# PARAMETRIC IDENTIFICATION OF STATE-SPACE DYNAMIC SYSTEMS: A TIME-DOMAIN PERSPECTIVE 

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#### Abstract

In this paper we have presented a time-domain approach to parametric identification of state-space dynamic models comprised both an equation of motion and a system potential (a performance measure). The proposed techniques have been elaborated in order to obtain high simulation and forecast properties and applied to systems of nonstationary accelerator, gradient systems, and linear-quadratic stationary systems. We have also demonstrated a new concept of system potential specification in case of linearquadratic stationary systems. It is based on the principle of its basis decomposition as an element of energy space. All models and algorithms have been approbated using real statistical data for models of macroeconomic dynamics.


Keywords: Parametric identification, Performance measure, Non-stationary accelerator, Gradient system, Macroeconomic dynamics

1. Introduction. Dynamic systems design is one of the most urgent and difficult issues in different branches of science. Basically, it has three essential practical implications: simulation, forecast, and optimization. The first one is purely retrospective, while the others allow anticipating system dynamics. It is especially useful for analysts. This paper deals with techniques aimed at simulation and forecast. Thus, the quality of applied methodology will be considered in terms of simulation or forecast properties inherent in the resulted models.
The core issue of dynamic systems design is identification methodology and its practical efficiency. Basically, two types of concepts are common in the field of system identification [12]. The first one assumes full or partial specification of the relationship between systems inputs, states and outputs, while a number of unknown parameters should be estimated. Such an approach is called a grey box model. It uses parametric identification techniques (lest squares, general method of moments, maximum likelihood technique) $[5,1]$. The other concept is a black box model, which assumes no prior specification and uses both parametric and non-parametric identification techniques (transient response analysis, Fourier analysis EFTE) [4, 21]. System identification can also be carried out in either the time (applying Markov Parameters [8], subspace system identification method [13]) or frequency domain (using singular value decomposition of Hankel matrix of Markov parameters [8]). This paper pertains only the time-domain identification within grey box model framework.
The main problems which arise in this area are computational difficulties connected with numerical realization of optimization algorithms; bias, inefficiency or inconsistency of estimators caused by failure to meet all preconditions of classical estimation techniques; insufficient precision and inadequacy of models. That results in low simulation
and forecast properties. Furthermore, standard time-domain identification methods deal exclusively with parameters of a state equation (a law of dynamic system motion). That allows one to skip the accompanying problem of a system potential (or a performance measure) identification. In this paper performance measure is specified as a complimentary internal characteristics of a dynamic system interconnected with a law of its motion. A performance measure can not be considered as additional state variable to avoid correlation problem which may apparently appear in such models. Moreover, presence of a performance measure provides the model with an extra advantage of using it in optimization problems.

Thus, we intend to demonstrate a new effective computational approach to time-domain parametric identification of state-space dynamic systems comprising both a law of motion and a system potential. Actually, the paper is organized as follows. Section 2 presents the general problem statement of state-space dynamic systems identification. The proposed identification techniques are elaborated in Section 3. Section 4 demonstrates results of approbation, based on real statistical data for models of macroeconomic dynamics. Eventually, we make a conclusion in Section 5.
2. Problem Statement and Preliminaries. We would like to confine ourselves to systems described by ordinary differential equations (in state variable form). Thus, let $\{x(t)\}=\left\{x(t) \in E^{n} \mid t_{0} \leq t \leq t_{1}\right\}$ be a state trajectory of a dynamic system, a continuous vector-valued function. Its values are state column-vectors $x(t)=\left(x_{i}\right)^{i=1, n}$ from $n$-dimensional Euclidean space at time $t$. Similarly, let $\{u(t)\}=\left\{u(t) \in E^{r} \mid t_{0} \leq\right.$ $\left.t \leq t_{1}\right\}$ be an input trajectory, a piecewise continuous vector-valued function. Its values are input column-vectors $u(t)=\left(u_{j}\right)^{j=1, r}$ from $r$-dimensional Euclidean space at time $t$. Here and further on the simplifying assumption is made that the states are all available for measurement. Then the dynamic system may be described by $n$ first-order differential equations in the matrix form:

$$
\begin{equation*}
\dot{x}(t)=f(x(t), u(t), t), \tag{1}
\end{equation*}
$$

where $f(\cdots)$ is a vector of continuously differentiable functions. Its specification depends on physical interpretation of (1) [7].

Let $G$ be a system potential, an empirical continuously differentiable function of the form:

$$
\begin{equation*}
G=G(x(t), \dot{x}(t), u(t), t) . \tag{2}
\end{equation*}
$$

Let us divide the continuous time domain $\left[t_{0}, t_{1}\right]$ into $N$ discrete points of time. Supposing at every point of $t=0,1, \ldots, N-1$ there is a statistical information $x_{t}$ about state column-vector $x(t)$ and statistical information $G_{t}$ about system potential $G(t)$. If $f(\cdots)$ and $G$ are specified with some functional forms, then the task is to identify system (1), (2) so that the congruencies

$$
x_{t} \cong x(t) \text { and } G_{t} \cong G(t) \text { for } t=0,1, \ldots, N-1
$$

are performed with a certain precision.
3. Time-domain Parametric Identification. Generally speaking, time-domain parametric identification requires discretization of the examined system. Thus (1), (2) should be transformed from the differential to difference form. As a result, for every $t$ from $t=0,1, \ldots, N-1$

$$
\begin{array}{r}
G=G(x(t), \Delta x(t), u(t), t), \\
\Delta x(t)=f(x(t), u(t), t), \tag{4}
\end{array}
$$

where differences $x(t)$ may be either forward

$$
\begin{equation*}
\Delta x(t)=x(t+1)-x(t), t=0,1, \ldots, N-2 \tag{5}
\end{equation*}
$$

or backward [17]

$$
\begin{equation*}
\Delta x(t)=x(t)-x(t-1), t=1,2, \ldots, N-1 \tag{6}
\end{equation*}
$$

