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# Asymptotic tracking with prescribed transient behaviour for linear systems\*

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## Abstract

The problem of asymptotic tracking of reference signals is considered in the context of  $m$ -input,  $m$ -output linear systems  $(A, B, C)$  with the following structural properties: (i)  $CB$  is sign definite (but possibly of unknown sign), (ii) the zero dynamics are exponentially stable. The class  $\mathcal{Y}_{\text{ref}}(\alpha)$  of reference signals is the set of all possible solutions of a fixed, stable, linear, homogeneous differential equation (with associated characteristic polynomial  $\alpha$ ). The first control objective is asymptotic tracking, by the system output  $y = Cx$ , of any reference signal  $r \in \mathcal{Y}_{\text{ref}}(\alpha)$ . The second objective is guaranteed error  $e = y - r$  transient performance:  $e$  should evolve within a prescribed performance funnel  $\mathcal{F}_\varphi$  (determined by a function  $\varphi$ ). Both objectives are achieved simultaneously by an internal model in series with a proportional time-varying error feedback  $t \mapsto u(t) = -k(t)e(t)$ . The time-varying proportional factor  $k(t)$  is generated via a nonlinear function of the product  $\|e(t)\|\varphi(t)$ . The feedback structure essentially exploits an intrinsic high-gain property of the system by ensuring that, if  $(t, e(t))$  approaches the funnel boundary, then the gain attains values sufficiently large to preclude boundary contact.

**Keywords:** Tracking, output feedback, transient behaviour, internal model, minimum phase

## 1 Introduction

In the precursor [4] to the present paper, the concept of a performance funnel was introduced in a context of tracking control for nonlinear systems. The basic problem addressed therein was that of approximate tracking (with prescribed transient behaviour), by the system output  $y$ , of any absolutely continuous and bounded function  $r$  with essentially bounded derivative: the terminology “approximate tracking” means that, for any prescribed  $\lambda > 0$ , a control structure can be determined which ensures that the tracking error  $e = y - r$  is ultimately bounded by  $\lambda$  (that is,  $\|e(t)\| \leq \lambda$  for all  $t$  sufficiently large); the terminology “with prescribed transient behaviour” means that,

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for some suitable prescribed function  $\varphi$ , the error function is required to satisfy  $\|e(t)\| \leq 1/\varphi(t)$  for all  $t > 0$ . The choice of  $\varphi$  determines the transient behaviour; moreover, by imposing the property  $\liminf_{t \rightarrow \infty} \varphi(t) \geq \lambda > 0$ , the approximate tracking objective is assured. For example, with  $\varphi : t \mapsto \min\{t/T, 1\}/\lambda$ , the approximate tracking objective is achieved in prescribed time  $T > 0$ . Figure 1 encapsulates the approach: the function  $\varphi$  determines the performance funnel  $\mathcal{F}_\varphi$ , which may be identified with the graph of the set-valued map  $t \mapsto \{v \mid \varphi(t)\|v\| < 1\}$ . Simply stated, the control objective is to maintain the evolution of the tracking error in the funnel  $\mathcal{F}_\varphi$ . For

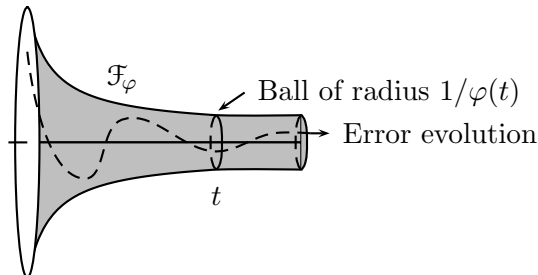


Figure 1: Performance funnel  $\mathcal{F}_\varphi$

the general class of reference signals considered in [4], the condition  $\liminf_{t \rightarrow \infty} \varphi(t) > 0$  cannot be relaxed, and so *exact* asymptotic tracking cannot be achieved. The purpose of the present note is to demonstrate that the condition may be relaxed if one confines attention to minimum-phase linear systems with sign-definite high frequency gain and restricts the class of reference signals to coincide with the set of solutions of a fixed, stable, linear, homogeneous differential equation. Under these restrictions, exact asymptotic tracking is achieved by adopting an internal model (capable of replicating the reference signals) in conjunction with a performance funnel with radius asymptotic to zero and an output feedback structure akin to that in [4, Section 6.3]. In an adaptive control context, the use of internal models in problems of asymptotic tracking for linear systems is well established (see, for example, [6, 7, 2, 3]). We emphasize that the approach adopted in the present paper is non-adaptive.

## 2 Class of systems

We consider the class of  $m$ -input ( $u(t) \in \mathbb{R}^m$ ),  $m$ -output ( $y(t) \in \mathbb{R}^m$ ) linear systems of the form

$$\left. \begin{aligned} \dot{x}(t) &= Ax(t) + Bu(t), & x(0) &= x^0 \in \mathbb{R}^n \\ y(t) &= Cx(t), \end{aligned} \right\} \quad (2.1)$$

where the triple  $(A, B, C) \in \mathbb{R}^{n \times n} \times \mathbb{R}^{n \times m} \times \mathbb{R}^{m \times n}$  has the following properties:

P1: *strict relative degree one with sign-definite high-frequency gain*, that is,

$$\langle x, CBx \rangle = 0 \iff x = 0,$$

P2: *minimum-phase*, that is,

$$\det \begin{bmatrix} sI - A & B \\ C & 0 \end{bmatrix} \neq 0 \quad \text{for all } s \in \mathbb{C} \text{ with } \operatorname{Re} s \geq 0.$$

Under assumption P1, the minimum-phase property P2 is equivalent to the assumption that the system (2.1) has exponentially stable zero dynamics (this equivalence can also be deduced from Lemma 3.3 below).

