## Text Supplementary 1 - A Critical Quantity for Noise Attenuation in Feedback Systems

Liming Wang, Jack Xin, Qing Nie
Center for Mathematical and Computational Biology, Center for Complex Biological Systems, and Department of Mathematics, University of California at Irvine, Irvine, CA 92697, United States of America

* E-mail: qnie@math.uci.edu


## 1 Single-Positive-Loop Module

The following ordinary differential equations can be used to describe the dynamics of a single-positive-loop module (inside the red dashed box of Figure 1A):

$$
\begin{align*}
\frac{d c^{*}}{d t} & =k_{1}^{*} b^{*}\left(c_{\mathrm{tot}}-c^{*}\right)-k_{2}^{*} c^{*}+k_{3}^{*},  \tag{1}\\
\frac{d b^{*}}{d t} & =\left(k_{c}^{*} s^{*} c^{*}\left(b_{\mathrm{tot}}-b^{*}\right)-k_{5}^{*} b^{*}+k_{4}^{*}\right) \tau_{b}^{*}
\end{align*}
$$

Here, $c^{*}$ and $b^{*}$ are the active levels of the corresponding components, varying between $\left[0, c_{\text {tot }}\right]$ and $\left[0, b_{\text {tot }}\right]$, respectively. The variable $s^{*}$ stands for the input, whose maximum is $s_{\max }$. For the convenience of analysis, we assume that the total amount of the active form and inactive form are conserved in each component, as in previous works [2, 13]. Later on, we test our conclusion in other models with degradation/dilution terms (e.g. a yeast cell polarization model and a polymyxin B resistance model in enteric bacteria).

Next, we rescale the variables by defining

$$
c=\frac{c^{*}}{c_{\mathrm{tot}}}, \quad b=\frac{b^{*}}{b_{\mathrm{tot}}}, \quad s=\frac{s^{*}}{s_{\max }} .
$$

Now, equations in (1) becomes

$$
\begin{align*}
& \frac{d c}{d t}=k_{1}^{*} b_{\mathrm{tot}} b(1-c)-k_{2}^{*} c+\frac{k_{3}^{*}}{c_{\mathrm{tot}}},  \tag{2}\\
& \frac{d b}{d t}=\left(\frac{k_{c}^{*}}{k_{5}^{*}} s \max c_{\mathrm{tot}} s c(1-b)-b+\frac{k_{4}^{*}}{k_{5}^{*} b_{\mathrm{tot}}}\right) \tau_{b}^{*} k_{5}^{*} .
\end{align*}
$$

Let

$$
k_{1}=k_{1}^{*} b_{\mathrm{tot}}, \quad k_{2}=k_{2}^{*}, \quad k_{3}=\frac{k_{3}^{*}}{c_{\mathrm{tot}}}, \quad k_{c}=\frac{k_{c}^{*}}{k_{5}^{*}} s_{\max } c_{\mathrm{tot}}, \quad k_{4}=\frac{k_{4}^{*}}{k_{5}^{*} b_{\mathrm{tot}}}, \quad \tau_{b}=\tau_{b}^{*} k_{5}^{*},
$$

and (2) becomes

$$
\begin{align*}
& \frac{d c}{d t}=k_{1} b(1-c)-k_{2} c+k_{3},  \tag{3}\\
& \frac{d b}{d t}=\left(k_{c} s c(1-b)-b+k_{4}\right) \tau_{b} .
\end{align*}
$$

Conditions for a "switch-like" response When the stimulus $s$ goes from an "off" state to an "on" state, $C$ usually responds in a similar fashion. In order for the two response states in $C$ to separate significantly like a "switch", the following three conditions must be satisfied:

- $k_{4} \ll 1$, and $k_{3} \ll k_{2}$. This set of constraints means the basal ativation levels of $C$ and $B$ must be relatively low compared to the deactivation rate of $C$ and $B$, respectively.
- $k_{1} / k_{2} \ll 1 / k_{4}$. This implies that when the signal is off, the activation from $B$ to $C$, which is a product of $k_{1}$ and $k_{4}$ (the level of $B$ at $s=0$ ) must be significantly less than the deactivation of $C$, $k_{2}$.
- $k_{1} / k_{2}>1 / k_{c}$. This suggests that the strength of activation from $B$ to $C$, measured by $k_{1} / k_{2}$, is greater than the deactivation of $B$ to produce $C$, measured by $1 / k_{c}$. That is, more active $C$ should be produced from $B$ than the $C$ that participates in the activation of $B$.

The second and third conditions together imply $k_{4} \ll k_{c}$, that is, $B$ should be mainly activated through
$C$ not from basal activation. The details on derivation of the three conditions are as follows. We use ( $\cdot \cdot$ ) to denote the steady steady value of a variable. At the steady state, the right hand side of system (3) equals zero, yielding

$$
\bar{b}=\frac{k_{c} s \bar{c}+k_{4}}{k_{c} s \bar{c}+1},
$$

where $\bar{c}$ satisfies

$$
\begin{equation*}
A_{2} \bar{c}^{2}+A_{1} \bar{c}-A_{0}=0 \tag{4}
\end{equation*}
$$

with

$$
A_{2}=\left(k_{1}+k_{2}\right) k_{c} s, \quad A_{1}=k_{2}-k_{3} k_{c} s-k_{1} k_{c} s+k_{1} k_{4}, \quad A_{0}=k_{1} k_{4}+k_{3} .
$$

Let $\bar{c}_{0}$ and $\bar{b}_{0}\left(\bar{c}_{1}\right.$ and $\left.\bar{b}_{1}\right)$ denote the steady states of $c$ and $b$ at $s=0(s=1)$, respectively. When $s=0$, equation (4) is linear, and thus

$$
\bar{c}_{0}=\frac{k_{1} k_{4}+k_{3}}{k_{1} k_{4}+k_{2}}, \quad \bar{b}_{0}=k_{4} .
$$

It is clear that $\bar{c}_{0}$ and $\bar{b}_{0}$ are close to zero provided

$$
\begin{equation*}
k_{4} \ll 1, \quad k_{3} \ll k_{2}, \quad k_{1} k_{4} \ll k_{2} . \tag{5}
\end{equation*}
$$

When $s \neq 0$, (4) is a quadratic equation with two real roots of different signs. The positive root can be written as

$$
\begin{equation*}
\bar{c}=\frac{-A_{1}+\sqrt{A_{1}^{2}+4 A_{2} A_{0}}}{2 A_{2}} . \tag{6}
\end{equation*}
$$

Under condition (5) and the assumption $k_{1} \gg k_{3}$, we have $A_{1} \gg A_{0}$ and that $A_{1}$ and $A_{2}$ are about the same order (they both contain terms like $k_{2}$ and $k_{1} k_{c} s$ ). Thus, $A_{1}^{2}+4 A_{2} A_{0} \approx A_{1}^{2}$, and depending on the sign of $A_{1}$ there are two cases of $\bar{c}$ at $s=1$.

1. $k_{1} k_{c} \leq k_{2}$. In this case, $\sqrt{A_{1}^{2}+4 A_{2} A_{0}} \approx A_{1}$, and by (6), $\bar{c}_{1} \approx 0$, which is close to the inactive steady state $\bar{c}_{0}$, violating our "switch-like" response assumption. This unrealistic case should be avoided.
2. $k_{1} k_{c}>k_{2}$. In this case, $\sqrt{A_{1}^{2}+4 A_{2} A_{0}} \approx-A_{1}$, and by (6),

$$
\begin{equation*}
\bar{c}_{1} \approx \frac{-2 A_{1}}{2 A_{2}} \approx \frac{k_{1} k_{c}-k_{2}}{\left(k_{1}+k_{2}\right) k_{c}} . \tag{7}
\end{equation*}
$$

We remark that the same condition is also obtained from the Fluctuation Dissipation Theorem approach in Section 5. The corresponding steady state value of $b$ at $s=1$ is

$$
\begin{equation*}
\bar{b}_{1} \approx \frac{k_{1} k_{c}-k_{2}}{k_{1} k_{c}+k_{1}} . \tag{8}
\end{equation*}
$$

Let us denote $J$ as the Jacobian matrix of system (3):

$$
J(s, c, b)=\left(\begin{array}{cc}
-k_{1} b-k_{2} & k_{1}(1-c) \\
k_{c} s(1-b) \tau_{b} & -\left(k_{c} s c+1\right) \tau_{b}
\end{array}\right) .
$$

### 1.1 The deactivation time scale

Our goal here is to estimate the time scale $t_{1 \rightarrow 0}$ and to find the parameters that controls $t_{1 \rightarrow 0}$. We start by defining the deactivation process. Suppose that system (3) is well stabilized at the active state, that is, the initial condition of system (3) is ( $\bar{c}_{1}, \bar{b}_{1}$ ), and then we apply the signal $s(t)=0, t \geq 0$. The process from $t=0$ till the output is stabilized around the equilibrium $\bar{c}_{0}$ is called the deactivation process. The deactivation time scale $t_{1 \rightarrow 0}$ is defined as the time when $c(t)$ reaches $\left(\bar{c}_{1}-\bar{c}_{0}\right) / e+\bar{c}_{0}$. During the whole deactivation process, system (3) sees only signal $s \equiv 0$, so we linearize (3) around the steady state at $s=0$ :

$$
\begin{equation*}
\binom{\delta c(t)}{\delta b(t)}^{\prime}=J\left(0, \bar{c}_{0}, \bar{b}_{0}\right)\binom{\delta c(t)}{\delta b(t)} \tag{9}
\end{equation*}
$$

where

$$
J\left(0, \bar{c}_{0}, \bar{b}_{0}\right)=\left(\begin{array}{cc}
-k_{1} \bar{b}_{0}-k_{2} & k_{1}\left(1-\bar{c}_{0}\right) \\
0 & -\tau_{b}
\end{array}\right) .
$$

Clearly, the eigenvalues of $J\left(0, \bar{c}_{0}, \bar{b}_{0}\right)$ are

$$
\lambda_{c}=-\left(k_{1} k_{4}+k_{2}\right), \quad \lambda_{b}=-\tau_{b}
$$

and the corresponding eigenvectors are

$$
\xi_{c}=\binom{1}{0}, \quad \xi_{b}=\binom{\frac{k_{1}\left(1-\bar{c}_{0}\right)}{k_{1} k_{4}+k_{2}-\tau_{b}}}{1}
$$

Solutions of system (9) can be written as linear combination of $\xi_{c} e^{-k_{2} t}, \xi_{b} e^{-\tau_{b} t}$, that is,

$$
\begin{align*}
\delta c(t) & =l_{c} \xi_{c}^{1} e^{-\left(k_{1} k_{4}+k_{2}\right) t}+l_{b} \xi_{b}^{1} e^{-\tau_{b} t}  \tag{10}\\
\delta b(t) & =l_{b} \xi_{b}^{2} e^{-\tau_{b} t} \tag{11}
\end{align*}
$$

for some constants $l_{c}$ and $l_{b}$ depending on initial conditions [11]. Here, $\xi_{b}^{1}$ denotes the first coordinate of $\xi_{b}$.

