Fast Fourier Color Constancy

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Abstract

We present Fast Fourier Color Constancy (FFCC), a color constancy algorithm which solves illuminant estimation by reducing it to a spatial localization task on a torus. By operating in the frequency domain, FFCC produces lower error rates than the previous state-of-the-art by 13−20% while being 250−3000× faster. This unconventional approach introduces challenges regarding aliasing, directional statistics, and preconditioning, which we address. By producing a complete posterior distribution over illuminants instead of a single illuminant estimate, FFCC enables better training techniques, an effective temporal smoothing technique, and richer methods for error analysis. Our implementation of FFCC runs at ~700 frames per second on a mobile device, allowing it to be used as an accurate, real-time, temporally-coherent automatic white balance algorithm.

1. Intro

A fundamental problem in computer vision is that of estimating the underlying world that resulted in some observed image \[1,5\]. One subset of this problem is color constancy: estimating the color of the illuminant of the scene and the colors of the objects in the scene viewed under a white light. Despite its apparent simplicity, this problem has yielded a great deal of depth and challenge for both the human vision and computer vision communities \[17,22\]. Color constancy is also a practical concern in the camera industry: producing a natural looking photograph without user intervention requires that the illuminant be automatically estimated and discounted, a process referred to as “auto white balance” among practitioners. Though there is a profound historical connection between color constancy and consumer photography (exemplified by Edwin Land, the inventor of both Retinex theory \[25\] and the Polaroid instant camera), “color constancy” and “white balance” have come to mean different things — color constancy aims to recover the veridical world behind an image, while white balance aims to give an image a pleasant appearance consistent with some aesthetic or cultural norm. But with the current ubiquity of learning-based techniques in computer vision, both problems reduce to just estimating the “best” illuminant from an image, and the question of whether that illuminant is objectively true or subjectively attractive is just a matter of the data used during training.

Despite their accuracy, modern learning-based color constancy algorithms are not immediately suitable as practical white balance algorithms, as practical white balance has several requirements besides accuracy:

- **Speed** - An algorithm running in a camera’s viewfinder must run at 30 FPS on mobile hardware. But a camera’s compute budget is precious: demosaicing, face detection, auto exposure, etc, must also run simultaneously and in real time. Spending more than a small fraction (say, 5−10%) of a camera’s compute budget on white balance is impractical, suggesting that our speed requirement is closer to 1−5 milliseconds per frame.

- **Impoverished Input** - Most color constancy algorithms are designed for full resolution, high bit-depth input images, but operating on such large images is challenging and costly in practice. To be fast, the algorithm must work well on the small, low bit-depth “preview” images (32×24 or 64×48 pixels, 8-bit) which are usually computed by specialized camera hardware for this task.

- **Uncertainty** - In addition to the illuminant, the algorithm should produce some confidence measure or a complete posterior distribution over illuminants, thereby enabling convenient downstream integration with hand-engineered heuristics or external sources of information.

- **Temporal Coherence** - The algorithm should allow the estimated illuminant to be smoothed over time, to prevent color composition in videos from varying erratically.

In this paper we present a novel color constancy algorithm, which we call “Fast Fourier Color Constancy” (FFCC). Viewed as a color constancy algorithm, FFCC is 13−20% more accurate than the state of the art on standard benchmarks. Viewed as a prospective white balance algorithm, FFCC addresses our previously described requirements: Our technique is 250−3000× faster than the state of the art, and is capable of running at 1.44 milliseconds per frame on a standard consumer mobile platform using the thumbnail images already produced by that camera’s hard-
Figure 1: CCC [4] reduces color constancy to a 2D localization problem similar to object detection (1a). FFCC repeatedly wraps this 2D localization problem around a small torus (1b), which creates challenges but allows for faster illuminant estimation. See the text for details.

ware. FFCC produces a complete posterior distribution over illuminants which allows us to reason about uncertainty and enables simple and effective temporal smoothing.

We build on the “Convolutional Color Constancy” (CCC) approach of [4], which is currently one of the top-performing techniques on standard color constancy benchmarks [12, 20, 30]. CCC works by observing that applying a per-channel gain to a linear RGB image is equivalent to inducing a 2D translation of the log-chroma histogram of that image, which allows color constancy to be reduced to the task of localizing a signature in log-chroma histogram space. This reduction is at the core of the success of CCC and, by extension, our FFCC technique; see [4] for a thorough explanation. The primary difference between FFCC is that instead of performing an expensive localization on a large log-chroma plane, we perform a cheap localization on a small log-chroma torus.