As investigation of forecast properties is essential in this paper, we will use backward differences (6) for further research.
The dynamic system presented by (3), (4) may be identified while using two different concepts. According to the first one (Section 3.1), a potential $G$ identification is primary. According to the second concept (Section 3.2), a law of motion $f(\cdots)$ is proposed to be identified firstly. Section 3.1 is devoted to non-stationary accelerators and gradient systems identification. Stationary linear-quadratic (LQ) systems are examined in Section 3.2.
3.1. Non-stationary accelerators and gradient systems. In this section we will derive effective techniques for two classes of time-variant dynamic systems identification. They are frequently used in macroeconomic modeling [18] and differential games theory [11]. Following the works [15, 20], let it be that the law of motion (1) and the potential (2) form a circuited (closed) system. Besides, we will specify $G$ only with state variables.

Then non-stationary accelerator is a system of the form:

$$
\begin{align*}
& G=G(x(t)), \\
& \dot{x}(t)=u(t) G . \tag{7}
\end{align*}
$$

Here the dimension $r$ of the input column-vector $u(t)$ equals to $n$. Model (7) belongs to zero-order differentiating elements [10].

We will use the following model for gradient systems:

$$
\begin{align*}
& G=G(x(t)) \\
& \dot{x}(t)=U(t)\left(\frac{\partial G}{\partial x}\right)^{\operatorname{tr}} \tag{8}
\end{align*}
$$

where $U$ is time-variant $n \times n$ diagonal matrix with the coordinates $\left\{u_{j}\right\}_{j=1, n}$ on its diagonal. From here on we assume that a derivative of a scalar over a column (row) vector is a row (column) vector [7].

Although both non-stationary accelerators and gradient systems are quite different physically, technically, either of them has the same identification technique. It should be noted that the problem of systems (7) and (8) identification naturally gets transformed into inverse dynamic problem. Given the measurements of state variables and identified potential, it deals with input identification [16].

Let vector $g=\left(g_{l}\right)^{l=1, n}$ state for the $n$-dimensional column-vector $(G, G, \ldots, G)^{\text {tr }}$ in case of non-stationary accelerator and for the $n$-dimensional column-vector $\left(\frac{\partial G}{\partial x_{1}}, \frac{\partial G}{\partial x_{2}}, \ldots, \frac{\partial G}{\partial x_{n}}\right)^{\operatorname{tr}}$ in terms of gradient system. Then the general dynamic process will take the form of

$$
\begin{equation*}
\dot{x}(t)=U(t) g \tag{9}
\end{equation*}
$$

At the beginning the potential $G$ should be specified. Firstly, it may be any smooth function to provide the existence of the first derivative. Secondly, it should be convenient for further research. Basically, log-linear form will be sufficient for the analysis of the first-order effects. The second order effects are usually examined by trans-log forms [6]. $G$ is specified here as a multiplicative Cobb-Douglas function [19], i.e.

$$
\begin{equation*}
G\left(x_{1}, x_{2}, \ldots, x_{n}\right)=a_{0} x_{1}^{a_{1}} x_{2}^{a_{2}} \cdots x_{n}^{a_{n}}, \tag{10}
\end{equation*}
$$

or in the log-linear form

$$
\begin{equation*}
\ln G\left(x_{1}, x_{2}, \ldots, x_{n}\right)=\ln a_{0}+a_{1} \ln x_{1}+a_{2} \ln x_{2}+\cdots+a_{n} \ln x_{n} . \tag{11}
\end{equation*}
$$

Parameters from (10) may be easily identified. Ordinary least squares (OLS) estimation for the linearized model (11) can be applied here.

It is also convenient to use polynomial specification of inputs $u_{j}(t)$ :

$$
\begin{equation*}
u_{j}(t)=b_{j 0}+b_{j 1} t+b_{j 2} t^{2}+\cdots b_{j k_{j}} t^{k_{j}}, j=1,2, \ldots, n \tag{12}
\end{equation*}
$$

Generally speaking, the order $k_{j}$ may be chosen upon different considerations: to ensure high simulation or forecast properties, to increase adequacy of the model, etc. Nevertheless, the substantial restriction to achieve all this goals is that the higher the order $k_{j}$ in (11) the lower the degree of freedom. The latter is undesirable for practical application. Further on this issue will be tackled with in details.