## 2.1 Control objectives, class of reference signals, and performance funnel

Let  $\mathcal{M}$  denote the set of square real matrices having no eigenvalue with positive real part and such that every eigenvalue on the imaginary axis is semi-simple. The *reference signals* to be tracked are all functions  $r : \mathbb{R}_+ \rightarrow \mathbb{R}^m$  the components  $r_i$  of which are solutions of the scalar differential equation  $\alpha(\frac{d}{dt})r_i(\cdot) = 0$ , where  $\alpha \in \mathbb{R}[s]$  is the characteristic polynomial of some  $M \in \mathcal{M}$  (and so every such function  $r$  is bounded). We denote this reference signal class by

$$\mathcal{Y}_{\text{ref}}(\alpha) := \left\{ r \in C^\infty(\mathbb{R}_+, \mathbb{R}^m) \mid \alpha\left(\frac{d}{dt}\right)r(\cdot) = 0, \alpha(s) = \det[sI - M], M \in \mathcal{M} \right\}.$$

The first control objective is asymptotic (output) tracking of any reference signal  $r \in \mathcal{Y}_{\text{ref}}(\alpha)$ . By this we mean a (dynamic) output feedback strategy which incorporates an internal model (capable of replicating the reference signal) and which ensures that  $\lim_{t \rightarrow \infty} (y(t) - r(t)) = 0$  whilst maintaining boundedness all the other signals. The second control objective is prescribed transient behaviour of the error signal  $e = y - r$ . We capture both objectives in the concept of a performance funnel

$$\mathcal{F}_\varphi := \left\{ (t, e) \in \mathbb{R}_+ \times \mathbb{R}^m \mid \varphi(t) \|e\| < 1 \right\} \quad (2.2)$$

associated with a function  $\varphi$  (the reciprocal of which determines the funnel boundary) with the following properties

$$\left. \begin{array}{l} \text{(a) } \varphi : \mathbb{R}_+ \rightarrow \mathbb{R}_+ \text{ is absolutely continuous and non-decreasing;} \\ \text{(b) } \varphi(t) = 0 \iff t = 0; \\ \text{there exists } c > 1 \text{ such that:} \\ \text{(c) } \varphi(t) \leq c\varphi(t/2) \quad \text{for all } t \in \mathbb{R}_+; \\ \text{(d) } \dot{\varphi}(t) \leq c[1 + \varphi(t)] \quad \text{for almost all } t \in \mathbb{R}_+. \end{array} \right\} \quad (2.3)$$

For example,  $t \mapsto \varphi(t) = t^a$ ,  $a \geq 1$ , satisfies (2.3) with  $c = 2^a$ . We record the following observation for later use.

**Proposition 2.1** *Let  $\varphi$  be such that (2.3) holds. For every  $p \geq \ln c / \ln 2$ ,*

$$0 < \varphi(t) \leq \varphi(1) [1 + ct^p] \quad \text{for all } t > 0. \quad (2.4)$$

**Proof:** Since  $\varphi$  is non-decreasing with property (b), we have  $0 < \varphi(t) \leq \varphi(1)$  for all  $t \in (0, 1]$ . Now, let  $t \in (1, \infty)$  be arbitrary and choose  $n \in \mathbb{N}$  such that  $2^{n-1} \leq t \leq 2^n$  or, equivalently,  $1/2 \leq t/2^n \leq 1$ . Then, by (b), (c) and the non-decreasing property,

$$0 < \varphi(t) \leq c\varphi(t/2) \leq \dots \leq c^n \varphi(t/2^n) \leq c^n \varphi(1) = c\varphi(1) 2^{(n-1)\ln c/\ln 2} \leq c\varphi(1) t^p.$$

The claim (2.4) follows.  $\square$

Proposition 2.1 implies, in particular, that exponentially contracting funnels are excluded.

### 3 The control

Let  $(A, B, C) \in \mathbb{R}^{n \times n} \times \mathbb{R}^{n \times m} \times \mathbb{R}^{m \times n}$  be such that P1 and P2 hold, and define

$$s(CB) := \begin{cases} +1, & \text{if } CB \text{ is positive definite, i.e. } \langle x, CBx \rangle > 0 \quad \forall x \neq 0 \\ -1, & \text{if } CB \text{ is negative definite, i.e. } \langle x, CBx \rangle < 0 \quad \forall x \neq 0. \end{cases} \quad (3.5)$$

We will have occasion to consider the two possible cases:  $s(CB)$  known or unknown *a priori* (the latter case is largely of academic interest).

#### 3.1 Internal model

A body of work by Francis and Wonham in the 1970s (see, for example, [1, 10]) led to the so-called *Internal Model Principle*, succinctly summarized in the context of linear systems in [11, p. 210] as “every good regulator must incorporate a model of the outside world”. Recent extensions of this “principle” to a nonlinear setting are contained in [9].