At a high level, CCC reduces color constancy to object detection — in the computability theory sense of “reduce”. FFCC reduces color constancy to localization on a torus instead of a plane, and because this task has no intuitive analogue in computer vision we will attempt to provide one. Given a large image $A$ on which we would like to perform object detection, imagine constructing a smaller $n \times n$ image $B$ in which each pixel in $B$ is the sum of all values in $A$ separated by a multiple of $n$ pixels in either dimension:

$$B(i, j) = \sum_{k,l} A(i + nk, j + nl) \quad (1)$$

This amounts to taking $A$ and repeatedly wrapping it around a small torus (see Figure [1]). Detecting objects this way may yield a speedup as the image being processed is smaller, but it also raises new problems: 1) pixel values are corrupted with superimposed shapes that make detection difficult, 2) detections must “wrap” around the edges of this toroidal image, and 3) instead of an absolute, global location we can only recover an aliased, incomplete location. FFCC works by taking the large convolutional problem of CCC (ie, face detection on $A$) and aliasing that problem down to a smaller size where it can be solved efficiently (ie, face detection on $B$). We will show that we can learn an effective color constancy model in the face of the difficulty and ambiguity introduced by aliasing. This convolutional classifier will be implemented and learned using FFTs, because the naturally periodic nature of FFT convolutions resolves the problem of detections “wrapping” around the edge of toroidal images, and produces a significant speedup.

Our approach to color constancy introduces a number of issues. The aforementioned periodic ambiguity resulting from operating on a torus (which we dub “illuminant aliasing”) requires new techniques for recovering a global illuminant estimate from an aliased estimate (Section [5]). Localizing the centroid of the illuminant on a torus is difficult, requiring that we adopt and extend techniques from the directional statistics literature (Section [4]). But our approach presents a number of benefits. FFCC improves accuracy relative to CCC by $17 - 24\%$ while retaining its flexibility, and allows us to construct priors over illuminants (Section [5]). By learning in the frequency-domain we can construct a novel method for fast frequency-domain regularization and preconditioning, making FFCC training $20\times$ faster than CCC (Section [6]). Our model produces a complete unimodal posterior over illuminants as output, allowing us to construct a Kalman filter-like approach for processing videos instead of independent images (Section [7]).

### 2. Convolutional Color Constancy

Let us review the assumptions made in CCC and inherited by our model. Assume that we have a photometrically linear input image $I$ from a camera, with a black level of zero and with no saturated pixels. Each pixel $k$’s RGB value in image $I$ is assumed to be the product of that pixel’s “true” white-balanced RGB value $W^{(k)}$ and some global RGB illumination $L$ shared by all pixels:

$$\forall_k \begin{bmatrix} I^{(k)}_r \\ I^{(k)}_g \\ I^{(k)}_b \end{bmatrix} = \begin{bmatrix} W^{(k)}_r \\ W^{(k)}_g \\ W^{(k)}_b \end{bmatrix} \circ \begin{bmatrix} L_r \\ L_g \\ L_b \end{bmatrix} \quad (2)$$

The task of color constancy is to use the input image $I$ to estimate $L$, and with that produce $W^{(k)} = I^{(k)}/L$.

Given a pixel from our input RGB image $I^{(k)}$, CCC de-

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1 We cannot speak to the merit of this idea in the context of object detection, and we present it here solely to provide an intuition of our work on color constancy.