Given the identified system potential $G$ and specified input $u(t)$, the equation of motion (9) identification is proceeded to. Using backward difference scheme (4), (6), we arrive at the following difference equations for $t=1,2, \ldots, N-1$ and $i=1,2, \ldots, n$ :

$$
\begin{equation*}
x_{i}(t)=x_{i}(t-1)+\left.b_{i 0} \cdot g_{i}\right|_{t}+\left.b_{i 1} \cdot t \cdot g_{i}\right|_{t}+\cdots+\left.b_{i k_{i}} \cdot t^{k_{i}} \cdot g_{i}\right|_{t}+\varepsilon_{i}(t), \tag{13}
\end{equation*}
$$

where $\varepsilon_{k_{i}}(t)$ is a random disturbance, or if presented in the matrix form:

$$
\begin{gathered}
\left(\begin{array}{cccc}
\left.g_{i}\right|_{t=1} & \left.g_{i}\right|_{t=1} & \cdots & \left.g_{i}\right|_{t=1} \\
\left.g_{i}\right|_{t=2} & \left.2 g_{i}\right|_{t=2} & \cdots & \left.2^{k_{i}} g_{i}\right|_{t=2} \\
\vdots & \vdots & \ddots & \vdots \\
\left.g_{i}\right|_{t=N-1} & \left.(N-1) g_{i}\right|_{t=N-1} & \cdots & \left.(N-1)^{k_{i}} g_{i}\right|_{t=N-1}
\end{array}\right)\left(\begin{array}{c}
b_{i 0} \\
b_{i 1} \\
\vdots \\
b_{i k_{i}}
\end{array}\right)+ \\
\\
+\left(\begin{array}{c}
\varepsilon_{i}(1) \\
\varepsilon_{i}(2) \\
\vdots \\
\varepsilon_{i}(N-1)
\end{array}\right)=\left(\begin{array}{c}
x_{i}(1)-x_{i}(0) \\
x_{i}(2)-x_{i}(1) \\
\vdots \\
x_{i}(N-1)-x_{i}(N-2)
\end{array}\right)
\end{gathered}
$$

The other identification methodology is based on integral specification of the equation of motion (9):

$$
\begin{equation*}
x(t)=x\left(t_{0}\right)+\int_{t_{0}}^{t} U(t) g d t \tag{14}
\end{equation*}
$$

Substituting (10), (12) in (14) and using left-hand Riemann sums approximation for integrating, we arrive at a discrete counterpart of the model (14) for $t=0,1, \ldots, N-1$ and $i=1,2, \ldots, n$ :

$$
\begin{equation*}
x_{i}(t)=x_{i}^{*}\left(t_{0}\right)+\left.b_{i 0} \sum_{j=0}^{t-1} g_{i}\right|_{t=j}+\left.b_{i 1} \sum_{j=0}^{t-1} j g_{i}\right|_{t=j}+\cdots+\left.b_{i k_{i}} \sum_{j=0}^{t-1} j^{k_{i}} g_{i}\right|_{t=j}+\nu_{i}(t), \tag{15}
\end{equation*}
$$

where $\nu_{i}(t)$ is a random disturbance, or if presented in the matrix form:

$$
\left(\begin{array}{ccccc}
1 & 0 & 0 & \cdots & 0 \\
1 & \left.g_{i}\right|_{t=1} & \left.g_{i}\right|_{t=1} & \cdots & \left.g_{i}\right|_{t=1} \\
1 & \left.\sum_{t=0}^{1} g_{i}\right|_{t} & \left.\sum_{t=0} t g_{i}\right|_{t} & \cdots & \left.\sum_{t=0} t^{k_{i}} g_{i}\right|_{t} \\
\vdots & \vdots & \vdots & \ddots & \vdots \\
1 & \left.\sum_{t=0}^{N-2} g_{i}\right|_{t} & \left.\sum_{t=0}^{N-2} t g_{i}\right|_{t} & \cdots & \left.\sum_{t=0}^{N-2} t^{k_{i}} g_{i}\right|_{t}
\end{array}\right)\left(\begin{array}{c}
x_{i}^{*} \\
b_{i 0} \\
b_{i 1} \vdots \\
b_{i k_{i}}
\end{array}\right)+\left(\begin{array}{c}
\nu_{i}^{*}(0) \\
\nu_{i}(1) \\
\nu_{i}(2) \\
\vdots \\
\nu_{i}(N-1)
\end{array}\right)=\left(\begin{array}{c}
x_{i}(0) \\
x_{i}(1) \\
x_{i}(2) \\
\vdots \\
x_{i}(N-1)
\end{array}\right)
$$

Notice: $x\left(t_{0}\right)$ from (14) is not equal to $x^{*}\left(t_{0}\right)$ from (15), and $x^{*}\left(t_{0}\right)$ is assumed to be unknown. This allows obtaining regressions with constant terms represented by (15). Model (13) does not contain a constant term though. The latter requires some adjustments when applying classical techniques to parameters identification [22]. Further we will use both scheme (13) and (15) to compare them from the view-point of simulation and forecast properties.

But we still have to identify the order $k_{j}$ from (12). The order $k_{i}$ will be obtained under conditions of forecast confidence intervals minimization, while forecast properties of regression models are especially substantial for practical applications. The latter is defined as follows for the forecast value of state variable $x_{i}(N)$ :

$$
\begin{equation*}
\hat{x}_{i}(N)-\delta t_{\alpha}<x_{i}(N)<x_{i}(N)+\delta t_{\alpha}, \tag{16}
\end{equation*}
$$

where $\hat{x}_{i}(N)$ is a point forecast; $t_{\alpha}$ is a percentile of Student's distribution with $\alpha$, level of significance; $\delta$ is a standard error of forecasting [14].
3.2. Stationary LQ systems. Let us consider the general form of a linear stationary process:

$$
\begin{equation*}
\dot{x}(t)=A x(t)+B u(t), \tag{17}
\end{equation*}
$$

where $A$ and $B$ are fixed $n \times n$ and $n \times r$ matrices with unknown elements.
In this paper we suggest supposing $B=I$ and $A=\left\{a_{i j}\right\}_{i=1, n}^{j=1, n}$ as a symmetric matrix. This allows putting an interpretation on input $u(t)$ and considering $A$ as a matrix of some quadratic form [7]. In this case we may also elaborate an energy approach to identification of $A$.