Let  $\alpha \in \mathbb{R}[s]$  be the characteristic polynomial of some  $M \in \mathcal{M}$  (and so every  $r \in \mathcal{Y}_{\text{ref}}(\alpha)$  is bounded). Let  $\beta \in \mathbb{R}[s]$  be a monic Hurwitz polynomial (i.e. all zeros of  $\beta$  lie in the open left half complex plane) and such that  $\alpha$  and  $\beta$  are coprime of degree  $p := \deg \beta = \deg \alpha$ . Then

$$\lim_{s \rightarrow \infty} \beta(s)/\alpha(s) = 1. \quad (3.6)$$

The *internal model* is now defined to be the  $m$ -input,  $m$ -output linear system with transfer function

$$G_m(s) := \frac{\beta(s)}{\alpha(s)} I_m. \quad (3.7)$$

Let  $(\hat{A}, \hat{b}, \hat{c}, 1) \in \mathbb{R}^{p \times p} \times \mathbb{R}^{p \times 1} \times \mathbb{R}^{1 \times p} \times \mathbb{R}$  be a minimal state space realization of  $\beta(s)/\alpha(s)$ . Then a minimal state space realization of the internal model is given by

$$\left. \begin{aligned} \dot{\xi}(t) &= \tilde{A} \xi(t) + \tilde{B} w(t), & \xi(0) &= \xi^0 \\ u(t) &= \tilde{C} \xi(t) + I_m w(t) \end{aligned} \right\} \quad (3.8)$$

with

$$\begin{aligned} \tilde{A} &= \text{diag}\{\hat{A}, \dots, \hat{A}\} \in \mathbb{R}^{mp \times mp}, & \tilde{B} &= \text{diag}\{\hat{b}, \dots, \hat{b}\} \in \mathbb{R}^{mp \times m}, \\ \tilde{C} &= \text{diag}\{\hat{c}, \dots, \hat{c}\} \in \mathbb{R}^{m \times mp}, & \xi^0 &\in \mathbb{R}^{mp}. \end{aligned}$$

We will refer to  $(\tilde{A}, \tilde{B}, \tilde{C}, I_m)$  as the internal model (although, strictly speaking, the use of “the” here is incorrect as any quadruple in the similarity orbit of  $(\tilde{A}, \tilde{B}, \tilde{C}, I_m)$  also qualifies for the title “internal model”).

### 3.2 Feedback

Let  $\varphi$  be such that (2.3) holds and let  $\mathcal{F}_\varphi$  be the associated performance funnel given by (2.2). Let  $\nu: \mathbb{R} \rightarrow \mathbb{R}$  be any  $C^\infty$  function such that, for some strictly-increasing, unbounded sequence  $(k_j)$  in  $(1, \infty)$ ,

$$\nu(k_j) s(CB) \rightarrow \infty \quad \text{as } j \rightarrow \infty. \quad (3.9)$$

If  $s(CB)$  is known *a priori*, then  $\nu: k \mapsto k s(CB)$  suffices. If  $s(CB)$  is unknown *a priori*, then any  $C^\infty$  function  $\nu$  with the following properties suffices:

$$\limsup_{k \rightarrow \infty} \nu(k) = +\infty \quad \text{and} \quad \liminf_{k \rightarrow \infty} \nu(k) = -\infty. \quad (3.10)$$

A simple example of a function satisfying (3.10) is  $\nu: k \mapsto k \cos k$ . In the latter case of unknown  $s(CB)$ , the rôle of the function  $\nu$  is similar to the concept of a “Nussbaum” function in adaptive control. Note, however, that the requisite properties (3.10) are less restrictive than (a) the “Nussbaum property”

$$\limsup_{k \rightarrow \infty} \frac{1}{k} \int_0^k \nu(\kappa) d\kappa = \infty, \quad \liminf_{k \rightarrow \infty} \frac{1}{k} \int_0^k \nu(\kappa) d\kappa = -\infty,$$

as required in [12], for example, or (b) the stronger “scaling invariant” Nussbaum property, as required in [5], for example.

The control strategy is given by

$$w(t) = -\nu(k(t)) [y(t) - r(t)], \quad k(t) = \frac{1}{1 - (\varphi(t) \|y(t) - r(t)\|)^2}, \quad (3.11)$$

in series with the internal model (3.8).

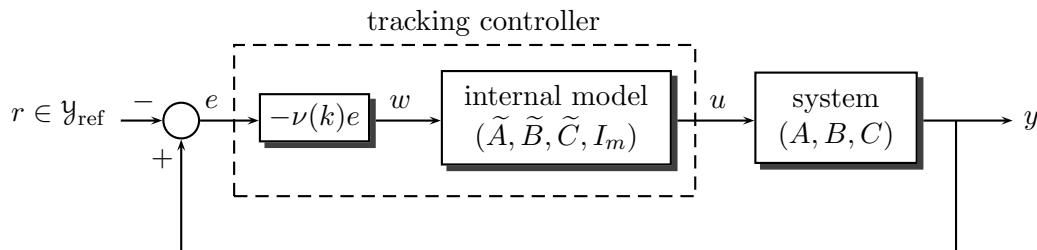


Figure 2: Asymptotic tracking controller with internal model

### 3.3 Closed-loop system

The conjunction of (2.1), (3.8), and (3.11) yields the closed-loop initial-value problem

$$\left. \begin{aligned} \dot{\bar{x}}(t) &= \bar{A}\bar{x}(t) - \nu(k(t))\bar{B}[y(t) - r(t)], & \bar{x}^0 &= \begin{bmatrix} x^0 \\ \xi^0 \end{bmatrix} \\ y(t) &= \bar{C}\bar{x}(t), \\ k(t) &= \frac{1}{1 - (\varphi(t)\|\bar{C}\bar{x}(t) - r(t)\|)^2}, \end{aligned} \right\} \quad (3.12)$$

where

$$\bar{A} = \begin{bmatrix} A & B\tilde{C} \\ 0 & \tilde{A} \end{bmatrix}, \quad \bar{B} = \begin{bmatrix} B \\ \tilde{B} \end{bmatrix}, \quad \bar{C} = [C, 0], \quad \bar{x}(t) = \begin{bmatrix} x(t) \\ \xi(t) \end{bmatrix}. \quad (3.13)$$