2 In practice, saturated pixels are identified and removed from all downstream computation, similarly to how color checker pixels are ignored.
fines two log-chroma measures:

\[ u(k) = \log \left( \frac{I_g^{(k)}}{I_r^{(k)}} \right) \quad v(k) = \log \left( \frac{I_b^{(k)}}{I_r^{(k)}} \right) \]  

The absolute scale of \( L \) is assumed to be unrecoverable, so estimating \( L \) simply requires estimating its log-chroma:

\[ L_u = \log \left( \frac{L_g}{L_r} \right) \quad L_v = \log \left( \frac{L_g}{L_b} \right) \]  

After recovering \((L_u, L_v)\), assuming that \( L \) has a magnitude of 1 lets us recover the RGB values of the illuminant:

\[ L_r = \frac{\exp(-L_u)}{z} \quad L_g = \frac{1}{z} \quad L_b = \frac{\exp(-L_v)}{z} \]

\[ z = \sqrt{\exp(-L_u)^2 + \exp(-L_v)^2 + 1} \]  

Framing color constancy in terms of predicting log-chroma has several small advantages over the standard RGB approach (2 unknowns instead of 3, better numerical stability, etc) but the primary advantage of this approach is that using log-chroma turns the multiplicative constraint relating \( W \) and \( I \) into an additive constraint \[15\], and this in turn enables a convolutional approach to color constancy. As shown in \[4\], color constancy can be framed as a 2D spatial localization task on a log-chroma histogram \( N \), where some sliding-window classifier is used to filter that histogram and the centroid of that filtered histogram is used as the log-chroma of the illuminant.

### 3. Illuminant Aliasing

We assume the same convolutional premise of CCC, but with one primary difference to improve quality and speed: we use FFTs to perform the convolution that filters the log-chroma histogram, and we use a small histogram to make that convolution as fast as possible. This change may seem trivial, but the periodic nature of FFT convolution combined

with the properties of natural images has a significant effect, as we will demonstrate.

Similarly to CCC, given an input image \( I \) we construct a histogram \( N \) from \( I \), where \( N(i, j) \) is the number of pixels in \( I \) whose log-chroma is near the \((u, v)\) coordinates corresponding to histogram position \((i, j)\):

\[ N(i, j) = \sum_k \left( \text{mod} \left( \frac{u(k) - u_{lo}}{h} - i, n \right) < 1 \right. \]

\[ \text{mod} \left( \frac{v(k) - v_{lo}}{h} - j, n \right) < 1 \]  

Where \( i, j \) are 0-indexed, \( n = 64 \) is the number of bins, \( h = 1/32 \) is the bin size, and \((u_{lo}, v_{lo})\) is the starting point of the histogram. Because our histogram is too small to contain the wide spread of colors present in most natural images, we use modular arithmetic to cause pixels to “wrap around” with respect to log-chroma (any other standard boundary condition would violate our convolutional assumption and would cause many image pixels to be ignored). This means that, unlike standard CCC, a single \((i, j)\) coordinate in the histogram no longer corresponds to an absolute \((u, v)\) color, but instead corresponds to an infinite family of \((u, v)\) colors. Accordingly, the centroid of a filtered histogram no longer corresponds to the color of the illuminant, but instead is an infinite set of illuminants. We will refer to this phenomenon as illuminant aliasing. Solving this problem requires that we use some technique to de-alias an aliased illuminant estimate\[1\] A high-level outline of

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\[1\]It is tempting to refer to resolving the illuminant aliasing problem as “anti-aliasing”, but anti-aliasing usually refers to preprocessing a signal to prevent aliasing during some resampling operation, which does not appear possible in our framework. “De-aliasing” suggests that we allow aliasing to happen to the input, but then remove the aliasing from the output.
our FFCC pipeline that illustrates illuminant (de-)aliasing can be seen in Fig. 2.

De-aliasing requires that we use some external information (or some external color constancy algorithm) to disambiguate between illuminants. An intuitive approach is to select the illuminant that causes the average image color to be as neutral as possible, which we call “gray world de-aliasing”. We compute average log-chroma values $(\bar{\hat{u}}, \bar{\hat{v}})$ for the entire image and use this to turn an aliased illuminant estimate $(\hat{L}_u, \hat{L}_v)$ into a de-aliased illuminant $(\hat{L}'_u, \hat{L}'_v)$:

$$\bar{\hat{u}} = \log \left( \sum_k u^{(k)} \right) \quad \bar{\hat{v}} = \log \left( \sum_k v^{(k)} \right)$$

$$\begin{bmatrix} \hat{L}'_u \\ \hat{L}'_v \end{bmatrix} = \begin{bmatrix} \hat{L}_u \\ \hat{L}_v \end{bmatrix} \left( \frac{1}{nh} \begin{bmatrix} \hat{L}_u - \bar{\hat{u}} \\ \hat{L}_v - \bar{\hat{v}} \end{bmatrix} + \frac{1}{2} \right)$$

