous variables. It is well known that several “identification-robust” tests for \( \beta \) [such as those proposed by Kleibergen (2002) and Moreira (2003)] are not robust to allowing a general reduced
form for \( Y \) such as the one in (2.2); see Dufour and Taamouti (2007) and Doko Tchatoka (2015b).

We also make the following rank assumption on the matrices \([Y, X]\) and \([\bar{P}[X]Y, X_1]\):

\[
[Y, X] \text{ and } [\bar{P}[X]Y, X_1] \text{ have full-column rank with probability one (conditional on } X). \tag{2.6}
\]

This (fairly standard) condition ensures that the matrices \( X, \bar{M}[X_1]Y \) and \( \bar{M}[X]Y \) have full column rank, hence the unicity of the least-squares (LS) estimates when each column of \( Y \) is regressed on \( X \), as well as the existence of a unique two-stage-least-squares (2SLS) estimate for \( \beta \) and \( \gamma \) based on \( X \) as the instrument matrix. Clearly, (2.6) holds when \( X \) has full column rank and the conditional distribution of \( Y \) given \( X \) is absolutely continuous (with respect to the Lebesgue measure).

A common additional maintained hypothesis in this context consists in assuming that \( g(\cdot) \) is a linear equation of the form

\[
Y = X_1 \Pi_1 + X_2 \Pi_2 + V = X \Pi + V \tag{2.7}
\]

where \( \Pi_1 \in \mathbb{R}^{k_1 \times G} \) and \( \Pi_2 \in \mathbb{R}^{k_2 \times G} \) are matrices of unknown reduced-form coefficients. In this case, the reduced form for \( y \) is

\[
y = X_1 \pi_1 + X_2 \pi_2 + v \tag{2.8}
\]

where \( \pi_1 = \gamma + \Pi_1 \beta, \pi_2 = \Pi_2 \beta \), and \( v = u + V \beta \). When the errors \( u \) and \( V \) have mean zero (though this assumption may also be replaced by another “location assumption”, such as zero medians), the usual necessary and sufficient condition for identification of this model is

\[
\text{rank}(\Pi_2) = G. \tag{2.9}
\]

If \( \Pi_2 = 0 \), the instruments \( X_2 \) are irrelevant, and \( \beta \) is completely unidentified. If \( 1 \leq \text{rank}(\Pi_2) < G \), \( \beta \) is not identifiable, but some linear combinations of the elements of \( \beta \) are identifiable [see Dufour and Hsiao (2008) and Doko Tchatoka (2015b)]. If \( \Pi_2 \) is close not to have full column rank [e.g., if some eigenvalues of \( \Pi_2^T \Pi_2 \) are close to zero], some linear combinations of \( \beta \) are ill-determined by the data, a situation often called “weak identification” in this type of setup [see Dufour (2003), Andrews and Stock (2007)].

We study here, from a finite-sample viewpoint, the size and power properties of the exogeneity tests of the type proposed by Durbin (1954), Wu (1973), Hausman (1978), and Revankar and Hartley (1973) for assessing the exogeneity of \( Y \) in (2.1)-(2.7) when: (a) instruments may be weak; (b) \([u, V]\) may not follow a Gaussian distribution [e.g., heavy-tailed distributions which may lack moments (such as the Cauchy distribution) are allowed]; and (c) the usual reduced-form specification (2.7) is misspecified, and \( Y \) follows the more general model (2.2) which allows for omitted instruments, an unspecified nonlinear form and heterogeneity. To achieve this, we consider the following distributional assumptions on model disturbances (where \( \mathbb{P}[\cdot] \) refers to the relevant probability measure).

**Assumption 2.1** CONDITIONAL SCALE MODEL FOR THE STRUCTURAL ERROR DISTRIBUTION.

*For some fixed vector \( a \) in \( \mathbb{R}^G \), we have:

\[
u = Va + e, \tag{2.10} \]
where $\sigma_1(\bar{X})$ is a (possibly random) function of $\bar{X}$ such that $\mathbb{P}[\sigma_1(\bar{X}) \neq 0 | \bar{X}] = 1$, and the conditional distribution of $\varepsilon$ given $\bar{X}$ is completely specified.

**Assumption 2.2** CONDITIONAL MUTUAL INDEPENDENCE OF $e$ AND $V$. $V$ and $\varepsilon$ are independent, conditional on $\bar{X}$.

In the above assumptions, possible dependence between $u$ and $V$ is parameterized by $a$, while $\varepsilon$ is independent of $V$ (conditional on $\bar{X}$), and $\sigma_1(\bar{X})$ is an arbitrary (possibly random) scale parameter which may depend on $\bar{X}$ (except for the non-degeneracy condition $\mathbb{P}[\sigma_1(\bar{X}) \neq 0 | \bar{X}] = 1$). So we call $a$ the “endogeneity parameter” of the model. Assumption 2.1 is quite general and allows for heterogeneity in the distributions of the reduced-form disturbance $V_t$, $t = 1, \ldots, T$. In particular, the rows of $V$ need not be identically distributed or independent. Further, non-Gaussian distributions are covered, including heavy-tailed distributions which may lack second moments (such as the Cauchy distribution). In such cases, $\sigma_1(\bar{X})^2$ does not represent a variance. Since the scale factor may be random, we can have $\sigma_1(\bar{X}) = \bar{\sigma}(\bar{X}, V, e)$. Of course, these conditions hold when $u = \sigma \varepsilon$, where $\sigma$ is an unknown positive constant and $\varepsilon$ is independent of $X$ with a completely specified distribution. In this context, the standard Gaussian assumption is obtained by taking: $\varepsilon \sim \mathcal{N}[0, I_T]$. The distributions of $\varepsilon$ and $\sigma_1$ may also depend on a subset of $\bar{X}$, such as $X = [X_1, X_2]$. Note also the parameter $a$ is not presumed to be identifiable, and $e$ may not be independent of $V$ – though this would be a reasonable additional assumption to consider in the present context.

In this context, we consider the hypothesis that $Y$ can be treated as independent of $u$ in (2.1), deemed the (strict) exogeneity of $Y$ with respect to $u$, so no simultaneity bias would show up if (2.1) is estimated by least squares. Under the Assumptions 2.1 and 2.2, $a = 0$ is clearly a sufficient condition for $u$ and $e$ to be independent. Further, as soon as $V$ has full column rank with probability one, $a = 0$ is also necessary for the latter independence property. This leads one to test:

$$H_0 : a = 0.$$  \hspace{1cm} (2.12)

We stress here that “exogeneity” may depend on a set of conditioning variables ($\bar{X}$), though of course we can have cases where it does not depend on $\bar{X}$ or holds unconditionally. The setup we consider in this paper allows for both possibilities.

Before we move to describe tests of exogeneity, it will be useful to study how $H_0$ can be reinterpreted in the more familiar language of covariance hypotheses, provided standard second-moment assumptions are made.

**Assumption 2.3** HOMOSKEDASTICITY. The vectors $U_t = [u_t, V_t]$, $t = 1, \ldots, T$, have zero means and the same (finite) nonsingular covariance matrix:

$$\mathbb{E}[U_t U_t^\prime | \bar{X}] = \Sigma = \begin{bmatrix} \sigma_u^2 & \sigma_{vu}^\prime \\ \sigma_{vu} & \Sigma_V \end{bmatrix} > 0, \quad t = 1, \ldots, T.$$  \hspace{1cm} (2.13)

where $\sigma_u^2$, $\sigma_{vu}$ and $\Sigma_V$ may depend on $\bar{X}$. 

Assumption 2.4 Orthogonality between $e$ and $V$. \( \mathbb{E}[V_t e_t | \bar{X}] = 0, \mathbb{E}[e_t | \bar{X}] = 0 \) and \( \mathbb{E}[e_t^2 | \bar{X}] = \sigma_e^2, \) for $t = 1, \ldots, T$.

Under the above assumptions, the reduced-form disturbances
\[
W_t = [v_t, V_t']' = [u_t + V_t' \beta, V_t']', \quad t = 1, \ldots, T,
\]
also have a nonsingular covariance matrix (conditional on $\bar{X}$),
\[
\Omega = \begin{bmatrix}
\sigma_u^2 + \beta' \Sigma_V \beta + 2 \beta' \sigma_{Vu} & \beta' \Sigma_V + \sigma_V'
\end{bmatrix} \Sigma_V^{-1}.
\]

In this context, the exogeneity hypothesis of $Y$ can be formulated as
\[
H_0 : \sigma_{Vu} = 0.
\]
Further,
\[
\sigma_{Vu} = \Sigma_V a, \quad \sigma_u^2 = \sigma_e^2 + a' \Sigma_V a = \sigma_e^2 + \sigma_{Vu}' \Sigma_V^{-1} \sigma_{Vu},
\]
so $\sigma_{Vu} = 0 \iff a = 0$, and the exogeneity of $Y$ can be assessed by testing whether $a = 0$. Note, however, that Assumptions 2.3 and 2.4 will not be needed for the results presented in this paper.

In order to study the power of exogeneity tests, it will be useful to consider the following separability assumptions.

Assumption 2.5 Endogeneity-parameter distributional separability. $\bar{\Pi}$ is not restricted by $a$, and the conditional distribution of $[V, e]$ given $\bar{X}$ does not depend on the parameter $a$.

Assumption 2.6 Reduced-form linear separability for $Y$. $Y$ satisfies the equation
\[
Y = g(X_1, X_2, X_3, \bar{\Pi}) + V.
\]

Assumption 2.5 means that the distributions of $V$ and $e$ do not depend on the endogeneity parameter $a$, while Assumption 2.6 means that $V$ can be linearly separated from $g(X_1, X_2, X_3, \bar{\Pi})$ in (2.2).

3. Exogeneity tests

We consider the four statistics proposed by Wu (1973) [$T_l, l = 1, 2, 3, 4$], the statistic proposed by Hausman (1978) [$H_1$] as well as some variants [$H_2, H_3$] occasionally considered in the literature [see, for example, Hahn et al. (2010)], and the test suggested by Revankar and Hartley (1973, RH) [$R$]. These statistics can be formulated in two alternative ways: (1) as Wald-type statistics for the difference between the two-stage least squares (2SLS) and the ordinary least squares (OLS) estimators of $\hat{\beta}$ in equation (2.1), where different statistics are obtained by changing the covariance matrix; or (2) a $F$-type significance test on the coefficients of an “extended” version of (2.1), so
the different statistics can be written in terms of the difference between restricted and unrestricted residual sum of squares.

3.1. Test statistics

We now give a unified presentation of different available DWH-type statistics. The test statistics considered can be written as follows:

\[ \mathcal{R}_i = \kappa_i(\hat{\beta} - \hat{\beta})' \Sigma_i^{-1}(\hat{\beta} - \hat{\beta}), \quad i = 1, 2, 3, 4, \]
\[ \mathcal{R}_j = T(\hat{\beta} - \hat{\beta})' \Sigma_j^{-1}(\hat{\beta} - \hat{\beta}), \quad j = 1, 2, 3, \]
\[ \mathcal{R} = \kappa_R (y' \Psi_R y / \hat{\sigma}_R^2), \]

where \( \hat{\beta} \) and \( \hat{\beta} \) are the ordinary least squares (OLS) estimator and two-stage least squares (2SLS) estimators of \( \beta \), i.e.,

\[ \hat{\beta} = (Y'M_1 Y)^{-1} Y'M_1 y, \]
\[ \hat{\beta} = [(PY)'M_1 (PY)]^{-1} (PY)'M_1 y = (Y'N_1 Y)^{-1} Y'N_1 y, \]

while we denote \( \hat{\gamma} \) and \( \hat{\gamma} \) the corresponding OLS and 2SLS estimators of \( \gamma \), and

\[ M_1 = \bar{M}[X], \quad P = \bar{P}[X], \quad M = \bar{M}[X] = I_T - \bar{P}[X], \quad N_1 = M_1 P, \]
\[ \Sigma_1 = \hat{\sigma}_1^2 \hat{\Delta}, \quad \Sigma_2 = \hat{\sigma}_2^2 \hat{\Delta}, \quad \Sigma_3 = \hat{\sigma}_3^2 \hat{\Delta}, \quad \Sigma_4 = \hat{\sigma}_4^2 \hat{\Delta}, \]
\[ \hat{\Delta} = \hat{\Delta}_IV^1 - \hat{\Delta}_LS^1, \quad \hat{\Delta}_IV = \frac{1}{T} y'N_1 y, \quad \hat{\Delta}_LS = \frac{1}{T} Y'M_1 Y, \]
\[ \hat{\sigma}^2 = \frac{1}{T} \hat{u}' \hat{u} = \frac{1}{T} (y - Y \hat{\beta})'M_1 (y - Y \hat{\beta}), \quad \hat{\sigma}^2 = \frac{1}{T} \hat{u}' \hat{u} = \frac{1}{T} (y - Y \hat{\beta})'M_1 (y - Y \hat{\beta}), \]
\[ \hat{\sigma}_1^2 = \frac{1}{T} (y - Y \hat{\beta})'N_1 (y - Y \hat{\beta}) = \hat{\sigma}^2 - \hat{\sigma}_2^2, \quad \hat{\sigma}_2^2 = \frac{1}{T} (y - Y \hat{\beta})'M_1 (y - Y \hat{\beta}), \]
\[ \hat{\sigma}_3^2 = \hat{\sigma}^2 - (\hat{\beta} - \hat{\beta})' \hat{\Delta}^{-1}(\hat{\beta} - \hat{\beta}), \]
\[ \Psi_R = \frac{1}{T} (\bar{M}[\bar{Y}] - \bar{M}[Z]), \quad \hat{\sigma}_R^2 = y' \Lambda_r y, \quad \Lambda_r = \frac{1}{T} \bar{M}[Z], \]

\[ \kappa_1 = (k_2 - G) / G, \quad \kappa_2 = (T - k_1 - 2G) / G, \quad \kappa_3 = k_2 = T - k_1 - G, \quad \kappa_R = (T - k_1 - k_2 - G) / k_2. \]

Here, \( \hat{u} \) is the vector of OLS residuals from equation (2.1) and \( \hat{\sigma}^2 \) is the corresponding OLS-based estimator of \( \sigma_u^2 \) (without correction for degrees of freedom), while \( \hat{u} \) is the vector of the 2SLS residuals and \( \hat{\sigma}^2 \) the usual 2SLS-based estimator of \( \sigma_u^2 \); \( \hat{\sigma}_1^2 \), \( \hat{\sigma}_2^2 \), \( \hat{\sigma}_3^2 \) and \( \hat{\sigma}_R^2 \) may be interpreted as alternative IV-based scaling factors. Note also that \( P_1 P = P_1 P_1 = P_1, M_1 M = M_1 M = M_1, N_1 = N_1, M_1 = M_1, N_1 = N_1, M_1 = M, \) and

\[ N_1 = M_1 P = P M_1 = P M_1 P = M_1 P M_1 = N_1 M_1 = M_1 N_1 = N_1 N_1 \]
\[ = M_1 - M = P - P_1 = \bar{P}[X] - \bar{P}[X_1] = \bar{P}[M_1 X_2]. \]
Each one of the corresponding tests rejects $H_0$ when the statistic is “large”. We also set

$$\hat{\Sigma}_V =: \frac{1}{T} \hat{V}' \hat{V}, \quad \hat{V} =: MY,$$  \hspace{1cm} (3.16)

i.e. $\hat{\Sigma}_V$ is the usual sample covariance matrix of the LS residuals ($\hat{V}$) from the reduced-form linear model (2.7).

The tests differ through the use of different “covariance matrix” estimators. $\mathscr{H}_1$ uses two different estimators of $\sigma_u^2$, while the others resort to a single scaling factor (or estimator of $\sigma_u^2$). We think the expressions given here for $\hat{V}$, $i = 1, 2, 3, 4$, in (3.1) are easier to interpret than those of Wu (1973), and show more clearly the relation with Hausman-type tests. The statistic $\mathscr{H}_1$ can be interpreted as the statistic proposed by Hausman (1978), while $\mathscr{H}_2$ and $\mathscr{H}_3$ are sometimes interpreted as variants of $\mathscr{H}_1$ [see Staiger and Stock (1997) and Hahn et al. (2010)]. We use the above notations to better see the relation between Hausman-type tests and Wu-type tests. In particular, $\hat{\Sigma}_3 = \hat{\Sigma}_2$ and $\hat{\Sigma}_4 = \hat{\Sigma}_3$, so $\mathcal{F}_3 = (\kappa_3/T) \mathcal{H}_2$ and $\mathcal{F}_4 = (\kappa_4/T) \mathcal{H}_3$. Further, $\mathcal{F}_4$ is a nonlinear monotonic transformation of $\mathcal{F}_2$:

$$\mathcal{F}_4 = \frac{\kappa_4 \mathcal{F}_2}{\mathcal{F}_2 + \kappa_2} = \frac{\kappa_4}{(\kappa_2 / \mathcal{F}_2) + 1}. \hspace{1cm} (3.17)$$

Despite these relations, the tests based on $\mathcal{F}_3$ and $\mathcal{H}_2$ are equivalent only if exact critical values are used, and similarly for the pairs $(\mathcal{F}_4, \mathcal{H}_3)$ and $(\mathcal{F}_3, \mathcal{F}_2)$. We are not aware of a simple equivalence between $\mathcal{H}_1$ and $\mathcal{F}_i$, $i = 1, 2, 3, 4$, and similarly between $\mathcal{F}_1$ and $\mathcal{H}_j$, $j = 1, 2, 3$.