Noting the potential singularity in the function  $k$ , some care must be exercised in defining the concept of a solution of (3.12): a function  $\bar{x}: [0, \omega) \rightarrow \mathbb{R}^{n+mp}$ , with  $0 < \omega \leq \infty$ , is deemed a *solution* of (3.12) if, and only if, it is absolutely continuous, with  $\bar{x}(0) = \bar{x}^0$ , it satisfies the differential equations in (3.12) for almost all  $t \in [0, \omega)$ , and  $\varphi(t)\|\bar{C}\bar{x}(t) - r(t)\| < 1$  for all  $t \in [0, \omega)$ . A solution is *maximal* if, and only if, it has no proper right extension that is also a solution.

### 3.4 Main result

**Theorem 3.1** Let  $(A, B, C) \in \mathbb{R}^{n \times n} \times \mathbb{R}^{n \times m} \times \mathbb{R}^{m \times n}$  have strict relative degree one, sign-definite high-frequency gain, and be minimum-phase. Let  $\varphi$  satisfy (2.3), let  $\mathcal{F}_\varphi$  be the performance funnel (2.2) associated with  $\varphi$ , and let  $r \in \mathcal{Y}_{\text{ref}}(\alpha)$ . Then the feedback (3.11) applied in series with the internal model (3.8) yields the initial-value problem (3.12) which, for every  $x^0 \in \mathbb{R}^n$  and  $\xi^0 \in \mathbb{R}^{mp}$ , has a solution and every solution can be extended to a maximal solution. Every maximal solution  $\bar{x}: [0, \omega) \rightarrow \mathbb{R}^{n+mp}$  has the properties:

- (i)  $\omega = \infty$ ;
- (ii) the functions  $\bar{x}$ ,  $k$  and  $u$  are bounded;
- (iii) there exists  $\varepsilon \in (0, 1)$  such that, for all  $t \geq 0$ ,  $\varphi(t)\|y(t) - r(t)\| \leq 1 - \varepsilon$ ;
- (iv) if  $\varphi$  is unbounded, then  $(y(t) - r(t), u(t)) \rightarrow (0, 0)$  as  $t \rightarrow \infty$ .

**Remark 3.2** In the specific case of positive-definite  $CB$  and zero reference signal  $r \equiv 0$ , it is shown in [4] that the assertions of Theorem 3.1 hold for the feedback  $u = -ke$  without recourse to an internal model.

The proof of Theorem 3.1 invokes three lemmas; we briefly digress to present these.

### 3.5 Three technical lemmas

The first lemma is well known and is a re-statement of [3, Lemma 2.1.3].

**Lemma 3.3** Assume that  $(A, B, C) \in \mathbb{R}^{n \times n} \times \mathbb{R}^{n \times m} \times \mathbb{R}^{m \times n}$  has strict relative degree one. Let  $V \in \mathbb{R}^{n \times (n-m)}$  be such that  $\text{im } V = \ker C$  (of dimension  $n - m$ ) and write

$$N := (V^T V)^{-1} V^T [I_n - B(CB)^{-1} C].$$

Then

$$L = \begin{bmatrix} C \\ N \end{bmatrix}$$

is invertible, with inverse  $L^{-1} = [B(CB)^{-1}, V]$  and

$$LAL^{-1} = \begin{bmatrix} A_1 & A_2 \\ A_3 & A_4 \end{bmatrix}, \quad LB = \begin{bmatrix} CB \\ 0 \end{bmatrix}, \quad CL^{-1} = [I_m \quad 0]$$

where  $A_1 \in \mathbb{R}^{m \times m}$  (with  $A_2, A_3, A_4$  of conforming formats). Furthermore,  $(A, B, C)$  is minimum phase if, and only if,  $A_4$  is Hurwitz.

**Lemma 3.4** Let  $(A, B, C) \in \mathbb{R}^{n \times n} \times \mathbb{R}^{n \times m} \times \mathbb{R}^{m \times n}$  be minimum phase with strict relative degree one and sign-definite high-frequency gain. If  $(\tilde{A}, \tilde{B}, \tilde{C}, I_m)$  is a minimal realization of the internal model as specified in Subsection 3.1, then  $(\bar{A}, \bar{B}, \bar{C})$ , as defined in (3.13), is minimum phase with strict relative degree one and sign-definite high-frequency gain.

**Proof:** Clearly,  $\bar{C}\bar{B} = CB$  and so the system  $(\bar{A}, \bar{B}, \bar{C})$  has strict relative degree one and sign-definite high-frequency gain.

It remains to show that

$$\det \begin{bmatrix} sI - \bar{A} & \bar{B} \\ \bar{C} & 0 \end{bmatrix} \neq 0 \quad \text{for all } s \in \mathbb{C} \text{ with } \text{Re } s \geq 0.$$

Since  $(\hat{A}, \hat{b})$  is a controllable pair, the Hautus condition implies that  $[sI - \hat{A}, \hat{b}]$  has full rank  $p$  for all  $s \in \mathbb{C}$ , whence

$$\text{rank} \begin{bmatrix} sI - \tilde{A} & \tilde{B} \end{bmatrix} = mp \quad \text{for all } s \in \mathbb{C}.$$

By the minimum-phase property of  $(A, B, C)$ , we have

$$\text{rank} \begin{bmatrix} sI - A & B \\ C & 0 \end{bmatrix} = n + m \quad \text{for all } s \in \mathbb{C} \text{ with } \text{Re}(s) \geq 0,$$

and so

$$\text{rank} \begin{bmatrix} sI - \bar{A} & \bar{B} \\ \bar{C} & 0 \end{bmatrix} = \text{rank} \begin{bmatrix} sI - A & -B\tilde{C} & B \\ 0 & sI - \tilde{A} & \tilde{B} \\ C & 0 & 0 \end{bmatrix} = n + mp + m$$

for all  $s \in \mathbb{C}$  with  $\text{Re } s \geq 0$ , and the claim follows.  $\square$ .