Let system (17) take a form of

$$
\begin{equation*}
\dot{x}(t)=A x(t)+u(t) . \tag{18}
\end{equation*}
$$

Model (18) is an open lag model. It is overdetermined having $r$ degrees of freedom $(u(t) \neq 0)$. Similarly, the model

$$
\begin{equation*}
\dot{x}(t)=A x(t) \tag{19}
\end{equation*}
$$

is called a closed lag model [10]. Then vector-valued function $u(t)$ may be considered as a deviation of the open lag model (18) from the closed lag model (19). Obviously, $u(t)$ is characteristics of external impacts on the dynamic system. This allows treating input $u(t)$ as a control parameter.

Let us build up the following difference model using (6) to identify the symmetric matrix $A$ for the open lag model (18), :

$$
\begin{equation*}
\Delta x(t)=A x(t)+v(t), \tag{20}
\end{equation*}
$$

where $v(t)=\left(v_{i}(t)\right)^{i=1, n}$ is a discrete counterpart of $u(t)$. Then let (20) be transformed into regression model of the form

$$
\left\{\begin{align*}
\Delta x_{1} & =a_{10}+a_{11} x_{1}+a_{12} x_{2}+\cdots+a_{1 n} x_{n}+w_{1}  \tag{21}\\
\Delta x_{2} & =a_{20}+a_{12} x_{1}+a_{22} x_{2}+\cdots+a_{2 n} x_{n}+w_{2} \\
\quad & \quad \cdots \\
\Delta x_{n} & =a_{n 0}+a_{1 n} x_{1}+a_{2 n} x_{2}+\cdots+a_{n n} x_{n}+w_{n}
\end{align*}\right.
$$

where $a_{0}=\left(a_{i 0}\right)^{i=1, n}$ is a vector of fixed coefficients, required from econometric considerations; $w=\left(w_{i}\right)^{i=1, n}$ is a vector of random disturbances.

The following parametric iterative procedure for matrix $A$ identification has been established. Let us first estimate parameters of $n$ regressions from (21) separately. OLS as
an estimation method can be used. It is proposed that all non-diagonal elements of $A$ are averaged according to the formula

$$
a_{i j}^{\mathrm{av}}=\frac{a_{i j}^{*}+a_{j i}^{*}}{2}, i, j=1,2, \ldots, n
$$

where $\left\{a_{i j}^{*}\right\}_{j=1, n}^{i=1, n}$ are OLS-estimators of $\left\{a_{i j}\right\}_{j=1, n}^{i=1, n}$. Then diagonal elements of $A$ can be corrected. The models

$$
\left\{\begin{align*}
\Delta x_{1}-a_{12}^{\mathrm{av}} x_{2}-\cdots-a_{1 n}^{\mathrm{av}} x_{n} & =a_{10}^{*}+a_{11}^{*} x_{1}+\varepsilon_{1}  \tag{22}\\
\Delta x_{2}-a_{12}^{\mathrm{av}} x_{1}-\cdots-a_{2 n}^{\mathrm{av}} x_{n} & =a_{20}^{*}+a_{22}^{*} x_{1}+\varepsilon_{2} \\
& \cdots \\
\Delta x_{n}-a_{1 n}^{\mathrm{av}} x_{1}-\cdots-a_{n-1, n}^{\mathrm{av}} x_{n-1} & =a_{n 0}^{*}+a_{n n}^{*} x_{1}+\varepsilon_{n}
\end{align*}\right.
$$

may be OLS estimated. Here $\left\{\varepsilon_{i}\right\}_{i=1, n}$ are random disturbances.
Thus the estimated model corresponding to (4) takes the form of

$$
\left\{\begin{array}{l}
\Delta x_{1}=a_{11} x_{1}+a_{12}^{\mathrm{av}} x_{2}+\cdots+a_{1 n}^{\mathrm{av}} x_{n}+v_{1}  \tag{23}\\
\Delta x_{2}=a_{12}^{\mathrm{a}} x_{1}+a_{22} x_{2}+\cdots+a_{2 n}^{\mathrm{av}} x_{n}+v_{2} \\
\quad \quad \cdots \\
\Delta x_{n}=a_{1 n}^{\mathrm{av}} x_{1}+a_{2 n}^{\mathrm{av}} x_{2}+\cdots+a_{n n} x_{n}+v_{n}
\end{array}\right.
$$

where $\left\{v_{i} \mid v_{i}=\varepsilon_{i}+a_{i 0}\right\}$ are the coordinates of the vector-valued function $u(t) ;\left\{a_{i 0}\right\}_{i=1, n}$ are the OLS-estimators of $\left\{a_{i 0}^{*}\right\}$.

Coefficient of determination $R^{2}$ should be computed for every regression from (21) and (22), as well as significance of estimated coefficients should be verified. If regression analysis [22] indicates essential insignificancy of a state variable, then it may be excluded from the model with further recomputation.