The link between the formulation of Wu (1973) and the one above is discussed in Appendix A.\footnote{When the errors $U_1, \ldots, U_T$ are i.i.d. Gaussian [in which case Assumptions 2.3 and 2.4 hold], the $\mathcal{F}_2$ test of Wu (1973) can also be interpreted as the LM test of $a = 0$; see Smith (1983) and Engle (1982).} Condition (2.6) entails that $\hat{\Omega}_{IV}$, $\hat{\Omega}_{LS}$ and $\hat{\Sigma}_V$ are (almost surely) nonsingular, which in turn implies that $\Delta$ is invertible; see Lemma A.1 in Appendix. In particular, it is of interest to observe that

$$\Delta^{-1} = \hat{\Omega}_{IV} + \hat{\Omega}_{IV}(\hat{\Omega}_{LS} - \hat{\Omega}_{IV})^{-1}\hat{\Omega}_{IV} = \hat{\Omega}_{IV} + \hat{\Omega}_{IV} \hat{\Sigma}_V^{-1} \hat{\Omega}_{IV} = \hat{\Omega}_{LS} \hat{\Sigma}_V^{-1} \hat{\Omega}_{LS} - \hat{\Omega}_{LS}$$

$$= \frac{1}{T} Y'N_1 \left[ I_T + Y(Y'MY)^{-1}Y' \right] N_1 Y = \frac{1}{T} Y'M_1 [Y(Y'MY)^{-1}Y' - I_T]M_1 Y. \hspace{1cm} (3.18)$$

from which we see easily that $\Delta^{-1}$ is positive definite. Further, $\Delta^{-1}$ only depends on the least-squares residuals $M_1 Y$ and $MY$ from the regressions of $Y$ on $X_1$ and $X$ respectively, and $\Delta^{-1}$ can be bounded as follows:

$$\hat{\Omega}_{IV} \leq \Delta^{-1} \leq \hat{\Omega}_{LS} \hat{\Sigma}_V^{-1} \hat{\Omega}_{LS} \hspace{1cm} (3.19)$$

so that

$$(\hat{\beta} - \beta)' \hat{\Omega}_{IV} (\hat{\beta} - \beta) \leq (\hat{\beta} - \beta)' \Delta^{-1} (\hat{\beta} - \beta) \leq (\hat{\beta} - \beta)' \hat{\Omega}_{LS} \hat{\Sigma}_V^{-1} \hat{\Omega}_{LS} (\hat{\beta} - \beta). \hspace{1cm} (3.20)$$

To the best of our knowledge, the additive expressions in (3.18) are not available elsewhere.

Finite-sample distributional results are available for $\mathcal{F}_1$, $\mathcal{F}_2$ and $\mathcal{F}_4$ when the disturbances $u_t$ are
i.i.d. Gaussian. If \( u \sim N[0, \sigma^2 I_T] \) and \( X \) is independent of \( u \), we have:

\[
\mathcal{F}_1 \sim F(G, k_2 - G), \quad \mathcal{F}_2 \sim F(G, T - k_1 - 2G), \quad \mathcal{F} \sim F(k_2, T - k_1 - k_2 - G),
\]

(3.21)

under the null hypothesis of exogeneity. Furthermore, for large samples, we have under the null hypothesis (along with standard asymptotic regularity conditions):

\[
\mathcal{H}_l \xrightarrow{L} \chi^2(G), \quad i = 1, 2, 3 \quad \text{and} \quad \mathcal{F}_l \xrightarrow{L} \chi^2(G), \quad l = 3, 4,
\]

when \( \text{rank}(\Pi_2) = G \).

Finite-sample distributional results are not available in the literature for \( \mathcal{H}_l, \quad i = 1, 2, 3 \) and \( \mathcal{F}_l, \quad l = 3, 4 \), even when errors are Gaussian and usual full identification assumptions are made. Of course, the same remark applies when usual conditions for identification fail \( \text{rank}(\Pi_2) < G \) or get close to do so – e.g., some eigenvalues of \( \Pi_2^\prime \Pi_2 \) are close to zero (weak identification) – and disturbances may not be Gaussian. This paper provides a formal characterization of the size and power of the tests when IVs may be arbitrary weak, with and without Gaussian errors.

### 3.2. Regression-based formulations of exogeneity statistics

We now show that all the above test statistics can be computed from relatively simple linear regressions, which may be analytically revealing and computationally convenient. We consider again the regression of \( u \) on \( V \) in (2.10):

\[
u = Va + e \tag{3.22}\]

for some constant vector \( a \in \mathbb{R}^G \), where \( e \) has mean zero and variance \( \sigma_e^2 \), and is uncorrelated with \( V \) and \( X \). We can write the structural equation (2.1) in three different ways as follows:

\[
y = Y\beta + \gamma + \tilde{V}a + e, \quad \gamma = \theta + e, \tag{3.23}\]

\[
y = \tilde{Y}\beta + \gamma + \tilde{V}b + e, \quad \gamma = \theta_2 + e, \tag{3.24}\]

\[
y = \tilde{Y}b + X_1\tilde{\gamma} + X_2\tilde{a} + e, \quad \tilde{\gamma} = \tilde{\theta} + e. \tag{3.25}\]

where

\[
\tilde{Z} = [Y, X_1, \tilde{V}], \quad \theta = (\beta', \gamma', a')', \quad Z_s = [\tilde{Y}, X_1, \tilde{V}], \quad \theta_s = (\beta', \gamma', b')', \quad \tilde{Z}_s = [Y, X_1, X_2], \quad \tilde{\theta} = (b', \tilde{\gamma}', a')', \quad b = \beta + a, \quad \tilde{\gamma} = \gamma - \Pi_1 a, \quad \tilde{a} = -\Pi_2 a, \tag{3.26}\]

\[
\tilde{Y} = \tilde{P}[X]Y, \quad \tilde{V} = \tilde{M}[X]V, \quad e_s = \tilde{P}[X]Va + e. \tag{3.28}\]

Clearly, \( \beta = b \) if and only if \( a = 0 \). Equations (3.22) - (3.25) show that the endogeneity of \( Y \) in (2.1) - (2.7) can be interpreted as an omitted-variable problem [for further discussion of this view, see Dufour (1979, 1987) and Doko Tchatoka and Dufour (2014)]. The inclusion of \( \tilde{V} \) in equations (3.23) - (3.24) may also be interpreted as an application of control function methods [see Wooldridge
\[ y - \tilde{Y} \beta = X_1\tilde{Y} + X_2\tilde{a} + e_{ss} = X\theta_{ss} + e_{ss} \]  
\quad (3.29)

where \( \tilde{\beta} \) is the 2SLS estimator of \( \beta \).

Let \( \hat{\theta} \) be the OLS estimator of \( \theta \) and \( \hat{\theta}^0 \) the restricted OLS estimator of \( \theta \) under the constraint \( H_0 : a = 0 \) [in (3.23)], \( \hat{\theta}_r \) the OLS estimator of \( \theta_r \) and \( \hat{\theta}^0_r \) the restricted OLS estimate of \( \theta_r \) under \( H^0_r : \beta = b \) [in (3.24)], \( \hat{\theta} \) the OLS estimate of \( \hat{\theta} \) and \( \hat{\theta}^0 \) the restricted OLS estimate of \( \hat{\theta} \) under \( H^0 : \bar{a} = 0 \) [in (3.25)]. Similarly, the OLS estimate of \( \theta_{ss} \) based on (3.29) is denoted \( \hat{\theta}_{ss} \), while \( \hat{\theta}^0_{ss} \) represents the corresponding restricted estimate under \( H_0 : \bar{a} = 0 \). The sum of squared error functions associated with (3.23) - (3.25) are denoted:

\[ S(\theta) = ||y - \bar{Z}\theta||^2, \quad S_r(\theta_r) = ||y - Z_r\theta_r||^2, \quad \bar{S}(\hat{\theta}) = ||y - \bar{Z}\hat{\theta}||^2, \]  
\quad (3.30)
\[ \bar{S}(\hat{\theta}_{ss}) = ||y - \bar{Y}\hat{\beta} - X\theta_{ss}||^2. \]  
\quad (3.31)

Using \( Y = \hat{Y} + \hat{V} \), we see that:

\[ S(\hat{\theta}) = S_r(\hat{\theta}_r) = \bar{S}(\hat{\theta}^0), \quad S(\hat{\theta}^0) = S_r(\hat{\theta}^0_r) = \bar{S}(\hat{\theta}_{ss}), \]  
\quad (3.32)
\[ S(\hat{\theta}) = T \hat{\sigma}^2, \quad S(\hat{\theta}^0) = T \hat{\sigma}^2, \quad S_r(\hat{\theta}^0_r) = T \hat{\sigma}^2, \quad \bar{S}(\hat{\theta}_{ss}) = T \bar{\hat{\sigma}}^2. \]  
\quad (3.33)

We then get the following expressions for the statistics in (3.1) - (3.3):

\[ \mathcal{I}_1 = \kappa_1 \left( \frac{S(\hat{\theta}^0) - S(\hat{\theta})}{S_r(\hat{\theta}^0_r) - \bar{S}(\hat{\theta}_{ss})} \right), \]  
\quad (3.34)
\[ \mathcal{I}_2 = \kappa_2 \left( \frac{S(\hat{\theta}^0) - S(\hat{\theta})}{S(\hat{\theta})} \right), \quad \mathcal{I}_3 = \kappa_3 \left( \frac{S(\hat{\theta}^0) - S(\hat{\theta})}{S_r(\hat{\theta}^0_r)} \right), \quad \mathcal{I}_4 = \kappa_4 \left( \frac{S(\hat{\theta}^0) - S(\hat{\theta})}{S(\hat{\theta}^0)} \right), \]  
\quad (3.35)
\[ \mathcal{H}_2 = T \left( \frac{S(\hat{\theta}^0) - S(\hat{\theta})}{S_r(\hat{\theta}^0_r)} \right), \quad \mathcal{H}_3 = T \left( \frac{S(\hat{\theta}^0) - S(\hat{\theta})}{S(\hat{\theta}^0)} \right), \]  
\quad (3.36)
\[ \mathcal{R} = \kappa_\mathcal{R} \left[ \bar{S}(\hat{\theta}^0) - S(\hat{\theta}) / \bar{S}(\hat{\theta}) \right]. \]  
\quad (3.37)

Details on the derivation of the above formulas are given in Appendix B.

(3.36) - (3.37) provide simple regression formulations of the DWH and RH statistics in terms of restricted and unrestricted sum of squared errors in linear regressions. However, we did not find such a simple expression for the Hausman statistic \( \mathcal{H}_1 \). While DWH-type tests consider the null hypothesis \( H_0 : a = 0 \), the RH test focuses on the null hypothesis \( H^0_2 : \bar{a} = -\Pi_2 a = 0 \). If \( \text{rank}(\Pi_2) = G \), we have: \( a = 0 \) if and only if \( \bar{a} = 0 \). However, if \( \text{rank}(\Pi_2) < G \), \( \bar{a} = 0 \) does not imply \( a = 0 \): \( H_0 \) entails \( H^0_2 \), but the converse does not hold in this case.

The regression interpretation of the \( \mathcal{F}_2 \) and \( \mathcal{H}_2 \) statistics was mentioned earlier in Dufour (1979, 1987) and Nakamura and Nakamura (1981). The \( \mathcal{R} \) statistic was also derived as a standard regression test by Revankar and Hartley (1973). To our knowledge, the other regression interpretations given here are not available elsewhere.
4. **Incomplete models and pivotal properties**

In this section, we study the finite-sample null distributions of DWH-type and RH exogeneity tests under Assumption 2.1, allowing for the possibility of identification failure (or weak identification) and model incompleteness. The proofs of these results rely on two lemmas of independent interest (Lemmas C.1 - C.2) given in Appendix.

4.1. **Distributions of test statistics under exogeneity**

We first show that the exogeneity test statistics in (3.1) - (3.3) can be rewritten as follows, irrespective whether the null hypothesis holds or not.

**Proposition 4.1** QUADRATIC-FORM REPRESENTATIONS OF EXOGENEITY STATISTICS. The exogeneity test statistics in (3.1) - (3.3) can be expressed as follows:

\[
\mathcal{J}_l = \kappa_l \left( \frac{y'\Psi \bar{y}}{y'\Lambda_3 y} \right), \text{ for } l = 1, 2, 3, 4, \tag{4.1}
\]

\[
\mathcal{H}_1 = T \left( y' \Psi \psi_y y \right) = T \left( C_1 y \right) \left( (y' \Lambda_3 y) \hat{\Omega}_{IV}^{-1} - (y' \Lambda_4 y) \hat{\Omega}_{LS}^{-1} \right]^{-1} \left( C_1 y \right), \tag{4.2}
\]

\[
\mathcal{H}_2 = T \left( y' \Psi \psi_0 y \right), \quad \mathcal{H}_3 = T \left( y' \Psi \psi_3 y \right), \quad \mathcal{R} = \kappa_r \left( y' \psi_r y \right), \tag{4.3}
\]

where

\[
\Lambda_1 = \frac{1}{T} N_1 \bar{M}[N_1 Y] N_1, \quad \Lambda_2 = M_1 \left( \frac{1}{T} \bar{M}[M_1 Y] - \Psi_0 \right) M_1, \tag{4.4}
\]

\[
\Lambda_3 = \frac{1}{T} M_1 N_2 N_2 M_1, \quad \Lambda_4 = \frac{1}{T} \bar{M}[\bar{Y}] = \frac{1}{T} M_1 \bar{M}[M_1 Y] M_1, \tag{4.5}
\]

\[
\Psi[y] = C_1^\prime \hat{\Sigma}_1^{-1} C_1 = C_1^\prime \left[ (y' \Lambda_3 y) \hat{\Omega}_{IV}^{-1} - (y' \Lambda_4 y) \hat{\Omega}_{LS}^{-1} \right]^{-1} C_1, \tag{4.6}
\]

and \(\Psi_0, B_2, C_1, \Psi_r\) and \(\Lambda_3\) are defined as in Lemma C.1.

The following theorem characterizes the distributions of all exogeneity statistics under the null hypothesis of exogeneity (\(H_0 : \alpha = 0\)).

**Theorem 4.2** NULL DISTRIBUTIONS OF EXOGENEITY STATISTICS. Under the model described by (2.1) - (2.6), suppose Assumption 2.1 holds. If \(H_0 : \alpha = 0\) also holds, then the test statistics defined in (3.1) - (3.3) have the following representations:

\[
\mathcal{J}_l = \kappa_l \left( \frac{\varepsilon' \Psi \varepsilon}{\varepsilon' \Lambda_l \varepsilon} \right), \text{ for } l = 1, 2, 3, 4, \tag{4.7}
\]

\[
\mathcal{H}_1 = T \left( \varepsilon' \Psi \psi_y \varepsilon \right) = T \left( C_1 \varepsilon \right) \left[ (\varepsilon' \Lambda_3 \varepsilon) \hat{\Omega}_{IV}^{-1} - (\varepsilon' \Lambda_4 \varepsilon) \hat{\Omega}_{LS}^{-1} \right]^{-1} \left( C_1 \varepsilon \right), \tag{4.8}
\]

\[
\mathcal{H}_2 = T \left( \varepsilon' \Psi_0 \varepsilon \right) / (\varepsilon' \Lambda_3 \varepsilon), \quad \mathcal{H}_3 = T \left( \varepsilon' \Psi_3 \varepsilon \right) / (\varepsilon' \Lambda_3 \varepsilon), \quad \mathcal{R} = \kappa_r \left( \varepsilon' \Psi_r \varepsilon \right) / (\varepsilon' \Lambda_3 \varepsilon), \tag{4.9}
\]
where $\Psi_0, \Lambda_1, \ldots, \Lambda_4, \Psi_R$ and $\Lambda_R$ are defined as in Proposition 4.1. If Assumption 2.2 also holds, the distributions of the test statistics $T_1, T_2, T_3, T_4, H_1, H_2, H_3$ and $R$, conditional on $X$ and $Y$, only depend on the conditional distribution of $\varepsilon$ given $\bar{X}$, as specified in Assumption 2.1, and the values of $Y$ and $X$.