From (10)-(11), it is easy to see that without loop $B$ the time scale of $C$ is around $1 / k_{2}$, whereas with loop $B$ the time scale would be determined by the larger one between $1 / k_{2}$ and $1 / \tau_{b}$. Therefore, in order to obtain a much slower time scale of $C$ when $B$ is present, $\tau_{b}$ needs to be much smaller than $k_{2}$. From now on, we assume that $\tau_{b} \ll k_{2}$. Note that when $\tau_{b} \ll k_{2}$,

$$
\xi_{b}^{1} \approx \frac{k_{1}\left(k_{2}-k_{3}\right)}{\left(k_{1} k_{4}+k_{2}\right)^{2}}
$$

For an arbitrary initial condition $\left(\gamma_{0}^{*}, \beta_{0}^{*}\right)$, we have

$$
l_{c} \approx \gamma_{0}^{*}-\beta_{0}^{*} \frac{k_{1}\left(k_{2}-k_{3}\right)}{\left(k_{1} k_{4}+k_{2}\right)^{2}}, \quad l_{b}=\beta_{0}^{*}
$$

and thus

$$
\begin{equation*}
\delta c(t) \approx\left(\gamma_{0}^{*}-\beta_{0}^{*} \frac{k_{1}\left(k_{2}-k_{3}\right)}{\left(k_{1} k_{4}+k_{2}\right)^{2}}\right) e^{-\left(k_{1} k_{4}+k_{2}\right) t}+\beta_{0}^{*} \frac{k_{1}\left(k_{2}-k_{3}\right)}{\left(k_{1} k_{4}+k_{2}\right)^{2}} e^{-\tau_{b} t} . \tag{12}
\end{equation*}
$$

On the other hand, under the assumption $\tau_{b} \ll k_{2}$, there is a time scale separation in (10). Throughout the initial layer (the time interval on the order of $1 / k_{2}$ ), we can view $l_{b} \xi_{b}^{1} e^{-\tau_{b} t}$ as a constant $l_{b} \xi_{b}^{1}$, while the first term in (10) decays quickly. As a result, after the initial layer, $\delta c$ falls from its initial condition $c_{0}$ to a neighborhood of $l_{b} \xi_{b}^{1}$. If $l_{b} \xi_{b}^{1}$ is small comparing to $c_{0}$, then $\delta c$ has changed substantially during the initial layer, and thus the effect from $e^{-\tau_{b} t}$ is diminished. Consider the extreme case when $l_{b} \xi_{b}^{1}$ is below $\left(\bar{c}_{1}-\bar{c}_{0}\right) / e$. The deactivation time scale $t_{1 \rightarrow 0}$ is defined as the time when $c(t)$ reaches $\left(\bar{c}_{1}-\bar{c}_{0}\right) / e+\bar{c}_{0}$, which when linearized around $\bar{c}_{0}$ corresponds to $\delta c(t)$ being $\left(\bar{c}_{1}-\bar{c}_{0}\right) / e$. If $l_{b} \xi_{b}^{1}$ is below $\left(\bar{c}_{1}-\bar{c}_{0}\right) / e$, then $\delta c$ reaches $\left(\bar{c}_{1}-\bar{c}_{0}\right) / e$ during the initial layer, i.e., $t_{1 \rightarrow 0}$ is within the order of $1 / k_{2}$. So when $\tau_{b}$ is small, the deactivation process may not be significantly affected by the term $e^{-\tau_{b} t}$. Conversely, if $l_{b} \xi_{b}^{1} / \gamma_{0}^{*}$ is large, $e^{-\tau_{b} t}$ will have substantial contribution in determining $t_{1 \rightarrow 0}$. We remark that it is possible in general (for example, the positive-negative-loop case) to have $l_{b} \xi_{b}^{1}$ greater than the initial state (the active state), and in that case $\delta c(t)$ increases during the initial layer, resulting in a even longer period for $\delta c(t)$ to reach $\left(\bar{c}_{1}-\bar{c}_{0}\right) / e$. In order to estimate $l_{b} \xi_{b}^{1}$, notice that $l_{b} \xi_{b}^{2}=\beta_{0}^{*}$ and $l_{b} \xi_{b}^{1} / \gamma_{0}^{*}=\left(\beta_{0}^{*} \xi_{b}^{1}\right) /\left(\gamma_{0}^{*} \xi_{b}^{2}\right)$. Therefore, for any fixed initial condition $\left(\gamma_{0}^{*}, \beta_{0}^{*}\right)$ around the steady state $(0,0)$, the ratio $\xi_{b}^{1} / \xi_{b}^{2}$ determines how effectively $e^{-\tau_{b} t}$ slows down the dynamics of $\delta c$. Calculating this ratio, we obtain

$$
\frac{\xi_{b}^{1}}{\xi_{b}^{2}}=\frac{k_{1}\left(1-\bar{c}_{0}\right)}{k_{1} k_{4}+k_{2}-\tau_{b}} \stackrel{\tau_{b} \lll k_{2}}{\approx} \frac{k_{1}\left(k_{2}-k_{3}\right)}{\left(k_{1} k_{4}+k_{2}\right)^{2}}{ }^{k_{3} \ll k_{2}} \frac{k_{1} k_{2}}{\left(k_{1} k_{4}+k_{2}\right)^{2}} \stackrel{k_{1} k_{4} \ll k_{2}}{\approx} \frac{k_{1} k_{2}}{k_{2}^{2}}=K_{a} .
$$

The bigger $K_{a}$ is, the more contribution from the slow term $e^{-\tau_{b} t}$, and thus the slower $\delta c$ converges to the steady state at zero.

To estimate $t_{1 \rightarrow 0}$, we consider the initial condition $\left(\gamma_{0}^{*}, \beta_{0}^{*}\right)=\left(\bar{c}_{1}-\bar{c}_{0}, \bar{b}_{1}-\bar{b}_{0}\right)$, and compute the time when $\delta c$ reaches $\gamma_{0}^{*} / e$. If $l_{b} \xi_{b}^{1} / \gamma_{0}^{*}$ is large, the initial layer can be neglected, and $\delta c(t) \approx l_{b} \xi_{b}^{1} e^{-\tau_{b} t}$. The time
it takes for $\delta c$ to drop from $\gamma_{0}^{*}$ to $\gamma_{0}^{*} / e$ is approximately,

$$
\frac{1+\ln \left(\frac{l_{b} \xi_{b}^{1}}{\gamma_{0}^{*}}\right)}{\tau_{b}}=\frac{1+\ln \left(\frac{\xi_{b}^{1} \beta_{0}^{*}}{\xi_{b}^{2} \gamma_{0}^{*}}\right)}{\tau_{b}} .
$$

So, the quantity, $\frac{\xi_{b}^{1} \beta_{0}^{*}}{\xi_{b}^{2} \gamma_{0}^{*}}$ affects the duration of the deactivation process, with bigger $\frac{\xi_{b}^{1} \beta_{0}^{*}}{\xi_{b}^{2} \gamma_{0}^{*}}$ corresponding to longer deactivation. Notice that,

$$
\frac{\beta_{0}^{*}}{\gamma_{0}^{*}}=\frac{\bar{b}_{1}-\bar{b}_{0}}{\bar{c}_{1}-\bar{c}_{0}} \approx\left(\frac{K_{a}+1}{K a}\right)\left(\frac{k_{c}}{k_{c}+1}\right) .
$$

Thus,

$$
\begin{equation*}
\frac{\xi_{b}^{1} \beta_{0}^{*}}{\xi_{b}^{2} \gamma_{0}^{*}} \approx\left(K_{a}+1\right) \frac{k_{c}}{k_{c}+1} . \tag{13}
\end{equation*}
$$

Therefore, $t_{1 \rightarrow 0}$ is increasing in $K_{a}$ and $k_{c}$. The above linear approximation (10)-(11) is only valid around a small neighborhood of $(0,0)$, and the global dynamics could be different. However, demonstrated by numerical simulations, the linear analysis seems to provide good qualitative predictions of the dynamics and the dependence on $K_{a}$ and $k_{c}$ (Figures 3E-3F). For the sake of completeness, we also used the two-time-scale asymptotic expansion of the global solution (Section 4) and the Fluctuation Dissipation Theorem (Section 5). Both approaches reinforced the conclusions obtained by the linear analysis.

### 1.2 The activation time scale

Suppose that system (3) is well stabilized at the inactive state, at time zero, apply the signal $s(t)=1, t \geq 0$. The process from $t=0$ till the system is stabilized around the equilibrium $\bar{c}_{1}$ is defined as the activation process. During activation, the input signal is $s \equiv 1$, so we linearize (3) around the steady state at $s=1$ :

$$
\binom{\delta c(t)}{\delta b(t)}^{\prime}=J\left(1, \bar{c}_{1}, \bar{b}_{1}\right)\binom{\delta c(t)}{\delta b(t)}
$$

where

$$
J\left(1, \bar{c}_{1}, \bar{b}_{1}\right)=\left(\begin{array}{cc}
-k_{1} \bar{b}_{1}-k_{2} & k_{1}\left(1-\bar{c}_{1}\right) \\
k_{c}\left(1-\bar{b}_{1}\right) \tau_{b} & -\left(k_{c} \bar{c}_{1}+1\right) \tau_{b}
\end{array}\right)
$$

Let us assume that the term $k_{c}\left(1-\bar{b}_{1}\right) \tau_{b}$ is small enough $\left(\tau_{b} \ll k_{2}\right)$ so that the eigenvalues of $J\left(1, \bar{c}_{1}, \bar{b}_{1}\right)$ are

$$
\lambda_{c} \approx-\left(k_{1} \bar{b}_{1}+k_{2}\right), \quad \lambda_{b} \approx-\left(k_{c} \bar{c}_{1}+1\right) \tau_{b}
$$

and their corresponding eigenvectors are

$$
\xi_{c} \approx\binom{1}{0}, \quad \xi_{b} \approx\binom{\frac{k_{1}\left(1-\bar{c}_{1}\right)}{k_{1} \bar{b}_{1}+k_{2}-\left(k_{c} \bar{c}_{1}+1\right) \tau_{b}}}{1} .
$$