In case when the proposed technique gives unsatisfied results concerning $R^{2}$, some adjustments are necessary. Particularly, (22) should be transformed into

$$
\left\{\begin{align*}
\Delta x_{1}-\lambda_{12} x_{2}-\cdots-\lambda_{1 n} x_{n} & =a_{10}^{*}+a_{11}^{*} x_{1}+\varepsilon_{1}  \tag{24}\\
\Delta x_{2}-\lambda_{21} x_{1}-\cdots-\lambda_{2 n} x_{n} & =a_{10}^{*}+a_{11}^{*} x_{1}+\varepsilon_{1} \\
& \cdots \\
\Delta x_{n}-\lambda_{n 1} x_{1}-\cdots-\lambda_{n, n-1} x_{n-1} & =a_{n 0}^{*}+a_{n n}^{*} x_{1}+\varepsilon_{n}
\end{align*}\right.
$$

where adjustment parameters $\left\{\lambda_{i j}\right\}_{j=1, n}^{i=1, n}$ are iteratively chosen from the interval $\left[a_{i j}, a_{j i}\right]$. This arbitrariness in (24) is aimed at improving the resulted models (for example, increasing of coefficient of determination $R^{2}$ ).

Once high values of coefficients of determination $R^{2}$ for regressions (23) have been obtained, one can start considering the equation of motion (18), where the elements of $A$ are the corresponding coefficients of model (23) and vector-valued function $v(t)$ states for input $u(t)$. As it is shown in [15, 20], such an approach provides sufficiently high simulation and forecast properties of dynamic models.

After matrix $A$ has been fully identified, it enables one to find an analytical solution of (18). This issue belongs to the class of direct dynamic problems, i.e. to solve matrix differential equation (18) given $A$ and inputs $u(t)$. Notice: $u(t)$ is defined at discrete points of time $t=1,2, \ldots, N-1$. In order to obtain continuous control, $u(t)$ may be presented as a kinked line.

When considering real dynamic systems, it appears that state variables are monotonous over time. For instance, the main macroeconomic indices are increasing for developed countries. (Then $\dot{x}(t)>0, t \in\left[t_{0}, t_{1}\right]$.) In such cases matrix $A$ may be either positive or
negative definite. The latter allows our solving system of ODE (18) analytically. Specifically, if $A$ is negative definite, then the left boundary condition $x\left(t_{0}\right)=x_{0}$ is given and the solution is found in direct time:

$$
\begin{equation*}
x(t)=e^{A t} x_{0}+\int_{0}^{t} e^{A(t-\tau)} u(\tau) d \tau \tag{25}
\end{equation*}
$$

if $A$ is positive definite, then the right boundary condition $x\left(t_{1}\right)=x_{1}$ is given and the solution is found in inverse time:

$$
\begin{equation*}
x(t)=e^{A\left(t-t_{1}\right)} x_{1}+\int_{t_{1}}^{t} e^{A(t-\tau)} u(\tau) d \tau \tag{26}
\end{equation*}
$$

Solutions (25) and (26) are stable [3] and they may be employed for simulation and forecast.

Now we proceed to system potential (3) specification and identification. For this we elaborate and employ an energy approach. It is based on potential $G$ basis decomposition as an element of energy space. The following assumption and theorem are essential for further discussion.

Let model (2) be an autonomous system and the impact of all components on $G(\cdots)$ be characterized with the corresponding quadratic forms:

$$
Q_{P_{1}}(x)=\frac{1}{2} x^{\operatorname{tr}} P_{1} x, Q_{P_{2}}(\dot{x})=\frac{1}{2} \dot{x}^{\operatorname{tr}} P_{2} \dot{x}, Q_{P_{3}}(u)=\frac{1}{2} u^{\operatorname{tr}} P_{3} u
$$

The quadratic form $Q_{P_{1}}(x)$ states for potential energy of the open lag model (18) at time $t, Q_{P_{2}}(\dot{x})$ for kinetic energy of the open lag model (18) at time $t$, and, at last, $Q_{P_{3}}(u)$ for energy of external inputs.

Then considering $G$ as an element of energy space and energies $Q_{P_{1}}(x), Q_{P_{2}}(\dot{x}), Q_{P_{3}}(u)$ as components of the basis of this space, one can arrive at the following basis decomposition of potential $G$ :

$$
\begin{equation*}
G=C_{0}+C_{1} \frac{1}{2} Q_{P_{1}}(x)+C_{2} \frac{1}{2} Q_{P_{2}}(\dot{x})+C_{3} \frac{1}{2} Q_{P_{3}}(u), \tag{27}
\end{equation*}
$$

where $C_{0}, C_{1}, C_{2}$ and $C_{3}$ are unknown constants.
Generally, the matrices $P_{1}, P_{2}$ and $P_{3}$ are hard to identify. Nevertheless, further considerations will substantially simplify this issue.
Theorem 3.1. The identity

$$
\begin{equation*}
\frac{1}{2} x^{t r} A x+\frac{1}{2} \dot{x}^{t r} A^{-1} \dot{x}-\frac{1}{2} u^{t r} A^{-1} u \equiv \frac{1}{2} \frac{d}{d t}\left(x^{t r} x\right) \tag{28}
\end{equation*}
$$

holds true, if and only if the law of motion is described by (18).
Proof: It follows from (18) that

$$
\frac{1}{2} u^{\operatorname{tr}} A^{-1} u \equiv \frac{1}{2}\left(\dot{x}^{\operatorname{tr}}-x^{\operatorname{tr}} A\right) A^{-1}(\dot{x}-A x) \equiv \frac{1}{2} x^{\operatorname{tr}} A x+\frac{1}{2} \dot{x}^{\operatorname{tr}} A^{-1} \dot{x}-\frac{1}{2} \dot{x}^{\operatorname{tr}} x-\frac{1}{2} x^{\operatorname{tr}} \dot{x}
$$

Notice: $\dot{x}^{\operatorname{tr}} x \equiv x^{\operatorname{tr}} \dot{x} \equiv \frac{d}{d t}\left(x^{\operatorname{tr}} x\right)$. Then formula (28) obviously arises. Apparently, the equation of motion (18) is also derived from (28).