The last statement of Theorem 4.2 comes from the fact that the weighting matrices defined in (4.4) - (4.6) only depend on $X$, $Y$ and $\varepsilon$. Given $X$ and $Y$, the null distributions of the exogeneity test statistics only depend on the distribution of $\varepsilon$ given $\bar{X}$ can be simulated, exact tests can be obtained through the Monte Carlo test method [see Section 4.2]. Furthermore, the tests obtained in this way are robust to weak instruments in the sense that the level is controlled even if identification fails (or is weak). This result holds even if the distribution of $\varepsilon | \bar{X}$ does not have moments (the Cauchy distribution, for example). This may be useful, for example, in financial models with fat-tailed error distributions, such as the Student $t$ distribution. There is no further restriction on the distribution of $\varepsilon | \bar{X}$. For example, the distribution of $\varepsilon | \bar{X}$ may depend on $\bar{X}$, provided it can be simulated.

It is interesting to observe that the distribution of $V$ plays no role here, so the vectors $V_1, \ldots, V_T$ may follow arbitrary distributions with unspecified heterogeneity (or heteroskedasticity) and serial dependence. In addition to finite-sample validity of all the exogeneity tests in the presence of identification failure (or weak identification), Theorem 4.2 entails robustness to incomplete reduced forms and instrument exclusion under the null hypothesis of exogeneity. No further information is needed on the form of the reduced form for $Y$ in (2.2): $g(\cdot)$ can be an unspecified nonlinear function, $\Pi = [\Pi_1, \Pi_2]$ an unknown parameter matrix, and $V$ may follow an arbitrary distribution. This result extends to the exogeneity tests the one given in Dufour and Taamouti (2007) on Anderson-Rubin-type tests (for structural coefficients).

As long as the distribution of $\varepsilon$ (given $\bar{X}$ and $Y$) can be simulated, all tests remain valid under $H_0$, and test sizes are controlled conditional on $\bar{X}$ and $Y$, hence also unconditionally. In particular, Monte-Carlo test procedures remain valid even if the instrument matrix $X_3$ is not used by the test statistics. A similar property is underscored in Dufour and Taamouti (2007) for Anderson-Rubin tests in linear structural equation models. This observation is also useful to allow for models with structural breaks in the reduced form: exogeneity tests remain valid in such contexts without knowledge of the form and timing of breaks. In such contexts, inference on the structural form may be more reliable than inference on the reduced form, a question of great relevance for macroeconomic models; see Dufour et al. (2013). However, although the exclusion of instruments does not affect the null distributions of exogeneity test statistics, it may lead to power losses when the missing information is important.

### 4.2. Exact Monte Carlo exogeneity tests

To implement the exact Monte Carlo exogeneity tests of $H_0$ with level $\alpha$ ($0 < \alpha < 1$), we suggest the following methodology; for a more general discussion, see Dufour (2006). Suppose that the conditional distribution of $\varepsilon$ (given $\bar{X}$) is continuous, so that the conditional distribution, given $\bar{X}$, of all exogeneity statistics is also continuous. Let $\mathcal{W}$ denotes any of the DWH and RH statistic in (3.1) - (3.3). We can then proceed as follows:
1. choose $\alpha^*$ and $N$ so that
\[ \alpha = \frac{I[\alpha^*N] + 1}{N + 1} \] (4.10)
where for any nonnegative real number $x$, $I[x]$ is the largest integer less than or equal to $x$;

2. compute the test statistic $\mathcal{W}^{(0)}$ based on the observed data;

3. generate $N$ i.i.d. error vectors $\mathbf{e}^{(j)} = [\mathbf{e}_1^{(j)}, \ldots, \mathbf{e}_T^{(j)}]'$, $j = 1, \ldots, N$, according to the specified distribution of $\mathbf{e} | \bar{X}$, and compute the corresponding statistics $\mathcal{W}^{(j)}$, $j = 1, \ldots, N$, following Theorem 4.2; the distribution of each statistic does not depend on $\beta_0$ under the null hypothesis;

4. compute the empirical distribution function based on $\mathcal{W}^{(j)}$, $j = 1, \ldots, N$,
\[ \hat{F}_N(x) = \frac{\sum_{j=1}^N 1[\mathcal{W}^{(j)} \leq x]}{N + 1} \] (4.11)
or, equivalently, the simulated $p$-value function
\[ \hat{p}_N[x] = \frac{1 + \sum_{j=1}^N 1[\mathcal{W}^{(j)} \geq x]}{N + 1} \] (4.12)
where $1[C] = 1$ if condition $C$ holds, and $1[C] = 0$ otherwise;

5. reject the null hypothesis of exogeneity, $H_0$, at level $\alpha$ when $\mathcal{W}^{(0)} \geq \hat{F}_N^{-1} (1 - \alpha^*)$, where $\hat{F}_N^{-1}(q) = \inf \{x : \hat{F}_N(x) \geq q\}$ is the generalized inverse of $\hat{F}_N(\cdot)$, or (equivalently) when $\hat{p}_N[\mathcal{W}^{(0)}] \leq \alpha$.

Under $H_0$,
\[ \mathbb{P}[\mathcal{W}^{(0)} \geq \hat{F}_N^{-1} (1 - \alpha^*)] = \mathbb{P}[\hat{p}_N[\mathcal{W}^{(0)}] \leq \alpha] = \alpha \] (4.13)
so that we have a test with level $\alpha$. The property given by (4.13) is a finite-sample validity result which holds irrespective of the sample size $T$, and no asymptotic assumption is required. If the distributions of the statistics are not continuous, the Monte Carlo test procedure can easily be adapted by using “tie-breaking” method described in Dufour (2006).3

It is important to note here that the distributions of the exogeneity test statistics in Theorem 4.2 generally depend on the specific “instrument matrix” $X$ used by the tests (especially when $\mathbf{e}$ is not Gaussian), so no general valid “critical value” (independent of $X$) is available. The Monte Carlo test procedure transparently controls the level of the test irrespective of this complication, so there is no need to compute critical values.

3Without correction for continuity, the algorithm proposed for statistics with continuous distributions yields a conservative test, i.e. the probability of rejection under the null hypothesis is not larger than the nominal level ($\alpha$). Further discussion of this feature is available in Dufour (2006).
5. Block-triangular invariance and exogeneity canonical form

In this section, we establish invariance results for exogeneity tests which will be useful to study the distributions of the test statistics under the alternative hypothesis. This basic invariance property is given by the following proposition.

**Proposition 5.1** BLOCK-TRIANGULAR INVARIANCE OF EXOGENEITY TESTS. Let

\[
R = \begin{bmatrix}
    R_{11} & 0 \\
    R_{21} & R_{22}
\end{bmatrix}
\]

be a lower block-triangular matrix such that \( R_{11} \neq 0 \) is a scalar and \( R_{22} \) is a nonsingular \( G \times G \) matrix. If we replace \( y \) and \( Y \) by \( y^* = yR_{11} + YR_{21} \) and \( Y^* = YR_{22} \) in (3.1) - (3.14), the statistics \( \mathcal{T}_i \) (\( i = 1, 2, 3, 4 \)), \( \mathcal{H}_j \) (\( j = 1, 2, 3 \)) and \( \mathcal{E} \) do not change.

The above result is purely algebraic, so no statistical assumption is needed. However, when it is combined with our statistical model, it has remarkable consequences on the properties of exogeneity tests. For example, if the reduced-form errors \( V_1, \ldots, V_T \) for \( Y \) have the same nonsingular covariance matrix \( \Sigma \), the latter can be eliminated from the distribution of the test statistic by choosing \( R_{22} \) so that \( R_{22}' \Sigma R_{22} = I_G \). This entails that the distributions of the exogeneity statistics do not depend on \( \Sigma \) under both the null and the alternative hypotheses.

Consider now the following transformation matrix:

\[
R = \begin{bmatrix}
    1 & 0 \\
    -\beta + a & I_G
\end{bmatrix}
\]

Then, we have \([y^*, Y^*] = [y, Y]R \) with

\[
y^* = y - Y(\beta + a) = Y\beta + X_1\gamma + Va + e - Y(\beta + a) = \mu_{y^*}(a) + e,
\]

\[
Y^* = Y
\]

where \( \mu_{y^*}(a) \) is a \( T \times 1 \) vector such that

\[
\mu_{y^*}(a) = X_1\gamma + [V - g(X_1, X_2, X_3, V, \bar{\Pi})]a.
\]

The (invertible) transformation (5.3) - (5.4) yields the following “latent reduced-form” representation:

\[
y^* = X_1\gamma + [V - g(X_1, X_2, X_3, V, \bar{\Pi})]a + e,
\]

\[
Y = g(X_1, X_2, X_3, V, \bar{\Pi}).
\]

We say “latent” because the function \( g(\cdot) \) and the variables \( X_3 \) are unknown or unspecified. An important feature here is that the endogeneity parameter \( a \) can be isolated from other model parameters. This will allow us to get relatively simple characterizations of the power of exogeneity tests. For this reason, we will call (5.6) - (5.7), the “exogeneity canonical form” associated with model (2.1) - (2.2) along with Assumption 2.1.
In the important case where reduced-form error linear separability holds (Assumption 2.6) in addition to (2.1) - (2.2), we can write

\[ Y = g(X_1, X_2, X_3, \tilde{\Pi}) + V = \mu_Y + V \] (5.8)

which, by (2.1), entails

\[ y = \mu_Y(a) + (u + V \beta) = \mu_Y(a) + v \] (5.9)

where \( \mu_Y \) is a \( T \times G \) matrix and \( \mu_y \) is a \( T \times 1 \) vector, such that

\[ \mu_Y = g(X_1, X_2, X_3, \tilde{\Pi}), \quad \mu_y(a) = g(X_1, X_2, X_3, \tilde{\Pi}) \beta + X_1 \gamma, \] (5.10)

\[ v = u + V \beta = e + V(\beta + a). \] (5.11)

Then

\[ \mu_y(a) = \mu_y(a) - \mu_Y(\beta + a) = X_1 \gamma - g(X_1, X_2, X_3, \tilde{\Pi})a \] (5.12)

does not depend on \( V \), and the exogeneity canonical form is:

\[ y^* = X_1 \gamma - g(X_1, X_2, X_3, \tilde{\Pi})a + e, \] (5.13)

\[ Y = g(X_1, X_2, X_3, \tilde{\Pi}) + V. \] (5.14)

6. Power

In this section, we provide characterizations of the power of exogeneity tests. We first consider the general case where only Assumption 2.1 is added to the basic setup (2.1) - (2.6). To simplify the exposition, we use the following notation: for any \( T \times 1 \) vector \( x \) and \( T \times T \) matrix \( A \), we set

\[ S_T[x, A] = T x' Ax. \] (6.1)

**Theorem 6.1** **Exogeneity Test Distributions Under the Alternative Hypothesis.**

Under the model described by (2.1) - (2.6), suppose Assumption 2.1 holds. Then the test statistics defined in (3.1) - (3.3) have the following representations:

\[ \mathcal{T}_l = \kappa_l \left( \frac{S_T[u(\bar{a}), \Psi_0]}{S_T[u(\bar{a}), \Lambda_l]} \right), \quad \text{for } l = 1, 2, 3, 4, \] (6.2)

\[ \mathcal{H}_1 = T \{ u(\bar{a})' \Psi_1 [u(\bar{a})] u(\bar{a}) \}, \quad \mathcal{H}_2 = T \left( \frac{S_T[u(\bar{a}), \Psi_0]}{S_T[u(\bar{a}), \Lambda_3]} \right), \quad \mathcal{H}_3 = T \left( \frac{S_T[u(\bar{a}), \Psi_0]}{S_T[u(\bar{a}), \Lambda_4]} \right), \] (6.3)

\[ \mathcal{R} = \kappa_\pi \left( \frac{S_T[u(\bar{a}), \Psi_\pi]}{S_T[u(\bar{a}), \Lambda_\pi]} \right), \] (6.4)

where \( u(\bar{a}) = V \bar{a} + \epsilon, \bar{a} = \sigma(\tilde{X})^{-1} a, \)

\[ \Psi_1[u(\bar{a})] = C_1' \left( S_T[u(\bar{a}), \Lambda_3] \hat{\Omega}^{-1}_IV - S_T[u(\bar{a}), \Lambda_4] \hat{\Omega}^{-1}_{IV} \right)^{-1} C_1 \] (6.5)
and \( C_1, \Psi_0, \Psi_1, \Psi_R, \Lambda_x, \Lambda_1, \ldots, \Lambda_4 \) are defined as in Theorem 4.2. If Assumption 2.5 also holds, the distributions of the test statistics (conditional on \( \tilde{X} \)) depend on \( a \) only through \( \tilde{a} \) in \( u(\tilde{a}) \).

By Theorem 6.1, the distributions of all the exogeneity statistics depend on \( a \), though possibly in a rather complex way (especially when disturbances follow non-Gaussian distributions). If the distribution of \( \epsilon \) does not depend on \( \tilde{a} \) as would be typically the case – power depends on the way the distributions of the quadratic forms \( S_T[u(\tilde{a}), \Psi_1] \) and \( S_T[u(\tilde{a}), \Lambda_2] \) in (6.2) - (6.4) are modified when the value of \( \tilde{a} \) changes. Both the numerator and the denominator of the statistics in Theorem 6.1 may follow different distributions, in contrast to what happens in standard \( F \) tests in the classical linear model.

The power characterization given by Theorem 6.1 does not provide a clear picture of the parameters which determine the power of exogeneity tests. This can be done by exploiting the invariance result of Proposition 5.1, as follows.

**Theorem 6.2** Invariance-based distributions of exogeneity statistics. Under the model described by (2.1) - (2.6), suppose Assumption 2.1 holds. Then the test statistics defined in (3.1) - (3.3) have the following representations:

\[
\mathcal{H}_1 = S_T[y^*_+(\tilde{a}), \Psi_1[y^*_+(\tilde{a})]], \quad \mathcal{H}_2 = T \left( \frac{S_T[y^*_+(\tilde{a}), \Psi_2]}{S_T[y^*_+(\tilde{a}), \Lambda_2]} \right), \quad \mathcal{H}_3 = T \left( \frac{S_T[y^*_+(\tilde{a}), \Psi_3]}{S_T[y^*_+(\tilde{a}), \Lambda_4]} \right),
\]

\[
\mathcal{R} = \kappa \left( \frac{S_T[y^*_+(\tilde{a}), \Psi_4]}{S_T[y^*_+(\tilde{a}), \Lambda_R]} \right), \quad \text{for } l = 1, 2, 3, 4,
\]

where

\[
y^*_+(\tilde{a}) = \tilde{\mu}^*_y+(\tilde{a}) + M_1 \varepsilon,
\]

\[
\tilde{\mu}^*_y+(\tilde{a}) = M_1[V - g(X_1, X_2, X_3, V, \tilde{\Pi})] \tilde{a}, \quad \tilde{a} = \sigma(\tilde{X})^{-1} a,
\]

\[
\Psi_1[y^*_+(\tilde{a})] = C_1 \left( S_T[y^*_+(\tilde{a}), \Lambda_3] \tilde{\Omega}_R^{-1} - S_T[y^*_+(\tilde{a}), \Lambda_4] \tilde{\Omega}_L^{-1} \right)^{-1} C_1,
\]

and \( C_1, \Psi_0, \Psi_1, \Psi_R, \Lambda_x, \Lambda_1, \ldots, \Lambda_4 \) are defined as in Theorem 4.2. If Assumption 2.5 also holds, the distributions of the test statistics (conditional on \( \tilde{X} \) and \( V \)) depend on \( a \) only through \( \tilde{\mu}^*_y+(\tilde{a}) \) in \( y^*_+(\tilde{a}) \). If Assumption 2.6 also holds,

\[
\tilde{\mu}^*_y+(\tilde{a}) = -M_1 g(X_1, X_2, X_3, \tilde{\Pi}) \tilde{a}.
\]

Following Theorem 6.2, the powers of the different exogeneity tests are controlled by \( \tilde{\mu}^*_y+(\tilde{a}) \) in (6.10). Clearly \( a = 0 \) entails \( \tilde{\mu}^*_y+(\tilde{a}) = 0 \), which corresponds to the distribution under the null hypothesis [under Assumption 2.5]. Note however, the latter property also holds when

\[
M_1[V - g(X_1, X_2, X_3, V, \tilde{\Pi})] = 0
\]

(6.13)
Theorem 6.3

If \( a \neq 0 \).

Under Assumption 2.6, \( V \) is evacuated from \( \bar{\mu}_y^+ (\bar{a}) \) as given by (6.12). If Assumptions 2.5 and 2.6 hold, power is determined by this parameter. \( \bar{\mu}_y^+ (\bar{a}) = 0 \) when \( a = 0 \), but also when \( X_1 \) and \( g(X_1, X_2, X_3, \bar{\Pi}) \) are orthogonal. Note also the norm of \( \bar{\mu}_y^+ (\bar{a}) \) shrinks when \( \sigma(\bar{X}) \) increases, so power decreases when the variance of value of \( \varepsilon \) increases (as expected). Under Assumption 2.6, conditioning on \( \bar{X} \) and \( V \) also becomes equivalent to conditioning on \( \bar{X} \) and \( Y \).