Another approach, which we call “gray light de-aliasing”, is to assume that the illuminant is as close to the center of the histogram as possible. This de-aliasing approach simply requires carefully setting the starting point of the histogram $(u_{lo}, v_{lo})$ such that the true illuminants in natural scenes all lie within the span of the histogram, and setting $\hat{L}' = \hat{L}$. We do this by setting $u_{lo}$ and $v_{lo}$ to maximize the distance between the edges of the histogram and the bounding box surrounding the ground-truth illuminants in the training data. Gray light de-aliasing is trivial to implement but, unlike gray world de-aliasing, it will systematically fail if the histogram is too small to fit all illuminants within its span.

To summarize the difference between CCC [4] and our approach with regards to illuminant aliasing, CCC (approximately) performs illuminant estimation as follows:

$$\begin{bmatrix} \hat{L}_u \\ \hat{L}_v \end{bmatrix} = \begin{bmatrix} u_{lo} \\ v_{lo} \end{bmatrix} + h \left( \arg \max_{i,j} (Nh) \right)$$

Where $Nh$ is performed using a pyramid convolution. FFCC corresponds to this procedure:

$$P \leftarrow \text{softmax}(Nh)$$

$$(\boldsymbol{\mu}, \Sigma) \leftarrow \text{fit_bvm}(P)$$

$$\begin{bmatrix} \hat{L}_u \\ \hat{L}_v \end{bmatrix} \leftarrow \text{de_alias}(\boldsymbol{\mu})$$

Where $Nh$ is a small and aliased toroidal histogram, convolution is performed with FFTs, and the centroid of the filtered histogram is estimated and de-aliased as necessary. By constructing this pipeline to be differentiable we can train our model in an end-to-end fashion by propagating the gradients of some loss computed on the de-aliased illuminant prediction $\hat{L}$ back onto the learned filters $F$. The centroid fitting in Eq. 11 is performed by fitting a bivariate von Mises distribution to a PDF, which we will now explain.

### 4. Differentiable Bivariate von Mises

Our architecture requires some mechanism for reducing a toroidal PDF $P(i, j)$ to a single estimate of the illuminant. Localizing the center of mass of a histogram defined on a torus is difficult: fitting a bivariate Gaussian may fail when the input distribution “wraps around” the sides of the PDF, as shown in Fig. 3. Additionally, for the sake of temporal smoothing (Section 7) and confidence estimation, we want our model to predict a well-calibrated covariance matrix around the center of mass of $P$. This requires that our model be trained end-to-end, which therefore requires that our mean/covariance fitting be analytically differentiable and therefore usable as a “layer” in our learning architecture. To address these problems we present a variant of the bivariate von Mises distribution [27], which we will use to efficiently localize the mean and covariance of $P$ in a manner that allows for easy backpropagation.

The bivariate von Mises distribution (BVM) is a common parameterization of a PDF on a torus. There exist several parametrizations which mostly differ in how “concentration” is represented (“concentration” having a similar meaning to covariance). All of these parametrizations present problems in our use case: none have closed form expressions for maximum likelihood estimators [24], none lend themselves to convenient backpropagation, and all define concentration in terms of angles and therefore require “conversion” to covariance matrices during color de-aliasing. For these reasons we present an alternative parametrization in which we directly estimate a BVM as a mean $\boldsymbol{\mu}$ and covariance $\Sigma$ in a simple and differentiable closed form expression. Though necessarily approximate, our estimator is accurate when the distribution is well-concentrated, which is generally the case for our task.

Our input is a PDF $P(i, j)$ of size $n \times n$, where $i$ and $j$ are integers in $[0, n - 1]$. For convenience we define a mapping from $i$ or $j$ to angles in $[0, 2\pi)$ and the marginal distributions of $P$ with respect to $i$ and $j$:

$$\theta(i) = \frac{2\pi i}{n} \quad P_i(i) = \sum_j P(i, j) \quad P_j(j) = \sum_i P(i, j)$$

We also define the marginal expectation of the sine and cosine of the angle:

$$y_i = \sum_i P_i(i) \sin(\theta(i)) \quad x_i = \sum_i P_i(i) \cos(\theta(i))$$

With $x_j$ and $y_j$ defined similarly.
Our estimated covariance matrix is simply the sample co-

Figure 3: We fit a bivariate von Mises distribution (shown in solid blue) to toroidal PDFs $P(i, j)$ to produce an aliased illuminant estimate. Contrast this with fitting a bivariate Gaussian (shown in dashed red) which treats the PDF as if it lies on a plane. Both approaches behave similarly if the distribution lies near the center of the unwrapped plane (left) but fitting a Gaussian fails as the distribution begins to “wrap around” the edge (middle, right).