Therefore, it follows from the energy equality that $P_{1}, P_{2}$ and $P_{3}$ may be reasonably identified as $A, A^{-1}$ and $A^{-1}$ respectively. Finally, the unknown coefficients from (27) are to be estimated. Employing OLS, we arrive at the identified model of system potential:

$$
\begin{equation*}
\hat{G}=\hat{C}_{0}+\hat{C}_{1} \frac{1}{2}\left(x^{\mathrm{tr}} A x\right)+\hat{C}_{2} \frac{1}{2}\left(\Delta x^{\operatorname{tr}} A^{-1} \Delta x\right)+\hat{C}_{3} \frac{1}{2}\left(u^{\mathrm{tr}} A^{-1} u\right) \tag{29}
\end{equation*}
$$

where $\hat{C}_{0}, \hat{C}_{1}, \hat{C}_{2}$ and $\hat{C}_{3}$ are estimated coefficients. They may be obtained by solving the following matrix equation:

$$
\left(\begin{array}{cccc}
1 & \left.\frac{1}{2} Q_{A}(x)\right|_{t=1} & \left.\frac{1}{2} Q_{A^{-1}}(\Delta x)\right|_{t=1} & \left.\frac{1}{2} Q_{A^{-1}}(u)\right|_{t=1} \\
1 & \left.\frac{1}{2} Q_{A}(x)\right|_{t=2} & \left.\frac{1}{2} Q_{A^{-1}}(\Delta x)\right|_{t=2} & \left.\frac{1}{2} Q_{A^{-1}}(u)\right|_{t=2} \\
\vdots & \vdots & \vdots & \vdots \\
1 & \left.\frac{1}{2} Q_{A}(x)\right|_{t=N-1} & \left.\frac{1}{2} Q_{A^{-1}}(\Delta x)\right|_{t=N-1} & \left.\frac{1}{2} Q_{A^{-1}}(u)\right|_{t=N-1}
\end{array}\right)\left(\begin{array}{c}
\hat{C}_{0} \\
\hat{C}_{1} \\
\hat{C}_{2} \\
\hat{C}_{3}
\end{array}\right)=\left(\begin{array}{c}
G_{1} \\
G_{2} \\
\vdots \\
G_{N-1}
\end{array}\right)
$$

where $u_{t}=x_{t}-x_{t-1}-A x_{t}$ for $t=1,2, \ldots, N-1$.
4. Numerical Experiment. This section deals with numerical application of proposed identification techniques to macroeconomic modeling for USA, Japan and a number of EU countries.

In all models we use GDP as potential of macroeconomic system. This specification follows from the fact that GDP accumulates all information of the system's production capacity and is considered as one of the most essential macroeconomic indicator. As to state variables specification, it generally depends on the results of correlation and regression analysis. We will use fixed capital stock, property costs, labor force and final consumption expenditures as initial candidates for this role though. All statistical information is available on Internet (see, for example, http://epp.eurostat.ec.europa.eu).
4.1. A model of investment development. Let us consider an open macroeconomic system that consists of $n-1$ industries. As experiment reveals [15] state variables of such a system are fixed capital stocks $\left\{x_{i}\right\}_{i=1, n-1}$ per industry and foreign debt $x_{n}$.

The rate of change $\dot{x}_{i}(t)$ of fixed capital is the value of net investments into the $i$ th industry for $i=1,2, \ldots, n-1$. The rate of change $\dot{x}_{n}$ of foreign debt is the value of its gross accumulation. These values are used for the analysis of investment activity in economic theory [10], but mainly as ratios to gross output (e.g. GDP) rather than in absolute numbers. Therefore, specifying vector-valued function $f(\cdots)$ as a product of dimensionless value $u(t)$ and GDP $G$, we obtain a model of non-stationary accelerator (7). In this case economic interpretation of control variables is also obvious. Functions $\left\{u_{j}(t)\right\}_{j=1, n-1}$ may be considered as indices of investment activity per industry and $u_{n}(t)$ as an index of export-import non balance.

Approbation of (7) is based on two-industrial Danish economy in 1966-1997. Let the first industry consist of manufacturing and agricultural branches and the second one of services.

Using (10) Denmark's GDP is specified with two regressors: fixed capital stock $x_{1}+x_{2}$ and foreign debt $x_{3}$. The obtained OLS estimations are:

$$
\begin{equation*}
\underset{(\text { s.e. })}{\ln G}=\underset{(0.3195)}{0.3859}+\underset{(0.0428)}{0.8145} \ln \left(x_{1}+x_{2}\right)+\underset{(0.0252)}{0.1120} \ln x_{3}, R^{2}=0.9968, \tag{30}
\end{equation*}
$$

where numbers in brackets are standard errors of regression coefficients. All coefficients, except the first one, appear to be statistically significant (we use Student's test for the number of freedom $l=29$ and level of significance $\alpha=5 \%$ ). Such results are quite natural in econometric literature [6] when using log-linear functional forms for the regression analysis and with high coefficient of determination $R^{2}$ imply that the model is quite fulfilled with specified factors. It is also evident that the impact of foreign debt on Denmark's GDP was not essential in the analyzed period of time (although it appeared to be statistically significant). Another conclusion is that the diminishing return to scale was present $\left(a_{1}+a_{2}=0.9265<1\right)$ at that time.