Consider the special case of a complete linear model where equations (2.7) and (2.8) hold. We then have:

\[
g(X_1, X_2, X_3, \bar{\Pi}) = X_1 \Pi_1 + X_2 \Pi_2, \quad \bar{\mu}_{y+1} (\bar{a}) = -M_1 X_2 \Pi_2 \bar{a}. \tag{6.14}
\]

When \( \Pi_2 = 0 \) (complete non-identification of model parameters), or \( M_1 X_2 = 0 \) (non-orthogonal with \( X_1 \)), or more generally when \( M_1 X_2 \Pi_2 = 0 \), we have \( \bar{\mu}_{y+1} (\bar{a}) = 0 \). Then, under Assumption 2.5, the distributions of the exogeneity test statistics do not depend on \( a \), and the power function is flat (with respect to \( a \)).

Theorem 6.2 provides a conditional power characterization [given \( \bar{X} \) and \( V \) (or \( Y \))]. Even though the level of the test does not depend on the distribution of \( V \), power typically depends on the distribution of \( V \). Unconditional power functions can be obtained by averaging over \( V \), but this requires formulating specific assumptions on the distribution of \( V \).

When the disturbances \( \varepsilon_1, \ldots, \varepsilon_T \) are i.i.d. Gaussian, it is possible to express the power function in terms of non-central chi-square distributions. We denote by \( \chi^2[n; \delta] \) the non-central chi-square distribution with \( n \) degrees of freedom and noncentrality parameter \( \delta \), and by \( F[n_1, n_2; \delta_1, \delta_2] \) the doubly noncentral \( F \)-distribution with degrees of freedom \( (n_1, n_2) \) and noncentrality parameters \( (\delta_1, \delta_2) \), i.e. \( F \sim F[n_1, n_2; \delta_1, \delta_2] \) means that \( F \) can be written as \( F = [Q_1/m_1] / [Q_2/m_2] \) where \( Q_1 \) and \( Q_2 \) are two independent random variables such that \( Q_1 \sim \chi^2[n_1; \delta_1] \) and \( Q_2 \sim \chi^2[n_2; \delta_2] \); see Johnson, Kotz and Balakrishnan (1995, Ch. 30). When \( \delta_2 = 0, F \sim F[n_1, n_2; \delta_1] \) the usual noncentral \( F \)-distribution.

**Theorem 6.3 Invariance-Based Distributions of Exogeneity Statistics Components with Gaussian Errors.** Under the model described by (2.1) - (2.6), suppose Assumptions 2.1 and 2.2 hold. If \( \varepsilon \sim N[0, I_T] \), then, conditional on \( \bar{X} \) and \( V \), we have:

\[
S_T[y_y^+ (\bar{a}), \Psi_y] \sim \chi^2[G; \delta(\bar{a}, \Psi_y)], \quad S_T[y_y^+ (\bar{a}), \Lambda_1] \sim \chi^2[k_2 - G; \delta(\bar{a}, \Lambda_1)],
\]

\[
S_T[y_y^+ (\bar{a}), \Lambda_2] \sim \chi^2[T - k_1 - 2G; \delta(\bar{a}, \Lambda_2)], \quad S_T[y_y^+ (\bar{a}), \Lambda_4] \sim \chi^2[T - k_1 - G; \delta(\bar{a}, \Lambda_4)],
\]

\[
S_T[y_y^+ (\bar{a}), \Psi_{LR}] \sim \chi^2[k_2; \delta(\bar{a}, \Psi_{LR})], \quad S_T[y_y^+ (\bar{a}), \Lambda_{LR}] \sim \chi^2[T - k_1 - k_2 - G; \delta(\bar{a}, \Lambda_{LR})],
\]

where

\[
\delta(\bar{a}, \Psi_y) = S_T[\bar{\mu}_y^+ (\bar{a}), \Psi_y], \quad \delta(\bar{a}, \Lambda_1) = S_T[\bar{\mu}_y^+ (\bar{a}), \Lambda_1],
\]

\[
\delta(\bar{a}, \Lambda_2) = S_T[\bar{\mu}_y^+ (\bar{a}), \Lambda_2], \quad \delta(\bar{a}, \Lambda_4) = S_T[\bar{\mu}_y^+ (\bar{a}), \Lambda_4],
\]

\[
\delta(\bar{a}, \Psi_{LR}) = S_T[\bar{\mu}_y^+ (\bar{a}), \Psi_{LR}], \quad \delta(\bar{a}, \Lambda_{LR}) = S_T[\bar{\mu}_y^+ (\bar{a}), \Lambda_{LR}]
\]

and the other symbols are defined as in Theorem 6.2. Further, conditional on \( \bar{X} \) and \( V \), the random variable \( S_T[y_y^+ (\bar{a}), \Psi_y] \) is independent of \( S_T[y_y^+ (\bar{a}), \Lambda_1] \) and \( S_T[y_y^+ (\bar{a}), \Lambda_2] \), and \( S_T[y_y^+ (\bar{a}), \Psi_{LR}] \) is
independent of \( S_T[y^+_X(a), \Lambda_x] \).

Note we do not have a chi-square distributional result for \( S_T[y^+_X(a), \Lambda_3] \) which depends on the usual 2SLS residuals. On the other hand, \( S_T[y^+_X(a), \Lambda_4] \) follows a noncentral chi-square distribution, but it is not independent of \( S_T[y^+_X(a), \Psi^+_0] \).

The noncentrality parameters in Theorem 6.3 can be interpreted as concentration parameters. For example,

\[
\delta(a, \Psi^+_0) = T [\bar{\mu}^+_y(a') \Psi^+_0 \bar{\mu}^+_y(a)] = T [\bar{\mu}^+_y(a') C_1 \hat{\Delta}^{-1} C_1 \bar{\mu}^+_y(a)] \\
= \{M_1 [V - g(X_1, X_2, X_3, V, \hat{\Pi})] \bar{a}\}' C_1 (C_1 C_1)'^{-1} C_1 \{M_1 [V - g(X_1, X_2, X_3, V, \hat{\Pi})] \bar{a}\} \\
= \{M_1 [V - g(X_1, X_2, X_3, V, \hat{\Pi})] \bar{a}\}' \bar{\Pi}' X_2 M_1 \bar{\Pi} C_1' [M_1 [V - g(X_1, X_2, X_3, V, \hat{\Pi})] \bar{a}] \\
\]  

(6.21)

and, in the case of the simple complete linear model where (2.7) and (2.8) hold,

\[
\delta(a, \Psi^+_0) = (M_1 X_2 \Pi_2 \bar{a})' \bar{\Pi}' X_2 M_1 \bar{\Pi} C_1' [M_1 X_2 \Pi_2 \bar{a}] = a' \Pi_2 X_2 M_1 \bar{\Pi} C_1' [M_1 X_2 \Pi_2 \bar{a}] . \\
\]  

(6.22)

For \( \delta(a, \Psi^+_0) \) to be different from zero, we need \( M_1 X_2 \Pi_2 \bar{a} \neq 0 \). In particular, this requires that the instruments \( X_2 \) not be totally weak (\( \Pi_2 \neq 0 \)) and linearly independent of \( X_1 \) (\( M_1 X_2 \neq 0 \)). Similar interpretations can easily be formulated for the other centrality parameters. In particular, in the simple complete linear model, all noncentrality parameters are zero if \( M_1 X_2 \Pi_2 \bar{a} = 0 \). Note, however, this may not hold in the more general model described by (2.1)-(2.6), because of the nonlinear reduced form for \( Y \) and the presence of excluded instruments.

Theorem 6.3 allows us to conclude that \( \mathcal{R}_1, \mathcal{R}_2 \) and \( \mathcal{R} \) follow doubly noncentral \( F \)-distributions under the alternative hypothesis (conditional on \( X \) and \( V \)). This is spelled out in the following corollary.

**Corollary 6.4** Doubly noncentral distributions for exogeneity statistics. **Under the model described by (2.1)-(2.6), suppose Assumptions 2.1 and 2.2 hold. If \( \varepsilon \sim N[0, I_T] \), then conditional on \( X \) and \( V \), we have:**

\[
\mathcal{R}_1 \sim F[G, k_2 - G; \delta(a, \Psi^+_0), \delta(a, \Lambda_1)] , \\
\mathcal{R}_2 \sim F[G, T - k_1 - 2G; \delta(a, \Psi^+_0), \delta(a, \Lambda_2)] , \\
\mathcal{R}_4 = \frac{K_4}{K_2 \mathcal{F}_2^{-1} + 1} \leq \left( \frac{K_4}{K_2} \right) \mathcal{R}_2 , \\
\mathcal{R} \sim F[k_2, T - k_1 - k_2 - G; \delta(a, \Psi^+_R), \delta(a, \Psi^+_R)] ,
\]


where the noncentrality parameters are defined in Theorem 6.3.

In the special case where (2.7) and (2.8) hold, we have \( \Lambda_x M_1 g(X_1, X_2, X_3, \hat{\Pi}) = \Lambda_x g(X_1, X_2, X_3, \hat{\Pi}) = 0 \) and \( \delta(a, \Psi^+_R) = 0 \), so \( \mathcal{R} \sim F[k_2, T - k_1 - k_2 - G; \delta(a, \Psi^+_R)] \) the usual noncentral noncentrality \( F \)-distribution. When \( a = 0 \), the distributions of \( \mathcal{R}_1, \mathcal{R}_2 \) and \( \mathcal{R} \) reduce to the central chi-square in (3.21) originally provided by Wu (1973) and Revankar and Hartley (1973).
The setup under which these are obtained here is considerably more general than the usual linear reduced-form specification (2.7) considered by these authors.

Note $\mathcal{H}_3$ is proportional to a ratio of two noncentral chi-square distributions, but it is not doubly-noncentral chi-square due to the non-orthogonality of $\Psi_a$ and $A_4$ [$\Psi_a A_4 = T^{-1}\Psi_a$, see (C.50)]. This observation carries to $\mathcal{H}_3$ through the identity $\mathcal{H}_3 = (T/\kappa_4) \mathcal{H}_4$. The same applies to $\mathcal{H}_1$ and $\mathcal{H}_2$, because of the presence of $S_T[y_t^2(\bar{\eta} ), A_3]$ in these statistics.

7. Simulation experiment

We use simulation to analyze the finite-sample performances (size and power) of the standard and exact Monte Carlo DWH and RH tests. The DGP is described by equations (2.1) and (2.7) without included exogenous instruments variables $X_1$, $Y = [Y_1 : Y_2] \in \mathbb{R}^{T \times 2}$, the $T \times k_2$ instrument matrix $X_2$ is a such that $X_2 \overset{i.i.d.}{\sim} \mathcal{N}(0, I_{k_2})$ for all $t = 1, \ldots, T$, and is fixed within experiment. We set the true values of $\beta$ at $\beta_0 = (2, 5)'$ but the results are qualitatively the same for alternative choices of $\beta_0$. The matrix $\Lambda_2$ that describes the quality of the instruments in the first stage regression is such that $\Pi_2 = [\eta_1 \Pi_{01} : \eta_2 \Pi_{02}] \in \mathbb{R}^{k_2 \times 2}$, where $[\Pi_{01} : \Pi_{02}]$ is obtained by taking the first two columns of the identity matrix of order $k_2$. We vary both $\eta_1$ and $\eta_2$ in $\{0, 0.01, 0.5\}$, where $\eta_1 = \eta_2 = 0$ is a design of complete non-identification, $\eta_1 = \eta_2 = 0.01$ is a design of weak identification, $\eta_1 \in \{0, 0.01\}$ and $\eta_2 = 0.5$ or vice versa is a design of partial identification, and finally, $\eta_1 = \eta_2 = 0.5$ corresponding to strong identification (strong instruments).

The errors $u$ and $V$ are generated so that

$$u = V a + e = V_1 a_1 + V_2 a_2 + e$$ (7.1)

where $a_1$ and $a_2$ are fixed scalar coefficients. In this experiment, we set $a = (a_1, a_2)' = \lambda a_0$, where $a_0 = (0.5, 0.2)'$ and $\lambda \in \{-20, -5, 0, 1, 100\}$ but the results do not change qualitatively with alternative values of $a_0$ and $\lambda$. In the above setup, $\lambda$ controls the endogeneity of $Y$: $\lambda = 0$ corresponds to the exogeneity hypothesis (level), while values of $\lambda$ different from zero represent the alternative of endogeneity (power). We consider two specifications for the joint distribution of $[e, V]$. In the first one, $(e_t, V_t)' \sim \mathcal{N}(0, I_3)$ for all $t = 1, \ldots, T$ (Gaussian errors). In the second one, $e_t$ and $V_{jt}$, $j = 1, 2$, follow a $t(3)$ distribution and are uncorrelated for all $t = 1, \ldots, T$. In both cases, $V_1$ and $V_2$ are independent. The sample size is $T = 50$, and the Monte Carlo test $p$-values are computed with $N = 199$ pseudo-samples. The simulations are based on 10000 replications. The nominal level for both the MC critical values and the standard tests is set at 5%.

7.1. Size and power with the usual critical values

Tables 1-2 present the empirical rejections of the standard DWH and RH tests for both Gaussian errors (Table 1) and $t(3)$ errors (Table 2). The first column of each table reports the statistics, while the second column contains the values of $k_2$ (number of excluded instruments). The other columns

---

4We run the experiment where $[\Pi_{01} : \Pi_{02}]$ is the $k_2 \times 2$ matrix of ones, and we found similar results as those presented here.
First, the rejection frequencies of all tests under the null hypothesis of exogeneity ($\lambda = 0$) are equal or smaller than the nominal 5% level, whether identification is weak ($\eta_1, \eta_2 \in \{0, 0.01\}$), partial ($\eta_1 \in \{0, 0.01\}$ and $\eta_2 = 0.5$ or vice versa), or strong ($\eta_1 = \eta_2 = 0.5$), with or without Gaussian errors. Thus, all DWH-type and RH tests are valid in finite samples and robust to weak instruments (i.e., level is controlled). This confirms the analysis of Section 4. As expected, the tests $F_2$, $F_4$, $H_3$, and $R$ have rejections close to the 5% nominal level. Meanwhile, $F_3$, $H_1$ and $H_2$ are highly conservative when identification is weak $[\eta_1, \eta_2 \in \{0, 0.01\}$ in the tables].

Second, all tests have power when identification is partial (columns $\lambda \neq 0$ and $\eta_1 \in \{0, 0.01\}$ and $\eta_2 = 0.5$ or vice versa) or strong (columns $\lambda \neq 0$ and $\eta_1 = \eta_2 = 0.5$), with and without Gaussian errors. Their rejection frequencies are close to 100% when $\lambda \neq 0$ and identification is strong ($\eta_1 = \eta_2 = 0.5$), despite the relatively small sample size ($T = 50$). However, all tests have low power when all instruments are irrelevant ($\lambda \neq 0$ and $\eta_1, \eta_2 \in \{0, 0.01\}$). In particular, the rejection frequencies are close to 5% when $\lambda \neq 0$, with $\eta_1, \eta_2 \in \{0, 0.01\}$, thus confirming the results of Theorems 6.2 and 6.3. The simulations also suggest that the tests $F_2$, $H_3$, $F_4$, and $R$ have greater power than the others. However, this is not always the case after size correction through the exact Monte Carlo test method, as shown in the next subsection.

7.2. Performance of the exact Monte Carlo tests

We now examine the performance of the proposed exact Monte Carlo exogeneity tests. Tables 3-4 present the results for Gaussian errors (Table 3) and $t(3)$ errors (Table 4). The results confirm our theoretical findings.

First, the rejection frequencies under the null hypothesis of exogeneity ($\lambda = 0$) of all Monte Carlo tests are around 5% whether identification is weak ($\eta_1, \eta_2 \in \{0, 0.01\}$), partial ($\eta_1 \in \{0, 0.01\}$ and $\eta_2 = 0.5$ or vice versa), or strong ($\eta_1 = \eta_2 = 0.5$), with or without Gaussian errors. This represents a substantial improvement for the standard $F_3$, $H_2$ and Hausman (1978) $H_1$ statistics.