Estimating the mean $\mu$ of a BVM from a histogram just requires computing the circular mean in $i$ and $j$:

$$
\mu = \left[ \frac{u_{lo}}{v_{lo}} \right] + h \left[ \frac{\text{mod} \left( \frac{2\pi}{n} \text{atan2}(y_i, x_i), n \right)}{\text{mod} \left( \frac{2\pi}{n} \text{atan2}(y_j, x_j), n \right)} \right]
$$

(14)

Eq. (14) includes gray light de-aliasing, though gray world de-aliasing can also be applied to $\mu$ after fitting.

We can fit the covariance of our model by simply “un-

wrapping” the coordinates of the histogram relative to the estimated mean and treating these unwrapped coordinates as though we are fitting a bivariate Gaussian. We define the “unwrapped” $(i, j)$ coordinates such that the “wrap around” point on the torus lies as far away from the mean as possible, or equivalently, such that the unwrapped coordinates are as close to the mean as possible:

$$
\tilde{i} = \text{mod} \left( i - \left[ \frac{\mu_u - u_{lo}}{h} \right] + \frac{n}{2}, n \right)
$$

$$
\tilde{j} = \text{mod} \left( j - \left[ \frac{\mu_v - v_{lo}}{h} \right] + \frac{n}{2}, n \right)
$$

(15)

Our estimated covariance matrix is simply the sample co-

variance of $P(\tilde{i}, \tilde{j})$:

$$
\Sigma = h^2 \left[ \epsilon + \sum_i P(i, \tilde{i})^2 - E[\tilde{i}]^2 \sum_{i,j} P(i, j) \bar{j} - E[\tilde{j}] E[\tilde{j}] \right] \left[ \epsilon + \sum_j P(j, \tilde{j})^2 - E[\tilde{j}]^2 \right]
$$

(17)

We regularize the sample covariance matrix slightly by adding a constant $\epsilon = 1$ to the diagonal.

With our estimated mean and covariance we can com-

pute our loss: the negative log-likelihood of a Gaussian (ignoring scale factors and constants) relative to the true illuminant $L^*$:

$$
f(\mu, \Sigma) = \log |\Sigma| + \left( \begin{bmatrix} L_u^* \\ L_v^* \end{bmatrix} - \mu \right)^T \Sigma^{-1} \left( \begin{bmatrix} L_u^* \\ L_v^* \end{bmatrix} - \mu \right)
$$

(18)

Using this loss causes our model to produce a well-

calibrated complete posterior of the illuminant instead of just a single estimate. This posterior will be useful when processing video sequences (Section 7) and also allows us to attach confidence estimates to our predictions using the entropy of $\Sigma$ (see the supplement).

Our entire system is trained end-to-end, which requires that every step in BVM fitting and loss computation be analytically differentiable. See the supplement for the analytical gradients for Eqs. (14), (17), and (18) which can be chained together to backpropagate the gradient of $f(\cdot)$ onto the input PDF $P$.

5. Model Extensions

The system we have described thus far (compute a peri-

odic histogram of each pixel’s log-chroma, apply a learned FFT convolution, apply a softmax, fit a de-aliased bivariate von Mises distribution) works reasonably well (Model A in Table 1) but does not produce state-of-the-art results. This is likely because this model reasons about pixels independently, ignores all spatial information in the image, and does not consider the absolute color of the illuminant. Here we present extensions to the model which address these issues and improve accuracy accordingly.