Therefore,

$$
\begin{align*}
& \delta c(t) \approx l_{c} \xi_{c}^{1} e^{-\left(k_{1} \bar{b}_{1}+k_{2}\right) t}+l_{b} \xi_{b}^{1} e^{-\left(k_{c} \bar{c}_{1}+1\right) \tau_{b} t},  \tag{14}\\
& \delta b(t) \approx l_{b} \xi_{b}^{2} e^{-\left(k_{c} \bar{c}_{1}+1\right) \tau_{b} t}
\end{align*}
$$

for some constants $l_{c}$ and $l_{b}$ depending on initial conditions. If $\left(k_{c} \bar{c}_{1}+1\right) \tau_{b} \ll k_{2}$, there is a time scale separation in the solution of $\delta c$. Similar to the analysis for the deactivation time scale, if the ratio

$$
\frac{\xi_{b}^{1}}{\xi_{b}^{2}} \approx \frac{k_{1}\left(1-\bar{c}_{1}\right)}{k_{1} \bar{b}_{1}+k_{2}}=\frac{K_{a}}{\left(K_{a}+1\right)^{2}}\left(1+\frac{1}{k_{c}}\right)^{2},
$$

is large, then $e^{-\left(k_{c} \bar{c}_{1}+1\right) \tau_{b} t}$ determines the time scale of the convergence to zero. To estimate $t_{0 \rightarrow 1}$, we consider the initial condition $\left(\gamma_{0}^{*}, \beta_{0}^{*}\right)=\left(\bar{c}_{0}-\bar{c}_{1}, \bar{b}_{0}-\bar{b}_{1}\right)$, and define $t_{0 \rightarrow 1}$ as the time when $\delta c$ reaches $\gamma_{0}^{*} / e$. If $l_{b} \xi_{b}^{1} / \gamma_{0}^{*}$ is large, then the initial layer can be neglected, and $\delta c(t) \approx l_{b} \xi_{b}^{1} e^{-\left(k_{c} \bar{c}_{1}+1\right) \tau_{b} t}$. The time it takes for $\delta c$ to drop from $\gamma_{0}^{*}$ to $\gamma_{0}^{*} / e$ is

$$
\frac{1+\ln \left(\frac{\xi_{b}^{1} \beta_{0}^{*}}{\xi_{b}^{2} \gamma_{0}^{*}}\right)}{\tau_{b}}
$$

where

$$
\begin{equation*}
\frac{\xi_{b}^{1} \beta_{0}^{*}}{\xi_{b}^{2} \gamma_{0}^{*}} \approx \frac{1}{K_{a}+1}\left(1+\frac{1}{k_{c}}\right) . \tag{15}
\end{equation*}
$$

So, $t_{0 \rightarrow 1}$ is decreasing in $K_{a}$ and $k_{c}$.

## 2 Positive-Positive-Loop Module

Recall the positive-positive-loop model:

$$
\begin{align*}
& \frac{d c}{d t}=k_{1}(a+b)(1-c)-k_{2} c+k_{3}, \\
& \frac{d a}{d t}=\left(k_{c} s c(1-a)-a+k_{4}\right) \tau_{a},  \tag{16}\\
& \frac{d b}{d t}=\left(k_{c} s c(1-b)-b+k_{4}\right) \tau_{b} .
\end{align*}
$$

We have similar conditions for a "switch-like" response as in the single-positive-loop case:

- $k_{4} \ll 1$, and $k_{3} \ll k_{2}$. This set of constraints means the basal activation levels of $C$ and $B$ must be relatively low compared to the deactivation rate of $C$ and $B$, respectively.
- $2 k_{1} / k_{2} \ll 1 / k_{4}$. This implies that when the signal is off, the activation from $A$ and $B$ to $C$, which is a product of $2 k_{1}$ and $k_{4}$, the level of $A$ or $B$ at $s=0$, must be significantly less than the deactivation of $C, k_{2}$.
- $2 k_{1} / k_{2}>1 / k_{c}$. This suggests that the strength of activation from $A$ and $B$ to $C$, measured by $2 k_{1} / k_{2}$, is greater than the deactivation of $B$ to produce active $C$, measured by $1 / k_{c}$. That is, more $C$ should be activated from $A$ and $B$ than the $C$ that participates in the activation of $A$ and $B$.

The second and the third conditions together imply $k_{4} \ll k_{c}$, that is, $A$ and $B$ should be mainly activated through $C$ not from basal activation. The details on derivation of the three conditions are as follows. At
the steady state of system (16), we have

$$
\bar{a}=\bar{b}=\frac{k_{c} s \bar{c}+k_{4}}{k_{c} s \bar{c}+1},
$$

where $\bar{c}$ satisfies

$$
\begin{equation*}
A_{2} \bar{c}^{2}+A_{1} \bar{c}-A_{0}=0 \tag{17}
\end{equation*}
$$

with

$$
A_{2}=\left(2 k_{1}+k_{2}\right) k_{c} s, \quad A_{1}=k_{2}-k_{3} k_{c} s-2 k_{1} k_{c} s+2 k_{1} k_{4}, \quad A_{0}=2 k_{1} k_{4}+k_{3} .
$$

When $s=0$, equation (17) is linear, and the steady state is

$$
\bar{c}_{0}=\frac{2 k_{1} k_{4}+k_{3}}{2 k_{1} k_{4}+k_{2}}, \quad \bar{a}_{0}=k_{4}, \quad \bar{b}_{0}=k_{4} .
$$

The values of $\bar{c}_{0}, \bar{a}_{0}$, and $\bar{b}_{0}$ are close to zero provided

$$
\begin{equation*}
k_{4} \ll 1, \quad k_{3} \ll k_{2}, \quad 2 k_{1} k_{4} \ll k_{2} . \tag{18}
\end{equation*}
$$

When $s \neq 0,(17)$ is a quadratic equation with two real roots of different signs. The positive root can be written as

$$
\begin{equation*}
\bar{c}=\frac{-A_{1}+\sqrt{A_{1}^{2}+4 A_{2} A_{0}}}{2 A_{2}} . \tag{19}
\end{equation*}
$$

Similar to the single-positive-loop case, under condition (18), there are two cases at $s=1$.

1. $2 k_{1} k_{c} \leq k_{2}$. In this case, the positive root $\bar{c}_{1}$ as in (19) is close to zero, and thus can not be differentiated from $\bar{c}_{0}$. We discard this unrealistic case.
2. $2 k_{1} k_{c}>k_{2}$. In this case, $\sqrt{A_{1}^{2}+4 A_{2} A_{0}} \approx-A_{1}$, and thus,

$$
\bar{c}_{1} \approx \frac{2 k_{1} k_{c}-k_{2}}{\left(2 k_{1}+k_{2}\right) k_{c}} .
$$

The corresponding steady state values of $a$ and $b$ are given by

$$
\bar{a}_{1}=\bar{b}_{1} \approx \frac{2 k_{1} k_{c}-k_{2}}{2 k_{1}\left(k_{c}+1\right)} .
$$

The Jacobian matrix of system (16) is

$$
J(s, c, a, b)=\left(\begin{array}{ccc}
-k_{1}(a+b)-k_{2} & k_{1}(1-c) & k_{1}(1-c) \\
k_{c} s(1-a) \tau_{a} & -\left(k_{c} s c+1\right) \tau_{a} & 0 \\
k_{c} s(1-b) \tau_{b} & 0 & -\left(k_{c} s c+1\right) \tau_{b}
\end{array}\right) .
$$

### 2.1 The deactivation time scale

Similar to the single-positive-loop analysis, we linearize system (16) around the steady state ( $\bar{c}_{0}, \bar{a}_{0}, \bar{b}_{0}$ ) and obtain the following system:

$$
\left(\begin{array}{l}
\delta c(t)  \tag{20}\\
\delta a(t) \\
\delta b(t)
\end{array}\right)^{\prime}=J\left(0, \bar{c}_{0}, \bar{a}_{0}, \bar{b}_{0}\right)\left(\begin{array}{c}
\delta c(t) \\
\delta a(t) \\
\delta b(t)
\end{array}\right) .
$$

The eigenvalues of $J\left(0, \bar{c}_{0}, \bar{a}_{0}, \bar{b}_{0}\right)$ are

$$
\lambda_{c}=-\left(2 k_{1} k_{4}+k_{2}\right) \approx-k_{2}, \quad \lambda_{a}=-\tau_{a}, \quad \lambda_{b}=-\tau_{b}
$$

The corresponding eigenvectors are

$$
\xi_{c}=\left(\begin{array}{l}
1 \\
0 \\
0
\end{array}\right), \quad \xi_{a}=\left(\begin{array}{c}
\frac{k_{1}\left(1-\bar{c}_{0}\right)}{2 k_{1} k_{4}+k_{2}-\tau_{a}} \\
1 \\
0
\end{array}\right), \quad \xi_{b}=\left(\begin{array}{c}
\frac{k_{1}\left(1-\bar{c}_{0}\right)}{2 k_{1} k_{4}+k_{2}-\tau_{b}} \\
0 \\
1
\end{array}\right) .
$$

Solutions of system (20) can be written as linear combinations of

$$
\xi_{c} e^{-k_{2} t}, \quad \xi_{a} e^{-\tau_{a} t}, \quad \xi_{b} e^{-\tau_{b} t},
$$

that is,

$$
\begin{align*}
& \delta c(t)=l_{c} \xi_{c}^{1} e^{-k_{2} t}+l_{a} \xi_{a}^{1} e^{-\tau_{a} t}+l_{b} \xi_{b}^{1} e^{-\tau_{b} t}, \\
& \delta a(t)=l_{a} \xi_{a}^{2} e^{-\tau_{a} t},  \tag{21}\\
& \delta b(t)=l_{b} \xi_{b}^{3} e^{-\tau_{b} t},
\end{align*}
$$

for some constants $l_{c}, l_{a}$, and $l_{b}$ depending on initial conditions.
From (21), we see that the largest number among $1 / k_{2}, 1 / \tau_{a}$, and $1 / \tau_{b}$ indicates the time scale of $\delta c$. In order to observe a slow deactivation, one of $\tau_{a}$ and $\tau_{b}$ must be much smaller than $k_{2}$. Without loss of generality, we assume that $\tau_{b} \ll k_{2}$ and $\tau_{b} \leq \tau_{a}$ ( $\tau_{a}$ and $\tau_{b}$ are symmetric).