Second, when $\lambda \neq 0$ (endogeneity), the rejection frequencies of all tests improve in most cases. This is especially the case for $F_3$, $H_1$ and $H_2$. For example, with Gaussian errors and $k_2 = 5$ instruments, the rejection frequencies of $F_3$, $H_1$ and $H_2$ have increased from 34.1%, 20.9% and 36.8% (for the standard tests) to 60.7%, 56.5% and 60.7% (for the exact Monte Carlo tests); see the columns for $\lambda = 1$ ($\eta_1 = 0.5$ and $\eta_2 = 0$) in Tables 1 and 3. The results are more remarkable with $t(3)$ errors and $k_2 = 5$ instruments. In this case, the rejection frequencies of the exact Monte Carlo $F_3$, $H_1$ and $H_2$ tests have tripled those of their standard versions; see $\lambda = 1$ ($\eta_1 = 0.5$ and $\eta_2 = 0$) in Tables 2 and 4. The results are essentially the same for other values of $k_2$, $\lambda$ and IV strength ($\eta_1$ and $\eta_2$). Moreover, except for $F_1$, the other exact Monte Carlo tests exhibit power with or without Gaussian errors, including when identification is very weak ($\eta_1 = 0.01$, $\eta_2 = 0$) and endogeneity is large ($\lambda = 100$ for example). Note that the standard exogeneity tests (including $F_2$ and $R$) perform poorly in this case. Thus, size correction through the exact Monte Carlo test method yields a substantial improvement for the exogeneity tests considered. In addition, observe that after size correction, even the Hausman (1978) statistic ($H_1$) becomes attractive in terms of power. This
Table 1. Size and power of exogeneity tests with Gaussian errors at nominal level 5%

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<th>T4</th>
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Table 2 (continued). Size and Power of exogeneity tests with \( t(3) \) errors at nominal level 5%

| \( \lambda \) | \( \kappa \) | \( \eta_1 = 0 \) | \( \eta_1 = .01 \) | \( \eta_1 = .5 \) | \( \eta_2 = 0 \) | \( \eta_2 = .01 \) | \( \eta_2 = .5 \) | \( \eta_3 = 0 \) | \( \eta_3 = .01 \) | \( \eta_3 = .5 \) | \( \eta_4 = 0 \) | \( \eta_4 = .01 \) | \( \eta_4 = .5 \) |
|-----|-----|-----|-----|-----|-----|-----|-----|-----|-----|-----|-----|-----|-----|-----|
| \( \lambda = -10 \) | 5 | 47.0 | 47.6 | 67.0 | 26.4 | 27.2 | 59.0 | 4.5 | 4.8 | 5.4 | 6.6 | 7.1 | 18.3 | 50.6 | 49.9 | 68.3 |
| \( \lambda = -5 \) | 99.7 | 99.8 | 100.0 | 83.3 | 86.2 | 99.8 | 4.6 | 4.9 | 4.9 | 8.9 | 10.1 | 48.9 | 99.9 | 99.8 | 100.0 |
| \( \lambda = 0 \) | 89.8 | 89.9 | 97.1 | 51.0 | 54.9 | 95.9 | 0.5 | 0.4 | 0.7 | 1.4 | 1.6 | 26.1 | 91.1 | 91.3 | 97.7 |
| \( \lambda = 1 \) | 99.7 | 99.8 | 100.0 | 82.5 | 85.7 | 99.8 | 4.3 | 4.5 | 4.6 | 8.3 | 9.5 | 48.0 | 99.9 | 99.8 | 100.0 |
| \( \lambda = 10 \) | 82.6 | 83.4 | 91.7 | 38.5 | 42.5 | 88.2 | 0.3 | 0.2 | 0.3 | 0.7 | 0.8 | 16.0 | 84.6 | 85.3 | 91.9 |
| \( \lambda = 20 \) | 90.8 | 90.8 | 97.3 | 54.1 | 57.7 | 96.3 | 0.6 | 0.5 | 0.8 | 1.7 | 1.8 | 28.3 | 91.8 | 92.1 | 97.9 |
| \( \lambda = 30 \) | 99.7 | 99.8 | 100.0 | 83.8 | 86.7 | 99.8 | 4.8 | 5.1 | 5.2 | 9.3 | 10.7 | 50.0 | 99.9 | 99.8 | 100.0 |

Note: The table continues with similar entries for different values of \( \lambda \) and \( \kappa \), showing the size and power of exogeneity tests with \( t(3) \) errors at nominal level 5%.
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Table 3: Size and power of exact Monte Carlo tests with Gaussian errors at nominal level 5%
Table 3 (continued). Size and power of exact Monte Carlo tests with Gaussian errors at nominal level 5%
Table 4. Size and power of exact Monte Carlo tests with $t(3)$ errors at nominal level 5%

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27
Table 4 (Continued). Size and power of exact Monte Carlo tests with $t(3)$ errors at nominal level 5%

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</table>

| $\bar{T}_{mc}$ | 5                | 6.3             | 6.5             | 18.4           | 50.3           | 51.8           | 68.9           | 99.9          | 99.9          | 100.0         |
| $\bar{T}_{mc}$ | -                | 9.1             | 9.4             | 48.9           | 96.8           | 97.2           | 99.1           | 99.9          | 99.9          | 100.0         |
| $\bar{T}_{mc}$ | -                | 10.1            | 10.1            | 52.6           | 95.6           | 96.0           | 97.7           | 99.9          | 99.9          | 100.0         |
| $\bar{T}_{mc}$ | -                | 10.1            | 10.1            | 52.6           | 99.9           | 99.9           | 100.0          | 100.0         | 100.0         | 100.0         |
| $\bar{T}_{mc}$ | 10               | 3.7             | 7.7             | 40.2           | 100.0          | 100.0          | 100.0          | 100.0         | 100.0         | 100.0         |
| $\bar{T}_{mc}$ | -                | 8.6             | 8.8             | 34.6           | 91.2           | 91.5           | 98.6           | 99.9          | 99.9          | 100.0         |
| $\bar{T}_{mc}$ | -                | 10.9            | 11.2            | 53.0           | 99.9           | 99.9           | 100.0          | 100.0         | 100.0         | 100.0         |
| $\bar{T}_{mc}$ | -                | 10.9            | 10.8            | 48.3           | 99.9           | 99.9           | 100.0          | 100.0         | 100.0         | 100.0         |
is the case in particular for $t(3)$ errors when $k_2 = 10, 20$ and $\lambda = -5, 1$; see Table 4.

8. **Empirical illustrations**

We illustrate our theoretical results on exogeneity tests through two empirical applications related to important issues in macroeconomics and labor economics literature: (1) the relation between trade and growth [Irwin and Tervio (2002), Frankel and Romer (1999), Harrison (1996), Mankiw, Romer and Weil (1992)]; (2) the standard problem of measuring returns to education [Dufour and Taamouti (2007), Angrist and Krueger (1991), Angrist and Krueger (1995), Angrist, Imbens and Krueger (1999), Mankiw et al. (1992)].

8.1. **Trade and growth**

The trade and growth model studies the relationship between standards of living and openness. Frankel and Romer (1999) argued that trade share (ratio of imports or exports to GDP) which is the commonly used indicator of openness should be viewed as endogenous. So, instrumental variables method should be used to estimate the income-trade relationship. The equation studied is

$$ln(INC_i) = \beta_0 + \beta_1 Trade_i + \gamma_1 ln(Pop_i) + \gamma_2 ln(Area_i) + u_i, \ i = 1, \ldots, T \tag{8.1}$$

where $INC_i$ is the income per capita in country $i$, $Trade_i$ is the trade share (measured as a ratio of imports and exports to GDP), $Pop_i$ is the population of country $i$, and $Area_i$ is country $i$ area. The first stage model for Trade variable is given by

$$Trade_i = a + bX_i + c_1 ln(Pop_i) + c_2 ln(Area_i) + V_i, \ i = 1, \ldots, T \tag{8.2}$$

where $X_i$ is an instrument constructed on the basis of geographic characteristics. In this paper, we use the sample of 150 countries and the data include for each country: the trade share in 1985, the area and population (1985), per capita income (1985), and the fitted trade share (instrument).

We wish to assess the exogeneity of the trade share variable in (8.1). The $F$-statistic in the first stage regression (8.2) is around 13 [see Frankel and Romer (1999, Table 2, p.385) and Dufour and Taamouti (2007)], so the fitted instrument $X$ does not appear to be weak. Table 5 presents the $p$-values of the DWH and RH tests computed from the tabulated and exact Monte Carlo critical values. The Monte Carlo critical values are computed for Gaussian and $t(3)$ errors. Because the model contains one instrument and one (supposedly) endogenous variable, the statistic $T_1$ is not well defined and is omitted.

First, we note that the $p$-values based on the usual asymptotic distributions are close to the 5% nominal level for $H_3, F_2, F_4$ and $R$. So, there is evidence against the exogeneity of the trade share (at nominal level of 5%) when these statistics are applied. Meanwhile, the $p$-values of $H_1, H_2$, and $F_3$ are relatively large (around 12%) so that there is little evidence against trade share exogeneity at 5% nominal level using the latter statistics. Since the standard $H_1, H_2$, and $F_3$ tests are conservative when identification is weak, the latter result may be due to the fact that the fitted instrument is not very strong.
Table 5. Exogeneity in trade and growth model

<table>
<thead>
<tr>
<th>Statistics</th>
<th>Estimation</th>
<th>Standard p-value (%)</th>
<th>MC p-value (%) (Gaussian errors)</th>
<th>MC p-value (%) [t(3)-errors]</th>
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<td>5.66</td>
</tr>
</tbody>
</table>

Second, we observe the exact Monte Carlo tests yield $p$-values close to the 5% level in all cases, thus indicating that there is evidence of trade share endogeneity in this model. This is supported by the relatively large discrepancy between the OLS estimate of $\beta_1$ (0.28) and the 2SLS estimate (2.03). Overall, our results underscore the importance of size correction through the exact Monte Carlo procedures proposed.

8.2. Education and earnings

We now consider the well known example of estimating the returns to education [see Angrist and Krueger (1991); Angrist and Krueger (1995); and Bound, Jaeger and Baker (1995)]. The equation studies is a relationship where the log-weekly earning ($y$) is explained by the number of years of education ($E$) and several other covariates (age, age squared, 10 dummies for birth of year):

$$y = \beta_0 + \beta_1 E + \sum_{i=1}^{k_1} \gamma_i X_i + u.$$  \hspace{1cm} (8.3)

In this model, $\beta_1$ measures the return to education. Because education can be viewed as endogenous, Angrist and Krueger (1991) used instrumental variables obtained by interacting quarter of birth with the year of birth (in this application, we use 40 dummies instruments). The basic idea is that individuals born in the first quarter of the year start school at an older age, and can therefore drop out after completing less schooling than individuals born near the end of the year. Consequently, individuals born at the beginning of the year are likely to earn less than those born during the rest of the year. The first stage model for $E$ is then given by

$$E = \pi_0 + \sum_{i=1}^{k_2} \pi_i X_i + \sum_{i=1}^{k_1} \phi_i X_i + V$$ \hspace{1cm} (8.4)

where $X$ is the instrument matrix. It is well known that the instruments $X$ constructed in this way are very weak and explains very little of the variation in education; see Bound et al. (1995). The data set consists of the 5% public-use sample of the 1980 US census for men born between 1930.
Table 6. Exogeneity in education and earning model

<table>
<thead>
<tr>
<th>Statistics</th>
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<th>Standard p-value (%)</th>
<th>MC p-value (%) (Gaussian errors)</th>
<th>MC p-value (%) [t(3)-errors]</th>
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<td>1.35</td>
<td>24.54</td>
<td>24.26</td>
<td>24.30</td>
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<td>24.54</td>
<td>24.26</td>
<td>24.30</td>
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</table>

and 1939. The sample contains 329,509 observations.

As in Section 8.2, we want to assess the exogeneity of education in (8.3) - (8.4). Table 6 shows the results of the tests with both the usual and exact Monte Carlo critical values. As seen, the p-values of all tests are quite large, thus suggesting that there is little evidence against the exogeneity of the education variable, even at 15% nominal level. This means that either the education variable is effectively exogenous or the instruments used are very poor so that the power of the test is flat, as shown in Section 6. The latter scenario is highly plausible from the previous literature [for example, see Bound et al. (1995)]. This viewed is reinforced by the small discrepancy between the OLS estimate (0.07) and the 2SLS estimate (0.08) of \( \beta_1 \).

9. Conclusion

This paper develops a finite-sample theory of the distribution of standard Durbin-Wu-Hausman and Revankar-Hartley specification tests under both the null hypothesis of exogeneity (level) and the alternative hypothesis of endogeneity (power), with or without identification. Our analysis provides several new insights and extensions of earlier procedures.

Our study of the finite-sample distributions of the statistics under the null hypothesis shows that all tests are robust to weak instruments, missing instruments or misspecified reduced forms – in the sense that level is controlled. Indeed, we provided a general characterization of the structure of the test statistics which allows one to perform exact Monte Carlo tests under general parametric distributional assumptions, which are in no way restricted to the Gaussian case, including heavy-tailed distributions without moments. The tests so obtained are exact even in cases where identification fails (or is weak) and conventional asymptotic theory breaks down.

After proving a general invariance property, we provided a characterization of the power of the tests that clearly exhibits the factors which determine power. We showed that exogeneity tests have no power in the extreme case where all IVs are weak [similar to Staiger and Stock (1997), and Guggenberger (2010)], but typically have power as soon as we have one strong instrument. Consequently, exogeneity tests can detect an exogeneity problem even if not all model parameters
are identified, provided at least some parameters are identifiable.

Though the exact distributional theory given in this paper requires relatively specific distributional assumptions, the “finite-sample” procedures provided remain asymptotically valid in the same way (in the sense that test level is controlled) under standard asymptotic assumptions. We study this problem in a separate paper [Doko Tchatoka and Dufour (2016)]. Further, even if exogeneity hypotheses can have economic interest by themselves, we also show there how exogeneity tests can be fruitfully applied to build pretest estimators which generally dominate OLS and 2SLS estimators when the exogeneity of explanatory variables is in uncertain.


APPENDIX

A. Wu and Hausman test statistics

We show here that Durbin-Wu statistics can be expressed in the same way as alternative Hausman statistics. The statistics $T_l, l = 1, 2, 3, 4$ are defined in Wu (1973, eqs. (2.1), (2.18), (3.16), and (3.20)) as:

$$T_1 = \frac{Q^*}{Q_1}, \quad T_2 = \frac{Q^*}{Q_2}, \quad T_3 = \frac{Q^*}{Q_3}, \quad T_4 = \frac{Q^*}{Q_4}, \quad (A.1)$$

$$Q^* = (b_1 - b_2)' \left[ (Y' A_2 Y)^{-1} - (Y' A_1 Y)^{-1} \right]^{-1} (b_1 - b_2), \quad (A.2)$$

$$Q_1 = (y - Y b_2)' A_2 (y - Y b_2), \quad Q_2 = Q_4 - Q^*, \quad (A.3)$$

$$Q_4 = (y - Y b_1)' A_1 (y - Y b_1), \quad Q_3 = (y - Y b_2)' A_1 (y - Y b_2), \quad (A.4)$$

$$b_1 = (Y' A_i Y)^{-1} Y' A_i y, \quad i = 1, 2, A_1 = M_1, A_2 = M - M_1, \quad (A.5)$$

where $b_1$ is the ordinary least squares estimator of $\beta$, and $b_2$ is the instrumental variables method estimator of $\beta$. So, in our notations, $b_1 \equiv \hat{\beta}$ and $b_2 \equiv \tilde{\beta}$. From (3.8) - (3.13), we have:

$$Q^* = T (\hat{\beta} - \tilde{\beta})' \hat{\Sigma}^{-1} (\hat{\beta} - \tilde{\beta}), \quad (A.6)$$

$$Q_1 = T \hat{\sigma}^2, \quad Q_3 = T \tilde{\sigma}^2, \quad Q_4 = T \hat{\sigma}^2, \quad (A.7)$$

$$Q_2 = Q_4 - Q^* = T \hat{\sigma}^2 - T (\hat{\beta} - \tilde{\beta})' \hat{\Sigma}^{-1} (\hat{\beta} - \tilde{\beta}) = T \tilde{\sigma}^2. \quad (A.8)$$

Hence, we can write $T_l$ as:

$$T_l = \kappa_l (\hat{\beta} - \tilde{\beta})' \hat{\Sigma}^{-1} (\hat{\beta} - \tilde{\beta}), \quad l = 1, 2, 3, 4,$$

where $\kappa_l$ and $\hat{\Sigma}_l$ are defined in (3.8) - (3.13).

To obtain (3.17), set $T_0 = (\hat{\beta} - \tilde{\beta})' \hat{\Sigma}^{-1} (\hat{\beta} - \tilde{\beta})$. Then $\hat{\sigma}^2 = \sigma^2 - T_0$, $T_4 = \kappa_4 T_0 / \hat{\sigma}^2$, and

$$T_2 = \kappa_2 \frac{\hat{\sigma}^2}{\hat{\sigma}^2 - T_0} = \frac{\hat{\sigma}^2}{\hat{\sigma}^2 - T_0} = \kappa_2 \frac{(T_0 / \hat{\sigma}^2)}{1 - (T_0 / \hat{\sigma}^2)} = \kappa_2 \left( \frac{T_4 / \kappa_4}{1 - (T_4 / \kappa_4)} \right), \quad (A.9)$$

hence

$$\frac{T_4}{\kappa_4} = \frac{T_2 / \kappa_2}{T_2 / \kappa_2 + 1} = \frac{T_2}{T_2 + \kappa_2} = \frac{1}{(T_2 / \kappa_2 + 1)}. \quad (A.10)$$

In the sequel of this appendix, we shall use the following matrix formulas which are easily established by algebraic manipulations [on the invertibility of matrix differences, see Harville (1997, Theorem 18.2.4)].