As explored in [4], a CCC-like model can be generalized to a set of “augmented” images provided that these images are non-negative and “scale with intensity” [14]. This lets us apply certain filtering operations to image $I$ and, instead of constructing a single histogram from our image, construct a “stack” of histograms constructed from the image and its filtered versions. Instead of learning and applying one filter, we learn a stack of filters and sum across channels after convolution. The general family of augmented images used in [4] are expensive to compute, so we instead use just the input image $I$ and a local measure of absolute deviation in the input image:

$$
E(x, y, c) = \frac{1}{8} \sum_{i=-1}^{1} \sum_{j=-1}^{1} |I(x, y, c) - I(x+i, y+j, c)|
$$

(19)

These two features appears to perform similarly to the four features used in [4], while being cheaper to compute.

Just as a sliding-window object detector is often invariant to the absolute location of an object in an image, the convolutional nature of our baseline model makes it invariant to any global shift of the color of the input image. This means that our baseline model cannot rely on any statistical regularities of the illumination by, say, modeling black body radiation, the specific properties of commonly manufactured light bulbs, or any varying spectral sensitivity across cameras. Though CCC does not model illumination directly, it appears to indirectly reason about illumination by using the boundary conditions of its pyramid convolution to learn
a model which is not truly spatially varying and is therefore sensitive to absolute color. Because a torus has no boundaries, our model is invariant to global input color, so we must therefore introduce a mechanism for directly reasoning about illuminants. We use a per-illuminant “gain” map $G(i, j)$ and “bias” map $B(i, j)$, which together apply a per-illuminant affine transformation to the output of our previously-described convolution at (aliased) color $(i, j)$. The bias $B$ causes our model to prefer certain illuminants over others, while the gain $G$ causes the contribution of the convolution at certain colors to be amplified.

Our two extensions (an augmented edge channel and an illuminant gain/bias map) let us redefine the $P$ in Eq. 10 as

$$P = \text{softmax} \left( B + G \circ \sum_k (N_k \ast F_k) \right)$$

(20)

Where $\{F_k\}$ are the set of learned filters for each augmented channel’s histogram $N_k$, $G$ is our learned gain map, and $B$ is our learned bias map. In practice we actually parametrize $G_{\log}$ when training and define $G = \exp(G_{\log})$, which constraints $G$ to be non-negative. Visualizations of $G$ and $B$ and our learned filters can be seen in Fig. 4.

6. Fourier Regularization and Preconditioning

Our learned model weights ($\{F_k\}, G, B$) are all periodic $n \times n$ images. To improve generalization, we want these weights to be small and smooth. In this section we present the general form of the regularization used during training, and we show how this regularization lets us precondition the optimization problem solved during training to find lower-cost minima in fewer iterations. Because this frequency-domain optimization technique applies generally to any optimization problem concerning smooth and periodic images, we will describe it in general terms.

Let us construct an optimization problem with respect to a single $n \times n$ image $Z$ consisting of a data term $f(Z)$ and a regularization term $g(Z)$:

$$Z^* = \arg\min_Z \left( f(Z) + g(Z) \right)$$

(21)

We require that the regularization $g(Z)$ is the weighted sum of squared periodic convolutions of $Z$ with some filter bank. In our experiments $g(Z)$ is the weighted sum of the squared difference between adjacent values (similar to a total variation loss [29]) and the sum of squared values:

$$g(Z) = \lambda_1 \sum_{i,j} \left( (Z(i, j) - Z(\text{mod}(i + 1, n), j))^2 + (Z(i, j) - Z(i, \text{mod}(j + 1, n)))^2 \right) + \lambda_0 \sum_{i,j} Z(i, j)^2$$

(22)

Where $\lambda_1$ and $\lambda_0$ are hyperparameters that determine the strength of each smoothness term. We require that $\lambda_0 > 0$ to prevent divide-by-zero issues during preconditioning.