For the same reason as the single-positive-loop system, for any fixed initial condition in a small neighborhood of the stead state, the ratio of $\xi_{b}^{1}$ to $\xi_{b}^{3}$ determines how effectively $e^{-\tau_{b} t}$ slows down the dynamics of $\delta c$. Calculating the ratio, we obtain

$$
\frac{\xi_{b}^{1}}{\xi_{b}^{3}}=\frac{k_{1}\left(1-\bar{c}_{0}\right)}{2 k_{1} k_{4}+k_{2}-\tau_{b}} \stackrel{\tau_{b} \ll k_{2}}{\approx} \frac{k_{1}\left(k_{2}-k_{3}\right)}{\left(2 k_{1} k_{4}+k_{2}\right)^{2}}{ }^{k_{3} \ll k_{2}} \frac{k_{1} k_{2}}{\left(2 k_{1} k_{4}+k_{2}\right)^{2}} \stackrel{2 k_{1} k_{4} \ll k_{2}}{\approx} \frac{k_{1} k_{2}}{k_{2}^{2}}=K_{a} .
$$

Thus, the bigger $K_{a}$ is, the slower $\delta c$ converges to the zero steady state. To estimate $t_{1 \rightarrow 0}$, we consider the initial condition $\left(\gamma_{0}^{*}, \alpha_{0}^{*}, \beta_{0}^{*}\right)=\left(\bar{c}_{1}-\bar{c}_{0}, \bar{a}_{1}-\bar{a}_{0}, \bar{b}_{1}-\bar{b}_{0}\right)$, and take $t_{1 \rightarrow 0}$ as the time when $\delta c$ reaches $\gamma_{0}^{*} / e$. We use two extreme cases of $\tau_{a}$ to illustrate the effect of loop $A$ on the deactivation time scale $t_{1 \rightarrow 0}$.

Case $1 \tau_{a} \approx k_{2} \gg \tau_{b}$. After an initial period, $\delta c(t) \approx l_{b} \xi_{b}^{1} e^{-\tau_{b} t}$. If $l_{b} \xi_{b}^{1}$ is significant comparing to $\gamma_{0}^{*}$, then
the time it takes for $\delta c$ to drop from $\gamma_{0}^{*}$ to $\gamma_{0}^{*} / e$ is approximately

$$
\frac{1+\ln \left(\frac{l_{b} \xi_{b}^{1}}{\gamma_{0}^{*}}\right)}{\tau_{b}}=\frac{1+\ln \left(\frac{\xi_{b}^{1} \beta_{0}^{*}}{\xi_{b}^{3} \gamma_{0}^{*}}\right)}{\tau_{b}},
$$

where

$$
\frac{\xi_{b}^{1}}{\xi_{b}^{3}}=\frac{k_{1}\left(1-\bar{c}_{0}\right)}{2 k_{1} k_{4}+k_{2}-\tau_{b}} \stackrel{\tau_{b} \lll k_{2}}{\approx} \frac{k_{1}\left(k_{2}-k_{3}\right)}{\left(2 k_{1} k_{4}+k_{2}\right)^{2}} \stackrel{2 k_{1} k_{4} \ll k_{2}}{\approx} \frac{k_{1} k_{2}}{k_{2}^{2}}=K_{a},
$$

and

$$
\frac{\beta_{0}^{*}}{\gamma_{0}^{*}} \approx\left(\frac{2 K_{a}+1}{2 K_{a}}\right)\left(\frac{k_{c}}{k_{c}+1}\right) .
$$

Thus,

$$
\begin{equation*}
\frac{\xi_{b}^{1} \beta_{0}^{*}}{\xi_{b}^{3} \gamma_{0}^{*}} \approx\left(K_{a}+\frac{1}{2}\right) \frac{k_{c}}{1+k_{c}} . \tag{22}
\end{equation*}
$$

Comparing (22) with (13), we see that the $t_{1 \rightarrow 0}$ in the positive-positive-loop system is smaller. On the other hand, in both single-positive-loop and positive-positive-loop systems, $t_{1 \rightarrow 0}$ increases in $K_{a}$ and $k_{c}$.

Case $2 k_{2} \gg \tau_{a}=\tau_{b}$. The solutions in (21) can be rewritten as

$$
\begin{aligned}
& \delta c(t)=l_{c} \xi_{c}^{1} e^{-k_{2} t}+2 l_{b} \xi_{b}^{1} e^{-\tau_{b} t} \\
& \delta b(t)=l_{b} \xi_{b}^{3} e^{-\tau_{b} t} \\
& \delta a(t)=\delta b(t)
\end{aligned}
$$

After the initial layer, $\delta c(t) \approx 2 l_{b} \xi_{b}^{1} e^{-\tau_{b} t}$. If $2 l_{b} \xi_{b}^{1} / \gamma_{0}^{*}$ is large, then the time it takes for $\delta c$ to drop from $\gamma_{0}^{*}$ to $\gamma_{0}^{*} / e$ is

$$
\frac{1+\ln \left(\frac{2 \xi_{b}^{1} \beta_{0}^{*}}{\xi_{b}^{3} \tau_{0}^{*}}\right)}{\tau_{b}}
$$

where

$$
\begin{equation*}
\frac{2 \xi_{b}^{1} \beta_{0}^{*}}{\xi_{b}^{3} \gamma_{0}^{*}} \approx\left(2 K_{a}+1\right) \frac{k_{c}}{1+k_{c}} . \tag{23}
\end{equation*}
$$

Comparing (22) to (13), we see that the $t_{1 \rightarrow 0}$ in the positive-positive-loop system is bigger.

The analysis of the above two extreme cases suggests that with the additional loop $A$, the deactivation time scale $t_{1 \rightarrow 0}$ of the positive-positive-loop system can be either smaller or bigger than the single-positiveloop case depending on $\tau_{a}$ in a manner that bigger $\tau_{a}$ corresponding to smaller $t_{1 \rightarrow 0}$. In addition, $t_{1 \rightarrow 0}$ is increasing in $K_{a}$ and $k_{c}$.

### 2.2 The activation time scale

Linearizing system (16) around the steady state $\left(\bar{c}_{1}, \bar{a}_{1}, \bar{b}_{1}\right)$, we obtain the following system:

$$
\left(\begin{array}{l}
\delta c(t) \\
\delta a(t) \\
\delta b(t)
\end{array}\right)^{\prime}=J\left(1, \bar{c}_{1}, \bar{a}_{1}, \bar{b}_{1}\right)\left(\begin{array}{l}
\delta c(t) \\
\delta a(t) \\
\delta b(t)
\end{array}\right),
$$

where

$$
J\left(1, \bar{c}_{1}, \bar{a}_{1}, \bar{b}_{1}\right)=\left(\begin{array}{ccc}
-k_{1}\left(\bar{a}_{1}+\bar{b}_{1}\right)-k_{2} & k_{1}\left(1-\bar{c}_{1}\right) & k_{1}\left(1-\bar{c}_{1}\right) \\
k_{c}\left(1-\bar{a}_{1}\right) \tau_{a} & -\left(k_{c} \bar{c}_{1}+1\right) \tau_{a} & 0 \\
k_{c}\left(1-\bar{b}_{1}\right) \tau_{b} & 0 & -\left(k_{c} \bar{c}_{1}+1\right) \tau_{b}
\end{array}\right) .
$$

Let us assume that the term $k_{c}\left(1-\bar{b}_{1}\right) \tau_{b}$ is small enough so that $J\left(1, \bar{c}_{1}, \bar{a}_{1}, \bar{b}_{1}\right) \approx J^{*}$, where

$$
J^{*}=\left(\begin{array}{ccc}
-k_{1}\left(\bar{a}_{1}+\bar{b}_{1}\right)-k_{2} & k_{1}\left(1-\bar{c}_{1}\right) & k_{1}\left(1-\bar{c}_{1}\right) \\
k_{c}\left(1-\bar{a}_{1}\right) \tau_{a} & -\left(k_{c} \bar{c}_{1}+1\right) \tau_{a} & 0 \\
0 & 0 & -\left(k_{c} \bar{c}_{1}+1\right) \tau_{b}
\end{array}\right) .
$$

Denote the eigenvalues of $J^{*}$ as

$$
\lambda_{c}, \quad \lambda_{a}, \quad \lambda_{b}=-\left(k_{c} \bar{c}_{1}+1\right) \tau_{b},
$$

with their corresponding eigenvectors

$$
\xi_{c} \approx\left(\begin{array}{c}
\xi_{c}^{1} \\
\xi_{c}^{2} \\
0
\end{array}\right), \quad \xi_{a} \approx\left(\begin{array}{c}
\xi_{a}^{1} \\
\xi_{a}^{2} \\
0
\end{array}\right), \quad \xi_{b} \approx\left(\begin{array}{c}
\frac{k_{1}\left(1-\bar{c}_{1}\right)}{k_{1}\left(\bar{a}_{1}+\bar{b}_{1}\right)+k_{2}-\left(k_{c} \bar{c}_{1}+1\right) \tau_{b}} \\
0 \\
1
\end{array}\right)
$$

To estimate $t_{0 \rightarrow 1}$, we consider the initial condition $\left(\gamma_{0}^{*}, \alpha_{0}^{*}, \beta_{0}^{*}\right)=\left(\bar{c}_{0}-\bar{c}_{1}, \bar{a}_{0}-\bar{a}_{1}, \bar{b}_{0}-\bar{b}_{1}\right)$. Again, we use the previous two extreme cases of $\tau_{a}$ to elucidate the effect of the loop $A$ on $t_{0 \rightarrow 1}$.

Case $1 \tau \approx k_{2} \gg \tau_{b}$. After an initial period, $\delta c$ can be approximated by $l_{b} \xi_{b}^{1} e^{-\tau_{b} t}$. If $l_{b} \xi_{b}^{1} / \gamma_{0}^{*}$ is large, then $t_{1 \rightarrow 0}$ can be approximated by

$$
\frac{1+\ln \left(\frac{\xi_{b}^{1} \beta_{0}^{*}}{\xi_{b}^{3} \gamma_{0}^{*}}\right)}{\tau_{b}},
$$

where

$$
\frac{\xi_{b}^{1}}{\xi_{b}^{3}} \approx \frac{k_{1} k_{2}}{\left(2 k_{1}+k_{2}\right)^{2}}\left(1+\frac{1}{k_{c}}\right)^{2}=\frac{K_{a}}{\left(2 K_{a}+1\right)^{2}}\left(1+\frac{1}{k_{c}}\right)^{2}, \frac{\beta_{0}^{*}}{\gamma_{0}^{*}} \approx\left(\frac{2 K_{a}+1}{2 K_{a}}\right)\left(\frac{k_{c}}{k_{c}+1}\right) .
$$

Thus,

$$
\begin{equation*}
\frac{\xi_{b}^{1} \beta_{0}^{*}}{\xi_{b}^{3} \gamma_{0}^{*}} \approx \frac{1}{2\left(2 K_{a}+1\right)^{2}}\left(1+\frac{1}{k_{c}}\right) . \tag{24}
\end{equation*}
$$

The ratio in (24) is much smaller than that in (15). On the other hand, in both single-positive-loop and positive-positive-loop systems, $t_{0 \rightarrow 1}$ is decreasing in $K_{a}$ and $k_{c}$.