A proof of the following lemma can be found in [8], see also [3, Lem. 5.1.2].



**Lemma 3.5** Let  $(A, B, C) \in \mathbb{R}^{n \times n} \times \mathbb{R}^{n \times m} \times \mathbb{R}^{m \times n}$  be minimum phase with strict relative degree one and sign-definite high-frequency gain. If  $(\bar{A}, \bar{B}, \bar{C})$  is a minimal realization of the internal model as specified in Subsection 3.1, then, for any  $r \in \mathcal{Y}_{\text{ref}}(\alpha)$ , there exists  $\rho^0 \in \mathbb{R}^{n+mp}$  such that

$$\left. \begin{aligned} \dot{\rho}(t) &= \bar{A} \rho(t), & \rho(0) &= \rho^0 \\ r(t) &= \bar{C} \rho(t), \end{aligned} \right\} \quad (3.14)$$

where  $\bar{A}$  and  $\bar{C}$  are given by (3.13).

### 3.6 Proof of Theorem 3.1

By Lemma 3.5, there exists  $\rho^0 \in \mathbb{R}^{n+mp}$  such that  $r(\cdot) = \bar{C} \rho(\cdot)$ , where  $\rho: t \mapsto (\exp \bar{A} t) \rho^0$ . Writing

$$x_e(t) = \bar{x}(t) - \rho(t), \quad e(t) = y(t) - r(r),$$

the closed-loop initial-value problem (3.12) may be expressed in the equivalent form

$$\left. \begin{aligned} \dot{x}_e(t) &= \bar{A} x_e(t) - \nu(k(t)) \bar{B} \bar{C} e(t), & x_e(0) &= x_e^0 := \bar{x}^0 - \rho^0, \\ e(t) &= \bar{C} x_e(t), \\ k(t) &= \left[ 1 - (\varphi(t) \|e(t)\|)^2 \right]^{-1}. \end{aligned} \right\} \quad (3.15)$$

Introducing the open set

$$\mathcal{D} := \{(x_e, \eta) \in \mathbb{R}^{n+mp} \times \mathbb{R} \mid \varphi(|\eta|) \|\bar{C} x_e\| < 1\},$$

and defining the function

$$d: \mathcal{D} \rightarrow \mathbb{R}_+, \quad (x_e, \eta) \mapsto d(x_e, \eta) := \frac{1}{1 - (\varphi(|\eta|) \|\bar{C} x_e\|)^2}, \quad (3.16)$$

then the non-autonomous closed-loop initial-value problem (3.12) (equivalently, (3.15)) may be recast on  $\mathcal{D}$  as the following autonomous initial-value problem

$$\left. \begin{aligned} \dot{x}_e(t) &= \bar{A} x_e(t) - \nu(d(x_e(t), \eta(t))) \bar{B} \bar{C} x_e(t) \\ \dot{\eta}(t) &= 1 \\ (x_e(0), \eta(0)) &= (x_e^0, 0) \in \mathcal{D}. \end{aligned} \right\} \quad (3.17)$$

The standard theory of ordinary differential equations now applies to conclude the existence of a solution  $t \mapsto (x_e(t), \eta(t)) \in \mathcal{D}$  to (3.17) and, moreover, every solution can be extended to a maximal solution  $(x_e, \eta): [0, \omega) \rightarrow \mathcal{D}$ . We will make use of the following fact in due course: if there exists a compact set  $\mathcal{C} \subset \mathcal{D}$  such that  $(x_e(t), \eta(t)) \in \mathcal{C}$  for all  $t \in [0, \omega)$ , then  $\omega = \infty$ . To see this, assume that such a set  $\mathcal{C}$  exists. Then  $(x_e(\cdot), \eta(\cdot))$  and  $\nu(d(x_e(\cdot), \eta(\cdot)))$  are bounded functions which, together with (3.17), implies that  $(x_e(\cdot), \eta(\cdot))$  is uniformly continuous. Seeking a contradiction, suppose that  $\omega < \infty$ . By uniform continuity, it follows that the limit  $(x_e^*, \omega) = \lim_{t \nearrow \omega} (x_e(t), \eta(t))$  exists and, by compactness, lies in  $\mathcal{C} \subset \mathcal{D}$ . By the existence theory, the initial-value problem (3.17), with initial

data  $(x_e^*, \omega)$  replacing  $(x_e^0, 0)$  has a solution: concatenation of this solution with  $(x_e, \eta)$  yields a proper right extension of the latter, contradicting its maximality.

Clearly, if  $(x_e, \eta): [0, \omega) \rightarrow \mathcal{D}$  is a solution of (3.17), then  $\bar{x} = x_e + \rho: [0, \omega) \rightarrow \mathbb{R}^{n+mp}$  is a solution of (3.12); conversely, if  $\bar{x}: [0, \omega) \rightarrow \mathbb{R}^{n+mp}$  is a solution of (3.12), then  $(x_e, \eta): [0, \omega) \rightarrow \mathbb{R}^{n+mp} \times \mathbb{R}$ , with  $x_e = \bar{x} - \rho$  and component  $\eta$  given by  $\eta(t) = t$ , is a solution of (3.17). We may now conclude that, for each  $\bar{x}^0 \in \mathbb{R}^{n+mp}$ , (3.12) has a solution and every solution can be maximally extended.