Lemma A.1 DIFFERENCE OF MATRIX INVERSES. Let $A$ and $B$ be two nonsingular $r \times r$ matrices.
\[ A^{-1} - B^{-1} = B^{-1}(B - A)A^{-1} = A^{-1}(B - A)B^{-1} \]
\[ = A^{-1}(A - AB^{-1}A)A^{-1} = B^{-1}(BA^{-1}B - B)B^{-1}. \] 
(A.11)

Furthermore, \( A^{-1} - B^{-1} \) is nonsingular if and only if \( B - A \) is nonsingular. If \( B - A \) is nonsingular, we have:

\[ (A^{-1} - B^{-1})^{-1} = A(B - A)^{-1}B = A - A(B - A)^{-1}A = A + A(B - A)^{-1}A = A[A^{-1} + (B - A)^{-1}]A \]
\[ = B(B - A)^{-1}A = B(B - A)^{-1}B - B = B[(B - A)^{-1} - B^{-1}]B \]
\[ = A(A - AB^{-1}A)^{-1}A = B(BA^{-1}B - B)^{-1}B. \] 
(A.12)

It is easy to see from condition (2.6) that \( \hat{\Omega}_{IV}, \hat{\Omega}_{LS} \) and \( \hat{\Sigma}_{V} \) are nonsingular. On setting \( A = \hat{\Omega}_{IV} \) and \( B = \hat{\Omega}_{LS} \), we get:

\[ B - A = \hat{\Omega}_{LS} - \hat{\Omega}_{IV} = \frac{1}{T} Y'M_1 Y - \frac{1}{T} Y'N_1 Y = \frac{1}{T} Y'(M_1 - N_1)Y = \frac{1}{T} Y'MY = \frac{1}{T} \hat{\Sigma}^{V} = \hat{\Sigma}_{IV}, \] 
(A.13)

so \( \hat{\Omega}_{LS} - \hat{\Omega}_{IV} \) is nonsingular. By Lemma A.1, \( \hat{\Delta} = \hat{\Omega}_{IV}^{-1} - \hat{\Omega}_{LS}^{-1} = A^{-1} - B^{-1} \) is also nonsingular, and

\[ \hat{\Delta}^{-1} = A + A(B - A)^{-1}A = \hat{\Omega}_{IV} + \hat{\Omega}_{IV}(\hat{\Omega}_{LS} - \hat{\Omega}_{IV})^{-1}\hat{\Omega}_{IV} = \hat{\Omega}_{IV} + \hat{\Omega}_{IV} \hat{\Sigma}^{-1}_{V} \hat{\Omega}_{IV} \]
\[ = \frac{1}{T} \left[ Y'N_1 Y + Y'N_1 Y(Y'MY)^{-1}Y'N_1 Y \right] = \frac{1}{T} Y'N_1 \left[ I_T + Y(Y'MY)^{-1}Y' \right] N_1 Y. \] 
(A.14)

From the above form, it is clear that \( \hat{\Delta}^{-1} \) is positive definite. Note also that

\[ \hat{\Delta}^{-1} = B(B - A)^{-1}B - B = \hat{\Omega}_{LS}(\hat{\Omega}_{LS} - \hat{\Omega}_{IV})^{-1}\hat{\Omega}_{LS} - \hat{\Omega}_{LS} = \hat{\Omega}_{LS} \hat{\Sigma}^{-1}_{V} \hat{\Omega}_{LS} - \hat{\Omega}_{LS} \]
\[ = \frac{1}{T} \left[ (Y'M_1 Y)(Y'MY)^{-1}(Y'M_1 Y) - (Y'M_1 Y) \right] = \frac{1}{T} Y'M_1 [Y(Y'MY)^{-1}Y' - I_T]M_1 Y. \] 
(A.15)

The latter shows that \( \hat{\Delta}^{-1} \) only depends on the least-squares residuals \( M_1 Y \) and \( MY \).

**B. Regression interpretation of DWH test statistics**

Let us now consider the regressions (3.22) - (3.25). Using \( Y = \hat{Y} + \hat{\Sigma}, \hat{Y} = X\hat{\Pi} \) and \( \hat{\Pi} = (X'X)^{-1}X'Y \), we see that the 2SLS residual vector \( \tilde{u} \) for model (2.1) based on the instrument matrix \( X = [X_1, X_2] \) can be written as

\[ \tilde{u} = y - Y\hat{\beta} - X_1\hat{\gamma} = (y - \hat{Y}\hat{\beta} - X_1\hat{\gamma}) - \hat{\Sigma}\hat{\beta} = M_1(y - \hat{Y}\hat{\beta}) - \hat{\Sigma}\hat{\beta} \]
\[ = M_1(y - \hat{Y}\hat{\beta} - \hat{\Sigma}\hat{\beta}) = M_1(y - Y\hat{\beta}) \] 
(B.1)
where  $\bar{\beta}$ and $\bar{\gamma}$ are the 2SLS estimators of $\beta$ and $\gamma$, and the different sum-of-squares functions satisfy:

\[
S(\hat{\theta}) = S_s(\hat{\theta}_*) , \quad \bar{u}'u = S(\hat{\theta}^0) = S_s(\hat{\theta}_0^0) = \bar{S}(\hat{\theta}_*^0) , \quad \bar{S}(\hat{\theta}_*) = (y - Y\bar{\beta})'M(y - Y\bar{\beta}) , \quad \text{(B.2)} \\
S(\hat{\theta}^0) - S(\hat{\theta}) = S_s(\hat{\theta}_0^0) - S_s(\hat{\theta}_*) . \quad \text{(B.3)}
\]

Let $R = \begin{bmatrix} 0 & 0 & I_G \end{bmatrix}$, and $R_* = \begin{bmatrix} I_G & 0 & -I_G \end{bmatrix}$, so that $Rb = a$ and $R_\theta = \beta - a$. The null hypotheses $H_0 : a = 0$ and $H_0^* : \beta = b$ can thus be written as

\[
H_0 : R\theta = 0 , \quad H_0^* : R_*\theta = 0 . \quad \text{(B.4)}
\]

Further, $\hat{\theta}_* = [\hat{\beta}', \hat{\gamma}', \hat{\beta}']'$ and $\hat{\theta}_0^0 = [\bar{\beta}', \bar{\gamma}', \bar{\beta}']'$, where $\hat{\beta}$ and $\hat{\gamma}$ are the OLS estimators of $\beta$ and $\gamma$ based on the model (2.1), and

\[
R_*\hat{\theta} = \begin{bmatrix} I_G & 0 & -I_G \end{bmatrix} \begin{bmatrix} \hat{\beta} \\ \hat{\gamma} \\ \hat{\beta} \end{bmatrix} = \bar{\beta} - \bar{\beta} , \quad \text{(B.5)}
\]

\[
\hat{\theta}_0^0 = \hat{\theta}_* + (Z'_*Z_*)^{-1}R'_* [R_* (Z'_*Z_*)^{-1} R'_*]^{-1} (-R_* \hat{\theta}_*) , \quad \text{(B.6)}
\]

\[
S(\hat{\theta}_*) = (\hat{\theta}_0^0 - \hat{\theta}_*)'Z'_*Z_*(\hat{\theta}_0^0 - \hat{\theta}_*) = (R_* \hat{\theta}_*')' [R_* (Z'_*Z_*)^{-1} R'_*]^{-1} (R_* \hat{\theta}_*) , \quad \text{(B.7)}
\]

where $Z_* = [\hat{Y}, X_1, \hat{V}]$. On writing $Z_* = [\hat{X}_1, \hat{V}]$, where $\hat{X}_1 = [\hat{Y}, X_1]$, we get:

\[
Z'_*Z_* = \begin{bmatrix} (\hat{X}_1'\hat{X}_1) & 0 \\ 0 & (\hat{V}'\hat{V}) \end{bmatrix} , \quad (Z'_*Z_*)^{-1} = \begin{bmatrix} (\hat{X}_1'\hat{X}_1)^{-1} & 0 \\ 0 & (\hat{V}'\hat{V})^{-1} \end{bmatrix} , \quad \text{(B.8)}
\]

\[
(\hat{X}_1'\hat{X}_1)^{-1} = \begin{bmatrix} \hat{P}'\hat{P} \\ X_1'\hat{P} X_1' \end{bmatrix}^{-1} = \begin{bmatrix} W_{YY} & W_{Y1} \\ W_{Y1} & W_{11} \end{bmatrix} , \quad \text{(B.9)}
\]

where $W_{YY} = [(\hat{P}'\hat{P}) - \hat{Y}'X_1(X_1'X_1)^{-1}X_1'\hat{P}]^{-1} = \hat{Y}'M_1\hat{Y}^{-1} = [\hat{Y}'(M_1 - M)\hat{Y}]^{-1}$,

\[
(Z'_*Z_*)^{-1}R'_* = \begin{bmatrix} W_{YY} & W_{Y1} & 0 \\ W_{Y1} & W_{11} & 0 \\ 0 & 0 & (\hat{V}'\hat{V})^{-1} \end{bmatrix} \begin{bmatrix} I_G \\ 0 \\ -I_G \end{bmatrix} = \begin{bmatrix} W_{YY} \\ W_{11} \\ -(\hat{V}'\hat{V})^{-1} \end{bmatrix} , \quad \text{(B.10)}
\]

\[
R_* (Z'_*Z_*)^{-1} R'_* = W_{YY} + (\hat{V}'\hat{V})^{-1} , \quad \text{(B.11)}
\]

\[
\hat{\theta}_0^0 - \hat{\theta}_* = \begin{bmatrix} \hat{\beta} - \bar{\beta} \\ \hat{\gamma} - \bar{\gamma} \\ \hat{\beta} - \bar{\beta} \end{bmatrix} = \begin{bmatrix} W_{YY} \\ W_{11} \\ -(\hat{V}'\hat{V})^{-1} \end{bmatrix} [W_{YY} + (\hat{V}'\hat{V})^{-1}]^{-1} (\bar{\beta} - \bar{\beta}) . \quad \text{(B.12)}
\]

From the latter equation, we see that

\[
\hat{\beta} - \bar{\beta} = W_{YY} [W_{YY} + (\hat{V}'\hat{V})^{-1}]^{-1} (\bar{\beta} - \bar{\beta}) = W_{YY} [W_{YY} + (\hat{V}'\hat{V})^{-1}]^{-1} \tilde{a} , \quad \text{(B.13)}
\]
where $\bar{a} = \bar{b} - \bar{\beta}$ is the OLS estimate of $a$ in (3.23). Hence, we have

\[
\bar{a} = b - \bar{\beta} = \left[ W_{YY} + (\hat{V}'\hat{V})^{-1} \right] W_{YY}^{-1} (\hat{b} - \bar{\beta}) = \left\{ (\hat{V}'(M_1 - M)Y)^{-1} + (\hat{V}'\hat{V})^{-1} \right\} \left[ (\hat{V}'(M_1 - M)Y) (\hat{b} - \bar{\beta}) \right],
\]

which entails that

\[
S(\hat{0}) - S(\hat{\theta}_s) = (R, \hat{\theta}_s)' [R, (Z'Z)^{-1} R']^{-1} (R, \hat{\theta}_s) = (b - \bar{\beta}')(\hat{V}'(M_1 - M)Y)^{-1} + (\hat{V}'\hat{V})^{-1} (b - \bar{\beta}) = (\hat{b} - \bar{\beta})' \left\{ \hat{V}'(M_1 - M)Y \right\} (\hat{b} - \bar{\beta}) = (\hat{b} - \bar{\beta})' W_{YY}^{-1} [W_{YY} + (Y'M_1Y - W_{YY}^{-1})] W_{YY}^{-1} (\hat{b} - \bar{\beta}).
\]

Using Lemma A.1 with $A = W_{YY}^{-1}$ and $B = Y'M_1Y$ in (B.15), we then get:

\[
S(\hat{0}) - S(\hat{\theta}_s) = (\hat{b} - \bar{\beta})' W_{YY}^{-1} [W_{YY} + (Y'M_1Y - W_{YY}^{-1})] W_{YY}^{-1} (\hat{b} - \bar{\beta}) = T(\hat{b} - \bar{\beta})' (\hat{\theta}_s - \hat{\theta}_s) = T(\hat{b} - \bar{\beta})' (\hat{\theta}_s - \hat{\theta}_s).
\]

where $\hat{\theta}_s = \frac{1}{T} Y'(M_1 - M)Y$ and $\hat{\theta}_s = \frac{1}{T} Y'M_1Y$. Since we have $S(\hat{0}) - S(\hat{\theta}_s) = S(\hat{0}) - S(\hat{\theta}_s)$, we get from (B.16), (3.13) and (3.30):

\[
S(\hat{0}) = S(\hat{0}) - [S(\hat{0}) - S(\hat{\theta}_s)] = S(\hat{0}) - T(\hat{\beta} - \bar{\beta})' (\hat{\theta}_s - \hat{\theta}_s) = T\hat{\sigma}_2^2.
\]

It is also clear from (3.13) and (3.30) that

\[
S(\hat{0}) = T\hat{\sigma}_2^2, S(\hat{0}) = T\hat{\sigma}_2^2.
\]

Hence, except for $R_1$, the other statistics can be expressed as:

\[
\rho_2 = T \frac{S(\hat{0}) - S(\hat{\theta})}{S(\hat{\theta})}, \quad \rho_3 = T \frac{S(\hat{0}) - S(\hat{\theta})}{S(\hat{\theta})},
\]

\[
\rho_4 = \kappa_1 \left( \frac{S(\hat{0}) - S(\hat{\theta})}{S(\hat{\theta})} \right), \quad \rho_5 = \kappa_1 \left( \frac{S(\hat{0}) - S(\hat{\theta})}{S(\hat{\theta})} \right),
\]

\[
\rho_6 = \kappa_2 \left( \frac{S(\hat{0}) - S(\hat{\theta})}{S(\hat{\theta})} \right), \quad \rho_7 = \kappa_3 \left( \frac{S(\hat{0}) - S(\hat{\theta})}{S(\hat{\theta})} \right), \quad \rho_8 = \kappa_4 \left( \frac{S(\hat{0}) - S(\hat{\theta})}{S(\hat{\theta})} \right),
\]

\[
\rho_{R} = \kappa_R \left( \frac{S(\hat{0}) - S(\hat{\theta})}{S(\hat{\theta})} \right).
\]
C. Proofs

To establish Proposition 4.1, it will be useful to state some basic identities for the different components of alternative exogeneity test statistics.

**Lemma C.1** PROPERTIES OF EXOGENEITY STATISTICS COMPONENTS. The random vectors and matrices in (3.1) - (3.14) satisfy the following identities: setting

\[ B_1 =: (Y'M_1Y)^{-1}Y'M_1, \quad B_2 =: (Y'N_1Y)^{-1}Y'N_1, \]  
\[ C_1 =: B_2 - B_1, \quad \Psi_0 =: C_1^{-1}C_1, \quad N_2 =: I_T - M_1YA_2, \]  
we have

\[ B_1M_1 = B_1, \quad B_2M_1 = B_2N_1 = B_2, \quad B_1Y = B_2Y = I_G, \]  
\[ C_1Y = 0, \quad C_1X_1 = 0, \quad C_1\tilde{P}[M_1Y] = 0, \quad C_1M_1 = C_1\tilde{M}[M_1Y] = C_1, \]  
\[ M_1YA_1 = \tilde{P}[M_1Y], \quad M_1\Psi_0M_1 = M_1\Psi_0 = \Psi_0M_1 = \Psi_0, \]  
\[ M_1\Psi R M_1 = \Psi_R, \quad M_1\Lambda M_1 = M_1\Lambda M = \Lambda, \]  
\[ B_1B_1' = B_1B_2' = B_2B_2' = \frac{1}{T}\hat{\Omega}^{-1}_{IV}, \quad B_2B_2' = \frac{1}{T}\hat{\Omega}^{-1}_{LS}, \]  
\[ C_1C_1' = \frac{1}{T}(\hat{\Omega}^{-1}_{IV} - \hat{\Omega}^{-1}_{LS}), \quad C_1\Psi_0 = \frac{1}{T}C_1, \quad \Psi_0\Psi_0 = \frac{1}{T}\Psi_0, \]  
\[ \bar{\beta} - \hat{\beta} = (B_2 - B_1)y = C_1y = C_1(M_1y), \]  
\[ (\bar{\beta} - \hat{\beta})'\hat{\Delta}^{-1}(\bar{\beta} - \hat{\beta}) = y'\Psi_0y = (M_1y)'\Psi_0(M_1y), \]  
\[ y - \hat{\beta}\hat{\beta}' = [I_T - YB_1]y, \quad y - \hat{\beta}\hat{\beta}' = [I_T - YB_2]y, \]  
\[ \tilde{u} = M_1(y - \hat{\beta}) = \tilde{M}[\tilde{y}]y = M\tilde{M}[M_1Y]y = \tilde{M}[M_1Y](M_1y), \]  
\[ M(y - \hat{\beta}) = M\tilde{M}[M_1Y]y = M\tilde{M}[M_1Y](M_1y), \]  
\[ N_1(y - \hat{\beta}) = M_1P(y - \hat{\beta}) = M_1\tilde{M}[M_1PY]P\tilde{y} = \tilde{M}[N_1Y]N_1y = PM_1(y - \hat{\beta}) = \tilde{M}[PM_1Y]P(M_1y), \]  
\[ \tilde{u} = M_1(y - \hat{\beta}) = N_2(M_1y), \quad M(y - \hat{\beta}) = MN_2(M_1y), \]  
\[ \tilde{\sigma}^2 = \frac{1}{T}(M_1y)'N_2N_2(M_1y), \]  
\[ \tilde{\sigma}^2 = \frac{1}{T}y'\tilde{N}[\tilde{y}]y = \frac{1}{T}y'M_1\tilde{M}[M_1Y]y = \frac{1}{T}y'M_1\tilde{M}[M_1Y](M_1y), \]  
\[ \tilde{\sigma}^2 = \frac{1}{T}y'[N_1\tilde{M}[N_1Y]N_1y = \frac{1}{T}(M_1y)'PM_1\tilde{M}[PM_1Y]P(M_1y), \]  