Case $2 k_{2} \gg \tau_{a}=\tau_{b}$. After the initial layer, $\delta c \approx 2 l_{b} \xi_{b}^{1} e^{-\tau_{b} t}$. If $2 l_{b} \xi_{b}^{1} / \gamma_{0}^{*}$ is significant, then the time it takes for $\delta c$ to reach $\gamma_{0}^{*} / e$ is

$$
\frac{1+\ln \left(\frac{2 \xi_{\xi^{1}}^{*} b_{0}^{*}}{\xi_{b}^{3} \gamma_{0}^{*}}\right)}{\tau_{b}},
$$

where

$$
\begin{equation*}
\frac{2 \xi_{b}^{1} \beta_{0}^{*}}{\xi_{b}^{3} \gamma_{0}^{*}} \approx \frac{1}{\left(2 K_{a}+1\right)^{2}}\left(1+\frac{1}{k_{c}}\right) . \tag{25}
\end{equation*}
$$

Comparing (25) to (15), we see that the $t_{1 \rightarrow 0}$ in the positive-positive-loop system is smaller.

The analysis of the above two extreme cases suggests that with the additional loop $A$, the activation time scale $t_{1 \rightarrow 0}$ is smaller than that in the single-positive-loop case. Moreover, $t_{0 \rightarrow 1}$ is decreasing in $K_{a}$ and $k_{c}$.

## 3 Positive-Negative-Loop Module

Recall the equations for the positive-negative-loop module:

$$
\begin{align*}
& \frac{d c}{d t}=k_{1 b} b(1-c)-\left(k_{2}+k_{1 a} a\right) c+k_{3}, \\
& \frac{d a}{d t}=\left(k_{c a} s c(1-a)-a+k_{4}\right) \tau_{a},  \tag{26}\\
& \frac{d b}{d t}=\left(k_{c b} s c(1-b)-b+k_{4}\right) \tau_{b} .
\end{align*}
$$

For the purpose of analysis, we assume that $k_{c a}=k_{c b}:=k_{c}$, which significantly simplifies our computations.
The steady state at $s=0$ is

$$
\bar{c}_{0}=\frac{k_{3}+k_{1 b} k_{4}}{k_{2}+k_{4}\left(k_{1 a}+k_{1 b}\right)}, \quad \bar{a}_{0}=\bar{b}_{0}=k_{4} .
$$

When $s=1$, under the assumptions $k_{4} \ll 1, k_{3} \ll k_{2}$, and $k_{1 b} k_{c}>k_{2}$ (similar to those in the positive-positive-loop systems, Section 2), we have

$$
\bar{c}_{1} \approx \frac{k_{1 b} k_{c}-k_{2}}{k_{2} k_{c}+\left(k_{1 b}+k_{1 a}\right) k_{c}}, \quad \bar{a}_{1}=\bar{b}_{1} \approx \frac{k_{1 b} k_{c}-k_{2}}{k_{1 b} k_{c}+k_{1 b}+k_{1 a}} .
$$

The Jacobian matrix of (26) is

$$
J(s, c, a, b)=\left(\begin{array}{ccc}
-k_{1 b} b-\left(k_{2}+k_{1 a} a\right) & -k_{1 a} c & k_{1 b}(1-c) \\
k_{c a} s(1-a) \tau_{a} & -\left(k_{c a} s c+1\right) \tau_{a} & 0 \\
k_{c b} s(1-b) \tau_{b} & 0 & -\left(k_{c b} s c+1\right) \tau_{b}
\end{array}\right) .
$$

### 3.1 The deactivation time scale

Linearizing (26) around the steady state ( $\bar{c}_{0}, \bar{a}_{0}, \bar{b}_{0}$ ), we obtain the following system:

$$
\left(\begin{array}{l}
\delta c(t)  \tag{27}\\
\delta a(t) \\
\delta b(t)
\end{array}\right)^{\prime}=J\left(0, \bar{c}_{0}, \bar{a}_{0}, \bar{b}_{0}\right)\left(\begin{array}{c}
\delta c(t) \\
\delta a(t) \\
\delta b(t)
\end{array}\right) .
$$

The eigenvalues of $J\left(0, \bar{c}_{0}, \bar{a}_{0}, \bar{b}_{0}\right)$ are

$$
\lambda_{c}=-\left(\left(k_{1 b}+k_{1 a}\right) k_{4}+k_{2}\right) \approx-k_{2}, \quad \lambda_{a}=-\tau_{a}, \quad \lambda_{b}=-\tau_{b} .
$$

The corresponding eigenvectors are

$$
\xi_{c}=\left(\begin{array}{l}
1 \\
0 \\
0
\end{array}\right), \quad \xi_{a}=\left(\begin{array}{c}
-\frac{k_{1 a} \bar{c}_{0}}{\left(k_{1 b}+k_{1 a} k_{4}+k_{2}-\tau_{a}\right.} \\
1 \\
0
\end{array}\right), \quad \xi_{b}=\left(\begin{array}{c}
\frac{k_{1 b}\left(1-\bar{c}_{0}\right)}{\left(k_{1 b}+k_{1 a}\right) k_{4}+k_{2}-\tau_{b}} \\
0 \\
1
\end{array}\right) .
$$

Solutions of system (27) can be written as the linear combination of

$$
\xi_{c} e^{-k_{2} t}, \quad \xi_{a} e^{-\tau_{a} t}, \quad \xi_{b} e^{-\tau_{b} t},
$$

that is,

$$
\begin{align*}
& \delta c(t)=l_{c} \xi_{c}^{1} e^{-k_{2} t}+l_{a} \xi_{a}^{1} e^{-\tau_{a} t}+l_{b} \xi_{b}^{1} e^{-\tau_{b} t}, \\
& \delta a(t)=l_{a} \xi_{a}^{2} e^{-\tau_{a} t},  \tag{28}\\
& \delta b(t)=l_{b} \xi_{b}^{3} e^{-\tau_{b} t},
\end{align*}
$$

for some constants $l_{c}, l_{a}$, and $l_{b}$ depending on initial conditions. As done for the positive-positive-loop system, we consider the following two cases.

Case $1 \tau_{a} \approx k_{2} \gg \tau_{b}$, fast negative loop and slow positive loop. The contribution of the slow dynamics is through the term $l_{b} \xi_{b}^{1} e^{-\tau_{b} t}$. As before, we compute

$$
\frac{\xi_{b}^{1}}{\xi_{b}^{3}}=\frac{k_{1 b}\left(1-\bar{c}_{0}\right)}{\left(k_{1 b}+k_{1 a}\right) k_{4}+k_{2}-\tau_{b}} \stackrel{\tau_{b} \ll k_{2}}{\approx} \frac{k_{1 b}\left(k_{2}+k_{1 a} k_{4}-k_{3}\right)}{\left(\left(k_{1 b}+k_{1 a}\right) k_{4}+k_{2}\right)^{2}} \stackrel{k_{3}, k_{4} \ll k_{2}}{\approx} \frac{k_{1 b}}{k_{2}}:=K_{a}
$$

and

$$
\frac{\beta_{0}^{*}}{\gamma_{0}^{*}} \approx \frac{\left(1+K_{a}+K_{d}\right) k_{c}}{K_{a}+K_{a} k_{c}+K_{d}} .
$$

Thus,

$$
\begin{equation*}
\frac{\xi_{b}^{1} \beta_{0}^{*}}{\xi_{b}^{3} \gamma_{0}^{*}} \approx \frac{\left(1+K_{a}+K_{d}\right) k_{c}}{k_{c}+1+K_{d} / K_{a}} . \tag{29}
\end{equation*}
$$

A straightforward calculation shows that (29) is always greater than (13) under the assumption $k_{1 b} k_{c}>k_{2}$ (the same assumption that guarantees the positivity of $\bar{c}_{1}$ ). Equation (29) also shows that $t_{1 \rightarrow 0}$ increases in $K_{a}$ and $k_{c}$.

Case $2 k_{2} \gg \tau_{a}=\tau_{b}$, slow negative loop and slow positive loop. In this case, both terms $l_{a} \xi_{a}^{1} e^{-\tau_{a} t}$ and $l_{b} \xi_{b}^{1} e^{-\tau_{b} t}$ are slow. The overall contribution from the slow terms is indicated by $\frac{\left(\xi_{b}^{1}+\xi_{a}^{1}\right) \beta_{0}^{*}}{\xi_{b}^{2} \gamma_{0}^{*}}$. Thus,

$$
\begin{equation*}
\frac{\left(\xi_{b}^{1}+\xi_{a}^{1}\right) \beta_{0}^{*}}{\xi_{b}^{2} \gamma_{0}^{*}} \approx \frac{\left(1+K_{a}+K_{d}\right) k_{c}}{k_{c}+1+K_{d} / K_{a}}-K_{d}\left(\frac{k_{3}}{k_{2}}+K_{a} k_{4}\right) \frac{\left(1+K_{a}+K_{d}\right) k_{c}}{K_{a}+K_{a} k_{c}+K_{d}} . \tag{30}
\end{equation*}
$$

Notice that $k_{4} \ll 1$ and $k_{3} \ll k_{2}$, so the second term in (30) is very small. Thus, (30) is still bigger than (13), but smaller than (29). That is, the deactivation in this case is slower than that in the single-positive-loop case, but faster than the fast-negative-slow-positive case.