Let  $x^0 \in \mathbb{R}^n$  and  $\xi^0 \in \mathbb{R}^{mp}$  be arbitrary and let  $\bar{x}$  be a maximal solution of (3.12) with interval of existence  $[0, \omega)$ . Then, for  $x_e^0 = \bar{x}^0 - \rho^0$ , the function  $t \mapsto (x_e(t), \eta(t)) = (\bar{x}(t) - \rho(t), t)$  is a maximal solution of (3.17) with interval of existence  $[0, \omega)$ . By (3.16) and the first of equations (3.17), we now have

$$\dot{x}_e(t) = \bar{A}x_e(t) - \nu(k(t))\bar{B}\bar{C}e(t), \quad k(t) = \left[1 - (\varphi(t)\|e(t)\|)^2\right]^{-1} \quad \forall t \in [0, \omega). \quad (3.18)$$

By Lemma 3.4,  $(\bar{A}, \bar{B}, \bar{C})$  is minimum phase with strict relative degree one, and so, by Lemma 3.3, there exists  $N$  such that

$$L := \begin{bmatrix} \bar{C} \\ N \end{bmatrix}$$

is invertible and the transformation

$$\begin{bmatrix} \bar{C} \\ N \end{bmatrix} x_e(t) = \begin{bmatrix} e(t) \\ z(t) \end{bmatrix}$$

converts (3.18) into the equivalent form

$$\left. \begin{aligned} \dot{e}(t) &= A_1 e(t) + A_2 z(t) - \nu(k(t))CB e(t) \\ \dot{z}(t) &= A_3 e(t) + A_4 z(t) \\ k(t) &= \left[1 - (\varphi(t)\|e(t)\|)^2\right]^{-1} \end{aligned} \right\} \quad \forall t \in [0, \omega), \quad (3.19)$$

wherein  $A_4 \in \mathbb{R}^{(n+m(p-1)) \times (n+m(p-1))}$  is Hurwitz and we have invoked the equality  $\bar{C}\bar{B} = CB$ . Since  $(x_e(t), t) \in \mathcal{D}$  for all  $t \in [0, \omega)$ , we have

$$\varphi(t)\|e(t)\| < 1 \quad \forall t \in [0, \omega) \quad (3.20)$$

and so  $e$  is bounded, which, together with the Hurwitz property of  $A_4$  and the second of equations (3.19), implies that  $z$  is bounded. It immediately follows that  $x_e$  is bounded, whence boundedness of  $\bar{x} = x_e + \rho$ .

Writing  $e^0 = \bar{C}x_e^0$ ,  $z^0 = Nx_e^0$  and defining

$$q_0(t) := A_2 \exp(A_4 t)z^0, \quad q_1(t) := A_1 e(t) + A_2 \int_0^t \exp(A_4(t-s))A_3 e(s)ds, \quad \forall t \in [0, \omega), \quad (3.21)$$

then the first two equations in (3.19) are equivalent to

$$\dot{e}(t) = q_0(t) + q_1(t) - \nu(k(t))CB e(t) \quad \forall t \in [0, \omega). \quad (3.22)$$

Since  $A_4$  is Hurwitz, there exist  $c_1, \mu > 0$  such that

$$\|q_0(t)\| = \|A_2 \exp(A_4 t) z^0\| \leq c_1 e^{-\mu t} \quad \forall t \in [0, \omega] \quad (3.23)$$

and

$$\begin{aligned} \|q_1(t)\| &\leq \|A_1\| \|e(t)\| + c_1 \left( \int_0^{t/2} + \int_{t/2}^t \right) e^{-\mu(t-s)} \|e(s)\| ds \\ &\leq \|A_1\| \|e(t)\| + \frac{c_1}{\mu} \left[ e^{-\mu t/2} \max_{s \in [0, t/2]} \|e(s)\| + \max_{s \in [t/2, t]} \|e(s)\| \right] \quad \forall t \in [0, \omega]. \end{aligned} \quad (3.24)$$

By boundedness of  $e$ , together with (3.20) and invoking property (2.3d) of  $\varphi$ , we may infer the existence of  $c_2 > 0$  such that

$$\begin{aligned} \varphi(t) \dot{\varphi}(t) \|e(t)\|^2 &\leq c[1 + \varphi(t)] \varphi(t) \|e(t)\|^2 \\ &\leq c[1 + 2\varphi^2(t)] \|e(t)\|^2 \leq c[\|e(t)\|^2 + 2] \leq c_2 \quad \text{for almost all } t \in [0, \omega]. \end{aligned} \quad (3.25)$$

Since  $CB$  is sign definite, there exists  $c_3 > 0$  such that

$$\frac{1}{2} c_3 \|e\|^2 \leq |\langle e, CB e \rangle| \quad \forall e \in \mathbb{R}^m. \quad (3.26)$$

Now we are in a position to prove boundedness of  $k$ .