37
\[ \hat{\sigma}_2^2 = (M_1 y)' \left\{ \frac{1}{T} \tilde{M}[M_1 y] - \Psi_0 \right\} (M_1 y), \]  
(C.19)

\[ y' \Psi_R y = \frac{1}{T} y' \tilde{P}[\tilde{M}[\tilde{Y}] X_2] \tilde{M}[\tilde{Y}] y = \frac{1}{T} (M_1 y)' \tilde{P}[\tilde{M}[\tilde{Y}] X_2] (M_1 y), \]  
(C.20)

\[ \hat{\sigma}_R^2 = \frac{1}{T} y' \tilde{M}[Z] y = \frac{1}{T} (M_1 y)' \tilde{M}[Z] (M_1 y). \]  
(C.21)

**Proof of Lemma C.1** Using the idempotence of \( M_1 \) and (3.15), we see that:

\[ B_1 M_1 = (Y' M_1 Y)^{-1} Y' M_1 = (Y' M_1 Y)^{-1} Y' M_1 = B_1, \]  
(C.22)

\[ B_2 M_1 = [Y' N_1 Y]^{-1} Y' N_1 M_1 = [Y' N_1 Y]^{-1} Y' N_1 = B_2 = B_2 N_1 = B_2 (M_1 - M), \]  
(C.23)

\[ M_1 Y A_1 = M_1 Y (Y' M_1 Y)^{-1} Y' M_1 = \tilde{P}(M_1 Y), \]  
(C.24)

\[ C_1 M_1 = B_2 M_1 - B_1 M_1 = B_2 - B_1 = C_1, \quad C_1 X_1 = C_1 M_1 X_1 = 0, \]  
(C.25)

\[ B_1 Y = (Y' M_1 Y)^{-1} Y' M_1 Y = I_G = (Y' N_1 Y)^{-1} Y' N_1 Y = B_2 Y, \]  
(C.26)

\[ C_1 Y = B_2 Y - B_1 Y = 0, \]  
(C.27)

\[ C_1 \tilde{P}[M_1 Y] = \left[ (Y' N_1 Y)^{-1} Y' N_1 - (Y' M_1 Y)^{-1} Y' M_1 \right] M_1 Y (Y' M_1 Y)^{-1} Y' M_1 \]  
\[ = \left[ (Y' N_1 Y)^{-1} Y' N_1 - (Y' M_1 Y)^{-1} Y' M_1 \right] (Y' M_1 Y)^{-1} Y' M_1 \]  
\[ = (I_G - I_G) (Y' M_1 Y)^{-1} Y' M_1 = 0, \]  
(C.28)

\[ C_1 \tilde{M}[M_1 Y] = C_1 [I_T - \tilde{P}[M_1 Y] = C_1, \]  
(C.29)

\[ M_1 \tilde{M}[\tilde{Y}] M_1 = \tilde{M}[\tilde{Y}], \quad M_1 \tilde{M}[Z] M_1 = \tilde{M}[Z], \]  
(C.30)

\[ M_1 \Psi_R M_1 = \frac{1}{T} \left( \{M_1 \tilde{M}[\tilde{Y}] M_1 - M_1 \tilde{M}[Z] M_1 \} \right) = \Psi_R, \quad M_1 \Lambda_\alpha M_1 = \frac{1}{T} M_1 \tilde{M}[Z] M_1 = \Lambda_\alpha, \]  
(C.31)

so (C.3) - (C.6) are established. (C.7) and (C.8) follow directly from (3.15) and the definitions of \( B_1, B_2, C_1 \) and \( \Psi_0 \). We get (C.9) and (C.10) by using the definitions of \( \hat{\beta} \) and \( \tilde{\beta} \) in (3.4) - (3.5). (C.11) follows on using (3.4) and (3.5). (C.12) comes from the fact that the residuals \( M_1 (y - \hat{Y}) \) are obtained by minimizing \( ||y - \tilde{Y} - X_1 \gamma||^2 \) with respect to \( \gamma \), or equivalently \( ||y - \tilde{Y} \beta - X_1 \gamma||^2 \) with respect to \( \beta \) and \( \gamma \). (C.13) follows from (C.12) and noting that \( M = M_1 M_1 \). Similarly, the first identity in (C.14) comes from the fact that the residuals \( M_1 P (y - \hat{Y}) = M_1 (y - \tilde{Y} \beta) \) are obtained by minimizing \( ||y - \tilde{Y} \beta - X_1 \gamma||^2 \) with respect to \( \gamma \), or equivalently minimizing \( ||y - \tilde{Y} \beta - X_1 \gamma||^2 \) with respect to \( \beta \) and \( \gamma \). The others follow on noting that \( N_1 = M_1 P = PM_1 \) and

\[ M_1 \tilde{M}[M_1 P Y] = \tilde{M}[PM_1 Y] M_1 P = \tilde{M}[PM_1 Y] PM_1. \]  
(C.32)

To get (C.15) and (C.16), we note that

\[ \hat{u} = y - \tilde{Y} \beta - X_1 \gamma = M_1 (y - \tilde{Y} \beta) = M_1 [I_T - YA_2] y = [I_T - M_1 Y A_2] (M_1 y) = N_2 (M_1 y) \]  
(C.33)
hence
\[ \hat{\sigma}^2 = \frac{1}{T} \bar{u}' \bar{u} = \frac{1}{T} (y - Y \hat{\beta})' M_1 (y - Y \hat{\beta}) = \frac{1}{T} (M_1 y)' N_2' N_2 (M_1 y). \]  

(C.34)

Further, using (3.11) - (3.3), (C.12) and (C.14), we see that:

\[ \hat{\sigma}^2 = \frac{1}{T} (y - Y \hat{\beta})' M_1 (y - Y \hat{\beta}) = \frac{1}{T} y' \hat{\beta} [\hat{\beta}]' y = \frac{1}{T} y' M_1 \hat{\beta} [M_1 y]' y = \frac{1}{T} (M_1 y)' \hat{\beta} [M_1 y] (M_1 y), \]

(C.35)

\[ \sigma_1^2 = \frac{1}{T} (y - Y \hat{\beta})' N_1 (y - Y \hat{\beta}) = \frac{1}{T} (y - Y \hat{\beta})' P M_1 P (y - Y \hat{\beta}) \]

\[ = \frac{1}{T} y' N_1 [\hat{\beta} M_1 y] N_1 y = \frac{1}{T} (M_1 y)' \hat{\beta} [P M_1 y] P (M_1 y), \]

(C.36)

\[ \sigma_2^2 = (M_1 y)' \left\{ \frac{1}{T} \hat{\beta} [M_1 y] - \Psi_0 \right\} (M_1 y), \]

(C.37)

so (3.11) - (3.13) are established. Finally, (C.20) and (C.21) follow by observing that

\[ M_1 [\hat{\beta} M_1 y] = \hat{\beta} [\hat{\beta} M_1 y] M_1 \] and \[ M_1 [\hat{\beta} M_1 y] = \hat{\beta} [\hat{\beta} M_1 y] M_1 \] and \[ M_1 M_1 y] = \hat{\beta} [\hat{\beta} M_1 y] and \[ M_1 M_1 y] = \hat{\beta} [\hat{\beta} M_1 y]. \]

Using Lemma C.1, we can now prove Proposition 4.1.

**Proof of Proposition 4.1** We first note that

\[ \tilde{\beta} - \hat{\beta} = (B_2 - B_1) y = C_1 y, \]

(C.38)

\[ (\tilde{\beta} - \hat{\beta})' \hat{\beta}^{-1} (\tilde{\beta} - \hat{\beta}) = y' C_1 \hat{\beta}^{-1} C_1 y = y' \Psi_0 y, \]

(C.39)

so that, by the definitions (3.1) - (3.3),

\[ \mathcal{S}_i = \kappa_i (\tilde{\beta} - \hat{\beta})' \hat{\beta}^{-1} (\tilde{\beta} - \hat{\beta}) = \frac{\kappa_i (\tilde{\beta} - \hat{\beta})' \hat{\beta}^{-1} (\tilde{\beta} - \hat{\beta})}{\sigma_i^2} = \frac{y' \Psi_0 y}{\sigma_i^2}, \quad i = 1, 2, 3, 4, \]

(C.40)

\[ \mathcal{H}_i = T (\tilde{\beta} - \hat{\beta})' \hat{\beta}^{-1} (\tilde{\beta} - \hat{\beta}) = T \frac{(\tilde{\beta} - \hat{\beta})' \hat{\beta}^{-1} (\tilde{\beta} - \hat{\beta})}{\sigma_i^2} = \frac{y' \Psi_0 y}{\sigma_i^2}, \quad i = 2, 3, \]

(C.41)

where, using Lemma C.1,

\[ \sigma_1^2 = \frac{1}{T} (y - Y \hat{\beta})' N_1 (y - Y \hat{\beta}) = \frac{1}{T} y' N_1 \hat{\beta} M_1 y N_1 y = y' \Lambda_1 y, \]

(C.42)

\[ \sigma_2^2 = y' M_1 \left\{ \frac{1}{T} \hat{\beta} M_1 y - \Psi_0 \right\} (M_1 y) = y' \Lambda_2 y, \]

(C.43)
Lemma C.2

For $\mathcal{H}_1$, we have

$$\mathcal{H}_1 = T(\hat{\beta} - \hat{\beta}_0)' \hat{\Sigma}_1^{-1}(\hat{\beta} - \hat{\beta}_0) = T y' C_1 \hat{\Sigma}_1^{-1} C_1 y = T (y' \Psi_0[y] y)$$

where

$$\hat{\Sigma}_1 = \sigma^2 \hat{\Omega}_I^{-1} - \sigma^2 \hat{\Omega}_LS^{-1} = (y' \Lambda_3 y) \hat{\Omega}_I^{-1} - (y' \Lambda_4 y) \hat{\Omega}_LS^{-1}.$$  \hfill (C.48)

The result for $\mathcal{H}$ follows directly by using (3.3).

In order to characterize the null distributions of the test statistics (Theorem 4.2), it will be useful to first spell out some algebraic properties of the weighting matrices in Proposition 4.1. This is done by the following lemma.

**Lemma C.2** PROPERTIES OF WEIGHTING MATRICES IN EXOGONEITY STATISTICS. The matrices $\Psi_0$, $\Lambda_1$, $\Lambda_2$, $\Lambda_4$, $\Psi_R$ and $\Lambda_x$ in (4.1) - (4.6) satisfy the following identities:

$$\Lambda_2 = \Lambda_4 - \Psi_0, \quad C_1 \Lambda_1 = C_1 \Lambda_2 = \Psi_0 \Lambda_1 = \Psi_0 \Lambda_2 = \Psi_R \Lambda_x = 0,$$ \hfill (C.49)

$$C_1 \Lambda_4 = \frac{1}{T} C_1, \quad \Psi_0 \Lambda_4 = \frac{1}{T} \Psi_0,$$ \hfill (C.50)

$$M_1 \Lambda_i M_1 = \Lambda_i, \quad i = 1, \ldots, 4.$$ \hfill (C.51)

Further, the matrices $T \Psi_0$, $T \Lambda_1$, $T \Lambda_2$, $T \Lambda_4$, $T \Psi_R$ and $T \Lambda_x$ are symmetric idempotent.

**Proof of Lemma C.2** To get (C.49) - (C.50), we observe that:

$$\Lambda_2 = M_1 \left( \frac{1}{T} \bar{M}[M_1 Y] - \Psi_0 \right) M_1 = \Lambda_4 - M_1 \Psi_0 M_1 = \Lambda_4 - \Psi_0,$$ \hfill (C.52)

$$C_1 \Lambda_1 \bar{P}[N] = \frac{1}{T} [B_2 - B_1] N_1 N_1 Y \hat{\Omega}_IV^{-1} Y' N_1 = \frac{1}{T} [\hat{\Omega}_IV^{-1} Y' N_1 - \hat{\Omega}_LS^{-1} Y' M_1] N_1 Y \hat{\Omega}_IV^{-1} Y' N_1$$

$$= \frac{1}{T} [\hat{\Omega}_IV^{-1} Y' N_1 Y \hat{\Omega}_IV^{-1} Y' - \hat{\Omega}_LS^{-1} Y' N_1 Y \hat{\Omega}_IV^{-1} Y'] N_1 = \frac{1}{T} [\hat{\Omega}_IV^{-1} Y' - \hat{\Omega}_LS^{-1} Y'] N_1$$

$$= \frac{1}{T} [\hat{\Omega}_IV^{-1} Y' N_1 - \hat{\Omega}_LS^{-1} Y' M_1] N_1 = [B_2 - B_1] N_1 = C_1 \Lambda_1,$$ \hfill (C.53)

$$C_1 \Lambda_4 \bar{P}[N] = C_1 M_1 Y (Y' M_1 Y)^{-1} Y' M_1 = 0,$$ \hfill (C.54)

$$\bar{\bar{M}}[\bar{Y}] \bar{\bar{M}}[Z] = \bar{\bar{M}}[Z],$$ \hfill (C.55)

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hence
\[ C_1A_1 = C_1 \left( \frac{1}{T}N_1 \tilde{M}[N_1 Y] N_1 \right) = \frac{1}{T} C_1 N_1 \tilde{M}[N_1 Y] N_1 = \frac{1}{T} C_1 N_1 (I_T - \tilde{P}[N_1 Y]) N_1 = 0, \]  
(C.56)
\[ C_1A_2 = C_1M_1 \left( \frac{1}{T} \tilde{M}[M_1 Y] - \Psi_0 \right) M_1 = \frac{1}{T} C_1 M_1 \tilde{M}[M_1 Y] M_1 - C_1 \Psi_0 M_1 = \frac{1}{T} C_1 - \frac{1}{T} C_1 = 0, \]  
(C.57)
\[ C_1A_4 = \frac{1}{T} C_1 M_1 \tilde{M}[M_1 Y] M_1 = \frac{1}{T} C_1 M_1 \tilde{M}[M_1 Y] = \frac{1}{T} C_1, \]  
(C.58)
\[ \Psi_0 A_4 = \frac{1}{T} C'_1 \hat{\Delta}^{-1} C_1 M_1 \tilde{M}[M_1 Y] M_1 = \frac{1}{T} C'_1 \hat{\Delta}^{-1} C_1 M_1 \tilde{M}[M_1 Y] = \frac{1}{T} C'_1 \hat{\Delta}^{-1} C_1 = \frac{1}{T} \Psi_0, \]  
(C.59)
\[ \Psi_0 A_2 = \Psi_0 M_1 \left( \frac{1}{T} \tilde{M}[M_1 Y] - \Psi_0 \right) M_1 = \Psi_0 (A_4 - \Psi_0) = \frac{1}{T} \Psi_0 - \frac{1}{T} \Psi_0 = 0, \]  
(C.60)
\[ \Psi_R A_r = \frac{1}{T^2} \{ \tilde{M}[\hat{Y}] - \tilde{M}[Z] \} \tilde{M}[Z] = 0. \]  
(C.61)
(C.51) follow directly from the idempotence of \( M_1 \) and the definitions of \( A_l, l = 1, \ldots, 4 \). Finally, the idempotence and symmetry of the weight matrices can be checked as follows:
\[ (T \Psi_0')(T \Psi_0') = T C'_1 \hat{\Delta}^{-1} C_1 C'_1 \hat{\Delta}^{-1} C_1 = T^2 C'_1 \hat{\Delta}^{-1} \left( \frac{1}{T} \hat{\Delta} \right) \hat{\Delta}^{-1} C_1 = T C'_1 \hat{\Delta}^{-1} C_1 \]  
(C.62)
\[ (T A_1)(T A_1) = (N_1 \tilde{M}[N_1 Y] N_1) (N_1 \tilde{M}[N_1 Y] N_1) = N_1 \tilde{M}[N_1 Y] N_1 = T A_1 = T A_1', \]  
(C.63)
\[ (T A_4)(T A_4) = M_1 \tilde{M}[M_1 Y] M_1 M_1 \tilde{M}[M_1 Y] M_1 = M_1 \tilde{M}[M_1 Y] M_1 = T A_4 = T A_4', \]  
(C.64)
\[ (T A_2)(T A_2) = T^2 (A_4 - \Psi_0^0) (A_4 - \Psi_0^0) = T^2 (A_4 A_4 - A_4 \Psi_0^0 - \Psi_0^0 A_4 + \Psi_0^0 \Psi_0^0) = T^2 \left( \frac{1}{T} A_4 - \frac{2}{T} \Psi_0^0 + \frac{1}{T} \Psi_0^0 \right) = T (A_4 - \Psi_0^0) = T A_2 = T A_2', \]  
(C.65)
\[ (T \Psi_R)(T \Psi_R) = \{ \tilde{M}[\hat{Y}] - \tilde{M}[Z] \} \{ \tilde{M}[\hat{Y}] - \tilde{M}[Z] \} = \tilde{M}[\hat{Y}] - \tilde{M}[Z] = T \Psi_R = T \Psi_R', \]  
(C.66)
\[ (T A_r)(T A_r) = \tilde{M}[Z] \tilde{M}[Z] = \tilde{M}[Z] = T A_r = T A_r'. \]  
(C.67)