### 3.2 The activation time scale

Linearizing system (26) around the steady state $\left(\bar{c}_{1}, \bar{a}_{1}, \bar{b}_{1}\right)$, we obtain the following system:

$$
\left(\begin{array}{c}
\delta c(t) \\
\delta a(t) \\
\delta b(t)
\end{array}\right)^{\prime}=J\left(1, \bar{c}_{1}, \bar{a}_{1}, \bar{b}_{1}\right)\left(\begin{array}{c}
\delta c(t) \\
\delta a(t) \\
\delta b(t)
\end{array}\right),
$$

where

$$
J(1, c, a, b)=\left(\begin{array}{ccc}
-k_{1 b} \bar{b}_{1}-\left(k_{2}+k_{1 a} \bar{a}_{1}\right) & -k_{1 a} \bar{c}_{1} & k_{1 b}\left(1-\bar{c}_{1}\right) \\
k_{c a}\left(1-\bar{a}_{1}\right) \tau_{a} & -\left(k_{c a} \bar{c}_{1}+1\right) \tau_{a} & 0 \\
k_{c b}\left(1-\bar{b}_{1}\right) \tau_{b} & 0 & -\left(k_{c b} \bar{c}_{1}+1\right) \tau_{b}
\end{array}\right) .
$$

Let us assume that the term $k_{c b}\left(1-\bar{b}_{1}\right) \tau_{b}$ is small enough so that $J\left(1, \bar{c}_{1}, \bar{a}_{1}, \bar{b}_{1}\right) \approx J^{*}$, where

$$
J^{*}=\left(\begin{array}{ccc}
-k_{1 b} \bar{b}_{1}-\left(k_{2}+k_{1 a} \bar{a}_{1}\right) & -k_{1 a} \bar{c}_{1} & k_{1 b}\left(1-\bar{c}_{1}\right) \\
k_{c a}\left(1-\bar{a}_{1}\right) \tau_{a} & -\left(k_{c a} \bar{c}_{1}+1\right) \tau_{a} & 0 \\
0 & 0 & -\left(k_{c b} \bar{c}_{1}+1\right) \tau_{b}
\end{array}\right) .
$$

Denote the eigenvalues of $J^{*}$ as

$$
\lambda_{c}, \quad \lambda_{a}, \quad \lambda_{b}=-\left(k_{c b} \bar{c}_{1}+1\right) \tau_{b},
$$

with their corresponding eigenvectors

$$
\xi_{c} \approx\left(\begin{array}{c}
\xi_{c}^{1} \\
\xi_{c}^{2} \\
0
\end{array}\right), \quad \xi_{a} \approx\left(\begin{array}{c}
\xi_{a}^{1} \\
\xi_{a}^{2} \\
0
\end{array}\right), \quad \xi_{b} \approx\left(\begin{array}{c}
\frac{k_{1 b}\left(1-\bar{c}_{1}\right)}{k_{1 b} \bar{b}_{1}+\left(k_{2}+k_{1 a} \bar{a}_{1}\right)-\left(k_{c b} \bar{c}_{1}+1\right) \tau_{b}} \\
0 \\
1
\end{array}\right)
$$

To estimate $t_{0 \rightarrow 1}$, we consider the initial condition $\left(\gamma_{0}^{*}, \alpha_{0}^{*}, \beta_{0}^{*}\right)=\left(\bar{c}_{0}-\bar{c}_{1}, \bar{a}_{0}-\bar{a}_{1}, \bar{b}_{0}-\bar{b}_{1}\right)$. Again, we use the previous two extreme cases of $\tau_{a}$ to elucidate the effect of the loop $A$ on $t_{0 \rightarrow 1}$.

Case $1 \tau_{a} \approx k_{2} \gg \tau_{b}$, fast negative loop and slow positive loop. The slow term is $l_{b} \xi_{b}^{1} e^{-\tau_{b} t}$. After an initial period, $\delta c$ can be approximated by $l_{b} \xi_{b}^{1} e^{-\tau_{b} t}$. Calculating $\xi_{b}^{1} / \xi_{b}^{3}$, we have

$$
\frac{\xi_{b}^{1}}{\xi_{b}^{3}} \approx \frac{\left(k_{c}+k_{c} K_{d}+1\right)\left(K_{a}+K_{a} k_{c}+K_{d}\right)}{k_{c}^{2}\left(K_{a}+K_{d}+1\right)^{2}},
$$

which is consistent with of the single-positive-loop case (corresponding to $K_{d}=0$ ). On the other hand,

$$
\frac{\beta_{0}^{*}}{\gamma_{0}^{*}} \approx \frac{\left(1+K_{a}+K_{d}\right) k_{c}}{K_{a}+K_{a} k_{c}+K_{d}} .
$$

Thus

$$
\begin{equation*}
\frac{\xi_{b}^{1} \beta_{0}^{*}}{\xi_{b}^{3} \gamma_{0}^{*}} \approx \frac{k_{c}+k_{c} K_{d}+1}{k_{c}\left(K_{a}+K_{d}+1\right)} . \tag{31}
\end{equation*}
$$

It is easy to see that (31) is greater than (15) under the condition $k_{1 b} k_{c}>k_{2}$. So, in this case we expect slower activation in positive-negative-loop systems than in single-positive-loop systems.

Case $2 k_{2} \gg \tau_{a}=\tau_{b}$, slow negative loop and slow positive loop. Now both terms $l_{a} \xi_{a}^{1} e^{-\tau_{a} t}$ and $l_{b} \xi_{b}^{1} e^{-\tau_{b} t}$ are slow, and the contribution from the slow terms is indicated by

$$
\frac{\left(\xi_{b}^{1}+\xi_{a}^{1}\right) \beta_{0}^{*}}{\xi_{b}^{2} \gamma_{0}^{*}} \approx \frac{k_{c}+k_{c} K_{d}+1}{k_{c}\left(K_{a}+K_{d}+1\right)}-\frac{K_{d}\left(K_{a} k_{c}-1\right)}{K_{a} k_{c}\left(K_{a}+K_{d}+1\right)},
$$

which is less than (15). That is, the activation in this case is faster than that in the single-positive-loop system. Simulation shows that the are very close (Figures 5C-5D).

The analysis of the above two extreme cases suggests that with the additional negative loop $A$, the activation time scale in the positive-negative-loop system could be either slower or faster than that in the single-positive-loop system.

## 4 Two-Time-Scale Asymptotic Expansion

In this section, we derive an asymptotic expansion of the overall solution, including those far away from the steady state. We will show that the single-positive-loop system can function as a low-pass filter, which explains why the relative size of the noise's time scale and the system's intrinsic time scales is important to noise attenuation. We also show that our local solution closely resembles the global solution.

Different from the traditional expansion approach, the two-time-scale asymptotic expansion method [7], originated from the strained-corrdinate method [7], is mainly used to avoid the accumulative effect and to obtain a solution uniformly valid over a long time interval. In our case, it also provides an explicit relation between the solutions and the two separated time scales, suggesting that the single-positive-loop system can function as a low-pass filter.

### 4.1 Low-pass filter

This section is devoted to obtain the asymptotic expansion of solutions of system (3) and to show that the single-positive-loop module can function as a low-pass filter. Let $\varepsilon=\tau_{b}$. If $\varepsilon \ll k_{2}$, we could expand the solutions of system (3) in powers of $\varepsilon$ :

$$
\begin{aligned}
& c=c_{0}\left(t^{+}, \tilde{t}\right)+\varepsilon c_{1}\left(t^{+}, \tilde{t}\right)+\varepsilon^{2} c_{2}\left(t^{+}, \tilde{t}\right)+\cdots, \\
& b=b_{0}\left(t^{+}, \tilde{t}\right)+\varepsilon b_{1}\left(t^{+}, \tilde{t}\right)+\varepsilon^{2} b_{2}\left(t^{+}, \tilde{t}\right)+\cdots,
\end{aligned}
$$

involving the fast time scale $t^{+}=t$ and the slow time scale $\tilde{t}=\varepsilon t$. The initial conditions are

$$
c_{0}(0,0)=\gamma_{0}, \quad b_{0}(0,0)=\beta_{0}, \quad c_{i}(0,0)=b_{i}(0,0)=0, \quad i=1,2, \ldots
$$

Accordingly, system (3) becomes

$$
\begin{align*}
& \frac{\partial c_{0}}{\partial t^{+}}+\varepsilon\left(\frac{\partial c_{0}}{\partial \tilde{t}}+\frac{\partial c_{1}}{\partial t^{+}}\right)+O\left(\varepsilon^{2}\right)  \tag{32}\\
& =k_{1}\left(b_{0}\left(t^{+}, \tilde{t}\right)+\varepsilon b_{1}\left(t^{+}, \tilde{t}\right)+\cdots\right)\left(1-c_{0}\left(t^{+}, \tilde{t}\right)-\varepsilon c_{1}\left(t^{+}, \tilde{t}\right)-\cdots\right)-k_{2}\left(c_{0}\left(t^{+}, \tilde{t}\right)+\varepsilon c_{1}\left(t^{+}, \tilde{t}\right)+\cdots\right)+k_{3}, \\
& \frac{\partial b_{0}}{\partial t^{+}}+\varepsilon\left(\frac{\partial b_{0}}{\partial \tilde{t}}+\frac{\partial b_{1}}{\partial t^{+}}\right)+O\left(\varepsilon^{2}\right)  \tag{33}\\
& =\varepsilon\left(k_{c} s\left(t^{+}, \tilde{t}\right)\left(c_{0}\left(t^{+}, \tilde{t}\right)+\varepsilon c_{1}\left(t^{+}, \tilde{t}\right)+\cdots\right)\left(1-b_{0}\left(t^{+}, \tilde{t}\right)-\varepsilon b_{1}\left(t^{+}, \tilde{t}\right)-\cdots\right)\right. \\
& \left.-\left(b_{0}\left(t^{+}, \tilde{t}\right)+\varepsilon b_{1}\left(t^{+}, \tilde{t}\right)+\cdots\right)+k_{4}\right) .
\end{align*}
$$

Equating the same order on both sides of each equation, we have

$$
\begin{align*}
\frac{\partial b_{0}}{\partial t^{+}} & =0, \text { i.e. } b_{0}\left(t^{+}, \tilde{t}\right)=b_{0}(\tilde{t})  \tag{34}\\
\frac{d b_{0}}{d \tilde{t}}+\frac{\partial b_{1}}{\partial t^{+}} & =k_{c} s c_{0}\left(1-b_{0}\right)-b_{0}+k_{4} \tag{35}
\end{align*}
$$

and

$$
\begin{align*}
\frac{\partial c_{0}}{\partial t^{+}} & =k_{1} b_{0}(\tilde{t})\left(1-c_{0}\left(t^{+}, \tilde{t}\right)\right)-k_{2} c_{0}\left(t^{+}, \tilde{t}\right)+k_{3}  \tag{36}\\
\frac{\partial c_{0}}{\partial \tilde{t}}+\frac{\partial c_{1}}{\partial t^{+}} & =k_{1}\left(b_{1}\left(1-c_{0}\right)-b_{0} c_{1}\right)-k_{2} c_{1} \tag{37}
\end{align*}
$$