Define  $\tilde{\nu} : \mathbb{R} \rightarrow \mathbb{R}$  as follows

$$\tilde{\nu}(k) := \nu(k) s(CB).$$

By property (3.9) of  $\nu$ , there exists a strictly-increasing unbounded sequence  $(k_j)$  in  $(1, \infty)$  such that  $\tilde{\nu}(k_j) \rightarrow \infty$  as  $j \rightarrow \infty$ . Passing to a subsequence if necessary, we may assume that the sequence  $(\tilde{\nu}(k_j))$  is in  $(0, \infty)$  and is strictly increasing. Seeking a contradiction, suppose that  $k$  is unbounded. For each  $j \in \mathbb{N}$ , define

$$\begin{aligned} \tau_j &:= \inf\{t \in [0, \omega] \mid k(t) = k_{j+1}\} \\ \sigma_j &:= \sup\{t \in [0, \tau_j] \mid \tilde{\nu}(k(t)) = \tilde{\nu}(k_j)\} \\ \tilde{\sigma}_j &:= \sup\{t \in [0, \tau_j] \mid k(t) = k_j\} \leq \sigma_j. \end{aligned}$$

Observe that

$$k(\tau_j) > k(\sigma_j) \quad \forall j \in \mathbb{N}. \quad (3.27)$$

Furthermore, for all  $j \in \mathbb{N}$  and all  $t \in [\sigma_j, \tau_j]$ , we have  $k(t) \geq k_j$  and  $\tilde{\nu}(k(t)) \geq \tilde{\nu}(k_j)$ . Therefore,

$$1 > (\varphi(t) \|e(t)\|)^2 \geq 1 - \frac{1}{k_j} \geq 1 - \frac{1}{k_1} =: c_4 > 0 \quad \forall t \in [\sigma_j, \tau_j] \quad \forall j \in \mathbb{N}, \quad (3.28)$$

and since  $\varphi$  is non-decreasing, we arrive at

$$\max_{s \in [t/2, t]} \|e(s)\| < \frac{1}{\varphi(t/2)} \leq \frac{\varphi(t)}{\sqrt{c_4} \varphi(t/2)} \|e(t)\| \quad \forall t \in [\sigma_j, \tau_j] \quad \forall j \in \mathbb{N}. \quad (3.29)$$

By (3.24) and (3.29), together with boundedness of  $e$  and property (2.3c) of  $\varphi$ , we may infer the existence of  $c_5 > 0$  such that

$$\|q_1(t)\| \leq c_5 \left[ e^{-\mu t/2} + \|e(t)\| \right] \quad \forall t \in [\sigma_j, \tau_j] \quad \forall j \in \mathbb{N}. \quad (3.30)$$

Invoking (3.25), (3.23), (3.26), (3.30), (3.28), recalling that  $\varphi(t)\|e(t)\| < 1$  for all  $t \in [0, \omega)$ , and noting that, by Proposition 2.1, the functions  $t \mapsto \varphi(t)e^{-\mu t}$  and  $t \mapsto \varphi(t)e^{-\mu t/2}$  are bounded, we may conclude the existence of  $c_6 > 0$  such that

$$\begin{aligned} \frac{d}{dt} k(t) &= k^2(t) \left[ 2\varphi(t)\dot{\varphi}(t)\|e(t)\|^2 + 2\varphi^2(t)\langle e(t), q_0(t) + q_1(t) - \nu(k(t))CB e(t) \rangle \right] \\ &\leq k^2(t) \left[ 2c_2 + 2\varphi(t) [\|q_0(t)\| + \|q_1(t)\|] - 2\varphi^2(t)\tilde{\nu}(k(t))|\langle e(t), CB e(t) \rangle| \right] \\ &\leq k^2(t) \left[ 2c_2 + 2c_1\varphi(t)e^{-\mu t} + 2c_5\varphi(t) \left[ e^{-\mu t/2} + \|e(t)\| \right] - c_3\varphi^2(t)\tilde{\nu}(k(t))\|e(t)\|^2 \right] \\ &\leq k^2(t) [c_6 - c_3c_4\tilde{\nu}(k(t))] \quad \text{for almost all } t \in [\sigma_j, \tau_j] \text{ and all } j \in \mathbb{N}. \end{aligned}$$

Let  $j^* \in \mathbb{N}$  be sufficiently large to that  $c_6 - c_3c_4\tilde{\nu}(k_{j^*}) < 0$ . Then,

$$\frac{d}{dt} k(t) < 0 \quad \text{for almost all } t \in [\sigma_{j^*}, \tau_{j^*}],$$

which contradicts (3.27). This proves boundedness of  $k$ .

Next we show boundedness of  $u$ . Since  $k$  is bounded, there exists  $\varepsilon > 0$  such that  $\varphi(t)\|e(t)\| \leq 1 - \varepsilon$  for all  $t \in [0, \omega)$ . By boundedness of  $e$ ,  $z$ , and  $k$ , it follows that  $u$  is bounded.

We proceed to prove that  $\omega = \infty$ . Suppose that  $\omega$  is finite. Let  $c_7 > 0$  be such that  $\|x_e(t)\| \leq c_7$  for all  $t \in [0, \omega)$ , and set

$$\mathcal{C} := \{ (x_e, \eta) \in \mathcal{D} \mid \varphi(|\eta|)\|\tilde{C}x_e\| \leq 1 - \varepsilon, \quad \|x_e\| \leq c_7, \quad \eta \in [0, \omega] \}.$$

Then  $\mathcal{C}$  is a compact subset of  $\mathcal{D}$  and contains the trajectory of the maximal solution  $t \mapsto (x_e(t), t)$  of (3.17). Therefore, the supposition that  $\omega$  is finite is false. This completes the proof of Assertions (i)-(iii).