We use a variant of the standard FFT $F_\nu(\cdot)$ which bijectively maps from some real $n \times n$ image to a real $n^2$-dimensional vector, instead of the complex $n \times n$ image produced by a standard FFT (See the supplement for a formal description). With this, we can rewrite Eq. 22 as follows:

$$w = \frac{1}{n} \sqrt{\lambda_1 \left( |F_\nu([1, -1])|^2 + |F_\nu([1; -1])|^2 \right) + \lambda_0}$$

$$g(Z) = F_\nu(Z)^T \text{diag}(w)^2 F_\nu(Z)$$

(23)

where the vector $w$ is just some fixed function of the definition of $g(Z)$ and the values of the hyperparameters $\lambda_1$ and $\lambda_0$. The 2-tap difference filters in $F_\nu([1, -1])$ and $F_\nu([1; -1])$ are padded to size $(n \times n)$ before the FFT. With $w$ we can define a mapping between our 2D image space and a rescaled FFT vector space:

$$z = w \circ F_\nu(Z)$$

(24)

Where $\circ$ is an element-wise product. This mapping lets us rewrite the optimization problem in Eq. 21 as:

$$Z^* = F_\nu^{-1} \left( \frac{1}{w} \left( \arg\min_z \left( f \left( F_\nu^{-1} \left( \frac{Z}{w} \right) \right) + \|z\|^2 \right) \right) \right)$$

(25)

where $F_\nu^{-1}(\cdot)$ is the inverse of $F_\nu(\cdot)$, and division is element-wise. This reparametrization reduces the complicated regularization of $Z$ to a simple L2 regularization of $z$, which has a preconditioning effect.

We use this technique during training to reparameterize all model components ($\{F_k\}, G, B$) as rescaled FFT vectors, each with their own values for $\lambda_0$ and $\lambda_1$. The effect of this can be seen in Fig. 5, where we show the loss during our two training stages. We compare against naive time-domain optimization (Eq. 21) and non-preconditioned frequency-domain optimization (Eq. 25 with $w = 1$). Our preconditioned reformulation exhibits a significant speedup and finds minima with lower losses.
For all experiments (excluding our “deep” variants, see the supplement), training is as follows: All model parameters are initialized to 0, then we have a convex pre-training step which optimizes Eq. 22 where \( f(\cdot) \) is a logistic loss (described in the supplement) using LBFGS for 16 iterations, and then we optimize Eq. 25 where \( f(\cdot) \) is the non-convex BVM loss in Eq. 18 using LBFGS for 64 iterations.

### 7. Temporal Smoothing

Color constancy is usually studied in the context of individual images, which are assumed to be IID. But a practical white balance algorithm must run on a video sequence, and must enforce some temporal smoothing of the predicted illuminant to avoid presenting the viewer with an erratically-varying image in the viewfinder. This smoothing cannot be too aggressive or else the viewfinder may appear unresponsive when the illumination changes rapidly (a colorful light turning on, the camera quickly moving outdoors, etc). Additionally, when faced with multiple valid hypotheses (a blue wall under white light vs a white wall under blue light, etc) we may want to use earlier images to resolve ambiguities. These desiderata of stability, responsiveness, and robustness are at odds with each other, and so some compromise must be struck.

Our task of constructing a temporally coherent illuminant estimate is aided by the probabilistic nature of the output of our per-frame model, which produces a posterior distribution over illuminants parameterized as a bivariate Gaussian. Let us assume that we have some ongoing estimate of the illuminant and its covariance \((\mu_t, \Sigma_t)\). Given the observed mean and covariance \((\mu_o, \Sigma_o)\) provided by our model we update our ongoing estimate by first convolving it with an zero-mean isotropic Gaussian (encoding our prior belief that the illuminant may change over time) and then multiplying that “fuzzed” Gaussian by the observed Gaussian:

\[
\Sigma_{t+1} = \left( \left( \Sigma_t + \begin{bmatrix} \alpha & 0 \\ 0 & \alpha \end{bmatrix} \right) \right)^{-1} \Sigma_o + \Sigma_t \\
\mu_{t+1} = \left( \left( \Sigma_t + \begin{bmatrix} \alpha & 0 \\ 0 & \alpha \end{bmatrix} \right) \right)^{-1} \mu_t + \Sigma_o \mu_o 
\]

Where \( \alpha \) is a parameter that defines the expected variance of the illuminant over time. This update resembles a Kalman filter but with a simplified transition model, no control model, and variable observation noise.

This temporal smoothing is not used in our benchmarks, but its effect can be seen in the supplemental video.