In principal, there are different ways of choosing $b_{0}(\tilde{t})$. However, in order to obtain a uniform zero-order approximation, we need to choose $b_{0}$ such that the solution $b_{1}$ from (35) has no linear term in $t^{+}[7]$. Once $b_{0}$ is determined, solving $c_{0}$ from (36), we obtain

$$
\begin{align*}
c_{0}\left(t^{+}, \tilde{t}\right) & =e^{-\left(k_{1} b_{0}(\tilde{t})+k_{2}\right) t^{+}}\left(h_{0}(\tilde{t})+\frac{k_{1} b_{0}(\tilde{t})+k_{3}}{k_{1} b_{0}(\tilde{t})+k_{2}}\left(e^{\left(k_{1} b_{0}(\tilde{t})+k_{2}\right) t^{+}}-1\right)\right)  \tag{38}\\
& =\left(\gamma_{0}-\frac{k_{1} b_{0}(\tilde{t})+k_{3}}{k_{1} b_{0}(\tilde{t})+k_{2}}\right) e^{-\left(k_{1} b_{0}(\tilde{t})+k_{2}\right) t^{+}}+\frac{k_{1} b_{0}(\tilde{t})+k_{3}}{k_{1} b_{0}(\tilde{t})+k_{2}}
\end{align*}
$$

Thus, the question becomes how to choose an appropriate $b_{0}$ for a given noise input. Consider the following scenarios of the input signal:

1. Noise is on a fast time scale.

In this case, the signal can be written as the sum of a constant input and fast noise term $s_{1}\left(t^{+}\right)$with mean zero, i.e. $s=s_{0}+s_{1}\left(t^{+}\right)$. In order to eliminate the linear term of $t^{+}$in the solution of $b_{1}$, we choose $b_{0}$ satisfying:

$$
\begin{equation*}
\frac{d b_{0}}{d \tilde{t}}=k_{c} s_{0}\left(1-b_{0}\right) \frac{k_{1} b_{0}+k_{3}}{k_{1} b_{0}+k_{2}}-b_{0}+k_{4} . \tag{39}
\end{equation*}
$$

As a result, (35) becomes

$$
\frac{\partial b_{1}}{\partial t^{+}}=k_{c}\left(\left(s_{1}\left(t^{+}\right)+s_{0}\right) c_{0}-s_{0} \frac{k_{1} b_{0}+k_{3}}{k_{1} b_{0}+k_{2}}\right)\left(1-b_{0}\right),
$$

whose the solution $b_{1}$ does not contain any linear term of $t^{+}$. Notice that (39) is separable and can be solved analytically. Since none of $b_{1}, c_{0}$, and $\partial c_{0} / \partial \tilde{t}$ contains linear term in $t^{+}$(see (38)), the solution $c_{1}$ of (37),

$$
c_{1}\left(t^{+}, \tilde{t}\right)=e^{-\left(k_{1} b_{0}(\tilde{t})+k_{2}\right) t^{+}} \int_{0}^{t^{+}}\left(k_{1} b_{1}\left(1-c_{0}\right)-\frac{\partial c_{0}}{\partial \tilde{t}}\right) e^{\left(k_{1} b_{0}(\tilde{t})+k_{2}\right) t^{+}} d t^{+}
$$

has no linear term in $t^{+}$. Therefore,

$$
\begin{align*}
& c\left(t^{+}, \tilde{t}\right)=\left(\gamma_{0}-\frac{k_{1} b_{0}(\tilde{t})+k_{3}}{k_{1} b_{0}(\tilde{t})+k_{2}}\right) e^{-\left(k_{1} b_{0}(\tilde{t})+k_{2}\right) t^{+}}+\frac{k_{1} b_{0}(\tilde{t})+k_{3}}{k_{1} b_{0}(\tilde{t})+k_{2}}+O(\varepsilon)  \tag{40}\\
& b\left(t^{+}, \tilde{t}\right)=b_{0}(\tilde{t})+O(\varepsilon)
\end{align*}
$$

is a uniform zero-order approximation of (3) [7] (Figures S3A-S3B). Observe that the noise term $s_{1}\left(t^{+}\right)$does not show up in equation (39), so the zero-order approximations with and without noise are the same, suggesting that fast varying noises are filtered out through the system (Figures S3DS3E).
2. Noise is on the slow time scale $\tilde{t}$, i.e. $s=s(\tilde{t})$.

In this case, to eliminate the linear term of $t^{+}$in $b_{1}$, we take

$$
\begin{equation*}
\frac{d b_{0}}{d \tilde{t}}=k_{c} s(\tilde{t})\left(1-b_{0}\right) \frac{k_{1} b_{0}+k_{3}}{k_{1} b_{0}+k_{2}}-b_{0}+k_{4} . \tag{41}
\end{equation*}
$$

Similar to case 1, we have (40) as the uniform zero-order approximation of (3), but the $b_{0}(\tilde{t})$ in this case is the solution to (41) instead of (39). Notice that noises persist in the zero-order approximation through the term $s(\tilde{t})$ in equation (41), and thus the slow noise could significantly affect the output (Figure S3F). Nevertheless, the leading order in (40) again matches the whole solution well (Figure S3C).
3. Noise is decoupled into a sum of fast and slow noise terms, i.e., $s=s_{0}+s_{1}\left(t^{+}\right)+s_{2}(\tilde{t})$.

In this case, we take $b_{0}$ as the solution of

$$
\frac{d b_{0}}{d \tilde{t}}=k_{c} s_{2}(\tilde{t})\left(1-b_{0}\right) \frac{k_{1} b_{0}+k_{3}}{k_{1} b_{0}+k_{2}}-b_{0}+k_{4} .
$$

Thus, only the slow term enters in the equation of $b_{0}$, and the fast noise is filtered out through the system (Figures S3G-S3J).

In summary, the zero-order solutions match the whole solution very well (Figures S3A-S3C). System (3) is not significantly affected by noise on a fast time scale (Figures S3D-S3E). The single-positive-loop system can function as a low-pass filter for different nature of noises (Figure S3D-S3J).

### 4.2 Connection to the linearization approach

The zero-order approximation obtained in (40) is a global solution, while the linear stability analysis in Sections 1-3 focuses on local solutions around the steady states. In this section, we use the deactivation process of the single-positive-loop system to elucidate the connections between these two approaches. Comparison to other systems could be done in a similar fashioin, but much more complicated computationally.

During deactivation, $s \equiv 0$ and the solution of (39) is

$$
b_{0}(\tilde{t})=\left(\beta_{0}-k_{4}\right) e^{-\tilde{t}}+k_{4} .
$$

Writing in the new coordinate $b_{0}^{*}=b_{0}-k_{4}$, we have $b_{0}^{*}(\tilde{t})=\beta_{0}^{*} e^{-\tilde{t}}$, where $\beta_{0}^{*}=\beta_{0}-k_{4}$ is the initial condition. For the solution $c_{0}$ in (38), we linearize it around

$$
\bar{c}_{0}:=\frac{k_{1} k_{4}+k_{3}}{k_{1} k_{4}+k_{2}},
$$

obtaining

$$
\begin{align*}
c_{0}^{*}\left(t^{+}, \tilde{t}\right) & =\left(\gamma_{0}^{*}-\beta_{0}^{*} \frac{\left(k_{2}-k_{3}\right) k_{1}}{\left(k_{1} k_{4}+k_{2}\right)\left(k_{1}\left(\beta_{0}^{*} e^{-\tau_{b} t}+k_{4}\right)+k_{2}\right)} e^{-\tilde{t}}\right) e^{-\left(k_{1}\left(\beta_{0}^{*} e^{-\tilde{t}}+k_{4}\right)+k_{2}\right) t^{+}}  \tag{42}\\
& +\beta_{0}^{*} \frac{\left(k_{2}-k_{3}\right) k_{1}}{\left(k_{1} k_{4}+k_{2}\right)\left(k_{1}\left(\beta_{0}^{*} e^{-\tau_{b} t}+k_{4}\right)+k_{2}\right)} e^{-\tilde{t}},
\end{align*}
$$

where $\gamma_{0}^{*}=\gamma_{0}-\bar{c}_{0}$ is the initial condition of $c_{0}^{*}$. The dynamics of $c_{0}^{*}$ is determined by two exponential functions. One is on the fast time scale $t^{+}$, and the other one is on the slow time scale $\tilde{t}$. Notice that (42) is in the same form of the equation (12), obtained using linearization and eigenvalue analysis $\left(\tilde{t}=\tau_{b} t\right.$, $\left.t^{+}=t\right)$.

## 5 Fluctuation-Dissipation Theorem Approach

Let us first review the Fluctuation-Dissipation Theorem used in $[4,5,9,10]$. Consider a general inputoutput system

$$
\begin{equation*}
\frac{d x_{i}}{d t}=J_{i}^{+}(x)-J_{i}^{-}(x), \quad i=1, \ldots, n . \tag{43}
\end{equation*}
$$

Here, $J_{i}^{+}$and $J_{i}^{-}$are the fluxes of the production and degradation of $x_{i}$, respectively. We assume that $x_{0}$ is the only source of noise, and $x_{0}(t)=\left\langle x_{0}\right\rangle+\sigma_{0}(t)$, where $\sigma_{0}$ has mean zero and autocorrelation time $\tau_{0}$. The two quantities that we are interested are the sensitivity of the steady state and the noise amplification
rate. We define the sensitivity of the steady state for each $i=1, \ldots, n$ by susceptibility $[5,9,10,12]$

$$
s_{i}=\frac{\left\langle x_{0}\right\rangle}{\left\langle x_{i}\right\rangle} \frac{d\left\langle x_{i}\right\rangle}{d\left\langle x_{0}\right\rangle},
$$

and the noise amplification rate by [5]

$$
r_{i}=\frac{\eta_{i}}{\eta_{0}}=\frac{\operatorname{std}\left(x_{i}\right) /\left\langle x_{i}\right\rangle}{\operatorname{std}\left(x_{0}\right) /\left\langle x_{0}\right\rangle} .
$$