The exponential form of the estimated log-linear production function is

$$
G\left(x_{1}+x_{2}, x_{3}\right)=1.4710\left(x_{1}+x_{2}\right)^{0.8145} x_{3}{ }^{0.1120}
$$

Let us now turn to the polynomial (12) identification. The results of OLS estimation are given below:

$$
\underset{(\text { s.e.) })}{u_{1}(t)}=\underset{(0.0118)}{0.1469}-\underset{(0.0005)}{0.0037} t, \underset{(\text { s.e. })}{u_{2}(t)}=\underset{(0.0448)}{0.2219}-\underset{(0.0043)}{0.0066} t+\underset{(0.0001)}{0.000042} t^{2}
$$

for difference identification scheme (13) and

$$
\underset{(\text { s.e. })}{u_{1}(t)}=\underset{(0.0033)}{0.1543}-\underset{(0.0002)}{0.0041} t, \underset{(\text { s.e. })}{u_{2}(t)}=\underset{(0.0167)}{0.2278}-\underset{(0.0019)}{0.0062} t
$$

for integral identification scheme (15). All coefficients in estimated polynomials appear to be statistically significant and provide the high level of approximation for state variables $x_{1}$ and $x_{2}$. Particularly coefficients of determination $R^{2}$ are 0.9979 and 0.9966 for scheme (13) and 0.9986 and 0.9982 for scheme (15).

As experiments reveal, optimal order $k_{j}$ for each of polynomials (12) appears to be equal to 1 or 2. Interval forecasts (16) for state variables $x_{1}$ and $x_{2}$ are $1197674.539 \pm$ 20073.928 and $1800344.834 \pm 36369.454$ for scheme (13) versus $11883527.018 \pm 34599.568$ and $1764626.422 \pm 56578.455$ for scheme (15). Obviously confidence intervals for the forecasts based on difference identification algorithm are narrower: $1.68 \%$ and $2.02 \%$ versus $2.92 \%$ and $3.21 \%$ (as percentages of point forecasts). The results obtained in this subsection confirm that difference identification scheme (13) is more preferable for forecast purposes. Integral identification scheme (15) ensures better simulation properties.
4.2. A Gradient macroeconomic model. Let macroeconomic system be an aggregate market of $n$ participants (players), which impact a system potential $G$ acting either within competitive or cooperative framework. Let $x(t)$ be a vector of $n$ strategies (factors) available for each of $n$ players. All strategies can be divided into those that impact $G$ positively and those that impact it negatively. Thus a potential $G$ accumulates aggregate results of players' interaction.
Following [20], the aforementioned situation may be described using gradient model (8). Therefore coordinates of $u(t)$ may be considered as rates of each players investing into development of the corresponding factor. Then $\frac{\partial G}{\partial x}$ is a vector of investment efficiencies.
In this sub-section we will also employ Danish economy as an example. As experiment reveals [20], it is sufficient to consider three aggregate players: producers (firms) operated with fixed capital consumption $\left(x_{1}\right)$, households operated with labor force ( $x_{2}$ ) and government operated with final consumption expenditures $\left(x_{3}\right)$.

OLS-estimation of unknown parameters in linearized production function (11) gives:

$$
\begin{equation*}
\ln _{(\mathrm{s} . \mathrm{e} .)} G=\underset{(1.538)}{-1.714}+\underset{(0.162)}{0.450} \ln x_{1}+\underset{(0.183)}{0.434} \ln x_{2}+\underset{(0.177)}{0.487} \ln x_{3}, R^{2}=0.9986, \tag{31}
\end{equation*}
$$

or in the form of (10):

$$
G\left(x_{1}, x_{2}, x_{3}\right)=0.180 x_{1}{ }^{0.450} x_{2}{ }^{0.434} x_{3}{ }^{0.487} .
$$

Notice that all coefficients in (31) are statistically significant with $\alpha=0.05$, the level of significance. Therefore it can be employed for modeling of Danish macroeconomic system. It is also derived from (31) that all factors impact GDP positively and more or less evenly provide increasing return to the scale $\left(a_{1}+a_{2}+a_{3}=1.371>1\right)$.

Using difference scheme (13) for polynomials (12) identification, one can obtain:

$$
\begin{array}{r}
u_{1}(t)=\underset{(25.735)}{389.526}, \underset{(\text { s.e. })}{u_{2}(t)}=\underset{(0.201)}{0.466}-\underset{(0.0082)}{0.018} t \\
(\text { s.e. }) \\
u_{3}(t)=\underset{(1471.323)}{4937.955}+\underset{(1107.170)}{2504.817} t-\underset{(218.916)}{503.835} t^{2}+\underset{(11.831)}{26.465} t^{3} .
\end{array}
$$

Table 1. The results of modeling for USA, Japan, France, Italy, Netherlands, and Finland


Interval forecasts (16) of state variables are $31757.1 \pm 594.2$ for fixed capital consumption, $2749 \pm 53$ for labor force and $147740.5 \pm 3670.1$ for final consumption expenditures.

In its turn, integral identification scheme (15) results in:

$$
\begin{array}{r}
\underset{(\text { s.e. })}{u_{1}(t)}=\underset{(17.379)}{355.147}+\underset{(2.796)}{8.169} t \underset{(\text { s.e. })}{u_{2}(t)}=\underset{(0.288)}{-2.069}+\underset{(0.107)}{1.104} t, \underset{(0.0083)}{-0.0859} t^{2} \\
u_{3}(t)=\underset{(146.178)}{7266.0158 .}
\end{array}
$$

In this case interval forecasts (16) of state variables are $32327.9 \pm 934.7$ for fixed capital consumption, $2675 \pm 99$ for labor force and $146293.8 \pm 4581.1$ for final consumption expenditures.