**Proof of Theorem 4.2**  Using Lemma C.1, we first note the following identities:
\[ B_1 Y = (Y' M_1 Y)^{-1} Y' M_1 Y = I_G = (Y' N_1 Y)^{-1} Y' N_1 Y = B_2 Y, \]  
(C.68)
\[ \hat{M}[M_1Y]M_1Y = \bar{M}[N_1Y]N_1Y = 0, \quad B_1X_1 = B_2X_1 = 0, \quad N_1X_1 = M_1X_1 = 0, \quad (C.69) \]

\[ N_2M_1Y = (I_T - M_1YA_2)M_1Y = (M_1 - M_1YA_2)Y = M_1(Y - YA_2Y) = 0, \quad N_2M_1X_1 = 0, \quad (C.70) \]

\[ \tilde{M}[\tilde{Y}]Y = M[Z]Y = 0, \quad \tilde{M}[\tilde{Y}]X_1 = \bar{M}[Z]X_1 = 0, \quad \tilde{P}[\tilde{M}[\tilde{Y}]X_2]\tilde{M}[\tilde{Y}] = \bar{M}[\tilde{Y}]\tilde{P}[\tilde{M}[\tilde{Y}]X_2]\tilde{M}[\tilde{Y}], \quad (C.71) \]

Then

\[ C_1Y = (B_2 - B_1)(Y + X_1Y + u) = C_1u, \quad (C.72) \]

\[ y'\Psi_Ry = y'\tilde{M}[\tilde{Y}]X_2\tilde{M}[\tilde{Y}]y = \frac{1}{T}y'\tilde{M}[\tilde{Y}]\tilde{P}[\tilde{M}[\tilde{Y}]X_2]\tilde{M}[\tilde{Y}]y = \frac{1}{T}u'\tilde{M}[\tilde{Y}]\tilde{P}[\tilde{M}[\tilde{Y}]X_2]\tilde{M}[\tilde{Y}]u = u'\Psi_Ru, \quad (C.78) \]

\[ \hat{\sigma}_R^2 = \frac{1}{T}y'M[Z]u = \frac{1}{T}u'M[Z]u. \quad (C.79) \]

Further, when \( a = 0 \), we have \( u = \sigma_1(\bar{X}) \varepsilon \), and the expressions in (4.7) - (4.8) follow from (4.1) - (4.3) in Proposition 4.1 once \( u \) is replaced by \( \sigma_1(\bar{X}) \varepsilon \) in (C.72) - (C.79). \( \sigma_1(\bar{X}) \) disappears because it can be factorized in both the numerator and the denominator of each statistic. 

**Proof of Proposition 5.1** We must study how the statistics defined in (3.1) - (3.3) change when \( Y \) and \( Y \) are replaced by \( Y^* = yR_{11} + YR_{21} \) and \( Y^* = YR_{22} \). This can be done by looking at the way the relevant variables in (3.4) - (3.14) change. We first note that

\[ \hat{\Omega}_{IV}^* = \frac{1}{T}Y'^*N_1Y^* = (YR_{22})'N_1(YR_{22}) = R_{22}'\hat{\Omega}_{IV}R_{22}, \quad \hat{\Omega}_{LS}^* = \frac{1}{T}Y'^*M_1Y^* = R_{22}'\hat{\Omega}_{LS}R_{22}, \quad (C.80) \]

hence

\[ \hat{\Delta}^* = (\hat{\Omega}_{IV}^*)^{-1} - (\hat{\Omega}_{LS}^*)^{-1} = R_{22}^{-1}(\hat{\Omega}_{IV}^{-1} - \hat{\Omega}_{LS}^{-1})(R_{22}^{-1})' = R_{22}^{-1}\hat{\Delta}(R_{22}^{-1})'. \quad (C.81) \]

Using Lemma C.1, we also get:

\[ B_1^* = (Y'^*M_1Y^*)^{-1}Y'^*M_1 = [(YR_{22})M_1(YR_{22})]^{-1}(YR_{22})'M_1 = R_{22}^{-1}(Y'M_1Y)^{-1}Y'M_1 = R_{22}^{-1}B_1, \quad (C.82) \]
\[ B_2^* = (Y^* Y_{1*})^{-1} Y^* N_1 = R_{22}^{-1} (Y^* N_1) Y^{-1} Y^* N_1 = R_{22}^{-1} B_2, \]  
\[ C_1^* = B_2^* - B_1^* = R_{22}^{-1} C_1, \quad C_1^* Y = R_{22}^{-1} C_1 Y = 0, \]  
\[ \hat{\beta}^* = B_1^* y_* = R_{22}^{-1} B_1 (y R_{11} + Y R_{21}) = R_{11} R_{22}^{-1} \hat{\beta} + R_{22}^{-1} R_{21}, \]  
\[ \tilde{\beta}^* = B_2^* y_* = R_{11} R_{22}^{-1} \tilde{\beta} + R_{22}^{-1} R_{21}, \]  
\[ \hat{\beta}^* - \tilde{\beta}^* = C_1^* y_* = R_{11} R_{22}^{-1} (\hat{\beta} - \tilde{\beta}); \]  
\[ \hat{u}^* = M_1 (y^* - Y^* \hat{\beta}^*) = M_1 (y R_{11} + Y R_{21} - Y R_{22} (R_{11} R_{22}^{-1} \hat{\beta} + R_{22}^{-1} R_{21})), \]  
\[ \hat{u}^* = R_{11} M_1 (y - Y \hat{\beta}) = R_{11} \hat{u}, \]  
\[ \hat{u}^* = M_1 (y^* - Y^* \tilde{\beta}^*) = M_1 (y R_{11} + Y R_{21} - Y R_{22} (R_{11} R_{22}^{-1} \tilde{\beta} + R_{22}^{-1} R_{21}))) = R_{11} \tilde{u}, \]  
hence, since \( N_1 X_1 = 0, \)
\[ \hat{\sigma}^2 = \frac{1}{T} \hat{u}^T \hat{u} = R_{11}^2 \hat{\sigma}_2^2, \quad \tilde{\sigma}^2 = \frac{1}{T} \tilde{u}^T \tilde{u} = R_{11}^2 \tilde{\sigma}_2^2, \]  
\[ \hat{\sigma}_{1*}^2 = \frac{1}{T} (y^* - Y^* \hat{\beta}^*)^T N_1 (y^* - Y^* \hat{\beta}^*) = \frac{1}{T} (y^* - Y^* \tilde{\beta}^* - X_1 \gamma^*)^T N_1 (y^* - Y^* \tilde{\beta}^* - X_1 \gamma^*) \]  
\[ = \frac{1}{T} \hat{u}^T N_1 \hat{u} = R_{11}^2 \hat{\sigma}_1^2, \]  
\[ \tilde{\sigma}_{2*}^2 = \hat{\sigma}_{1*}^2 - (\hat{\beta}^* - \tilde{\beta}^*) (\hat{\Delta}^*)^{-1} (\hat{\beta}^* - \tilde{\beta}^*) \]  
\[ = R_{11}^2 \hat{\sigma}_{2*}^2 - (\hat{\beta} - \tilde{\beta}) (R_{11} R_{22}^{-1})^T R_{22}^{-1} \hat{\Delta}^{-1} R_{22}^{-1} (R_{11} R_{22}^{-1}) (\hat{\beta} - \tilde{\beta}) \]  
\[ = R_{11}^2 [\hat{\sigma}_{2*}^2 - (\hat{\beta} - \tilde{\beta}) (\hat{\Delta}^{-1} (\hat{\beta} - \tilde{\beta})) = R_{11}^2 \tilde{\sigma}_{2*}^2, \]  
\[ \hat{\Sigma}_i^* = \hat{\sigma}_{i*}^2 \hat{\Delta}^* = (R_{11} \hat{\sigma}_i^2) R_{22}^{-1} \hat{\Delta} (R_{22}^{-1})^T = R_{11} R_{22}^{-1} (\hat{\sigma}_i^2 \hat{\Delta} (R_{22}^{-1})^T \]  
\[ = R_{11} R_{22}^{-1} \hat{\Sigma}_i (R_{22}^{-1})^T, \quad i = 1, 2, 3, 4, \]  
\[ \tilde{\Sigma}_j^* = R_{11} R_{22}^{-1} \tilde{\Sigma}_j (R_{22}^{-1})^T, \quad j = 1, 2, 3. \]  
It follows that the \( \mathcal{T}_i \) and \( \mathcal{H}_j \) exogeneity test statistics based on the transformed data are identical to those based on the original data:
\[ \mathcal{T}_i^* = \kappa_i (\hat{\beta}^* - \tilde{\beta}^*) (\hat{\Sigma}_i^*)^{-1} (\hat{\beta}^* - \tilde{\beta}^*) \]  
\[ = (\hat{\beta} - \tilde{\beta}) (R_{11} R_{22}^{-1})^T (R_{11} R_{22}^{-1})^{-1} (R_{11} R_{22}^{-1}) (\hat{\beta} - \tilde{\beta}) \]  
\[ = \kappa_i (\hat{\beta} - \tilde{\beta}) (\hat{\Sigma}_i^{-1} (\hat{\beta} - \tilde{\beta}) = \mathcal{T}_i, \quad i = 1, 2, 3, 4, \]  
\[ \mathcal{H}_j^* = T (\hat{\beta}^* - \tilde{\beta}^*) (\hat{\Sigma}_j^*)^{-1} (\hat{\beta}^* - \tilde{\beta}^*) \]
Finally, the invariance of the statistic $\mathcal{R}$ is obtained by observing that
\[
y^* \tilde{M}[Z^*] y^* = R_{11}^2 y^* \tilde{M}[y] y^* = R_{11}^2 y^* \tilde{M}[\tilde{y}] y^*,
\]
where $Z^* = [Y^*, X_1, X_2]$ and $Y^* = [Y^*, X_1]$, so $R_{11}^2$ cancels out in $\mathcal{R}$.

**Proof of Theorem 6.1** Since $u = Va + \sigma_1(\tilde{X}) \varepsilon$, we can use the identities (C.72) - (C.79) and replace $y$ by $Va + \sigma_1(\tilde{X}) \varepsilon$ in (4.1) - (4.4). The expressions (6.2) - (6.4) then follow through division of the numerator and denominator of each statistic by $\sigma_1(\tilde{X})$.

**Proof of Theorem 6.2** This result follows by applying the invariance property of Proposition 5.1 with $R$ defined as in (5.2). $y$ is then replaced by $y^* = X_1 \gamma + [V - g(X_1, X_2, X_3, V, \tilde{I})]a + e$ [see (5.5)], and the identities (C.72) - (C.79) hold with $u$ replaced by
\[
u_* = [V - g(X_1, X_2, X_3, V, \tilde{I})]a + e.
\]
Further, in view of (C.5) and (4.4) - (3.14), each one of the matrices $\Psi_0, \Lambda_1, \ldots, \Lambda_4, \Psi_R$ and $\Lambda_R$ remains the same if it is pre- and postmultiplied by $M_1$, i.e.
\[
\Psi_0 = M_1 \Psi_0 M_1, \quad \Lambda_1 = M_1 \Lambda_1 M_1, \quad i = 1, 2, 3, 4,
\]
\[
\Psi_1 = M_1 \Psi_1 M_1, \quad \Psi_R = M_1 \Psi_R M_1, \quad \Lambda_R = M_1 \Lambda_R M_1,
\]
so $u_*$ can in turn be replaced by
\[
M_1 u_* = -M_1 [V - g(X_1, X_2, X_3, V, \tilde{I})]a + M_1 e
\]
in (C.72) - (C.79). Upon division of the numerator and denominator of each statistic by $\sigma_1(\tilde{X})$, we get the expressions (6.6) - (6.8).

**Proof of Theorem 6.3** The result follows from well known properties of the normal and chi-square distributions: if $x \sim N_n[\mu, I_n]$ and $A$ is a fixed idempotent $n \times n$ matrix of rank $r$, then $x'Ax \sim \chi^2[r; \mu' A \mu]$. Conditional on $\tilde{X}$ and $V, \Psi_0$ is fixed, and
\[
y^* (\tilde{a}) = \tilde{\mu}_i^{-1}(\tilde{a}) + M_1 \varepsilon = M_1 \{[V - g(X_1, X_2, X_3, V, \tilde{I})]a + e\} = M_1 (\mu + e)
\]
where $\mu = [V - g(X_1, X_2, X_3, V, \tilde{I})]a$ is fixed and $\varepsilon \sim N_n[\mu, I_n]$. By Lemmas C.1 and C.2, $T \Psi_0, T \Lambda_1, T \Lambda_2, T \Lambda_4, T \Psi_R$ and $T \Lambda_R$ are symmetric idempotent, and each of these matrices remain invariant.
through by pre- and post-multiplication by \(M_1 [M_1 \Psi_0 M_1 = \Psi_0, \text{ etc.}]\). Thus

\[
\begin{align*}
S_T[y^\perp_\ast (\bar{\alpha}) \Psi_0] &= T y^\perp_\ast (\bar{\alpha})'^T \Psi_0 y^\perp_\ast (\bar{\alpha}) = (\mu + \varepsilon)'M_1 (T \Psi_0) M_1 (\mu + \varepsilon) \\
&= (\mu + \varepsilon)'(T \Psi_0)(\mu + \varepsilon) \sim \chi^2[\text{rank}(T \Psi_0) ; \text{ } \mu'(T \Psi_0) \mu]
\end{align*}
\]

where

\[
\begin{align*}
\text{rank}(T \Psi_0) &= \text{tr}(T \Psi_0) = \text{tr}(T C_1' \tilde{\Delta}^{-1} C_1) = \text{tr}(T \tilde{\Delta}^{-1} T^{-1} \tilde{\Delta}) = G, \\
\mu'(T \Psi_0) \mu &= \mu'M_1 (T \Psi_0) M_1 \mu = \bar{\alpha}'(T \Psi_0)' \bar{\alpha} = S_T[\bar{\alpha}'(\bar{\alpha}) \Psi_0] = \delta(\bar{\alpha}, \Psi_0).
\end{align*}
\]

The proofs for the other quadratic forms are similar, with the following degrees of freedom vary:

\[
\begin{align*}
\text{rank}(T A_1) &= \text{tr}\{N_1 \tilde{\Psi}[N_1 Y \tilde{N}_1]\} = \text{tr}\{N_1 \} - \text{tr}\{\hat{\Psi}[N_1 Y]\} = \text{tr}\{M_1 - M\} - \text{tr}\{N_1 Y (Y'N_1)^{-1} Y'N_1\} \\
&= (T - k_1) - (T - k_1 - k_2) - \text{tr}\{(Y'N_1)^{-1} Y'N_1\} = k_2 - G,
\end{align*}
\]

\[
\begin{align*}
\text{rank}(T A_2) &= \text{tr}\{M_1 \{T^{-1} \tilde{M}[M_1 Y \tilde{M} Y]\} M_1\} = \text{tr}\{M_1 \tilde{M}[M_1 Y \tilde{M} Y]\} - \text{tr}\{T \Psi_0\} \\
&= \text{tr}\{M_1\} - \text{tr}\{\hat{\Psi}[M_1 Y]\} - \text{tr}\{T \Psi_0\} = T - k_1 - 2G,
\end{align*}
\]

\[
\begin{align*}
\text{rank}(T A_3) &= \text{tr}\{M_1 \tilde{M}[M_1 Y]\} M_1\} = \text{tr}\{M_1\} - \text{tr}\{\hat{\Psi}[M_1 Y]\} = T - k_1 - G, \\
\text{rank}(T \Psi_R) &= \text{tr}\{\tilde{\Psi}[N] - \tilde{M}[Z]\} = (T - k_1 - G) - (T - k_1 - G - k_2) = k_2, \\
\text{rank}(T A_R) &= \text{tr}\{T A_R\} = \text{tr}\{\tilde{M}[Z]\} = T - G - k_1 - k_2.
\end{align*}
\]

The independence properties follow from the orthogonalities given in (C.49) and the normality assumption. □

**PROOF OF COROLLARY 6.4** These results directly from Theorem 6.3 and the definition of the doubly noncentral \(F\)-distribution. □
References