### 8. Results

We evaluate our technique using two standard color constancy datasets: the Gehler-Shi dataset [20, 30] and the Cheng et al. dataset [12] (see Tables 1 and 2). For the Gehler-Shi dataset we present several ablations and variants of our model to show the effect of each design decision and to investigate trade-offs between speed and accuracy. Models labeled “full” were run on 384 × 256 16-bit images, while models labeled “thumb” were run on 48 × 32 8-bit images, which are the kind of images that a practical white-balance system embedded on a hardware device might use. Models labeled “4 channel” use the four feature channels used in [4], while models labeled “2 channel” use the two channels we present in Section 5. We also present models in which we only use the “pixel channel” \( I \) or the “edge channel” \( E \) as input. All models have a histogram size of \( n = 64 \) except for Models K and L where \( n \) is varied to show the impact of illuminant aliasing. Two models use “gray world” de-aliasing, and the rest use “gray light” de-aliasing. The former seems slightly less effective than the latter unless chroma histograms are heavily aliased, which is why we use it in Model K. Model C only has one training stage that minimizes logistic loss for 64 iterations, thereby removing the BVM fitting from training. Model E fixes \( G(i, j) = 1 \) and \( B(i, j) = 0 \), thereby removing the model’s ability to reason about the absolute color of the illuminant. Model B was trained only to minimize the data term (ie, \( \lambda_0 = \lambda_1 = 0 \) in Eq. 22) while Model D uses L2 regularization but not total variation (ie, \( \lambda_3 = 0 \) in Eq. 22).

Models N, O and P are variants of Model J in which, instead of learning a fixed model \( \{ F_k, G, B \} \) we express those model parameters as the output of a small 2-layer neural network. As inputs to this network we use image metadata, which allows the model to reason about exposure time and camera sensor type, and/or a CNN-produced feature vector.
which allows the model to reason about semantics (see the supplement for details). For each experiment we tune all λ hyperparameters to minimize the “average” error during cross-validation, using cyclic coordinate descent.

Model P achieves the lowest-error results, with a 20% reduction in error on Gehler-Shi compared to the previously best-performing published technique. This improvement in accuracy also comes with a significant speedup compared to previous techniques: ~30 ms/image for most models, compared to the 520 ms of CCC [4] or the 3 seconds (on a GPU) of Shi et al. [31]. Model Q (our fastest model) has an accuracy comparable to [4] and [31] but takes only 1.1 milliseconds to process an image, making it hundreds or millions of times faster than the current-state-of-the-art. Additionally, our model appears to be faster to train than the state-of-the-art, though training times for prior work are often not available. All runtimes in Table 1 for our model were computed on an Intel Xeon CPU E5-2680. Runtimes for the “full” model were produced using a Matlab implementation, while runtimes for the “thick” model were produced using a Halide [23] CPU implementation (our Matlab implementation of Model Q takes 2.37 ms/image). Runtimes for our “+semantic” models are not presented as we were unable to profile [34] accurately (CNN feature computation appears to dominate runtime).

To demonstrate that our model is a viable automatic white balance system for consumer photography, we ran our Halide code on a 2016 Google Pixel XL using the thumbnail images computed by the device’s camera stack. This implementation ran at 1.44ms per image, which is equivalent to 30 frames per second using < 5% of the total compute budget, thereby satisfying our previously-stated speed requirements. A video of our system running in real-time on a phone can be found in the supplement.

9. Conclusion

We have presented FFCC, a color constancy algorithm that produces a 13 – 20% reduction in error and a 250 – 3000× speedup relative to prior work. In doing so we have introduced the concept of convolutional color constancy on a torus, and we have introduced techniques for illuminant de-aliasing and differentiable bivariate von Mises fitting required for this toroidal approach. We have also presented a novel technique for fast Fourier-domain optimization subject to a certain family of regularizers. FFCC produces a complete posterior distribution over illuminants, which lets us assess the model’s confidence and also enables a Kalman filter-like temporal smoothing model. FFCC’s speed, accuracy, and temporal consistency allows it to be used for real-time white balance on a consumer camera.

Table 1: Performance on the Gehler-Shi dataset [20, 30]. We present five error metrics and their average (the geometric mean) with the lowest error metric highlighted in yellow. We present the time (in seconds) for training each model and for evaluating a single image, when available.

Table 2: Performance on the dataset from Cheng et al. [12], in the same format as Table 1 excluding runtimes. As was done in [4] we present the average performance (the geometric mean) over all 8 cameras in the dataset.
References