As experiments reveal, optimal order $k_{j}$ for each of polynomials (12) appears to range from 0 to 3 . Obviously, confidence intervals for the forecasts based on difference identification algorithm are narrower: $1.87 \%, 1.93 \%$ and $2.48 \%$ versus $4.52 \%, 5.89 \%$ and $6.31 \%$ (as percentages of point forecasts). All coefficients in estimated polynomials appear to be statistically significant, providing the high level of approximation for state variables $x_{1}$, $x_{2}$ and $x_{3}$. Computations of coefficients of determination $R^{2}$ give values $0.99414,0.91348$, 0.99841 for scheme (13) and $0.99832,0.98068,0.99517$ for scheme (15).

The results of similar modeling for USA, Japan, France, Italy, Netherlands and Finland are presented in Table 1. It consists of forecast values (numbers from above) and confidence intervals as percentages of point forecasts (numbers from below). System potential $G$ was initially specified with three factors. If it appeared that $G$ contained insignificant coefficients, then the recomputations would be held with corresponding two factors. Table 1 as well as aforementioned results coincided with those from Subsection 4.1: difference identification scheme (13) is more preferable for prospective analysis. Integral identification scheme (15) has better descriptive properties though.

The adequacy of the results from Subsections 4.1 and 4.2 was verified using GaussMarkov assumptions for classical regressions [22]. As computations [15, 20] reveal, only the assumption of non-autocorrelation failed. The presence of autocorrelation in models (13) and (15) is considered admissible, while input $u(t)$ is approximated by polynomials depended on time [6, 14].
4.3. An energy macroeconomic model. This subsection deals with approbation of stationary LQ model (18), (27) identification techniques. As before, it is proposed to investigate Denmark's economy. Let the dynamic system (18), (27) be specified with fixed capital stock $\left(x_{1}\right)$, property costs $\left(x_{2}\right)$, and final consumption expenditures $\left(x_{3}\right)$. All data are considered as ratios to initial time.

We obtain acceptable results right after the first iteration, when applying the proposed iterative procedure of matrix $A$ identification (see Subsection 3.2):

$$
A=\left(\begin{array}{ccc}
0.2042 & -0.2974 & -0.0421 \\
-0.2974 & 0.7154 & -0.1601 \\
-0.0421 & -0.1601 & 0.2527
\end{array}\right)
$$

All regressions from (23) appear to have high coefficients of determination $R^{2}$ (97$98 \%$ ). This argues in favor of models correspondence with real trends in macroeconomic dynamics.

Then, we obtain identified system potential (29), using OLS estimation:

$$
\begin{equation*}
\underset{(\mathrm{s} . \mathrm{e} .)}{\hat{G}}=\underset{(0.11)}{0.74}+\underset{(0.14)}{0.53} Q_{A}(x)+\underset{(0.14)}{0.39} Q_{A^{-1}}(\dot{x})-\underset{(0.22)}{0.64} Q_{A^{-1}}(u) \tag{32}
\end{equation*}
$$

According to Student's criterion with the level of significance $\alpha=0.05$, all components from specified energy space impact system potential (GDP) significantly. Coefficient of determination $R^{2}$ equal to $98 \%$ also argues in favor of proposed model. Figure 1(a) demonstrates real (designated with dots) and simulated (curves) values of potential $G$ dynamics. We may also make a forecast. We obtain a point forecast 3.02 and its standard error $\delta=0.05$, while extrapolating (32) for one period ahead.

Our investigation concerns dynamics of stable development. All eigenvalues of $A$ are positive: $\lambda_{1}=0.8787, \lambda_{2}=0.0166, \lambda_{3}=0.2770$. Therefore $A$ is defined positive. Therefore, in order to verify simulation and forecast properties, the formula (26) can be employed. Here are confidence intervals of forecasts for $x_{1}, x_{2}$ and $x_{3}: 4.95 \pm 0.20,3.32 \pm$ 0.09 and $3.04 \pm 0.08$.

Figure 1 demonstrates real and simulated vales of fixed capital stock (Figure 1(b)), property costs (Figure 1(c)), and final consumption expenditures (Figure 1(d)). It reveals high simulation quality of the models and thus efficiency of the proposed identification techniques. The sensitivity of output to the inputs has been tested using the analytical solution (26) and econometric model (20). Both lead to the same results and indicates the stability of obtained models.


Figure 1. Real and simulated values of (a) system potential $G$, (b) fixed capital consumption $x_{1}$, (c) property costs $x_{2}$ and (d) consumption expenditures $x_{3}$
5. Conclusions. A series of time-domain parametric identification techniques for a number of dynamic systems have been presented. On the one hand, two identification algorithms for models of non-stationary accelerator and gradient systems have been elaborated. Difference identification scheme has appeared to be more applicable for forecast analysis, while integral one has demonstrated better simulation properties. System potential for these models should be identified primarily. On the other hand, a new approach to specification and identification of stationary LQ systems has been established. In this case the equation of motion should be identified primarily. A parametric iterative procedure has been proposed for this purpose. It is based on analogues between overdetermined equations and regressions. System potential has been specified as an element of energy space with the components: potential, kinetic, and energy of inputs. The energy equality
has also been established. It has considerably simplified identification procedures. All identification techniques has been approbated using real statistical data for models of macroeconomic dynamics. As experiment reveals, all models have demonstrated good simulation, forecast properties and good correspondence with reality. Concerning identification time-span as a learning domain, further investigations may be connected with optimization problems in future periods.

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