Both quantities are closely related to the reaction flux elasticities [5, 9, 10], $H_{i j}$, which measures how the ratio of production to degradation changes with respect to the concentrations,

$$
\begin{equation*}
H_{i j}=-\frac{\left\langle x_{j}\right\rangle}{\left\langle J_{i}\right\rangle}\left(\frac{\partial\left\langle J_{i}^{+}\right\rangle}{\partial\left\langle x_{j}\right\rangle}-\frac{\partial\left\langle J_{i}^{-}\right\rangle}{\partial\left\langle x_{j}\right\rangle}\right)=\frac{\partial \ln \left\langle J_{i}^{-}\right\rangle}{\partial \ln \left\langle x_{j}\right\rangle}-\frac{\partial \ln \left\langle J_{i}^{+}\right\rangle}{\partial \ln \left\langle x_{j}\right\rangle}=-\frac{\partial \ln \left(\left\langle J_{i}^{+}\right\rangle /\left\langle J_{i}^{-}\right\rangle\right)}{\partial \ln \left\langle x_{j}\right\rangle} . \tag{44}
\end{equation*}
$$

The susceptibility $s_{i}$ can be solved from the equation [5, 9, 10]

$$
\begin{equation*}
H_{i 0}+H_{i 1} s_{1}+\cdots+H_{i n} s_{n}=0, \quad i=1, \ldots n . \tag{45}
\end{equation*}
$$

The $\eta_{i}$ 's are solutions to the following equation $[5,9,10]$ :

$$
\begin{equation*}
M \eta+\eta M^{t}+D=0 \tag{46}
\end{equation*}
$$

where

$$
\begin{gathered}
M_{i j}=-\frac{\left\langle J_{i}\right\rangle}{\left\langle x_{i}\right\rangle}\left(\frac{\partial \ln \left\langle J_{i}^{-}\right\rangle}{\partial \ln \left\langle x_{j}\right\rangle}-\frac{\partial \ln \left\langle J_{i}^{+}\right\rangle}{\partial \ln \left\langle x_{j}\right\rangle}\right)=-\frac{\left\langle J_{i}\right\rangle}{\left\langle x_{i}\right\rangle} H_{i j}, \quad \eta_{i i}=\eta_{i}^{2}, \\
D_{00}=\frac{2 \eta_{0}^{2}}{\tau_{0}}, \quad D_{i j}=0 \text { for }(i, j) \neq(0,0), i, j=0, \ldots, n .
\end{gathered}
$$

Applying to our single-positive-loop system (3), we have

$$
n=2, \quad x_{0}=s, \quad x_{1}=c, \quad x_{2}=b,
$$

$$
J_{1}^{+}=k_{1} b(1-c)+k_{3}, \quad J_{1}^{-}=k_{2} c, \quad J_{2}^{+}=\left(k_{c} s c(1-b)+k_{4}\right) \tau_{b}, \quad J_{2}^{-}=b \tau_{b} .
$$

By the definition of $H_{i j}$ in (44), we obtain

$$
H=\left(H_{i j}\right)=\left(\begin{array}{ccc}
-1 & B & -1 \\
0 & -1 & A / k_{2}
\end{array}\right), \quad i=1,2, \quad j=0,1,2,
$$

where $A=\langle\bar{b}\rangle k_{1}+k_{2}, B=1+k_{c}\langle\bar{c}\rangle\langle s\rangle$. Thus,

$$
\begin{gathered}
M=\left(M_{i j}\right)=\left(\begin{array}{ccc}
-\frac{1}{\tau_{0}} & 0 & 0 \\
\tau_{b} & -\tau_{b} B & \tau_{b} \\
0 & k_{2} & -A
\end{array}\right), \quad \eta=\left(\eta_{i j}\right)=\left(\begin{array}{ccc}
\eta_{0}^{2} & \eta_{01} & \eta_{02} \\
\eta_{10} & \eta_{1}^{2} & \eta_{12} \\
\eta_{20} & \eta_{21} & \eta_{22}
\end{array}\right) \\
D=\left(D_{i j}\right)=\left(\begin{array}{ccc}
-\frac{2 \eta_{0}^{2}}{\tau_{0}} & 0 & 0 \\
0 & 0 & 0 \\
0 & 0 & 0
\end{array}\right), \quad i=0,1,2, \quad j=0,1,2
\end{gathered}
$$

Notice that a key approximation in the FDT approach is that the average rates can be approximated by the rates at the average concentration [6, 9, 10]. Equations (45)-(46) can be solved explicitly, and we have

$$
s_{2}=\frac{k_{2}}{A B-k_{2}}, \quad r_{2}^{2}=\frac{\eta_{2}^{2}}{\eta_{0}^{2}}=\frac{k_{2}^{2} \tau_{0} \tau_{b}\left(1+A \tau_{0}+B \tau_{0} \tau_{b}\right)}{\left(A B-k_{2}\right)\left(A+B \tau_{b}\right)\left(1+A \tau_{0}+B \tau_{0} \tau_{b}+\tau_{0}^{2} \tau_{b}\left(A B-k_{2}\right)\right)} .
$$

When $\tau_{0} \tau_{b} \ll 1$, we have $A \tau_{0} \ll \tau_{0}^{2} \tau_{b} A B$, and thus,

$$
r_{2}^{2} \approx \frac{k_{2}^{2} \tau_{0} \tau_{b}}{\left(A B-k_{2}\right)\left(A+B \tau_{b}\right)}
$$

If further, $\tau_{b} \ll 1$,

$$
\begin{equation*}
r_{2}^{2} \approx \frac{k_{2}^{2} \tau_{0} \tau_{b}}{\left(A B-k_{2}\right) A} \approx \frac{\tau_{0} \tau_{b}}{\langle s\rangle\left(K_{a} k_{c}-1\right)\left(K_{a}+1\right) \frac{k_{c}}{k_{c}+1}} . \tag{47}
\end{equation*}
$$

Therefore, the amplification rate is a decreasing function in $K_{a}$ and $k_{c}$ (as $K_{a} k_{c}>1$ ). Moreover, the noise amplification rate not only depends on $\tau_{b}$, but also depends on $K_{a}$ and $k_{c}$, in a way that larger $K_{a}$ or $k_{c}$ leads to better noise attenuation. The proposed quantity $\left(t_{1 \rightarrow 0}-t_{0 \rightarrow 1}\right)$ incorporates both time scale of the $b$-equation $\left(\tau_{b}\right)$ and the key kinetic constants ( $k_{c}$ and $K_{a}$ ), and thus, serves as a better characterization than merely the time scale of the $b$-equation, $\tau_{b}$. In our application, the autocorrelation time $\tau_{0}$ is the inverse of the noise frequency, $\omega$. Therefore, (47) becomes

$$
r_{2}^{2} \approx \frac{\tau_{b}}{\omega\langle s\rangle\left(K_{a} k_{c}-1\right)\left(K_{a}+1\right) \frac{k_{c}}{k_{c}+1}},
$$

and the inverse relation between $r_{2}$ and $\omega$ arises naturally.

## 6 Models with Hill Functions

Hill functions are often used to model saturation effect and cooperativity in enzymatic reactions. In this section, we study the positive feedback module in Figure 1A based on Hill functions. First, we introduce Hill functions to the activation processes, as in the previous work of [2]:

$$
\begin{align*}
& \frac{d c}{d t}=k_{1}(a+b)(1-c)-k_{2} c+k_{3}, \\
& \frac{d a}{d t}=\left(k_{c} s \frac{c^{n}}{c^{n}+K^{n}}(1-a)-a+k_{4}\right) \tau_{a},  \tag{48}\\
& \frac{d b}{d t}=\left(k_{c} s \frac{c^{n}}{c^{n}+K^{n}}(1-b)-b+k_{4}\right) \tau_{b} .
\end{align*}
$$

When the Hill exponent is one, system (48) corresponds to Michaelis-Menten kinetics, whereas system (16) is derived following mass action kinetics. In [2], $n=3$ is used. Plotting the steady state response of the output as a function of constant input $s$ for system (16) and system (48) with $n=1$ and $n=3$ yields three close curves (Figure S5A). Note that when $n$ is large, for example, $n=10$, bistability could arise in system (48), which is expected $[1,3]$. Here, since we are solely interested in systems with a unique response to a constant signal, we focus on small $n$ 's, such as $n=3$.

The time evolution of the outputs of the above three systems in response to a noise-free signal (Figure S5B) and the noisy signal (Figures S5C-S5D) shows extensive overlapping, suggesting similar relationship of the noise attenuation and the deactivation and activation time scales between system (48) and system (16). This is confirmed by direct simulation (Figure S5E-S5G). However, the analytical study for the model using Hill equations becomes more difficult, in particular, for Hill exponents larger than two.

Next, we considered a variation of system (48),

$$
\begin{align*}
\frac{d c}{d t} & =k_{1}(a+b)(1-c)-k_{2} c+k_{3} \\
\frac{d a}{d t} & =\left(k_{c} s \frac{c^{n}}{c^{n}+K^{n}}(1-a)-\frac{a^{m}}{a^{m}+L^{m}}+k_{4}\right) \tau_{a}  \tag{49}\\
\frac{d b}{d t} & =\left(k_{c} s \frac{c^{n}}{c^{n}+K^{n}}(1-b)-\frac{b^{m}}{b^{m}+L^{m}}+k_{4}\right) \tau_{b}
\end{align*}
$$

with nonlinear deactivation terms in $A$ and $B$ loops. The simulations of (49) with $L=1$ and $m=1$ show similar outputs with system (48) (Figures S6A-S6B). Furthermore, the noise amplification rate displays the same trend as in all other models (Figures S6C-S6E).

## 7 A Polymyxin B Resistance Model in Enteric Bacteria

We use the same equations as in [8] to describe the feedforward connector loop (FCL) model:

$$
\begin{align*}
\frac{d x_{1}}{d t} & =k_{p b g P}\left(1-\frac{1}{\left(1+K_{2}[P h o P-P]^{2}\right)\left(1+K_{3}\left(x_{4}^{2}+x_{5}^{2}\right)\right)}\right)-k_{-p b g P} x_{1} \\
\frac{d x_{2}}{d t} & =k_{P m r D} \frac{K_{1}[P h o P-P]^{2}}{1+K_{1}[P h o P-P]^{2}}-k_{-P m r D} x_{2}-k_{c} x_{2} x_{4}+k_{-c} x_{5} \\
\frac{d x_{3}}{d t} & =k_{P m r A}-k_{P m r A} x_{3}+k_{-p} x_{4}+k_{p} x_{3}  \tag{50}\\
\frac{d x_{4}}{d t} & =k_{p} x_{3}-k_{-p} x_{4}+k_{-c} x_{5}-k_{c} x_{2} x_{4} \\
\frac{d x_{5}}{d t} & =k_{c} x_{2} x_{4}-k_{-c} x_{5} .
\end{align*}
$$

The parameter values used in simulations are listed in Table S1.

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