It remains only to establish Assertion (iv). Assume that  $\varphi$  is unbounded. Then  $\|e(t)\| < 1/\varphi(t) \rightarrow 0$  as  $t \rightarrow \infty$ . By boundedness of  $k$ , we have  $u(t) = -\nu(k(t))e(t) \rightarrow 0$  as  $t \rightarrow \infty$ .  $\square$

## 4 Example

Let  $(A, b, c)$  be a single-input, single-output minimum-phase system with positive high-frequency gain  $cb > 0$ . Assume that the class of reference signals  $r : \mathbb{R}_+ \rightarrow \mathbb{R}$  comprises all linear combinations of constant functions and sinusoidal functions of period  $2\pi$ . Choosing as internal model the linear system with transfer function

$$\frac{\beta(s)}{\alpha(s)} = \frac{(s+1)^3}{s(s^2+1)},$$

and selecting the funnel function  $t \mapsto \varphi(t) := t^2$ , then the feedback

$$u(t) = -k(t)e(t), \quad k(t) = \frac{1}{1 - (t^2 e(t))^2}, \quad e(t) = y(t) - r(t),$$

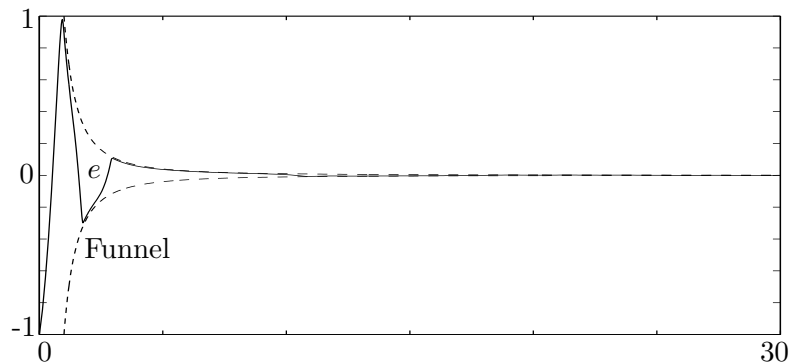
in series with the internal model, ensures asymptotic tracking of every admissible reference signal  $r$  and achieves a tracking error decay rate of the order  $t^{-2}$ . In the specific case

$$A = \begin{bmatrix} 1 & 1 & 1 \\ 1 & -1 & 0 \\ 1 & 0 & -1 \end{bmatrix}, \quad b = \begin{bmatrix} 1 \\ 0 \\ 0 \end{bmatrix}, \quad c = [1 \ 0 \ 0],$$

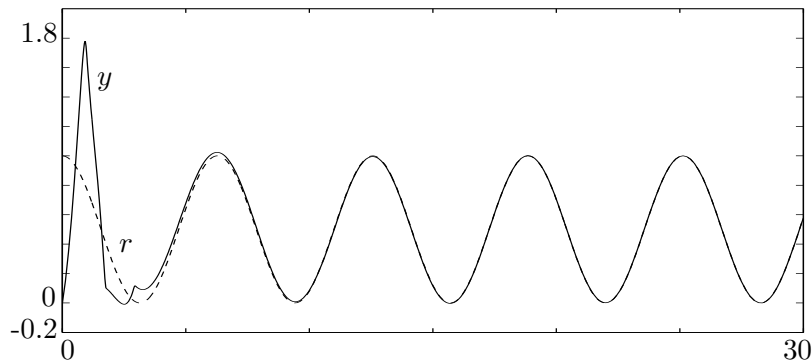
with zero initial conditions and reference signal

$$r : t \mapsto \frac{1}{2}[1 + \cos t],$$

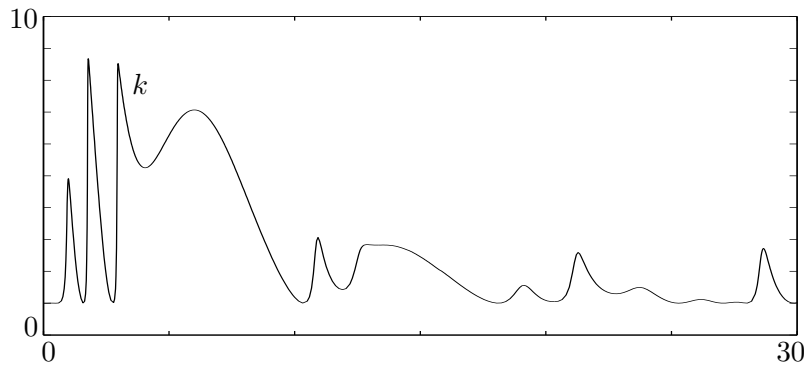
the behaviour of the feedback system is depicted in Figure 3(a-d).



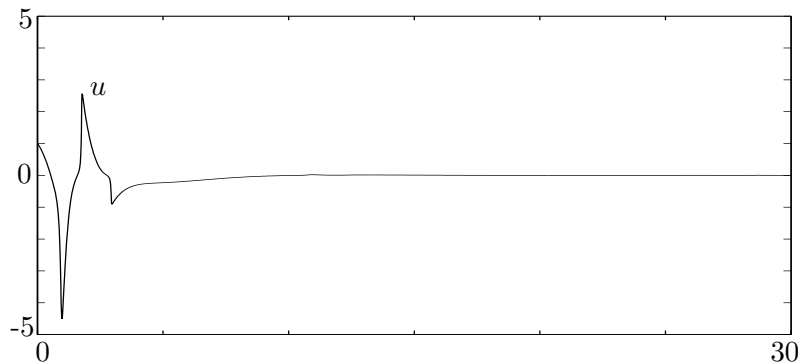
(a) The funnel and tracking error  $e$



(b) The reference  $r$  and output  $y$



(c) The gain function  $k$



(d) The control  $u$

Figure 3: Example

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